Régine Douady Adrien Douady

## Algebra and Galois Theories

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# Algebra and Galois Theories 

Springer

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ISBN 978-3-030-32795-8
ISBN 978-3-030-32796-5 (eBook)
https://doi.org/10.1007/978-3-030-32796-5
Mathematics Subject Classification (2010): 15-01, 18-01, 55-01, 12F10, 14E20, 14F35, 14H, 14H30

Translation from the French language edition: Algèbre et théories galoisiennes by Régine and Adrien Douady, © Published by Cassini 2004. All Rights Reserved.
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## Introduction

Similarities between Galois theory and the theory of covering spaces are so striking that algebraists use geometric language to talk of field extensions whereas topologists talk of Galois covers. Here we have tried to develop these theories in parallel, beginning with that of coverings. The reader will thereby be better able to visualize. This similarity can sometimes be found in specific formulations:
(4.5.5) Proof of the proposition. Let $\varphi$ : $\mathrm{S}(\mathrm{X}) \rightarrow \mathrm{S}(\mathrm{Y})$ be a morphism in G -Sets.
There is a morphism $\varphi_{*}: \mathrm{E} \times \mathrm{S}(\mathrm{X}) \rightarrow \mathrm{E} \times$ $\mathrm{S}(\mathrm{Y})$ defined by $\varphi_{*}(t, s)=(t, \varphi(s))$ corresponding to the morphism $\varphi$. The morphism $\varphi_{*}$ is compatible with the G-operation $\perp \ldots$
(5.7.4) Proof of the proposition. Let $\varphi$ : $\mathrm{S}(\mathrm{B}) \rightarrow \mathrm{S}(\mathrm{A})$ be a morphism in G -Sets. There is a morphism of L -algebras $\varphi^{*}$ : $\mathrm{L}^{\mathrm{S}(\mathrm{A})} \rightarrow \mathrm{L}^{\mathrm{S}(\mathrm{B})}$ defined by $\varphi^{*}(h)=h \circ \varphi$ corresponding to $\varphi$. The homomorphism $\varphi^{*}$ is compatible with the G-operation $\perp \ldots$

The presentation chosen highlights this resemblance, the intention being to transfer to Galois theory any eventual geometric intuition gained in the framework of coverings. This presentation is unlikely to be the best suited for each individual context. The reader will recognize the influence of Grothendieck... and of Bourbaki. The similarity is not a perfect correspondence: in the algebraic context, profinite groups occur whereas discrete groups do in coverings. In the context of Galois theory, there is nothing comparable to points of a space, nor to loops; the choice of an algebraic closure of a field can only very indirectly be considered as similar to the choice of a base point. In the study of Riemann surfaces in Chapter 6, this similarity takes the shape of an equivalence of categories. This enables us in Section 6.4 to obtain through transcendental means a purely algebraic result without any known purely algebraic proof (Proposition 6.5.6). The end of that chapter deals with the study of automorphisms of Riemann surfaces, and some Fuchsian groups are considered. In each of the Chapters 4, 5, 6, the central theorem is expressed as an equivalence of categories (4.5.3, 4.6.8, 5.7.3, 6.2.4). We have therefore included a preliminary chapter on categories: Chapter 2. To avoid difficulties stemming from the impossibility of considering the set of all sets, our approach is that of universes; so we only consider "small categories", namely, categories whose underlying set is indeed a set. This in no way impacts on our goal
as seen in "lemmas about cardinalities" (4.4.7, 5.2.5). Infinite Galois theory forces us to consider sets and profinite groups. Tychonoff's theorem (1.7.7), based on the axiom of choice and Zorn's theorem, is essential in the theory of profinite groups. These theorems are proved in an initial chapter (Chapter 1). In 1.5.3, we mention the places where Zorn's theorem occurs in mathematics. Excepting Sections 2.5.14, 2.5.15, and 2.7.7 and paragraphs 2.8 and 2.9 (profinite spaces and groups), written in view of infinite Galois theory, Chapter 2 is independent from Chapter 1. We have inserted a chapter on rings, modules, and tensor products (Chapter 3). This chapter is broader than strictly needed for the development of Galois theory.

In this new edition, we have added a chapter on Grothendieck's theory of "dessins d'enfants". Through topology (theory of ramified coverings), it reduces problems about the classification of number fields to combinatorics questions.

We would like to thank the readers of the manuscript for their comments and for the improvements they suggested, in particular, Antoine Chambert-Loir and Christian Houzel.

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## About the Authors

Régine Douady was born in 1934. After a Ph.D. in the didactics of mathematics, she became a Lecturer at the University Paris Denis-Diderot, as well as the head of the IREM (Institut de Recherche sur l' Enseignement des Mathematiques). She was made a chevalier dans l'ordre des palmes academiques. Now retired, she has ceaselessly endeavored to introduce into teaching the necessity to address concepts from different standpoints-the motivating idea behind this book.

Adrien Douady was born in 1935. After receiving his Ph.D. in mathematics in the field of complex analytic geometry, he later joined the University Paris-Sud (Orsay). A recipient of the Ampere prize, he was a member of the French Academy of Sciences, as well as of the highly influential informal Bourbaki group. Throughout his life, he remained interested in several areas. Yet his odyssey always brought him back to complex numbers. He passed away in 2006 and is notably survived by his wife, Regine Douady.

## Chapter 1 <br> Zorn's Lemma

## Introduction

The purpose we have in mind is infinite Galois theory. The notion of inverse limits of finite groups will be needed for this. By Tychonoff's theorem, these are compact groups. The example in 6.4.6 shows that countable products are not sufficient. Tychonoff's theorem is required in full, and hence so is Zorn's lemma, and thus the axiom of choice.

This chapter introduces the axiom of choice as well as the two results enabling its application: Zorn's Lemma and Zermelo's theorem. Both play more or less the same role, and either can usually be chosen. In our opinion, Zorn's lemma is more powerful, while Zermelo's theorem sometimes casts a more significant light.

Well-ordered sets form the framework for Zermelo's theorem. After outlining their theory, we proceed with the proof of Zorn's lemma. The quickest way to prove Tychonoff's theorem uses ultrafilters, and this is the method we follow.

### 1.1 Choice Functions

### 1.1.1

Definition Let E be a set. Set $\mathfrak{P}^{*}(\mathrm{E})$ be the set of its non-empty subsets. The choice function over E is a map $\tau: \mathfrak{P}^{*}(\mathrm{E}) \rightarrow \mathrm{E}$ such that $\forall \mathrm{X} \in \mathfrak{P}^{*}(\mathrm{E}), \tau(\mathrm{X}) \in \mathrm{X}$.

### 1.1.2

AXIOM OF CHOICE. For any set E, there exists a choice function defined on E .

### 1.1.3 Commentary

The axiom of choice is now accepted by most mathematicians. This was not always the case: in the early XX-th century, it remained highly controversial. It is only slowly that it was realized that all the axiom does is give a weaker meaning to the expression "there exists" (and to the quantifier $\exists$ ), at least weaker than "it can be explicitly constructed".

As Zorn's lemma is equivalent to it, the axiom of choice is a powerful tool which we do not wish to forgo. But given the monsters whose existence it ensures (see for example the Banach-Tarski paradox 1.1.5), the emergence of mathematical contradictions could be feared. Only in 1938 did Gödel prove that there is no such danger: he showed that if (ZF) set theory together with the axiom of choice is inconsistent, then it remains inconsistent without the axiom of choice. For example, let $x$ and $y$ be two real numbers explicitly computable with arbitrary precision. If a large inequality $x \leqslant y$ can be proved with the axiom of choice, it must then be true. Indeed, if their computation showed that $x>y$, this would contradict the axiom of choice. Gödel's theorem would then ensure that a contradiction can be obtained without the axiom of choice.

In 1963, Cohen showed, in the event that the axiom of choice is discarded, the axiom stating that there is no choice function over $\mathbb{R}$ can be added to set theory without making it inconsistent unless it already is so. This addition amounts to giving the expression "there exists" a strong meaning. Cohen's result tells us that an explicit algorithm giving a choice function over $\mathbb{R}$ would result in inconsistent mathematics. Even the stronger statement (see 1.1.4) that every subset of $\mathbb{R}$ is Lebesgue measurable can be accepted.

### 1.1.4 Example 1: Existence of Non-measurable Sets

A non-measurable set can be constructed by means of a choice function $\tau$ over $\mathbb{T}=\mathbb{R} / \mathbb{Z}:$

Let $\alpha \in \mathbb{T}$ be an irrational element. Define $\varpi: \mathbb{T} \rightarrow \mathbb{T}$ by $\varpi(x)=$ $\tau\left(\{x+n \alpha\}_{n \in \mathbb{Z}}\right)$. Set X to be the image of $\varpi$ and $\mathrm{R}_{\alpha}$ the map $t \mapsto t+\alpha$. Then $\mathbb{T}=\bigcup_{n \in \mathbb{Z}} \mathrm{R}_{\alpha}{ }^{n}(\mathrm{X})$, and these sets are disjoint. Hence X is not measurable for the Lebesgue measure $\mu$. Indeed, $\mu(\mathbf{X})=0$ would imply $\mu(\mathbb{T})=0$, and $\mu(\mathbf{X})>0$ would imply $\mu(\mathbb{T})=\infty$.

### 1.1.5 Example 2: The Banach-Tarski Paradox

Theorem (Banach-Tarski) There exist a finite partition $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ of the unit sphere $\mathrm{S}^{2} \subset \mathbb{R}^{3}$, isometries $\left(\mathrm{F}_{i}\right)_{i \in \mathrm{I}}$ of $\mathrm{S}^{2}$ and a partition of the index set I into two subsets $\mathrm{I}^{\prime}$ and $\mathrm{I}^{\prime \prime}$, such that $\mathrm{F}_{i}\left(\mathrm{~A}_{i}\right)$ for $i \in \mathrm{I}^{\prime}$ form a partition of $\mathrm{S}^{2}$, as well as for $\mathrm{I}^{\prime \prime}$. The same holds for unit balls.

See 1.1, Exercise 5, for a proof.
Commentary. This shows that there can be no "finitely additive measure", preserved by rotations, on every subset of $\mathrm{S}^{2}$, with respect to which the latter has measure 1 . Indeed, the union $\mathrm{X}^{\prime}$ of $\mathrm{A}_{i}$ for $i \in \mathrm{I}^{\prime}$ should have measure 1 , similarly for $\mathrm{I}^{\prime \prime}$, which would yield measure 2 for $S^{2}$. By contrast, there are such "measures" on the circle (1.1, Exercise 4), and contradictions arise only if countable additivity is insisted upon. It is now known that for a coherent theory of integration, countable additivity is of prime importance; this was not so obvious at the time (1926).

### 1.1.6

Proposition Every surjection $f: \mathrm{X} \rightarrow \mathrm{Y}$ between two sets X and Y has a section, i.e. a map $\sigma: \mathrm{Y} \rightarrow \mathrm{X}$ such that $f \circ \sigma=\mathrm{I}_{\mathrm{Y}}$.

Proof Let $\tau$ be a choice function on X. The map defined by $\sigma(y)=\tau\left(f^{-1}(y)\right)$ is a right inverse.

Remark The axiom of choice also follows from this proposition. Indeed, let E be a set, and $\Gamma=\left\{(x, \mathrm{X}) \in \mathrm{E} \times \mathfrak{P}^{*}(\mathrm{E}) \mid x \in \mathrm{X}\right\}$. The second projection $\pi: \Gamma \rightarrow \mathfrak{P}^{*}(\mathrm{E})$ is surjective. By the proposition, it has a section of the form $\mathrm{X} \mapsto(\tau(\mathrm{X}), \mathrm{X})$, and the function $\tau$ defined thereby is a choice function. This proposition is therefore equivalent to the axiom of choice.

## Exercises 1.1. (Choice Functions)

1.-(a) Construct a choice function on $\mathbb{Q}$, and on $\mathbb{Q}^{2}$.
(b) Let X and Y be two sets, each equipped with a choice function. Construct a choice function on $\mathrm{X} \times \mathrm{Y}$.
2.-Construct a choice function on the set of non-empty open subsets of $\mathbb{R}^{2}$ (i.e, a function $\tau$ which associates to all non-empty open subsets $U$ of $\mathbb{R}^{2}$ a point of U ). Do the same for closed subsets. Is it possible to repeat this for the locally closed subsets? for the Borel subsets? (see 1.1.4).
3.-Let $E$ denote the set of $\mathbb{R}$-valued continuous functions on $[0,1]$, and $F$ the quotient of $E$ by the following equivalence relation:

$$
f \sim g \Longleftrightarrow f=g \text { or } f=-g .
$$

Define a section $\sigma$ of the canonical quotient map $\chi: \mathrm{E} \rightarrow \mathrm{F}$, i.e. a map $\sigma: \mathrm{F} \rightarrow \mathrm{E}$ such that $\chi \circ \sigma=1_{\mathrm{F}}$.

Can this be repeated for the set of all functions $[0,1] \rightarrow \mathbb{R}$ ?
4. (Finitely additive measure on $\mathbb{T}$.) -Set $\mathbb{T}=\mathbb{R} / \mathbb{Z}$. For $\alpha \in \mathbb{T}$, let $\mathrm{R}_{\alpha}$ denote the map $x \mapsto x+\alpha$, and $\mathrm{R}_{\alpha}^{*}$ the map $f \mapsto f \circ \mathrm{R}_{\alpha}$ from the space E of all bounded functions $\mathbb{T} \rightarrow \mathbb{R}$ into itself.
(a) Show that the constant function 1 cannot be written as follows:

$$
\begin{equation*}
\left(\sum_{1}^{n} f_{i}-\mathrm{R}_{\alpha_{i}}^{*} f_{i}\right)-g \text { with } f_{i} \in \mathrm{E} \text { and } \mathrm{g} \geqslant 0 ; \tag{*}
\end{equation*}
$$

(for such an $h$, there is an inequality of the form:

$$
\left.\sum_{0 \leqslant k_{i} \leqslant \mathrm{~N}} h\left(\sum k_{i} \alpha_{i}\right) \leqslant \mathrm{C} \cdot \mathrm{~N}^{n-1}\right) .
$$

(b) Show that there is a linear form $\mu$ on E such that $\mu(h) \leqslant 0$, where $h$ can be expressed as $(*)$, and $\mu(1)=1$.
(c) For $x \in \mathbb{T}$, write $\mu(\mathrm{X})$ pour $\mu\left(\chi_{\mathrm{X}}\right)$, where $\chi_{\mathrm{X}}(x)=1$ if $x \in \mathrm{X}$ and 0 if $x \notin \mathrm{X}$. Then $\mu$ is a finitely additive measure on $\mathbb{T}$, invariant under all $\mathbf{R}_{\alpha}$.
5. (Banach-Tarski Paradox.) -Let $S^{2}$ be the unit sphere of $\mathbb{R}^{3}$, and $G$ the group $\mathrm{SO}_{3}$ of direct isometries of $\mathbb{R}^{3}$, and $e$ its identity element. Define the following equivalence relation on the set of subsets of $S^{2}: X \sim Y$ if and only if there exist finite partitions $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ and $\left(\mathrm{B}_{i}\right)_{i \in \mathrm{I}}$ of X and Y respectively (with the same index set I , and a family $\left(f_{i}\right)_{i \in \mathrm{I}}$ of elements of G such that for all $i \in \mathrm{I}, f_{i}\left(\mathrm{~A}_{i}\right)=\mathrm{B}_{i}$.

The aim is to show, using the axiom of choice, that there exists a partition of $S^{2}$ into two subsets $X, Y$ such that $X \sim Y \sim S^{2}$.
A. (a) Let $f \in \mathrm{G}$ be an irrational rotation (counting the angle in turns), and $\mathrm{E} \subset \mathrm{S}^{2}$ be a set intersecting each orbit at most at one point, and only containing fixed points of $f$. Show that $\mathrm{S}^{2}-\mathrm{E} \sim \mathrm{E}$.
(b) Show that for each countable subset $E \subset S^{2}, S^{2}-E \sim S^{2}$.
B. (a) Find a subset $\Gamma$ in G isomorphic to the "free product" $\mathbb{Z} / 2 * \mathbb{Z} / 3$ (a sum in the category of groups (2.6.2, Example 3)), i.e. two rotations $a$ and $b$ satisfying relations $a^{2}=b^{3}=1$, and not satisfying any relation deducible from these.
(b) Show that $\mathrm{G}=\mathrm{A} \sqcup \mathrm{B} \sqcup \mathrm{C}$ so that

$$
a \mathrm{~A}=\mathrm{B} \cup \mathrm{C}, \quad b \mathrm{~A}=\mathrm{B}, \quad b^{2} \mathrm{~A}=\mathrm{C} .
$$

The elements of $\Gamma$ can be written as words using the letters $a, b, b^{2}$ using the following rule contingent on the parity of the length $n$ of the word and on its first letter: $e \in \mathrm{~A}$, and

$$
\text { parity de } n \begin{cases}0 & \text { first letter }\left\{\begin{array}{ll}
a & \mathrm{~A} \\
b & \mathrm{~A} \\
b^{2} & \mathrm{~B} \\
a & \mathrm{C} \\
b & \text { (irst letter } \\
b^{2} & \mathrm{C}
\end{array}, \$ \mathrm{l}\right.\end{cases}
$$

C. Let E be the set of fixed points of $\Gamma$, and set $\Omega=S^{2}-E$.
(a) Show that $\Gamma$ acts freely on $\Omega$.
(b) Using the axiom of choice, show that there is a partition of $\Omega$ into three sets $\mathscr{A}, \mathscr{B}, \mathscr{C}$ satisfying $\mathscr{A} \sim \mathscr{B} \sim \mathscr{C} \sim \mathscr{B} \cup \mathscr{C}$.
(c) Show that the sets $\mathrm{S}^{2}, \Omega, \mathscr{A}, \mathscr{B}, \mathscr{C}$ are all equivalent.
D. Prove the Banach-Tarski theorem (1.1.5).

### 1.2 Well-Ordered Sets

### 1.2.1 Ordered Sets

An ordered set is a set $E$ equipped with an order relation, i.e. a reflexive, transitive and antisymmetric binary relation, denoted by $x \leqslant y$, in other words, satisfying:
$\left(\mathrm{O}_{1}\right)$ (reflexivity) $\quad \forall x \in \mathrm{E}, x \leqslant x$;
$\left(\mathrm{O}_{2}\right)$ (transitivity) $\quad \forall x \in \mathrm{E}, \forall y \in \mathrm{E}, \forall z \in \mathrm{E},(x \leqslant y$ and $\mathrm{y} \leqslant \mathrm{z}) \Rightarrow \mathrm{x} \leqslant \mathrm{z}$;
$\left(\mathrm{O}_{3}\right)$ (antisymmetry) $\forall x \in \mathrm{E}, \forall y \in \mathrm{E},(x \leqslant y$ and $\mathrm{y} \leqslant \mathrm{x}) \Rightarrow \mathrm{x}=\mathrm{y}$.
An ordered set is said to be totally ordered if its order relation satisfies:
$\left(\mathrm{O}_{4}\right) \forall x \in \mathrm{E}, \forall y \in \mathrm{E}, x \leqslant y$ or $y \leqslant x$.
A chain in an ordered set E is a totally ordered subset X , endowed with the induced order.

### 1.2.2 Upper and Lower Bounds

Let E be an ordered set, X a subset of E and $a \in \mathrm{E}$. The element $a$ is said to be an upper bound of X if, $\forall x \in \mathrm{X}, a \geqslant x$. The element $a$ is said to be the greatest element (resp. a strict upper bound) of X if $a$ is an upper bound of X and $a \in \mathrm{X}$ (resp. $a \notin \mathrm{X}$ ). Lower bounds, least element and strict lower bounds are defined likewise. The element $a$ is said to be the supremum (resp. infimum) of X and is written $a=\sup \mathrm{X}$ (resp. $a=\inf \mathrm{X}$ ) if $a$ is the least upper bound (resp. the greatest lower bound) of X. The subset X is an upper set of E if

$$
\forall x \in \mathrm{X}, \forall y \in \mathrm{E}, \quad y \leqslant x \Rightarrow y \in \mathrm{X}
$$

Lower sets are defined likewise.
A subset $Y$ of $E$ is said to be cofinal if every element of $E$ is bounded above by an element of Y .

### 1.2.3

Definition A well-ordered set E is an ordered set all of whose non-empty subsets have a least element.

A well-order is a total order. Indeed, for $x, y \in \mathrm{E}$, the set $\{x, y\}$ must have a least element.

### 1.2.4 Examples

(1) A totally ordered set is well-ordered.
(2) $\mathbb{N}$ is well-ordered.
(3) $\overline{\mathbb{N}}=\mathbb{N} \cup\{\infty\}$ is well-ordered. More generally, if E is a well-ordered set, then so is the set $\overline{\mathrm{E}}$ obtained by adjoining an element $\omega \notin \mathrm{E}$ to E such that $\forall x \in \mathrm{E}, \omega>x$.
(4) Let A and B be two well-ordered sets with non-trivial intersection. Equip $\mathrm{E}=\mathrm{A} \cup \mathrm{B}$ with the order inducing the given orders on A and B , and such that $\forall x \in \mathrm{~A}, \forall y \in \mathrm{~B}, y>x$. The set E is well-ordered. Denote it by $\mathrm{A}+\mathrm{B}$. Note that $\mathrm{A}+\mathrm{B}$ is not necessarily isomorphic to $\mathrm{B}+\mathrm{A}$.
(5) Every subset of a well-ordered set is well-ordered by the induced order.
(6) The image in an ordered set of a well-ordered set under an increasing map is well-ordered.
(7) The lexicographic order on $\mathbb{N} \times \mathbb{N}$ is defined by

$$
(x, y) \leqslant\left(x^{\prime}, y^{\prime}\right) \Longleftrightarrow x<x^{\prime} \text { or }\left(x=x^{\prime} \text { and } y \leqslant y^{\prime}\right)
$$

Equipped with this order, $\mathbb{N} \times \mathbb{N}$ is well-ordered. More generally, the lexicographic order can be defined on the product of two arbitrary ordered sets. The product of two ordered sets, equipped with the lexicographic order, is well-ordered.

### 1.2.5

Proposition Let E be a well-ordered set, A and B two upper sets of E and $f$ an isomorphism from A onto B . Then $\mathrm{A}=\mathrm{B}$ and $f=1_{\mathrm{A}}$.

Proof Suppose this is false. Let $x$ be the least element of A such that $f(x) \neq x$. Write $\mathrm{S}_{x}$ for the set of strict lower bounds of $x$. Then $x$ is the least element of $\mathrm{E}-\mathrm{S}_{x}$, and so is $f(x)$, and so $f(x)=x$; a contradiction.

### 1.2.6

Proposition Let E and F be two well-ordered sets. Then at least one of the following statements holds:
(i) E is isomorphic to an upper set of F ;
(ii) F is isomorphic to an upper set of E .

Proof Let $\Gamma$ be the set of isomorphisms from an upper set of E onto an upper set of F. Its order by extension (i.e. inclusion order on graphs) is a total order. Indeed, E being totally ordered, if A and $\mathrm{A}^{\prime}$ are two upper sets, then $\mathrm{A} \subset \mathrm{A}^{\prime}$ or $\mathrm{A}^{\prime} \subset \mathrm{A}$. Let $f: \mathrm{A} \rightarrow \mathrm{B}$ and $f^{\prime}: \mathrm{A}^{\prime} \rightarrow \mathrm{B}^{\prime}$ be two elements of $\Gamma$, where $\mathrm{A} \subset \mathrm{A}^{\prime}$ and $\left.f^{\prime}\right|_{\mathrm{A}} \neq f$. Writing $x_{0}$ for the least $x \in \mathrm{~A}$ such that $f^{\prime}(x) \neq f(x)$,

$$
f\left(x_{0}\right)=\inf \left(\mathrm{F}-f\left[0, x_{0}[)=\inf \left(\mathrm{F}-f^{\prime}\left[0, x_{0}[)=f^{\prime}\left(x_{0}\right),\right.\right.\right.\right.
$$

a contradiction.
The union of graphs of elements $g: \mathrm{A}_{g} \rightarrow \mathrm{~B}_{g}$ of $\Gamma$ is again the graph of an element $\bar{g} \in \Gamma$. Indeed, setting $\overline{\mathrm{A}}=\bigcup \mathrm{A}_{g}$ and $\overline{\mathrm{B}}=\bigcup \mathrm{B}_{g}$, for all $x \in \overline{\mathrm{~A}}$ there exists $g \in \Gamma$ such that $x \in \mathrm{~A}_{g}$, and the element $\bar{g}(x)=g(x)$ does not depend on the choice of $g$; The map $\bar{g}$ thus defined is clearly strictly increasing and its image is $\overline{\mathrm{B}}$.

Let us prove by contradiction that $\overline{\mathrm{A}}=\mathrm{E}$ or $\overline{\mathrm{B}}=\mathrm{F}$. Otherwise, let $x=\inf (\mathrm{E}-\overline{\mathrm{A}})$ and $y=\inf (\mathrm{F}-\overline{\mathrm{B}})$. The map $\overline{\bar{g}}: \overline{\mathrm{A}} \cup\{x\} \rightarrow \overline{\mathrm{B}} \cup\{y\}$ which extends $\bar{g}$ by $\overline{\bar{g}}(x)=y$ is an element of $\Gamma$ strictly extending $\bar{g}$, a contradiction.

## Exercises 1.2. (Well-ordered Sets)

1. (Farey Isomorphism.)—Given two rationals $x_{1}=p_{1} / q_{1}$ and $x_{2}=p_{2} / q_{2}$, write $\mathrm{M}\left(x_{1}, x_{2}\right)$ the Farey mean $\left(p_{1}+p_{2}\right) /\left(q_{1}+q_{2}\right)$. Let $\mathbf{D}_{2}$ denote the set of dyadic numbers, i.e. of the form $m / 2^{k}$, and I the interval [0, 1]. Define $\varphi(x)$ for $x=m / 2^{k} \in$ $\mathbf{D}_{2} \cap \mathrm{I}$ by recursion on $k$ by setting:

$$
\varphi(0)=0, \varphi(1)=1, \varphi\left((2 m+1) / 2^{k+1}\right)=\mathrm{M}\left(\varphi\left(m / 2^{k}\right), \varphi\left((m+1) / 2^{k}\right) .\right.
$$

Show that this gives an isomorphism of well-ordered sets from $\mathbf{D}_{2} \cap \mathrm{I}$ onto $\mathbb{Q} \cap \mathrm{I}$.
Show that $\varphi$ can be extended to a homeomorphism $\Phi$ of I. Is the function $\Phi$ Hölderian (i.e. does it satisfy an equality of the form $|\Phi(y)-\Phi(x)| \leqslant c .|y-x|^{\beta}$ )?
2.-Let I be an ordered set and $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ a family of ordered sets. Let $\mathrm{E}=\bigsqcup_{i \in \mathrm{I}} \mathrm{E}_{i}$ be the set of pairs $(i, x)$ such that $i \in \mathrm{I}$ and $x \in \mathrm{E}_{i}$ (summation set, or disjoint union).

Define an order on E by $(i, x) \leqslant(j, y) \Longleftrightarrow i<j$ or $(i=j)$ and $x \leqslant y$. Show that if I and all $\mathrm{E}_{i}$ are well ordered, then so is E .
3.-Fix and integer $k$, and let $\mathrm{A}_{k}$ the set of numbers of type

$$
a_{0}-\frac{1}{a_{1}-\frac{1}{a_{2}-\frac{1}{\cdots-\frac{1}{a_{k}}}}}
$$

where the $a_{i}$ are integers $\geqslant 2$. Show that $\mathrm{A}_{k}$ is a well-ordered set (with respect to the order induced by $\mathbb{R}$ ), isomorphic to lexicographic $\mathbb{N}^{k+1}$. Is the set $\bigcup_{k} \mathrm{~A}_{k}$ wellordered?
4.-(a) Show that a well-ordered set extends to $\mathbb{R}$ (with the induced order) if and only if it is countable. In this case, can it always be embedded into $\mathbb{Q}$ ?
(b) Let A be a well-ordered subset of $\mathbb{R}$, show that its closure $\overline{\mathrm{A}}$ with respect to the topology of $\mathbb{R}$ is well-ordered.
5.-Let E be a well-ordered set. For every $x \in \mathrm{E}$, denote by $\mathrm{S}_{x}$ the set of strict lower bounds of $x$ in E .
(a) The principle of transfinite induction: Let F be a subset of E . Suppose that

$$
(\forall x \in \mathrm{E}) \quad\left(\mathrm{S}_{x} \subset \mathrm{~F} \Rightarrow x \in \mathrm{~F}\right) .
$$

Show that $\mathrm{F}=\mathrm{E}$.
(b) Construction by transfinite induction: Let $U$ be a set, $\Gamma$ the set of maps from the upper sets of E to U and $\varphi$ a map from $\Gamma$ to U .

Show that there exists a unique map $f: \mathrm{E} \rightarrow \mathrm{U}$ such that for all $x \in \mathrm{E}$, $f(x)=\varphi\left(\left.f\right|_{\mathrm{S}_{x}}\right)$. (The proof given should be formally correct).
(c) Construction by transfinite induction with choice and stopping: let $\Phi$ be a map from $\Gamma$ to $\mathfrak{P}(\mathrm{U})$. Show that there is an upper set $\mathbf{J}$ of E and a map $f$ from $\mathbf{J}$ to U such that the following hold simultaneously
(i) $(\forall j \in \mathrm{~J}) \quad f(j) \in \Phi\left(\left.f\right|_{\mathrm{S}_{j}}\right)$;
(ii) $\mathrm{J}=\mathrm{E}$ or $\Phi(f)=\varnothing$.
6.-Let $U$ be a set. Denote by $\Omega$ the set of the isomorphism classes of well-ordered sets whose underlying set is contained in U . Equip $\Omega$ with the order defined by: $\alpha \leqslant \beta$ if $\alpha$ has as representative the upper set of a representative of $\beta$. Show that this gives a well-ordering of $\Omega$.
7.-Let $X$ be an uncountable set, denote by $E$ the set of finite subsets of $X$. Show that:
(a) every totally ordered subset of E is well-ordered;
(b) there are no totally ordered cofinal subsets.
8.-Let I be a well-ordered set and $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ a family of totally ordered sets. The lexicographic product of the sets $\mathrm{E}_{i}$ is the set $\mathrm{E}=\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$ equipped with the following order relation: " $\left(x_{i}\right)_{i \in \mathrm{I}} \leqslant\left(y_{i}\right)_{i \in \mathrm{I}}$ " if $\left(x_{i}\right)_{i \in \mathrm{I}}=\left(y_{i}\right)_{i \in \mathrm{I}}$ or $\left(x_{i}\right) \neq\left(y_{i}\right)$ and $x_{i_{0}}<y_{i_{0}}$, where $i_{0}=\inf \left\{i \mid x_{i} \neq y_{i}\right\}$. Show that:
(a) E is totally ordered.
(b) If the sets $\mathrm{E}_{i}$ are well-ordered, then E is not necessarily well-ordered.
(c) Suppose that I is countable and that each $\mathrm{E}_{i}$ satisfies the following property: every well-ordered subset of $E_{i}$ is countable. Then so does $E$.
(d) As an ordered set, the lexicographic product $\mathbb{N}^{\mathbb{N}}$ is isomorphic to $\mathbb{R}_{+}$.
9.-Show that the lexicographic product $\mathbb{R} \times \mathbb{R}$ is not embedded in $\mathbb{R}$. For this, show that $\mathbb{R}$ satisfies the following property: if $\Lambda$ is a totally ordered set, $\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ and $\left(y_{\lambda}\right)_{\lambda \in \Lambda}$ increasing families of real numbers such that $(\forall \lambda \in \Lambda) x_{\lambda}<y_{\lambda}$ and $(\forall \lambda \in \Lambda)(\forall \mu \in \Lambda) \lambda<\mu \Rightarrow y_{\lambda}<x_{\mu}$, then $\Lambda$ is countable.
10.-Let K be a compact space and $f: \mathrm{K} \rightarrow \mathbb{R}$ a continuous map. Suppose that $f$ has a local maximum at some point of K . Show that $f(\mathrm{~K})$ is well-ordered (with respect to the order induced by that of $\mathbb{R}$ ). Deduce that it is countable.

## $1.3 \tau$-Chains

### 1.3.1

Definition Let E be a well-ordered set, $\tau$ a choice function on E and A a subset of E . The set A is said to be a $\tau$-chain if, for any upper set $\mathrm{C}(\neq \mathrm{A})$ of A , the least element of the set $\mathrm{A}-\mathrm{C}$ is $\tau\left(\mathrm{M}_{\mathrm{C}}\right)$, where $\mathrm{M}_{\mathrm{C}}$ is the set of strict upper bounds of C in E .

Remark If a $\tau$-chain A is not empty, then its least element is $\tau(\mathrm{E})$. If moreover A does not reduce to $\tau(\mathrm{E})$, then the least element of $\mathrm{A}-\tau(\mathrm{E})$ is $\tau\left(\mathrm{M}_{\{\tau(\mathrm{E})\}}\right)$.

The remark shows that $\tau$-chains need to satisfy strong conditions. This will be addressed in greater detail in 1.3.4.

### 1.3.2 Example

Let $\tau$ be a choice function defined on the ordered set $\mathbb{R}$. There is a unique $\tau$-chain $\left(a_{0}, a_{1}, \ldots, a_{n}, \ldots\right)$ of $\mathbb{R}$ isomorphic to $\mathbb{N}$ : it is defined by

$$
a_{0}=\tau(\mathbb{R}), a_{1}=\tau(] a_{0},+\infty[), \ldots, a_{n}=\tau(] a_{n-1},+\infty[), \ldots
$$

### 1.3.3

Proposition Let E be an ordered set and $\tau$ a choice function on E . Then every $\tau$-chain of E is well-ordered.

Proof Let A be a $\tau$-chain of E and X a non-empty subset of A . We show that X has a least element. The set $C$ of strict lower bounds of $X$ in $A$ is an upper set of $A$, and $C$ is distinct from A since $\mathrm{A}-\mathrm{C}$ contains the non-empty subset X . Since A is a $\tau$-chain, $x_{0}=\tau\left(\mathrm{M}_{\mathrm{C}}\right)$ is the least element of $\mathrm{A}-\mathrm{C}$, and so $x_{0}$ is a lower bound of X , but not a strict one. Consequently, $x_{0}$ is the least element of X and A is well-ordered.

### 1.3.4

Proposition Let E be an ordered set, $\tau$ a choice function, A and $\mathrm{A}^{\prime}$ two $\tau$-chains. Then, either A is an upper set of $\mathrm{A}^{\prime}$, or $\mathrm{A}^{\prime}$ is an upper set of A .

Proof Denote by C the union of upper sets common to A and $\mathrm{A}^{\prime}$ : it is the largest among them. Suppose that $\mathrm{C} \neq \mathrm{A}$ and $\mathrm{C} \neq \mathrm{A}^{\prime}$. Then $\tau\left(\mathrm{M}_{\mathrm{C}}\right)=\inf (\mathrm{A}-\mathrm{C})=$ $\inf \left(\mathrm{A}^{\prime}-\mathrm{C}\right)$ necessarily holds. But then $\mathrm{C} \cup\left\{\tau\left(\mathrm{M}_{\mathrm{C}}\right)\right\}$ is an upper set common to A and $\mathrm{A}^{\prime}$ strictly containing C , contradicting the definition of C , and so $\mathrm{C}=\mathrm{A}$ or $\mathrm{C}=\mathrm{A}^{\prime}$,

### 1.3.5

Proposition Let E be an ordered set and $\tau$ a choice function on E .
(a) The union A of all $\tau$-chains of E is a $\tau$-chain.
(b) The $\tau$-chain $\overline{\mathrm{A}}$ has no strict upper bound.

Proof (a) Let C be an upper set of $\overline{\mathrm{A}}$ distinct from $\overline{\mathrm{A}}$, and $x \in \overline{\mathrm{~A}}-\mathrm{C}$. There is a $\tau$-chain A containing $x$. We show that C is an upper set of A . It suffices to show that C is contained in A . Let $y \in \mathrm{C}$ and B a $\tau$-chain containing $y$. By 1.3.4, either B is an upper set of A and $y \in \mathrm{~A}$, or A is an upper set of B . The elements $x$ and $y$ belong to the totally ordered set B . Since $y \in \mathrm{C}$ and $x \notin \mathrm{C}$, it necessarily follows that $y<x$, and if $x \in \mathrm{~A}$ then $y \in \mathrm{~A}$.
(b) Suppose that the set X of strict upper bounds of $\overline{\mathrm{A}}$ is not empty. Adjoining $\tau(\mathrm{X})$ to $\overline{\mathrm{A}}$ gives a $\tau$-chain of E strictly containing $\overline{\mathrm{A}}$, which is impossible. Hence X is empty.

## Exercises 1.3. ( $\tau$-chains)

1.-Let E be a set and $\tau$ a choice function on E . There is a well-order on E such that, for any non-empty subset X of $\mathrm{E}, \tau(\mathrm{X})$ is the least element of X if and only if for all subsets X and Y of E with $\mathrm{Y} \subset \mathrm{X}$ and $\tau(\mathrm{X}) \in \mathrm{Y}, \tau(\mathrm{Y})=\tau(\mathrm{X})$.
2.-Show that every totally ordered set has a well-ordered cofinal subset.
3. (Direct proof of Zermelo's theorem.)—Let E be a set, $\tau: \mathfrak{P}(\mathrm{E})-\{\varnothing\} \rightarrow \mathrm{E}$ a choice function on E . In this exercise, by a $\tau$-chain, we will mean a pair $(\mathrm{F}, \omega)$ where $F$ is a subset of $E$ and $\omega$ an order on $F$ such that for any upper set $A \neq F$ of $F$, $\tau(\mathrm{E}-\mathrm{A})$ is the least element of $\mathrm{F}-\mathrm{A}$. (This notion is not the same as in the text since E is not ordered.) Show that:
(a) if $(\mathrm{F}, \omega)$ is a $\tau$-chain, F is well-ordered by $\omega$;
(b) every upper set of a $\tau$-chain equipped with the induced order is a $\tau$-chain;
(c) if A and $\mathrm{A}^{\prime}$ are two $\tau$-chains, then one of the them is the upper set of the other with respect to the induced order;
(d) let $\left(\mathrm{F}_{i}, \omega_{i}\right)$ be the family of all $\tau$-chains. There is a unique order $\omega$ on $\mathrm{F}=\bigcup \mathrm{F}_{i}$ inducing $\omega_{i}$ on each $\mathrm{F}_{i}$. The pair $(\mathrm{F}, \omega)$ is a $\tau$-chain and $\mathrm{F}=\mathrm{E}$.

Deduce that E is well-ordered.

### 1.4 Inductive Sets and Zorn's Lemma

### 1.4.1

Definition An ordered set is inductive if every ordered subset is bounded above.

### 1.4.2 Examples

(0) The empty set is not inductive: the empty chain is not bounded above.
(1) The set of subsets of a set ordered by inclusion is inductive.
(2) Let E be an inductive set. The set of upper bounds of any $x \in \mathrm{E}$ is inductive.
(3) Let X be a set, E a non-empty subset of the set of its subsets. The set E is said to be of finite character if for any subset $A$ of $X$ the following conditions are equivalent:
(i) $\mathrm{A} \in \mathrm{E}$;
(ii) For any $\mathrm{B} \subseteq \mathrm{A}, \mathrm{B} \in \mathrm{E}$.

Examples of inductive sets are given by the next result:
Every set E of finite character ordered by inclusion is inductive.
Indeed, let $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ be a chain of E , set $\mathrm{A}=\bigcup_{i \in \mathrm{I}} \mathrm{A}_{i}$. Any finite subset B of A is contained in some $\mathrm{A}_{i}$, hence in E .

Example of sets of finite character: the set $\Gamma(\mathrm{X}, \mathrm{Y})$ of graphs of maps from subsets of a set X to a set Y is a subset of $\mathfrak{P}(\mathrm{X} \times \mathrm{Y})$ of finite character. Indeed, for any $\mathrm{A} \subseteq \mathrm{X} \times \mathrm{Y}$, the following properties are equivalent:
(i) $\mathrm{A} \in \Gamma(\mathrm{X}, \mathrm{Y})$;
(ii) the projection $\pi_{\mathrm{A}}: \mathrm{A} \rightarrow \mathrm{X}$ is injective;
(iii) For any finite $\mathrm{B} \subseteq \mathrm{A}$, the projection $\pi_{\mathrm{B}}: \mathrm{B} \rightarrow \mathrm{X}$ is injective.

### 1.4.3

Theorem (ZORN's LEMMA) Every inductive ordered set has a maximal element.

Proof Let E be an inductive ordered set, $\tau$ a choice function on E . The union $\overline{\mathrm{A}}$ of all $\tau$-chains of E has an upper bound $m$. The element $m$ is maximal in E. Indeed, if $m$ was strictly bounded above by $m^{\prime}$, so would $\overline{\mathrm{A}}$, which is impossible by 1.3 .5 (b).

### 1.4.4

Corollary Let E an inductive ordered set. Every element of E is bounded above by a maximal element of E .

Indeed, the set of upper bounds of a point $x$ of E is inductive, and so has a maximal element which is maximal in E .

## Exercises 1.4. (Inductive sets and Zorn's lemma)

1. (Path- and arc-connectedness.) -Let I denote the interval [ 0,1 ]. Let X be a Hausdorff topological space, $a$ and $b$ points of X. Suppose there is a path from $a$ to $b$ in X, i.e. a continuous map $\gamma: \mathrm{I} \rightarrow \mathrm{X}$ such that $\gamma(0)=a$ and $\gamma(1)=b$. The aim is to show that there is an injective path $\eta: \mathrm{I} \rightarrow \mathrm{X}$ from $a$ to $b$.
(a) Consider the set $\Omega$ of open subsets U of I not containing 0 or 1 and such that, for any connected component $\mathrm{J}=] s, s^{\prime}$ [ of $\mathrm{U}, \gamma(s)=\gamma\left(s^{\prime}\right)$. Show that $\Omega$ ordered by inclusion is inductive.
(b) Show that, if the open subset W of I is a maximal element of $\Omega$, the only isolated points of the closed subset $\mathrm{I}-\mathrm{W}$ are 0 or 1 .
(c) Let $U$ be an open subset of I for which 0 or 1 are the only isolated points of $\mathrm{I}-\mathrm{U}$. Show that there is an increasing continuous function $h: \mathrm{I} \rightarrow \mathrm{I}$, constant on each connected component of U , satisfying the following condition: if $h(t)=h\left(t^{\prime}\right)$ with $t \neq t^{\prime}$, then there is a connected component J of U such that $t$ and $t^{\prime}$ belong to the closure of J .
(d) Conclude.

### 1.5 Applications of Zorn's Lemma

### 1.5.1

Theorem Let X and Y be two sets. There is an injection from X into Y or an injection from Y into X .

Proof Let E be the set of graphs of bijections from subsets of X onto subsets of Y ordered by inclusion. It is also the set of subsets $G$ of $X \times Y$ such that the projections $\pi_{1}: \mathrm{G} \rightarrow \mathrm{X}$ and $\pi_{2}: \mathrm{G} \rightarrow \mathrm{Y}$ are injective. The set E is of finite character and so is inductive. By Zorn's lemma, E has as maximal element the graph H of a bijection from a subset A of X onto a subset B of Y . Then $\mathrm{A}=\mathrm{X}$ or $\mathrm{B}=\mathrm{Y}$. Otherwise $\mathrm{H} \cup\{x, y\}$, where $x \in \mathrm{X}-\mathrm{A}$ and $y \in \mathrm{Y}-\mathrm{B}$, would be a strict upper bound of H in E.

### 1.5.2

Theorem (ZERMELO's THEOREM) Every set can be well-ordered.

Proof Let E be a set, and $\Omega$ denote the set of pairs ( $\mathrm{F}, \omega$ ), where $\mathrm{F} \subseteq \mathrm{E}$ is wellordered by $\omega$. The set $\Omega$ can be ordered as follows: " $(\mathrm{F}, \omega) \leqslant\left(\mathrm{F}^{\prime}, \omega^{\prime}\right)$ " if F is an upper set of $\mathrm{F}^{\prime}$ with respect to $\omega^{\prime}$ and if $\omega=\left.\omega^{\prime}\right|_{\mathrm{F}}$. With this order, $\Omega$ is inductive. Indeed, let $\left(\mathrm{F}_{i}, \omega_{i}\right)_{i \in \mathrm{I}}$ be a totally ordered family of elements of $\Omega$. There is an order $\omega$ on $\mathrm{F}=\bigcup \mathrm{F}_{i}$ such that, for all $i \in \mathrm{I},\left.\omega\right|_{\mathrm{F}_{i}}=\omega_{i}$. Equipped with $\omega$, the set F is well-ordered. Let $(\overline{\mathrm{F}}, \bar{\omega})$ be a maximal element of $\Omega$. We show by contradiction that $\overline{\mathrm{F}}=\mathrm{E}$. Let $a \in \mathrm{E}-\overline{\mathrm{F}}$. Let $\omega^{\prime}$ be the order on $\mathrm{F}^{\prime}=\overline{\mathrm{F}} \cup\{a\}$ inducing $\bar{\omega}$ on $\overline{\mathrm{F}}$ and such that $\forall x \in \overline{\mathrm{~F}}, a>x$. Then $\left(\mathrm{F}^{\prime}, \omega^{\prime}\right)$ is a strict upper bound of $(\overline{\mathrm{F}}, \bar{\omega})$, contradicting the maximality of $(\overline{\mathrm{F}}, \bar{\omega})$.

### 1.5.3

Zorn's lemma is also applied to prove the following results:

1. Every commutative ring $A \neq 0$ contains an ideal $\mathfrak{m}$ such that $A / \mathfrak{m}$ is a field. (Krull's Theorem, 3.1.5.)
2. Let $A$ be a commutative ring. The intersection of prime ideals of $A$ is the set of nilpotent elements of A (3.1.6).
3. Every vector space has a basis (3.4.2).
4. Let A be a principal ring. Every submodule of a free A-module is free (3.5.1). Every divisible A-module is injective (3.10.11, Corollary, b).
5. Every product of compact spaces is compact. (Tychonoff's Theorem, 1.7.7.)
6. Every open cover of a paracompact space has a shrinking (N. Bourbaki, General Topology, Chap. 5 § 10 [1], Chap. 9, § 4, cor. 1).
7. Let E be a normed vector space over $\mathbb{R}$ or $\mathbb{C}$ and $x \in \mathrm{E}$ with norm 1 . There is a linear form $\xi$ on E such that $\|\xi\|=1$ and $\xi(x)=1$. (Hahn-Banach Theorem.)

## Exercises 1.5 (Applications of Zorn's lemma)

1. (Transfinite line.)-(a) Show that there is a well-ordered set $\Omega$, unique up to isomorphism, such that $\Omega$ is not countable but every upper set other than $\Omega$ is countable.
(b) The transfinite line (or rather transfinite half-line) is the set $\mathrm{L}=\Omega \times[0,1[$ ordered by the lexicographic product. Endow L with the order topology: a fundamental system of neighbourhoods of $x$ is given by the intervals ] $x^{\prime}, x^{\prime \prime}[$ where $x^{\prime}<x<x^{\prime \prime}$. Set 0 to be the least element $(0,0)$ and $\mathrm{L}^{*}=\mathrm{L}-\{0\}$. Show that, for all $x \in \mathrm{~L}$, the interval $[0, x]$ is homeomorphic to $[0,1]$, but that L is not homeomorphic to $\mathbb{R}_{+}$. Show that L is neither a metric space, nor the countable union of compact spaces. Show that any increasing function $f: \mathrm{L} \rightarrow \mathbb{R}$ is stationary (i.e. $(\exists a)(\forall x>a) f(x)=f(a))$.
(c) Show that every increasing sequence in L converges. Show that L is not compact, but that every sequence in $L$ has a convergent subsequence.

Show that, if A and B are two disjoint closed subsets of L, then one of them is bounded (suppose this does not hold and construct an increasing sequence ( $x_{n}$ ) such that $x_{n} \in \mathrm{~A}$ for odd $n$ and $x_{n} \in \mathrm{~B}$ for even $n$ ). Deduce that L is normal (two disjoint closed subsets have disjoint neighbourhoods). Show that every continuous function $f: \mathrm{L} \rightarrow \mathbb{R}$ is stationary.
(d) Show that there is a continuous function $f: \mathrm{L} \rightarrow \mathrm{L}$ such that $(\forall x) f(x) \geqslant x$ and such that the set of $x$ satisfying $f(x)>x$ is unbounded, but that there is no such function $f$ with $f(x)>x$ for all $x$. Show that every continuous map $f: \mathrm{L} \rightarrow \mathrm{L}$ such that $(\forall x) f(x)<x$ is stationary.
(e) Set $\mathrm{M}=\left\{(x, y) \in \mathrm{L}^{2} \mid 0<x<y\right\}$. Show that M , endowed with the projection $(x, y) \rightarrow x$, is a non-trivial fiber bundle (4.1.7) with fiber $\mathbb{R}$ over $\mathrm{L}^{*}$, but that its restriction to any compact subspace of $\mathrm{L}^{*}$ is a trivial bundle.

Show that L is not contractible.
(f) Show that $L^{*}$ has a structure of a real-analytic manifold (assume that every real-analytic manifold homeomorphic to $\mathbb{R}$ is isomorphic to $\mathbb{R}$ ).

### 1.6 Noetherian Ordered Sets

### 1.6.1

Proposition and Definition Let E be an ordered set. The following properties are equivalent:
(i) every increasing sequence of elements of E is stationary;
(ii) Every non-empty subset X of E has a maximal element in X .

If these properties hold, the ordered set E is said to be Noetherian.
Proof (ii) $\Rightarrow$ (i): if $\left(x_{n}\right)$ is an increasing sequence of elements of E , then the set of these elements has a maximal element $x_{n_{0}}$; for $n \geqslant n_{0}, x_{n} \geqslant x_{n_{0}}$ and so $x_{n}=x_{n_{0}}$.
(i) $\Rightarrow$ (ii): Suppose that the nonempty subset X of E has no maximal element and let $\tau$ be a choice function on X . Define a sequence of elements $\left(x_{n}\right) \in \mathrm{X}$ inductively by setting $x_{0}=\tau(\mathrm{X})$ and $x_{n+1}=\tau\left(\mathrm{M}\left(x_{n}\right)\right)$, where $\mathrm{M}\left(x_{n}\right)$ is the set of strict upper bounds for $x_{n}$ in X . The sequence $\left(x_{n}\right)$ is strictly increasing, a contradiction.

An ordered set will be said to be co-Noetherian if it is Noetherian with respect to the dual order.

Example The ordered set $\mathbb{N}$ is co-Noetherian. More generally, every well-ordered set is co-Noetherian.

### 1.6.2

Proposition Let E be an ordered set. If there is a strictly increasing map $\varphi: \mathrm{E} \rightarrow \mathbb{N}$ (i.e. such that $x<y \Rightarrow \varphi(x)<\varphi(y)$ ), then E is co-Noetherian.

Proof If there were a strictly decreasing sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ in E , then the sequence $\left(\varphi\left(x_{n}\right)\right)$ would be strictly decreasing in $\mathbb{N}$, which is impossible.

Examples (1) Let X be a set. The set of finite subsets of X is co-Noetherian with respect to inclusion.
(2) Let E be a vector space. The set of finite dimensional vector subspaces of E is co-Noetherian with respect to inclusion.

### 1.6.3

Proposition (Principle of Noetherian Induction.) Let E be a Noetherian ordered set, and for all $x \in \mathrm{E}, \mathrm{M}(x)$ the set of strict upper bounds for $x$ in E . Let F be a subset of E such that $(\forall x \in \mathrm{E}) \mathrm{M}(x) \subset \mathrm{F} \Rightarrow x \in \mathrm{~F}$. Then $\mathrm{F}=\mathrm{E}$.

Proof Set $\mathrm{X}=\mathrm{E}-\mathrm{F}$. If $\mathrm{X} \neq \varnothing$, let $x$ be a maximal element of X . Then $\mathrm{M}(x) \subset \mathrm{F}$ and $x \notin \mathrm{~F}$, a contradiction.
(This proposition will be needed in 3.6.4 and 3.6.7.)

### 1.7 Tychonoff's Theorem

### 1.7.1 Filters

Let X be a topological space. A filter on X is a subset $\mathscr{F}$ of $\mathfrak{P}(\mathrm{X})$ satisfying the following conditions:
$\left(\mathrm{F}_{1}\right) \mathscr{F}$ is closed under finite intersection.
$\left(\mathrm{F}_{2}\right) \mathrm{F} \in \mathscr{F}$ and $\mathrm{G} \supset \mathrm{F} \Rightarrow \mathrm{G} \in \mathscr{F}$.
$\left(\mathrm{F}_{3}\right) \varnothing \notin \mathscr{F}$.
Example For all $x \in \mathrm{X}$, the set $\mathscr{V}_{x}$ of neighbourhoods of $x$ is a filter on X .
Given two filters $\mathscr{F}$ and $\mathscr{G}, \mathscr{F}$ is said to be finer than $\mathscr{G}$ if $\mathscr{F} \supset \mathscr{G}$. If $\mathscr{H}$ is finer than both $\mathscr{F}$ and $\mathscr{G}$, then $\mathscr{F}$ and $\mathscr{G}$ are said to be compatible. Equivalently, $\forall \mathrm{F} \in \mathscr{F}, \forall \mathrm{G} \in \mathscr{G}, \mathrm{F} \cap \mathrm{G} \neq \varnothing$. The supremum of $\mathscr{F}$ and $\mathscr{G}$ is then called the intersection filter and is written filter $\mathscr{F} \wedge \mathscr{G}=\{\mathrm{F} \cap \mathrm{G}\}_{\mathrm{F} \in \mathscr{F}, \mathrm{G} \in \mathscr{G}}$. The filter $\mathscr{F}$ is convergent to $x \in \mathrm{X}$ (resp. has $x$ as cluster point) if $\mathscr{F}$ is finer than (resp. compatible with) $\mathscr{V}_{x}$. The filter $\mathscr{F}$ has $x$ as cluster point if and only if for all F of $\mathscr{F}, x$ belongs to the closure $\overline{\mathrm{F}}$ of F .

Commentary. A property that may or may not hold for a point of X is described by using a subset of X : the set of $x$ satisfying this property. A filter is meant to describe a more subtle condition that may be variably satisfied, each set of the filter (or of a a base of this filter) satisfying it to some degree.

### 1.7.2 Ultrafilters

An ultrafilter is a maximal filter, namely one for which there is no finer one.
Example For all $x \in \mathrm{X}$, the set $\mathcal{U}_{x}$ of $\mathrm{F} \subset \mathrm{X}$ such that $x \in \mathrm{~F}$ is an ultrafilter.
If $x$ is a cluster point of an ultrafilter $\mathscr{U}$, then $\mathscr{U}$ converges to $x$. Indeed $\mathscr{U} \wedge \mathscr{V}_{x}$ is finer than $\mathscr{U}$, and so equal to $\mathscr{U}$, and $\mathscr{U}$ is finer than $\mathscr{V}_{x}$.

Proposition For every filter $\mathscr{F}$ on X, there is an ultrafilter finer than $\mathscr{F}$.
Proof It readily follows that the set of of filters, ordered by the fineness relation, is inductive. The proposition is then a consequence of 1.4.4.

### 1.7.3 Characterization of Compact Sets

Proposition Let X be a Hausdorff topological space. The following conditions are equivalent:
(i) X is compact;
(ii) every filter on X has a cluster point;
(iii) every ultrafilter on X is convergent.

Proof (i) $\Leftrightarrow$ (ii). By definition, (i) means that every open cover of X has a finite subcover. Considering their complements in X , (i) is seen to be equivalent to
(i') if $\left(\mathrm{F}_{i}\right)_{i \in \mathrm{I}}$ is a closed family of X such that $\bigcap_{j \in \mathrm{~J}} \mathrm{~F}_{j} \neq \varnothing$ for all finite $\mathrm{J} \subset \mathrm{I}$, then $\bigcap_{i \in \mathrm{I}} \mathrm{F}_{i} \neq \varnothing$.
(ii) $\Rightarrow$ ( $\mathrm{i}^{\prime}$ ) follows by taking the filter generated by the subsets $\mathrm{F}_{i}$, and ( $\mathrm{i}^{\prime}$ ) $\Rightarrow$ (ii) by taking the family $(\overline{\mathrm{F}})_{\mathrm{F} \in \mathscr{F}}$.
(ii) $\Leftrightarrow$ (iii). Since an ultrafilter having a cluster point is convergent, (ii) $\Rightarrow$ (iii). Suppose (iii) holds and let $\mathscr{F}$ be a filter. By Proposition 1.7.2 there is an ultrafilter $\mathscr{U}$ finer than $\mathscr{F}$, and if $\mathscr{U}$ converges to $x$, then $x$ is a cluster point of the filter $\mathscr{F}$.

### 1.7.4 Image Filter

Let $f: \mathrm{X} \rightarrow \mathrm{Y}$ be a map and $\mathscr{F}$ a filter on X . The set $f_{*} \mathscr{F}$ of subsets $\mathrm{G} \subset \mathrm{Y}$ such that $f^{-1}(\mathrm{G}) \in \mathscr{F}$ is a filter on Y , called the image filter of $\mathscr{F}$ by $f$.

If X and Y are topological spaces, if $f$ is continuous and if $\mathscr{F}$ converges to a point $x \in \mathrm{X}$, the image filter $f_{*} \mathscr{F}$ converges to $f(x)$.

### 1.7.5

Proposition Let $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ be a family of topological spaces. Set $\mathrm{X}=\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$ and $p_{i}$ the projection $\mathrm{X} \rightarrow \mathrm{X}_{i}$. A filter $\mathscr{F}$ on X is convergent if and only if for all $i \in \mathrm{I}$, so is $\left(p_{i}\right)_{*} \mathscr{F}$.

Proof The condition is obviously necessary. Let us show it is sufficient. For all $i$, let $x_{i} \in \mathrm{X}_{i}$ be such that $\left(p_{i}\right)_{*} \mathscr{F}$ is convergent to $x_{i}$, and set $x=\left(x_{i}\right)$.

If J is a finite subset of I and $\left(\mathrm{V}_{j}\right)_{j \in \mathrm{~J}}$ a family where $\mathrm{V}_{j}$ is a neighbourhood of $x_{j}$, the set $\bigcap_{j \in \mathrm{~J}} p_{j}^{-1} \mathrm{~V}_{j}$ belongs to $\mathscr{F}$. Since such sets form a fundamental system of neighbourhoods of $x$, the filter $\mathscr{F}$ is convergent to $x$.

Remark Even when, for all $i$, the filter $\left(p_{i}\right)_{*} \mathscr{F}$ has a cluster point, this does not necessarily imply that this is also the case for $\mathscr{F}$.

### 1.7.6

Proposition The image filter of an ultrafilter is an ultrafilter.
Proof A filter $\mathscr{F}$ on X is an ultrafilter if and only if, for every $\mathrm{A} \subseteq \mathrm{X}, \mathrm{A} \in \mathscr{F}$ or $\mathrm{X}-\mathrm{A} \in \mathscr{F}$ (consider $\{\mathrm{F} \cap \mathrm{A}\}_{\mathrm{F} \in \mathscr{F}}$ and $\{\mathrm{F} \cap(\mathrm{X}-\mathrm{A})\}_{\mathrm{F} \in \mathscr{F}}$ ). If $\mathscr{U}$ is an ultrafilter on X , then, for all $\mathrm{B} \subset \mathrm{Y}$,

$$
f^{-1}(\mathrm{~B}) \in \mathscr{U} \quad \text { or } \quad f^{-1}(\mathrm{Y}-\mathrm{B})=\mathrm{X}-f^{-1}(\mathrm{~B}) \in \mathscr{U} .
$$

### 1.7.7

## Theorem (TYCHONOFF'S THEOREM) An arbitrary product of compact spaces is compact.

Proof Let $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ be a family of compact spaces and set $\mathrm{X}=\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$. Let $\mathscr{U}$ be an ultrafilter on X . For all $i$, the filter $p_{i *} \mathscr{U}$ is an ultrafilter 1.7.6, and so is convergent 1.7.3.

Hence, $\mathscr{U}$ is convergent 1.7.5. Since this holds for all ultrafilters on $X$, the space X is compact 1.7.3,

## Exercises 1.7 (Filters and ultrafilters, Tychonoff's theorem)

1. (Moore-Smith convergence.)—An ordered set is said to be directed if all finite subsets are bounded above, in other words, if it is nonempty and every pair of elements has an upper bound. Any totally ordered set is directed, in particular $\mathbb{N}$. The following generalizes the notion of sequence.

Given a set X , a Moore-Smith sequence in X is a family $\left(x_{i}\right)_{i \in \mathrm{I}}$ where I is a directed set. If X is a topological space and $a \in \mathrm{X}$, then $\left(x_{i}\right)$ is said to converge to $a$ if for all neighbourhoods V of $a, \exists k \in \mathrm{I}, \forall i \geqslant k, x_{i} \in \mathrm{~V}$.

Given two sets X and Y , a map $f: \mathrm{X} \rightarrow \mathrm{Y}$ and a filter $\mathscr{F}$ on X , the subsets V of Y such that $f^{-1}(\mathrm{~V}) \in \mathscr{F}$ form a filter on Y written $f_{*} \mathscr{F}$ and called image filter. If X is a topological space, for $a \in \mathrm{X}$, denote by $\mathscr{V}_{\mathrm{X}, a}$ the neighbourhood filter of $a$ in X .
(a) Let X and Y be two topological spaces, $f: \mathrm{X} \rightarrow \mathrm{Y}$ a map, let $a$ be a point of X and $b=f(a)$. Show that the following conditions are equivalents:
(i) $f$ is continuous at $a$;
(ii) $f_{*} \mathscr{V}_{\mathrm{X}, a} \supset \mathscr{V}_{\mathrm{Y}, b}$;
(iii) For every Moore-Smith sequence ( $x_{i}$ ) convergent to $a$ in X , the (Moore-Smith) sequence $\left(f\left(x_{i}\right)\right)$ is convergent to $b$ in Y .
(b) Let (iii ${ }_{\mathbb{N}}$ ) be the condition analogous to (iii) for ordinary sequences. Give an example showing that (iii $\mathbb{N}$ ) does not imply (iii). Show that Lebesgue's dominated convergence theorem provides such an example.
(The convergence of Moore-Smith sequences provides a language which, like that of filters, enables us to tackle non-metric topological spaces.)
2.-Let X be an infinitely countable set. Set A to be the $\operatorname{ring} \mathbb{R}^{\mathrm{X}}$. For $f \in \mathrm{~A}$, denote by $\mathrm{V}(f)$ the set of $x \in \mathrm{X}$ such that $f(x)=0$.
(a) For $f, g \in \mathrm{~A}$, show that $f$ divides $g$ if and only if $\mathrm{V}(f) \subset \mathrm{V}(g)$.
(b) Let $\mathrm{I}(\neq \mathrm{A})$ be an ideal of A. Show that $\mathscr{F}_{\mathrm{I}}=\{\mathrm{V}(f)\}_{f \in \mathrm{I}}$ is a filter on X , and that $\mathrm{I}=\left\{f ; \mathrm{V}(f) \in \mathscr{F}_{\mathrm{I}}\right\}$.
(c) Show that, if $\mathscr{F}$ is a filter on X , the set $\mathrm{I}_{\mathscr{F}}$ of $f \in \mathrm{~A}$ such that $\mathrm{V}(f) \in \mathscr{F}$ is a strict ideal of A.

Show that $\mathrm{I} \mapsto \mathscr{F}_{\mathrm{I}}$ and $\mathscr{F}_{\mathrm{I}} \mapsto \mathrm{I}_{\mathscr{F}}$ are mutually inverse bijections, from the set of filters on X onto the set of strict ideals of A.
(d) Let $\mathscr{F}$ be a filter on X and $\mathrm{I}=\mathrm{I}_{\mathscr{F}}$. Show that the following three conditions are equivalent:
(i) I is a maximal ideal;
(ii) I est a prime ideal;
(iii) $\mathscr{F}$ is an ultrafilter.
3. (Hyper-reals.) -We keep the notation of the previous exercise. Let $\mathscr{U}$ be an ultrafilter on X and $\mathbb{R}_{\mathscr{U}}$ the field $\mathrm{A} / \mathrm{I}_{\mathscr{U}}$. The aim is to prove properties stating that $\mathbb{R}_{\mathscr{U}}$ is a model of the real field in non-standard analysis.
(a) $\mathbb{R}_{\mathcal{U}}$ is a totally ordered field and there is a natural embedding of $\mathbb{R}$ in $\mathbb{R}_{\mathscr{U}}$. For any subset E of $\mathbb{R}$, set $\mathrm{E}_{\mathscr{U}}$ to be the image of $\mathrm{E}^{\mathrm{X}}$ in $\mathbb{R}_{\mathscr{U}}$. For a function $\varphi: \mathbb{R}^{n} \rightarrow \mathbb{R}$, let $\varphi_{\mathscr{U}}$ be the map $\mathbb{R}_{\mathcal{U}}^{n} \rightarrow \mathbb{R}_{\mathcal{U}}$ induced by $\left(f_{1}, \ldots, f_{n}\right) \rightarrow \varphi \circ\left(f_{1}, \ldots, f_{n}\right)$ from $\mathrm{A}^{n}$ to A. If $\varphi$ is $x \rightarrow|x|,(x, y) \rightarrow x+y$, etc., write $x \rightarrow|x|,(x, y) \rightarrow x+y$, etc. the map $\varphi_{\varkappa}$.

There are elements $\omega \in \mathbb{R}_{\mathscr{U}}$ such that $\forall r \in \mathbb{R}, \omega>r$. They are called infinitely large, while $\varepsilon \in \mathbb{R}_{\mathcal{U}}$ is infinitely small if $(\forall r \in \mathbb{R}, r>0)|\varepsilon|<r$. If $\xi-x$ is infinitely small, $\xi \in \mathbb{R}_{\mathcal{U}}$ is said to have shadow $x \in \mathbb{R}$.
(b) For $\quad \mathrm{E}, \mathrm{E}^{\prime} \subset \mathbb{R}, \quad\left(\mathrm{E} \cap \mathrm{E}^{\prime}\right)_{\mathscr{U}}=\mathrm{E}_{\mathscr{U}} \cap \mathrm{E}_{\mathscr{U}}^{\prime}, \quad\left(\mathrm{E} \cup \mathrm{E}^{\prime}\right)_{\mathcal{U}}=\mathrm{E}_{\mathscr{U}} \cup \mathrm{E}_{\mathscr{U}}^{\prime}$. For $f: \mathbb{R} \rightarrow \mathbb{R}, f_{\mathscr{U}}\left(\mathrm{E}_{\mathscr{U}}\right)=(f(\mathrm{E}))_{\varkappa}$.
(c) Let $\mathrm{E} \subset \mathbb{R}$. The set $\mathrm{E}_{\mathscr{U}}$ contains all infinitely large elements if and only if $\exists r \in \mathbb{R}, \mathrm{E} \supset[r, \infty[$.
(d) Let $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ be a function. $\varphi$ is continuous if and only if $\forall x \in \mathbb{R}, \forall \xi \in \mathbb{R}_{\mathscr{U}}$, ( $\xi$ has shadow $x) \Rightarrow\left(\varphi_{\varkappa}(\xi)\right.$ has shadow $\left.\varphi(x)\right)$.

## Reference

## Chapter 2 <br> Categories and Functors

## Introduction

Some mathematical constructions are "natural" because they do not involve any arbitrary choice. These constructions can be transferred from one model to another representing the same situation. Category theory has been elaborated so as give this vague but powerful idea a precise meaning which can be used in mathematical proofs.

A logical difficulty arises. It is well-known that "the set of all sets" is an ill-defined concept as it leads to a contradiction: For any property $\mathscr{P}$, one should be able to define the set of sets satisfying $\mathscr{P}$, in particular $\mathrm{Y}=\{\mathrm{X} \mid \mathrm{X} \notin \mathrm{X}\}$. But then $\mathrm{Y} \in \mathrm{Y}$ implies $\mathrm{Y} \notin \mathrm{Y}$, and $\mathrm{Y} \notin \mathrm{Y}$ implies $\mathrm{Y} \in \mathrm{Y}$ (Russell's paradox). Likewise, it is not possible to talk of the set of all groups, etc.

This difficulty can be overcome in two ways. One involves the notion of "classes" that are not sets. This however is not altogether satisfactory. The categories constructed from these classes are not veritably objects of the theory, and can only be manipulated with great caution. The use of quantifiers cannot be justified, statements do not really amount to theorems that can be formally applied by substituting to an object a letter governed by the quantifier " $\forall$ ". No proof of Theorem 2.3.5 is known in this framework (the proof we give uses the axiom of choice).

Another way to avoid the problem, and which we have adopted here, amounts to considering "small" categories, i.e. categories whose set of objects is truly a set. The category of groups is not a well-defined concept, but that of groups whose underlying set belongs to a predetermined set U is. Although "statements" are not as general as one would wish, they are truly statements on which formal reasoning can be built with total peace of mind. Drawbacks are minor, and they completely disappear in practice if the "rag-bag" U is a universe (2.1.3). Admitting the universe axiom, whose role is not critical, greatly simplifies the language.

### 2.1 Categories

### 2.1.1

Definition A category consists of a set $\mathscr{C}$ whose elements are called objects,
(a) for all objects $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, of a set $\operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})$ or simply $\operatorname{Hom}(\mathrm{X}, \mathrm{Y})$, called the set of morphisms of $\mathscr{C}$ from $X$ to Y;
(b) and, for any three objects $\mathrm{X}, \mathrm{Y}, \mathrm{Z} \in \mathscr{C}(f, g) \mapsto g \circ f$ from $\operatorname{Hom}(\mathrm{X}, \mathrm{Y}) \times \operatorname{Hom}(\mathrm{Y}, \mathrm{Z})$ to $\operatorname{Hom}(\mathrm{X}, \mathrm{Z})$ called composition; satisfying the following conditions:
$\left(\mathrm{C}_{1}\right)(\forall \mathrm{X}, \mathrm{Y}, \mathrm{Z}, \mathrm{T} \in \mathscr{C})(\forall f \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y}))$
$(\forall g \in \operatorname{Hom}(\mathrm{Y}, \mathrm{Z}))(\forall h \in \operatorname{Hom}(\mathrm{Z}, \mathrm{T}))(h \circ g) \circ f=h \circ(g \circ f)$.
$\left(\mathrm{C}_{2}\right)(\forall \mathrm{X} \in \mathscr{C})(\exists e \in \operatorname{Hom}(\mathrm{X}, \mathrm{X}))(\forall \mathrm{Y} \in \mathscr{C})$

$$
((\forall f \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y})) f \circ e=f \text { and }(\forall g \in \operatorname{Hom}(\mathrm{Y}, \mathrm{X})) e \circ g=g) .
$$

The element $e$ whose existence is predicated in $\left(\mathrm{C}_{2}\right)$ is unique. It is called the identity of X and is written $1_{\mathrm{X}}$.

A morphism $f \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y})$ is often written $f: \mathrm{X} \rightarrow \mathrm{Y}$.
An endomorphism of X is a morphism from X to X and $\operatorname{End}(\mathrm{X})=\operatorname{Hom}(\mathrm{X}, \mathrm{X})$.
A morphism $f: \mathrm{X} \rightarrow \mathrm{Y}$ from X to Y is said to be an isomorphism from X onto Y if there is a morphism $g: \mathrm{Y} \rightarrow \mathrm{X}$ such that $g \circ f=1_{\mathrm{X}}$ and $f \circ g=1_{\mathrm{Y}}$. Then $g$ is unique and is called the inverse of $f$. If $f$ is an isomorphism from X onto X , then $f$ is said to be an automorphism of X .

Two objects X and Y de $\mathscr{C}$ are isomorphic if there is an isomorphism from X onto Y . We then write $\mathrm{X} \approx \mathrm{Y}$.

### 2.1.2 Examples

Let $U$ be a set.
(1) The set $\mathscr{C}$ of all sets belonging to $U$ with $\operatorname{Hom}(X, Y)$ the set of maps from X to Y , for $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, composition being map composition, is a category called the category of sets of $U$ and of maps, written $\mathscr{E} n s_{U}$.
(2) The set $\mathscr{C}$ of groups (resp. commutative unitary rings, resp. modules over a given ring A, resp. topological spaces) whose underlying set belongs to $U$, with the morphisms being group (resp. unitary ring, resp. A-linear, resp. continuous) homomorphisms or maps, is a category $\mathscr{E}_{\boldsymbol{\gamma}_{U}}$ called the category of groups and group homomorphisms of $U$ (resp. the category $\mathscr{A} \boldsymbol{n} \boldsymbol{n}_{U}$ of commutative unitary rings and unitary ring homomorphisms, resp. the category A- $\mathscr{M} \boldsymbol{\sigma} \boldsymbol{d}_{\mathrm{U}}$ of A-modules and A-linear maps, resp. the category $\mathscr{T} \sigma \mathcal{p}_{\mathrm{U}}$ of topological spaces and continuous maps).
(3) Let G be a group. A G-set is a set X endowed with a left action of G on X , i.e. a map $(g, x) \mapsto g \cdot x$ from $\mathrm{G} \times \mathrm{X}$ to X such that $\forall x \in \mathrm{X} e \cdot x=x$ where $e$ is the identity element of G and $\forall g \in \mathrm{G} \forall h \in \mathrm{G} \forall x \in \mathrm{X} g \cdot(h \cdot x)=(g \cdot h) \cdot x$.

This amounts to defining a homomorphism $\rho_{\mathrm{X}}$ from $G$ to the group $\mathfrak{S}(\mathrm{X})$ of permutations of X . If X and Y are two G -sets, a G-morphism from X to Y is a map $f: \mathrm{X} \rightarrow \mathrm{Y}$ such that $\forall g \in \mathrm{G} \forall x \in \mathrm{X} f(g \cdot x)=g \cdot f(x)$. Equivalently, $\forall g \in \mathrm{G}$, the diagram

is commutative.
A category G- $\mathscr{E} \boldsymbol{n s}_{\mathrm{U}}$ is defined by taking as its objects the G-sets whose underlying sets belong to U and as its morphisms the G-morphisms.
(4) Let X and Y be two topological spaces, $f_{0}$ and $f_{1}$ two continuous maps from X to Y . The maps $f_{0}$ and $f_{1}$ are said to be homotopic if there is a continuous map $h$ from $[0,1] \times \mathrm{X} \rightarrow \mathrm{Y}$ such that, for all $x \in \mathrm{X}, h(0, x)=f_{0}(x)$ and $h(1, x)=f_{1}(x)$. The homotopy relation is an equivalence relation on the set of continuous maps from X to Y .

Let $\mathscr{C}$ be the set of topological spaces whose underlying sets belong to $U$ and with the homotopy classes of continuous maps as morphisms, composition being induced by map composition on passing to the quotient. This gives a category called the category of topological spaces and homotopy classes of continuous maps.
(5) Let I be a preordered set (i.e. endowed with a reflexive and transitive binary relation, see 1.2 .1 ). The set I becomes a category by setting $\operatorname{Hom}_{\mathrm{I}}(i, j)$ to be a singleton if $i \leq j$ and empty otherwise, composition being the only possible one.
(6) Let $M$ be a monoid (i.e. a set endowed with an associative composition law with an identity element). A category $\mathscr{M}$ can be defined by taking a singleton with the element $x$ and setting $\operatorname{Hom}(x, x)=\mathrm{M}$, composition being that of M .

### 2.1.3 Universe

In Examples 1-4 of 2.1.2 and other similar ones, if $U$ is taken to be the set to which all the sets in question belong, then, for the sake of simplicity, any mention of $U$ is omitted, and the category is said to be that of sets, groups, etc., and is written $\mathscr{E}$ ns, $\mathscr{G}_{r}, \mathscr{T} \mathscr{O}_{\boldsymbol{p}}$, etc.

The notion of universe provides a framework where the drawbacks of this simplification are alleviated. The set U is said to be a universe if it satisfies the following conditions:
$\left(\mathrm{U}_{1}\right) \varnothing \in \mathrm{U}$.
$\left(U_{2}\right) \forall X \in U,\{X\} \in U$.
$\left(\mathrm{U}_{3}\right) \forall \mathrm{X} \in \mathrm{U}, \forall \mathrm{Y} \in \mathrm{X}, \mathrm{Y} \in \mathrm{U}$.
$\left(U_{4}\right)$ For all sets $X$ and $Y$ belonging to $U$, the product set $X \times Y$ belongs to $U$.
$\left(\mathrm{U}_{5}\right)$ For any set X belonging to $\mathrm{U}, \mathfrak{P}(\mathrm{X}) \in \mathrm{U}$.
$\left(\mathrm{U}_{6}\right)$ If $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ is a family of sets

$$
\left(\mathrm{I} \in \mathrm{U} \text { and }(\forall i \in \mathrm{I}) \mathrm{X}_{i} \in \mathrm{U}\right) \Rightarrow \bigcup_{i \in \mathrm{I}} \mathrm{X}_{i} \in \mathrm{U}
$$

The following axiom can be added to the axioms of the theory of sets:

$$
\text { Axiom of universes: }(\forall X)(\exists \mathrm{U} \text { universe }) \mathrm{X} \in \mathrm{U} \text {. }
$$

It follows that, if U is a universe, the objects constructed from an object $\mathrm{X} \in \mathrm{U}$ using usual mathematical methods belong to U . For example, if E is a uniform space belonging to $U$, its completion $\widehat{\mathrm{E}}$ belongs to U . Indeed, each element of $\widehat{\mathrm{E}}$ is a set of filters on E , and so $\widehat{\mathrm{E}} \in \mathfrak{P}(\mathfrak{P}(\mathfrak{P}(\mathfrak{P}(\mathrm{E})))$ ). Likewise, let A be a commutative ring, $E$ and $F$ two A-modules. If $A, E$ and $F$ belong to $U$, so does $E \otimes_{A} F$, which as a quotient of $\mathrm{A}^{(\mathrm{E} \times \mathrm{F})}$ is an element of $\mathfrak{P}(\mathfrak{P}(\mathfrak{P}(\mathrm{E} \times \mathrm{F} \times \mathrm{A})))$. However, beware that U , $\mathscr{E} n \jmath_{\mathrm{U}}, \mathscr{E}_{\boldsymbol{\gamma}_{\mathrm{U}}}$, etc. are not elements of U.

As far as we know, as a result of adding the axiom of universes, the theory of sets may become inconsistent (assuming it is not already the case) since the converse has not been proved.

### 2.1.4

Definition Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories, $\mathscr{C}^{\prime}$ is said to be a subcategory of $\mathscr{C}$ if

1. $\mathscr{C}^{\prime} \subset \mathscr{C}$;
2. $\left(\forall \mathrm{X}, \mathrm{Y} \in \mathscr{C}^{\prime}\right) \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{X}, \mathrm{Y}) \subset \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})$;
3. Composition in $\mathscr{C}^{\prime}$ is inherited from composition in $\mathscr{C}$;
4. For any object X of $\mathscr{C}^{\prime}$, the identity element of X in $\mathscr{C}^{\prime}$ is the identity element of X in $\mathscr{C}$.
$\mathscr{C}^{\prime}$ is said to be a full subcategory of $\mathscr{C}$ if, for all $\mathrm{X}, \mathrm{Y} \in \mathscr{C}^{\prime}, \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{X}, \mathrm{Y})=$ $\operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})$.

### 2.1.5 Examples

(1) The category of metric spaces and continuous maps is a full subcategory of $\mathscr{T} 0$ р.
(2) The category of local rings and local homomorphisms is a (not full) subcategory of $\mathscr{A} n \boldsymbol{n}$. The category $\mathscr{A} n \boldsymbol{n}$ is a subcategory of the category of (not necessarily unitary or commutative) rings and homomorphisms.
(3) The category of complex Banach spaces and $\mathbb{C}$-linear continuous maps is a full subcategory of the category of topological vector spaces over $\mathbb{C}$ and $\mathbb{C}$-linear continuous maps.
(4) Let $\mathscr{C}$ be a category. A subcategory of $\mathscr{C}$ can be defined by taking its objects to be those of $\mathscr{C}$ and for its morphisms the isomorphisms of $\mathscr{C}$.

### 2.1.6 Definition

Let $\mathscr{C}^{\prime}$ and $\mathscr{C}^{\prime \prime}$ be two categories. The product category is the category $\mathscr{C}$ defined by $\mathscr{C}=\mathscr{C}^{\prime} \times \mathscr{C}^{\prime \prime}, \quad \operatorname{Hom}_{\mathscr{C}}\left(\left(\mathrm{X}^{\prime} \times \mathrm{X}^{\prime \prime}\right),\left(\mathrm{Y}^{\prime} \times \mathrm{Y}^{\prime \prime}\right)\right)=\operatorname{Hom}_{\mathscr{C}^{\prime}}\left(\mathrm{X}^{\prime}, \mathrm{Y}^{\prime}\right) \times \operatorname{Hom}_{\mathscr{C}^{\prime \prime}}\left(\mathrm{X}^{\prime \prime}, \mathrm{Y}^{\prime \prime}\right)$ and $\left(g^{\prime}, g^{\prime \prime}\right) \circ\left(f^{\prime}, f^{\prime \prime}\right)=\left(\left(g^{\prime} \circ f^{\prime}\right),\left(g^{\prime \prime} \circ f^{\prime \prime}\right)\right)$.

### 2.1.7 Definition

For any category $\mathscr{C}$, another category $\mathscr{C}^{\text {o }}$, called the opposite category, can be defined by $\mathscr{C}^{\circ}=\mathscr{C}, \operatorname{Hom}_{\mathscr{C}^{\circ}}(\mathrm{X}, \mathrm{Y})=\operatorname{Hom}_{\mathscr{C}}(\mathrm{Y}, \mathrm{X})$ and $g \circ \mathscr{C}^{\circ} f=f \circ \mathscr{C} g$.

## Exercises 2.1. (Categories)

1.-Show that a category $\mathscr{C}$ can be defined by taking rings for its objects, isomorphism classes of (A, B)-bimodules for morphisms from A to B and (M,N) $\mapsto$ $M \otimes_{B} N$ for composition, where $M$ is a (A, B)-bimodule and $N$ a ( $B, C$ )-bimodule.
2.-Let $\mathscr{C}$ be a category. An object $\mathrm{A} \in \mathscr{C}$ is initial (resp. terminal) if $\forall \mathrm{X} \in \mathscr{C}$, $\operatorname{Hom}_{\mathscr{C}}(\mathrm{A}, \mathrm{X})\left(\mathrm{resp} . \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{A})\right)$ is a singleton.

Give an initial object and a terminal one in each of the following categories: $\mathscr{E} \mathbf{n s}$, $\mathscr{G}_{\boldsymbol{\gamma}}, \mathscr{A} \boldsymbol{n n}, \mathscr{T} \boldsymbol{\sigma} \boldsymbol{p}, \mathrm{G}-\mathscr{E} n s, \mathrm{~A}-\mathscr{M} \boldsymbol{\sigma} \boldsymbol{d}$, and the category of Banach algebras over $\mathbb{C}$. Which of these categories contains an object which is both initial and terminal?
3.-Let $\alpha$ be an initial object and $\omega$ a terminal object in a category $\mathscr{C}$. Show that, if there is a morphism from $\omega$ to $\alpha$, it is unique and is an isomorphism.

### 2.2 Functors

### 2.2.1 Definition

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories. A covariant functor (or simply functor) from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ is a map F from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ such that, for all objects $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, there is a map $\mathrm{F}_{\mathrm{X}, \mathrm{Y}}: \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y}) \rightarrow \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{F}(\mathrm{X}), \mathrm{F}(\mathrm{Y}))$ satisfying the following conditions:
$\left(\mathrm{F}_{1}\right)(\forall \mathrm{X} \in \mathscr{C}) \mathrm{F}_{\mathrm{X}, \mathrm{X}}\left(1_{\mathrm{X}}\right)=1_{\mathrm{F}(\mathrm{X})}$.
$\left(\mathrm{F}_{2}\right)(\forall \mathrm{X}, \mathrm{Y}, \mathrm{Z} \in \mathscr{C})(\forall f: \mathrm{X} \rightarrow \mathrm{Y})(\forall g: \mathrm{Y} \rightarrow \mathrm{Z})$

$$
\mathrm{F}_{\mathrm{X}, \mathrm{Z}}(g \circ f)=\mathrm{F}_{\mathrm{Y}, \mathrm{Z}}(g) \circ \mathrm{F}_{\mathrm{X}, \mathrm{Y}}(f)
$$

A contravariant functor from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ is a map F from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ such that, for all $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, there is a map $\mathrm{F}_{\mathrm{X}, \mathrm{Y}}: \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y}) \rightarrow \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{F}(\mathrm{Y}), \mathrm{F}(\mathrm{X}))$ satisfying condition $\left(\mathrm{F}_{1}\right)$ and

$$
\begin{aligned}
\left(\mathrm{F}_{2}^{*}\right)(\forall \mathrm{X}, \mathrm{Y}, \mathrm{Z} \in \mathscr{C})(\forall f: \mathrm{X} \rightarrow \mathrm{Y})(\forall g: \mathrm{Y} \rightarrow \mathrm{Z}) \\
\mathrm{F}_{\mathrm{X}, \mathrm{Z}}(g \circ f)=\mathrm{F}_{\mathrm{X}, \mathrm{Y}}(f) \circ \mathrm{F}_{\mathrm{Y}, \mathrm{Z}}(g) .
\end{aligned}
$$

In other words, a contravariant functor from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ is a covariant functor from $\mathscr{C}^{\circ}$ to $\mathscr{C}^{\prime}$, or, a covariant functor from $\mathscr{C}$ to $\mathscr{C}^{\prime 0}$. The map $\mathrm{F}_{\mathrm{X}, \mathrm{Y}}$ is often written $f \mapsto f_{*}$ (resp. $f \mapsto f^{*}$ ) when F is a covariant (resp. contravariant) functor.

When defining a functor, often only F is given, leaving it to the reader to specify all the $\mathrm{F}_{\mathrm{X}, \mathrm{Y}}$.

### 2.2.2 Examples

(1) Forgetful functor. Consider the categories $\mathscr{A} n \boldsymbol{n}$ and $\mathscr{G} r$. Assigning to every ring, the underlying additive group, and to every ring homomorphism $f, f$ considered as a group homomorphism, defines a functor called the forgetful functor from $\mathscr{A} n \boldsymbol{n}$ to $\mathscr{E} r$.

Let $\mathscr{C}$ be the category of topological vector spaces and $\mathbb{R}$-linear continuous maps, the category $\mathscr{V} e c t_{\mathbb{R}}$ of $\mathbb{R}$ - vector spaces and linear maps. Define a functor from $\mathscr{C}$ to $\mathscr{V} e c t_{\mathbb{R}}$ by assigning to every object of $\mathscr{C}$ its underlying vector space.

Functors obtained in this manner are called forgetful functors. Thus there are forgetful functors $\mathscr{E}_{\boldsymbol{r}} \rightarrow \mathscr{E} n s, \mathscr{T} a p \rightarrow \mathscr{E} n s$, etc.
(2) Inclusion functor. Let $\mathscr{C}$ be a category, $\mathscr{C}^{\prime}$ a subcategory of $\mathscr{C}$. The canonical inclusions $\iota: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$ and $\iota_{\mathrm{X}, \mathrm{Y}}: \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{X}, \mathrm{Y}) \rightarrow \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})$ define a functor from $\mathscr{C}^{\prime}$ to $\mathscr{C}$.
(3) For any topological space X write $\mathscr{C}(\mathrm{X}, \mathbb{R})$ for the ring of continuous functions defined on X with values in $\mathbb{R}$. For any continuous map $f: \mathrm{X} \rightarrow \mathrm{Y}$ write $f^{*}$ for the map from $\mathscr{C}(\mathrm{Y}, \mathbb{R})$ to $\mathscr{C}(\mathrm{X}, \mathbb{R})$ defined by $f^{*}(h)=h \circ f$.

This defines a contravariant functor $\mathrm{X} \rightarrow \mathscr{C}(\mathrm{X}, \mathbb{R})$ from $\mathscr{T}$ ap to $\mathscr{A} \boldsymbol{n} n$.
(4) Duality functor. Let A be a commutative ring. Define a contravariant functor from the category $\mathrm{A}-\mathscr{M o d}$ to itself by assigning to every A -module X , its dual $\mathrm{X}^{\top}=\operatorname{Hom}(\mathrm{X}, \mathrm{A})$ and to every A-linear map its transpose.
(5) Bidual functor. Define a covariant functor from the category A- $\mathscr{M} \sigma d$ to itself by assigning to every A-module its bidual and to every A-linear map its bitranspose.
(6) Let $\mathscr{C}$ be a category, X an object of $\mathscr{C}$. Define a covariant functor $\overline{\mathrm{X}}: \mathscr{C} \rightarrow \mathscr{E}$ ns by setting

$$
\begin{aligned}
& (\forall \mathrm{Y} \in \mathscr{C}) \quad \check{\mathrm{X}}(\mathrm{Y})=\operatorname{Hom}(\mathrm{X}, \mathrm{Y}) \\
& (\forall \mathrm{Y}, \mathrm{Z} \in \mathscr{C})(\forall f \in \operatorname{Hom}(\mathrm{Y}, \mathrm{Z}))(\forall h \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y})) \quad f_{*}(h)=f \circ h .
\end{aligned}
$$

(7) Let $\mathscr{C}$ be a category, X an object of $\mathscr{C}$. Define a contravariant functor $\widehat{\mathrm{X}}: \mathscr{C} \rightarrow \mathscr{E} n s$ by setting

$$
\begin{aligned}
& (\forall \mathrm{Y} \in \mathscr{C}) \quad \widehat{\mathrm{X}}(\mathrm{Y})=\operatorname{Hom}(\mathrm{Y}, \mathrm{X}) \\
& (\forall \mathrm{Y}, \mathrm{Z} \in \mathscr{C})(\forall f \in \operatorname{Hom}(\mathrm{Y}, \mathrm{Z}))(\forall h \in \operatorname{Hom}(\mathrm{Z}, \mathrm{X})) \quad f^{*}(h)=h \circ f .
\end{aligned}
$$

### 2.2.3 Definition

Let $\mathscr{C}, \mathscr{C}^{\prime}, \mathscr{C}^{\prime \prime}$ be categories, $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ and $\mathrm{G}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}^{\prime \prime}$ functors. Define the composite functor $\mathrm{G} \circ \mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime \prime}$ as follows:

$$
\begin{aligned}
& (\forall \mathrm{X} \in \mathscr{C}) \mathrm{G} \circ \mathrm{~F}(\mathrm{X})=\mathrm{G}(\mathrm{~F}(\mathrm{X})) ; \\
& (\forall \mathrm{X}, \mathrm{Y} \in \mathscr{C})\left(\forall f \in \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})\right)(\mathrm{G} \circ \mathrm{~F})_{\mathrm{X}, \mathrm{Y}}(f)=\mathrm{G}_{\mathrm{F}(\mathrm{X}), \mathrm{F}(\mathrm{Y})}\left(\mathrm{F}_{\mathrm{X}, \mathrm{Y}}(f)\right) .
\end{aligned}
$$

The composite of two functors, one covariant and the other contravariant, as well as of two contravariant functors can be likewise defined. The former is contravariant, while the latter is covariant.

### 2.2.4 Examples

(1) The composite of the contravariant duality functor from A- $\mathscr{M} \sigma d$ to itself (2.2.2, Example 4) and of the forgetful functor from A- $\mathscr{M o d}$ to $\mathscr{E n s}$ (2.2.2, Example $1)$ is a contravariant functor $\widehat{\mathrm{A}}: \mathrm{A}-\mathscr{M} \boldsymbol{O} \rightarrow \mathscr{E}$ ns (2.2.2, Example 7).
(2) The contravariant biduality functor is the composition of the duality functor with itself.
(3) The composite of the contravariant functor $\mathrm{X} \mapsto \mathscr{C}(\mathrm{X}, \mathbb{R})$ from the category of compact spaces and continuous maps to the category $\mathscr{B}$ of Banach spaces over $\mathbb{R}$ and of the contravariant functor $\mathrm{E} \mapsto \mathrm{E}^{\top}$ from $\mathscr{B}$ to itself is a covariant functor
$\mathrm{X} \rightarrow \mathscr{M}(\mathrm{X})$. The space $\mathscr{M}(\mathrm{X})$ is the space of measures on X endowed with the norm topology, and for $f: \mathrm{X} \rightarrow \mathrm{Y}$ and $\mu \in \mathscr{M}(\mathrm{X}), f_{*}(\mu) \in \mathscr{M}(\mathrm{Y})$ is the direct image under $f$ of the measure $\mu$.

### 2.2.5 Category of Categories

Let U be a universe. Denote by $\mathscr{C} \mathfrak{a} \boldsymbol{t}_{\mathrm{U}}$ the category whose objects are categories $\mathscr{C}$ with the following property: the set $\mathscr{C}$ belongs to U and for all objects $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, $\operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y}) \in \mathrm{U}$. Morphisms from $\mathscr{C} \boldsymbol{a} \boldsymbol{t}_{\mathrm{U}}$ are covariant functors.

Note that $\mathscr{E} n \boldsymbol{s}_{U}$ does not belong to $\mathscr{C} \boldsymbol{a} \boldsymbol{t}_{\mathrm{U}}$; however if $\mathrm{U}^{\prime}$ is a universe such that $\mathrm{U} \in \mathrm{U}^{\prime}$, then $\mathscr{E} n s_{\mathrm{U}} \in \mathscr{C} \boldsymbol{a} t_{\mathrm{U}^{\prime}}$.

### 2.2.6 Another Example of a Functor

Let $f: \mathrm{G} \rightarrow \mathrm{H}$ be a group homomorphism and X a H -set.
Denote by $f^{*} \mathrm{X}$ the G -set obtained by endowing the set X with the action $(g, x) \mapsto f(g) \cdot x$. This defines a functor $f^{*}$ from H- $\mathscr{E} n s$ to G- $\mathscr{E} n s$. More precisely, if U is a universe, then $f^{*}$ is a functor from $\mathrm{H}-\mathscr{E} n s_{\mathrm{U}}$ to $\mathrm{G}-\mathscr{E} n s_{\mathrm{U}}$.

Let U and $\mathrm{U}^{\prime}$ be two universes such that $\mathrm{U} \in \mathrm{U}^{\prime}, \mathrm{G} \rightarrow \mathrm{G}$ - $\mathscr{E n s} \boldsymbol{s}_{\mathrm{U}}$ from $\mathscr{E} \mathscr{r}_{\mathrm{U}}$ to $\mathscr{C} \boldsymbol{a} \boldsymbol{t}_{U^{\prime}}$ is a contravariant functor.

### 2.2.7

Proposition Let G and H be two groups, $\omega_{\mathrm{G}}: \mathrm{G}-\mathscr{E} n s \rightarrow \mathscr{E} n s$ and $\omega_{\mathrm{H}}: \mathrm{H}-\mathscr{E} n s \rightarrow \mathscr{E} n s$ the forgetful functors. Let $\Phi: \mathrm{H}-\mathscr{E} n s \rightarrow \mathrm{G}-\mathscr{E} n s$ be a functor such that the diagram

commutes. Then there is a unique homomorphism $f: \mathrm{G} \rightarrow \mathrm{H}$ such that $\Phi=f^{*}$.
Proof For any H-set X, its image $\Phi(\mathrm{X})$ is the set X endowed with an action $(g, x) \mapsto g \perp x$ from $\mathrm{G} \times \mathrm{X}$ to X . If $\Phi=f^{*}$, then $g \perp x=f(g) \cdot x$. In particular for $\mathrm{X}=\mathrm{H}$ endowed with left translations, $f(g)=g \perp e$, hence uniqueness.

Define $f: \mathrm{G} \rightarrow \mathrm{H}$ by $f(g)=g \perp e$. Let X be a H -set, then for all $g \in \mathrm{G}$ and all $x \in \mathrm{X}, g \perp x=f(g) x$. Indeed, the map $\delta_{x}: h \mapsto h \cdot x$ from H to X is a H-morphism, its image $\Phi\left(\delta_{x}\right)$ is a G-morphism, but the maps $\Phi\left(\delta_{x}\right)$ and $\delta_{x}$ agree, and
so $g \perp x=g \perp \delta_{x}(e)=\delta_{x}(g \perp e)=\delta_{x}(f(g))=f(g) \cdot x$. In particular, for $\mathrm{X}=\mathrm{H}$ and $x=f\left(g^{\prime}\right)$,

$$
f(g) \cdot f\left(g^{\prime}\right)=g \perp f\left(g^{\prime}\right)=g \perp\left(g^{\prime} \perp e\right)=g g^{\prime} \perp e=f\left(g g^{\prime}\right)
$$

and hence $f$ is a homomorphism, and $\Phi=f^{*}$.

### 2.3 Morphisms of Functors

### 2.3.1 Definition

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories, F and G two functors from $\mathscr{C}$ to $\mathscr{C}^{\prime}$. A functorial morphism or morphism of functors $\Phi$ from $F$ to $G$, denoted by $\Phi: F \rightarrow G$, is defined by assigning to each object $X$ of $\mathscr{C}$ a morphism $\Phi_{X}: F(X) \rightarrow G(X)$ such that, for any morphism $f: X \rightarrow Y$ of $\mathscr{C}$ the following diagram commutes:


If $F$ and $G$ are contravariant functors from $\mathscr{C}$ to $\mathscr{C}^{\prime}$, define the morphisms from F to G by considering F and G as covariant functors from $\mathscr{C}^{\circ}$ to $\mathscr{C}^{\prime}$ (not from $\mathscr{C}$ to $\left.\mathscr{C}^{\prime \mathrm{O}}\right)$.

### 2.3.2 Examples

(1) Biduality morphism. Let $\mathscr{C}=\mathscr{C}^{\prime}=\mathrm{A}-\mathscr{M} \boldsymbol{O}$. Write $1_{\mathscr{C}}$ for the identity functor of $\mathscr{C}$ and G for the biduality functor. Let X be an A-module, $x \in \mathrm{X}$ and $\delta_{x}: \mathrm{X}^{\top} \rightarrow \mathrm{A}$ the linear form $h \mapsto h(x)$. The map $\delta_{\mathrm{X}}: x \mapsto \delta_{x}$ from X to $\mathrm{X}^{\top \top}$ is A-linear, and so is a morphism of A- $\mathscr{M} \boldsymbol{O} \boldsymbol{d}$. Define a functorial morphism from $1_{\mathscr{C}}$ to G by assigning to each A -module X the morphism $\delta_{\mathrm{X}}$.
(2) Kronecker morphism. Let $\mathscr{C}$ be the category $(\mathrm{A}-\mathscr{M} \boldsymbol{\sigma})^{\circ} \times \mathrm{A}-\mathscr{M} \boldsymbol{\sigma}$. For every object $(X, Y)$ of $\mathscr{C}$, set $F(X, Y)=X^{\top} \otimes Y$ and $G(X, Y)=\operatorname{Hom}(X, Y)$. This defines two functors $F$ and $G$ from $\mathscr{C}$ to $A-\mathscr{M} \neq d$. For every object $(X, Y)$ of $\mathscr{C}$, the Kronecker morphism $\mathrm{K}_{\mathrm{X}, \mathrm{Y}}$ from $\mathrm{X}^{\top} \otimes \mathrm{Y}$ to $\operatorname{Hom}(\mathrm{X}, \mathrm{Y})$ defined by $\mathrm{K}(\alpha \otimes y)(x)=$ $\alpha(x) \cdot y$ is a morphism from $\mathrm{F}(\mathrm{X}, \mathrm{Y})$ to $\mathrm{G}(\mathrm{X}, \mathrm{Y})$. The family $\left(\mathrm{K}_{\mathrm{X}, \mathrm{Y}}\right)_{(\mathrm{X}, \mathrm{Y}) \in \mathscr{C}}$ is a functorial morphism from $F$ to $G$. This will be checked in (3.9).

### 2.3.3 Functor Category

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories. Construct a category by taking as its objects the functors from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ and for its morphisms the morphisms of functors. If $\Phi: F \rightarrow G$ and $\Psi: \mathrm{G} \rightarrow \mathrm{H}$ are morphisms of functors, define the composition $\Psi \circ \Phi: \mathrm{F} \rightarrow \mathrm{H}$ by $(\Psi \circ \Phi)_{\mathrm{X}}=\Psi_{\mathrm{X}} \circ \Phi_{\mathrm{X}}$ for any object X of $\mathscr{C}$.

A morphism of functors $\Phi: \mathrm{F} \rightarrow \mathrm{G}$ is an isomorphism if there is a morphism $\Psi: \mathrm{G} \rightarrow \mathrm{F}$ such that $\Psi \circ \Phi$ is the identity morphism $1_{\mathrm{F}}$ of F and $\Phi \circ \Psi=1_{\mathrm{G}}$.

A morphism of functors $\Phi$ is an isomorphism if and only if for any object $\mathrm{X} \in \mathscr{C}$, the morphism $\Phi_{X}: F(X) \rightarrow G(X)$ is an isomorphism. Indeed the morphisms $\Phi_{X}^{-1}$ then define a morphism from $G$ to $F$.

Remark $\mathrm{F}(\mathrm{X})$ may be isomorphic to $\mathrm{G}(\mathrm{X})$ for any X even when the functors F and $G$ are not (2.3, Exercise 1).

### 2.3.4 Definitions

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories, $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ a functor.

1. The functor F is essentially surjective if for any $\mathrm{X}^{\prime} \in \mathscr{C}^{\prime}$, there exists $\mathrm{X} \in \mathscr{C}$ such that $X^{\prime}$ is isomorphic to $F(X)$; in other words if the isomorphism induced by $F$ from the set of isomorphism classes of objects of $\mathscr{C}$ to the set of isomorphism classes of objects of $\mathscr{C}^{\prime}$ is surjective.
2. The functor F is said to be faithful (resp. fully faithful) if, for all $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, the $\operatorname{map} \mathrm{F}_{\mathrm{X}, \mathrm{Y}}: \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y}) \rightarrow \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{F}(\mathrm{X}), \mathrm{F}(\mathrm{Y}))$ is injective (resp. bijective).

### 2.3.5 Equivalence of Categories

Theorem and Definition Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories, $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ a functor. The following conditions are equivalent:
(i) F is fully faithful and essentially surjective;
(ii) there is a functor $\mathrm{G}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$ such that $\mathrm{G} \circ \mathrm{F}$ is isomorphic to $1_{\mathscr{C}}$ and $\mathrm{F} \circ \mathrm{G}$ is isomorphic to $1_{\mathscr{C}^{\prime}}$.

The functor F is said to be an equivalence of categories if it satisfies one of the two conditions.

A functor G satisfying (ii) is called a quasi-inverse of F .
Proof (i) $\Rightarrow$ (ii). Since F is essentially surjective, by the axiom of choice, there is a map $\mathrm{G}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$ and a family $\left(\alpha_{\mathrm{X}^{\prime}}\right)_{\mathrm{X}^{\prime} \in \mathscr{C}^{\prime}}$ such that, for all $\mathrm{X}^{\prime}$, the morphism $\alpha_{\mathrm{X}^{\prime}}$ is an isomorphism from $\mathrm{F}\left(\mathrm{G}\left(\mathrm{X}^{\prime}\right)\right)$ onto $\mathrm{X}^{\prime}$. For any morphism $f: \mathrm{X}^{\prime} \rightarrow \mathrm{Y}^{\prime}$ of $\mathscr{C}^{\prime}$, set $\mathrm{G}_{\mathrm{X}^{\prime}, \mathrm{Y}^{\prime}}(f)=\mathrm{F}_{\mathrm{X}, \mathrm{Y}}^{-1}\left(\alpha_{\mathrm{Y}^{\prime}}^{-1} \circ f \circ \alpha_{\mathrm{X}^{\prime}}\right)$, where $\mathrm{X}=\mathrm{G}\left(\mathrm{X}^{\prime}\right)$ and $\mathrm{Y}=\mathrm{G}\left(\mathrm{Y}^{\prime}\right)$, which
is possible since $F$ is fully faithful. The functor $G$ together with $\left(G_{X^{\prime}, Y^{\prime}}\right)$ can be checked to be a functor from $\mathscr{C}^{\prime}$ to $\mathscr{C}$ satisfying the conditions of (ii).
(ii) $\Rightarrow$ (i). Since $\mathrm{F} \circ \mathrm{G}$ is isomorphic to $1_{\mathscr{C}^{\prime}}$, every object $\mathrm{X}^{\prime}$ of $\mathscr{C}^{\prime}$ is isomorphic to $F\left(G\left(X^{\prime}\right)\right)$. Hence $F$ is essentially surjective.

We show that F is fully faithful. Let $\alpha=\left(\alpha_{\mathrm{X}}\right)_{\mathrm{X} \in \mathscr{C}}$ be an isomorphism from $\mathrm{G} \circ \mathrm{F}$ onto $1_{\mathscr{C}}$ and $\beta=\left(\beta_{\mathrm{X}^{\prime}}\right)_{\mathrm{X}^{\prime} \in \mathscr{C}^{\prime}}$ an isomorphism from $\mathrm{F} \circ \mathrm{G}$ onto $1_{\mathscr{C}^{\prime}}$. Let $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$; set $\mathrm{X}^{\prime}=\mathrm{F}(\mathrm{X}), \mathrm{Y}^{\prime}=\mathrm{F}(\mathrm{Y}), \mathrm{X}^{\prime \prime}=\mathrm{G}\left(\mathrm{X}^{\prime}\right), \mathrm{Y}^{\prime \prime}=\mathrm{G}\left(\mathrm{Y}^{\prime}\right), \varphi_{\mathrm{X}, \mathrm{Y}}(f)=\alpha_{\mathrm{Y}}^{-1} \circ f \circ \alpha_{\mathrm{X}}$ for $f \in \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{Y})$ and $\psi_{\mathrm{X}^{\prime}, \mathrm{Y}^{\prime}}(g)=\beta_{\mathrm{Y}^{\prime}}^{-1} \circ g \circ \beta_{\mathrm{X}^{\prime}}$ for $g \in \operatorname{Hom}_{\mathscr{C}^{\prime}}\left(\mathrm{X}^{\prime} \mathrm{Y}^{\prime}\right)$.

The following diagram commutes:


The maps $\varphi_{X, Y}$ and $\psi_{X^{\prime}, Y^{\prime}}$ are bijective. Hence $G_{X^{\prime}, Y^{\prime}}$ is surjective and injective and consequently, $\mathrm{F}_{\mathrm{X}, \mathrm{Y}}=\mathrm{G}_{\mathrm{X}^{\prime}, \mathrm{Y}^{\prime}}^{-1} \circ \varphi_{\mathrm{X}, \mathrm{Y}}$ is itself bijective.

An anti-equivalence of categories is a contravariant factor $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ which is an equivalence from $\mathscr{C}^{\circ}$ to $\mathscr{C}^{\prime}$.

### 2.3.6

Proposition Let $f: \mathrm{G} \rightarrow \mathrm{H}$ be a group homomorphism. The functor $f^{*}: \mathrm{H}-\mathscr{E} n s \rightarrow$ G- $\mathscr{n} s$ is an equivalence of categories if and only if $f$ is an isomorphism.

Proof (a) Sufficiency. $\left(f^{-1}\right)^{*}$ is an inverse of $f^{*}$.
(b) $f^{*}$ is essentially surjective $\Rightarrow f$ injective. If $f(g)=e$, then $g$ acts trivially on any G-set $f^{*} \mathrm{Y}$, hence on any G-set, in particular on G itself by left translations, and so $g=e$.
(c) $f^{*}$ fully faithful $\Rightarrow f$ surjective. Let $\mathrm{X} \in \mathrm{H}-\mathscr{E} n s$ be a singleton with H acting trivially. For all $\mathrm{Y} \in \mathrm{H}-\mathscr{E}$ ns, $\operatorname{Hom}_{\mathrm{H}}(\mathrm{X}, \mathrm{Y})$ can be identified with the set of fixed points of Y under the action of H . Take Y to be the set $\mathrm{H} / f(\mathrm{G})$ together with the action of H induced by left translations. Then $\operatorname{Hom}_{\mathrm{H}}(\mathrm{X}, \mathrm{Y})$ is empty unless $f$ is surjective, but $\operatorname{Hom}_{\mathrm{G}}\left(f^{*} \mathrm{X}, f^{*} \mathrm{Y}\right)$ is not empty, since the image of the identity element in Y is kept fixed by G and so $f$ is necessarily surjective.

### 2.3.7

Let U and $\mathrm{U}^{\prime}$ be two sets and $\mathrm{U} \subset \mathrm{U}^{\prime}$. Suppose U is hereditary, i.e. for any set $\mathrm{X} \in \mathrm{U}$, $\mathfrak{P}(X) \subset U$ (a far weaker condition than the existence of universes, which supposes $\mathfrak{P}(\mathrm{X}) \in \mathrm{U})$.

Let $\mathscr{C}$ be a category whose set of objects is contained in $\mathrm{U}^{\prime}$. Suppose that $\mathscr{C}$ is a category with underlying sets and transfer of structure (in $\mathrm{U}^{\prime}$ ). This means there is a functor (the forgetful functor) $\Omega: \mathscr{C} \rightarrow \mathscr{E} n s_{\mathrm{U}}$, and that for any object X of $\mathscr{C}$ and any bijection $f: \Omega(\mathrm{X}) \rightarrow \mathrm{F}$ with $\mathrm{F} \in \mathrm{U}^{\prime}$, there is an object Y of $\mathscr{C}$ and a morphism $\varphi: \mathrm{X} \rightarrow \mathrm{Y}$ such that $\Omega(\mathrm{Y})=\mathrm{F}$ and $\Omega(\varphi)=f$ (this is the case for $\mathscr{E} \boldsymbol{n s}_{\mathrm{U}^{\prime}}, \mathscr{E}_{\boldsymbol{r}_{\mathrm{U}^{\prime}}}$, $\mathscr{T} \sigma \mathcal{R}_{U^{\prime}}, \mathrm{A}-\mathscr{M} \sigma d_{\mathrm{U}^{\prime}}$, etc.).

Denote by $\mathscr{C}_{\mathrm{U}}$ the full subcategory of $\mathscr{C}$ consisting of objects whose underlying sets belong to U . Suppose there is a set $\mathrm{E} \in \mathrm{U}$ such that, for any object X of $\mathscr{C}$, $\operatorname{Card} \Omega(\mathrm{X}) \leq \operatorname{Card} \mathrm{E}$.

Then the inclusion morphism $\mathscr{C}_{\mathrm{U}} \rightarrow \mathscr{C}$ is an equivalence of categories.
Indeed it is fully faithful since it is the inclusion morphism of a full subcategory, and essentially surjective since any object of $\mathscr{C}$ is isomorphic to an object whose underlying set is a subset of $E$.

Thanks to the lemmas about cardinalities (4.4.7, 5.2.5), this is why, up to equivalence, the category of coverings of a given space, or that of algebraic extensions of a given field, do not depend on the choice of the universe.

## Exercises 2.3. (Morphisms of functors)

1.-Let $\mathscr{C}$ be the category of finite dimensional vector spaces over a field $k$ and of isomorphisms. Let F be the identity functor of $\mathscr{C}$ and $\mathrm{G}: \mathscr{C} \rightarrow \mathscr{C}$ the functor defined by $\mathrm{G}(\mathrm{X})=\mathrm{X}$ and $\mathrm{G}(f)=\left(f^{\top}\right)^{-1}$. Show that, for any object X of $\mathscr{C}$, the objects $F(X)$ and $G(X)$ are isomorphic but that the functors $F$ are $G$ not.
2.-Show that the set of endomorphisms of the identity functor of the category of abelian groups can be identified with $\mathbb{Z}$.
3.-Is there a map $f: \mathbb{N} \rightarrow \mathbb{N}$ such that

$$
\forall x \in \mathbb{N}, \forall y \in \mathbb{N} \quad f(y)^{f(x)}=x^{y} ?
$$

Is the category $\operatorname{Setf}$ of finite set its opposite category?
What about the category $\mathscr{E} n s$.
4.-Let $\mathscr{C}$ be the category of finite commutative groups. The aim is to show that $\mathscr{C}$ is equivalent to its opposite category. For this, assume that every object of $\mathscr{C}$ is a product of cyclic groups (3.5.8, Corollary 3 ).
(a) Show that for any group $\mathrm{X} \in \mathscr{C}$, the group $\mathrm{X}^{\vee}=\operatorname{Hom}_{\mathscr{G}_{\boldsymbol{\gamma}}}(\mathrm{X}, \mathbb{Q} / \mathbb{Z})$ is finite and isomorphic to X . (First consider the case where X is a cyclic group.)
(b) Show that the natural map $\mathrm{X} \mapsto \mathrm{X}^{\vee \vee}$ is an isomorphism. Conclude.
(c) Let $\mathscr{C}^{\prime}$ be the category of objects of $\mathscr{C}$ and of isomorphisms. Show that the identity functor $1_{\mathscr{C}^{\prime}}$ and the functor $\mathrm{X} \mapsto \mathrm{X}^{\vee}$ from $\mathscr{C}^{\prime}$ to itself are not isomorphic (apply Exercise 1 to $k=\mathbb{Z} /(p)$ ).
5.-The category of locally compact commutative groups can be shown to be equivalent to its opposite category (Pontryagin duality, the natural framework of the Fourier transform). The aim is to show that there is a subcategory for which this is the case.

Let $\mathscr{C}$ be the category of commutative topological groups G satisfying the following properties:
(i) the connected component $\mathrm{G}_{0}$ of the identity element of G is isomorphic to the quotient of a real finite dimensional vector space by a discrete subgroup;
(ii) $G_{0}$ is open and the quotient $G / G_{0}$ is a finitely generated $\mathbb{Z}$-module.
(a) Show that every object of $\mathscr{C}$ is a finite product of groups of the following type: $\mathbb{R}, \mathbb{R} / \mathbb{Z}, \mathbb{Z}, \mathbb{Z} /(n)$.
(b) Let $\mathrm{G} \in \mathscr{C}$. Show that $\mathrm{G}^{\vee}=\operatorname{Hom}_{\mathscr{C}}(\mathrm{G}, \mathbb{R} / \mathbb{Z})$ equipped with the topology of compact convergence (2.4.9, Example 2 ) is an object of $\mathscr{C}$.

Determine $\mathrm{G}^{\vee}$ where G is successively $\mathbb{R}, \mathbb{R} / \mathbb{Z}, \mathbb{Z}$ and $\mathbb{Z} /(n)$. Show that $\mathrm{G}^{\vee \vee}=\mathrm{G}$. Conclude.
6.-Let U and $\mathrm{U}^{\prime}$ be two universes with $\mathrm{U} \in \mathrm{U}^{\prime}$. Show that the categories $\mathscr{E} n s_{\mathrm{U}}$ and $\mathscr{E} \boldsymbol{n s}_{U^{\prime}}$ are not equivalent.

Denote by $\mathbb{R}$-vectfin ${ }_{U}$ the category of finite dimensional vector spaces over $\mathbb{R}$ belonging to U . Show that $\mathbb{R}$-vectfin $\mathrm{U}_{\mathrm{U}}$ and $\mathbb{R}$-vectfin $\mathrm{U}^{\prime}$ are equivalent.

### 2.4 Representable Functors

### 2.4.1 Notation

For each object X of a category $\mathscr{C}$, set $\widehat{\mathrm{X}}$ to be the contravariant functor from $\mathscr{C}$ to $\mathscr{E}$ ns defined by $\widehat{\mathrm{X}}(\mathrm{Y})=\operatorname{Hom}(\mathrm{Y}, \mathrm{X})$.

### 2.4.2 Definition

Let $\mathscr{C}$ be a category and F a contravariant functor from $\mathscr{C}$ to $\mathscr{E}$ ns. Then F is said to be representable if there is an object X of $\mathscr{C}$ such that F is isomorphic to $\widehat{\mathrm{X}}$.

### 2.4.3

The previous definition implicitly involves an isomorphism from $\widehat{X}$ onto $F$. We next give a description of the morphisms from $\widehat{X}$ to $F$.

Scholium. Let $\mathscr{C}$ be a category, F a contravariant functor from $\mathscr{C}$ to $\mathscr{E}$ ns and X an object of $\mathscr{C}$. The morphisms from $\widehat{\mathrm{X}}$ to F are in bijective correspondence with the elements of $\mathrm{F}(\mathrm{X})$.

More precisely, for $\xi \in \mathrm{F}(\mathrm{X}), \mathrm{Y} \in \mathscr{C}$, and $f \in \operatorname{Hom}(\mathrm{Y}, \mathrm{X})$, write $\hat{\xi}_{\mathrm{Y}}(f)$ for the element $f^{*}(\xi)$ of $\mathrm{F}(\mathrm{Y})$. This defines a map $\hat{\xi}_{\mathrm{Y}}$ from $\widehat{\mathrm{X}}(\mathrm{Y})$ to $\mathrm{F}(\mathrm{Y})$ and $\hat{\xi}=\left(\hat{\xi}_{\mathrm{Y}}\right)_{\mathrm{Y} \in \mathscr{C}}$ is a morphism of functors from $\widehat{X}$ to $F$.

Proposition (Yoneda Lemma) The map $\xi \mapsto \hat{\xi}$ is a bijection from $\mathrm{F}(\mathrm{X})$ onto the set of morphisms from $\widehat{\mathrm{X}}$ to F .

Proof We first show that $\hat{\xi}$ is a morphism of functors. Consider the following diagram:

where $\mathrm{Y}, \mathrm{Z} \in \mathscr{C}$ and $g \in \operatorname{Hom}(\mathrm{Y}, \mathrm{Z})$. Then, $\left(\hat{\xi}_{\mathrm{Y}} \circ \widehat{\mathrm{X}}(g)\right)(\alpha)=\mathrm{F}(\alpha \circ g)(\xi)$ and $\mathrm{F}(g)\left(\hat{\xi}_{\mathrm{Z}}(\alpha)\right)=\mathrm{F}(g) \circ \mathrm{F}(\alpha)(\xi)=\mathrm{F}(\alpha \circ g)(\xi)$. So the diagram commutes.

For any morphism $\Phi$ from $\widehat{\mathrm{X}}$ to F , set $\eta(\Phi)=\Phi_{\mathrm{X}}\left(1_{\mathrm{X}}\right) \in \mathrm{F}(\mathrm{X})$.
We show that the maps $\xi \mapsto \hat{\xi}$ and $\Phi \mapsto \eta(\Phi)$ are mutual inverses. $\eta(\hat{\xi})=$ $\hat{\xi}_{\mathrm{X}}\left(1_{\mathrm{X}}\right)=\xi$ and, for all $\mathrm{Y} \in \mathscr{C}$ and all $f \in \operatorname{Hom}(\mathrm{Y}, \mathrm{X})$

$$
\widehat{\eta(\Phi)}_{\mathrm{Y}}(f)=\mathrm{F}(f)(\eta(\Phi))=\mathrm{F}(f)\left(\Phi_{\mathrm{X}}\left(1_{\mathrm{X}}\right)\right)=\Phi_{\mathrm{Y}} \circ \widehat{\mathrm{X}}(f)\left(1_{\mathrm{X}}\right)=\Phi_{\mathrm{Y}}(f)
$$

### 2.4.4 Definition

Keeping the notation of 2.4 .3 , the functor F is said to be represented by $(\mathrm{X}, \xi)$ if $\hat{\xi}$ is an isomorphism from $\widehat{X}$ onto $F$. In other words, $F$ is represented by $(X, \xi)$ if and only if

$$
(\forall \mathrm{Y} \in \mathscr{C})(\forall \eta \in \mathrm{F}(\mathrm{Y})) \quad(\exists \mid f: \mathrm{Y} \rightarrow \mathrm{X}) f^{*}(\xi)=\eta
$$

### 2.4.5 Comments. Universal Problems

Choosing a pair $(\mathrm{X}, \xi)$ representing F amounts to identifying $\mathrm{F}(\mathrm{Y})$ with $\widehat{\mathrm{X}}(\mathrm{Y})$. Hence for simplicity's sake, for any object Y of $\mathscr{\mathscr { C }}$, one writes $\operatorname{Hom}(\mathrm{Y}, \mathrm{X})=\mathrm{F}(\mathrm{Y})$. Such a relation is called a universal property of X .

The search for a pair $(\mathrm{X}, \xi)$ representing F is called a universal problem defined by the functor F .

### 2.4.6 Definition

A covariant functor F from a category $\mathscr{C}$ to $\mathscr{E} n s$ is representable if F is a representable functor from $\mathscr{C}^{\circ}$ to $\mathscr{E}$ ns .

### 2.4.7 Example of a Universal Problem

Tensor product. Let A be a commutative ring, E and F A-modules. We want to find an A-module T such that, for any A-module $G$, the set of bilinear maps from $\mathrm{E} \times \mathrm{F}$ to $G$ can be identified with the set of linear maps from $T$ to $G$.

More precisely, we want to find an A-module T and a bilinear map $\theta: \mathrm{E} \times \mathrm{F} \rightarrow \mathrm{T}$ such that, for any A-module G and any bilinear map $f: \mathrm{E} \times \mathrm{F} \rightarrow \mathrm{G}$, there is a unique linear map $g: \mathrm{T} \rightarrow \mathrm{G}$ such that $f=g \circ \theta$.

Put differently, let B be the covariant functor from A- $\mathscr{M} O \boldsymbol{d}$ to $\mathscr{E} n s$ which assigns to a module G the set of bilinear maps from $\mathrm{E} \times \mathrm{F}$ to G . We want to find a pair ( $\mathrm{T}, \theta$ ) where T is an A-module and $\theta \in \mathrm{B}(\mathrm{T})$ such that, for any A-module G the map $g \mapsto g^{*}(\theta)=g \circ \theta$ from $\operatorname{Hom}(\mathbf{T}, \mathrm{G})$ to $\mathbf{B}(\mathrm{G})$ is bijective. In other words, we are trying to represent the covariant functor B .

This universal problem always has a solution: the tensor product $\mathrm{E} \otimes \mathrm{F}$ of the A-modules E and F (3.8.1).

### 2.4.8 Other Examples of Representable Covariant Functors

(1) Quotient. Let X be a set, R an equivalence relation on X . We want to find a set Y and a map $\chi: \mathrm{X} \rightarrow \mathrm{Y}$ such that
(i) $\left(\forall x, x^{\prime} \in \mathrm{X}\right) \quad x \mathrm{R} x^{\prime} \Rightarrow \chi(x)=\chi\left(x^{\prime}\right)$;
(ii) for any set Z and any map $f: \mathrm{X} \rightarrow \mathrm{Z}$ such that $x \mathrm{R} x^{\prime} \Rightarrow f(x)=f\left(x^{\prime}\right)$, there is a unique map $g: \mathrm{Y} \rightarrow \mathrm{Z}$ satisfying $f=g \circ \chi$.

This universal problem has a solution: the quotient set $Y=X / R$ endowed with the canonical quotient map $\chi: \mathrm{X} \rightarrow \mathrm{X} / \mathrm{R}$.

In other words, the quotient represents the covariant functor which assigns to a set Z the set of maps $f: \mathrm{X} \rightarrow \mathrm{Z}$ for which $x \mathrm{R} x^{\prime} \Rightarrow f(x)=f\left(x^{\prime}\right)$.
(2) Quotient modules. Let A be a commutative ring, E an A-module, and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of E . Denote by $\Phi(\mathrm{G})$ the set of linear maps from E to an A-module

G which vanish at $x_{i}$ for all $i \in \mathrm{I}$. The functor $\mathrm{G} \mapsto \Phi(\mathrm{G})$ from $\mathrm{A}-\mathscr{M} \circ \boldsymbol{d}$ to $\mathscr{E} n \boldsymbol{n}$ is representable. The pair $(\mathrm{F}, \chi)$, where F is the quotient module of E by the submodule generated by all $x_{i}$ and $\chi: \mathrm{E} \rightarrow \mathrm{F}$ the canonical quotient map, represents $\Phi$.
(3) Free module. Let I be a set. We want to find an A-module L such that for any A-module F,

$$
\operatorname{Hom}(L, F)=F^{I}
$$

where $F^{I}$ is the set of maps from $I$ to $F$.
This universal problem has a solution: the free A-module $\mathrm{A}^{(\mathrm{I})}$, constructed over I (3.4.6).
(4) Exterior products. Let E be an A-module, $p$ a positive integer. We want to find an A-module T such that, for any A-module F, the $\operatorname{Hom}(T, F)$ can be identified to the set of $p$-linear alternating maps from $\mathrm{E}^{p}$ to F .

This universal problem has a solution: the $p$ th exterior power ${ }^{1} \bigwedge^{p}(\mathrm{E})$.
(5) Completion. Let E be a normed vector space. We want to find a Banach space $\widehat{\mathrm{E}}$ and an embedding $\iota: \mathrm{E} \rightarrow \widehat{\mathrm{E}}$ such that, for any Banach space F , every continuous linear map from E to F has a unique continuous extension from $\widehat{\mathrm{E}}$ to F .

In other words, we want to represent the functor $\Phi$ defined on the category of Banach spaces and continuous linear maps with values in $\mathscr{E}$ ns by $\Phi(\mathrm{F})=$ $\{f: \mathrm{E} \rightarrow \mathrm{F} f$ linear continuous $\}$.

The completion of E , the space $\widehat{\mathrm{E}}$, endowed with the canonical embedding from $E$ into $\widehat{E}$ is the desired solution.
(6) Action of a group on itself. Let $G$ be a group. The forgetful functor $\mathrm{G}-\mathscr{E} n \boldsymbol{n} \rightarrow \mathscr{E} n \boldsymbol{n}$ is represented by ( ${ }^{\circ} \mathrm{G}, e$ ), where ${ }^{\circ} \mathrm{G}$ is the set G on which the group G acts by left translations, and $e$ is the identity element of G .

### 2.4.9 Examples of Representable Contravariant Functors

(1) Induced topology. Let $X$ be a topological space, $Y$ a subset of $X$. We want to find a topological space $\mathrm{Y}^{\prime}$ such that, for any topological spaces T , the set of continuous maps from T to $\mathrm{Y}^{\prime}$ can be identified with the set of continuous maps from T to X whose images are contained in Y. In other words, we want to represent the functor $F$ from $\mathscr{T} a p$ to $\mathscr{E} n s$ defined by

$$
\mathrm{F}(\mathrm{~T})=\{f \in \operatorname{Hom}(\mathrm{~T}, \mathrm{X}) \mid f(\mathrm{~T}) \subset \mathrm{Y}\} .
$$

This universal problem has a solution: the set Y equipped with the topology induced by that of X.

[^0](2) Topology of compact convergence. Let X and Y be two topological spaces. We want to find a topological space $F$ such that, for any topological space $T$,
$$
\operatorname{Hom}_{\mathscr{\sigma _ { \boldsymbol { p } }}}(\mathrm{T} \times \mathrm{X}, \mathrm{Y})=\operatorname{Hom}_{\mathscr{T}_{\bullet p}}(\mathrm{~T}, \mathrm{~F}) .
$$
 $\operatorname{Hom}_{\mathscr{T a p}}(\mathrm{X}, \mathrm{Y})$ endowed with a topology called the topology of compact convergence, more precisely by ( $\mathrm{F}, \varphi$ ) where, for all $f \in \mathrm{~F}$ and all $x \in \mathrm{X}, \varphi(f, x)=f(x)$.

### 2.4.10

Let $\mathscr{C}$ be a category, F a contravariant functor from $\mathscr{C}$ to $\mathscr{E} n s,(\mathrm{X}, \xi)$ and $\left(\mathrm{X}^{\prime}, \xi^{\prime}\right)$ two pairs representing F . There is a unique morphism $h: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ such that $h^{*} \xi^{\prime}=\xi$.

Definition The morphism $h$ is the canonical morphism from X onto $\mathrm{X}^{\prime}$.

### 2.4.11 Uniqueness of the Solution of the Universal Problem

Proposition keeping the above notation,
(a) the canonical morphism $h$ is an isomorphism;
(b) let $(\mathrm{X}, \xi),\left(\mathrm{X}^{\prime}, \xi^{\prime}\right)\left(\mathrm{X}^{\prime \prime}, \xi^{\prime \prime}\right)$ be three pairs representing F , and $h: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$, $h^{\prime}: \mathrm{X}^{\prime} \rightarrow \mathrm{X}^{\prime \prime}, h^{\prime \prime}: \mathrm{X} \rightarrow \mathrm{X}^{\prime \prime}$ the canonical morphisms.

Then $h^{\prime \prime}=h^{\prime} \circ h$.
Proof (b) follows from the uniqueness of canonical morphisms. Apply (b) to the case $\left(\mathrm{X}^{\prime \prime}, \xi^{\prime \prime}\right)=(\mathrm{X}, \xi)$. Then $h^{\prime} \circ h=1_{\mathrm{X}}$ and $h \circ h^{\prime}=1_{\mathrm{X}^{\prime}}$, which proves (a).

In practice, the two solutions of a same universal problem are identified, i.e. the two representatives of a same functor, and for the sake of simplicity, any of them is referred to as the solution to the universal problem.

### 2.4.12

Let F and G be representable contravariant functors, (X, $\xi$ ) and (Y, $\eta$ ) respective representatives. The functorial morphisms from F to G are in bijective correspondence with morphisms from $X$ to $Y$. More precisely, for any morphism $\Phi: F \rightarrow G$, there is a unique morphism $\varphi: \mathrm{X} \rightarrow \mathrm{Y}$ such that $\Phi_{\mathrm{X}}(\xi)=\varphi^{*}(\eta)$.

This result is an immediate consequence of 2.4.3. It can also be stated as follows:
Proposition Let $\mathscr{C}$ be a category, and $\widehat{\mathscr{C}}$ the category of contravariant functors from $\mathscr{C}$ to $\mathscr{E} n s$. The functor $\mathrm{X} \rightarrow \widehat{\mathrm{X}}$ from $\mathscr{C}$ to $\widehat{\mathscr{C}}$ is fully faithful.

## Exercices 2.4. (Representable functors)

1.-Let X be a compact metric space. Denote by F the set of closed subsets of X . For closed subsets A and B of X , set $\sigma(\mathrm{A}, \mathrm{B})=\sup _{x \in \mathrm{~A}} d(x, \mathrm{~B})$ and $d(\mathrm{~A}, \mathrm{~B})=$ $\sup (\sigma(\mathrm{A}, \mathrm{B}), \sigma(\mathrm{B}, \mathrm{A}))$. Show that
(a) $d$ is a distance on F ;
(b) F is compact with respect to the topology defined by this distance;
(c) F represents the functor $\Phi: \mathscr{T} 0 \boldsymbol{p} \rightarrow \mathscr{E} n s$ which assigns to each topological space S the set of closed subsets of $\mathrm{S} \times \mathrm{X}$ whose projection onto S is open.

Is the functor $\psi: \mathscr{T o p} \rightarrow \mathscr{E} n s$ which assigns to each space $S$ the set of closed subsets of $\mathrm{S} \times \mathrm{X}$ representable? If yes, is it represented by a Hausdorff space?
2.-Let the locally compact metric space $X$ be the union of countable compact spaces. Define two functors from $\mathscr{T} \sigma \nsim$ to $\mathscr{E} n s$ by assigning respectively to each topological space $S$, the set of closed subsets of $S \times X$ whose projection onto $S$ is open, and the set of closed subsets of $S \times X$ whose projection onto $S$ is both proper (4.1.5) and open. Are these functors representable? Can the assumption that X is metrizable be removed?
3.-Let $E$ and $F$ be normed vector spaces, $B_{E}$ and $B_{F}$ the unit balls of $E$ and $F$ respectively, and $\mathrm{B}_{\mathrm{E}} \otimes \mathrm{B}_{\mathrm{F}}$ the set of $x y$ for $x \in \mathrm{~B}_{\mathrm{E}}$ and $y \in \mathrm{~B}_{\mathrm{F}}$.
(a) Show that the convex envelope from $\mathrm{B}_{\mathrm{E}} \otimes \mathrm{B}_{\mathrm{F}}$ to $\mathrm{E} \otimes \mathrm{F}$ is the unit ball open with respect to the semi-norm $\pi(\mathrm{E}, \mathrm{F})$ on $\mathrm{E} \otimes \mathrm{F}$.
(b) Give an expression for $\pi(\mathrm{E}, \mathrm{F})$.
(c) Formulate the universal problem having $(\mathrm{E} \otimes \mathrm{F}, \pi(\mathrm{E}, \mathrm{F}))$ as solution.
(d) Let $\mathrm{E}_{1}$ be a subspace of E equipped with the induced semi-norm (resp. a quotient of $E$ by a closed subspace equipped with the quotient semi-norm). Is the semi-norm $\pi\left(\mathrm{E}_{1}, \mathrm{~F}\right)$ on $\mathrm{E}_{1} \otimes \mathrm{~F}$ induced by (resp. quotient of) the semi-norm $\pi(\mathrm{E}, \mathrm{F})$ ?
(e) Is the semi-norm $\pi(\mathrm{E}, \mathrm{F})$ a norm?

### 2.5 Products and Inverse Limits

### 2.5.1 Definition

Let $\mathscr{C}$ be a category, X and Y two objects of $\mathscr{C}$. Suppose that the functor $\mathrm{T} \mapsto \operatorname{Hom}(\mathrm{T}, \mathrm{X}) \times \operatorname{Hom}(\mathrm{T}, \mathrm{Y})$ from $\mathscr{C}$ to $\mathscr{E} n \boldsymbol{n}$ is representable and let $(\mathrm{P},(p, q))$ be a representative. Equipped with $(p, q), \mathrm{P}$ is said to be a product of X and Y in $\mathscr{C}$ and, for the sake of simplicity, we write $\mathrm{P}=\mathrm{X} \times \mathrm{Y}, p=\mathrm{pr}_{1}, q=\mathrm{pr}_{2}$.
$\mathscr{C}$ is said to be a product category if there is a product in $\mathscr{C}$ for any two arbitrary objects of $\mathscr{C}$.

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two product categories and F a functor from $\mathscr{C}$ to $\mathscr{C}^{\prime} . \mathrm{F}$ is said to commute with the products of two objects if

$$
(\forall X, Y \in \mathscr{C}) \quad F(X \times Y)=F(X) \times F(Y)
$$

more precisely, if $\mathrm{F}(\mathrm{X} \times \mathrm{Y})$ equipped with $\left(\mathrm{F}\left(\mathrm{pr}_{1}\right), \mathrm{F}\left(\mathrm{pr}_{2}\right)\right)$ is a product of $\mathrm{F}(\mathrm{X})$ and $\mathrm{F}(\mathrm{Y})$ in $\mathscr{C}^{\prime}$.

### 2.5.2 Examples

Let $\mathscr{C}$ be one of the following categories: $\mathscr{G r}, \mathscr{A} \boldsymbol{n} \boldsymbol{n}, \mathrm{A}-\mathscr{M} \boldsymbol{O} d, \mathscr{T} a p$. Then $\mathscr{C}$ is a product category and the forgetful functor from $\mathscr{C}$ to $\mathscr{E}$ ns commutes with products.

In other words, this gives the product $\mathrm{X} \times \mathrm{Y}$ in $\mathscr{C}$ of two objects of $\mathscr{C}$ by equipping the product set with the underlying sets of a particular structure.

### 2.5.3

Let $\mathscr{C}$ be a category, $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ a family of objects of $\mathscr{C}$. Suppose the functor $\mathrm{T} \mapsto \prod_{i} \operatorname{Hom}\left(\mathrm{~T}, \mathrm{X}_{i}\right)$ from $\mathscr{C}$ to $\mathscr{E} n s$ is representable and let $\left(\mathrm{P},\left(p_{i}\right)_{i \in \mathrm{I}}\right)$ be a representative. Together with the family $\left(p_{i}\right)_{i \in \mathrm{I}}, \mathrm{P}$ is said to be a product of the family $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$. We write $\mathrm{P}=\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$ and $\mathrm{pr}_{i}=p_{i}$ is called the projection of index $i$.

Let U be a universe and $\mathscr{C}$ be one of the U-categories $\mathscr{G}_{\boldsymbol{r}}, \mathscr{A} n \boldsymbol{n}, \mathrm{~A}-\mathscr{M} \boldsymbol{\sim}, \mathscr{T} a p$. Then any family of objects indexed by a set belonging to $U$ has a product in $\mathscr{C}$, and the forgetful functor $\mathscr{C} \rightarrow \mathscr{E}$ ns commutes with products.

In $\mathscr{T} a p$, the product $\mathrm{P}=\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$ is the product set endowed with the topology generated by the sets $\mathrm{pr}_{i}^{-1}(\mathrm{U})$, for $i \in \mathrm{I}$ and U open in $\mathrm{X}_{i}$. This topology is called the product topology. A base for this topology is given by $\prod_{i \in \mathrm{I}} \mathrm{U}_{i}$ where, for all $i$, ( $\mathrm{U}_{i}$ ) is a family such that $\mathrm{U}_{i}$ is open in $\mathrm{X}_{i}$ and $\mathrm{U}_{i}=\mathrm{X}_{i}$ except for finitely many $i$. Let $x=\left(x_{i}\right) \in \mathrm{P}$, and for all $i$, let $\left(\mathrm{V}_{i, \alpha}\right)_{\alpha \in \mathrm{A}}$ be a fundamental system of neighbourhoods of $x_{i}$ in $\mathrm{X}_{i}$. The sets $\mathrm{V}_{i_{1}, \alpha_{1}} \times \cdots \times \mathrm{V}_{i_{k}, \alpha_{k}} \times \prod_{j \in \mathrm{I}-\left\{i_{1}, \ldots, i_{k}\right\}} \mathrm{X}_{j}$ form a fundamental system of neighbourhoods of $x$.

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories. A functor F from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ is said to commute with products when these exist if, for any family $\left(\mathrm{X}_{i}\right)$ of objects of $\mathscr{C}$ having a product P in $\mathscr{C}$, the object $\mathrm{F}(\mathrm{P})$ is a product of $\mathrm{F}\left(\mathrm{X}_{i}\right)$ in $\mathscr{C}^{\prime}$.

### 2.5.4 Definition

Let $\mathrm{X}, \mathrm{Y}, \mathrm{S}$ be sets, $f: \mathrm{X} \rightarrow \mathrm{S}$ and $g: \mathrm{Y} \rightarrow \mathrm{S}$ maps. The fibered product of X and Y over S with respect to $f$ and $g$ is the set $\mathrm{X} \times_{\mathrm{S}} \mathrm{Y}=\{(x, y) \mid(x, y) \in$
$\mathrm{X} \times \mathrm{Y}$ and $f(x)=g(y)\}$. Next let $\mathscr{C}$ be a category, $\mathrm{X}, \mathrm{Y}, \mathrm{S}$ objects of $\mathscr{C}, f: \mathrm{X} \rightarrow \mathrm{S}$ and $g: \mathrm{Y} \rightarrow \mathrm{S}$ morphisms. Consider the contravariant functor $\mathrm{T} \mapsto \widehat{\mathrm{X}}(\mathrm{T}) \times \widehat{\mathrm{S}}(\mathrm{T})$ $\widehat{\mathrm{Y}}(\mathrm{T})$ from $\mathscr{C}$ to $\mathscr{E}$ ns. If this functor is representable and if $(\mathrm{P},(p, q))$ is a representative, then P together with $(p, q)$ is said to be a fibered product of X and Y over S with respect to $\grave{a} f$ and $g$, and is written $\mathrm{P}=\mathrm{X} \times_{\mathrm{S}} \mathrm{Y}$.

### 2.5.5 Definition

Let X and Y be sets, $f$ and $g$ maps from X to Y . The kernel of a double arrow $(f, g)$ is the set

$$
\operatorname{Ker}(f, g)=\{x \mid x \in \mathrm{X} \text { and } f(x)=g(x)\} .
$$

Let $\mathscr{C}$ be a category, $\mathrm{X}, \mathrm{Y} \in \mathscr{C}, f$ and $g$ morphisms from X to Y . Consider the functor $\mathrm{F}: \mathrm{T} \mapsto \operatorname{Ker}\left(f_{*}, g_{*}\right)$, where $f_{*}: \widehat{\mathrm{X}}(\mathrm{T}) \rightarrow \widehat{\mathrm{Y}}(\mathrm{T})$ is defined by $f_{*}(\alpha)=f \circ \alpha$ and $g_{*}$ is defined likewise. If this functor is representable, then a representative of this functor is called a kernel of a double arrow $(f, g)$.

### 2.5.6 Definition

An ordered set I is said to be right directed or simply directed if

$$
(\forall i, j \in \mathrm{I})(\exists k \in \mathrm{I}) \quad k \geq i \text { and } k \geq j .
$$

Let I be a directed set and $\mathscr{C}$ a category. An inverse system in $\mathscr{C}$ indexed by I consists of a family $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ of objects of $\mathscr{C}$ and, for any pair $(i, j) \in \mathrm{I}^{2}$ such that $i \leq j$, of a morphism $f_{i}^{j}: \mathrm{X}_{j} \rightarrow \mathrm{X}_{i}$ satisfying the following conditions:
(i) $(\forall i \in \mathrm{I}) f_{i}^{i}=1_{\mathrm{X}_{i}}$;
(ii) $(\forall i, j, k \in \mathrm{I} \mid i \leq j \leq k) f_{i}^{j} \circ f_{j}^{k}=f_{i}^{k}$.

In other words, an inverse system in $\mathscr{C}$ indexed by I is a contravariant functor from $\mathscr{J}$ to $\mathscr{C}$, where $\mathscr{J}$ is the category associated to the ordered set I (2.1.2, Example 5).

### 2.5.7 Definition

Let I be a directed set and $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ an inverse system in $\mathscr{E}$ ns indexed by I. The inverse limit of the inverse system of the given sets is the set

$$
\left\{x=\left(x_{i}\right)_{i \in \mathrm{I}} \in \prod_{i \in \mathrm{I}} \mathrm{X}_{i} \mid(\forall i \in \mathrm{I})(\forall j \in \mathrm{I})_{i \leq j} f_{i}^{j}\left(x_{j}\right)=x_{i}\right\} .
$$

It is written $\underset{i \in \mathrm{I}}{\lim } \mathrm{X}_{i}$, or simply $\underset{\leftarrow}{\lim } \mathrm{X}_{i}$.

### 2.5.8 Definition

Let $\mathscr{C}$ be a category, I a directed set and $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)\right)$ an inverse system in $\mathscr{C}$ indexed by I. If the functor $\mathrm{T} \mapsto \lim _{\leftarrow} \operatorname{Hom}\left(\mathrm{T}, \mathrm{X}_{i}\right)$ is representable, a representative $\left(\mathrm{X},\left(p_{i}\right)\right)$, i.e. a pair $\left(\mathrm{X},\left(p_{i}\right)\right)$ where X is an object of $\mathscr{C}$ and $\left(p_{i}: \mathrm{X} \rightarrow \mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ a family of morphisms such that, for any pair $(i, j) i \leq j, p_{i}=f_{i}^{j} \circ p_{j}$ is called the inverse limit of the inverse system $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$.

In other words, the inverse limit $\lim X_{i}$ equipped with morphisms $p_{i}: \lim _{\leftarrow} \mathrm{X}_{i} \rightarrow \mathrm{X}_{i}$ is characterized by the following universal property:

$$
\begin{aligned}
(\forall \mathrm{T} \in \mathscr{C})\left(\forall\left(g_{i}\right)_{i \in \mathrm{I}}\right. & \left.\in \underset{\leftarrow}{\lim } \operatorname{l\exists } \operatorname{Hom}\left(\mathrm{T}, \mathrm{X}_{i}\right)\right) \\
& \left.\operatorname{Hom}\left(\mathrm{T}, \lim _{\leftarrow} \mathrm{X}_{i}\right)\right)(\forall i \in \mathrm{I}) \quad p_{i} \circ g=g_{i} .
\end{aligned}
$$

### 2.5.9 Example

Let E be a set, I a non-empty directed set and $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ a decreasing sequence of subsets of E. For $i \leq j$, let $f_{i}^{j}$ be the canonical injection from $\mathbf{X}_{j}$ to $\mathbf{X}_{i}$. This defines an inverse system whose inverse limit can be identified with the intersection of the subsets $\mathrm{X}_{i}$.

### 2.5.10

Proposition Let $\mathscr{C}$ be a category and $\mathbb{ש}$ an infinite cardinal. If every family $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ with Card $\mathrm{I} \leq \boldsymbol{ש}$ has a product in $\mathscr{C}$ and if every double arrow $\mathrm{X} \rightarrow \mathrm{Y}$ has a kernel, then every inverse system indexed by an ordered set I of cardinality Card $\mathrm{I} \leq \boldsymbol{\omega}$ has an inverse limit.

Proof Let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ be an inverse system indexed by an ordered set $e ́ \mathrm{I}$ such that Card $\mathrm{I} \leq \boldsymbol{ש}$. For every pair of indices $(i, j)$ with $i \leq j$, set $\mathrm{Y}_{i j}=\mathrm{X}_{i}, \alpha_{i j}=\mathrm{pr}_{i}$ where $\grave{u} \mathrm{pr}_{i}: \prod_{i \in \mathrm{I}} \mathrm{X}_{i} \rightarrow \mathrm{X}_{i}$ is the projection of index $i$ and $\beta_{i j}=f_{i}^{j} \circ p_{j}$. This defines two morphisms $\alpha=\left(\alpha_{i j}\right)$ and $\beta=\left(\beta_{i j}\right)$ from $\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$ to $\prod_{\substack{(i, j) \in ⿺^{2} \\ i \leq j}} \mathrm{Y}_{i j}$.

The kernel for the double arrow $(\alpha, \beta)$ is the inverse limit of the inverse system $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$.

### 2.5.11 Examples

Let $\mathscr{C}$ be one of the following U-categories $\mathscr{E}_{\boldsymbol{r}}, \mathscr{A} \boldsymbol{n} \boldsymbol{n}, \mathrm{A}-\mathscr{M} \boldsymbol{a} \boldsymbol{d}, \mathscr{T} \boldsymbol{\sigma} \boldsymbol{p}$. Then every inverse system (indexed by a set belonging to U ) has an inverse limit and the forgetful functor from $\mathscr{C}$ to $\mathscr{E}$ ns commutes with inverse limits.

Remark In $\mathscr{T} a p$, let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ be an inverse system, and setting $\mathrm{X}=\underset{\leftarrow}{\lim } \mathrm{X}_{i}$, write $p_{i}$ for the canonical map $\mathrm{X} \rightarrow \mathrm{X}_{i}$. Then the topology of X is induced by the product topology on $\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$. Indeed, let $x=\left(x_{i}\right)$ be a point of X and, for all $i$, let $\left(\mathrm{V}_{i, \alpha}\right)_{\alpha \in \mathrm{A}_{i}}$ be a fundamental system of neighbourhoods of $x_{i}$ in X. Set $\mathrm{W}_{i, \alpha}=p_{i}^{-1}\left(\mathrm{~V}_{i, \alpha}\right)$. The subsets $\mathrm{W}_{i_{1}, \alpha_{1}} \cap \cdots \cap \mathrm{~W}_{i_{k}, \alpha_{k}}$ clearly form a fundamental system of neighbourhoods of $x$. Let $i \in \mathrm{I}$ be a common upper bound of $i_{1}, \ldots, i_{k}$. Let $\left(f_{i_{1}}^{i}\right)^{-1}\left(\mathrm{~V}_{i_{1}, \alpha_{1}}\right) \cap \cdots \cap\left(f_{i_{k}}^{i}\right)^{-1}\left(\mathrm{~V}_{i_{k}, \alpha_{k}}\right)$ be a neighbourhood of $x_{i}$ in $\mathrm{X}_{i}$; it thus contains some $\mathrm{V}_{i, \alpha}$, with $\alpha \in \mathrm{A}_{i}$, and $\mathrm{W}_{i, \alpha} \subset \mathrm{~W}_{i_{1}, \alpha_{1}} \cap \cdots \cap \mathrm{~W}_{i_{k}, \alpha_{k}}$.

### 2.5.12 Morphisms of Inverse Systems

Let $\mathscr{C}$ be a category, I a directed set, $\mathscr{X}=\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ and $\mathscr{Y}=\left(\left(\mathrm{Y}_{i}\right),\left(g_{i}^{j}\right)\right)$ inverse systems in $\mathscr{C}$ indexed by I. A morphism from $\mathscr{X}$ to $\mathscr{Y}$ is a family $\left(\varphi_{i}\right)$ of morphisms such that, for $i \leq j$, the following diagram is commutative


If $\mathscr{X}$ and $\mathscr{Y}$ have inverse limits $\mathrm{X}_{\infty}$ and $\mathrm{Y}_{\infty}$, then for any morphism $\left(\varphi_{i}\right)$ from $\mathscr{X}$ to $\mathscr{Y}$, there is a unique morphism $\varphi_{\infty}: \mathrm{X}_{\infty} \rightarrow \mathrm{Y}_{\infty}$ such that $\varphi_{i} \circ p_{i}=q_{i} \circ \varphi_{\infty}$ for all $i$, where $p_{i}: \mathrm{X}_{\infty} \rightarrow \mathrm{X}_{i}$ and $q_{i}: \mathrm{Y}_{\infty} \rightarrow \mathrm{Y}_{i}$ are canonical morphisms. The morphism $\varphi_{\infty}$ is said to be induced by the morphisms $\varphi_{i}$ on passing to the inverse limit.

### 2.5.13 Cofinal Sets

Let I be a directed set and J a subset of I . We remind the reader (1.2.2) that J is said to be cofinal in I if $(\forall i \in \mathrm{I})(\exists j \in \mathrm{~J}) i \leq j$. If J is cofinal in I , then J is a directed set under the induced order.

Let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ be an inverse system of sets. If $\mathbf{J}$ is a cofinal subset of I , the


It follows that, if $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)\right)$ is an inverse system in a category $\mathscr{C}$ and if J is a cofinal subset of I , the inverse system $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{~J}},\left(f_{i}^{j}\right)\right)$ has an inverse limit if and only if so does $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)\right)$, and these limits are identical if they exist.

A cofinal sequence in I is an increasing sequence $\left(i_{n}\right)$ such that the set of $i_{n}$ is cofinal in I.

Proposition Every countable non-empty directed set has a cofinal sequence.
Proof Let $n \mapsto x_{n}$ be a surjective map from $\mathbb{N}$ over I. Define $i_{n}$ inductively by setting $i_{0}=x_{0}$ and choosing for $i_{n+1}$ a common upper bound of $i_{n}$ and $x_{n+1}$. The sequence $\left(i_{n}\right)$ is cofinal.

### 2.5.14 Inverse Limits of Compact Spaces

Theorem In $\mathscr{T o p}$, an inverse limit of non-empty compact is non-empty and compact.

Proof Let $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)\right)$ be non-empty compact projective spaces. By Tychonoff's Theorem (1.7.7), the space $\mathrm{P}=\prod_{i \in \mathrm{I}} \mathrm{X}_{i}$ is compact. For $k \in \mathrm{I}$, set $\mathrm{L}_{k}$ to be the subspace of P of elements $x=\left(x_{i}\right)$ such that $x_{i}=f_{i}^{k}\left(x_{k}\right)$ for $i \leq k$; it is closed in P . If $\tau$ is a choice function (1.1.2) on $\bigcup \mathrm{X}_{i}$, then the element $x$ defined by $x_{k}=\tau\left(\mathrm{X}_{k}\right)$, $x_{i}=f_{i}^{k}\left(x_{k}\right)$ if $i \leq k$ and $x_{i}=\tau\left(\mathrm{X}_{i}\right)$ otherwise, belongs to $\mathrm{L}_{k}$. Then $\left(\mathrm{L}_{k}\right)_{k \in \mathrm{I}}$ are nonempty directed closed sets in P , and so $\lim _{\leftarrow} \mathrm{X}_{i}=\cap_{k \in \mathrm{I}} \mathrm{L}_{k}$ is non-empty and compact.

### 2.5.15

Corollary In $\mathscr{T a p}$, let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ be an inverse system of compact spaces, and set $\mathrm{X}_{\infty}=\lim \mathrm{X}_{i}$. If every $f_{i}^{j}$ is surjective, then so are $f_{i}^{\infty}: \mathrm{X}_{\infty} \rightarrow \mathrm{X}_{i}$.

Proof Let $x \in \mathrm{X}_{i}$, and for $j \geq i$ set $\mathrm{Y}_{j}=\left(f_{i}^{j}\right)^{-1}\left(x_{i}\right)$. By the theorem, $\lim _{\leftarrow} \mathrm{Y}_{j} \neq \varnothing$.

## Exercises 2.5. (Products and inverse limits)

1.-Show that in the category of locally compact spaces and proper maps there are fibered products but no products.
2.-Let $\mathscr{C}$ be a category of finite products, containing a terminal object E . Let $\mathrm{X} \in \mathscr{C}$. A group law on X in $\mathscr{C}$ is a morphism $\mu: \mathrm{X} \times \mathrm{X} \rightarrow \mathrm{X}$ satisfying the following properties:
(i)

commutes.
(ii) There is a map $\varepsilon: \mathrm{E} \rightarrow \mathrm{X}$ such that the diagrams

commute, $e: \mathrm{X} \rightarrow \mathrm{E}$ being the unique morphism from X to E .
(iii) There is a morphism $\sigma: \mathrm{X} \rightarrow \mathrm{X}$ such that the diagrams

commutate.
(a) Show that, if $\mu$ is a group law, then the morphisms $\varepsilon$ and $\sigma$ are unique.
(b) Let $\mu: \mathrm{X} \times \mathrm{X} \rightarrow \mathrm{X}$ be a morphism. It is a group law if and only if for any object T of $\mathscr{C}$, the composition law defined by $\mu$ on $\operatorname{Hom}_{\mathscr{C}}(\mathrm{T}, \mathrm{X})$ is a group law.
(c) Show that defining a group law on X amounts to defining a functor $\Phi: \mathscr{C} \rightarrow \mathscr{G}_{\boldsymbol{\gamma}}$ such that $\omega \circ \Phi=\widehat{\mathrm{X}}$ where $\omega: \mathscr{G}_{\boldsymbol{\gamma}} \rightarrow \mathscr{E} n \boldsymbol{n}$ is the forgetful functor.
(d) Which groups have a group law in the category $\mathscr{E}_{\boldsymbol{\gamma}}$ ?
(e) Let X be an object of $\mathscr{C}$ and $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{G}_{\boldsymbol{\mu}}$ a functor commuting with products. Show that if there is a group law $\mu$ on X in $\mathscr{C}$, then $\mathrm{F}(\mathrm{X})$ is commutative and $\mathrm{F}(\mu)$ agrees with the group law of $\mathrm{F}(\mathrm{X})$.
3.-Let $\mathrm{E}_{0}$ be a vector space $\mathbb{R}[\mathrm{X}]$, and for all $n$, let $\mathrm{E}_{n}$ denote the subspace $\mathrm{X}^{n} \mathbb{R}[\mathrm{X}]$. The map $f: \mathrm{E}_{0} \rightarrow \mathbb{R}$ defined by $f(\mathrm{P})=\mathrm{P}(1)$ induces a surjective map $f_{n}: \mathrm{E}_{n} \rightarrow \mathbb{R}$ for each $n$. Show that the induced map on passing to the inverse limit is not surjective.
4.-Let $\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)_{j \geq i}\right)$ be an inverse system of sets such that all $f_{i}^{j}$ are surjective. Set $\mathrm{X}_{\infty}=\lim _{\longleftarrow} \mathrm{X}_{i}$.
(a) Show that if I is countable, $f_{i}^{\infty}: \mathrm{X}_{\infty} \rightarrow \mathrm{X}_{i}$ is surjective for all $i$.
(b) The aim is to show that the countability assumption can be omitted: let E be an uncountable infinite set and F a countable infinite set. For any countable subset $A$ of $E$, let $X_{A}$ be the set of injective maps from $A$ to $F$ for which the complement of the image is infinite; for $\mathrm{B} \subset \mathrm{A}$, denote by $p_{\mathrm{B}}^{\mathrm{A}}$ the restriction $\mathrm{X}_{\mathrm{A}} \rightarrow \mathrm{X}_{\mathrm{B}}$. Show that the maps $p_{\mathrm{B}}^{\mathrm{A}}$ are surjective, but that the inverse limit of the sets $\mathrm{X}_{\mathrm{A}}$ is empty.
5.-The aim is to prove the following result:

Theorem (Mittag-Leffler) Let $\left(0 \rightarrow \mathrm{E}_{i}^{\prime} \xrightarrow{u_{i}} \mathrm{E}_{i} \xrightarrow{v_{i}} \mathrm{E}_{i}^{\prime \prime} \rightarrow 0,\left(f_{i}^{\prime j}, f_{i}^{j}, f^{\prime \prime j}{ }_{i}\right)\right)$ be an inverse system of exact sequences of Fréchet spaces (i.e. complete and metrizable vector spaces), indexed by a countable directed set I . Suppose that, for any pair ( $i, j$ ) such that $j \geq i$, the image of the map ${f^{\prime}}_{i}^{\prime j}: \mathrm{E}_{j}^{\prime} \rightarrow \mathrm{E}_{i}^{\prime}$ is dense. Then, passage to the inverse limit gives an exact sequence $0 \rightarrow \mathrm{E}_{\infty}^{\prime} \rightarrow \mathrm{E}_{\infty} \rightarrow \mathrm{E}_{\infty}^{\prime \prime} \rightarrow 0$.
(a) Show that $0 \rightarrow \mathrm{E}_{\infty}^{\prime} \rightarrow \mathrm{E}_{\infty} \rightarrow \mathrm{E}_{\infty}^{\prime \prime}$ is an exact sequence.
(b) Show that the question reduces to the case $\mathrm{I}=\mathbb{N}$, where $\mathrm{E}_{n}$ is equipped with a distance $(x, y) \mapsto p_{n}(x-y), \mathrm{E}_{n}^{\prime}$ being a closed subspace of $\mathrm{E}_{n}$ equipped with the induced distance and all $f_{n}^{n+1}$ are contraction maps (i.e. $p_{n}\left(f_{n}^{n+1}(x)\right) \leq p_{n+1}(x)$ for $x \in \mathrm{E}_{n+1}$ ).
(c) Let $x^{\prime \prime}=\left(x_{n}^{\prime \prime}\right)$ be an element of $\mathrm{E}^{\prime \prime}$. Construct inductively a sequence $\left(y_{n}\right)$ such that $y_{n} \in \mathrm{E}_{n}, v_{n}\left(y_{n}\right)=x_{n}^{\prime \prime}$ and $p_{n}\left(f_{n}^{n+1}\left(y_{n+1}\right)-y_{n}\right) \leq \frac{1}{2^{n}}$.
(d) Set $x_{n, k}=f_{n}^{n+k}\left(y_{n+k}\right)$. Show that, for fixed $n,\left(x_{n, k}\right)_{k \in \mathbb{N}}$ is a Cauchy sequence in $\mathrm{E}_{n}$. Conclude.
(e) Show that condition "the images of $f_{i}^{\prime j}$ are dense" can be replaced by the following weaker one:

$$
(\forall i)(\exists j \geq i)(\forall k \geq j) \operatorname{Im} f_{i}^{\prime k} \text { dense in } \operatorname{Im} f_{i}^{\prime j} .
$$

6.-Let X be a set with an anti-reflexive symmetric relation $\mathscr{R}$. A $k$-colouring of X is a map $f$ from X to a set with $k$ elements such that $\mathscr{R}(x, y) \Rightarrow f(x) \neq f(y)$. (For instance, the set X of the regions of a geographical map and $\mathscr{R}(x, y) \Leftrightarrow x$ and $y$ are adjacent.)

Show that, if every finite subset of X has a $k$-colouring (for fixed $k$ ), so has X.
7.-Let A be a group (resp. a ring) and ( $\mathrm{I}_{n}$ ) a decreasing sequence of normal subgroups (resp. ideals) With $\mathrm{I}_{0}=\mathrm{A}$. Let $\left(\varepsilon_{n}\right)$ be a sequence approaching 0 . A function $h: \mathrm{A} \rightarrow \mathbb{R}_{+}$can be defined by $h(x)=\varepsilon_{n}$ if $x \in \mathrm{I}_{n}-\mathrm{I}_{n+1}$ and $h(x)=0$ if $x \in \cap_{n} \mathrm{I}_{n}$. Let $d$ be the mapping $\mathrm{A} \times \mathrm{A}$ to $\mathbb{R}_{+}$defined by $d(x, y)=h\left(x y^{-1}\right)($ resp. $h(x-y))$. Show that
(a) $d$ is a distance on A , the composition law (resp. laws) on A is continuous with respect to the topology defined by $d$.
(b) The Hausdorff completion of A with respect to this distance can be identified with the inverse limit $\lim \mathrm{A} / \mathrm{I}_{n}$.
(c) When $\mathrm{A}=\mathbb{Z}$ and $\overleftarrow{\mathbf{I}}_{n}$ is the ideal generated by $p^{n}$ where $p \in \mathbb{Z}$ :

- for prime $p, \widehat{\mathbb{Z}}_{p}=\lim _{\longleftarrow} \mathbb{Z} / p^{n} \mathbb{Z}$ is an integer ring;
- for $p=p_{1} \times p_{2}$, with g.c.d. $\left(p_{1}, p_{2}\right)=1, \widehat{\mathbb{Z}}_{p}=\widehat{\mathbb{Z}}_{p_{1}} \times \widehat{\mathbb{Z}}_{p_{2}}$;
- for $p=p_{1}^{r}, \widehat{\mathbb{Z}}_{p}=\widehat{\mathbb{Z}}_{p_{1}}$.
(d) The elements of $\widehat{\mathbb{Z}}_{10}$ can be represented by positive non terminating decimal expansions to the left.

What is the expansion of -1 ?
Give an algorithm for the decimal expansions of the idempotents of $\widehat{\mathbb{Z}}_{10}$.
(e) Show that the profinite completion $2.9 .5 \widehat{\mathbb{Z}}$ of $\mathbb{Z}$ can be identified with the product of $\widehat{\mathbb{Z}}_{p}$ over primes $p$.

### 2.6 Sums and Direct Limits

### 2.6.1 Definition

Let $\mathscr{C}$ be a category, $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ a family of objects of $\mathscr{C}$. The sum of $X_{i}$ in $\mathscr{C}$, written $\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}$, is the product, if it exists, of all $\mathrm{X}_{i}$ in the opposite category $\mathscr{C}^{\circ}$, i.e. an object $S$ of $\mathscr{C}$ equipped with a family of morphisms $\left(s_{i}: X_{i} \rightarrow S\right)_{i \in \mathrm{I}}$ such that, for any object T of $\mathscr{C}$ and for any family $\left(f_{i}: \mathrm{X}_{i} \rightarrow \mathrm{~T}_{i}\right)_{i \in \mathrm{I}}$ of morphisms, there is a unique morphism $f: \mathrm{S} \rightarrow \mathrm{T}$ for all $i \in \mathrm{I}$ satisfying $f \circ s_{i}=f_{i}$.

Let $\mathrm{X}, \mathrm{Y}, \mathrm{S} \in \mathscr{C}, f: \mathrm{S} \rightarrow \mathrm{X}$ and $g: \mathrm{S} \rightarrow \mathrm{Y}$ morphisms. An amalgamated sum of X and Y over S with respect to $f$ and $g$, written $\mathrm{X} \sqcup_{\mathrm{S}} \mathrm{Y}$, is the fibered product (if it exists) in the opposite category $\mathscr{C}^{\circ}$ (2.5.4 and 2.1.7).

Let $\mathrm{X}, \mathrm{Y} \in \mathscr{C}, f$ and $g$ morphisms from X to Y . The cokernel of a double arrow $(f, g)$ is the kernel of the double arrow $(f, g)$ (if it exist) in the opposite category $\mathscr{C}^{\circ}$ (2.5.5).

### 2.6.2 Examples

The categories considered in the following examples are U-categories, and the index sets are elements of $U$.
(1) Let $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ be a family of sets. If the sets $\mathrm{X}_{i}$ are disjoint, then the set $\bigcup_{i \in \mathrm{I}} \mathrm{X}_{i}$ equipped with canonical injections is the sum of all $\mathrm{X}_{i}$ in $\mathscr{E}$ ns. In the general case, $\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}=\bigcup_{i \in \mathrm{I}}\{i\} \times \mathrm{X}_{i}$ equipped with maps $s_{i}$ defined by $s_{i}(x)=(i, x)$ is the sum of all $X_{i}$ in $\mathscr{E} n s$.
(2) In $\mathscr{T} \sigma \not \mathcal{R}$ every family of topological spaces has a sum and the forgetful functor from $\mathscr{T} a p$ to $\mathscr{E} n s$ commutes with the sum.

Let $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$ be a family of topological spaces. The topological space $\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}$ is given by the summation set $\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}$ equipped with the topology whose open subsets are $\bigsqcup_{i \in \mathrm{I}} \mathrm{U}_{i}$, where, for each $i, \mathrm{U}_{i}$ is open in $\mathrm{X}_{i}$.

In the following examples, the sums exist but the forgetful functors do not commute with them. To avoid confusion, the sums will be written using a symbol other than $\sqcup$.
(3) In $\mathscr{G}_{\boldsymbol{\mu}}$ any family $\left(\mathrm{G}_{i}\right)_{i \in \mathrm{I}}$ of groups has a sum also called a free product of the groups $\mathrm{G}_{i}$. It is written $*_{i \in \mathrm{I}} \mathrm{G}_{i}$, where $\mathrm{G}_{1} * \mathrm{G}_{2} * \cdots * \mathrm{G}_{k}$ if $\mathrm{I}=\{1, \ldots, k\}$.

The free product of a family $\left(\mathrm{G}_{i}\right)_{i \in \mathrm{I}}$ of groups is obtained as follows: let S be the summation set of the underlying sets of the groups $\mathrm{G}_{i}, s_{i}: \mathrm{G}_{i} \rightarrow \mathrm{~S}$ the canonical injections, and $M$ the set $\bigsqcup_{n \in \mathbb{N}} S^{n}$. Define a composition law on $M$ by assigning to the pair $\left(\left(x_{1}, \ldots, x_{p}\right),\left(y_{1}, \ldots, y_{q}\right)\right)$ the element $\left(x_{1}, \ldots, x_{p}, y_{1}, \ldots, y_{q}\right) \in \mathrm{S}^{p+q}$. This law is associative and has an identity element: the empty set, which is the unique element of $S^{0}$. For any group H and any family $f=\left(f_{i}: \mathrm{G}_{i} \rightarrow \mathrm{H}\right)_{i \in \mathrm{I}}$ of morphisms, define a map $\alpha_{f}: \mathrm{S} \rightarrow \mathrm{H}$ by $\alpha_{f} \circ s_{i}=f_{i}$ for all $i \in \mathrm{I}$ and a map $\beta_{f}: \mathrm{M} \rightarrow \mathrm{H}$ par $\beta_{f}\left(x_{1}, \ldots, x_{p}\right)=\alpha_{f}\left(x_{1}\right) \cdots \alpha_{f}\left(x_{p}\right)$. The map $\beta$ is a homomorphism. In M, write $x \sim y$ if, for any group H and any family of morphisms $f=\left(f_{i}: \mathrm{G}_{i} \rightarrow \mathrm{H}\right)_{i \in \mathrm{I}}$, $\beta_{f}(x)=\beta_{f}(y)$. This is an equivalence relation compatible with the composition law of M , and hence there is an associative quotient law with an identity element on $\mathrm{L}=$ $\mathrm{M} / \sim$; moreover, setting $\chi$ for the canonical map $\mathrm{M} \rightarrow \mathrm{L}$, the element $\chi\left(x_{1}, \ldots, x_{p}\right)$ of L has as inverse $\chi\left(x_{p}^{-1}, \ldots, x_{1}^{-1}\right)$. Hence L is a group. Identify S with $\mathrm{S}^{1}$ and write $\lambda_{i}$ for the morphism $\chi \circ \iota_{1} \circ s_{i}$ where $\iota_{1}$ is the canonical injection from $S^{1}$ to $\bigsqcup S^{n}$. The group L equipped with the morphisms $\lambda_{i}$ is the desired free product.
(4) In the category of Abelian groups, any family $\left(\mathrm{G}_{i}\right)_{i \in \mathrm{I}}$ of objects has a sum, called the direct sum of the groups $\mathrm{G}_{i}$ and written $\bigoplus_{i \in \mathrm{I}} \mathrm{G}_{i}$. If I is finite, the direct sum can be identified with the product $\prod_{i \in \mathrm{I}} \mathrm{G}_{i}$. In the general case, $\bigoplus_{i \in \mathrm{I}} \mathrm{G}_{i}$ can be identified with the subgroup of $\prod_{i \in \mathrm{I}} \mathrm{G}_{i}$ consisting of elements having finitely many non-trivial coordinates.

In general, the free product $* \mathrm{G}_{i}$ is not commutative (for example $\mathbb{Z} * \mathbb{Z}$ ); the direct sum of the groups $\mathrm{G}_{i}$ can be identified with a quotient of the free product.
(5) Let A be a commutative ring. In the category of unital associative commutative A-algebras, the sum of two algebras E and F is the tensor product algebra $\mathrm{E} \otimes_{\mathrm{A}} \mathrm{F}$ (3.8.15).

### 2.6.3 Definition

Let $\mathscr{C}$ be a category, and I a directed set.
A direct system in $\mathscr{C}$ indexed by I is an inverse system in $\mathscr{C}^{0}$ indexed by I , i.e. for every $i \in \mathrm{I}$, consists of an object $\mathrm{X}_{i}$ of $\mathscr{C}$ and, for every pair $(i, j)$ such that $i \leq j$, of a morphism $f_{i}^{j}: \mathrm{X}_{i} \rightarrow \mathrm{X}_{j}$ such that, for all $i$,

$$
f_{i}^{i}=1_{\mathrm{X}_{i}},
$$

and for all $i \leq j \leq k$,

$$
f_{j}^{k} \circ f_{i}^{j}=f_{i}^{k}
$$

### 2.6.4 Definition

Let $\mathscr{C}$ be a category, I a directed set, and $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ a direct system in $\mathscr{C}$. If $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ has a direct limit $\left(\mathrm{X}_{\infty},\left(\sigma_{i}\right)\right)$ in $\mathscr{C}^{0}$, then $\left(\mathrm{X}_{\infty},\left(\sigma_{i}\right)\right)$ is said to be the direct limit in $\mathscr{C}$ of the given direct system. It is written $X_{\infty}=\lim X_{i}$.

In other words, the direct limit $\mathrm{X}_{\infty}$ equipped with morphisms $\overrightarrow{\sigma_{i}}: \mathrm{X}_{i} \rightarrow \mathrm{X}_{\infty}$ is characterized by the following universal property: for all $\mathrm{T} \in \mathscr{C}$, for all families of morphisms $\left(f_{i}: \mathrm{X}_{i} \rightarrow \mathrm{~T}\right)_{i \in \mathrm{I}}$ such that $f_{i}=f_{j} \circ f_{i}^{j}$, for all pairs $(i, j)$ with $i \leq j$, there is a unique morphism $g: \mathrm{X}_{\infty} \rightarrow \mathrm{T}$ such that, for all $i, g \circ \sigma_{i}=f_{i}$.

### 2.6.5 Direct Limits of Sets

Let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ be a direct system of sets indexed by a directed set I. On the summation set $\mathrm{S}=\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}$ define an equivalence relation by setting the class of $(i, x)$ to consist of all $(j, y)$ for which there is an upper bound $k$ of $i$ and $j$ such that $f_{i}^{k}(x)=f_{j}^{k}(y)$. The quotient set $\mathrm{X}_{\infty}=\mathrm{S} / \sim$, together with the maps $\sigma_{i}=\chi \circ s_{i}$, where $\chi$ is the canonical map $\mathrm{S} \rightarrow \mathrm{X}_{\infty}$ and $s_{i}$ the canonical map $\mathrm{X}_{i} \rightarrow \mathrm{~S}$, is the direct limit of the given system.

The set $\mathrm{X}_{\infty}$ together with the maps $\sigma_{i}$ is characterized by the following properties:
(i) $\sigma_{i}=\sigma_{j} \circ f_{i}^{j}$ pour tout $(i, j)$;
(ii) $\mathrm{X}_{\infty}=\bigcup_{i \in \mathrm{I}} \sigma_{i}\left(\mathrm{X}_{i}\right)$;
(iii) $(\forall i \in \mathrm{I})\left(\forall x, y \in \mathrm{X}_{i}\right) \sigma_{i}(x)=\sigma_{i}(y) \Rightarrow(\exists j \geq i) f_{i}^{j}(x)=f_{i}^{j}(y)$.

Let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ and $\left(\left(\mathrm{Y}_{i}\right),\left(g_{i}^{j}\right)\right)$ be two direct systems of sets indexed by the same set, and $\left(\mathrm{X}_{\infty},\left(\sigma_{i}\right)\right)$ and $\left(\mathrm{Y}_{\infty},\left(\tau_{i}\right)\right)$ their direct limits. Then, $\lim \left(\mathrm{X}_{i} \times \mathrm{Y}_{i}\right)=$ $\mathrm{X}_{\infty} \times \mathrm{Y}_{\infty}$. Indeed, $\mathrm{X}_{\infty} \times \mathrm{Y}_{\infty}$, equipped with the maps $\sigma_{i} \times \tau_{i}: \mathrm{X}_{i} \overrightarrow{\times} \mathrm{Y}_{i} \rightarrow \mathrm{X}_{\infty} \times$ $Y_{\infty}$, satisfy above conditions (i), (ii) and (iii).

### 2.6.6

In the categories $\mathscr{E}_{\boldsymbol{r}}, \mathscr{A} \boldsymbol{n} \boldsymbol{n}, \mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}, \mathscr{T} \sigma \boldsymbol{p}$, any direct system has a direct limit, and the forgetful functors commute with the direct limits.

In $\mathscr{T} \sigma p$, the direct $\operatorname{limit} \lim X_{i}$ is obtained by equipping the direct limit $X_{\infty}$ with the finest topology for which the canonical maps $\sigma_{i}: \mathrm{X}_{i} \rightarrow \mathrm{X}_{\infty}$ remain continuous. A subset U of $\mathrm{X}_{\infty}$ is open with respect to this topology if and only if, for any $i$, $\sigma_{i}^{-1}(\mathrm{U})$ is open in $\mathrm{X}_{i}$.

Note that the space $\mathrm{X}_{\infty}$ need not be Hausdorff even if the spaces $\mathrm{X}_{i}$ are.

### 2.6.7

Let $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ and $\left(\left(\mathrm{Y}_{i}\right),\left(g_{i}^{j}\right)\right)$ be two direct systems of sets indexed by the same set I, and $\left(\varphi_{i}: X_{i} \rightarrow \mathrm{Y}_{i}\right)_{i \in \mathrm{I}}$ a morphism of direct systems (i.e. a family of maps commuting with all $f_{i}^{j}$ and $g_{i}^{j}$ ). Write $\varphi_{\infty}$ for the map from $\mathrm{X}_{\infty}$ to $\mathrm{Y}_{\infty}$ induced by the morphisms the $\varphi_{i}$ on passing to the direct limit. If the $\varphi_{i}$ are injective (resp. surjective), $\varphi_{\infty}$ is injective (resp. surjective); in general,

$$
\varphi_{\infty}\left(\mathrm{X}_{\infty}\right)=\underset{\longrightarrow}{\lim } \varphi_{i}\left(\mathrm{X}_{i}\right) .
$$

Remark If $\left(\left(\mathrm{X}_{i}\right),\left(f_{i}^{j}\right)\right)$ and $\left(\left(\mathrm{Y}_{i}\right),\left(g_{i}^{j}\right)\right)$ are inverse systems of sets, and if $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ is a morphism of inverse systems such that, for all $i, \varphi_{i}$ is injective, then $\varphi_{\infty}$ is also injective; but if, for all $i, \varphi_{i}$ is surjective, $\varphi_{\infty}$ need not be surjective (2.5, Exercise 4).

### 2.6.8 Direct Limit of Categories

If $\mathscr{C}$ is a category, write $\mathscr{F} \mathscr{C}(\mathscr{C})$ for the set of morphisms of $\mathscr{C}$, namely

$$
\mathscr{F} \mathscr{C}(\mathscr{C})=\bigsqcup_{(\mathrm{X}, \mathrm{Y}) \in \mathscr{C} \times \mathscr{C}} \operatorname{Hom}(\mathrm{X}, \mathrm{Y}),
$$

and define the maps $\alpha$ and $\beta$ from $\mathscr{F} \mathscr{C}(\mathscr{C})$ to $\mathscr{C}$ by $\alpha(f)=\mathrm{X}$ and $\beta(f)=\mathrm{Y}$ if $f \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y})$. The composition law is a map $(\mathscr{F} \mathscr{C}(\mathscr{C}), \beta) \times_{\mathscr{C}}(\mathscr{F} \mathscr{C}(\mathscr{C}), \alpha) \rightarrow$ $\mathscr{F} \mathscr{C}(\mathscr{C})$.

Let $\left(\left(\mathscr{C}_{i}\right),\left(\mathrm{F}_{i}^{j}\right)\right)$ be a direct system in $\mathscr{C} \boldsymbol{a} t$ : all $\mathscr{C}_{i}$ are categories and $\mathrm{F}_{i}^{j}$ functors. Define a category $\mathscr{C}_{\infty}$ by setting its objects to be the direct limits $\mathscr{C}_{i}$ and by setting $\mathscr{F} \mathscr{C}\left(\mathscr{C}_{\infty}\right)=\lim \mathscr{F} \mathscr{C}\left(\mathscr{C}_{i}\right)$, the maps $\alpha$ and $\beta$ and the composition law being induced


The category $\mathscr{C}_{\infty}$ obtained is the direct limit of the given system in $\mathscr{C} \boldsymbol{a} t$.
Let X and Y be objects of $\mathscr{C}_{\infty}$, and let $(i, \mathrm{~A})$ and $(j, \mathrm{~B})$ be representatives of X and $Y$ respectively. Then

$$
\operatorname{Hom}_{\mathscr{C}_{\infty}}(\mathrm{X}, \mathrm{Y})=\underset{k \geq \overrightarrow{, k \geq j}}{\lim } \operatorname{Hom}_{\mathscr{C}_{k}}\left(\mathrm{~F}_{i}^{k}(\mathrm{X}), \mathrm{F}_{j}^{k}(\mathrm{Y})\right)
$$

## Exercises 2.6. (Sums and direct limits)

1.-Let $\mathscr{C}$ be a category with sum. Let $\mathrm{X} \in \mathscr{C}$, and $i_{1}$ and $i_{2}$ the canonical morphisms from X to $\mathrm{X} \sqcup \mathrm{X}$. The codiagonal morphism of X is the unique morphism $\nabla: \mathrm{X} \sqcup$ $\mathrm{X} \rightarrow \mathrm{X}$ such that $\nabla \circ i_{1}=\nabla \circ i_{2}=1_{\mathrm{X}}$. Describe the codiagonal morphism in the following categories: $\mathscr{E} n \boldsymbol{n}$, the category $\mathscr{T} a p$ of point spaces (topological spaces
with what is known as a basepoint, the morphisms from X to Y being the continuous maps from X to Y sending the basepoint of X to the basepoint of Y$), \mathrm{A}-\mathscr{M}-\boldsymbol{d}$, and the category $\mathscr{E}_{\mu}$ of unital associative commutative algebras over a ring A.
2.-Let $\mathscr{C}$ be a product category with sums, having an object O which is both initial and terminal. For $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, let $\mathrm{O}_{\mathrm{X}, \mathrm{Y}}$ be the unique morphism from X to Y which can be factorized through O , and define $\varphi_{X, Y}: \mathrm{X} \sqcup \mathrm{Y} \rightarrow \mathrm{X} \times \mathrm{Y}$ by $\varphi_{\mathrm{X}, \mathrm{Y}} \circ i_{1}=\left(1_{\mathrm{X}}, \mathrm{O}_{\mathrm{X}, \mathrm{Y}}\right)$ and $\varphi_{\mathrm{X}, \mathrm{Y}} \circ i_{2}=\left(\mathrm{O}_{\mathrm{Y}, \mathrm{X}}, 1_{\mathrm{Y}}\right)$.
(a) Suppose that, for all X and Y , the morphism $\varphi_{X, Y}$ is an isomorphism. For every object X , define $\alpha_{\mathrm{X}}: \mathrm{X} \times \mathrm{X} \rightarrow \mathrm{X}$ by $\alpha_{\mathrm{X}}=\nabla_{\mathrm{X}} \circ \varphi_{\mathrm{X}, \mathrm{X}}^{-1}$, where $\nabla_{\mathrm{X}}$ is the codiagonal morphism of X . Show that $\alpha_{\mathrm{X}}$ is a commutative and associative law in $\mathscr{C}$. For all $\mathrm{X}, \mathrm{Y}$, the morphism $\alpha_{\mathrm{Y}}$ gives a composition law $\alpha_{\mathrm{X}, \mathrm{Y}}$ on $\operatorname{Hom}(\mathrm{X}, \mathrm{Y})$. This law is associative, commutative and has $\mathrm{O}_{\mathrm{X}, \mathrm{Y}}$ as identity element. Let $\mathrm{X}, \mathrm{Y}, \mathrm{Z} \in \mathscr{C}$. For fixed $f \in \operatorname{Hom}(\mathbf{X}, \mathrm{Y})$ (resp. for $g \in \operatorname{Hom}(\mathrm{Y}, \mathrm{Z})$ ), the map $g \mapsto g \circ f$ (resp. $f \mapsto g \circ f)$ is a homomorphism from $\operatorname{Hom}(\mathrm{Y}, \mathrm{Z})(\operatorname{resp} . \operatorname{Hom}(\mathrm{X}, \mathrm{Y}))$ to $\operatorname{Hom}(\mathrm{X}, \mathrm{Z})$.
(b) Assume that, for all $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$, there is a composition law $\mu_{\mathrm{X}, \mathrm{Y}}$ on $\operatorname{Hom}(\mathrm{X}, \mathrm{Y})$ with $\mathrm{O}_{\mathrm{X}, \mathrm{Y}}$ the identity element. Suppose that, as above, the maps $g \mapsto g \circ f$ and $f \mapsto g \circ f$ are homomorphisms. Show that $\varphi_{\mathrm{X}, \mathrm{Y}}$ is an isomorphism for all X and Y , and that $\mu_{\mathrm{X}, \mathrm{Y}}=\alpha_{\mathrm{X}, \mathrm{Y}}$. In particular $\mu_{\mathrm{X}, \mathrm{Y}}$ is associative and commutative.
3.-Let X be a Hausdorff topological space.
(a) Show that the following conditions are equivalent:
(i) X is the direct limit of compact subspaces;
(ii) A subset A of X is closed if and only if, for any compact subspace K of X , $\mathrm{A} \cap \mathrm{K}$ is closed.
(b) Show that every metrizable space is the direct limit of compact subspaces and that the same holds for every locally compact space.
(c) Show that in the category $\mathscr{C}$ of Hausdorff direct limits of compact subspaces, there are products and that the product $\mathrm{X} \times \mathrm{Y}$ in $\mathscr{C}$ of $\mathrm{X}, \mathrm{Y} \in \mathscr{C}$ is the set $\mathrm{X} \times \mathrm{Y}$ equipped with a topology finer than the product topology. Give examples where these two topologies agree.
(d) In $\mathbb{R}^{(\mathbb{N})}$, write $\mathrm{X}_{n}$ for the union $\bigcup_{k \leq n}\left[0, e_{k}\right]$ (where the elements $e_{n}$ form the canonical basis of $\mathbb{R}^{(\mathbb{N})}$ ) equipped with the topology induced by that of $\mathbb{R}^{\mathbb{N}}$. Set $\mathrm{X}=\underset{\rightarrow}{\lim } \mathrm{X}_{n}$ (i.e. $\cup \mathrm{X}_{n}$ ) equipped with the direct limit topology. For any sequence $\varepsilon=\left(\overrightarrow{\varepsilon_{n}}\right)$ of strictly positive numbers, set $\mathrm{V}_{\varepsilon}=\bigcup\left[0, \varepsilon_{n} e_{n}[\right.$.

Show that the sets $\mathrm{V}_{\varepsilon}$ form a fundamental system of neighbourhoods in X , and that any compact subset of $X$ is contained in some $X_{n}$.
(e) Let Y be the unit ball of $\ell^{2}$. Write $\left(f_{p}\right)_{p \in \mathbb{N}^{*}}$ for the canonical Hilbert basis of Y . In $\mathrm{X} \times \mathrm{Y}$, for any $n \in \mathbb{N}^{*}$, consider the set $\mathrm{A}_{n}=\left\{\left(\frac{1}{p} e_{n}, \frac{1}{n} f_{p}\right)\right\}$. Show that $\mathrm{A}_{n}$ is closed in $\mathrm{X} \times \mathrm{Y}$. Let $\mathrm{A}=\bigcup \mathrm{A}_{n}$. Show that the intersection of A and of a compact subset of $\mathrm{X} \times \mathrm{Y}$ is closed, but that every neighbourhood of 0 meets A .

Deduce that the topology of the product $\mathrm{X} \times \mathrm{Y}$ in the category $\mathscr{C}$ is strictly finer than the product topology.
4. (Alternative proof of the existence of the free product of a family of groups.) Let $\left(\mathrm{G}_{i}\right)_{i \in \mathrm{I}}$ be a family of groups.
(a) (Lemma about cardinality.) Let H be a group and $\left(f_{i}: \mathrm{G}_{i} \rightarrow \mathrm{H}\right)_{i \in \mathrm{I}}$ a family of homomorphisms. Show that the cardinality of the subgroup of H generated by all $f_{i}\left(\mathrm{G}_{i}\right)$ is bounded above by the supremum $\boldsymbol{ש}$ of the cardinalities of $\mathrm{I}, \mathrm{G}_{i}$ and $\aleph_{0}$.
(b) Let X be a set with cardinality $\boldsymbol{ש}$, and $\Lambda$ the set of pairs $(\mathrm{H}, f)$, where H is a group whose underlying set is contained in X and $f=\left(f_{i}\right)_{i \in \mathrm{I}}$ a family of homomorphisms $f_{i}: \mathrm{G}_{i} \rightarrow \mathrm{H}$.

Set $\mathrm{P}=\prod_{(\mathrm{H}, f) \in \Lambda} \mathrm{H}$ and write $\varphi_{i}$ for the homomorphism $\left(f_{i}\right)_{(\mathrm{H}, f) \in \Lambda}$ from $\mathrm{G}_{i}$ to P. Show that the subgroup L of P generated by all $\varphi_{i}\left(\mathrm{G}_{i}\right)$, equipped with the family of homomorphisms $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ is a free product of the groups $\mathrm{G}_{i}$.
5.-With the notation of 2.6.2, Example 3, an element of $S^{n}$ is called a word of length $n$. A word ( $x_{1}, \ldots, x_{n}$ ) is said to be reduced if no $x_{k}$ is the identity element for $k \in\{1, \ldots, n\}$, and if $x_{k}$ and $x_{k+1}$ belong to different groups for $k \in\{1, \ldots, n-1\}$.
(a) Show that every element of the free product L can be represented by a reduced word.
(b) Give an algorithm associating to any two reduced words $\left(x_{1}, \ldots, x_{p}\right)$ and $\left(y_{1}, \ldots, y_{q}\right)$ a reduced word of length $\grave{a} \leq p+q$ representing the same element of L as $\left(x_{1}, \ldots, x_{p}, y_{1}, \ldots, y_{q}\right)$. Show that the law thereby defined in the set of reduced words is associative.
(c) Deduce that every element of L can be uniquely represented by a reduced word.

### 2.7 Adjoint Functors

### 2.7.1 Definition

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be categories, $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ and $\mathrm{G}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$ functors. G is said to be a right adjoint to F , or F a left adjoint to G , if the functors

$$
(\mathrm{X}, \mathrm{Y}) \mapsto \operatorname{Hom}_{\mathscr{B}^{\prime}}(\mathrm{F}(\mathrm{X}), \mathrm{Y}) \text { and }(\mathrm{X}, \mathrm{Y}) \mapsto \operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{G}(\mathrm{Y}))
$$

from $\mathscr{C}^{0} \times \mathscr{C}^{\prime}$ to $\mathscr{E}$ ns are isomorphic.
When this is the case, for the sake of simplicity, one writes:

$$
\operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{F}(\mathrm{X}), \mathrm{Y})=\operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{G}(\mathrm{Y}))
$$

### 2.7.2

Proposition Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be categories. A functor $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ has a right adjoint if and only if, for all $\mathrm{Y} \in \mathscr{C}^{\prime}$, the contravariant functor $\mathrm{X} \mapsto \operatorname{Hom}(\mathrm{F}(\mathrm{X}), \mathrm{Y})$ from $\mathscr{C}$ to $\mathscr{E}$ ns is representable.
Proof The condition is obviously necessary. We show that it is sufficient. Let $\widehat{\mathscr{C}}_{0}$ be the category of representable contravariant functors from $\mathscr{C}$ to $\mathscr{E}$ ns. By definition, the functor $\Lambda: \mathrm{X} \mapsto \widehat{\mathrm{X}}$ from $\mathscr{C}$ to $\widehat{\mathscr{C}}_{0}$ is fully faithful (2.4.12, proposition) and essentially surjective. Hence it is an equivalence of categories (2.3.5). Let $\mathrm{M}: \widehat{\mathscr{C}}_{0} \rightarrow \mathscr{C}$ be a functor such that $\Lambda \circ \mathrm{M}$ is isomorphic to the identity functor of $\widehat{\mathscr{C}}_{0}$. Let $\mathrm{H}: \mathscr{C}^{\prime} \rightarrow \widehat{\mathscr{C}}_{0}$ be the functor defined by $\mathrm{H}(\mathrm{Y})(\mathrm{X})=\operatorname{Hom}(\mathrm{F}(\mathrm{X}), \mathrm{Y})$. Set $\mathrm{G}=\mathrm{M} \circ \mathrm{H}$. Then $\Lambda \circ \mathrm{G}(\mathrm{Y})(\mathrm{X})=\operatorname{Hom}(\mathrm{X}, \mathrm{G}(\mathrm{Y}))$, and $\Lambda \circ \mathrm{G}=\Lambda \circ \mathrm{M} \circ \mathrm{H} \approx \mathrm{H}$. So F and G are adjoint functors.

### 2.7.3

Proposition Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be two categories and G a functor from $\mathscr{C}^{\prime}$ to $\mathscr{C}$. If G admits a left adjoint, then G commutes with products, kernels of double arrows, fibered products and inverse limits when these exist.

Proof We only give the proof for inverse limits, the other cases being similar. Let F be a left adjoint to G. Let $\left(\mathrm{Y}_{i}\right)$ be an inverse system of objects of $\mathscr{C}^{\prime}$ having an inverse limit. Then, for all $\mathrm{X} \in \mathscr{C}$

$$
\begin{aligned}
\operatorname{Hom}\left(\mathrm{X}, \mathrm{G}\left(\lim \mathrm{Y}_{i}\right)\right)= & \operatorname{Hom}\left(\mathrm{F}(\mathrm{X}), \stackrel{\lim }{\longleftarrow} \mathrm{Y}_{i}\right)=\lim \operatorname{Hom}\left(\mathrm{F}(\mathrm{X}), \mathrm{Y}_{i}\right) \\
= & \lim _{\leftarrow}^{\leftarrow} \operatorname{Hom}\left(\mathrm{X}, \mathrm{G}\left(\mathrm{Y}_{i}\right)\right)=\operatorname{Hom}\left(\mathrm{X}, \lim _{\leftarrow} \mathrm{G}\left(\mathrm{Y}_{i}\right)\right)
\end{aligned}
$$

and so $\mathrm{G}\left(\lim _{\leftarrow} \mathrm{Y}_{i}\right)$ is an inverse limit of $\left(\mathrm{G}\left(\mathrm{Y}_{i}\right)\right)$.
Corollary Let F be a functor from $\mathscr{C}$ to $\mathscr{C}^{\prime}$. If F admits a right adjoint, then F commutes with sums, amalgamated sums, cokernels of double arrows and direct limits when these exist.

Proof Apply the proposition to opposite categories.

### 2.7.4 Examples

(1) Let $\mathscr{C}$ be the category of normed vector spaces over $\mathbb{R}$ and $\mathscr{C}^{\prime}$ the category of Banach spaces. The functor from $\mathscr{C}$ to $\mathscr{C}^{\prime}$ which assigns to every normed space its completion is a left adjoint to the inclusion functor from $\mathscr{C}^{\prime}$ to $\mathscr{C}$. Indeed, for any normed space $E$ and any Banach space $F, L(\widehat{E}, F)=L(E, F)$.
(2) Let A be a ring. The functor from $\mathscr{E}$ ns to $\mathrm{A}-\mathscr{M} \boldsymbol{O} \boldsymbol{d}$ which assigns to every set X the free module $\mathrm{A}^{(\mathrm{X})}$ constructed on X is a left adjoint to the forgetful functor $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d} \rightarrow \mathscr{E} n s$. Indeed, for any set X and any A-module E , $\operatorname{Hom}_{\mathrm{A}-\operatorname{Mod}}\left(\mathrm{A}^{(\mathrm{X})}, \mathrm{E}\right)=\operatorname{Hom}_{\mathscr{\circ} n \boldsymbol{s}}(\mathrm{X}, \mathrm{E})$ (3.4.6).

Similarly, the functor from $\mathscr{E}$ ns to the category A- $\mathscr{A C g}$ of unital associative commutative algebras which assigns to each set $X$ the polynomial algebra $A[X]$ is a left adjoint to the forgetful functor from A- $\mathscr{A l} \mathscr{g}$ to $\mathscr{E} n s$.
(3) Let A and B be two rings and $\mathrm{A} \rightarrow \mathrm{B}$ a ring homomorphism. The extension of scalars functor $\mathrm{E} \mapsto \mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}$ from $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$ to $\mathrm{B}-\mathscr{M} \sigma \boldsymbol{d}$ is a left adjoint to the restriction of scalars functor $\mathrm{B}-\mathscr{M} \sigma \boldsymbol{d} \rightarrow \mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$. Indeed, for any A-module E and any B-module F, $\operatorname{Hom}_{\mathrm{B}}\left(\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}, \mathrm{F}\right)=\operatorname{Hom}_{\mathrm{A}}(\mathrm{E}, \mathrm{F})$.
(4) Let A be a commutative ring and F an A -module. The functor $\mathrm{E} \mapsto \mathrm{E} \otimes_{\mathrm{A}} \mathrm{F}$ from the category A- $\mathbb{M} \sigma d$ to itself is a left adjoint to the functor $G \mapsto \operatorname{Hom}_{A}(F, G)$. Indeed, if E and G are A-modules, then $\operatorname{Hom}_{\mathrm{A}}\left(\mathrm{E} \otimes_{\mathrm{A}} \mathrm{F}, \mathrm{G}\right)=\operatorname{Hom}_{\mathrm{A}}\left(\mathrm{E}, \operatorname{Hom}_{\mathrm{A}}(\mathrm{F}, \mathrm{G})\right)$.
(5) Let Y be a locally compact space. The functor $\mathrm{X} \mapsto \mathrm{X} \times \mathrm{Y}$ from the category $\mathscr{T} 0 \not 2$ to itself has for right adjoint $\grave{a}$ the functor $\mathrm{Z} \mapsto \mathrm{C}(\mathrm{Y}, \mathrm{Z})$, where $\mathrm{C}(\mathrm{Y}, \mathrm{Z})$ is the set of continuous functions from Y to Z equipped with the topology of compact convergence.
(6) Let $\mathscr{C}$ be the category of point spaces, i.e. of pairs $\left(\mathrm{X}, x_{0}\right)$ where X is a topological space and $x_{0} \in \mathrm{X}$, with

$$
\operatorname{Hom}_{\mathscr{C}}\left(\left(\mathrm{X}, x_{0}\right),\left(\mathrm{Y}, y_{0}\right)\right)=\left\{f: \mathrm{X} \rightarrow \mathrm{Y} \mid f \text { continuous and } f\left(x_{0}\right)=y_{0}\right\}
$$

For $\left(\mathrm{X}, x_{0}\right) \in \mathscr{C}$, let $\Omega\left(\mathrm{X}, x_{0}\right)$ be the loop space of $\left(\mathrm{X}, x_{0}\right)$, i.e. the set of continuous maps $\gamma:[0,1] \rightarrow \mathrm{X}$ such that $\gamma(0)=\gamma(1)=x_{0}$, equipped with the topology of compact convergence and basepoint $\bar{x}_{0}: t \mapsto x_{0}$. This defines a functor $\Omega$ from the category $\mathscr{C}$ to itself. This functor admits a left adjoint $S$ defined as follows: $\mathrm{S}\left(\mathrm{X}, x_{0}\right)$ is the quotient space of $\mathrm{X} \times[0,1]$ by the equivalence relation contracting the set

$$
(\mathrm{X} \times\{0\}) \cup(\mathrm{X} \times\{1\}) \cup\left(\left\{x_{0}\right\} \times[0,1]\right)
$$

to a point, equipped with this point as basepoint.
In other words, for $\mathrm{X}, \mathrm{Y} \in \mathscr{C}, \operatorname{Hom}(\mathrm{S}(\mathrm{X}), \mathrm{Y})=\operatorname{Hom}(\mathrm{X}, \Omega(\mathrm{Y}))$. The space $\mathrm{S}(\mathrm{X})$ is called the reduced suspension of X .

### 2.7.5

Let $\mathscr{C}$ and $\mathscr{C}^{\prime}$ be categories, $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ and $\mathrm{G}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$ functors such that G is the right adjoint to $F$. Write FX for $\mathrm{F}(\mathrm{X}), \mathrm{GF}$ for $\mathrm{G} \circ \mathrm{F}$, etc. Define a functorial morphism $\alpha: 1_{\mathscr{C}} \rightarrow$ GF by setting $\alpha_{\mathrm{X}}$ to be the element of $\operatorname{Hom}_{\mathscr{C}}(\mathrm{X}, \mathrm{GFX})$ corresponding to $1_{\mathrm{FX}} \in \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{FX}, \mathrm{FX})$. If $u: \mathrm{X}^{\prime} \rightarrow \mathrm{X}$ is a morphism of $\mathscr{C}$, the morphism $\mathrm{F}(u) \in$ $\operatorname{Hom}_{\mathscr{C}^{\prime}}\left(\mathrm{FX}^{\prime}, \mathrm{FX}\right)$ is the one that can be identified with $\alpha_{\mathrm{X}} \circ u \in \operatorname{Hom}_{\mathscr{C}}\left(\mathrm{X}^{\prime}, \mathrm{GFX}\right)$.

Likewise define $\beta_{\mathrm{Y}} \in \operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{FGY}, \mathrm{Y})$ corresponding to $1_{\mathrm{GY}}$, and for $v \in$ $\operatorname{Hom}_{\mathscr{E}}\left(\mathrm{Y}, \mathrm{Y}^{\prime}\right)$ the morphism $\mathrm{G}(v) \in \operatorname{Hom}_{\mathscr{C}}\left(\mathrm{GY}, \mathrm{GY}^{\prime}\right)$ can be identified with $v \circ$ $\beta_{\mathrm{Y}} \in \operatorname{Hom}_{\mathscr{G}^{\prime}}\left(\mathrm{FGY}, \mathrm{Y}^{\prime}\right)$.

Proposition With this notation, F is fully faithful if and only if $\alpha$ is an isomorphism.
Proof (a) Suppose $\alpha$ is an isomorphism. For $X$ and $X^{\prime}$ in $\mathscr{C}$, the map $u \mapsto \alpha_{\mathrm{X}} \circ u$ is a bijection of $\operatorname{Hom}_{\mathscr{C}}\left(\mathrm{X}^{\prime}, \mathrm{X}\right)$ onto $\operatorname{Hom}_{\mathscr{C}}\left(\mathrm{X}^{\prime}, \mathrm{GFX}\right)=\operatorname{Hom}_{\mathscr{C}^{\prime}}\left(\mathrm{FX}^{\prime}, \mathrm{FX}\right)$, which must be $\mathrm{F}_{\mathrm{X}^{\prime}, \mathrm{X}}$.
(b) Conversely, suppose F is fully faithful. For $\mathrm{X} \in \mathscr{C}$, the map $u \mapsto \alpha_{\mathrm{X}} \circ u=$ $\mathrm{F}(u)$ from $\operatorname{Hom}_{\mathscr{E}}(\mathrm{GFX}, \mathrm{X})$ to $\operatorname{Hom}_{\mathscr{C}}(\mathrm{GFX}, \mathrm{GFX})=\operatorname{Hom}_{\mathscr{C}^{\prime}}(\mathrm{FGFX}, \mathrm{FX})$ is a bijection. Hence there is a unique $w \in \operatorname{Hom}_{\mathscr{E}}(\mathrm{GFX}, \mathrm{X})$ such that $\alpha_{\mathrm{X}} \circ w=1_{\mathrm{GFX}}$. Then $\alpha_{\mathrm{X}} \circ w \circ \alpha_{\mathrm{X}}=\alpha_{\mathrm{X}}$, namely $\mathrm{F}\left(w \circ \alpha_{\mathrm{X}}\right)=\mathrm{F}\left(1_{\mathrm{X}}\right)$, and so $w \circ \alpha_{\mathrm{X}}=1_{\mathrm{X}}$, and $\alpha_{\mathrm{X}}$ is an isomorphism with $\alpha_{\mathrm{X}}^{-1}=w$,

Corollary G is fully faithful if and only if $\beta$ is an isomorphism.
Apply the proposition to opposite categories.
Comments. Hence if F is fully faithful, the functor GF is isomorphic to $1_{\mathscr{C}}$. In other words, for a fully faithful functor, a right adjoint functor is a left inverse up to isomorphism (the side change is only due to conventions regarding composition).

Saying that $F$ is fully faithful means that $F$ is an equivalence between $\mathscr{C}$ and the full subcategory of $\mathscr{C}$ ' A right or left adjoint then gives a "retraction up to isomorphism" from $\mathscr{C}^{\prime}$ onto $\mathscr{C}$. Note that, if F admits a right and a left adjoint, they need not necessarily be isomorphic (2.7, Exercise 3).

For examples, see 2.7.4, Example 1 for the corollary, and 2.7, Exercises 1 and 2 for the proposition.

### 2.7.6 Topology and Algebra

For any compact space X , let $\mathrm{C}(\mathrm{X})$ denote the algebra of continuous functions $\mathrm{X} \rightarrow \mathbb{C}$. This defines a contravariant functor C from the category $\mathscr{K}$ of compact spaces and continuous maps to the category $\mathscr{A}$ of unital associative commutative Banach algebras over $\mathbb{C}$ and unital continuous algebra homomorphisms. To complete the definition of C as a functor, it remains to specify that $f^{*} h=h \circ f$ for $f: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ continuous and $h \in \mathrm{C}\left(\mathrm{X}^{\prime}\right)$.

We will show that the contravariant functor C is fully faithful 2.7.8, and admits an adjoint $S$ (2.7.7, where we will indicate in what sense).
Comments. The functor C thus establishes an analogy between the category $\mathscr{K}$ and the opposite category $\mathscr{A}^{0}$. This is our first encounter with a very general and fruitful situation, where topological or geometric objects are made to correspond to algebraic ones (in Banach algebras, the norm does not really play any role, see 2.7.7, Lemma 2). Thus, the algebra of $\mathrm{C}^{\infty}$-functions from M to $\mathbb{R}$ corresponds to the $\mathrm{C}^{\infty}$-manifold

M , and the algebra of functions induced on X by polynomials, etc. to an algebraic variety.

Concepts or viewpoints can thereby be transferred from a topological or geometric framework to an algebraic one, where their domain of validity is extended. In turn, in a given algebraic situation, thanks to versions of the functor $S$, topological objects can be constructed, enabling us to apply geometric reasoning. This is what Grothendieck's "theory of schemes" consists in (an important part of the initial algebraic data is retained in the structure to avoid any loss of information). More generally, this method underscores every important recent progress in algebra and arithmetic.

As the correspondence is not perfect, the geometric framework is enhanced by "forcing" it (algebraic geometry with nilpotent elements, and more recently "noncommutative geometry").

Chapter 6 addresses situations where this correspondence is reflected by an equivalence of categories (6.2.4 and 6.3.5).

### 2.7.7 Spectra and Gelfand Transforms

We keep the notation of 2.7.6. Let A be a unital associative commutative Banach algebra over $\mathbb{C}$, and $S(A)$ the set of all unital algebra homomorphisms $A \rightarrow \mathbb{C}$, equipped with the weak topology. Then $S(A)$ is a compact space (Lemma 3 below) called the spectrum of A. For $a \in \mathrm{~A}$, define $\hat{a}: \mathrm{S}(\mathrm{A}) \rightarrow \mathbb{C}$ by $\hat{a}(\xi)=\xi(a)$. The function $\hat{a}$ is continuous, and $a \mapsto \hat{a}$ is a homomorphism $\Lambda_{\mathrm{A}}: \mathrm{A} \rightarrow \mathrm{C}(\mathrm{S}(\mathrm{A}))$, called the Gelfand transform.

Proposition The compact space $\mathrm{S}(\mathrm{A})$, together with $\Lambda_{\mathrm{A}}$, satisfies the following universal property:

For all compact spaces X and unital algebra homomorphisms $\Phi: \mathrm{A} \rightarrow \mathrm{C}(\mathrm{X})$, there is a unique continuous map $\varphi: \mathrm{X} \rightarrow \mathrm{S}(\mathrm{A})$ such that $\Phi=\varphi^{*} \circ \Lambda_{\mathrm{A}}$.

This can be restated in terms of categories as follows:
Corollary For $\mathrm{X} \in \mathscr{K}$ and $\mathrm{A} \in \mathscr{A}$, the map $\Psi_{\mathrm{X}, \mathrm{A}}: \varphi \mapsto \varphi^{*} \circ \Lambda_{\mathrm{A}}$ from $\operatorname{Hom}_{\mathscr{K}}(\mathrm{X}, \mathrm{S}(\mathrm{A}))$ to $\operatorname{Hom}_{\mathscr{A}}(\mathrm{A}, \mathrm{C}(\mathrm{X}))$ is a bijection. The family of maps $\Psi_{\mathrm{X}, \mathrm{A}}$ is an isomorphism of functors from $\mathscr{K}^{o} \times \mathscr{A}^{o}$ to $\mathscr{E n s}$. Thus $\mathrm{S}: \mathscr{A}^{o} \rightarrow \mathscr{K}$ is a right adjoint to $\mathrm{C}: \mathscr{K} \rightarrow \mathscr{A}^{o}$.

Comments. Since C and S are contravariant functors, to use the given definition of adjoint functors, they must be made covariant by replacing one of the two categories $\mathscr{K}$ and $\mathscr{A}$ by its opposite. We have chosen to do so for $\mathscr{A}$. Were $\mathscr{K}$ chosen, S would become the left adjoint of C .

The corollary follows from the proposition by a direct application of Proposition 2.7.2. The proof of the proposition uses the following lemmas:

Lemma 2.1 Let $\mathrm{A} \in \mathscr{A}$. Every algebra homomorphism $\varphi: \mathrm{A} \rightarrow \mathbb{C}$ is continuous and has norm $\leq 1$.

Proof Since $\varphi(1)^{2}=\varphi(1), \varphi(1)=1$ or 0 . If $\varphi(1)=0$, then $\varphi=0$; we may assume that $\varphi(1)=1$.

Were $\varphi$ not continuous and its norm $\not \approx 1$, then there would exist $u \in \mathrm{~A}$ such that $|\varphi(u)|>\|u\|$. We may assume that $\varphi(u)=1$. As $\|u\|<1,1-u$ is invertible and $\varphi(1-u)=0$, we get a contradiction.

Lemma 2.2 Let $\mathrm{A} \in \mathscr{A}$ and $\mathrm{X} \in \mathscr{K}$. Every algebra homomorphism $\varphi: \mathrm{A} \rightarrow \mathrm{C}(\mathrm{X})$ is continuous and has norm $\leq 1$.

Proof By Lemma 1, for all $x \in \mathrm{X}$, the map $a \mapsto \varphi(a)(x)$ is continuous and has norm $\leq 1$. So $|\varphi(a)(x)| \leq\|a\|$, and

$$
(\forall a \in \mathrm{~A}) \quad\|\varphi(a)\|=\sup _{x \in \mathrm{X}}|\varphi(a)(x)| \leq\|a\|
$$

Lemma 2.3 Let $\mathrm{A} \in \mathscr{A}$. The set $\mathrm{S}(\mathrm{A})=\operatorname{Hom}_{\mathscr{A}}(\mathrm{A}, \mathbb{C})$, equipped with the weak topology, is compact.

Proof Let $\mathrm{D}_{r}$ to be the closed disk of radius $r\left(\in \mathbb{R}_{+}\right)$in $\mathbb{C}$. By Lemma 1, the space $\mathrm{S}(\mathrm{A})$ is a subspace of $\prod_{a \in \mathrm{~A}} \mathrm{D}_{\|a\|}$ which, by Tychonoff's theorem, is compact. The set $S(A)$ being defined by the equalities

$$
\begin{gathered}
\varphi(a+b)=\varphi(a)+\varphi(b) \\
\varphi(a b)=\varphi(a) \cdot \varphi(b), \quad \varphi(1)=1
\end{gathered}
$$

is closed in $\prod_{a \in \mathrm{~A}} \mathrm{D}_{\|a\|}$, and so is compact.
Proof of the Proposition The set $\operatorname{Hom}_{\mathscr{H}}(\mathrm{X}, \mathrm{S}(\mathrm{A}))$ of all continuous maps from X to $S(A)$ equipped with the weak topology, can be identified with the set of functions $f: \mathrm{X} \times \mathrm{A} \rightarrow \mathbb{C}$ satisfying:
(i) for $a \in \mathrm{~A}$, the map $x \mapsto f(x, a)$ is continuous;
(ii) for $x \in \mathrm{X}, a \mapsto f(x, a)$ is a (necessarily continuous) unital algebra homomorphism.

This set can be identified to the set of (necessarily continuous) unital algebra homomorphisms $\mathrm{A} \rightarrow \mathrm{C}(\mathrm{X})$. Hence

$$
\operatorname{Hom}_{\mathscr{L}}(\mathrm{X}, \mathrm{~S}(\mathrm{~A}))=\operatorname{Hom}_{\mathscr{L}}(\mathrm{A}, \mathrm{C}(\mathrm{X}))
$$

and, for $\mathrm{X}=\mathrm{S}(\mathrm{A})$, the homomorphism corresponding to the identity element is readily seen to be the Gelfand transform.

### 2.7.8

## Proposition The contravariant functor $\mathrm{C}: \mathscr{K} \rightarrow \mathscr{A}$ is fully faithful.

Proof Thanks to Proposition 2.7.7, all we need to prove is that, for a compact space X , the natural map $\mathrm{X} \rightarrow \mathrm{S}\left(\mathrm{C}(\mathrm{X})\right.$ ) is a homeomorphism. This map is $x \mapsto \delta_{x}$, where $\delta_{x}(f)=f(x)$. It is clearly continuous. Since X is compact, it is sufficient to show that it is bijective, i.e. that every homomorphism $\xi: \mathrm{C}(\mathrm{X}) \rightarrow \mathbb{C}$ is of type $\delta_{x}$ for some unique $x \in \mathrm{X}$.

Hence let $\xi: \mathrm{C}(\mathrm{X}) \rightarrow \mathbb{C}$ be a homomorphism, $\mathfrak{m}$ its kernel and $\mathrm{V}(\mathfrak{m})$ the intersection of all $h^{-1}(0)$ for $h \in \mathfrak{m}$. Were the set $\mathrm{V}(\mathfrak{m})$ empty, there would be a finite family $h_{1}, \ldots, h_{k}$ in $\mathfrak{m}$ without any common roots; then $h_{1} h_{1}+\cdots+h_{k} h_{k}$ would be a function everywhere $>0$, hence invertible, and so $1 \in \mathfrak{m}$, an absurdity.

Let $x \in \mathrm{~V}(\mathfrak{m})$. Then $\operatorname{Ker} \xi \subset \operatorname{Ker} \delta_{x}$, but both are hyperplanes, and so $\operatorname{Ker} \xi=$ $\operatorname{Ker} \delta_{x}$. Since $\xi(1)=\delta_{x}(1)=1, \xi=\delta_{x}$. Hence the map $x \mapsto \delta_{x}$ is surjective.

Injectivity follows from Urysohn's theorem: if $x$ and $y$ are two distinct points of X , then there is a continuous function $h: \mathrm{X} \rightarrow \mathbb{C}$ for which $h(x) \neq h(y)$. So $\delta_{x}(h) \neq \delta_{y}(h)$, which implies that $\delta_{x} \neq \delta_{y}$.

## Exercises 2.7. (Adjoint functors)

1.-Let G be a group. The commutator group of G is the subgroup $\mathrm{G}^{\prime}$ of G generated by all elements $x y x^{-1} y^{-1}$, where $x, y \in \mathrm{G}$.
(a) Show that $\mathrm{G}^{\prime}$ is normal and that $\mathrm{A}(\mathrm{G})=\mathrm{G} / \mathrm{G}^{\prime}$ is commutative.
(b) Show that the functor $\mathrm{G} \rightarrow \mathrm{A}(\mathrm{G})$ from $\mathscr{G} \boldsymbol{r}$ to the category $\mathscr{A} \mathscr{\mathscr { C }}$ of Abelian groups is the adjoint to the inclusion functor $\mathscr{A} \mathscr{C}$ in $\mathscr{\mathscr { G }}$; determine whether it is a left or right one; state this property in elementary terms.
(c) Show that there is no adjoint functor on the other side.
2.-Let $X$ be a topological space, I the segment $[0,1]$ and $E$ the set of continuous maps from X to I . Define a map $\delta: \mathrm{X} \rightarrow \mathrm{I}^{\mathrm{E}}$ by $\delta(x)=(f(x))_{f \in \mathrm{E}}$. Equip $\mathrm{I}^{\mathrm{E}}$ with the product topology. Then $\delta$ is continuous. The Cech compactification of X , written $\check{\mathrm{X}}$, is the closure of $\delta(\mathrm{X})$ in $\mathrm{I}^{\mathrm{E}}$.
(a) The space X is said to be completely regular if, for all all neighbourhoods U of all elements $x \in \mathrm{X}$, there is a continuous function $f: \mathrm{X} \rightarrow \mathrm{I}$ such that $f(x)=1$ and $\left.f\right|_{\mathrm{Cu}}=0$. Show that:

- if X is completely regular, then $\delta$ induces a homeomorphism from X onto a subset of X̌;
- if X is compact, then $\check{\mathrm{X}}=\mathrm{X}$;
- if $X$ is locally compact, then $X$ is open in $\check{X}$.
(b) Let K be a compact space, and $f: \mathrm{X} \rightarrow \mathrm{K}$ a continuous map. Show that there is a continuous map $\check{f}: \check{\mathrm{X}} \rightarrow \mathrm{K}$ such that $f=\check{f} \circ \delta$. State this property in terms of adjoint functors.
(c) Is there an adjoint functor on the other side?
3.-Let $\mathscr{T} a p l c$ be the category of locally connected topological spaces, and identify $\mathscr{E} n s$ with the full subcategory of discrete spaces. Show that the inclusion functor $\iota: \mathscr{E n s} \rightarrow \mathscr{T a p l} \subset$ has both a left and a right adjoint, and that these are not equal.

4. (Fourier series.) -Consider the Banach space $\ell^{1}(\mathbb{Z})$ consisting of complex sequences $\mathbf{a}=\left(a_{n}\right)_{n \in \mathbb{Z}}$ such that $\sum\left|a_{n}\right|<\infty$, equipped with the norm defined by $\|\mathbf{a}\|=\sum_{n \in \mathbb{Z}}\left|a_{n}\right|$. Set $a(n)=a_{n}$, and define the basis element $\mathbf{e}_{\mathbf{n}}$ by $\mathbf{e}_{\mathbf{n}}(p)=1$ if $p=n$ and 0 otherwise.
(a) Show that there is a unique continuous bilinear multiplication law $*$ on $\ell^{1}(\mathbb{Z})$ such that $\mathbf{e}_{\mathbf{p}} * \mathbf{e}_{\mathbf{q}}=\mathbf{e}_{\mathbf{p}+\mathbf{q}}$ for $p, q \in \mathbb{Z}$. This law is called convolution. Let A be the algebra obtained by equipping $\ell^{1}(\mathbb{Z})$ with it. Write $c_{n}$ in terms of the elements $a_{p}$ and $b_{q}$ when $\mathbf{c}=\mathbf{a} * \mathbf{b}$.
(b) Show that if $\xi: \mathrm{A} \rightarrow \mathbb{C}$ is a unital homomorphism, then $\left\|\xi\left(\mathbf{e}_{\mathbf{1}}\right)\right\|=1$. Show that $S(A)$ can be identified with the unit circle of $\mathbb{C}$, i.e. with $\mathbb{T}=\mathbb{R} / \mathbb{Z}$.

Show that, for $\mathbf{a} \in \mathrm{A}$, the Gelfand transform $\widehat{\mathbf{a}}: \mathbb{T} \rightarrow \mathbb{C}$ is given by the formula for Fourier series $\widehat{\mathbf{a}}(t)=\sum_{n \in \mathbb{Z}} a_{n} e^{2 i \pi n t}$.
(c) Show that the elements $a_{n}$ can be retrieved from $\widehat{\mathbf{a}}$ by the inverse Fourier formula $a_{n}=\int_{\mathbb{T}} \widehat{\mathbf{a}}(t) e^{-2 i \pi n t} d t$. Deduce that the Gelfand transform of $\Lambda: \mathrm{A} \rightarrow \mathrm{C}(\mathbb{T})$ is injective.

For $g \in \mathbb{L}^{1}(\mathbb{T}, \mathbb{C})$, i.e. integrable over $\mathbb{T}$ with values in $\mathbb{C}$, set $\hat{g}(n)=\int_{\mathbb{T}} g(t) e^{-2 i \pi n t}$ $d t$ ( $n$th Fourier coefficient of $g$ ).
(d) The aim is to show that $\Lambda$ is not surjective. We remind the reader of the

Banach Theorem The inverse of a continuous bijective linear map between Banach spaces is continuous.

Let $g: \mathbb{T} \rightarrow \mathbb{C}$ be the function taking value 1 on $[0,1 / 2]$ and -1 on $] 1 / 2,1]$. Show that $\sum|\hat{g}(n)|=\infty$.

For $k>4$, let $f_{k}$ be the continuous function, affine on $[-1 / k, 1 / k]$ and on $[1 / 2-1 / k, 1 / 2+1 / k]$, and agreeing with $g$ elsewhere. Show that $\sum_{n} \widehat{f}_{k}(n)$ tends to $\hat{g}(n)$ as $k \rightarrow \infty$ for fixed $n$. Deduce that the upper bounds of $\sum_{n}\left|\hat{f}_{k}(n)\right|$ depend on $k$. Conclude.
(e) Give an explicit example of a continuous function on $\mathbb{T}$ which is not the Gelfand (here Fourier) transform of $\ell^{1}(\mathbb{Z})$.

### 2.8 Profinite Spaces

In this section and the next one, we study profinite spaces and groups with infinite Galois theory in mind.

### 2.8.1 Pro-objects of a Category

Let $\mathscr{C}$ be a category. Define a category $\mathscr{C}$ as follows: an object of $\mathscr{C}$ is an inverse system $^{2}\left(\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)_{i \leq j}\right)$ in $\mathscr{C}$. If $\mathrm{X}=\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}$, and $\mathrm{Y}=\left(\mathrm{Y}_{j}\right)_{j \in \mathrm{~J}}$ are two objects of
 $\left(\mathrm{Z}_{k}\right)_{k \in \mathrm{~K}}$ be inverse systems, $\varphi=\left(\varphi_{j}\right) \in \operatorname{Hom}(\mathrm{X}, \mathrm{Y})$, where $\varphi_{j} \in \underset{\lim }{\operatorname{Hom}}\left(\mathrm{X}_{i}, \mathrm{Y}_{j}\right)$, and $\psi=\left(\psi_{k}\right) \in \operatorname{Hom}(Y, Z)$. Define $\psi \circ \varphi$ as follows: for all $k \in \overrightarrow{K, \text { let }}(j, v)$, where $j \in \mathrm{~J}$ and $v \in \operatorname{Hom}\left(\mathrm{Y}_{j}, \mathrm{Z}_{k}\right)$, be a representative of $\psi_{k}$, and define $\omega_{k}=\nu \circ \varphi_{l} \in$ $\underset{\longrightarrow}{\lim } \operatorname{Hom}\left(\mathrm{X}_{l}, \mathrm{Z}_{k}\right)$. The element $\omega_{k}$ is independent of the choice of $(j, v)$.

Indeed, if $\left(j^{\prime}, v^{\prime}\right)$ is another representative of $\psi_{k}$, then there exists $j^{\prime \prime}$ greater than $j$ and $j^{\prime}$ such that $v \circ g_{j}^{j^{\prime}}=v^{\prime} \circ g_{j^{\prime}}^{j^{\prime \prime}}$, and so $v \circ \varphi_{j}=v \circ g_{j}^{j^{\prime}} \circ \varphi_{j^{\prime \prime}}=v^{\prime} \circ g_{j^{\prime}}^{j^{\prime \prime}} \circ$ $\varphi_{j^{\prime \prime}}=v^{\prime} \circ \varphi_{j^{\prime}}$. Then $\omega_{k}=h_{k}^{k^{\prime}} \circ \omega_{k^{\prime}}$. Thus the family $\left(\omega_{k}\right)_{k \in \mathrm{~K}}$ defines an element $\psi \circ \varphi=\omega \in \operatorname{Hom}(\mathrm{X}, \mathrm{Z})$. Checking that $\mathscr{C}$ is a category is tedious but not difficult.

The objects of $\mathscr{C}$ are called pro-objects de $\mathscr{C}$. Considering an object of $\mathscr{C}$ as an inverse system indexed by a singleton identifies $\mathscr{C}$ with a full subcategory of $\mathscr{C}$. A functor $\underset{\leftarrow}{F}: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$ corresponds to a functor $F: \mathscr{C} \rightarrow \mathscr{C}^{\prime}$. If $F$ is an equivalence of categories, so is $\underset{\leftarrow}{\leftarrow}$.

### 2.8.2 Totally Disconnected Spaces

A topological space X is said to be totally disconnected if for any two distinct elements $x$ and $y$ of X , there is a clopen (both open and closed) subset U of X such that $x \in \mathrm{U}$ and $y \notin \mathrm{U}$.

Examples The discrete spaces, the field $\mathbb{Q}$ of rationals, the Cantor set $\{0\} \cup\{1 / n\}_{n \in \mathbb{N}}$, are totally disconnected spaces.

Every subspace of a totally disconnected space is totally disconnected. Every product $\mathrm{E}=\prod \mathrm{E}_{i}$ of totally disconnected spaces is totally disconnected. Indeed, let $x$ and $y$ be distinct elements of E . Their projections on at least some $\mathrm{E}_{i}$ must be different. Let $i$ be such that $x_{i} \neq y_{i}$ and U be a clopen subset in $\mathrm{E}_{i}$ containing $x$ but not $y$. Then the inverse image of U in E is clopen in E , contains $x$ but not $y$.

Every inverse limit of totally disconnected spaces is totally disconnected.

[^1]
### 2.8.3

Proposition and Definition Let X be a topological space. The following conditions are equivalent:
(i) X is compact and totally disconnected;
(ii) X is homeomorphic to the inverse limit of an inverse system of finite discrete sets.

If these conditions are satisfied, then X is said to be a profinite space.
Proof (ii) $\Rightarrow$ (i). Let $\left(\mathrm{X}_{i}\right)$ be an inverse system of finite discrete sets such that $\lim X_{i}=X$. Then X is closed in the product $\prod X_{i}$, which is compact by Tychonoff's theorem (1.7.7) and totally disconnected since so are all $X_{i}$. Hence X is compact and totally disconnected.
(i) $\Rightarrow$ (ii). Let $\Lambda$ be the set of equivalence relations $\lambda$ on X such that the quotient $\mathrm{X}_{\lambda}$ is finite and discrete, and for all $x \in \mathrm{X}$ denote by $x_{\lambda}$ the class of $x$ for the relation $\lambda$. Write $\lambda \leq \mu$ if $\mu$ is finer than $\lambda$. The set $\Lambda$ is a directed set. Indeed, $\forall \lambda \in \Lambda$ $\forall \lambda^{\prime} \in \Lambda$ the map $x \mapsto\left(x_{\lambda}, x_{\lambda^{\prime}}\right)$ from X to $\mathrm{X}_{\lambda} \times \mathrm{X}_{\lambda^{\prime}}$ defines an equivalence relation $\mu$ on X finer than $\lambda$ and $\lambda^{\prime}$.

Let $\mathrm{X}_{\lambda}$ form an inverse system. For all $x \in \mathrm{X},\left(x_{\lambda}\right)_{\lambda \in \Lambda} \in \lim \mathrm{X}_{\lambda}$. We show that the map $\theta: x \mapsto\left(x_{\lambda}\right)_{\lambda \in \Lambda}$ from X to $\lim _{\leftarrow} \mathrm{X}_{\lambda}$ is a homeomorphism. It is obviously continuous. Since X is compact, it is image-closed and if $\theta$ is injective, $\theta$ is a homeomeorphism onto its image. Hence it is sufficient to show that $\theta$ is injective and image-dense. Let $x \neq y \in \mathrm{X}$ and let ( $\mathrm{U}, \mathrm{V}$ ) be a partition of X into two open subsets such that $x \in \mathrm{U}$ and $y \in \mathrm{~V}$. For the relation $\lambda$ associated to this partition, $x_{\lambda} \neq y_{\lambda}$, and so the images of $x$ and $y$ in $\lim _{\longleftarrow} X_{\lambda}$ are distinct. Consequently $\theta$ is injective.

The homomorphism $\theta$ is image-dense because of the following lemma:
Lemma Let X be a set, $\left(\mathrm{Y}_{\lambda}\right)$ an inverse system of topological spaces, $f: \mathrm{X} \rightarrow \lim _{\leftarrow} \mathrm{Y}_{\lambda}$ a map. If, for all $\lambda$, the map $f_{\lambda}: \mathrm{X} \rightarrow \mathrm{Y}_{\lambda}$ induced by $f$ is surjective, then $\overleftarrow{f}$ is image-dense.

Let $a=\left(a_{\lambda}\right)_{\lambda \in \Lambda} \in \lim \mathrm{Y}_{\lambda}$. For every finite family $\lambda_{1}, \ldots, \lambda_{k} \in \Lambda$ the set $\mathrm{V}_{\lambda_{1}, \ldots, \lambda_{k}}$ $(a)=\left\{b=\left(b_{\lambda}\right)_{\lambda \in \Lambda},(\forall i=1, \ldots, k) b_{\lambda_{i}}=a_{\lambda_{i}}\right\}$ is a neighbourhood of $a$ and the sets $\mathrm{V}_{\lambda_{1}, \ldots, \lambda_{k}}(a)$ form a fundamental system of neighbourhoods of $a$. Let $\mathrm{L}=\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ be a finite subset of $\Lambda$ and $\mu$ an upper bound of L in $\Lambda$. Let $x \in \mathrm{X}$ be such that $f_{\mu}(x)=a_{\mu}$; for all $\lambda_{i} \in \mathrm{~L}, f_{\lambda_{i}}(x)=a_{\lambda_{i}}$ and $\operatorname{Im} f \cap \mathrm{~V}_{\lambda_{1}, \ldots, \lambda_{k}}(a)$ is non-empty. Hence $f$ is image-dense.

### 2.8.4

Let $\mathscr{E}$ be the category of profinite spaces and continuous maps and $\operatorname{set} \mathcal{f}$ that of finite sets. The category $\mathscr{E}$ is a full subcategory of $\mathscr{T o p}$.

Theorem The functor $\mathrm{F}:\left(\mathrm{X}_{i}\right) \rightarrow \lim _{\leftarrow} \mathrm{X}_{i}$ from $\operatorname{Setf}$ to $\mathscr{E}$ is an equivalence of categories.

Definition 2.8 .3 states that F is essentially surjective. We show (2.8.6) that F is fully faithful. For this the following lemma is needed:

### 2.8.5

Lemma Let $\left(\mathrm{X}_{i}, f_{i}^{j}\right)$ be an inverse system of finite sets. Denote by $\left(\mathrm{X}_{\infty}, f_{i}^{\infty}\right)$ its inverse limit. Then

$$
(\forall i)(\exists j \geq i) \operatorname{Im} f_{i}^{\infty}=\operatorname{Im} f_{i}^{j} .
$$

Proof Let $x \in \mathrm{X}_{i}$. For all $j \geq i$, set $\mathrm{Y}_{j}=\left(f_{i}^{j}\right)^{-1}(x)$. These form an inverse system of finite sets and $\lim _{\longleftarrow} \mathrm{Y}_{j}=\left(f_{i}^{\infty}\right)^{-1}(x)$. Since

$$
x \in \cap \operatorname{Im} f_{i}^{j} \Longleftrightarrow \forall j \geq i \mathrm{Y}_{j} \neq \varnothing \Longleftrightarrow \lim _{\hookleftarrow} \mathrm{Y}_{j} \neq \varnothing \Longleftrightarrow x \in \operatorname{Im} f_{i}^{\infty}
$$

by $2.5 .14, \bigcap_{j \geq i} \operatorname{Im} f_{i}^{j}=\operatorname{Im} f_{i}^{\infty}$. The family $\left(\operatorname{Im} f_{i}^{j}\right)_{j \geq i}$ is a decreasing family of finite sets, whose intersection is equal to one of them.

### 2.8.6

Proposition Let $\left(\mathrm{X}_{i}, f_{i}^{j}\right)$ be an inverse system of finite sets and Y a finite set. Then,

$$
\left.\operatorname{Hom}_{\mathscr{T} a p}\left(\underset{\leftrightarrows}{\left(\lim X_{i}\right.} ; \mathrm{Y}\right)=\underset{\mathscr{E n s}}{\lim } \operatorname{Hom}_{i} ; \mathrm{Y}\right) .
$$

Proof Let $\varphi \in \lim _{\operatorname{Hom} \mathscr{g n s}^{\prime}}\left(\mathrm{X}_{i} ; \mathrm{Y}\right)$. Represent $\varphi$ by $\varphi_{i}: \mathrm{X}_{i} \rightarrow \mathrm{Y}$. The map $x=$ $\left(x_{k}\right) \mapsto \varphi_{i}\left(x_{i}\right)$ from $\lim X_{i}$ to Y does not depend on the choice of the representative $\varphi_{i}$ of $\varphi$. Denote this map by $\alpha(\varphi)$. We show that $\alpha$ is injective and surjective.

Let $\varphi$ and $\psi$ be such that $\alpha(\varphi)=\alpha(\psi)$. Then $\varphi_{i}$ and $\psi_{i}$ are equal on $\operatorname{Im} f_{i}^{\infty}$, hence, for some $j$, also on $\operatorname{Im} f_{i}^{j}$ by Lemma 2.8.5. So $\varphi=\psi$.

Let $\varphi: \mathrm{X}_{\infty} \rightarrow \mathrm{Y}$ be a continuous map, $\Gamma_{\varphi}\left(\right.$ resp. $\left.\Gamma_{i}\right)$ the graph of the equivalence relation defined by $\varphi$ (resp. $f_{i}^{\infty}$ ) i.e. the inverse image of the diagonal of $\mathrm{Y} \times \mathrm{Y}$ under the map $(\varphi, \varphi): \mathrm{X}_{\infty} \times \mathrm{X}_{\infty} \rightarrow \mathrm{Y} \times \mathrm{Y}$ (resp. of the diagonal of $\mathrm{X}_{i} \times \mathrm{X}_{i}$ under $\left(f_{i}^{\infty}, f_{i}^{\infty}\right)$ ). These graphs are all clopen in $\mathrm{X}_{\infty} \times \mathrm{X}_{\infty}$ and, for all $i$ and $j \geq i, \Gamma_{j} \subset$ $\Gamma_{i}$. Let $\Delta$ be the diagonal of $\mathrm{X}_{\infty} \times \mathrm{X}_{\infty}$. Then, $\Delta \subset \Gamma_{\varphi}$, and $\Delta=\cap_{i} \Gamma_{i}$ is closed.

Therefore

$$
\varnothing=\left(\bigcap_{i} \Gamma_{i}\right) \cap \complement \Gamma_{\varphi}=\bigcap_{i}\left(\Gamma_{i} \cap \complement \Gamma_{\varphi}\right) .
$$

The intersection of the decreasing closed sets $\Gamma_{i} \cap \subset \Gamma_{\varphi}$ is empty, hence so is one of them. Thus, there is an index $j$ such that $\Gamma_{j} \subset \Gamma_{\varphi}$. Thus $f_{j}^{\infty}(x)=f_{j}^{\infty}(y) \Rightarrow \varphi(x)=\varphi(y)$.

Therefore there is a map from $\operatorname{Im} f_{j}^{\infty}$ to Y making the diagram

commutative.
By Lemma 2.8.5 there exists $k \geq j$ such that $\operatorname{Im} f_{j}^{\infty}=\operatorname{Im} f_{j}^{k}$. Composing with $f_{j}^{k}$ then gives a map $\varphi^{k}$ from $\mathrm{X}_{k}$ to Y such that the diagram

commutes.
Corollary The functor F of theorem 2.8 .4 is fully faithful.
Indeed, if $\left(\mathrm{X}_{i}\right)$ and $\left(\mathrm{Y}_{j}\right)$ are two inverse systems of finite sets, then

$$
\begin{aligned}
& \operatorname{Hom}_{\mathscr{T a p}}\left(\lim _{\leftarrow} X_{i} ; \lim _{\leftarrow} \mathrm{Y}_{j}\right)=\lim _{\overleftarrow{j}} \operatorname{Hom}_{\mathscr{T a p}}\left(\lim _{\leftarrow} X_{i} ; \mathrm{Y}_{j}\right)
\end{aligned}
$$

## Exercises 2.8. (Profinite spaces)

1.-Let $X$ be a compact space.
(a) Let Y be closed in X and let $\left(\mathrm{V}_{i}\right)_{i \in \mathrm{I}}$ be a directed family of closed sets of X such that $\cap \mathrm{V}_{i}=\mathrm{Y}$. Show that every neighbourhood of Y contains some $\mathrm{V}_{i}$.
(b) Let $a \in \mathrm{X}$. Show that the intersection of the clopen sets of X containing $a$ is the connected component of $a$ in X .
(c) A compact space is totally discontinuous as defined in (2.8.2) if and only if every connected component is reduced to a point.
(d) For the equivalence given in (c), can the assumption of compactness of X be omitted?
2. (Cantor set.)—Define a sequence $\left(\mathrm{A}_{n}\right)_{n \in \mathbb{N}}$ of subsets of $[0,1]$ as follows:

$$
\begin{aligned}
& \mathrm{A}_{0}=[0,1], \\
& \left.\mathrm{A}_{1}=[0,1]-\right] 1 / 3,2 / 3[, \\
& \mathrm{A}_{2}=\mathrm{A}_{1}-(] 1 / 9,2 / 9[\cup] 7 / 9,8 / 9[) .
\end{aligned}
$$

Each $\mathrm{A}_{n}$ is a finite union of closed disjoint intervals and $\mathrm{A}_{n+1}$ is the third of the open middle of $\mathrm{A}_{n}$. The set $\mathrm{K}=\cap_{n \in \mathbb{N}} \mathrm{~A}_{n}$ is called the triadic Cantor set.
(a) Show that K is a profinite space.
(b) Show that K is the set of elements $x$ of $[0,1]$ having a base 3 expansion without the number 1.
(c) Show that $K$ is homeomorphic to $\{0,1\}^{\mathbb{N}}$.
(d) Write $[0,1]-\mathrm{K}$ as the disjoint union of intervals $] \alpha_{i}, \beta_{i}[i \in \mathrm{I}$, where I is a countable set. Show that the three sets $\left\{\alpha_{i}\right\},\left\{\beta_{i}\right\}$ and $\mathrm{K}-\left(\left\{\alpha_{i}\right\} \cup\left\{\beta_{i}\right\}\right)$ are dense in K and that the interior of K in $[0,1]$ is empty.
(e) Show that the quotient of K by the equivalence relation which equates $\alpha_{i}$ and $\beta_{i}$ is homeomorphic to $[0,1]$.
(f) Let $\left(\mathrm{X}_{n}\right)_{n \in \mathbb{N}}$ be an inverse system of finite sets. Suppose that for all $n, f_{n}^{n+1}$ : $\mathrm{X}_{n+1} \rightarrow \mathrm{X}_{n}$ is such that, $\forall x \in \mathrm{X}_{n},\left(f_{n}^{n+1}\right)^{-1}(x)$ contains at least two point. Show that K is homeomorphic to $\lim _{\leftrightarrows} \mathrm{X}_{i}$.
(g) Show that every metrizable profinite space without isolated points is homeomorphic to K .
3.-Let K be the Cantor set and $\gamma \in \mathrm{K}-\left(\left\{\alpha_{i}\right\} \cup\left\{\beta_{i}\right\}\right)$ (see Exercise 2 for notation $)$. Set $X=((K \cap[0, \gamma]) \times\{0\}) \cup((K \cap[\gamma, 1]) \times\{1\})$. Show that the projection $(k, i) \mapsto k$ from X to K has no continuous section.
4.-Let $\mathscr{C}$ be a category having inverse limits.
(a) Define a functor $\lambda:\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}} \mapsto \lim _{\leftarrow} \mathrm{X}_{i}$ from $\mathscr{C}$ to $\mathscr{C}$. Let $\iota$ be the inclusion functor $\mathscr{C} \rightarrow \mathscr{C}$. Show that $\lambda \circ \iota \approx \mathrm{I}_{\mathscr{C}}$, but that in general $\lambda$ and $\iota$ are not equivalences of categories.
(b) Let $\mathscr{C}^{\prime}$ be another category. Show that for every functor $\mathrm{F}: \mathscr{C}^{\prime} \rightarrow \mathscr{C}$, the functor $\underset{\leftarrow}{\leftarrow}$ from $\mathscr{C}^{\prime}$ to $\mathscr{C}$ commuting with inverse limits of objects of $\mathscr{C}$ is its unique extension. Suppose that $\mathscr{C}^{\prime}$ has inverse limits: does the diagram

necessarily commute?
(c) Show that, if $\mathscr{C}=\operatorname{Set}$, the inclusion $\mathscr{C} \rightarrow \mathscr{C}$ is not an equivalence of categories.

### 2.9 Profinite Groups

### 2.9.1

Proposition and Definition Let G be a topological group. The following conditions are equivalent:
(i) G is the inverse limit of finite discrete groups;
(ii) the underlying topological of G is profinite.

## If these conditions are satisfied, G is said to be a profinite group.

Proof (i) $\Rightarrow$ (ii) is obvious. We show that (ii) $\Rightarrow$ (i).
For any open normal subgroup $H$ of $G$, the quotient $G / H$ is discrete (because the inverse image of each point under the canonical projection is open) and compact (because it is the continuous image of the compact space G in the Hausdorff space
 limit of the set of open normal subsets, relative to the inclusion order) is imagedense by Lemma 2.8.3; as G is compact, $\gamma(\mathrm{G})$ is closed and equal to its closure, so that $\gamma$ is surjective. We show it is injective. Suppose that $x \in \mathrm{G}$ is not the identity element $e$. There is an open normal subgroup H not containing $x$. Indeed, as G is totally discontinuous, there is an open and closed subset U in G such that $e \in \mathrm{U}$ and $x \notin \mathrm{U}$. Set $\mathrm{V}=\{g \in \mathrm{G}, g \mathrm{U}=\mathrm{U}\}, \mathrm{K}=\{g \in \mathrm{G}, g \mathrm{U} \not \subset \mathrm{U}\}$ and $\mathrm{K}^{\prime}=\{g \in \mathrm{G}, \mathrm{U} \not \subset$ $g \mathrm{U}\}$. Then $\mathrm{V}=\mathrm{G}-\left(\mathrm{K} \cup \mathrm{K}^{\prime}\right)$ and $\mathrm{V} \subset \mathrm{U}$. As the images of $\mathrm{U} \times(\mathrm{G}-\mathrm{U})$ under the respective continuous maps $(x, y) \mapsto y x^{-1}$ and $x y^{-1}$, the subspaces K and $\mathrm{K}^{\prime}$ of G are compact. Hence V is open in G . The subgroup $\mathrm{H}=\cap_{g \in \mathrm{G}} g \mathrm{~V} g^{-1}$ is normal, its complement $\mathrm{G}-\mathrm{H}$ is the image of $\mathrm{G} \times(\mathrm{G}-\mathrm{V})$ under $(g, x) \mapsto g x g^{-1}$, i.e. the image of a compact set under a continuous map. $\mathrm{So} \mathrm{G}-\mathrm{H}$ is compact and H is open. Since G is compact, $\gamma$ is a homeomorphism.

### 2.9.2

 groups.

Proposition The functor $\mathrm{G}:\left(\mathrm{G}_{i}\right) \rightarrow \underset{\longleftarrow}{\lim \mathrm{G}_{i}}$ is an equivalence from $\mathscr{G}_{\wedge \mathcal{R}}$ onto the category of profinite groups (full subcategory of $\mathscr{E r ® O}_{\boldsymbol{T}}$ ).

Proof By 2.9.1, this functor is essentially surjective. We show that it is fully faithful,
 Passing to the inverse limit, this reduces to the following lemma:

Lemma Let $\left(\mathrm{G}_{i}\right)$ be an inverse system of finite groups, $\mathrm{G}=\lim \mathrm{G}_{i}$ and H a finite group. Then, $\operatorname{Hom}_{\mathscr{G}_{r} \mathscr{F o p}_{\boldsymbol{p}}}(\mathrm{G}, \mathrm{H})=\underline{\longrightarrow} \lim _{\mathscr{G}_{\boldsymbol{r}}}\left(\mathrm{G}_{i}, \mathrm{H}\right)$.

Proof of the Lemma The canonical map $\underset{\rightarrow}{\lim } \operatorname{Hom}\left(\mathrm{G}_{i}, \mathrm{H}\right) \rightarrow \operatorname{Hom}_{\mathcal{G}_{r} \mathscr{F}_{\text {op }}}(\mathrm{G}, \mathrm{H})$ is injective since it is induced by the canonical bijection $\xrightarrow{\lim } \operatorname{Hom}_{\mathscr{\varepsilon} n s}\left(\mathrm{G}_{i}, \mathrm{H}\right) \rightarrow$ $\operatorname{Hom}_{\mathscr{T o p p}}(\mathrm{G}, \mathrm{H})$ (2.8.6). We show that it is surjective, i.e. that for any continuous homomorphism $\varphi: \mathrm{G} \rightarrow \mathrm{H}$, there is an index $j$ and a continuous homomorphism $\varphi_{j}: \mathrm{G}_{j} \rightarrow \mathrm{H}$ such that $\varphi=\varphi_{j} \circ f_{j}^{\infty}$. By 2.8 .6 there is an index $i$ and a continuous map $\varphi_{i}: \mathrm{G}_{i} \rightarrow \mathrm{H}$ such that $\varphi=\varphi_{i} \circ f_{i}^{\infty}$; the map $\varphi_{i}$ is not necessarily a group homomorphism, but its restriction to $f_{i}^{\infty}(\mathrm{G})$ is. By $2.8 .5, f_{i}^{j}\left(\mathrm{G}_{j}\right)=f_{i}^{\infty}(\mathrm{G})$ for some $j \geq i$. So $\varphi_{j}=\varphi_{i} \circ f_{i}^{j}: \mathrm{G}_{j} \rightarrow \mathrm{H}$ is a homomorphism and $\varphi=\varphi_{j} \circ f_{j}^{\infty}$,

### 2.9.3

Let G be a topological group and X a topological space. A continuous operation from G onto X is a continuous map $(g, x) \mapsto g \cdot x$ from $\mathrm{G} \times \mathrm{X}$ to X which is an action from $G$ onto $X$ (2.1.2, Example 3).

Remark If X is finite and discrete, an action from G onto X is continuous if and only if, for all $x \in \mathrm{X}$, the stabilizer of $x$ is an open subgroup of G or equivalently that the homomorphism $g \mapsto(x \mapsto g \cdot x)$ from G to the group $\mathfrak{S}(\mathrm{X})$ of permutations of X equipped with the discrete topology is continuous.

### 2.9.4

Let $G$ be a profinite group. Denote by G- $\operatorname{set} \boldsymbol{f}$ the category of finite sets on which $G$ acts continuously. If $G$ is finite, the category G- $\operatorname{Setf}$ is a full subcategory of the category G-E्Ens (2.1.2, Example 3). Let $\left(\left(\mathrm{G}_{i}\right),\left(f_{i}^{j}\right)\right)$ be an inverse system of finite groups. There is a direct system $\left(\mathrm{G}_{i}-\boldsymbol{\mathcal { S } e t} \boldsymbol{\mathcal { F }}\right)$ of categories corresponding to it (2.2.6). Set $\mathrm{G}=\underset{\longleftarrow}{\lim } \mathrm{G}_{i}$.

Proposition The category G - $\operatorname{Set} \mathcal{f}$ is the direct limit of $\mathrm{G}_{i}-\operatorname{Set} \mathcal{P}$.
Proof An object of G-Setf is a pair ( $\mathrm{X}, \rho$ ), where $\mathrm{X} \in \operatorname{Setf}$ and $\rho \in \operatorname{Hom}_{\mathscr{E}_{r} \mathscr{T} a p}$ (G, $\mathfrak{S}(\mathrm{X})$ ). Applying Lemma 2.9.2 to $\mathrm{H}=\mathfrak{S}(\mathrm{X})$ gives $\operatorname{Hom}_{\mathscr{g}_{r \mathcal{G}}(\mathrm{G}}(\mathrm{G}, \mathfrak{S}(\mathrm{X}))=$ $\underset{i}{\lim } \operatorname{Hom}_{\mathscr{G}_{r}}\left(\mathrm{G}_{i}, \mathfrak{S}(\mathrm{X})\right.$ ). Therefore,

$$
\begin{aligned}
& \bigsqcup_{\mathrm{X} \in \mathcal{S e t f}} \operatorname{Hom}_{\mathscr{G}_{r} \mathscr{T} a p}(\mathrm{G}, \mathfrak{S}(\mathrm{X}))=\bigsqcup_{\mathrm{X}} \underset{i}{\underset{i}{\lim } \operatorname{Hom}\left(\mathrm{G}_{i}, \mathfrak{S}(\mathrm{X})\right)} \\
& =\underset{i}{\lim } \bigsqcup_{\mathrm{X}} \operatorname{Hom}\left(\mathrm{G}_{i}, \mathfrak{S}(\mathrm{X})\right),
\end{aligned}
$$

i.e. the objects of G- $\operatorname{Set} \mathcal{f}$ and of $\underset{\rightarrow}{\lim } \mathrm{G}_{i}-\operatorname{Set} \mathcal{f}$ are the same.

We show that, for all $i$ and $\mathrm{G}_{i}$-sets X and $\mathrm{Y}, \operatorname{Hom}_{\mathrm{G}-\operatorname{setf}}(\mathrm{X}, \mathrm{Y})=\underset{\overrightarrow{j \geq i}}{\lim } \operatorname{Hom}_{\mathscr{G}_{\boldsymbol{r}} \operatorname{setf}}$ (X, Y), X and Y being G -sets by $f_{i}^{\infty}$ and $\mathrm{G}_{j}$-sets by $f_{i}^{j}$. Note that, for all
 $\operatorname{Hom}_{\mathscr{G}_{\mu} \operatorname{Setf}_{f}}(\mathrm{X}, \mathrm{Y})$ amounts to a union. Let us show that a G-morphism $\varphi: \mathrm{X} \rightarrow \mathrm{Y}$ is a $\mathrm{G}_{j}$-morphism for some $j$. Now, if $\varphi$ is a G-morphism, then $(\forall g \in \mathrm{G})(\forall x \in \mathrm{X})$ $\varphi(g \cdot x)=g \cdot \varphi(x)$ or $\left(\forall g \in \operatorname{Im} f_{i}^{\infty}\right) \varphi(g \cdot x)=g \cdot \varphi(x)$. Choose an index $j$ for which $\operatorname{Im} f_{i}^{\infty}=\operatorname{Im} f_{i}^{j}$ (2.8.5). Then, $\varphi$ is a $\mathrm{G}_{j}$-morphism.

### 2.9.5

Let $G$ be a group, which can be endowed with the discrete topology. A profinite completion of G is the inverse limit of the groups $\mathrm{G} / \mathrm{N}$, where N runs over all normal subgroups of finite index in G ordered by reverse inclusion.
Proposition Let G be a discrete group. Denote its profinite completion by $\widehat{\mathrm{G}}$. Then the categories G-Setf and $\widehat{\mathrm{G}}$-Setf are equivalent.

Proof Let X be a finite set. The kernel of a homomorphism from $G$ to $\mathfrak{S}(X)$ is a normal subgroup of finite index, and so $\operatorname{Hom}_{\mathscr{G}_{x}}(G, \mathfrak{S}(X))=\xrightarrow{\lim } \operatorname{Hom}(G / N, \mathfrak{S}(X))=$ $\operatorname{Hom}_{\mathscr{G}_{r} \mathscr{T}_{o p}}(\widehat{\mathrm{G}}, \mathfrak{S}(\mathrm{X}))$. If X and Y are finite $\mathrm{G} / \mathrm{N}$-sets, then

$$
\operatorname{Hom}_{G}(X, Y)=\operatorname{Hom}_{G / N}(X, Y)=\underset{\underset{N^{\prime} \subset N}{l}}{\underset{\mathrm{lim}}{\mathrm{~N}}} \operatorname{Hom}_{G / \mathrm{N}^{\prime}}(\mathrm{X}, \mathrm{Y})
$$

Therefore, $\mathrm{G}-\operatorname{Set} \mathcal{P}=\underset{\longrightarrow}{\lim }(\mathrm{G} / \mathrm{N}-\operatorname{Set} \mathcal{f})=\widehat{\mathrm{G}}-\operatorname{Set} f$.

### 2.9.6

Let $G$ be a profinite group and $\left(X_{i}\right)$ an inverse system of finite G-sets. Then $X=\lim _{i} X_{i}$ is a profinite space on which $G$ acts continuously. Indeed, $G \times X=$ $\lim _{\leftarrow}\left(\overleftarrow{G} \times X_{i}\right)$, and the action $\mathrm{G} \times \mathrm{X} \rightarrow \mathrm{X}$ is induced by the actions $\mathrm{G} \times \mathrm{X}_{i} \rightarrow \mathrm{X}_{i}$.

Let G- $\mathscr{P r o f}$ denote the category of profinite spaces on which $G$ acts continuously.

Proposition Let G be a profinite group. The category G-Prof is equivalent to G-Setf.

More precisely, the functor G-Setf $\rightarrow$ G-PPaf which assigns to an inverse system its inverse limit is an equivalence.

Proof (a) It is essentially surjective. Let X be a profinite space on which G acts continuously. We know (2.8.3, pf. (i) $\Rightarrow$ (ii)) that X is the inverse limit of the inverse
system $\left(X_{\lambda}\right)_{\lambda \in \Lambda}$, where $\Lambda$ is the set of equivalence relations on $X$ with open classes and $X_{\lambda}$ the quotient $X / \lambda$. In general, $X_{\lambda}$ is not a G-set. We show that the subset $\Lambda^{\prime}$ of $\Lambda$ consisting of equivalence relations compatible with the action of G on X is cofinal in $\Lambda$, i.e.

$$
(\forall \lambda \in \Lambda)(\exists \mu \geq \lambda) \quad x \sim_{\mu} y \Rightarrow(\forall g \in \mathrm{G}) \quad g x \sim_{\mu} g y .
$$

Let $\lambda \in \Lambda$, and define $\mu$ by $x \sim_{\mu} y \Longleftrightarrow(\forall g \in \mathrm{G}) g x \sim_{\lambda} g y$.
The relation $\mu$ is finer than $\lambda$ and is compatible with the action of G. We show that its classes are open. It suffices to show that the graph $\Gamma_{\mu}$ of $\mu$ is open in $\mathrm{X} \times \mathrm{X}$. Consider the maps $\mathrm{X} \times \mathrm{X} \times \mathrm{G} \underset{q}{\stackrel{p}{\longrightarrow}} \mathrm{X} \times \mathrm{X}$ defined by $p(x, y, g)=(x, y)$ and $q(x, y, g)=(g x, g y)$.

Then $(\mathrm{X} \times \mathrm{X})-\Gamma_{\mu}=p\left(q^{-1}\left((\mathrm{X} \times \mathrm{X})-\Gamma_{\lambda}\right)\right)$, where $\hat{u} \Gamma_{\lambda}$ is the graph of $\lambda$. Since the classes of $\lambda$ are open, $\Gamma_{\lambda}$ is open; since $X$ and $G$ are compact, and $p$ and $q$ continuous, $(\mathrm{X} \times \mathrm{X})-\Gamma_{\mu}$ is closed. So $\Gamma_{\mu}$ is open.

Hence $\mathrm{X}=\lim _{\grave{\lambda} \in \Lambda} \mathrm{X}_{\lambda}=\lim _{\check{\overleftarrow{\prime}} \mathrm{\Lambda}^{\prime}} \mathrm{X}_{\mu}$, where all $\mathrm{X}_{\mu}$ are finite G -sets and the identification is compatible with the actions of G.
(b) It is fully faithful. Let $\left(\mathrm{X}_{i}\right)$ be an inverse system of finite G-sets, and Y a finite G-set. We show that the natural map $\alpha: \lim _{\rightarrow} \operatorname{Hom}_{\mathrm{G}-\operatorname{setf}}\left(\mathrm{X}_{i}, \mathrm{Y}\right) \rightarrow$ $\operatorname{Hom}_{\mathrm{G}-\mathscr{P r a f}}\left(\lim _{\mathrm{X}}, \mathrm{Y}\right)$ is a bijection.

In the following commutative diagram:

the arrow on the right is a bijection by 2.8.6, and so $\alpha$ is injective.
We show that it is surjective. Let $\varphi: \lim \mathrm{X}_{i} \rightarrow \mathrm{Y}$ be a morphism. There is an index $i$ and a map $\varphi_{i}: \mathrm{X}_{i} \rightarrow \mathrm{Y}$ making the diagram

commutative. The map $\varphi_{i}$ is compatible with the action of G on $\operatorname{Im} f_{i}^{\infty}$. Choose $j$ such that $\operatorname{Im} f_{i}^{j}=\operatorname{Im} f_{i}^{\infty}$. Then $\varphi_{i} \circ f_{i}^{j}$ is a G-morphism.

Let $\left(\mathrm{Y}_{j}\right)$ be an inverse system of finite G-sets. Then,

$$
\begin{aligned}
& \underset{j}{\lim } \underset{i}{\lim } \operatorname{Hom}_{\mathrm{G}}\left(\mathrm{X}_{i}, \mathrm{Y}_{j}\right) \approx \underset{\leftarrow}{\lim } \operatorname{Hom}_{\mathrm{G}-\mathscr{P r a f}}\left(\underset{\longleftarrow}{\lim } \mathrm{X}_{i}, \mathrm{Y}_{i}\right) \\
& \approx \operatorname{Hom}_{\mathrm{G}-\mathscr{P r a f}}\left(\lim _{\leftarrow} \mathrm{X}_{i}, \lim _{\underset{i}{ }}\right)
\end{aligned}
$$

### 2.9.7

Let G be a profinite group, $e$ its identity element, and for $g \in \mathrm{G}, \rho_{g}$ the right translation $x \mapsto x \cdot g$. Let ${ }^{\circ} \mathrm{G}$ be the profinite G -set obtained by making G act on itself by left translations. The maps $\rho_{g}$ are the automorphisms of ${ }^{\circ} \mathrm{G}$.

The pair ( $\mathrm{G}, e$ ) represents the forgetful functor G- $\mathscr{P} \mathscr{\circ} \boldsymbol{\mathcal { F }} \rightarrow \mathscr{E}$ ns: for any profinite G-set X, the map $x \mapsto \delta_{x}$, where $\delta_{x}(g)=g \cdot x$ is a bijection from the set X onto Hom( ${ }^{\circ}$; X ), and its inverse is $\phi \mapsto \phi(e)$. Then $\delta_{g \cdot x}=\delta_{x} \circ \rho_{g}$.

If X is finite, $\operatorname{Hom}(G ; \mathrm{X})$ is finite; in the general case X can be written $\underset{\leftrightarrows}{\lim } \mathrm{X}_{i}$ with $\mathrm{X}_{i}$ finite. Then $\operatorname{Hom}\left({ }^{\circ} \mathrm{G} ; \mathrm{X}\right)=\underset{\leftarrow}{\lim } \operatorname{Hom}\left({ }^{\circ} \mathrm{G} ; \mathrm{X}_{i}\right)$ is the inverse limit of finite sets and can be considered a profinite set.

The map $(g, \phi) \mapsto \phi \circ \rho_{g}$ defines a left-action of G on Hom('G; X) which can be identified with the action of G on X defining the G -set structure of X .

### 2.9.8

Let G and $\mathrm{G}^{\prime}$ be two profinite groups and $\phi: \mathrm{G}^{\prime} \rightarrow \mathrm{G}$ a morphism (i.e. a continuous homomorphism). Define a functor $\phi^{*}: \mathrm{G}-\mathscr{P r} \boldsymbol{\not r} \rightarrow \mathrm{G}^{\prime}-\mathscr{P r} \not \circ \mathcal{\not}$ as follows: if X is a profinite G-set, $\phi^{*} \mathrm{X}$ is the space X on which G-acts by $\left(g^{\prime}, x\right) \mapsto \phi\left(g^{\prime}\right) \cdot x$; if $f: \mathrm{X} \rightarrow \mathrm{Y}$ is a morphism of G-sets, then $\phi^{*}(f)=f$. This functor induces a functor G- $\mathscr{E} n s \rightarrow \mathrm{G}^{\prime}-\mathscr{E} n s$ also written $\phi^{*}$.

If $\phi$ is an isomorphism, the functor $\phi^{*}$ is an equivalence of categories with $\left(\phi^{-1}\right)^{*}$ as inverse.

Conversely:
Proposition Let G and $\mathrm{G}^{\prime}$ be two profinite groups. If the categories $\mathrm{G}-\mathscr{P r a f}$ and $\mathrm{G}^{\prime}-\mathscr{P} r a f$ are equivalent, the profinite groups G and $\mathrm{G}^{\prime}$ are isomorphic.

Every equivalence of categories $\Phi: \mathrm{G}-\mathscr{P r o f} \rightarrow \mathrm{G}^{\prime}$-Prof is isomorphic to a functor $\phi^{*}$, where $\phi: \mathrm{G}^{\prime} \rightarrow \mathrm{G}$ is an isomorphism.

Proof Let $\Phi: \mathrm{G}-\mathscr{P} r \nsim \mathcal{F} \rightarrow \mathrm{G}^{\prime}-\mathscr{P r} \nsim \mathcal{\beta}$ be an equivalence of categories. The object $\mathrm{O}=\varnothing$ of $\mathrm{G}-\mathscr{P} r a f$ is characterized by $\operatorname{Hom}(\mathrm{X} ; \mathrm{O})=\varnothing$ if $\mathrm{X} \not \approx \mathrm{O}$. Likewise for $\mathrm{G}^{\prime}$. Hence $\Phi(\varnothing)=\varnothing$.

The group G acts transitively on an object X of G- $\mathscr{P} \mathscr{O} \not \mathcal{\not}$ if and only if there is no object $\mathrm{Y} \sqcup \mathrm{Z}$ with Y and Z non-empty isomorphic to X . Therefore, if G acts transitively on X , the group $\mathrm{G}^{\prime}$ acts transitively on $\Phi(\mathrm{X})$.

If X is an object of $\mathrm{G}-\mathscr{P} r \mathscr{O}$ on which G acts transitively, then $\mathrm{X} \approx \mathrm{G}$ if and only if $\mathrm{X} \neq \varnothing$ and $\operatorname{Hom}(\mathrm{X}, \mathrm{Y}) \neq \varnothing$ for all $\mathrm{Y} \neq \varnothing$. Likewise for $\mathrm{G}^{\prime}$. Therefore $\Phi\left({ }^{\circ} \mathrm{G}\right) \approx{ }^{\circ} \mathrm{G}^{\prime}$.

Let $\alpha:{ }^{\circ} \mathrm{G}^{\prime} \rightarrow \Phi\left({ }^{\circ} \mathrm{G}\right)$ be an isomorphism in G- $\mathscr{P r} \circ \boldsymbol{\mathcal { F }}$. Then $\alpha_{!}: f \mapsto \alpha \circ f \circ$ $\alpha^{-1}$ is a group isomorphism from $\operatorname{Aut}\left({ }^{\circ} \mathrm{G}^{\prime}\right)$ onto $\operatorname{Aut}\left(\Phi\left({ }^{\circ} \mathrm{G}\right)\right)$. The map $\Phi \cdot{ }_{\mathrm{G}}$ : $\operatorname{Aut}\left({ }^{\circ} \mathrm{G}\right) \rightarrow \operatorname{Aut}\left(\Phi\left({ }^{\circ} \mathrm{G}\right)\right)$ is also an isomorphism, and so is $\Phi \cdot_{\mathrm{G}}^{-1} \circ \alpha_{!}: \operatorname{Aut}\left({ }^{\circ} \mathrm{G}^{\prime}\right) \rightarrow$
$\operatorname{Aut}\left({ }^{\circ} \mathrm{G}\right)$. The latter is of the form $\rho_{g^{\prime}} \mapsto \rho_{\phi\left(g^{\prime}\right)}$, where $\phi: \mathrm{G}^{\prime} \rightarrow \mathrm{G}$ is an isomorphism of profinite groups.

If X is a profinite G -set, define a bijection $\beta=\beta_{\mathrm{X}}: \mathrm{X} \rightarrow \Phi(\mathrm{X})$ by the commutativity of the diagram:


This bijection is a homeomorphism, as can be seen by writing $X$ as an inverse limit of finite G -sets. We show that $\beta_{\mathrm{X}}$ is a $\phi^{-1}$-morphism:

Let $g \in \mathrm{G}$ and $g^{\prime} \in \mathrm{G}^{\prime}$ be such that $g=\phi\left(g^{\prime}\right)$, and $x \in \mathrm{X}$. We need to show that $\beta(g \cdot x)=g^{\prime} \cdot \beta(x)$, i.e. that $\delta_{\beta(g \cdot x)}=\delta_{g^{\prime} \cdot \beta(x)}=\delta_{\beta(x)} \circ \rho_{g^{\prime}}$. By definition of $\beta_{\mathrm{X}}$, $\delta_{\beta(x)}=\Phi\left(\delta_{x}\right) \circ \alpha$. Hence it amounts to showing that $\Phi\left(\delta_{g \cdot x}\right) \circ \alpha=\Phi\left(\delta_{x}\right) \circ \alpha \circ \rho_{g^{\prime}}$, i.e. that $\Phi\left(\delta_{g \cdot x}\right)=\Phi\left(\delta_{x}\right) \circ \alpha_{!}\left(\rho_{g^{\prime}}\right)$. However, $\alpha_{!}\left(\rho_{g^{\prime}}\right)=\Phi\left(\rho_{g}\right)$ by definition of $\phi$ since $g=\phi\left(g^{\prime}\right)$, and $\Phi$ being a functor, $\Phi\left(\delta_{g \cdot x}\right)=\Phi\left(\delta_{x} \circ \rho_{g}\right)=\Phi\left(\delta_{x}\right) \circ \Phi\left(\rho_{g}\right)$.

Indeed, $\beta(g \cdot x)=\beta\left(\phi\left(g^{\prime}\right) \cdot x\right)=g^{\prime} \cdot \beta(x)$. Hence $\beta=\beta_{\mathrm{X}}$ can be considered an isomorphism in $\mathrm{G}^{\prime}-\mathscr{P} \mathscr{\sigma} \boldsymbol{\circ}$ from $\phi^{*}(\mathrm{X})$ onto $\Phi(\mathrm{X})$. The family $\left(\beta_{\mathrm{X}}\right)$ where X runs through the objects of G-Prof is a functorial isomorphism from $\phi^{*}$ onto $\Phi$.

Corollary Let G and $\mathrm{G}^{\prime}$ be two profinite groups. If the categories G -Setf and $\mathrm{G}^{\prime}$-Setf are equivalent, then the profinite groups G and $\mathrm{G}^{\prime}$ are isomorphic.

Every equivalence of categories $\Phi: \mathrm{G}-\operatorname{Set} \boldsymbol{f} \rightarrow \mathrm{G}^{\prime}$-Setf is isomorphic to a functor $\phi^{*}$, where $\phi: \mathrm{G}^{\prime} \rightarrow \mathrm{G}$ is an isomorphism.

Indeed every equivalence G - $\boldsymbol{\operatorname { e t }} \boldsymbol{f} \rightarrow \mathrm{G}^{\prime}$ - $\operatorname{Set} \boldsymbol{f}$ can be extended to an equivalence from G- $\mathscr{P r o f}$ onto $\mathrm{G}^{\prime}-\mathscr{P}^{\prime} \not \circ \mathscr{\not}$ by passing to the inverse limit.

## Exercises 2.9. (Profinite Groups)

1.-Let G be a compact group and H an open subgroup of G . Show that H contains an open normal subgroup.
2.-Show that the category of profinite Abelian groups is anti-equivalent to the category of torsion Abelian groups (i.e. where every element has finite order) (use 2.3, Exercise 4).
3.-Let $\Sigma$ be the set of prime numbers. $\mathbb{N}^{*}$ can be identified with $\Sigma^{(\mathbb{N})}$. Set $\widetilde{\mathbb{N}}^{*}=\Sigma^{\widetilde{\mathbb{N}}}$ o $\dot{u} \widetilde{\mathbb{N}}=\mathbb{N} \cup\{\infty\}$. The elements of $\widetilde{\mathbb{N}}^{*}$ are called supernatural numbers. They are formal products $n=\prod_{p \in \Sigma} p^{r_{p}}$, where $r_{p} \in \overline{\mathbb{N}}$. The product $n_{1} n_{2}$ is defined in the obvious manner. The support of $n$ is the set of $p$ such that $r_{p} \neq 0 . n_{1}$ and $n_{2}$ are said to be relatively prime if they have disjoint supports. Let G be a profinite group. Define the order \#G of $G$ to be the supernatural number sup Card $(G / N)$, where $N$ is subgroup of G of finite index. If H is a closed subgroup of G , define the index $(G: H)$ of $H$ to be the supernatural number $\sup C \operatorname{ard}\left(G / H^{\prime}\right)$ where $H^{\prime}$ is of finite index and contains H .
(a) Show that $\# \mathrm{G}=\# \mathrm{H} \cdot(\mathrm{G}: \mathrm{H})$.

Deduce that there is no non-zero homomorphism from $G$ to $G_{1}$ when the orders of $G$ and $G_{1}$ are relatively prime.
(b) If $\mathrm{H}_{1}$ and $\mathrm{H}_{2}$ are two subgroups of G , let $\left(\mathrm{H}_{1}, \mathrm{H}_{2}\right)$ denote the closed subgroup generated by all $a b a^{-1} b^{-1}$ with $a \in \mathrm{H}_{1}$ and $b \in \mathrm{H}_{2}$. Define the lower central series $\left(\mathrm{C}_{i}\right)$ by $\mathrm{C}_{1}=\mathrm{G}, \mathrm{C}_{n+1}=\left(\mathrm{G}, \mathrm{C}_{n}\right)$. For $p \in \Sigma, \mathrm{G}$ is said to be a $p$-group if \#G is a (finite or infinite) power of $p$.

Show that if G is a profinite $p$-group, then $\cap \mathrm{C}_{i}=\{e\}$. (Assume this result holds for finite groups.)
(c) Let G be a profinite group, $p \in \Sigma$. A Sylow p-subgroup of G is a closed subgroup which is also a $p$-group whose index is relatively prime to $p$. Extend the following two Sylow theorems to profinite groups:

Theorem 1 Every finite group has a Sylow p-subgroup.
Theorem 2 All Sylow p-subgroups are conjugate.
(The reader may write G as $\lim \mathrm{G}_{i}$, where all $\mathrm{G}_{j} \rightarrow \mathrm{G}_{i}$ are surjective, and that for each $i$, and consider the set $\overleftarrow{\mathrm{X}_{i}}$ of Sylow $p$-subgroups of $\mathrm{G}_{i}$.)
4.-Let $G$ be a compact group, $S_{\alpha}$ a descending filtration of closed subgroups of $G$. Set $S=\cap S_{\alpha}$. Show that $G / S=\lim _{\leftarrow} G / S_{\alpha}$.
5.-Let $G$ be a profinite group, $F$ a closed subgroup of $G$. The aim is to show the existence of a continuous section from $G / F$ to $G$ (which may not be a group homomorphism even when F is normal).
(a) Show that, if S is a closed subgroup of F of finite index in F , then there is a continuous section G/F $\rightarrow$ G/S.
(b) Let $\mathscr{E}$ be the set of pairs $(\mathrm{S}, \sigma)$, where S is a closed subgroup of F and $\sigma: \mathrm{G} / \mathrm{F} \rightarrow \mathrm{G} / \mathrm{S}$ is a continuous section. Define an order with respect to which $\mathscr{E}$ is direct.
(c) Using (a), show that $S=\{e\}$ for any maximal element of $\mathscr{E}$. Conclude. Compare with 2.8, Exercise 3.

## Reference

1. N. Bourbaki, Algèèbre, ch. 1 à 3, Hermann, 1970; èAlgèbre, ch. 10, algèbre homologique, Masson, 1980

## Chapter 3 <br> Linear Algebra

Unless stated otherwise, all rings considered in this chapter are assumed to be commutative, to have an identity element, and all homomorphisms to be unital (i.e. to send the identity onto the identity).

### 3.1 Integral and Principal Ideal Domains, Reduced Rings

### 3.1.1 Integral Domains

Let A be a ring. An element $a \in \mathrm{~A}$ is said to be a zero divisor in A if the map $x \mapsto a x$ from A to A is not injective.

The ring A is an integral domain if $\mathrm{A} \neq 0$ and 0 is its only zero divisor.

### 3.1.2 Examples of Integral Domains

(1) All subrings of an integral domain.
(2) The ring $\mathbb{Z}$ of all integers.
(3) The ring $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ of polynomials in $n$ variables, and the ring $\mathrm{K}\left[\left[\mathrm{X}_{1}, \ldots\right.\right.$, $\mathrm{X}_{n}$ ]] of formal series over an integral domain $\mathrm{K} .{ }^{1}$
(4) The ring $\mathbb{C}\left\{\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right\}$ of convergent power series in $n$ variables.
(5) The ring $\mathcal{O}(\mathrm{U})$ of all holomorphic functions on any open and connected subset U of $\mathbb{C} .{ }^{2}$

[^2]
## Counterexamples

The following rings are not integral domains.
(1) The ring $C(X, \mathbb{R})$ of continuous functions on a non-zero metric space X having more than one point.
(2) The ring $\mathscr{E}(\mathrm{I})$ of $\mathrm{C}^{\infty}$-functions on a non-empty interval I of $\mathbb{R}$ with more than one point (3.1, Exercise 2).
(3) The product ring of two non-zero rings $A$ and $B$. Indeed $(1,0) \cdot(0,1)=(0,0)$.

### 3.1.3 Field of Fractions

Let A be an integral domain. Set $\mathrm{A}^{*}=\mathrm{A}-\{0\}$. For $(a, b)$ and $(c, d) \in \mathrm{A} \times \mathrm{A}^{*}$, write $(a, b) \sim(c, d)$ if $a d=b c$. This defines an equivalence relation on $\mathrm{A} \times \mathrm{A}^{*}$; let K denote the quotient of $\mathrm{A} \times \mathrm{A}^{*}$ by this relation, and $a / b$ be the class of $(a, b)$. For any two elements $x$ and $y$ of K, there are representatives $(a, b)$ and $(c, d)$ of $x$ and $y$ such that $b=d$ (reduction to the same denominator). Addition in K is defined by $x+y=(a+c) / b$, which is an element only dependent on $x$ and $y$. Multiplication in $\mathrm{A} \times \mathrm{A}^{*}$ is compatible with the equivalence relation and defines a multiplication on K . Together with these two operations, K is a field called the field of fractions of A . The map $a \mapsto a / 1$ from A to K is an injective homomorphism embedding A in K .

Examples (1) $\mathbb{Q}$ is the field of fractions of $\mathbb{Z}$.
(2) If K is a field, the field of fractions of the polynomial ring $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ is the field of rational fractions $\mathrm{K}\left(\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right)$.
(3) If $U$ is open and connected in $\mathbb{C}$, the field of fractions of the ring $\mathcal{O}(U)$ of holomorphic functions on U is the field $\mathscr{M}(\mathrm{U})$ of meromorphic functions on $\mathrm{U}(3.1$, Exercise 17).

### 3.1.4 Prime Ideals and Maximal Ideals

Let A be a ring and I an ideal of A . I is said to be prime if the quotient $\mathrm{A} / \mathrm{I}$ is integral. For the sake of simplicity, I is called maximal if I is maximal among all strict ideals of $A$, i.e. if $I \neq A$ and there is no ideal $J$ such that $I \supsetneqq J \supsetneqq A$.

Proposition Let A be a ring and I an ideal of A . Then $\mathrm{A} / \mathrm{I}$ is a field if and only if I is maximal.

Proof The ring $\mathrm{B}=\mathrm{A} / \mathrm{I}$ is a field if and only if $\mathrm{B} \neq 0$ and its only ideals are 0 and B. As there is a bijection between the ideals of A containing I, and the ideals of A/I, the result follows.

Corollary Every maximal ideal is prime.

Examples (1) The prime ideals in the ring $\mathbb{Z}$ are 0 and $\mathbb{Z} p$, where $p$ is a prime number. The non-zero prime ideals are all maximal. This will be seen to always be the case in a principal ideal domain.
(2) Let A and B be two rings and $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ a homomorphism. If B is an integral domain, then the ideal $\operatorname{Ker} \varphi=\varphi^{-1}(0)$ in A is prime.
(3) Let K be a field. In the ring of polynomials $\mathrm{K}[\mathrm{X}, \mathrm{Y}]$, the ideal $(\mathrm{Y})$ generated by Y is prime. Indeed, it is the kernel of the homomorphism $\rho: \mathrm{K}[\mathrm{X}, \mathrm{Y}] \rightarrow \mathrm{K}[\mathrm{X}]$ defined by $\rho(\mathrm{X})=\mathrm{X}, \rho(\mathrm{Y})=0$ (see 3.7.2); in terms of polynomial functions, $\rho(f)$ is the function $x \mapsto f(x, 0)$. This ideal is not maximal: it is strictly contained in the ideal (X, Y).
(4) Let X be a metric space and A the ring $\mathrm{C}(\mathrm{X}, \mathbb{R})$ or $\mathrm{C}(\mathrm{X}, \mathbb{C})$. For $x \in \mathrm{X}$, the ideal $\mathfrak{m}_{x}$ of functions vanishing at $x$ is maximal: it is the kernel of the homomorphism $\delta_{x}$ from A onto $\mathbb{R}$ or $\mathbb{C}$, where $\delta_{x}(f)=f(x)$. If X is not discrete, then A contains non-maximal ideals (3.1, Exercise 18).

### 3.1.5

Proposition (Krull's Theorem) A nonzero ring A contains at least one maximal ideal.

Proof The set E of strict ideals of A is direct. Indeed, if $\left(\mathrm{I}_{\lambda}\right)$ is a totally ordered family of strict ideals, then $I=\bigcup I_{\lambda}$ is a strict ideal of A since $(\forall \lambda) 1 \notin I_{\lambda}$, and so $1 \notin \mathrm{I}$. By Zorn's theorem (1.4.3), E has a maximal element,

Corollary Let A be a ring. Every strict ideal of A is contained in a maximal ideal.
Indeed, the ideals of A containing I correspond to the ideals of $\mathrm{A} / \mathrm{I}$.

### 3.1.6 Nilpotent Elements, Nilradical

Let A be a ring. An element $x$ of A is said to be nilpotent if there exists $n \in \mathbb{N}$ such that $x^{n}=0$.

Examples (1) In an integral domain, only 0 is nilpotent.
(2) Let X be a topological space. In $\mathrm{C}(\mathrm{X}, \mathbb{R})$, only 0 is nilpotent.
(3) In the ring A of $\mathrm{C}^{\infty}$-functions $f: \mathbb{R} \rightarrow \mathbb{R}$, the functions $f$ such that $f(0)=$ $f^{\prime}(0)=0$ form and ideal I . The ring $\mathrm{B}=\mathrm{A} / \mathrm{I}$ is isomorphic to $\mathbb{R}[\mathrm{X}] /\left(\mathrm{X}^{2}\right)$. In this ring, X is nilpotent.
(4) In the ring $\mathbb{Z} /(1000)$, the class of 150 is nilpotent. More generally, in $\mathbb{Z} /(n)$, the class of $m$ is a nilpotent element if and only if every prime divisor of $n$ divides $m$.

Theorem and Definition Let A be a ring. The set $\mathfrak{N}$ of nilpotent elements of A is the intersection of prime ideals of A and is called the nilradical of A .

Proof If $x \in \mathrm{~A}$ is a nilpotent element then, for every prime ideal $\mathfrak{p}, \chi(x)=0$, where $\chi: \mathrm{A} \rightarrow \mathrm{A} / \mathfrak{p}$ is the quotient map, since $\chi(x)$ is nilpotent in the integral domain $\mathrm{A} / \mathfrak{p}$. So $x \in \mathfrak{N}$.

Conversely, let $x \in \mathrm{~A}$ be a non-nilpotent element, and $\mathrm{S}=\left\{x^{n}\right\}_{n \in \mathbb{N}}$. The set E of ideals I de A such that $\mathrm{I} \cap \mathrm{S}=\varnothing$ is directed, and so, by Zorn's lemma, has a maximal element $\mathfrak{p}$. We show that $\mathfrak{p}$ is prime. Let $\mathrm{A}^{\prime}=\mathrm{A} / \mathfrak{p}$ and $x^{\prime}$ be the image of $x$ in $\mathrm{A}^{\prime}$. The element $x^{\prime}$ is not nilpotent in $\mathrm{A}^{\prime}$ since $\mathfrak{p} \cap \mathrm{S}=\varnothing$. For any nonzero $y \in \mathrm{~A}^{\prime}$, there exists $n \in \mathbb{N}$ such that $x^{\prime n} \in \mathrm{~A}^{\prime} y$; otherwise, the inverse image of $\mathrm{A}^{\prime} y$ in A would be an ideal $\mathbf{J}$ strictly containing $\mathfrak{p}$ and such that $\mathbf{J} \cap \mathrm{S}=\varnothing$. If $y, z \in \mathrm{~A}^{\prime}$ are nonzero, then there exists $n$ and $p$ such that $x^{\prime n} \in \mathrm{~A}^{\prime} y$ and $x^{\prime p} \in \mathrm{~A}^{\prime} z$; so $x^{\prime n+p} \in \mathrm{~A}^{\prime} y z$, and thus $y z \neq 0$. Therefore $\mathfrak{p}$ is prime. As $x \notin \mathfrak{p}, x \notin \mathfrak{N}$.

### 3.1.7 Reduced Rings

Proposition and Definition Let A be a ring. The following conditions are equivalent:
(i) A does not have any nonzero nilpotent element;
(ii) A is isomorphic to a subring of a product of fields.

If these conditions hold, then A is said to be reduced.
Proof Let $\Delta$ be the homomorphism $x \mapsto\left(\chi_{\mathfrak{p}}(x)\right)$ from A to $\prod_{\mathfrak{p} \in S} K_{\mathfrak{p}}$, where S is the set of prime ideals of A and where, for $\mathfrak{p} \in S, K_{\mathfrak{p}}$ is the field of fractions of the integral domain $\mathrm{A} / \mathfrak{p}$ and $\chi_{\mathfrak{p}}: \mathrm{A} \rightarrow \mathrm{A} / \mathfrak{p}$ is the quotient map. The kernel of $\Delta$ is the nilradical of A. If (i) holds, $\Delta$ is injective, implying (ii). The converse is immediate.

### 3.1.8

Let $A$ be a ring. The spectrum of $A$ is the set $S$ of all prime ideals of $A$. For all $\mathfrak{p} \in S$, let $K_{\mathfrak{p}}$ denote the field of fractions of the ring $A / \mathfrak{p}$. If $x \in A$, define the function $\hat{x}$ on S by $\hat{x}(\mathfrak{p})=\chi_{\mathfrak{p}}(x) \in \mathrm{K}_{\mathfrak{p}}$ (its values are in varying fields). Then $\widehat{x+y}=\hat{x}+\hat{y}$, $\widehat{x y}=\hat{x} \hat{y}$.

If A is reduced, the map $x \mapsto \hat{x}$ is injective. Assume A is integral and let K be a field of fractions of A. For any $x=\frac{a}{b} \in \mathrm{~K}$, let $\hat{x}$ be the function defined by $\hat{x}=\frac{\hat{a}}{\hat{b}}$ on the elements $\mathfrak{q}$ of $S$ for which $\hat{b}(\mathfrak{q}) \neq 0$. It is said to have a zero of order $\geqslant k$ in $\mathfrak{p}$ if $x$ of the form $\frac{a}{b}$ with $a \in \mathfrak{p}^{k}$ and $b \notin \mathfrak{p}$. The functions $\hat{x}$ and $\hat{y}$ are said to have the same polar part in $\mathfrak{p}$ if $\widehat{x-y}$ is defined in $\mathfrak{p}$.

Considering functions $\hat{x}$ and using the above language enable us to draw on methods of analysis to solve algebra problems (3.1, Exercises 6 and 16); compare also the proofs of the Chinese theorem (5.1.5) and of Lagrange's interpolation formula. ${ }^{3}$

Let $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism. For every prime ideal $\mathfrak{q}$ of $\mathrm{B}, \varphi^{-1}(\mathfrak{q})$ is a prime ideal of A . This gives a map $\varphi^{*}$ from the spectrum of B to that of A , and thus defines a contravariant functor from $\mathscr{A} n n$ to $\mathscr{E} n s$. This functor factorizes naturally through the functor Spec from $\mathscr{A} n n$ to the category $\mathscr{T} 0 \boldsymbol{R}$ of (not necessarily Hausdorff) topological spaces (see 3.1, Exercise 7).

Comments As in 2.7.6, thanks to the functor Spec : $\mathscr{A} \boldsymbol{n} \boldsymbol{n} \rightarrow \mathscr{T} a p$, the opposite category of rings and $\mathscr{T} a p$ bear some resemblance.

### 3.1.9 Principal Ideal Domains

An ideal I in a ring A is said to be principal if $\mathrm{I}=\mathrm{A} x$ for some $x \in \mathrm{~A}$. An integral domain in which every ideal is principal is a principal ideal domain (PID).

As will be seen, in a principal ideal domain, every nonzero prime ideal is maximal (3.2.10, Lemma).

### 3.1.10 Euclidean Rings

As integral domain A is said to be Euclidian if there is a map $w$ from $\mathrm{A}^{*}=\mathrm{A}-\{0\}$ to $\mathbb{N}$ such that
$\left(\mathrm{S}_{1}\right)\left(\forall x, y \in \mathrm{~A}^{*}\right) w(x y) \geqslant w(y) ;$
$\left(\mathrm{S}_{2}\right)(\forall a \in \mathrm{~A})\left(\forall b \in \mathrm{~A}^{*}\right)(\exists q, r \in \mathrm{~A}) a=b q+r$ and $(r=0$ or $w(r)<w(b))$.
Proposition Every Euclidean ring is principal.
Proof Let A be a Euclidean ring, $w: \mathrm{A}^{*} \rightarrow \mathbb{N}$ satisfying $\left(\mathrm{S}_{1}\right)$ and $\left(\mathrm{S}_{2}\right)$ and I a nonzero ideal of A. Choose an element $b \neq 0$ in I such that $w(b)$ is minimum. For all $a \in \mathrm{I}$, there exists $q$ and $r$ in A satisfying the conditions of $\left(\mathrm{S}_{2}\right)$; then $r=a-b q \in \mathrm{I}$, and soù $r=0$ since $w(b)$ is minimum. Hence I is generated by the $b$.

### 3.1.11 Examples

(1) $\mathbb{Z}$ is a Euclidian ring: take $w(x)=|x|$. So it is principal. In fact, every additive subgroup of $\mathbb{Z}$ is a principal ideal domain.

[^3](2) If K is a ring, the ring $\mathrm{K}[\mathrm{X}]$ is Euclidean (for $w(f)$ take the degree of $f$ ).

For other examples see 3.1, Exercises 13 and 14. For an example of a nonEuclidean principal ideal domain, see 3.7, Exercise 9.

## Exercises 3.1. (Integral and principal ideal domains, reduced rings)

1.-(a) Let A be a ring. Show that A has idempotents $u \neq 0,1$ (i.e. $u^{2}=u$ ) if and only if $\mathrm{A}=\mathrm{B} \times \mathrm{C}$, for some nonzero rings B and C .
(b) Let X be a non-empty topological space. Show that $\mathrm{C}(\mathrm{X}, \mathbb{R})$ has idempotents different from 0 and 1 if and only if X is not connected.
2.-(a) Show that the function $f: \mathbb{R} \rightarrow \mathbb{R}$ defined by $f(x)=e^{-1 / x}$ for $x>0$ and $f(x)=0$ for $x \leqslant 0$ is infinitely differentiable. Show that the ring $\mathrm{C}^{\infty}(\mathbb{R}, \mathbb{R})$ of infinitely differentiable functions on $\mathbb{R}$ is not an integral domain.
(b) Show that the convolution ring of distributions with compact support on $\mathbb{R}$ is an integral domain. (Use the fact that the Fourier transform of a distribution with compact support is analytic, or the existence of the Laplace transform.)
3.-Let $\varphi$ be an analytic map from $\mathbb{R}$ to $\mathbb{R}^{2}, \mathrm{~V}=\varphi(\mathbb{R})$, and I be the ideal of $\mathbb{R}[\mathrm{X}, \mathrm{Y}]$ consisting of polynomials $f$ such that $f(x, y)=0$ for all $(x, y) \in \mathrm{V}$.
(a) Show that $I$ is prime.
(b) Give a generating family for I in the following cases:

$$
\begin{aligned}
\varphi(t)= & (\cos t, \sin t) ; \quad(\operatorname{ch} t, \operatorname{sh} t) ; \quad(\cos 2 t, \cos 3 t) \\
& (\cos t, \sin 2 t) ; \quad(\cos 2 t, \sin t) ; \quad\left(t, e^{t}\right) \\
& \left(\frac{2 t}{1+t^{2}}, \frac{1-t^{2}}{1+t^{2}}\right) ; \quad\left(t^{2}-1, t^{3}-t^{2}\right)
\end{aligned}
$$

(c) In each case describe the set W of $(x, y)$ such that $f(x, y)=0$ for all $f \in \mathrm{I}$. When does $\mathrm{W}=\mathrm{V}$ ?
(d) Give an example showing that, if $\varphi$ is only assumed to be $\mathrm{C}^{\infty}$, the ideal I need not be prime.
4.-Show that a direct (resp. inverse) limit of integral domains is an integral domain. Can the field of fractions of the limit be identified with the limit of the field of fractions? Investigate the case where the homomorphisms of the system are injective.
5.-Let $A$ be the quotient ring $\mathbb{R}[X, Y] / I$, where $I=\left(\mathrm{X}^{2}+\mathrm{X}^{3}-\mathrm{Y}^{2}\right)$.
(a) Construct the curve $\Gamma=\left\{(x, y) \in \mathbb{R}^{2} \mid x^{2}+x^{3}-y^{2}=0\right\}$.
(b) Let $f \in \mathbb{R}[\mathrm{X}, \mathrm{Y}]$. Show that $f \in \mathrm{I}$ if and only if $\left.f\right|_{\Gamma}=0$. Deduce that A can be identified with the ring of functions on $\Gamma$ induced by the polynomials.
(c) Show that A is an integral domain (see Exercise 3).
(d) Let $\mathfrak{m}$ be the ideal of functions of A such that $f(0)=0$. The sequence $\left(\mathfrak{m}^{k}\right)$, $k \in \mathbb{N}$ gives a fundamental system of neighbourhoods 0 with respect to the topology on A called the $\mathfrak{m}$-adic topology. Show that the completion $\widehat{\mathrm{A}}$ of A with respect to this topology is not an integral domain.
6.-(a) Prove that any rational fraction $f \in \mathbb{C}(X)$ can be uniquely written as $f(z)=\sum_{i \in \mathrm{I}} \frac{c_{i}}{\left(z-a_{i}\right)^{r_{i}}}+\mathrm{P}(z)$, where I is a finite set and P a polynomial. (Expand in the neighbourhood of each pole.)
(b) Show that the rational number $x \in \mathbb{Q}$ can be uniquely written as $\sum_{i \in \mathrm{I}} \frac{a_{i}}{p_{i}^{i}}+n$, where I is finite, $p_{i}$ primes, $r_{i} \in \mathbb{N}, 0<a_{i}<p_{i}, n \in \mathbb{Z}$.
7.-Let $A$ be a ring and $S$ its spectrum. If $I$ is an ideal of $A$, denote by $V(I)$ the set of $\mathfrak{p} \in \mathrm{S}$ such that $\hat{x}(\mathfrak{p})=0$ for all $x \in \mathrm{I}$.
(a) Show that $\mathrm{V}\left(\sum \mathrm{I}_{\lambda}\right)=\bigcap \mathrm{V}\left(\mathrm{I}_{\lambda}\right)$ if $\left(\mathrm{I}_{\lambda}\right)$ is a family of ideals of A .
(b) Show that $\mathrm{V}(\mathrm{I} \cap \mathrm{J})=\mathrm{V}(\mathrm{I} \cdot \mathrm{J})=\mathrm{V}(\mathrm{I}) \cup \mathrm{V}(\mathrm{J})$ if I and J are ideals of A .
(c) Deduce that the sets of type $\mathrm{V}(\mathrm{I})$ are the closed sets of a topology on S . This topology is called the Zariski topology. Show that, if $f: \mathrm{A} \rightarrow \mathrm{B}$ is a ring homomorphism, then the map $f^{*}$ from the spectrum of B to that of A is continuous when both spectra are equipped with the Zariski topology.
(d) Describe the Zariski topology on $\operatorname{Spec}(\mathbb{Z})$. Show that it is not Hausdorff.
(e) Show that, for any ring A, the space $\operatorname{Spec}(\mathrm{A})$ is quasi-compact.
(f) Let I be an ideal of A. Show that the set of $x \in A$ such that $\hat{x}(\mathfrak{p})=0$ for all $\mathfrak{p} \in \mathrm{V}(\mathrm{I})$ is the set of $x$ such that $(\exists r) x^{r} \in \mathrm{I}$.
8.-For every open subset $U$ of $\mathbb{C}$, let $\mathcal{O}(U)$ be the ring of holomorphic functions on U , and for every compact subset K of $\mathbb{C}$ and open neighbourhood U of $\mathrm{K}, \mathcal{O}(\mathrm{K})=$ $\xrightarrow[\mathrm{U} \supset \mathrm{K}]{\lim } \mathcal{O}(\mathrm{U})$.

Let K be a compact connected set; the aim is to show that $\mathcal{O}(\mathrm{K})$ is a principal ideal domain.
(a) Show that every nonzero $f \in \mathcal{O}(\mathrm{~K})$ can be uniquely written as the product of a polynomial all of whose roots are in K and of an invertible element.
(b) Suppose all the roots of $\mathrm{P} \in \mathbb{C}[\mathrm{Z}]$ are in K . For all $h \in \mathcal{O}(\mathrm{~K})$, show that there is a unique pair $(g, \mathrm{R}) \in \mathscr{O}(\mathrm{K}) \times \mathbb{C}[\mathrm{Z}]$ such that $h=g \mathrm{P}+\mathrm{R}$ with $\mathrm{d}^{\circ} \mathrm{R}<\mathrm{d}^{\circ} \mathrm{P}$.
(c) Deduce that $\mathcal{O}(\mathrm{K})$ is Euclidian. Conclude.
(d) What are the ideals of $\mathcal{O}(\mathrm{K})$ ?
(e) Let $f_{1}, \ldots, f_{n}$ be elements of $\mathcal{O}(\mathrm{K})$ represented by functions without common roots. Show that there exist $u_{1}, \ldots, u_{n} \in \mathcal{O}(\mathrm{~K})$ such that $\sum_{i=1}^{n} u_{i} f_{i}=1$.
9.-Let $A$ be a ring, but not a field. Prove that $A[X]$ is not a principal ideal domain. For this, the ideal of $\mathrm{A}[\mathrm{X}]$ consisting of the polynomials whose constant terms are in a non-zero ideal $\mathrm{I} \neq 0$, A in A may be shown not to be principal.
10. (Discrete valuation ring)-(a) Let $K$ be a field. Prove that the ring $K[[X]]$ of formal sequences is a principal ideal domain.
(b) More generally, let A be a ring and $w: \mathrm{A} \rightarrow \overline{\mathbb{N}}$ a map satisfying the following properties:
(i) $w(x y)=w(x)+w(y)$ for $x$ and $y$ in A;
(ii) $w(x+y) \geqslant \inf (w(x), w(y))$;
(iii) $w(x)<\infty$ for $x \neq 0$;
(iv) $w(x)>0$ for non-invertible $x$.

Prove that A is a principal ideal domain with a unique maximal ideal.
(c) Conversely, let A be a principal ideal domain with a unique maximal ideal $\mathfrak{m}$. For all $x$, let $w(x)$ be the largest integer $k$ such that $x \in \mathfrak{m}^{k}$. Show that $w$ satisfies conditions (i) to (iv) (to check (iii), let $x$ be a generator of $\mathfrak{m}$ and $y$ a generator of $\bigcap \mathfrak{m}^{k}$ and consider $w(y / x)$ ). (Rings satisfying these properties are called discrete valuation rings)
(d) Prove that the ring E of germs at 0 of $\mathrm{C}^{\infty}$-functions on $\mathbb{R}$ has a unique maximal ideal $\mathfrak{m}$, and that this ideal is principal. However $E$ is not a discrete valuation ring: $\bigcap \mathfrak{m}^{k} \neq 0$.
11.-Show that every nonzero prime ideal in a principal ideal domain is maximal.
12.-Let A be the ring of analytic functions $f: \mathbb{R} \rightarrow \mathbb{R}$ satisfying $(\forall x) f(x+2 \pi)=$ $f(x)$.
(a) Show that A is an integral domain.
(b) Show that every maximal ideal of A is of the form $\mathfrak{m}_{x}=\{f \mid f(x)=0\}$ with $x \in \mathbb{R}$.
(c) Show that every nonzero ideal is of the form $\mathfrak{m}_{x_{1}}^{k_{1}} \cdots \mathfrak{m}_{x_{r}}^{k_{r}}$. In particular every nonzero prime ideal is maximal.
(d) Let $x \in \mathbb{R}$ and $f \in \mathrm{~A}$ a function having $x$ as a simple root. Show that, for sufficiently small $\varepsilon>0, f$ has a root between $x+\varepsilon$ and $x+2 \pi-\varepsilon$. Deduce that the ideal $\mathfrak{m}_{x}$ is not principal.
(e) The ideal $\mathfrak{m}_{x_{1}}^{k_{1}} \cdots \mathfrak{m}_{x_{r}}^{k_{r}}$ is principal if and only if $k_{1}+\cdots+k_{r}$ is even.
(f) Let B be the ring of $\mathbb{R}$-analytic functions $f: \mathbb{R} \rightarrow \mathbb{C}$ such that ( $\forall x$ ) $f(x+2 \pi)=f(x)$. Show that B is principal (note that $t \mapsto e^{i t}-e^{i x}$ has a simple root at $x$ ). Is the ring B Euclidian?
13.-Show that the subrings $\mathbb{Z}[i]$ and $\mathbb{Z}\left[\frac{1+i \sqrt{3}}{2}\right]$ of $\mathbb{C}$ are principal (they are Euclidian for $z \mapsto|z|^{2}$ ).
14.-Let A be an integral domain. Show that if A is Euclidean, then there is a non-invertible element $x \in$ A such that every class of $\mathrm{A} /(x)$ contains an invertible invertible or zero (reduce to the case where $w\left(\mathrm{~A}^{*}\right)$ is an upper set of $\mathbb{N}$ and take $x$ such that $w(x)=1$, defining $w$ as in 3.1.10). Show that the ring $\mathbb{Z}\left[\frac{1+i \sqrt{d}}{2}\right]$ is not Euclidian if $d \geqslant 5$. ( $\mathbb{Z}\left[\frac{1+i \sqrt{19}}{2}\right]$ may be shown to be a principal ideal domain. For another example of a non-Euclidean principal ideal domain, see 3.7, Exercise 9.)
15.-Let A be a principal ideal domain. Show that, for any multiplicatively stable subset $S \subset A^{*}$, the set $S^{-1} A$ of fractions with numerator in $A$ and denominator in $S$ is a principal ideal domain.
16.-(a) Give a degree 2 expansion of $\sqrt{1+t}$ in the neighbourhood of $t=0$.
(b) Let A be a commutative ring in which 2 is not invertible, and $q \in$ A a multiple of $p \in \mathrm{~A}$. Show that $\left(1+\frac{q}{2}-\frac{q^{2}}{8}\right) \equiv 1+q\left(\bmod p^{3}\right)$.
(c) Find $u \in \mathbb{Z}$ such that $u^{2} \equiv 11(\bmod 125)$.
17.-The aim is to prove some algebraic properties of the ring $\mathcal{O}(\mathrm{U})$ of holomorphic functions on $U$, where $U$ is a connected open subset of $\mathbb{C}$. The reader interested in
the case $\mathrm{U}=\mathrm{D}_{\mathrm{R}}$ (or more generally, given Riemann's uniformization theorem, U simply connected) can omit part B.
A. (a) For $n, k \in \mathbb{N}$, let $\eta_{n, k}$ be the primitive of $((n+1) \ldots(n+k+1) / k!) z^{n}(1-z)^{k}$ vanishing at 0 . Show that $1-\eta_{n, k}$ has a root of order $\geqslant k$ at 1 , and that, as $n$ tends to infinity for fixed $k, \eta_{n, k}$ converges uniformly to 0 on every disc of radius $<1$ centered at 0 .
(b) Let $a \in \mathbb{C}^{*}, r<|a|, h$ a holomorphic function on $\mathrm{D}_{r^{\prime}}$ with $r^{\prime}>r, \varepsilon>0, g$ a holomorphic function in the neighbourhood of $a$ and $k \in \mathbb{N}$. Show that there is a polynomial $f$ on $\mathrm{D}_{r}$ such that $|f-h| \leqslant \varepsilon$ and that $f-g$ has a root of order $\geqslant k$ at $a$.
(c) Let $\mathrm{R} \in] 0, \infty]$, and set $\mathrm{U}=\mathrm{D}_{\mathrm{R}}$ (hence $\mathrm{U}=\mathbb{C}$ if $\mathrm{R}=\infty$ ). Let $\left(a_{n}\right)$ be a set of distinct points of U such that $\left|a_{n}\right|$ tends to R ; for all $n$, let $g_{n}$ be a holomorphic function in the neighbourhood of $a_{n}$ and $k_{n} \in \mathbb{N}$. Show that there is a holomorphic function $f \in \mathcal{O}(\mathrm{U})$ such that, for all $n$, the function $f-g_{n}$ has a root of order $\geqslant k_{n}$ at $a_{n}$.
B. (a) Let K be a compact subset of the Riemann sphere $\widehat{\mathbb{C}}=\mathbb{C} \cup\{\infty\}, a_{1}$ and $a_{2}$ two points in the same connected component of $\widehat{\mathbb{C}}-\mathrm{K}, f_{1}: \widehat{\mathbb{C}} \rightarrow \widehat{\mathbb{C}}$ a rational function such that $f_{1}^{-1}(\infty)=a_{1}$ and $\varepsilon>0$. Show that there is a rational function $f_{2}$ such that $f_{2}^{-1}(\infty)=a_{2}$ and $\left|f_{2}-f_{1}\right|<\varepsilon$ on K (first consider the case where there is a disk in $\widehat{\mathbb{C}}-\mathrm{K}$ containing $a_{1}$ and $a_{2}$, with $a_{2}$ mapped to $\infty$ by a homography).
(b) Let U be an open subset of $\mathbb{C}$. A compact subset $\mathrm{K} \subset \mathrm{U}$ is said to be complete (in U ) if $\mathrm{U}-\mathrm{K}$ does not have any relatively compact connected component in U . Prove the existence of a sequence ( $\mathrm{K}_{n}$ ) of complete compact subsets in U such that any compact subset of U is contained in some $\mathrm{K}_{n}$.
(c) Let $a=\left(a_{n}\right)$ be a sequence of distinct points of U tending to infinity in U (i.e. such that compact subsets of U only contain finitely many $a_{n}$ ); for all $n$, let $g_{n}$ be a holomorphic function in the neighbourhood of $a_{n}$ and $k_{n} \in \mathbb{N}$. Show that there is a holomorphic function $f \in \mathcal{O}(\mathrm{U})$ such that, for all $n$, the function $f-g_{n}$ has a root of order $\geqslant k_{n}$ at $a_{n}$.
C. Let $U$ be a connected open subset of $\mathbb{C}$.
(a) Let $f$ be a meromorphic function on U (see 6.1.3). Prove that $f=g / h$ for some $g, h \in \mathcal{O}(\mathrm{U})$ with $h \neq 0$. This identifies the field $\mathscr{M}(\mathrm{U})$ of meromorphic functions on U with the field of fractions of $\mathcal{O}(\mathrm{U})$.
(b) Let $a=\left(a_{n}\right)$ be a sequence of distinct points of U tending to infinity in U . Let $\mathfrak{U}$ be an ultrafilter on $\mathbb{N}$, and set

$$
\mathrm{I}(a, \mathfrak{U})=\left\{f \in \mathcal{O}(\mathrm{U}) \mid(\exists \mathrm{A} \in \mathfrak{U})(\forall n \in \mathrm{~A}) f\left(a_{n}\right)=0\right\} .
$$

Show that $\mathrm{I}(a, \mathfrak{U})$ is a maximal ideal of $\mathcal{O}(\mathrm{U})$ whose codimension as a vector subspace over $\mathbb{C}$ is not 1 (1.7, Exercise 2 ). Show that any ideal of $\mathcal{O}(\mathrm{U})$ that cannot be written as $\operatorname{Ker} \delta_{x}\left(\right.$ where $\left.\delta_{x}(f)=f(x)\right)$ is of the form $\mathrm{I}(a, \mathfrak{U})$.
(c) Let $\mathrm{J}(a, \mathfrak{U})$ be the ideal of $\mathcal{O}(\mathrm{U})$ consisting of $f$ such that, $\lim _{\mathfrak{U}} k_{n}=\infty$, where $k_{n}$ denotes the order of vanishing of $f$ at $a_{n}$. Show that $\mathrm{J}(a, \mathfrak{U})$ is a prime ideal but not a maximal one. Construct a strictly decreasing sequence of prime ideals in $\mathcal{O}(\mathrm{U})$.
18.-Let X be a metric space and A the ring $\mathrm{C}(\mathrm{X}, \mathbb{R})$ or $\mathrm{C}(\mathrm{X}, \mathbb{C})$. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence of distinct points of X converging to a point $a \in \mathrm{X}$, with $x_{n} \neq a$ for all $n$. For a filter $\mathscr{F}$ on $\mathbb{N}$ tending to $\infty$, set $\mathrm{I}_{\mathscr{F}}$ to be the set of $f \in \mathrm{~A}$ such that $\left\{n \in \mathbb{N} \mid f\left(x_{n}\right)=0\right\} \in \mathscr{F}$.

Show that $\mathrm{I}_{\mathscr{F}}$ is an ideal of A strictly contained in the maximal ideal $\mathfrak{m}_{a}$, and that it is a prime ideal if and only if $\mathscr{F}$ is an ultrafilter.

### 3.2 Unique Factorization Domains

### 3.2.1 Monoids

A monoid is a set with an associative composition law that has an identity element. The law is written either additively or multiplicatively, the additive notation being usually used for commutative monoids. The morphisms are the maps preserving the composition law and the identity element.

If $M$ is a monoid, the set $G$ of invertible elements of $M$ is a group for the induced law. If M is commutative, the relation

$$
(\exists g \in \mathrm{G}) \quad y=g x,
$$

between elements $x$ and $y$ of M is an equivalence relation compatible with the law of composition.

A monoid M is said to be regular if, for all $x, y, z \in \mathrm{M}$,

$$
(x z=y z) \Rightarrow(x=y) \text { and }(z x=z y) \Rightarrow(x=y)
$$

Let M be a monoid, $e$ its identity element. A subset $\mathrm{M}^{\prime}$ of M is said to be a submonoid of M if $\mathrm{M}^{\prime}$ is stable under the composition law and $e \in \mathrm{M}^{\prime}$. If $\mathrm{M}^{\prime}$ is stable under the law of M and has an identity element $e^{\prime} \neq e$ for the induced law, then $\mathrm{M}^{\prime}$ is not a submonoid of M .

Let M be a monoid and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of M . A submonoid M generated by the family $\left(x_{i}\right)$ is the smallest submoid $\mathrm{M}^{\prime}$ of M such that $(\forall i \in \mathrm{I})$ $x_{i} \in \mathrm{M}^{\prime}$. If M is commutative and written additively, the submonoid generated by the family $\left(x_{i}\right)$ is the set of finite linear combinations of $x_{i}$ with coefficients in $\mathbb{N}$.

### 3.2.2 Divisibility

Let M be a commutative monoid. For $x, y \in \mathrm{M}, x$ is said to divide $y$, in which case it is written ${ }^{4} x<y$ if $(\exists z \in \mathbf{M}) y=z x$, and to strictly divide $y$ if $x\langle y$ and $y\langle x$. If $\mathbf{M}$

[^4]is regular and $e$ its only invertible element, then the divisibility relation is an order on M . The identity element is the smallest element for this order. An element of M is irreducible if it is minimal in $\mathrm{M}-\{e\}$, i.e. if it cannot be written $x y$ with $x \neq e$, $y \neq e$.

### 3.2.3 Support of a Family

Let $x=\left(x_{i}\right)_{i \in \mathrm{I}}$ be a family of elements in a monoid M with identity element $e$. The support of $x$ is the set of indices $i \in \mathrm{I}$ such that $x_{i} \neq e$.

If M is commutative and written additively, and if $\left(x_{i}\right)_{i \in \mathrm{I}}$ is a family of elements of $M$ with finite support, then set

$$
\sum_{i \in \mathrm{I}} x_{i}=\sum_{i \in \mathrm{~J}} x_{i}
$$

where J is the support of the family $\left(x_{i}\right)$, or an arbitrary finite family of I containing this support.

### 3.2.4 Free Commutative Monoids on a Set

Let I be a set, and $\mathbb{N}^{(\mathrm{I})}$ the additive monoid of families $\left(n_{i}\right)_{i \in \mathrm{I}}$ of elements of $\mathbb{N}$ with finite support; $\mathbb{N}^{(\mathrm{I})}$ is said to be the free commutative monoid on I . For $i \in \mathrm{I}$, let $e_{i}$ be the element $\left(\delta_{i j}\right)_{j \in \mathrm{I}}$ of $\mathbb{N}^{(\mathrm{I})}$ defined by $\delta_{i i}=1$ and $\delta_{i j}=0$ if $i \neq j$. The $\left(e_{i}\right)$ are the irreducible elements of $\mathbb{N}^{(\mathrm{I})}$.

UnIVERSAL PROPERTY. Let M be a commutative monoid and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of M . There is a unique homomorphism from $\mathbb{N}^{(\mathrm{I})}$ to M such that $f\left(e_{i}\right)=$ $x_{i}$ for all $i \in \mathrm{I}$.

Proof Write M additively. If $f$ satisfies the stated property, then for all $n=\left(n_{i}\right)_{i \in \mathrm{I}} \in$ $\mathbb{N}^{(1)}$,

$$
f(n)=f\left(\sum n_{i} e_{i}\right)=\sum n_{i} x_{i},
$$

and uniqueness follows. The map $f$ defined by

$$
f(n)=\sum n_{i} x_{i}
$$

for $n=\left(n_{i}\right)_{i \in \mathrm{I}} \in \mathbb{N}^{(\mathrm{I})}$ is the desired one.
The map $f$ is surjective if and only if the family $\left(x_{i}\right)$ generates M. If $f$ is an isomorphism, the family $\left(x_{i}\right)$ is said to be a basis of M.

### 3.2.5 Free Commutative Monoids

A monoid M is said to be free if it has a basis. Unlike other structures, a commutative monoid has a unique basis, up to permutation: the basis elements are necessarily irreducible.

A free commutative monoid $M$ is a lattice for the divisibility order: every finite subset has an upper bound called l.c.m., and, if it is not empty, a lower bound called g.c.d.. Writing M multiplicatively, for $x, y \in \mathrm{M}, x y=$ g.c.d. $(x, y) \cdot$ l.c.m. $(x, y)$.

### 3.2.6

Proposition Let M be a regular commutative monoid and e its only invertible element. Then M is free if and only if it satisfies the following two conditions:
(i) any decreasing sequence of elements of M for the divisibility order is stationary;
(ii) for any irreducible $m \in \mathrm{M}$ and $x, y \in \mathrm{M}, m\langle x y \Rightarrow m\langle x$ or $m\langle y$ (writing M multiplicatively).

Proof To show that conditions (i) and (ii) are necessary, we may assume that $\mathbf{M}=\mathbb{N}^{(\mathrm{I})}$. We use additive notation, and define $\psi: \mathbf{M} \rightarrow \mathbb{N}$ by $\psi(x)=\sum x_{i}$. The homomorphism $\psi$ is increasing and

$$
(x \leqslant y \text { and } \psi(x)=\psi(y)) \Rightarrow x=y .
$$

So (i) follows. If $m=e_{i}$ divides $x+y$, then $x_{i}+y_{i}>0$, and so $x_{i}>0$ or $y_{i}>0$. So (ii) holds.

Shifting back multiplicative notation, we show by contradiction that condition (i) implies that M is generated by the irreducible elements. Let $\mathrm{M}^{\prime}$ be the submonoid of M generated by the irreducible elements, and let $x \in \mathrm{M}-\mathrm{M}^{\prime}$. We inductively define a strictly decreasing sequence $\left(x_{n}\right)$ of elements of $\mathrm{M}-\mathrm{M}^{\prime}$ such that $x_{0}=x$. Suppose that $x_{n}$ has been defined. As $x_{n}$ is not irreducible, $x=y z$ for some $y \neq e$ and $z \neq e$, and we may assume that $y \notin \mathrm{M}^{\prime}$. Set $x_{n+1}=y$; then $x_{n+1}$ strictly divides $x_{n}$ since $z \neq e$. This contradicts condition (i).

Suppose condition (ii) holds. We show by induction on $r$ that if $m_{1} \ldots m_{r}=$ $m_{1}^{\prime} \ldots m_{r}^{\prime}$, where all $m_{i}$ and $m_{i}^{\prime}$ are irreducible, then $r=r^{\prime}$ and there is a permutation $\sigma \in \mathfrak{S}_{r}$ such that $m_{i}=m_{\sigma(i)}^{\prime}$ for $i \in\{1, \ldots, r\}$. This readily follows for $r=0$ or 1 . Suppose the property holds for $r-1$. Then $m_{r}$ divides some $m_{i}^{\prime}$, and hence is equal to it. We may assume this to be $m_{r^{\prime}}^{\prime}$. Then $m_{1} \ldots m_{r-1}=m_{1}^{\prime} \ldots m_{r^{\prime}-1}^{\prime}$ since M is regular, and hence by induction $r-1=r^{\prime}-1$ and there is a permutation $\sigma^{\prime} \in \mathfrak{S}_{r-1}$ such that $m_{i}=m_{\sigma^{\prime}(i)}^{\prime}$ for $i \leqslant r-1$. Extending $\sigma^{\prime}$ to $\sigma \in \mathfrak{S}_{r}$ by $\sigma(r)=r$, the result follows.

Hence, if M satisfies (i) and (ii), it is free.

### 3.2.7 Monoids Associated to Integral Domains

Let A be an integral domain. Write $\operatorname{Mon}(A)$ for the multiplicative quotient monoid of A by the equivalence relation identifying $x$ and $u x$ if $u$ is invertible. Let $\chi$ be the canonical map $A \rightarrow \operatorname{Mon}(A)$ and $\operatorname{Mon}^{*}(A)$ the image of $A^{*}=A-\{0\}$ in $\operatorname{Mon}(A)$. Then $\operatorname{Mon}^{*}(\mathrm{~A})$ is called the monoid associated to A.

For $x, y \in \mathrm{~A}^{*}, \chi(x)=\chi(y)$ if and only if $x$ and $y$ generate the same ideal. An element $x \in \mathrm{~A}^{*}$ is said to be irreducible if $\chi(x)$ is irreducible in $\operatorname{Mon}^{*}(\mathrm{~A})$, i. e. if $x$ is not invertible and if, for all pairs $(u, v)$ such that $x=u v, u$ or $v$ is invertible.

### 3.2.8 Unique Factorization Domains

An integral domain A is said to be a unique factorization domain (UFD) if $\mathrm{Mon}^{*}$ (A) is a free commutative monoid. This amounts to saying that every element $x \in \mathrm{~A}^{*}$ can be written as a product $u \cdot p_{1} \ldots p_{r}$ of irreducible elements $p_{i}$ and some invertible $u$, and this representation is unique up to permutation of the elements $p_{i}$ and multiplication of $u$ and $p_{i}$ by invertible elements.

Thanks to Proposition 3.2.6, A is a UFD if and only if it satisfies the following two conditions:
(i) every increasing sequence of principal ideals is stationary;
(ii) for any irreducible element $m \in \mathrm{~A}^{*}$, the ideal $\mathrm{A} m$ is prime.

If A is a UFD, every finite subset of Mon*(A) has a lower bound (g.c.d.) and an upper bound (l.c.m.) ${ }^{5}$ for divisibility. If $\left(x_{i}\right)$ is a family of finite elements of A , an element $x$ of A is a g.c.d. (resp. a l.c.m.) of $x_{i}$ if $\chi(x)$ is the g.c.d. (resp. the l.c.m.) of $\chi\left(x_{i}\right)$. If $x$ is a l.c.m. of $x_{i}$, then $\mathrm{A} x=\bigcap \mathrm{A} x_{i}$.

### 3.2.9

Examples The polynomial ring $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ for a field K will be seen to be a UFD. Likewise for the ring of formal series $\mathrm{K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$. The ring $\mathbb{C}\left\{\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right\}$ of convergent series can be shown to be a UFD; the ring of germs at 0 of analytic functions on a Banach space is a UFD.
Counterexample Let K be a field. The subring of $\mathrm{K}[\mathrm{X}, \mathrm{Y}]$ consisting of polynomials containing only even powers is not a UFD. Indeed, the element XY of A is irreducible because it has minimal degree. It divides neither $\mathrm{X}^{2}$, nor $\mathrm{Y}^{2}$, but does divide $\mathrm{X}^{2} \mathrm{Y}^{2}$, and so A does not satisfy condition (ii) of 3.2.8.

[^5]
### 3.2.10

Theorem Every PID is a UFD.
Proof Let A be a PID. We show that it satisfies conditions (i) and (ii) of 3.2.8.
(i) Let $\left(\mathrm{I}_{n}\right)$ be an increasing sequence of ideals and set $\mathrm{I}=\bigcup \mathrm{I}_{n}$. There exists $x$ such that $\mathrm{I}=\mathrm{A} x$. This element $x$ is in $\mathrm{I}_{n}$; then $\mathrm{A} x \subset \mathrm{I}_{n} \subset \mathrm{I}=\mathrm{A} x$. So $\mathrm{I}_{n}=\mathrm{I}$, and $\mathrm{I}_{n^{\prime}}=\mathrm{I}_{n}$ for $n^{\prime} \geqslant n$.
(ii) follows from the following lemma:

Lemma Let A be a PID and $\mathrm{I}=\mathrm{Ax}$ a nonzero ideal of A . The following conditions are equivalent:
(i) the element $x$ is irreducible;
(ii) the ideal I is prime;
(iii) the ideal I is maximal.

Proof (iii) $\Rightarrow$ (ii) by 3.1.4, Corollary.
(ii) $\Rightarrow$ (i): Suppose that I is prime and $x=y z$. Then $y$ or $z$ is in I. Without loss of generality, assume $y \in \mathrm{I}$. Then $x$ divides $y$ and $z$ is invertible. Hence $x$ is irreducible.
(i) $\Rightarrow$ (iii): Let $\mathbf{J}=\mathrm{A} y$ be such that $\mathrm{I} \supsetneqq \mathbf{J} \supsetneqq \mathrm{A}$. Then $y$ is not invertible and $x$ can be written $y z$ where $z$ is not invertible. This is impossible if $x$ is irreducible.

Remark For invertible $x$, none of the conditions (i), (ii), (iii) are satisfied.

### 3.2.11

Proposition Let A be a PID, $\left(x_{i}\right)$ a finite family of elements of A . Then $x \in \mathrm{~A}$ is a g.c.d. of $x_{i}$ if and only if $\mathrm{A} x=\sum \mathrm{A} x_{i}$.

Proof Let $x \in \mathrm{~A}$ be such that $\mathrm{A} x=\sum \mathrm{A} x_{i}$. For $y \in \mathrm{~A},(\forall i) y\left\langle x_{i} \Leftrightarrow(\forall i) \mathrm{A} y \supset\right.$ $\mathrm{A} x_{i} \Leftrightarrow \mathrm{~A} y \supset \sum \mathrm{~A} x_{i} \Leftrightarrow \mathrm{~A} y \supset \mathrm{~A} x \Leftrightarrow y\left\langle x\right.$. So $x$ is a g.c.d. of $x_{i}$. If $x^{\prime}$ is any other g.c.d. of $x_{i}$, then $x=u x$ for some invertible $u$. The proposition follows.

Corollary (Bezout's identity) Under the assumptions of the proposition, if $x$ is a g.c.d. of $x_{i}$, there is a family $\left(a_{i}\right)$ of elements of A such that $x=\sum a_{i} x_{i}$.

## Exercises 3.2. (UFDs)

1.-Let K be a field. The subring $\mathrm{K}\left[t^{2}, t^{3}\right]$ of $\mathrm{K}[t]$ is not a UFD.
2.-The ring $\mathbb{R}[X, Y, Z] /\left(X^{2}+Y^{2}+Z^{2}\right)$ is a UFD, but not $\mathbb{C}[X, Y, Z] /\left(X^{2}+Y^{2}+Z^{2}\right)$ (it is isomorphic to the ring given in the counterexample in 3.2.9).
3.-Let A be an integral domain, and $K$ its field of fractions. Then $A$ is said to be integrally closed if, whenever $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$ are monic polynomials in $\mathrm{K}[\mathrm{X}]$ then $P_{1} P_{2} \in A[X]$ implies $P_{1} \in A[X]$ and $P_{2} \in A[X]$. (this condition can be shown to be
equivalent to the apparently weaker one: "any root of a monic polynomial of $\mathrm{A}[\mathrm{X}]$ in K belongs to A".)

Show that, if A is integrally closed, then the multiplicative monoid M of monic polynomials of $\mathrm{A}[\mathrm{X}]$ is a free commutative monoid.
4.-Let E be a complex Banach space, and $\mathcal{O}_{\mathrm{E}}$ the ring of germs at the origin of holomorphic functions in the neighbourhood of 0 in E . The aim is to show that the ring $\mathcal{O}_{\mathrm{E}}$ is a UFD by using Exercise 3 and assuming the following results, where F denotes a hyperplane in E :

Removable singularity theorem. The ring $\mathcal{O}_{\mathrm{F}}$ is integrally closed.
A polynomial $\mathrm{P} \in \mathcal{O}_{\mathrm{F}}[\mathrm{Z}]$ is normal if it is monic and all its coefficients except the leading one are functions vanishing at 0 . Let R be the set of $f \in \mathcal{O}_{\mathrm{F} \times \mathbb{C}}$ such that $f(0, Z) \in \mathbb{C}[Z]$ is not zero.

Preparation Theorem. Every element $f$ of R can be uniquely written as $f=$ $u \cdot \mathrm{P}$, where $u$ is an invertible element of $\mathcal{O}_{\mathrm{F} \times \mathbb{C}}$ and $\mathrm{P} \in \mathcal{O}_{\mathrm{F}}[\mathrm{Z}]$ is a normal polynomial.
(a) Show that the normal polynomials of $\mathcal{O}_{\mathrm{F}}[\mathrm{Z}]$ together with multiplication form a free commutative monoid. Deduce that $\mathrm{M}=\mathrm{R} / \mathrm{G}_{m}\left(\mathcal{O}_{\mathrm{F} \times \mathbb{C}}\right)$ is a free commutative monoid. Let $\mathrm{G}_{m}(\mathrm{~A})$ be the group of invertible elements of A .
(b) For any line D in E passing through 0 , let $\mathrm{R}^{\mathrm{D}}$ be the monoid of functions $f \in \mathcal{O}_{\mathrm{E}}$ that do not vanish everywhere in the neighbourhood of 0 in D . Set $\mathrm{M}^{\mathrm{D}}=\mathrm{R}^{\mathrm{D}} / \mathrm{G}_{m}\left(\mathcal{O}_{\mathrm{E}}\right)$. Show that $\mathrm{M}^{\mathrm{D}}$ is a free monoid. Deduce that $\mathcal{O}_{\mathrm{E}}$ is a UFD (criterion 3.2.6 may be used).
5.-Let $E$ be a complex Banach space, $U$ a connected open subset of $E$ and $K$ a compact subset of $\mathbb{C}$. Let $\mathcal{O}(\mathrm{U})$ be the ring of holomorphic functions on U . A polynomial $\mathrm{P} \in \mathcal{O}(\mathrm{U})[\mathrm{Z}]$ is said to be K-normal if, for all $x \in \mathrm{U}$, every root of the polynomial $\mathrm{P}(x, \mathrm{Z}) \in \mathbb{C}[Z]$ is in the interior K of K . Show that the K -normal polynomials of $\mathcal{O}(\mathrm{U})[\mathrm{Z}]$ together with multiplication form a free commutative monoid.
6.-Let A be a UFD all of whose nonzero prime ideals are maximal. The aim is to show that A is principal.
(a) Show that, for all $x, y \in \mathrm{~A}^{*}$, there exist $u, v \in \mathrm{~A}$ such that $u x+v y$ is a g.c.d. of $x$ and $y$. (First consider the case where $x$ is irreducible, then the case where $y /$ g.c.d. $(x, y)=$ l.c.m. $(x, y) / x$ is irreducible, and finally the general case.)
(b) Let $\mathrm{I} \subset \mathrm{A}$ be a non-zero ideal. Show that there is a g.c.d. $d \in \mathrm{~A}^{*}$ of all the elements of $\mathrm{I} \cap \mathrm{A}^{*}$. Conclude.
(c) Show that the ring A of real analytic functions with period $2 \pi$ (3.1, Exercise 12) is not a UFD. Give an example of an element of $\operatorname{Mon}^{*}(\mathrm{~A})$ whose decomposition into irreducible elements is not unique.

### 3.3 Modules

Throughout this section, A denotes a (commutative) ring.

### 3.3.1

Definition A module over A or A-module is a set E together with an internal composition law called addition and written $(x, y) \mapsto x+y$ and an external law $\mathrm{A} \times \mathrm{E} \rightarrow \mathrm{E}$ written $(\lambda, x) \mapsto \lambda x$, satisfying the following conditions:
$\left(\mathrm{M}_{1}\right)$ The additive set E is a commutative group;
$\left(\mathrm{M}_{2}\right)(\forall \lambda, \mu \in \mathrm{A})(\forall x \in \mathrm{E}) \quad \lambda(\mu x)=(\lambda \mu) x$
$(\forall \lambda, \mu \in \mathrm{A})(\forall x \in \mathrm{E}) \quad(\lambda+\mu) x=\lambda x+\mu x$
$(\forall \lambda \in \mathrm{A})(\forall x, y \in \mathrm{E}) \quad \lambda(x+y)=\lambda x+\lambda y$
$(\forall x \in \mathrm{E}) \quad 1 x=x$.
Examples (1) All A-modules for a field A are vectorial spaces over A.
(2) Every commutative group has a unique $\mathbb{Z}$-module structure.
(3) Let $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism. Setting $a \cdot b=\varphi(a) b$ for $a \in \mathrm{~A}$ and $b \in \mathrm{~B}$ gives B an A-module structure.
(4) The space $\mathscr{D}^{\prime}$ of distributions on $\mathbb{R}$ together with multiplication is a module over the ring $\mathscr{E}$ of $\mathrm{C}^{\infty}$-functions. The space $\mathscr{E}$ together with convolution is a module over the convolution ring of distributions with compact support.
(5) Let X be a topological space and E a vector bundle over X . The space of continuous sections of E is a module over the ring of continuous functions on X .

### 3.3.2

Definition Let E and F be two A -modules. A map $f$ from E to F is said to be a homomorphism or A-linear if

$$
\begin{aligned}
& \left(\mathrm{L}_{1}\right)(\forall x, y \in \mathrm{E}) \quad f(x+y)=f(x)+f(y) \\
& \left(\mathrm{L}_{2}\right)(\forall x \in \mathrm{E})(\forall \lambda \in \mathrm{A}) \quad f(\lambda x)=\lambda f(x)
\end{aligned}
$$

More generally, let $h: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism, E an A-module, F a B -module and $f$ a map from E to F . Then $f$ is said to be a $h$-morphism or $h$-linear if $\left(L_{1}\right)$ holds and
$\left(\mathrm{L}_{2}^{\prime}\right)(\forall x \in \mathrm{E})(\forall \lambda \in \mathrm{A}) \quad f(\lambda x)=h(\lambda) f(x)$.
Let $h: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism, E a B -module. Define an A-module structure on E by setting for $x \in \mathrm{E}$ and $\lambda \in \mathrm{A}$

$$
\lambda \cdot x=h(\lambda) x .
$$

The A-module structure on E is said to be obtained by scalar restriction. The identity map from the A-module E onto the B -module E is a $h$-morphism.

### 3.3.3 Submodules

Let E be an A -module. A subset $\mathrm{E}^{\prime}$ of E is said to be an A -submodule (or simply a submodule) of E if
$\left(\mathrm{SM}_{1}\right) \mathrm{E}^{\prime}$ is an additive subgroup of E ;
$\left(\mathrm{SM}_{2}\right)\left(\forall x \in \mathrm{E}^{\prime}\right)(\forall \lambda \in \mathrm{A}) \quad \lambda x \in \mathrm{E}^{\prime}$.
Every submodule $\mathrm{E}^{\prime}$ of an A -module E is an A -module for the induced laws, and the canonical injection $\iota: \mathrm{E}^{\prime} \rightarrow \mathrm{E}$ is a homomorphism.

Examples (1) The A-submodules of A are the ideals of A.
(2) Let E and F be two A-modules, $f: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism. The kernel of $f$, written $\operatorname{Ker} f$, is the submodule $f^{-1}(0)$ of E . The image of $f$, written $\operatorname{Im} f$, is the submodule $f(\mathrm{E})$ of F .

### 3.3.4 Torsion

Suppose that A is integral and let E be an A-module. An element $x$ of E is said to be a torsion element if the homomorphism $\lambda \mapsto \lambda x$ from A to E is not injective. The set of torsion elements of E is a submodule T of E called the torsion submodule of E ; E is torsion-free if $\mathrm{T}=\{0\}$, and a torsion module if $\mathrm{T}=\mathrm{E}$.

Let E and F be two modules over an integral domain. Every homomorphism $f: \mathrm{E} \rightarrow \mathrm{F}$ maps the torsion submodule of E to that of F .
Examples Suppose that A is integral and let I be a nonzero ideal in A. As an Amodule, the quotient ring $\mathrm{A} / \mathrm{I}$ is a torsion module.

The $\mathbb{Z}$-module $\mathbb{Q} / \mathbb{Z}$ is a torsion module.

### 3.3.5 Generating Families

Let E be an A -module and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of E . An A-submodule of E generated by the family $\left(x_{i}\right)_{i \in \mathrm{I}}$ is the smallest submodule $\mathrm{E}^{\prime}$ of E such that

$$
(\forall i \in \mathrm{I}) \quad x_{i} \in \mathrm{E}^{\prime} .
$$

It is also the set of finite linear combinations of $x_{i}$ with coefficients in A. If the submodule generated by $\left(x_{i}\right)_{i \in \mathrm{I}}$ is E , the latter is a generating family of E .

The A-module E is finitely generated if it has a finite generating family, and it is cyclic if it is generated by one element.

Let $\left(\mathrm{E}_{i}\right)$ be a family of submodules of E . The submodule generated by the union of all $\mathrm{E}_{i}$ is the subgroup $\sum \mathrm{E}_{i}$ generated by this union.

### 3.3.6 Quotient Modules

Let E be an A -module, $\mathrm{E}^{\prime}$ a submodule of E . A congruence modulo $\mathrm{E}^{\prime}$ is the relation $x-y \in \mathrm{E}^{\prime}$ between elements $x$ and $y$ of E . It is written

$$
x \equiv y \quad\left(\bmod \mathrm{E}^{\prime}\right),
$$

and is an equivalence relation. Relations $x^{\prime} \equiv y^{\prime}\left(\bmod \mathrm{E}^{\prime}\right)$ and $x^{\prime \prime} \equiv y^{\prime \prime}\left(\bmod \mathrm{E}^{\prime}\right)$ imply $x^{\prime}+x^{\prime \prime} \equiv y^{\prime}+y^{\prime \prime}\left(\bmod \mathrm{E}^{\prime}\right)$ and $\lambda x^{\prime} \equiv \lambda y^{\prime}\left(\bmod \mathrm{E}^{\prime}\right)$ for all $\lambda \in \mathrm{A}$.

The quotient set of E by congruence modulo $\mathrm{E}^{\prime}$ has a unique A -module structure for which the quotient map is A-linear, and with this structure, is called the quotient module of E by $\mathrm{E}^{\prime}$ and is written $\mathrm{E} / \mathrm{E}^{\prime}$.
Universal property. Let E be an A-module, $\mathrm{E}^{\prime}$ a submodule of E . For all $\mathrm{A}-$ modules F and homomorphisms $f: \mathrm{E} \rightarrow \mathrm{F}$ such that $\left.f\right|_{\mathrm{E}^{\prime}}=0$, there is a unique homomorphism $\bar{f}: \mathrm{E} / \mathrm{E}^{\prime} \rightarrow \mathrm{F}$ satisfying $f=\bar{f} \circ \chi$, where $\chi$ is the canonical map $\mathrm{E} \rightarrow \mathrm{E} / \mathrm{E}^{\prime}$.

In other words, the map $g \mapsto g \circ \chi$ is a bijection from $\operatorname{Hom}\left(\mathrm{E} / \mathrm{E}^{\prime}, \mathrm{F}\right)$ onto the set of $f \in \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ for which $\left.f\right|_{\mathrm{E}^{\prime}}=0$.

Example Every cyclic A-module is isomorphic to a module A/I, where I is an ideal of A . Indeed, if $\mathrm{E}=\mathrm{A} x$, then the map $\lambda \mapsto \lambda x$ from A to E is surjective, giving an isomorphism $\mathrm{A} / \mathrm{I} \rightarrow \mathrm{E}$, where I is the l'annihilator of $x$, i.e. the ideal consisting of elements $\lambda \in A$ such that $\lambda x=0$.

### 3.3.7 Canonical Factorization

Let E and F be two A-modules, $f: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism. There is a unique isomorphism $\tilde{f}: \mathrm{E} / \operatorname{Ker} f \rightarrow \operatorname{Im} f$ such that the diagram

commutes.
This diagram gives the canonical factorization of $f$.
The cokernel of $f$, written Coker $f$, is the quotient module $\mathrm{F} / \operatorname{Im} f$.

### 3.3.8 Exact Sequences

Let $\mathrm{E}_{m} \rightarrow \cdots \rightarrow \mathrm{E}_{i+1} \xrightarrow{f_{i+1}} \mathrm{E}_{i} \xrightarrow{f_{i}} \mathrm{E}_{i-1} \rightarrow \cdots \rightarrow \mathrm{E}_{n}$ be a sequence of A-module homomorphisms (the decreasing order of indices has been chosen with homology in mind, see 3.10.2). It is said to be an exact sequence if $\operatorname{Im} f_{i+1}=\operatorname{Ker} f_{i}$ for $m-1 \geqslant i \geqslant n+1$.

Example Let E and F be two A-modules, $f: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism. $0 \rightarrow$ Ker $f \rightarrow \mathrm{E} \xrightarrow{f} \mathrm{~F} \rightarrow$ Coker $f \rightarrow 0$ is an exact sequence.

In particular $f$ is injective (resp. surjective) if and only if $0 \rightarrow \mathrm{E} \xrightarrow{f} \mathrm{~F}$ (resp. $\mathrm{E} \xrightarrow{f} \mathrm{~F} \rightarrow 0)$ is an exact sequence.

If $0 \rightarrow \mathrm{E}^{\prime} \rightarrow \mathrm{E} \rightarrow \mathrm{E}^{\prime \prime} \rightarrow 0$ is an exact sequence of A -modules, then $\mathrm{E}^{\prime}$ can be identified with a submodule of E and $\mathrm{E}^{\prime \prime}$ to the quotient module $\mathrm{E} / \mathrm{E}^{\prime}$. The exact sequence is then written $0 \rightarrow \mathrm{E}^{\prime} \rightarrow \mathrm{E} \rightarrow \mathrm{E} / \mathrm{E}^{\prime} \rightarrow 0$.

An exact sequence of the form $0 \rightarrow \mathrm{E}^{\prime} \rightarrow \mathrm{E} \rightarrow \mathrm{E}^{\prime \prime} \rightarrow 0$ is called a short exact sequence. If $\mathrm{E}_{m} \rightarrow \cdots \rightarrow \mathrm{E}_{i} \rightarrow \cdots \rightarrow \mathrm{E}_{n}$ is an exact sequence, there are short exact sequences $0 \rightarrow \mathrm{E}_{i}^{\prime} \rightarrow \mathrm{E}_{i} \rightarrow \mathrm{E}_{i-1}^{\prime} \rightarrow 0$ for $m \geqslant i \geqslant n$, with $\mathrm{E}_{i}^{\prime}=\operatorname{Ker} f_{i}=\operatorname{Im} f_{i+1} \approx$ Coker $f_{i+2}$ for $m-2 \geqslant i \geqslant n+1, \mathrm{E}_{m-1}^{\prime}=\operatorname{Ker} f_{m-1}=\operatorname{Im} f_{m}, \mathrm{E}_{m}^{\prime}=\operatorname{Ker} f_{m}, \mathrm{E}_{n}^{\prime}=$ $\operatorname{Im} f_{n-1} \approx$ Coker $f_{n+2}, \mathrm{E}_{n-1}^{\prime}=\operatorname{Coker} f_{n+1}$.

### 3.3.9

Proposition Let $0 \rightarrow \mathrm{E}^{\prime} \xrightarrow{u} \mathrm{E} \xrightarrow{v} \mathrm{E}^{\prime \prime} \rightarrow 0$ be an exact sequence of A -modules, $\left(x_{i}^{\prime}\right)_{i \in \mathrm{I}^{\prime}}$ and $\left(x_{i}^{\prime \prime}\right)_{i \in \mathrm{I}^{\prime \prime}}$ generating families of $\mathrm{E}^{\prime}$ and $\mathrm{E}^{\prime \prime}$ respectively. Suppose that $\mathrm{I}^{\prime} \cap \mathrm{I}^{\prime \prime}=\varnothing$ and set $\mathrm{I}=\mathrm{I}^{\prime} \cup \mathrm{I}^{\prime \prime}$. For $i \in \mathrm{I}^{\prime}$, let $x_{i}=u\left(x_{i}^{\prime}\right)$ and, for each $i \in \mathrm{I}^{\prime \prime}$, choose $x_{i} \in \mathrm{E}$ such that $v\left(x_{i}\right)=x_{i}^{\prime \prime}$. Then the family $\left(x_{i}\right)_{i \in \mathrm{I}}$ generates E .

Proof Let $x \in \mathrm{E}$. There is a family $\left(a_{i}\right)_{i \in \mathrm{I}^{\prime \prime}}$ with finite scalar support such that $v(x)=\sum_{i \in \mathrm{I}^{\prime \prime}} a_{i} x_{i}^{\prime \prime}$. Then $x-\sum_{i \in \mathrm{I}^{\prime \prime}} a_{i} x_{i}$ is in Ker $v$, and so is of the form $u\left(x^{\prime}\right)$. There is a family $\left(a_{i}\right)_{i \in I^{\prime}}$ of finite scalar support such that $x^{\prime}=\sum_{i \in I^{\prime}} a_{i} x_{i}^{\prime}$. Then $x=\sum_{i \in \mathrm{I}} a_{i} x_{i}$.

Corollary Let $0 \rightarrow \mathrm{E}^{\prime} \rightarrow \mathrm{E} \rightarrow \mathrm{E}^{\prime \prime} \rightarrow 0$ be an exact sequence of A -modules. If $\mathrm{E}^{\prime}$ and $\mathrm{E}^{\prime \prime}$ are of finitely generated, then so is E .

### 3.3.10 Complement Submodules

Let E be an A -module, the submodules $\mathrm{E}_{1}$ and $\mathrm{E}_{2}$ are complements, or E is a direct sum of the submodules $\mathrm{E}_{1}$ and $\mathrm{E}_{2}$, and is written $\mathrm{E}=\mathrm{E}_{1} \oplus \mathrm{E}_{2}$, if any element of E
is of the form $x_{1}+x_{2}$, with $x_{1} \in \mathrm{E}_{1}$ and $x_{2} \in \mathrm{E}_{2} . \mathrm{E}_{1}$ and $\mathrm{E}_{2}$ are complements if and only if $\mathrm{E}_{1}+\mathrm{E}_{2}=\mathrm{E}$ and $\mathrm{E}_{1} \cap \mathrm{E}_{2}=0$.

A submodule of an A-module E is a direct factor of E if it has a complement. There are modules with submodules that are not direct factors. For example, the submodule $2 \mathbb{Z}$ is not a direct factor of $\mathbb{Z}$ because the quotient being a torsion module, there are no injections from $2 \mathbb{Z}^{6}$ to $\mathbb{Z}$.

More generally, let $\left(\mathrm{E}_{i}\right)$ be a family of submodules of E . Then E is said to be a direct sum of the submodules $\mathrm{E}_{i}$ if all $x \in \mathrm{E}$ can be uniquely written in the form $\sum x_{i}$, with $x_{i} \in \mathrm{E}_{i}$ for all $i$, the family $\left(x_{i}\right)_{i \in \mathrm{I}}$ having finite support.

### 3.3.11 Projections

Let E be an A -module. An endomorphism $p$ of E is a projection if $p \circ p=p$.
Proposition The map $p \mapsto(\operatorname{Im} p, \operatorname{Ker} p)$ is a bijection from the set of projections on E onto the set of pairs consisting of complement submodules of E .

Proof If $p$ is a projection of E , then $\mathrm{E}=\operatorname{Im} p \oplus \operatorname{Ker} p$. Indeed, any $x \in \mathrm{E}$ can be written $p(x)+(x-p(x))$, with $p(x) \in \operatorname{Im} p$ and $x-p(x) \in \operatorname{Ker} p$. So $\operatorname{Im} p+\operatorname{Ker} p=\mathrm{E}$; on the other hand, $p(x)=x$ for $x \in \operatorname{Im} p$. Thus $\operatorname{Im} p \cap \operatorname{Ker} p=$ 0 . Next let $\mathrm{E}_{1}$ and $\mathrm{E}_{2}$ be two complement submodules of E . A projection having image $\mathrm{E}_{1}$ and kernel $\mathrm{E}_{2}$ is necessarily the map $x_{1}+x_{2} \mapsto x_{1}$; it is the desired one, proving bijectivity.

### 3.3.12 Split Exact Sequences

Proposition and Definition Let $0 \rightarrow \mathrm{E}^{\prime} \xrightarrow{u} \mathrm{E} \xrightarrow{v} \mathrm{E}^{\prime \prime} \rightarrow 0$ be an exact sequence. The following conditions are equivalent:
(i) $u\left(\mathrm{E}^{\prime}\right)$ is a direct factor of E ;
(ii) there is a linear retraction of $u$, i.e. a homomorphism $\rho: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ such that $\rho \circ u=1_{\mathrm{E}^{\prime}}$;
(iii) there is a linear section of $v$, i.e. a homomorphism $\sigma: \mathrm{E}^{\prime \prime} \rightarrow \mathrm{E}$ such that $v \circ \sigma=$ $1_{\mathrm{E}^{\prime \prime}}$.

If these conditions hold, the exact sequence is said to be split.
Proof (i) implies (ii). Identify $\mathrm{E}^{\prime}$ with its image in E under $u$. Let S be a complement of $\mathrm{E}^{\prime}$ in E . The projection of E onto $\mathrm{E}^{\prime}$ with kernel S is a retraction.

[^6](i) implies (iii). Let S be a complement of $\operatorname{Ker} v=u\left(\mathrm{E}^{\prime}\right)$, the homomorphism $v$ induces an isomorphism from $S$ onto $\mathrm{E}^{\prime \prime}$ and the composite of the inverse isomorphism and the canonical injection from S to E is a section.
(ii) implies (i) (resp. (iii) implies (i)). If $\rho$ is a retraction of $u$ (resp. $\sigma$ is a section of $v$ ), then $u \circ \rho($ resp. $\sigma \circ v)$ is a projection of E whose kernel (resp. image) is $u\left(\mathrm{E}^{\prime}\right)$.

### 3.3.13 Direct Products and Sums

Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ be a family of A-modules. Define an A-module structure on the product set $\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$ by setting

$$
\left(x_{i}\right)_{i \in \mathrm{I}}+\left(y_{i}\right)_{i \in \mathrm{I}}=\left(x_{i}+y_{i}\right)_{i \in \mathrm{I}}, \quad a\left(x_{i}\right)_{i \in \mathrm{I}}=\left(a x_{i}\right)_{i \in \mathrm{I}}
$$

The module obtained is called the direct product of $\mathrm{E}_{i}$. The submodule of $\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$ consisting of families with finite support is called the direct sum of $\mathrm{E}_{i}$ and is written $\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$. If I is finite, then $\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}=\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$. For $k \in \mathrm{I}, p_{k}: \prod_{i \in \mathrm{I}} \mathrm{E}_{i} \rightarrow \mathrm{E}_{k}$ defined by $p_{k}\left(\left(x_{i}\right)_{i \in \mathrm{I}}\right)=x_{k}$ is a module homomorphism called the coordinate map of index $k$, and the map $\iota_{k}: \mathrm{E}_{k} \rightarrow \bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$ defined by $\iota_{k}(x)=\left(x_{i}\right)_{i \in \mathrm{I}}$, where $x_{k}=x$ and $x_{i}=0$ if $i \neq k$ is a module homomorphism called the canonical injection of index $k$.

The module $\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$ is a direct sum of submodules $\iota_{i}\left(\mathrm{E}_{i}\right)$. When no confusion arises, $\mathrm{E}_{i}$ is identified with $\iota_{i}\left(\mathrm{E}_{i}\right)$. We avoid doing so when several $\mathrm{E}_{i}$ are equal.

The module $\prod_{i \in \mathrm{I}} \mathrm{E} i$ (resp. $\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$ ), together with the maps $p_{i}$ (resp. $\iota_{i}$ ), is a product (resp. sum) of the modules $\mathrm{E}_{i}$ in the category of A-modules and linear maps.

In particular, let E be an A -module and $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ a family of submodules of E , the family of canonical inclusions of $\mathrm{E}_{i}$ in E defines a homomorphism $f$ from $\bigoplus \mathrm{E}_{i}$ to E given by $\left(x_{i}\right)_{i \in \mathrm{I}} \mapsto \sum x_{i}$.

It is an isomorphism (resp. surjective) if and only if E is the direct sum of the submodules $\mathrm{E}_{i}$ (resp. generated by $\left.\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}\right)$. When $\mathrm{I}=\{1,2\}, f$ is injective if and only if $\mathrm{E}_{1} \cap \mathrm{E}_{2}=\{0\}$.

### 3.3.14 Direct and Inverse Limits

Let $\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}},\left(f_{i}^{j}\right)\right)$ be a direct (resp. inverse) system of A-modules. The set $\mathrm{E}_{\infty}=$ $\underset{\longrightarrow}{\lim } \mathrm{E}_{i}$ (resp. $\lim _{\leftrightarrows} \mathrm{E}_{i}$ ) has a unique A-module structure with respect to which the $\overrightarrow{\text { canonical maps }} \stackrel{\overleftarrow{\mathrm{E}_{i}} \rightarrow \mathrm{E}_{\infty}}{ }$ (resp. $\mathrm{E}_{\infty} \rightarrow \mathrm{E}_{i}$ ) are linear. With this structure, $\mathrm{E}_{\infty}$ is the direct (resp. inverse) limit of the $\mathrm{E}_{i}$ in the category $\mathrm{A}-\mathscr{M} \boldsymbol{O}$.

Proposition Let I be a directed set, $\mathscr{E}=\left(\left(\mathrm{E}_{i}\right),\left(f_{i}^{j}\right)\right)$ and $\mathscr{F}=\left(\left(\mathrm{F}_{i}\right),\left(g_{i}^{j}\right)\right)$ direct systems of modules and $\left(\varphi_{i}\right)$ a morphism from $\mathscr{E}$ to $\mathscr{F}$. Let $\mathrm{E}=\underset{\longrightarrow}{\lim } \mathrm{E}_{i}, \mathrm{~F}=\underset{\longrightarrow}{\lim } \mathrm{F}_{i}$,
and $\varphi: \mathrm{E} \rightarrow \mathrm{F}$ be the morphism induced by $\varphi_{i}$ by passing to the direct limit. Then $\operatorname{Ker} \varphi=\xrightarrow{\lim } \operatorname{Ker} \varphi_{i}, \operatorname{Im} \varphi=\xrightarrow{\lim } \operatorname{Im} \varphi_{i}$ and $\operatorname{Coker} \varphi=\underline{\lim } \operatorname{Coker} \varphi_{i}$.

Proof It suffices to check that $\operatorname{Ker} \varphi($ resp. $\operatorname{Im} \varphi$, resp. $\operatorname{Coker} \varphi)$ with the canonical maps $\operatorname{Ker} \varphi_{i} \rightarrow \operatorname{Ker} \varphi$ (resp. etc.), satisfies conditions (i), (ii) and (iii) of 2.6.5. Condition (i) is immediate in all three cases. We show that $\operatorname{Ker} \varphi$ satisfies (ii). Let $x \in \operatorname{Ker} \varphi$. There exist $i \in \mathrm{I}$ and $x_{i} \in \mathrm{E}_{i}$ such that $\sigma_{i}\left(x_{i}\right)=x$, where $\sigma_{i}: \mathrm{E}_{i} \rightarrow \mathrm{E}$ is the canonical map. Then $\tau_{i}\left(\varphi_{i}\left(x_{i}\right)\right)=\varphi\left(\sigma_{i}\left(x_{i}\right)\right)=0$, where $\tau_{i}: \mathrm{F}_{i} \rightarrow \mathrm{~F}$ is a canonical map, and so there exists $j \geqslant i$ such that $g_{i}^{j}\left(\varphi_{i}\left(x_{i}\right)\right)=0$. Hence $x_{j}=f_{i}^{j}\left(x_{i}\right) \in \operatorname{Ker} \varphi_{j}$, and $x=\sigma_{j}\left(x_{j}\right) \in \sigma_{j}\left(\operatorname{Ker} \varphi_{j}\right)$.

Condition (ii) for $\operatorname{Im} \varphi$ and Coker $\varphi$ is immediate. So is condition (iii) for $\operatorname{Ker} \varphi$ and $\operatorname{Im} \varphi$; for Coker $\varphi$, it follows from condition (ii) for $\operatorname{Im} \varphi$,

Remark Let $\mathscr{E}$ and $\mathscr{F}$ be inverse systems, $\left(\varphi_{i}\right)$ a morphism from $\mathscr{E}$ to $\mathscr{F}$ and $\varphi$ the morphism induced by $\varphi_{i}$ by passing to the inverse limit. Then $\operatorname{Ker} \varphi=\lim \operatorname{Ker} \varphi_{i}$, but in general $\operatorname{Im} \varphi \neq \lim \operatorname{Im} \varphi_{i}$ and $\operatorname{Coker} \varphi \neq \lim \operatorname{Coker} \varphi_{i}(2.5$, Exercise 3$)$.

Corollary Let I be a directed set. For $m \geqslant k \geqslant n$, let $\mathscr{E}_{k}=\left(\left(\mathrm{E}_{k, i}\right)_{i \in \mathrm{I}},\left(f_{k, i}^{j}\right)\right)$ be a direct system of modules; for $m \geqslant k>n$, let $\left(u_{k, i}\right)_{i \in \mathrm{I}}$ be a morphism from $\mathscr{E}_{k}$ to $\mathscr{E}_{k-1}$. If, for all $i$,

$$
\mathrm{E}_{m, i} \rightarrow \cdots \rightarrow \mathrm{E}_{k, i} \xrightarrow{u_{k, i}} \mathrm{E}_{k-1, i} \rightarrow \cdots \rightarrow \mathrm{E}_{n, i}
$$

is an exact sequence, passing to the direct limit gives an exact sequence

$$
\mathrm{E}_{m} \rightarrow \cdots \rightarrow \mathrm{E}_{k} \xrightarrow{u_{k}} \mathrm{E}_{k-1} \rightarrow \cdots \rightarrow \mathrm{E}_{n} .
$$

### 3.3.15 Modules of Homomorphisms

Let E and F be two A-modules. The set $\operatorname{Hom}(\mathrm{E}, \mathrm{F})$ of homomorphisms from E to F is a submodule of $\mathrm{F}^{\mathrm{E}}$. With the induced structure it is an A-module. ${ }^{7}$

The composition law $(f, g) \mapsto g \circ f$ from $\operatorname{Hom}(\mathrm{E}, \mathrm{F}) \times \operatorname{Hom}(\mathrm{F}, \mathrm{G})$ to $\operatorname{Hom}(\mathrm{E}, \mathrm{G})$ is bilinear.

Let $u: \mathrm{E}^{\prime} \rightarrow \mathrm{E}$ be a homomorphism of A-modules. Define ${ }^{8}$ a homomorphism $u^{*}:$ $\operatorname{Hom}(\mathrm{E}, \mathrm{F}) \rightarrow \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}\right)$ by $u^{*}(f)=f \circ u$. Let $v: \mathrm{F} \rightarrow \mathrm{F}^{\prime}$ be a homomorphism of A-modules. Define a homomorphism $v_{*}: \operatorname{Hom}(\mathrm{E}, \mathrm{F}) \rightarrow \operatorname{Hom}\left(\mathrm{E}, \mathrm{F}^{\prime}\right)$ by $v_{*}(f)=$ $v \circ f$. Then there is a commutative diagram

[^7]

The map $(u, v) \mapsto v_{*} u^{*}$ from $\operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{E}\right) \times \operatorname{Hom}\left(\mathrm{F}, \mathrm{F}^{\prime}\right)$ to $\operatorname{Hom}(\operatorname{Hom}(\mathrm{E}, \mathrm{F})$, $\left.\operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)\right)$ is bilinear.

Then $(\mathrm{E}, \mathrm{F}) \mapsto \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is a functor from $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{\sigma} \times \mathrm{A}-\mathscr{M} \sigma d$ to $\mathrm{A}-\mathscr{M} \sigma d$. Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ and $\left(\mathrm{F}_{j}\right)_{j \in \mathrm{~J}}$ be two families of A-modules. Then

$$
\operatorname{Hom}\left(\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}, \prod_{j \in \mathrm{~J}} \mathrm{~F}_{j}\right)=\prod_{(i, j) \in \mathrm{I} \times \mathrm{J}} \operatorname{Hom}\left(\mathrm{E}_{i}, \mathrm{~F}_{j}\right) .
$$

Therefore, when I and J are finite, a homomorphism from $\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$ to $\bigoplus_{j \in \mathrm{~J}} \mathrm{~F}_{j}$ is represented by a matrix $\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ where, for all $(i, j)$ the elements $a_{i}^{j}$ are homomorphisms from $\mathrm{E}_{i}$ to $\mathrm{F}_{j}$.

Remark It may happen that $\operatorname{Hom}(\mathrm{E}, \mathrm{F})=0$ even when $\mathrm{E} \neq 0$ and $\mathrm{F} \neq 0$. For example, this is the case if A is integral and E is a torsion module, and F is torsion-free. Other examples:

$$
\operatorname{Hom}_{\mathbb{Z}}(\mathbb{Z} /(p), \mathbb{Z} /(q))=0
$$

if $p$ and $q$ are relatively prime and $\operatorname{Hom}_{\mathbb{Z}}(\mathbb{Q}, \mathbb{Z})=0$.

### 3.3.16 Left Exactness of the Functor Hom

Proposition (a) Let E be an $\mathrm{A}-m o d u l e ~ a n d ~ 0 \rightarrow \mathrm{~F}^{\prime} \xrightarrow{u} \mathrm{~F} \xrightarrow{v} \mathrm{~F}^{\prime \prime}$ an exact sequence of A-modules. Then there is an exact sequence

$$
0 \rightarrow \operatorname{Hom}\left(\mathrm{E}, \mathrm{~F}^{\prime}\right) \xrightarrow{u_{*}} \operatorname{Hom}(\mathrm{E}, \mathrm{~F}) \xrightarrow{v_{*}} \operatorname{Hom}\left(\mathrm{E}, \mathrm{~F}^{\prime \prime}\right) .
$$

(b) Let F be an $\mathrm{A}-m o d u l e$ and $\mathrm{E}^{\prime} \xrightarrow{u} \mathrm{E} \xrightarrow{v} \mathrm{E}^{\prime \prime} \rightarrow 0$ an exact sequence of $\mathrm{A}-m o d u l e s$. Then there is an exact sequence

$$
0 \rightarrow \operatorname{Hom}\left(\mathrm{E}^{\prime \prime}, \mathrm{F}\right) \xrightarrow{v_{*}} \operatorname{Hom}(\mathrm{E}, \mathrm{~F}) \xrightarrow{u_{*}} \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}\right) .
$$

Proof (a) Identifying $\mathrm{F}^{\prime}$ with $\operatorname{Ker} v$, $\operatorname{Ker} v_{*}$ is the set of $f: \mathrm{E} \rightarrow \mathrm{F}$, where $v \circ f=0$, i.e. $f(\mathrm{E}) \subset \operatorname{Ker} v$, and so $\operatorname{Ker} v_{*}=\operatorname{Hom}\left(\mathrm{E}, \mathrm{F}^{\prime}\right)$.
(b) Identifying $\mathrm{E}^{\prime \prime}$ with Coker $u, \operatorname{Ker} u^{*}$ is the set of $f: \mathrm{E} \rightarrow \mathrm{F}$ such that $f \circ u=$ 0 , i.e. $\left.f\right|_{\operatorname{Im} u}=0$, and so $\operatorname{Ker} u^{*}$ is $\operatorname{Hom}(\operatorname{Coker} u, \mathrm{~F})=\operatorname{Hom}\left(\mathrm{E}^{\prime \prime}, \mathrm{F}\right)$.

### 3.3.17 Multilinear Maps

Let $\mathrm{E}_{1}, \ldots, \mathrm{E}_{n}$ and F be A-modules. A map $f$ from $\mathrm{E}_{1} \times \cdots \times \mathrm{E}_{n}$ to F is said to be $n$-linear if, for all $i \in\{1, \ldots, n\}$ and $x_{1}, \ldots, x_{i-1}, x_{i+1}, \ldots x_{n}$, the partial map $x_{i} \mapsto f\left(x_{1}, \ldots, x_{n}\right)$ from $\mathrm{E}_{i}$ to F is linear. The module of $n$-linear maps from $\mathrm{E}_{1} \times \cdots \times \mathrm{E}_{n}$ to F can be identified with $\operatorname{Hom}\left(\mathrm{E}_{1}, \operatorname{Hom}\left(\mathrm{E}_{2}, \ldots\left(\operatorname{Hom}\left(\mathrm{E}_{n}, \mathrm{~F}\right)\right) \ldots\right)\right.$.

A $n$-linear map $f$ from $\mathrm{E}_{n}$ to F is said to be symmetric if $f\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right)=$ $f\left(x_{1}, \ldots, x_{n}\right)$ for all permutations $\sigma \in \mathfrak{S}_{n}$. It is alternating if $x_{i}=x_{j}$ with $i \neq j$ implies $f\left(x_{1}, \ldots, x_{n}\right)=0$.
If the $n$-linear map is alternating, then $f\left(x_{2}, x_{1}, x_{3}, \ldots, x_{n}\right)=-f\left(x_{1}, x_{2}, \ldots, x_{n}\right)$; more generally the sign of $f\left(x_{1}, \ldots, x_{n}\right)$ is reversed if any two $x_{i}$ are interchanged.

### 3.3.18 Signature of a Permutation

Let X be a finite set, and Y the set of 2-element subsets of X . The binary choice function on X is a map $\tau: \mathrm{Y} \rightarrow \mathrm{X}$ such that $\tau(y) \in y$ for $y \in \mathrm{Y}$. As X is finite, the axiom of choice is not needed to ensure the existence of such a choice function.

Let $\sigma$ be a permutation of X, i.e. a bijection from X to itself. Write $\sigma_{*}$ for the map $\left\{x, x^{\prime}\right\} \mapsto\left\{\sigma(x), \sigma\left(x^{\prime}\right)\right\}$ from Y to itself; and for every binary choice function $\tau$, define $\sigma_{*} \tau$ by $\sigma_{*} \tau=\sigma \circ \tau \circ \sigma^{-1}$. If $\tau$ and $\tau^{\prime}$ are two binary choice functions, write $\tau^{\prime}-\tau$ for the function $\mathrm{Y} \rightarrow \mathbb{Z} /(2)$ with value 0 at $y$ if $\tau^{\prime}(y)=\tau(y)$ and 1 otherwise. If $f$ is a function $\mathrm{Y} \rightarrow \mathbb{Z} /(2)$, write $\Sigma(f)$ for $\sum_{y \in \mathrm{Y}} f(y)$ and set $\sigma_{*} f=f \circ \sigma_{*}^{-1}$. For functions with values in $\mathbb{Z} /(2)$, the signs + and - are synonymous.

Let $\sigma$ be a permutation of X and $\tau: \mathrm{Y} \rightarrow \mathrm{X}$ abinary choice function. Set $i(\sigma, \tau)=$ $\Sigma\left(\sigma_{*} \tau-\tau\right)$. In fact, $i(\sigma, \tau)$ does not depend on the choice of $\tau$. Indeed, if $\tau^{\prime}$ is another binary choice function, then $\left(\sigma_{*} \tau^{\prime}-\tau^{\prime}\right)-\left(\sigma_{*} \tau-\tau\right)=\sigma_{*} h-h$, where $h=\tau^{\prime}-\tau$ and $\Sigma\left(\sigma_{*} h\right)=\Sigma(h)$. Therefore, $(-1)^{i(\sigma, \tau)}$ does not depend on $\tau$, and is called the signature of the permutation $\sigma$. It is written $\varepsilon(\sigma)$.

If $\sigma$ and $\sigma^{\prime}$ are two permutations of X , then $\varepsilon\left(\sigma \sigma^{\prime}\right)=\varepsilon(\sigma) \cdot \varepsilon\left(\sigma^{\prime}\right)$. In other words, $\varepsilon$ is a homomorphism from the group $\mathfrak{S}(X)$ of permutations of $X$ to the multiplicative group $\{+1,-1\}$.

Let X and $\mathrm{X}^{\prime}$ be two finite sets and $f: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ a bijection. Let $f_{*}$ be the isomorphism $\sigma \mapsto f \circ \sigma \circ f^{-1}$ from $\mathfrak{S}(\mathrm{X})$ onto $\mathfrak{S}\left(\mathrm{X}^{\prime}\right)$. Then $\varepsilon\left(f_{*}(\sigma)\right)=\varepsilon(\sigma)$. More generally, let $f: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ be an injection, with finite $\mathrm{X}^{\prime}$. Define an injective homomorphism $f_{*}$ from $\mathfrak{S}(\mathrm{X})$ to $\mathfrak{S}\left(\mathrm{X}^{\prime}\right)$ by setting $f_{*}(\sigma)$ to be the map equal to $f_{*}(\sigma)$ on $f(\mathrm{X})$ as defined previously, and to the identity on $\mathrm{X}^{\prime}-f(\mathrm{X})$. Also $\varepsilon\left(f_{*}(\sigma)\right)=\varepsilon(\sigma)$.

A transposition of X is a permutation interchanging two elements of X and keeping the others unchanged. If $\theta$ is a transposition, then $\varepsilon(\theta)=-1$. The transpositions of X generate $\mathfrak{S}(\mathrm{X})$, but if $\sigma$ is a product of $k$ transpositions, the number $k$ is necessarily even if $\varepsilon(\sigma)=1$ and odd if $\varepsilon(\sigma)=-1$.

Let E and F be A -modules and $f: \mathrm{E}^{n} \rightarrow \mathrm{~F}$ a $n$-linear alternating map. For any permutation $\sigma$ of $\{1, \ldots, n\}, f\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right)=\varepsilon(\sigma) f\left(x_{1}, \ldots, x_{n}\right)$.

### 3.3.19 Algebras

An A-algebra is an A-module E with a bilinear map $\mu: \mathrm{E} \times \mathrm{E} \rightarrow \mathrm{E}$, called multiplication. It is said to be an associative (resp. commutative, resp. unital) if the composition law $\mu$ on E is associative (resp. commutative, resp. has an identity element).

Examples (1) If B is a ring and $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ a ring homomorphism, B , together with its A-module structure (3.3.1, Example 3) and multiplication, is a unital associative commutative A-algebra.

More generally, if B is a not necessarily commutative ring and $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ a ring homomorphism such that $\varphi(\mathrm{A})$ is contained in the centre of B , then B is a unital associative A-algebra. All unital associative A-algebras are uniquely defined in this way.
(2) Let M be an A-module. The module $\operatorname{End}(\mathrm{M})=\operatorname{Hom}(\mathrm{M}, \mathrm{M})$ with composition $(f, g) \mapsto f \circ g$, is a unital associative algebra, and is not generally commutative.

## Exercises 3.3. (Modules)

1.-(a) Show that there are no minimal generating families in the $\mathbb{Z}$-module $\mathbb{Q}$.
(b) Give an example of a module having two minimal generating families of distinct finite cardinalities.
(c) Show that, if there is a minimal generating infinite family, then all generating families are infinite (see proof of 3.4.5).
(d) For all A-modules E, set $g(\mathrm{E})$ to be the minimal cardinality of a generating family of E . Show that if $0 \rightarrow \mathrm{E} \rightarrow \mathrm{F} \rightarrow \mathrm{G} \rightarrow 0$ is a short exact sequence, then $g(\mathrm{~F}) \leqslant g(\mathrm{E})+g(\mathrm{G})$. Give an example where the inequality is strict.
2.-Let $K$ be a field; set $A=K[X, Y], E=A^{2}$ and let $F$ be a submodule de $E$ generated by (X, Y). Show that:
(a) $\mathrm{E} / \mathrm{F}$ is torsion-free;
(b) There are no retractions from E onto F . Deduce that F is not a direct factor of $E$.
3.-Let E be an A-module and $\mathscr{F}$ an ideal in A . Let $\mathscr{J} \cdot \mathrm{E}$ be the submodule of E consisting of all elements $\sum_{i \in \mathrm{I}} a_{i} x_{i}$, where I is a finite set, $a_{i} \in \mathscr{J}$ and $x_{i} \in \mathrm{E}$.
(a) Let F and G be two submodules of E . Does the following necessarily hold:

$$
\begin{aligned}
& \mathscr{J} \cdot(\mathrm{F} \cap \mathrm{G})=\mathscr{J} \cdot \mathrm{F} \cap \mathscr{J} \cdot \mathrm{G} ? \\
& \mathscr{J} \cdot(\mathrm{~F}+\mathrm{G})=\mathscr{J} \cdot \mathrm{F}+\mathscr{J} \cdot \mathrm{G} ?
\end{aligned}
$$

(b) Let $\left(\mathrm{E}_{\lambda}\right)$ be a family of A-modules. Show that

$$
\mathscr{J} \cdot \bigoplus \mathrm{E}_{\lambda}=\bigoplus \mathscr{J} \cdot \mathrm{E}_{\lambda}
$$

(c) Show that $\mathscr{J} \cdot \prod \mathrm{E}_{\lambda} \subset \prod \mathscr{J} \cdot \mathrm{E}_{\lambda}$, and that the equality holds if $\mathscr{J}$ is finitely generated. Show that, if $\mathscr{F}$ is not finitely generated, then there is a family $\left(\mathrm{E}_{\lambda}\right)$ such that the inclusion is strict.
4.-Let E be an A-module, F and G two submodules of E . Construct an exact sequence $0 \rightarrow \mathrm{~F} \cap \mathrm{G} \rightarrow \mathrm{F} \oplus \mathrm{G} \rightarrow \mathrm{F}+\mathrm{G} \rightarrow 0$.
5.-Let E be an A-module, F and G -submodules of E such that $\mathrm{G} \subset \mathrm{F}$. Show that, if $G$ is a direct factor of $E$, the $G$ is a direct factor of $F$. Show that, if $F$ is a direct factor of E , then $\mathrm{F} / \mathrm{G}$ is a direct factor of $\mathrm{E} / \mathrm{G}$.
6.-Let E and F be two A-modules, and $u: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism.
(a) Show that, $\operatorname{Ker} u$ and $\operatorname{Im} u$ are direct factors of E and F respectively if and only if there is a homomorphism $v: \mathrm{F} \rightarrow \mathrm{E}$ such that $u \circ v \circ u=u$.
(b) Show that $u \circ v \circ u=u$ does not imply $v \circ u \circ v=v$, but that there exists $v$ such that $u \circ v \circ u=u$ and $w$ such that $u \circ w \circ u=u$ and $w \circ u \circ w=w$.
7. (Chart changes in Grassmannians)-(a) Let F and G be two A-modules. Show that the graphs of linear maps from $F$ to $G$ are the complements of $0 \times G$ in $F \times G$.
(b) Let E be an A-module, $\mathrm{G}_{0}$ a submodule of $\mathrm{E}, \mathrm{F}_{0}$ and $\mathrm{F}_{1}$ two complements of $\mathrm{G}_{0}$. Every complement F of $\mathrm{G}_{0}$ can be considered the graph of a map $f_{0}: \mathrm{F}_{0} \rightarrow \mathrm{G}_{0}$ or the graph of a map $f_{1}: \mathrm{F}_{1} \rightarrow \mathrm{G}_{0}$. Give the relation between $f_{0}$ and $f_{1}$.
(c) Let E be an A-module, $\mathrm{F}_{0}$ a submodule of $\mathrm{E}, \mathrm{G}_{0}$ and $\mathrm{G}_{1}$ two complements of $\mathrm{F}_{0}$. A complement F of $\mathrm{G}_{0}$ can be considered the graph of a map $f_{0}: \mathrm{F}_{0} \rightarrow \mathrm{G}_{0}$. What condition should $f_{0}$ satisfy for F to be the complement of $\mathrm{G}_{1}$. Then F can be considered the graph of a map $f_{1}: \mathrm{F}_{0} \rightarrow \mathrm{G}_{1}$. Express $f_{1}$ in terms of $f_{0}$.
8.-Let $\mathrm{E}, \mathrm{F}$, G be three modules, $i_{1}: \mathrm{E} \rightarrow \mathrm{F}$ and $i_{2}: \mathrm{E} \rightarrow \mathrm{G}$ injective homomorphisms. Embed E in $\mathrm{F} \oplus \mathrm{G}$ by $\left(i_{1}, i_{2}\right)$.
(a) Show that the following are exact sequences:

$$
0 \rightarrow \mathrm{~F} \rightarrow(\mathrm{~F} \oplus \mathrm{G}) / \mathrm{E} \rightarrow \mathrm{G} / \mathrm{E} \rightarrow 0, \quad 0 \rightarrow \mathrm{G} \rightarrow(\mathrm{~F} \oplus \mathrm{G}) / \mathrm{E} \rightarrow \mathrm{~F} / \mathrm{E} \rightarrow 0
$$

(b) Show that the first sequence is split if and only if there is a homomorphism from $G$ to $F$ inducing the identity of E . Give an example where one of the sequences is split and the other one is not.
9.-For every prime $p$, let $\mathrm{M}_{p}$ be the subgroup in $\mathbb{Q} / \mathbb{Z}$ of elements whose orders are $p$-powers.
(a) Show that $\mathbb{Q} / \mathbb{Z}$ is the direct sum of the subgroups $M_{p}$ (see 3.1, Exercise 6, b).
(b) Show that none of the $\mathrm{M}_{p}$ is the direct sum of two nonzero submodules.
(c) Show that $\operatorname{End}(\mathbb{Q} / \mathbb{Z})$ can be identified with the product of rings $\operatorname{End}\left(\mathbf{M}_{p}\right)$.
(d) Show that $\operatorname{End}\left(\mathbf{M}_{p}\right)$ can be identified with the ring $\widehat{\mathbb{Z}}_{p}$ of $p$-adic integers.
(e) Show that $\operatorname{End}(\mathbb{Q} / \mathbb{Z})$ can be identified with the profinite completion $\widehat{\mathbb{Z}}$ of $\mathbb{Z}$. We thus recover the result of (2.5, Exercise 7, e).
10.-(This exercise uses the results of the two preceding ones.) Consider the exact sequence:

$$
\begin{equation*}
0 \rightarrow \mathbb{Z} \xrightarrow{\iota} \mathbb{Q} \xrightarrow{\chi} \mathbb{Q} / \mathbb{Z} \rightarrow 0 \tag{E}
\end{equation*}
$$

where $\iota$ and $\chi$ are the canonical maps. Applying the functor $\mathrm{M} \mapsto \operatorname{Hom}(\mathbb{Q}, \mathrm{M})$ to (E) gives an exact sequence:

$$
0 \rightarrow \operatorname{Hom}(\mathbb{Q}, \mathbb{Z}) \xrightarrow{\iota_{*}} \operatorname{Hom}(\mathbb{Q}, \mathbb{Q}) \xrightarrow{\chi_{*}} \operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z}) .
$$

The aim is to describe $\operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z})$ and Coker $\chi_{*}$.
(a) Determine $\operatorname{Hom}(\mathbb{Q}, \mathbb{Z})$ and $\operatorname{Hom}(\mathbb{Q}, \mathbb{Q})$.
(b) Show that applying the functor $\mathrm{M} \mapsto \operatorname{Hom}(\mathrm{M}, \mathbb{Q} / \mathbb{Z})$ to $(\mathrm{E})$ gives a non-split exact sequence:

$$
0 \rightarrow \widehat{\mathbb{Z}} \rightarrow \operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z}) \rightarrow \mathbb{Q} / \mathbb{Z} \rightarrow 0
$$

(c) Define a homomorphism $w: \mathbb{Q} \oplus \widehat{\mathbb{Z}} \rightarrow \operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z})$ by setting $w(a, 0)=$ $\chi \circ a \cdot 1_{\mathbb{Q}}$ and $w(0, u)=u \circ \chi$, where $u$ is considered an endomorphism of $\mathbb{Q} / \mathbb{Z}$. Determine Ker $w$ and Coker $w$.
(d) Show that there is an exact sequence:

$$
0 \rightarrow \mathbb{Q} \rightarrow \operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z}) \rightarrow \widehat{\mathbb{Z}} / \mathbb{Z} \rightarrow 0
$$

(Use Exercise 8.) Determine Coker $\chi_{*}$. (The above exact sequence can be shown to be split.)
(e) Show that $\mathbf{M}_{p}$ can be identified with $\mathbb{Z}[1 / p] / \mathbb{Z}$, or to $\widehat{\mathbb{Q}}_{p} / \widehat{\mathbb{Z}}_{p}$, where $\mathbb{Z}[1 / p]$ denotes the ring of fractions whose denominator is a power of $p$, and $\widehat{\mathbb{Q}}_{p}$ the $p$-adic field, the field of fractions of $\widehat{\mathbb{Z}}_{p}$. Define $w_{p}: \mathbb{Z}[1 / p] \oplus \widehat{\mathbb{Z}}_{p} \rightarrow \operatorname{Hom}\left(\mathbb{Q}, \mathbf{M}_{p}\right)$ as in (c). Determine $\operatorname{Ker} w_{p}$ and Coker $w_{p}$. Deduce that $\operatorname{Hom}\left(\mathbb{Q}, \mathbf{M}_{p}\right)$ can be identified with $\widehat{\mathbb{Q}}_{p}$.
(f) Show that $\operatorname{Hom}(\mathbb{Q}, \mathbb{Q} / \mathbb{Z})$ can be identified with the set of $x=\left(x_{p}\right) \in \prod \widehat{\mathbb{Q}}_{p}$, where all but finitely many $x_{p}$ are in $\widehat{\mathbb{Z}}_{p}$.
11.-Let A be a ring. Show that, for all ideals $\mathscr{F}$ of A, the category $(\mathrm{A} / \mathscr{F})-\mathscr{M} \mathrm{od}$ is equivalent to a complete subcategory of A- $M \mathscr{O}$. Show that the same holds for $\left(S^{-1} A\right)-\mathscr{M} \sigma d$ for any multiplicatively stable subset $S$ of $A$, setting $S^{-1} A$ to be the ring of fractions with numerator in A and denominator in S .
12.-(a) Show that the set of endomorphisms of the identity functor ofA- $\mathscr{M} \boldsymbol{O}$ can be identified with A. Give an interpretation for addition (see 2.6, Exercise 2) and multiplication in A .
(b) Deduce that if A and B are two rings such that the categories $\mathrm{A}-\mathscr{M} \sigma$ and $B-\mathscr{M} O d$ are equivalent, then $A$ and $B$ are isomorphic.
(c) the previous result ((b)) does not hold for non-commutative rings: let A be the ring of $2 \times 2$ matrices with entries in $\mathbb{R}$. For any vector space $\mathrm{E}, \mathrm{A}$ can be made to act on $E^{2}$ by

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\binom{x}{y}=\binom{a x+b y}{c x+d y}
$$

Show that the functors $\mathrm{E} \mapsto \mathrm{E}^{2}$ from $\mathbb{R}-\mathscr{M} \sigma d$ to $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$ and $\mathrm{F} \mapsto \operatorname{Hom}_{\mathrm{A}}\left(\mathbb{R}^{2}, \mathrm{~F}\right)$ from A- $\mathscr{M} \sigma d$ to $\mathbb{R}-\mathscr{M} \sigma d$ are equivalences of categories (note that, if A is not commutative, $\operatorname{Hom}_{\mathrm{A}}(\mathrm{F}, \mathrm{G})$ is not a natural A-module for A-modules F and G).
13.-Let A be an integral domain where $2 \neq 0, \mathrm{E}$ and F torsion-free $\mathrm{A}-$ modules and $f: \mathrm{E}^{n} \rightarrow \mathrm{~F}$ a $n$-linear map. Show that, if $f\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right)=\varepsilon(\sigma) f\left(x_{1}, \ldots, x_{n}\right)$ for all sequences $\left(x_{1}, \ldots, x_{n}\right)$ in E and all permutations $\sigma$, then the map $f$ is alternating.

Taking for A a field of characteristic 2, construct a counterexample.

### 3.4 Free Modules, Matrices

Throughout this section, A denotes a (commutative) ring.

### 3.4.1 Free Modules

Let E be an A -module and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of E . The family $\left(x_{i}\right)_{i \in \mathrm{I}}$ is a free family if for any family $\left(\lambda_{i}\right)_{i \in \mathrm{I}}$ of elements of A with finite support, $\sum_{i \in \mathrm{I}} \lambda_{i} x_{i}=0$ implies $\lambda_{i}=0$ for all $i \in \mathrm{I}$.

A basis of E is a free generating family of E , and E is a free A -module, if it has a basis.

It is free of rank $n$ (or dimension $n$ ) the basis contains $n$ elements.
Every free module over an integral domain is torsion-free. The converse does not hold. However, as will be seen, every finitely generated torsion-free module over a PID is free (3.5.8. Corollary 3.2).

### 3.4.2 Examples

The module $\mathrm{A}^{r}$ is a free A-module. The ring $\mathrm{A}[\mathrm{X}]$ is a free A -module, ${ }^{9}$ the rings $\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ are free A-modules. ${ }^{10}$ If K is a field, every finite dimensional vector space is a free K -module ${ }^{11}$; more generally:
Proposition Every vector space is free.
Proof Let E be a vector space. The set $\mathscr{L}$ of free subsets of E, ordered by inclusion, is direct since it is of finite character (1.4.2, Example 3).

By Zorn's theorem, there is a maximal free subset, B say, and let M be the vector subspace of E generated by B . We prove by contradiction that $\mathrm{M}=\mathrm{E}$. Suppose that

[^8]$\mathrm{M} \neq \mathrm{E}$ and let $y \in \mathrm{E}-\mathrm{M}$. It follows that $\mathrm{B} \cup\{y\}$ is free, which contradicts the maximality of $\mathbf{B}$. Indeed, suppose that $\sum \lambda_{i} x_{i}+\mu y=0$ for $x_{i} \in \mathbf{B}$.

If $\mu=0, \sum \lambda_{i} x_{i}=0$ implies $\lambda_{i}=0$ for all $i$.
If $\mu \neq 0$, then $y=-\sum \frac{\lambda_{i}}{\mu} x_{i}$ and $y$ is in M , which is impossible since $y \in \mathrm{E}-\mathrm{M}$.

### 3.4.3 Counterexamples

As a $\mathbb{Z}$-module, $\mathbb{Q}$ is torsion-free but not free.
If A is an integral domain, a non-principal ideal of A is not a free module, despite being torsion-free as well as the submodule of a free module. As will be seen, every submodule of a free module over a PID is free (3.5.1), and that every finitely generated torsion-free module over an integral domain is isomorphic to a submodule of a free module (3.8.14).

### 3.4.4 Alternating n-Linear Forms

An alternating $n$-linear form on an A-module E is an alternating $n$-linear map from $\mathrm{E}^{n}$ to A .

If E has a generating family $\left(e_{1}, \ldots, e_{p}\right)$ with $p<n$, every alternating $n$-linear form $f$ on E is null. Indeed, if $x_{i}=\sum_{1 \leqslant j \leqslant p} x_{i}^{j} e_{j}$, then $f\left(x_{1}, \ldots, x_{n}\right)=\sum_{j_{1}, \ldots, j_{n}} x_{1}^{j_{1}} \ldots x_{n}^{j_{n}}$ $f\left(e_{j_{1}}, \ldots, e_{j_{n}}\right)=0$ because for each term there exists $k$ and $k^{\prime}$ with $k \neq k^{\prime}$ and $j_{k}=$ $j_{k^{\prime}}$.

Proposition Let E be a free A-module of rank $n$. Then the module $\Lambda$ of alternating $n$-linear forms on E is free of rank 1 . More precisely, if $\left(e_{1}, \ldots, e_{n}\right)$ is a basis for E , there is a unique form $\Delta \in \Lambda$ such that $\Delta\left(e_{1}, \ldots, e_{n}\right)=1$, and $\Delta$ is a basis of $\Lambda$.

Proof Let $\left(e_{1}, \ldots, e_{n}\right)$ be a basis for E. If $x_{i}=\sum x_{i}^{j} e_{j}$, then

$$
\begin{equation*}
f\left(x_{1}, \ldots, x_{n}\right)=\sum_{\sigma \in \mathfrak{S}_{n}} \varepsilon(\sigma) x_{1}^{\sigma(1)} \ldots x_{n}^{\sigma(n)} f\left(e_{1}, \ldots, e_{n}\right) \tag{3.1}
\end{equation*}
$$

for any alternating $n$-linear form $f$.
The map $\varphi: f \mapsto f\left(e_{1}, \ldots, e_{n}\right)$ from $\Lambda$ to A is linear. By (1), $f\left(e_{1}, \ldots, e_{n}\right)=$ $0 \Rightarrow f=0$, and so $\varphi$ is injective.

Define $\Delta: \mathrm{E}^{n} \rightarrow \mathrm{~A}$ by

$$
\Delta\left(x_{1}, \ldots, x_{n}\right)=\sum_{\sigma \in \mathfrak{S}_{n}} \varepsilon(\sigma) x_{1}^{\sigma(1)} \ldots x_{n}^{\sigma(n)}
$$

The map $\Delta$ is an alternating $n$-linear form on E , and $\Delta\left(e_{1}, \ldots, e_{n}\right)=1$, i.e. $\varphi(\Delta)=1$. Then $\varphi$ is surjective, and $\Lambda$ is isomorphic to A.

### 3.4.5 Uniqueness of the Dimension

Proposition and Definition Let A be a nonzero ring and E a free A-module. All bases for E have the same cardinality, called the dimension (or rank) of E .

Proof Let $\left(e_{i}\right)_{i \in \mathrm{I}}$ and $\left(f_{j}\right)_{j \in \mathrm{~J}}$ be two bases for E. For every $j \in \mathrm{~J}$, there is a family $\left(\lambda_{i j}\right)_{i \in \mathrm{I}}$ of elements of A such that $f_{j}=\sum_{i \in \mathrm{I}} \lambda_{i j} e_{i}$. Let $\mathrm{S}_{j}$ be the support of $\left(\lambda_{i j}\right)_{i \in \mathrm{I}}$. For all $j \in \mathrm{~J}$, the support $\mathrm{S}_{j}$ is a finite subset of I . Set $\mathrm{I}^{\prime}=\cup_{j \in \mathrm{~J}} \mathrm{~S}_{j}$. All $f_{j}$ are in the submodule $\mathrm{E}^{\prime}$ of E generated by $\left(e_{i}\right)_{i \in \mathrm{I}^{\prime}}$, and so $\mathrm{E}^{\prime}=\mathrm{E}$, whence $\mathrm{I}^{\prime}=\mathrm{I}$, and

$$
\operatorname{Card} \mathrm{I}=\operatorname{Card} \mathrm{I}^{\prime} \leqslant \sum_{j \in \mathrm{~J}} \operatorname{Card} \mathrm{~S}_{j} \leqslant \operatorname{Card} \mathrm{~J} \times \aleph_{0}
$$

If the set J is finite, then so is I . Let $p$ and $q$ be the number of elements of I and J respectively. There is an alternating $p$-linear form of E vanishing everywhere. Since E has a basis with $q$ elements, $p \leqslant q$ because every alternating $r$-linear form on E where $r>q$ vanishes everywhere.

If the set $\mathbf{J}$ is infinite, then Card $\mathbf{J}=\operatorname{Card} \mathbf{J} \times \aleph_{0} \geqslant$ Card I. Hence Card $\mathrm{I} \leqslant \operatorname{Card} \mathbf{J}$ in all cases. Interchanging the roles of I and J gives $\operatorname{Card} \mathrm{I}=\operatorname{Card} \mathrm{J}$,

Remark This result does not hold for non-commutative rings (3.4, Exercise 2), but does for skew fields. ${ }^{12}$

### 3.4.6 Free Module on a Set

Let $A$ be a ring, I a set. Let $A^{I}$ be the set of families $\left(\lambda_{i}\right)_{i \in I}$ of elements in A indexed by I. Define an A-module structure on $\mathrm{A}^{\mathrm{I}}$ by setting

$$
\begin{gathered}
\left(\lambda_{i}\right)_{i \in \mathrm{I}}+\left(\mu_{i}\right)_{i \in \mathrm{I}}=\left(\lambda_{i}+\mu_{i}\right)_{i \in \mathrm{I}} \\
a\left(\lambda_{i}\right)_{i \in \mathrm{I}}=\left(a \lambda_{i}\right)_{i \in \mathrm{I}} .
\end{gathered}
$$

Write $\mathrm{A}^{(\mathrm{I})}$ for the A -submodule of $\mathrm{A}^{\mathrm{I}}$ consisting of families with finite support. For all $i \in \mathrm{I}$, let $e_{i}$ be the element $\left(\delta_{i j}\right)_{j \in \mathrm{I}}$ of $\mathrm{A}^{(\mathrm{I})}$. The family $\left(e_{i}\right)_{i \in \mathrm{I}}$ is a basis for $\mathrm{A}^{(\mathrm{I})}$. Indeed, for all $\lambda=\left(\lambda_{i}\right)_{i \in \mathrm{I}}$ in $\mathrm{A}^{(\mathrm{I})}, \lambda=\sum \lambda_{i} e_{i}$, and so $\lambda$ is in the submodule of $\mathrm{A}^{(\mathrm{I})}$ generated by the elements $\left(e_{i}\right)$, and the relation $\sum \lambda_{i} e_{i}=0$ implies $\left(\lambda_{i}\right)_{i \in \mathrm{I}}=0$.

[^9]Hence $\lambda_{i}=0$ for all $i \in \mathrm{I} . \mathrm{A}^{(\mathrm{I})}$ is said to be the free A-module on I and $\left(e_{i}\right)_{i \in \mathrm{I}}$ the canonical basis for $\mathrm{A}^{(\mathrm{I})}$.
UnIVERSAL PROPERTY. Let E be an arbitrary A-module and $\left(x_{i}\right)_{i \in \mathrm{I}}$ an arbitrary family of elements of E . There is a unique homomorphism $f$ from $\mathrm{A}^{(\mathrm{I})}$ to E such that $f\left(e_{i}\right)=x_{i}$ for all $i \in \mathrm{I}$.

Indeed, if $f$ is such a homomorphism, then for all $\lambda=\left(\lambda_{i}\right)_{i \in \mathrm{I}}$ of $\mathrm{A}^{(\mathrm{I})}$,

$$
f(\lambda)=f\left(\sum \lambda_{i} e_{i}\right)=\sum \lambda_{i} f\left(e_{i}\right)=\sum \lambda_{i} x_{i}
$$

and uniqueness follows. The map $f$ defined by $f(\lambda)=\sum \lambda_{i} x_{i}$ for $\lambda=\left(\lambda_{i}\right)_{i \in \mathrm{I}}$ of $\mathrm{A}^{(\mathrm{I})}$ satisfies the desired condition.

Remarks (1) The universal property can also be stated as follows: the map $f \mapsto\left(f\left(e_{i}\right)\right)_{i \in \mathrm{I}}$ is a bijection from $\operatorname{Hom}\left(\mathrm{A}^{(\mathrm{I})}, \mathrm{E}\right)$ onto $\mathrm{E}^{\mathrm{I}}$.
(2) These definitions imply that the image of $f$ is the submodule of E generated by the family $\left(x_{i}\right)_{i \in \mathrm{I}}$. Then $f$ is surjective (resp. injective, resp. bijective) if and only if $\left(x_{i}\right)_{i \in \mathrm{I}}$ is a generating family (resp. free, resp. a basis).

An A-module E is free if and only if it is isomorphic to a module of the form $\mathrm{A}^{(\mathrm{I})}$. Every module is isomorphic to a quotient of a free module.

### 3.4.7 Projectivity of Free Modules

Proposition Let E be a free $\mathrm{A}-m o d u l e ~ a n d ~ 0 \rightarrow \mathrm{~F}^{\prime} \xrightarrow{u} \mathrm{~F} \xrightarrow{v} \mathrm{~F}^{\prime \prime} \rightarrow 0$ an exact sequence of A-modules. Then

$$
\begin{equation*}
0 \rightarrow \operatorname{Hom}\left(\mathrm{E}, \mathrm{~F}^{\prime}\right) \xrightarrow{u_{*}} \operatorname{Hom}(\mathrm{E}, \mathrm{~F}) \xrightarrow{v_{*}} \operatorname{Hom}\left(\mathrm{E}, \mathrm{~F}^{\prime \prime}\right) \rightarrow 0 \tag{3.2}
\end{equation*}
$$

is an exact sequence.
Proof We may assume that $\mathrm{E}=\mathrm{A}^{(\mathrm{I})}$. Then (1) becomes: $0 \rightarrow \mathrm{~F}^{\mathrm{I}} \rightarrow \mathrm{F}^{\mathrm{I}} \rightarrow \mathrm{F}^{\prime \prime \mathrm{I}} \rightarrow 0$ and $\mathrm{F}^{\mathrm{I}}=\operatorname{Ker}\left(v_{*}: \mathrm{F}^{\mathrm{I}} \rightarrow \mathrm{F}^{\prime / \mathrm{I}}\right)$ is immediate; the surjectivity of $v_{*}: \mathrm{F}^{\mathrm{I}} \rightarrow \mathrm{F}^{\prime / \mathrm{I}}$ follows from the axiom of choice.
Corollary If, in a short exact sequence $0 \rightarrow \mathrm{E}^{\prime} \xrightarrow{u} \mathrm{E} \xrightarrow{v} \mathrm{E}^{\prime \prime} \rightarrow 0$ of A-modules, $\mathrm{E}^{\prime \prime}$ is a free A-module, the exact sequence splits.
Proof The homomorphism $v_{*}: \operatorname{Hom}\left(\mathrm{E}^{\prime \prime}, \mathrm{E}\right) \rightarrow \operatorname{Hom}\left(\mathrm{E}^{\prime \prime}, \mathrm{E}^{\prime \prime}\right)$ is surjective, and so there exists $\sigma \in \operatorname{Hom}\left(\mathrm{E}^{\prime \prime}, \mathrm{E}\right)$ such that $v \circ \sigma=1_{\mathrm{E}^{\prime \prime}}$.

### 3.4.8 Projective Modules

Definition An A-module P is said to be projective if, for all surjective A-module homomorphisms $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime}$ and homomorphisms $g: \mathrm{P} \rightarrow \mathrm{M}^{\prime}$, there exists a
homomorphism $h: \mathrm{P} \rightarrow \mathrm{M}$ making the diagram

commutative.
In other words, P is projective if and only if the functor $\mathrm{M} \mapsto \operatorname{Hom}(\mathrm{P}, \mathrm{M})$ from the category A-Mod to itself transforms surjective homomorphisms into surjective homomorphisms.

All free modules are projective (3.4.7).
Proposition A module is projective if and only if it is the direct factor of a free module.

Proof (a) Let P and Q be two modules such that $\mathrm{L}=\mathrm{P} \oplus \mathrm{Q}$ is free, and $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime} \quad$ a surjective module homomorphism. Then $\operatorname{Hom}(\mathrm{L}, \mathrm{M})=$ $\operatorname{Hom}(\mathrm{P}, \mathrm{M}) \times \operatorname{Hom}(\mathrm{Q}, \mathrm{M})$. In the commutative diagram

the map $f_{*}$ is surjective and $\mathrm{pr}_{1}$ is surjective. So $f_{*} \circ \mathrm{pr}_{1}=\mathrm{pr}_{1} \circ f_{*}$ is surjective, and $f_{*}: \operatorname{Hom}(\mathrm{P}, \mathrm{M}) \rightarrow \operatorname{Hom}\left(\mathrm{P}, \mathrm{M}^{\prime}\right)$ is surjective.
(b) Conversely, suppose that P is a projective module. Let L be a free module and $f: \mathrm{L} \rightarrow \mathrm{P}$ a surjective homomorphism. There is a homomorphism $h: \mathrm{P} \rightarrow \mathrm{L}$ such that the diagram

commutes. Therefore (3.3.12) P can be identified with a direct factor of $L$.
For examples of non-free projective modules see 3.4, Exercise 8.

### 3.4.9 Matrices

Let A be a ring, I and J finite sets, $\left(e_{i}\right)_{i \in \mathrm{I}}$ and $\left(f_{j}\right)_{j \in \mathrm{~J}}$ canonical bases for $\mathrm{A}^{\mathrm{I}}$ and $\mathrm{A}^{\mathrm{J}}$ respectively. For any A-linear map $\varphi: \mathrm{A}^{\mathrm{I}} \rightarrow \mathrm{A}^{\mathrm{J}}$ there is a unique family $\mathrm{M}(\varphi)=$ $\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ of elements of A such that for all $i \in \mathrm{I}, \varphi\left(e_{i}\right)=\sum_{j \in \mathrm{~J}} a_{i}^{j} f_{j}$. Then $\mathrm{M}(\varphi)$ is said to be the matrix representing $\varphi$.

For $i_{0} \in \mathrm{I},\left(a_{i_{0}}^{j}\right)_{j \in \mathrm{~J}}$ is called the $i-0$-th column of $\mathrm{M}(\varphi)$; for $j_{0} \in \mathrm{~J}$ the family $\left(a_{i}^{j_{0}}\right)_{i \in \mathrm{I}}$ is called the $j_{0}$-th row.

Assigning to each A-linear map $\varphi$ its matrix $\mathrm{M}(\varphi)$ defines a bijection from $\operatorname{Hom}\left(\mathrm{A}^{\mathrm{I}}, \mathrm{A}^{\mathrm{J}}\right)$ onto $\mathrm{A}^{\mathrm{I} \times \mathrm{J}}$. Let $\alpha=\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}, \alpha^{\prime}=\left(a_{i}^{\prime j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}} \in \mathrm{A}^{\mathrm{I} \times \mathrm{J}}$, and let $\varphi$ and $\varphi^{\prime}$ be the maps represented by these matrices. By definition, set $\alpha+\alpha^{\prime}=$ $\mathrm{M}\left(\varphi+\varphi^{\prime}\right)$. Then,

$$
\begin{equation*}
\alpha+\alpha^{\prime}=\left(a_{i}^{j}+a_{i}^{\prime j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}} \tag{3.3}
\end{equation*}
$$

Let K be a finite set, $\left(g_{k}\right)_{k \in \mathrm{~K}}$ the canonical basis for $\mathrm{A}^{\mathrm{K}}, \beta=\left(b_{j}^{k}\right)_{(j, k) \in \mathrm{J} \times \mathrm{K}} \in \mathrm{A}^{\mathrm{J} \times \mathrm{K}}$, and $\psi$ the map represented by $\beta$. By definition, set

$$
\beta \alpha=\mathrm{M}(\psi \circ \varphi)
$$

Then,

$$
\begin{equation*}
\beta \alpha=\left(c_{i}^{k}\right)_{(i, k) \in \mathrm{I} \times \mathrm{K}} \text { with } c_{i}^{k}=\sum_{j \in \mathrm{~J}} b_{j}^{k} a_{i}^{j} . \tag{3.4}
\end{equation*}
$$

In particular, the set $\mathscr{M}_{\mathrm{I}}(\mathrm{A})=\mathrm{A}^{\mathrm{I} \times \mathrm{I}}=\operatorname{Hom}\left(\mathrm{A}^{\mathrm{I}}, \mathrm{A}^{\mathrm{I}}\right)$ together with the two composition laws defined above is a ring, which is not commutative if Card $\mathrm{I}>1$ and $\mathrm{A} \neq 0$.

This ring has an identity element: the matrix $\left(\delta_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ of the identity of $\mathrm{A}^{\mathrm{I}}$.

### 3.4.10 Matrices Representing a Homomorphism with Respect to Given Bases

Let E and F be finitely generated free A -modules, $x=\left(x_{i}\right)_{i \in \mathrm{I}}$ and $y=\left(y_{j}\right)_{j \in \mathrm{~J}}$ bases of E and F respectively, $\xi: \mathrm{A}^{\mathrm{I}} \rightarrow \mathrm{E}$ and $\eta: \mathrm{A}^{\mathrm{J}} \rightarrow \mathrm{F}$ isomorphisms such that for all $i \in \mathrm{I}$ and $j \in \mathrm{~J}, \xi\left(e_{i}\right)=x_{i}$ and $\eta\left(f_{j}\right)=y_{j}$. For every A-linear map $\varphi: \mathrm{E} \rightarrow \mathrm{F}$, $\mathbf{M}_{\mathscr{Y}}^{\mathscr{X}}(\varphi)=\mathrm{M}\left(\eta^{-1} \circ \varphi \circ \xi\right)$ is said to be the matrix representing $\varphi$ with respect to bases $\mathscr{X}$ and $\mathscr{Y}$. Its $i$-th column is the sequence of coordinates of $\varphi\left(x_{i}\right)$ with respect to the basis $\mathscr{Y}$. If $\mathscr{X}$ is a basis for an A-module G and $\psi: \mathrm{F} \rightarrow \mathrm{G}$ a homomorphism, then $\mathbf{M}_{\mathscr{E}}^{\mathscr{X}}(\psi \circ \varphi)=\mathrm{M}_{\mathscr{E}}^{\mathscr{Y}}(\psi) \cdot \mathrm{M}_{\mathscr{Y}}^{\mathscr{X}}(\varphi)$.

If $X$ and $X^{\prime}$ are two bases of E , the transition matrix from $\mathscr{X}$ to $X^{\prime}$ is $\mathrm{M}_{X^{\prime}}^{X^{\prime}}\left(1_{\mathrm{E}}\right)$. Its $i$-th column is the sequence of coordinates $x_{i}^{\prime}$ with respect to basis $\mathscr{X}$.

Let E and F be a finitely generated free A-modules, $\mathscr{X}$ and $\mathscr{X}^{\prime}$ bases of $\mathrm{E}, \mathscr{Y}$ and $\mathscr{Y}^{\prime}$ bases of $\mathrm{F}, \varphi$ a homomorphism from E to $\mathrm{F}, \mathrm{M}=\mathrm{M}_{\mathscr{Y}}^{x}(\varphi), \mathrm{M}^{\prime}=\mathrm{M}_{y^{\prime}}^{x^{\prime}}(\varphi)$, P the transition matrix from $\mathscr{X}$ to $X^{\prime}$ and Q the transition matrix from $\mathscr{y}$ to $\mathscr{Y}^{\prime}$. Then

$$
\mathrm{M}^{\prime}=\mathrm{Q}^{-1} \mathrm{MP}
$$

### 3.4.11 Determinants

Let E be a free A-module of rank $n$ and $\varphi$ an endomorphism of E. Let $\Lambda$ be the module of alternating $n$-linear forms on E and define the endomorphism $\varphi^{*}$ de $\Lambda$ by $\varphi^{*} f\left(x_{1}, \ldots, x_{n}\right)=f\left(\varphi\left(x_{1}\right), \ldots, \varphi\left(x_{n}\right)\right)$. It follows from Proposition 3.4.4 that there is a unique scalar $\delta \in \mathrm{A}$ such that $\varphi^{*} f=\delta f$ for all $f \in \Lambda$. It is called the determinant of $\varphi$ and is written $\delta=\operatorname{det} \varphi$.

If $\psi$ is another endomorphism of E , then $(\psi \circ \varphi)^{*}=\varphi^{*} \circ \psi^{*}$ and $\operatorname{det}(\psi \circ \varphi)=$ $\operatorname{det} \psi \cdot \operatorname{det} \varphi$.

Let I be a finite set and $\alpha=\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ a square matrix. The determinant of $\alpha$, written $\operatorname{det} \alpha$, is the determinant of the endomorphism of $\mathrm{A}^{\mathrm{I}}$ represented by $\alpha$. If $\alpha$ is the matrix representing an endomorphism $\varphi$ of a free module E with respect to a basis for E (the same at both ends), then $\operatorname{det} \varphi=\operatorname{det} \alpha$.
Proposition Let I be a finite set and $\alpha=\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ a square matrix. Then $\operatorname{det} \alpha=$ $\sum_{\sigma \in \mathfrak{S}(\mathrm{I})} \varepsilon(\sigma) \prod_{i \in \mathrm{I}} a_{i}^{\sigma(i)}$.

Proof Suppose that $\mathrm{I}=\{1, \ldots, n\}$. Let $\left(e_{1}, \ldots, e_{n}\right)$ be a canonical basis for $\mathrm{A}^{n}, \Delta$ the alternating $n$-linear form on $\mathrm{A}^{n}$ such that $\Delta\left(e_{1}, \ldots, e_{n}\right)=1$ and $\varphi$ the endomorphism of $\mathrm{A}^{n}$ represented by $\alpha$. Then $\operatorname{det} \alpha=\operatorname{det} \varphi=\Delta\left(\varphi\left(e_{1}\right), \ldots, \varphi\left(e_{n}\right)\right)=$ $\sum a_{1}^{j_{1}} \cdots a_{n}^{j_{n}} \Delta\left(e_{j_{1}}, \ldots, e_{j_{n}}\right)$.

Now, $\Delta\left(e_{j_{1}}, \ldots, e_{j_{n}}\right)=\varepsilon(\sigma)$ if $\sigma: k \mapsto j_{k}$ is a permutation of $\{1, \ldots, n\}$ and 0 otherwise; so

$$
\operatorname{det} \alpha=\sum_{\sigma \in \mathfrak{S}_{n}} \varepsilon(\sigma) a_{1}^{\sigma(1)} \cdots a_{n}^{\sigma(n)}
$$

### 3.4.12 Minors

Let $\alpha=\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ be a matrix, and $\left(i_{1}, \ldots, i_{r}\right),\left(j_{1}, \ldots, j_{r}\right)$ sequences of elements of I and J respectively. Denote by $\operatorname{Min}_{i_{1}, \ldots, i_{r}}^{j_{1}, \ldots, j_{r}}(\alpha)$ the determinant of the matrix $\left(a_{i_{k}}^{j_{l}}\right)_{k, l \in\{1, \ldots, r\}}$; it is a minor of order $r$ of the matrix $\alpha$. Such a minor is null if any two of $i_{k}$, or any two of $j_{l}$, are equal. If $\mathrm{I}=\mathrm{J}=\{1, \ldots, n\}$, for $i, j \in\{1, \ldots, n\}$, write $\operatorname{Min}_{\hat{i}}^{\hat{J}}(\alpha)$ for the minor $\operatorname{Min}_{1, \ldots, i-1, i+1, \ldots, n}^{1, \ldots, j-1, \ldots+1, \ldots}(\alpha)$ of order $n-1$.

Proposition (Expansion of a determinant along a row) Let $\alpha=\left(a_{i}^{j}\right)_{i, j \in\{1, \ldots, n\}}$ be a square matrix of order $n$ and $j \in\{1, \ldots, n\}$. Then

$$
\operatorname{det} \alpha=\sum_{i \in\{1, \ldots, n\}}(-1)^{i-j} \operatorname{Min}_{\hat{\imath}}^{\hat{\jmath}}(\alpha) a_{i}^{j}
$$

## Proof

$$
\operatorname{det} \alpha=\sum_{\sigma \in \mathfrak{S}_{n}} \varepsilon(\sigma) a_{1}^{\sigma(1)} \cdots a_{n}^{\sigma(n)}=\sum_{\substack{i \in\{1, \ldots, n\}}} \sum_{\substack{\sigma \in \mathfrak{S}_{n} \\ \sigma(i)=j}} \varepsilon(\sigma) a_{1}^{\sigma(1)} \cdots a_{i}^{j} \cdots a_{n}^{\sigma(n)}
$$

Let $i$ and $j \in\{1, \ldots, n\}$; write $\left(a_{i^{\prime}}^{\prime j^{\prime}}\right)_{i^{\prime}, j^{\prime} \in\{1, \ldots, n-1\}}$ for the matrix obtained by deleting the $i$-th column and $j$-th row in $\alpha$ and by renumbering the indices from 1 to $n-1$. For every permutation $\sigma \in \mathfrak{S}_{n}$ such that $\sigma(i)=j$, define $\sigma^{\prime} \in \mathfrak{S}_{n-1}$ as the composite

$$
\{1, \ldots, n-1\} \rightarrow\{1, \ldots, \widehat{\imath}, \ldots, n\} \rightarrow\{1, \ldots, \widehat{\jmath}, \ldots, n\} \rightarrow\{1, \ldots, n-1\}
$$



$$
\begin{aligned}
\sum_{\substack{\sigma \in \mathfrak{S}_{n} \\
\sigma(i)=j}} \varepsilon(\sigma) a_{1}^{\sigma(1)} \cdots a_{i}^{j} \cdots a_{n}^{\sigma(n)} & =(-1)^{i-j} a_{i}^{j} \sum_{\sigma^{\prime} \in \mathfrak{S}_{n-1}} \varepsilon\left(\sigma^{\prime}\right) a_{1}^{\prime \sigma^{\prime}(1)} \cdots a_{n-1}^{\prime \sigma^{\prime}(n-1)} \\
& =(-1)^{i-j} a_{i}^{j} \operatorname{Min}_{i}^{\widehat{\jmath}}(\alpha) .
\end{aligned}
$$

The proposition follows.

### 3.4.13

Theorem (Cramer) Let E be a free A -module of rank n. An endomorphism $\varphi$ of E is an automorphism if and only if $\operatorname{det} \varphi$ is invertible.

If $\varphi$ is invertible, then $\operatorname{det} \varphi \cdot \operatorname{det} \varphi^{-1}=1$, and so $\operatorname{det} \varphi$ is invertible. We give two proofs for the converse.
$1^{s t}$ proof Let $\left(e_{1}, \ldots, e_{n}\right)$ be a basis for E and $\Delta$ the $n$-linear form on E such that $\Delta\left(e_{1}, \ldots, e_{n}\right)=1$. Denote by $\Lambda^{\prime}$ the module of alternating $(n-1)$-linear forms on E and let $\widetilde{\Delta}: \mathrm{E} \rightarrow \Lambda^{\prime}$ be the homomorphism defined by $\widetilde{\Delta}(x)\left(y_{1}, \ldots, y_{n-1}\right)=$ $\Delta\left(x, y_{1}, \ldots, y_{n-1}\right)$.
Lemma $\widetilde{\Delta}$ is an isomorphism.
Proof of the Lemma If $x=\sum x^{i} e_{i}$, then

$$
\begin{equation*}
\widetilde{\Delta}(x)\left(e_{1}, \ldots, \widehat{e}_{i}, \ldots, e_{n}\right)=(-1)^{i-1} x^{i} \tag{3.5}
\end{equation*}
$$

(in this formula, the term with $\widehat{\sim}$ is omitted). It follows from (1) that $\widetilde{\Delta}$ is injective. We show it is surjective. Let $f \in \Lambda^{\prime}, x^{i}=(-1)^{i-1} f\left(e_{1}, \ldots, \widehat{e_{i}}, \ldots, e_{n}\right)$ and $x=\sum x^{i} e_{i}$. The forms $f$ and $\widetilde{\Delta}(x)$ agree on all $(n-1)$-uples $\left(e_{1}, \ldots, \widehat{e_{i}}, \ldots, e_{n}\right)$, and so are equal.

End of the 1st Proof of the Theorem Define the endomorphism $\varphi^{*}$ of $\Lambda^{\prime}$ by $\varphi^{*} f\left(x_{1}, \ldots, x_{n-1}\right)=f\left(\varphi\left(x_{1}\right), \ldots, \varphi\left(x_{n-1}\right)\right)$. In the diagram

$\varphi^{*} \circ \widetilde{\Delta} \circ \varphi=\operatorname{det} \varphi \cdot \widetilde{\Delta} . \quad$ Indeed, $\quad\left(\varphi^{*} \circ \widetilde{\Delta} \circ \varphi(x)\right)\left(y_{1}, \ldots, y_{n-1}\right)=\Delta(\varphi(x)$, $\left.\varphi\left(y_{1}\right), \ldots, \varphi\left(y_{n-1}\right)\right)=\operatorname{det} \varphi \cdot \Delta\left(x, y_{1}, \ldots, y_{n-1}\right)$.

If $\operatorname{det} \varphi$ is invertible, setting $\psi=(\operatorname{det} \varphi)^{-1} \widetilde{\Delta}^{-1} \circ \varphi^{*} \circ \widetilde{\Delta}, \psi \circ \varphi=1_{\mathrm{E}}$, and $\varphi$ is left invertible. As $\operatorname{det} \psi \cdot \operatorname{det} \varphi=1$, $\operatorname{det} \psi$ is invertible. Likewise, $\varphi_{1}$ is a left inverse of $\psi$. Now, $\varphi_{1}=\varphi_{1} \circ \psi \circ \varphi=\varphi$, and so $\psi$ is an inverse of $\varphi$.

2nd Proof of the Theorem We may assume that $\mathrm{E}=\mathrm{A}^{n}$. The endomorphism $\varphi$ is then represented by a matrix $\alpha$. Set $b_{i}^{j}=(-1)^{j-i} \operatorname{Min}_{\hat{\imath}}^{\widehat{J}}(\alpha)$ and $\beta=\left(b_{i}^{j}\right)$. By Proposition 3.4.12, $\sum_{i} a_{i}^{j} b_{i}^{j}=\operatorname{det} \alpha$, and $\sum_{i} a_{i}^{k} b_{i}^{j}=0$ for $k \neq j$ (expansion of the determinant of a matrix with two equal rows). Hence $\sum_{i} a_{i}^{k} b_{i}^{j}=\delta_{j, k} \operatorname{det} \alpha$, i.e. $\alpha \cdot \beta=\operatorname{det} \alpha \cdot 1_{\mathrm{A}^{n}}$. Similarly, the expansion of a determinant along a column (analogous to that of Proposition 3.4.12 by interchanging rows and columns), gives $\beta \cdot \alpha=\operatorname{det} \alpha \cdot 1_{\mathrm{A}^{n}}$. If $\operatorname{det} \alpha$ is invertible, $\frac{1}{\operatorname{det} \alpha} \beta$ is then an inverse of $\alpha$.

### 3.4.14

Proposition Let E be a free A-module of rank $n$ and $\varphi$ an endomorphism of E . The following conditions are equivalent:
(i) $\varphi$ is an automorphism;
(ii) $\varphi$ is surjective;
(iii) $\operatorname{det} \varphi$ is invertible.

Proof (iii) $\Rightarrow$ (i) by the preceding theorem. (i) $\Rightarrow$ (ii) is obvious. We show that (ii) $\Rightarrow$ (iii). Suppose that $\varphi$ is surjective; by (3.4.7, Corollary), there is a homomorphism $\psi: \mathrm{E} \rightarrow \mathrm{E}$ such that $\varphi \circ \psi=1_{\mathrm{E}}$, and $\operatorname{so} \operatorname{det} \varphi \cdot \operatorname{det} \psi=1$.

### 3.4.15

Proposition Let E be a free A-module of rank $n$ and $\varphi$ an endomorphism of E . The following conditions are equivalent:
(i) $\varphi$ is injective;
(ii) $\operatorname{det} \varphi$ is not a zero divisor.

Proof (ii) $\Rightarrow$ (i). Let $\left(e_{1}, \ldots, e_{n}\right)$ be a basis for E and $\Delta$ an alternating $n$-linear form on E such that $\Delta\left(e_{1}, \ldots, e_{n}\right)=1$. For $x=\sum x^{i} e_{i} \in \operatorname{Ker} \varphi$, as $x^{i}=(-1)^{i-1} \Delta\left(x, e_{1}, \ldots\right.$, $\widehat{e_{i}}, \ldots, e_{n}$ ) for all $i, \operatorname{det} \varphi \cdot x^{i}=(-1)^{i-1} \Delta(\varphi(x), \ldots)=0$. Hence $x^{i}=0$ if $\operatorname{det} \varphi$ is not a zero divisor.

Not (ii) $\Rightarrow$ not (i). We show that if $\operatorname{det} \varphi$ is a zero divisor, then $\varphi$ is not injective. Let $h \neq 0$ be such that $h \operatorname{det} \varphi=0$. Let $\alpha=\left(a_{i}^{j}\right)$ be a matrix representing $\varphi$, and $\mu=$ $\operatorname{Min}_{i_{1}, \ldots, i_{r}}^{j_{1}, \ldots, j_{r}}(\alpha)$ a minor such that $h \mu \neq 0$ and whose order is as maximal as possible for this property. ${ }^{13}$ Since $h \operatorname{det} \varphi=0, r<n$. Let $i_{0} \notin\left\{i_{1}, \ldots, i_{r}\right\}$, and $x=\sum x^{i} e_{i}$, where $x^{i_{k}}=(-1)^{k} h \operatorname{Min}_{i_{0}, \ldots, \hat{i}_{k}, \ldots, i_{r}}^{j_{1}, \ldots, i_{r}}(\alpha)$ and $x^{i}=0$ if $i \neq i_{k}$ for any $k$. For all $j$,

$$
\sum a_{i}^{j} x^{i}=h \operatorname{Min}_{i_{0}, \ldots, i_{r}}^{j, j_{1}, \ldots, j_{r}}(\alpha)=0
$$

and so $x \in \operatorname{Ker} \varphi$. As $x_{i_{0}}=h \mu \neq 0, x \neq 0$ and $\varphi$ is not injective.

### 3.4.16 Presentations

Let M be an A-module. A presentation of M is an exact sequence $\mathrm{A}^{(\mathrm{J})} \rightarrow \mathrm{A}^{(\mathrm{I})} \rightarrow$ $\mathrm{M} \rightarrow 0$.

Proposition Every module has a presentation.
Proof Let M be an A-module and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a generating family of M . The homomorphism $u: \mathrm{A}^{(\mathrm{I})} \rightarrow \mathrm{M}$ defined by $u\left(e_{i}\right)=x_{i}$ is surjective; let N be its kernel and $\left(y_{j}\right)_{j \in \mathrm{~J}}$ its generators. The homomorphism $v: \mathrm{A}^{(\mathrm{J})} \rightarrow \mathrm{A}^{(\mathrm{I})}$ defined by $v\left(e_{j}\right)=y_{j}$ has image N , and $\mathrm{A}^{(\mathrm{J})} \xrightarrow{v} \mathrm{~A}^{(\mathrm{I})} \xrightarrow{u} \mathrm{M} \rightarrow 0$ is an exact sequence.

Definition An A-module $M$ is finitely presented if it has a presentation $\mathrm{A}^{(\mathrm{J})} \rightarrow \mathrm{A}^{(\mathrm{I})} \rightarrow \mathrm{M} \rightarrow 0$, where I et J are finite sets.

### 3.4.17

Proposition Every module is the direct limit of finitely presented modules.
Proof Let M be an A-module, $\mathrm{A}^{(\mathrm{J})} \xrightarrow{v} \mathrm{~A}^{(\mathrm{I})} \xrightarrow{u} \mathrm{M} \rightarrow 0$ a presentation of M and H the set of pairs ( $\mathrm{I}^{\prime}, \mathrm{J}^{\prime}$ ), where $\mathrm{I}^{\prime}$ and $\mathrm{J}^{\prime}$ are finite subsets of I and J respectively and $v\left(\mathrm{~A}^{\left(\mathrm{J}^{\prime}\right)}\right) \subset \mathrm{A}^{\left(\mathrm{I}^{\prime}\right)}$. The set H is directed with the order defined by $\left(\mathrm{I}^{\prime}, \mathrm{J}^{\prime}\right) \leqslant\left(\mathrm{I}^{\prime \prime}, \mathrm{J}^{\prime \prime}\right) \Leftrightarrow$

[^10]$\mathrm{I}^{\prime} \subset \mathrm{I}^{\prime \prime}$ and $\mathrm{J}^{\prime} \subset \mathrm{J}^{\prime \prime}$. Then $\bigcup_{\left(\mathrm{I}^{\prime}, \mathrm{J}^{\prime}\right) \in \mathrm{H}} \mathrm{I}^{\prime}=\mathrm{I}$ because it is always possible to take $\mathrm{J}^{\prime}=\varnothing$, and $\bigcup_{\left(\mathrm{I}^{\prime}, \mathrm{J}^{\prime}\right) \in \mathrm{H}} \mathrm{J}^{\prime}=\mathrm{J}$ because, for $j \in \mathrm{~J}, v\left(e_{j}\right)$ has finite support in I, and so is in some $\mathrm{A}^{\mathrm{I}^{\prime}}$. For $h=\left(\mathrm{I}^{\prime}, \mathrm{J}^{\prime}\right) \in \mathrm{H}$, let $v_{h}: \mathrm{A}^{\mathrm{J}^{\prime}} \rightarrow \mathrm{A}^{\mathrm{I}^{\prime}}$ be the map induced by $v$ and set $\mathrm{M}_{h}=$ Coker $v_{h}$. Then $\mathrm{A}^{(\mathrm{I})}=\bigcup \mathrm{A}^{\mathrm{I}^{\prime}}=\underline{\lim } \mathrm{A}^{\mathrm{I}^{\prime}}$ and $\mathrm{A}^{(\mathrm{J})}=\bigcup \mathrm{A}^{\mathrm{J}^{\prime}}=\underline{\lim } \mathrm{A}^{\mathrm{J}^{\prime}}$. Hence $\mathrm{M}=$ Coker $v=\underset{\longrightarrow}{\lim } \mathrm{M}_{h}$, and each $\mathrm{M}_{h}$ is finitely presented.

### 3.4.18

Proposition Let M be a finitely presented A-module and $\left(\left(\mathrm{E}_{\lambda}\right)\right.$, $\left.\left(r_{\lambda, \mu}\right)\right)$ be a direct system of A-modules. Then

$$
\operatorname{Hom}\left(\mathrm{M}, \underset{\longrightarrow}{\lim } \mathrm{E}_{\lambda}\right)=\underset{\longrightarrow}{\lim } \operatorname{Hom}\left(\mathrm{M}, \mathrm{E}_{\lambda}\right) .
$$

Proof If $L=A^{I}$ with I finite, then for any module $\mathrm{E}, \operatorname{Hom}(\mathrm{L}, \mathrm{E})=\mathrm{E}^{\mathrm{I}}=\mathrm{E}^{(\mathrm{I})}$. Hence the functor $\mathrm{E} \mapsto \operatorname{Hom}(\mathrm{L}, \mathrm{E}) \quad$ commutes with direct limits. If $\mathrm{M}=\operatorname{Coker}\left(v: \mathrm{L}_{1} \rightarrow \mathrm{~L}_{0}\right)$, where $\mathrm{L}_{0}$ and $\mathrm{L}_{1}$ are free and finitely generated, then $\operatorname{Hom}(\mathrm{M}, \mathrm{E})=\operatorname{Ker}\left(v^{*}: \operatorname{Hom}\left(\mathrm{L}_{0}, \mathrm{E}\right) \rightarrow \operatorname{Hom}\left(\mathrm{L}_{1}, \mathrm{E}\right)\right)$. So, by Proposition 3.3.14, the functor $\mathrm{E} \mapsto \operatorname{Hom}(\mathrm{M}, \mathrm{E})$ commutes with direct limits.

## Exercises 3.4. (Free modules and matrices)

1.-Show that, in the vector space $\mathbb{C}^{\mathbb{N}}$ of complex sequences, the sequences $\left(n^{k} r^{n}\right)_{n \in \mathbb{N}}$ with $k \in \mathbb{N}$ and $r \in \mathbb{C}$ form a free family. (For each sequence $\left(a_{n}\right)$, consider the series $\sum a_{n} z^{n}$.)
2.-(a) Let A be a not necessarily commutative ring and E a left A-module. Show that, if E has a finite basis, every basis for E is finite. Show that if E has an infinite basis, all bases have the same cardinality.
(b) Let E be a vector space over a field or skew field K. Let $\left(a_{1}, \ldots, a_{p}\right)$ be a free family in E , and G a generating family with $n$ elements. Show by induction that, for all $i \leqslant p$, there is a generating subset $\mathrm{G}_{i}$ of the form $\left\{a_{1}, \ldots, a_{p}\right\} \cup \mathrm{G}_{i}^{\prime}$, where $\mathrm{G}_{i}^{\prime}$ has $n-i$ elements. In particular $n \geqslant p$.

Show that any two bases for E necessarily have the same number of elements.
(c) Let H be an infinite dimensional Hilbert space and A the ring $\mathrm{L}(\mathrm{H}, \mathrm{H})$ of continuous endomorphisms of H . Show that, as a left A-module, A is isomorphic to $\mathrm{A}^{2}$ (if $\mathrm{H}_{1}$ and $\mathrm{H}_{2}$ are two complement subspaces of $H$, then $L(H, H)=L\left(\mathrm{H}_{1}, H\right) \oplus L\left(\mathrm{H}_{2}, H\right)$ ).

Show that A contains elements $u_{1}, v_{1}, u_{2}, v_{2}$, such that $v_{1} \cdot u_{1}=1, v_{2} \cdot u_{2}=1$ and $v_{2} \cdot u_{1}=0$. Show that there is no unital homomorphism from A to a field or skew field, nor to any not necessarily commutative nonzero ring where left multiplication by a nonzero element is injective.
3.-Let A be a (commutative) ring.
(a) Let I and J be not necessarily finite sets. Show that the homomorphisms from $\mathrm{A}^{(\mathrm{I})}$ to $\mathrm{A}^{(\mathrm{J})}$ correspond to matrices $\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ whose columns have finite support.

Extend the results of 3.4.9 and 3.4.10 to this context. What are the homomorphisms for which the matrix only has finitely many nonzero rows?
(b) Let E be a free infinite dimensional A-module. Using the formula of Proposition 3.4.11, define $\operatorname{det}\left(1_{\mathrm{E}}+u\right)$ for an endomorphism $u$ of E whose image is contained in a finitely generated submodule of E . Prove the formula $\operatorname{det}\left(\left(1_{\mathrm{E}}+u\right)\left(1_{\mathrm{E}}+v\right)\right)=$ $\operatorname{det}\left(1_{\mathrm{E}}+u\right) \operatorname{det}\left(1_{\mathrm{E}}+v\right)$.
4.-(a) Does Proposition 3.4.11 hold when $I=\varnothing$ ?
(b) Does the proof of 3.4.15 still hold without the footnote?
5.-(a) Let $\mathrm{E}=\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}$ and $\mathrm{F}=\bigoplus_{j \in \mathrm{~J}} \mathrm{~F}_{j}$ be A-modules expressible as a finite direct sum of submodules. Show that, for all homomorphisms $f: \mathrm{E} \rightarrow \mathrm{F}$, there is a unique family $\left(a_{i}^{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ with $a_{i}^{j} \in \operatorname{Hom}\left(\mathrm{E}_{i}, \mathrm{~F}_{j}\right)$ such that $f(x)=\sum_{j} a_{i}^{j}(x)$ for $x \in \mathrm{E}_{i}$.

The family $\left(a_{i}^{j}\right)$ is said to be the matrix representing $f$ with respect to the given decompositions.

How is the matrix representing the composite of two homomorphisms computed?
(b) Let $f: \mathrm{E}_{1} \oplus \mathrm{E}_{2} \rightarrow \mathrm{~F}_{1} \oplus \mathrm{~F}_{2}$ be a homomorphism represented by a matrix $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ with $b$ invertible. Determine the homomorphisms $u: \mathrm{E}_{1} \rightarrow \mathrm{E}_{2}$ and $v: \mathrm{E}_{1} \rightarrow \mathrm{~F}_{2}$ such that $\operatorname{Ker} f$ is the graph of the restriction from $u$ to $\operatorname{Ker} v$.
6.-(a) Let $\mathrm{E}=\left(\bigoplus \mathrm{E}_{n}, d\right)$ be a bounded complex (cf. 3.10) of A-modules, and for all $n$, let $\mathrm{E}_{n}$ be the direct sum of three submodules $\mathrm{T}_{n}, \mathrm{~F}_{n}, \mathrm{~S}_{n}$. The homomorphism $d_{n}: \mathrm{E}_{n} \rightarrow \mathrm{E}_{n-1}$ is represented by a matrix:

$$
\left(\begin{array}{ccc}
a_{n} & e_{n} & h_{n} \\
b_{n} & f_{n} & i_{n} \\
c_{n} & g_{n} & j_{n}
\end{array}\right)
$$

Assume that, for all $n, h_{n}: \mathrm{S}_{n} \rightarrow \mathrm{~T}_{n-1}$ is an isomorphism. Show that, for all $n$, there is another decomposition $\left(\mathrm{T}_{n}^{\prime}, \mathrm{F}_{n}^{\prime}, \mathrm{S}_{n}^{\prime}\right)$ of $\mathrm{E}_{n}$, where $\mathrm{T}_{n}^{\prime}, \mathrm{F}_{n}^{\prime}, \mathrm{S}_{n}^{\prime}$ are respectively isomorphic to $\mathrm{T}_{n}, \mathrm{~F}_{n}, \mathrm{~S}_{n}$, with respect to which the matrix representing $d_{n}$ is

$$
\left(\begin{array}{ccc}
0 & 0 & h_{n}^{\prime} \\
0 & f_{n}^{\prime} & 0 \\
0 & 0 & 0
\end{array}\right)
$$

$h_{n}^{\prime}$ being an isomorphism. (First set $\mathrm{T}_{n}^{\prime}=d_{n+1}\left(\mathrm{~S}_{n+1}\right)$, so that $a, b, c, i$ and $j$ become zero, then successively modify each $\mathrm{F}_{n}$ in such a way that $e_{n}$ becomes null; show that the homomorphisms $g_{n}$ are not necessarily null.)
(b) Does the result of (a) extend to unbounded complexes?
7.-Let $f: \mathrm{A}^{p} \rightarrow \mathrm{~A}^{q}$ be a A-linear map. Show that $f$ is surjective if and only if the ideal generated by the minors of order $q$ of the matrix representing $f$ is the whole of the ring A.
8. (Projective modules)—Let E be an A-module.
(a) Show that, if E is projective and finitely generated, then it is a direct factor of a finitely generated free module.
(b) Show that, if E is projective, for all exact sequences $\mathrm{M}^{\prime} \rightarrow \mathrm{M} \rightarrow \mathrm{M}^{\prime \prime}$ of Amodules, $\operatorname{Hom}\left(E, M^{\prime}\right) \rightarrow \operatorname{Hom}(E, M) \rightarrow \operatorname{Hom}\left(E, M^{\prime \prime}\right)$ is an exact sequence.
(c) Show that the space E of analytic functions $g: \mathbb{R} \rightarrow \mathbb{R}$ such that $g(x+2 \pi)=$ $-g(x)$ is a non-free projective module over the ring A of 3.1, Exercise 12.

Show that all maximal ideals $\mathfrak{m}_{x}$ in this ring are modules isomorphic to E .
9. (Projective modules and vector bundles)—Let X be a topological space. A real (resp. complex) vector bundle over X is a topological space E equipped with
(i) a continuous map $\pi: \mathrm{E} \rightarrow \mathrm{X}$,
(ii) the structure of a real (resp. complex) vector space on the fibre $\mathrm{E}(x)=\pi^{-1}(x)$ for all $x \in \mathrm{X}$,
satisfying the following condition:
(FV) For all $x \in \mathrm{X}$, there is a neighbourhood U of $x$, an integer $n$ and a homeomorphism $\varphi$ from $\left.\mathrm{E}\right|_{\mathrm{U}}=\pi^{-1}(\mathrm{U})$ onto $\mathrm{U} \times \mathbb{R}^{n}$ (resp. $\mathrm{U} \times \mathbb{C}^{n}$ ) inducing an isomorphism of vector spaces from $\mathrm{E}(x)$ onto $\mathbb{R}^{n}$ (resp. $\left.\mathbb{C}^{n}\right)$ for all $x \in \mathrm{U}(\varphi$ is called a trivialization from E onto $\mathrm{U} ; \mathrm{E}$ is said to be trivial if there is a trivialization from E over X).

A continuous section of E is a continuous map $s: \mathrm{X} \rightarrow \mathrm{E}$ such that $\pi \circ s=1_{\mathrm{X}}$.
If E and F are two vector bundles over X , a morphism from E to F is a continuous map from E to F inducing a linear map from $\mathrm{E}(x)$ to $\mathrm{F}(x)$ for all $x \in \mathrm{X}$.
(a) Show that the continuous sections of E form a natural module $\mathrm{S}(\mathrm{E})$ over the ring $\mathrm{C}(\mathrm{X})=\mathrm{C}(\mathrm{X}, \mathbb{R})$ (resp. $\mathrm{C}(\mathrm{X}, \mathbb{C})$ ), and that for every morphism $f: \mathrm{E} \rightarrow \mathrm{F}$, the map $f_{*}: s \mapsto f \circ s$ from $\mathrm{S}(\mathrm{E})$ to $\mathrm{S}(\mathrm{F})$ is $\mathrm{C}(\mathrm{X})$-linear.
(b) Henceforth, X is assumed to be compact and metrizable. Show that, if E is a vector bundle over X passing through every point de E , then there is a continuous section (i.e. $(\forall t \in \mathrm{E})(\exists s \in \mathrm{~S}(\mathrm{E})) s(\pi(t))=t)$.
(c) Let $\mathscr{U}=\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ be a finite open cover of X . Show that there is a family $\left(f_{i}\right)_{i \in \mathrm{I}}$ of continuous functions on X such that for all $i$, the function $f_{i}$ has support ${ }^{14}$ in $\mathrm{U}_{i}$, and that $\sum f_{i}$ are strictly positive at all points. Deduce that there is a family of continuous functions $\left(\eta_{i}\right)_{i \in \mathrm{I}}$ on X such that $\eta_{i}$ have support in $\mathrm{U}_{i}$ for all $i$, and that $\sum \eta_{i}=1$ (such a family is called a partition of unity subject to $\mathscr{U}$ ).
(d) Let E and F be two vector bundles over X , and $f: \mathrm{E} \rightarrow \mathrm{F}$ a surjective morphism. Show that, for all $x \in \mathrm{X}$, there is a neighbourhood U of $x$ and a morphism $g_{\mathrm{U}}:\left.\left.\mathrm{F}\right|_{\mathrm{U}} \rightarrow \mathrm{E}\right|_{\mathrm{U}}$ such that $\left.g_{\mathrm{U}} \circ f\right|_{\mathrm{U}}=1_{\left.\mathrm{E}\right|_{\mathrm{U}}}$. Using a partition of unity show that there is a morphism $g: \mathrm{F} \rightarrow \mathrm{E}$ such that $g \circ f=1_{\mathrm{E}}$. Show that $f_{*}: \mathrm{S}(\mathrm{E}) \rightarrow \mathrm{S}(\mathrm{F})$ is surjective and has a $\mathrm{C}(\mathrm{X})$-linear section.
(e) Let E be a vector bundle over X . Show that a finite family $\left(s_{i}\right)$ of continuous sections of E such that, for all $x \in \mathrm{X},\left(s_{i}(x)\right)$ generate $\mathrm{E}(x)$. Deduce that E is a direct factor of a trivial bundle, and that $S(E)$ is a projective $C(X)$-module.

[^11](f) Show that the functor $\mathrm{S}: \mathrm{E} \mapsto \mathrm{S}(\mathrm{E})$ from the category of vector bundles over X to the category of $\mathrm{C}(\mathrm{X})$-modules is fully faithful. (To show that $\operatorname{Hom}(\mathrm{E}, \mathrm{F}) \rightarrow \operatorname{Hom}(\mathrm{S}(\mathrm{E}), \mathrm{S}(\mathrm{F}))$ is bijective, E may first be considered trivial.)
(g) Let E be a vector bundle over X and $p$ a projection of E (i.e. a morphism from E to itself such that $p \circ p=p$ ). Show that the image of $p$ vector is a subbundle of E.
(h) Show that the functor $S$ defines an equivalence between the category of vector bundles over X and the category of projective $\mathrm{C}(\mathrm{X})$-modules.

This reduces the study and the classification of projective $\mathrm{C}(\mathrm{X})$-modules to questions of algebraic topology. For rings that are not of the form $\mathrm{C}(\mathrm{X})$, the study of projective modules draws methods from algebraic topology.
(i) In the various parts of this problem, can the assumption of the compactness and metrizability of X be weakened? (For (c), only suppose that the open cover $\mathscr{U}$ is locally finite.)
10.-(a) Let $A$ be a ring, and let $J(A)$ be the set of isomorphism classes of finitely generated projective A-modules. Find an upper bound for the cardinality of $J(A)$ in terms of that of A.
(b) The action ( $\mathrm{E}, \mathrm{F}$ ) $\mapsto \mathrm{E} \oplus \mathrm{F}$ defines a structure of commutative monoid on $J(A)$. Show that there is a commutative group $K(A)$ and a homomorphism $\varepsilon: J(A) \rightarrow K(A)$ satisfying the following universal property: for any group $G$ and any homomorphism $f: \mathrm{J}(\mathrm{A}) \rightarrow \mathrm{G}$, there is a unique homomorphism $\bar{f}: \mathrm{K}(\mathrm{A}) \rightarrow \mathrm{G}$ such that $f=\bar{f} \circ \varepsilon$.
(c) Let $\xi, \eta \in \mathrm{J}(\mathrm{A})$ with representatives E and F . Show that $\varepsilon(\xi)=\varepsilon(\eta)$ if and only if $\mathrm{E} \oplus \mathrm{A}^{n}$ is isomorphic to $\mathrm{F} \oplus \mathrm{A}^{n}$ for sufficiently large $n$. (For an example where $\varepsilon$ is not injective, see 4.7, Exercise 7.)
11.-(a) Show that in the Banach space $\ell^{1}$ of real sequences $x_{n}$ such that $\sum\left|x_{n}\right|<\infty$, the sequences $\left(e^{-\lambda n}\right)_{n \in \mathbb{N}}$ form a free family as $\lambda$ runs through $] 0, \infty[$.
(b) Using Baire's theorem show that an infinite dimensional Banach space does not have a countable algebraic basis.
(c) Show that, for an infinite dimensional Banach space E, there is a continuous injective linear map from $\ell^{1}$ to E .
(d) Show that, for any Hausdorff Banach space E (i.e. with a countable subset generating a dense subspace), the cardinality of E is $2^{\aleph_{0}}$. Deduce that all Hausdorff Banach spaces are algebraically isomorphic. Give examples where they are not isomorphic as topological vector spaces.

### 3.5 Modules over Principal Ideal Domains

Throughout this section, A denotes a PID; for $a \in \mathrm{~A},(a)$ the ideal generated by $a$, and $\chi$ the canonical map from $A$ to $\operatorname{Mon}(A)$ (3.2.7).

### 3.5.1

Theorem Every subring of a free A-module is free.
AdDEndum. Let E be an A-module with basis $\left(e_{i}\right)_{i \in \mathrm{I}}$ and F a submodule of E . Then F has a basis $\left(f_{j}\right)_{j \in \mathrm{I}}$, where the cardinality of $\mathrm{I}^{\prime}$ is less than or equal to that of I .

For every subset J of I , let $\mathrm{E}_{\mathrm{J}}$ be the submodule de E generated by $\left(e_{i}\right)_{i \in \mathrm{~J}}$ and set $\mathrm{F}_{\mathrm{J}}=\mathrm{F} \cap \mathrm{E}_{\mathrm{J}}$. Let $\Omega$ be the set of triplets $\left(\mathrm{J}, \mathrm{K},\left(f_{i}\right)_{i \in \mathrm{~K}}\right)$, where J is a subset of $\mathrm{I}, \mathrm{K}$ a subset of J , and $\left(f_{i}\right)_{i \in \mathrm{~K}}$ a basis for $\mathrm{F}_{\mathrm{J}}$. For $\alpha=\left(\mathrm{J}, \mathrm{K},\left(f_{i}\right)_{i \in \mathrm{~K}}\right)$ and $\alpha^{\prime}=$ $\left(\mathbf{J}^{\prime}, \mathrm{K}^{\prime},\left(f_{i}^{\prime}\right)_{i \in \mathrm{~K}^{\prime}}\right)$ in $\Omega$, write $\alpha \leqslant \alpha^{\prime}$ if $\mathrm{J} \subset \mathrm{J}^{\prime}, \mathrm{K}=\mathrm{K}^{\prime} \cap \mathrm{J}$ and $(\forall i \in \mathrm{~K}) f_{i}^{\prime}=f_{i}$.

This is an order relation on $\Omega$.
Lemma Let $\alpha=\left(\mathrm{J}, \mathrm{K},\left(f_{i}\right)\right) \in \Omega$. If $\alpha$ is maximal, then $\mathrm{J}=\mathrm{I}$.
Proof of the Lemma Suppose that $\mathbf{J} \neq \mathbf{I}$. We define $\alpha^{\prime} \in \Omega$ such that $\alpha^{\prime}>\alpha$. Let $i_{0} \in \mathrm{I}-\mathrm{J}$, and set $\mathrm{J}^{\prime}=\mathrm{J} \cup\left\{i_{0}\right\}$. Then $\mathrm{E}_{\mathrm{J}^{\prime}}=\mathrm{E}_{\mathbf{J}} \oplus \mathrm{A} e_{i_{0}}$. Denote by $p$ the projection $\mathrm{E}_{\mathrm{J}^{\prime}} \rightarrow \mathrm{A} e_{i_{0}}$ and set $\mathrm{Q}=p\left(\mathrm{~F}_{\mathrm{J}^{\prime}}\right)$. There is an exact sequence

$$
\begin{equation*}
0 \rightarrow \mathrm{~F}_{\mathrm{J}} \rightarrow \mathrm{~F}_{\mathrm{J}^{\prime}} \rightarrow \mathrm{Q} \rightarrow 0 . \tag{*}
\end{equation*}
$$

The submodule Q of $\mathrm{A} e_{i_{0}}$ is of the form $\mathfrak{a} e_{i_{0}}$, where $\mathfrak{a}$ is an ideal in A , hence of the form Aae $i_{0}$ with $a \in \mathrm{~A}$ since A is principal. If $a=0$, then $\mathrm{F}_{J^{\prime}}=\mathrm{F}_{\mathrm{J}}$ and $\alpha^{\prime}=$ $\left(\mathrm{J}^{\prime}, \mathrm{K},\left(f_{i}\right)_{i \in \mathrm{~K}}\right)$ is a strict upper bound of $\alpha$. If $a \neq 0$, then $a e_{i_{0}}$ is a basis for Q , and the sequence $(*)$ splits. Set $\mathrm{K}^{\prime}=\mathrm{K} \cup\left\{i_{0}\right\}$ and choose $f_{i_{0}} \in \mathrm{~F}^{\prime}$ such that $p\left(f_{i_{0}}\right)=a e_{i_{0}}$. Then $\alpha$ is strictly bounded above by $\alpha^{\prime}=\left(\mathrm{J}^{\prime}, \mathrm{K}^{\prime},\left(f_{i}\right)_{i \in \mathrm{~K}^{\prime}}\right)$.

Proof of the Theorem The ordered set $\Omega$ is inductive. Indeed, let $\left(\mathrm{J}_{\lambda}\right.$, $\left.\mathrm{K}_{\lambda},\left(f_{i, \lambda}\right)_{i \in \mathrm{~K}_{\lambda}}\right)_{\lambda \in \Lambda}$ be a totally ordered family of elements of $\Omega$; set $\mathbf{J}=\bigcup \mathbf{J}_{\lambda}$, $\mathrm{K}=\bigcup \mathrm{K}_{\lambda}$, and for $i \in \mathrm{~K}$ let $f_{i}$ be the common value of $f_{i, \lambda}$ for $\lambda$ such that $i \in \mathrm{~J}_{\lambda}$; the family $\left(f_{i}\right)_{i \in \mathrm{~K}}$ is a basis for $\mathrm{F}_{\mathrm{J}}$ since all elements of $\mathrm{F}_{\mathrm{J}}$ are in some $\mathrm{F}_{\mathrm{J}_{\lambda}}$, so are linear combinations of the elements $f_{i}$, and a relation between these $f_{i}$ only involves finitely many $i$; therefore the family is bounded above by $\left(\mathrm{J}, \mathrm{K},\left(f_{i}\right)_{i \in \mathrm{~K}}\right)$.

By Zorn's theorem, $\Omega$ has a maximal element $\left(\mathrm{J}, \mathrm{K},\left(f_{i}\right)_{i \in \mathrm{~K}}\right)$, and $\mathrm{J}=\mathrm{I}$ necessarily holds. So $\left(f_{i}\right)_{i \in \mathrm{~K}}$ is a basis for F .

### 3.5.2 Adapted Bases for Submodules

Let E be a free A -module, F a submodule of E and $\left(e_{i}\right)_{i \in \mathrm{I}}$ a basis for E . This basis is said to be adapted F if there is a subset J of I and a family $\left(a_{i}\right)_{i \in \mathrm{~J}}$ in $\mathrm{A}^{*}$ such that $\left(a_{i} e_{i}\right)_{i \in \mathrm{~J}}$ is a basis for F .

If there is a basis for E adapted to F , then the quotient module $\mathrm{E} / \mathrm{F}$ is the direct sum of a free module and a torsion module, the latter being the direct sum of cyclic modules. Indeed, with the above notation, $\mathrm{E} / \mathrm{F}$ is isomorphic to $\mathrm{A}^{(\mathrm{I}-\mathrm{J})} \oplus \oplus_{i \in \mathrm{~J}} \mathrm{~A} /\left(a_{i}\right)$.

### 3.5.3 Content

We use the notation of 3.2.7. Let E be a free A -module and $x \in \mathrm{E}$. Let $\left(x_{i}\right)$ be the coordinates of $x$ with respect to the basis $\left(e_{i}\right)$. The g.c.d. $d$ of $\chi\left(x_{i}\right)$ in $\operatorname{Mon}(\mathrm{A})$ is independent of the choice of the basis $\left(e_{i}\right)$ since, for $a \in \mathrm{~A}, \chi(a)\langle d$ if and only if $(\exists y \in \mathrm{E}) x=a y$. It is said to be the content of $x$ and is written $c(x)$. An element $a \in \mathrm{~A}$ is a content of $x$ if $\chi(a)=c(x)$.

### 3.5.4

Proposition Let E be a free A-module, $x \in \mathrm{E}$ and a a content of $x$. Then, for all linear forms $\varphi: \mathrm{E} \rightarrow \mathrm{A}, a\langle\varphi(x)$, and there is a linear form $\varphi$ such that $\varphi(x)=a$.

Proof There exists $y \in \mathrm{E}$ such that $x=a y$; so $\varphi(x)=a \varphi(y)$, and $a\langle\varphi(x)$. We may assume that $\mathrm{E}=\mathrm{A}^{(\mathrm{I})}$ and $x=\left(x_{i}\right)$. Then $a$ is a g.c.d. of the $x_{i}$ and, by Bezout's formula (3.2.11, Corollary), there are $u_{i} \in \mathrm{~A}$ such that $a=\sum u_{i} x_{i}$, and the linear form $\varphi:\left(y_{i}\right) \mapsto \sum u_{i} y_{i}$ on E satisfies $\varphi(x)=a$.

Remark The proposition can also be stated as follows: Let $E^{\top}=\operatorname{Hom}(E, A)$ and define $\delta_{x}: \mathrm{E}^{\top} \rightarrow \mathrm{A}$ by $\delta_{x}(\varphi)=\varphi(x)$. Then $\operatorname{Im} \delta_{x}=\mathrm{A} \cdot c(x)$.

### 3.5.5 Elements of Content 1

Proposition Let E be a free A-module and $x \in \mathrm{E}$. The following conditions are equivalent:
(i) $c(x)=1$;
(ii) $(\exists \varphi: \mathrm{E} \rightarrow \mathrm{A}) \varphi(x)=1$;
(iii) $x \neq 0$ and $\mathrm{A} x$ is a direct factor of E ;
(iv) there is a basis $\left(e_{i}\right)_{i \in \mathrm{I}}$ such that $e_{i_{0}}=x$ for some $i_{0} \in \mathrm{I}$.

Proof By 3.5.4, (i) $\Rightarrow$ (ii). If (ii) holds, then $x \neq 0$, and $y \mapsto \varphi(y) \cdot x$ is a linear retraction from E onto $\mathrm{A} x$, implying (iii). Assume (iii) holds and that F such that $\mathrm{E}=\mathrm{A} x \oplus \mathrm{~F}$. By 3.5.1, F is free. Let $\left(e_{i}\right)_{i \in \mathrm{~J}}$ be a basis for F , and $i_{0} \notin \mathrm{~J}$. Set $\mathrm{I}=\mathrm{J} \cup\left\{i_{0}\right\}$ and $e_{i_{0}}=x$. As $x$ is torsion-free, $\left(e_{i}\right)_{i \in \mathrm{I}}$ is a basis for E , and so (iv) follows. Finally, (iv) $\Rightarrow$ (i) is immediate.

### 3.5.6 Elements of Minimal Content in a Submodule

Let E be a free A -module and F a submodule of E . The ring A , being a PID, is a a UFD, and so every decreasing sequence of elements of Mon(A) is stationary. By 1.6.1, there exists $f \in \mathrm{~F}$ such that $c(f)$ is a minimal element of $\{c(y)\}_{y \in \mathrm{~F}}$.

Proposition Let E be a free A-module, F a nonzero submodule of $\mathrm{E}, f \in \mathrm{~F}$ an element of minimal content, $a \in \mathrm{~A}^{*}$ a content of $f$ and $e \in \mathrm{E}$ such that $f=a e$. Then
(a) $\mathrm{A} e$ is a direct factor of E ;
(b) Let $\mathrm{E}_{1}$ be a complement of Ae in E and set $\mathrm{F}_{1}=\mathrm{F} \cap \mathrm{E}_{1}$; then $\mathrm{F}=\mathrm{A} f \oplus \mathrm{~F}_{1}$;
(c) $c(f)$ is the smallest element of $\{c(y)\}_{y \in \mathrm{~F}}$.

Proof (a) $c(f)=a c(e)$ and soù $c(e)=1$, and the assertion follows from Proposition 3.5.5.
(b) Let $p: \mathrm{E} \rightarrow \mathrm{A} e$ be the projection with kernel $\mathrm{E}_{1}$. The submodule $p(\mathrm{~F})$ of Ae is of the form $\mathfrak{b e}$ where $\mathfrak{b}$ is an ideal, and so of the form Abe, where $b \in \mathrm{~A}$ since A is principal. Since $f=a e \in p(\mathrm{~F})=\mathrm{A} b e, b\langle a$. There exists $y \in \mathrm{~F}$ such that $p(y)=b e$. Hence $b=\varphi(y)$ for some linear form $\varphi$ on E , and so $c(y)\langle\chi(b)$ by 3.5.4.
$c(y)\langle\chi(b)\langle\chi(a)=c(f)$; as $c(f)$ is minimal, $c(y)=c(f)$ and $\chi(b)=\chi(a)$. Hence $p(\mathrm{~F})=\mathrm{A} b e=\mathrm{A} a e=\mathrm{A} f \subset \mathrm{~F}$, and $p$ induces a projection of F with image $\mathrm{A} f$ and kernel $\mathrm{F} \cap \mathrm{E}_{1}=\mathrm{F}_{1}$. So $\mathrm{F}=\mathrm{A} f \oplus \mathrm{~F}_{1}$ by 3.3.11.
(c) $c(f)$ is minimal by assumption. We show that, for all $y \in \mathrm{~F}, c(f)\langle c(y)$. Let $y \in \mathrm{~F}$, write $y=\lambda f+y_{1}$ with $y_{1} \in \mathrm{~F}_{1}$. Then $c(y)=$ g.c.d. $\left(\lambda c(f), c\left(y_{1}\right)\right)$ follows by considering a basis for E consisting of $e$ and of a basis for $\mathrm{E}_{1}$. To show that $c(f)\left\langle c(y)\right.$ it suffices to show that $c(f)\left\langle c\left(y_{1}\right)\right.$. Set $y^{\prime}=f+y_{1}$. Then $c\left(y^{\prime}\right)=$ g.c.d. $\left(c(f), c\left(y_{1}\right)\right)\left\langle c(f)\right.$, and so $c\left(y^{\prime}\right)=c(f)$ since $c(f)$ is minimal. Hence $c(f)=c\left(y^{\prime}\right)\left\langle c\left(y_{1}\right)\right.$.

### 3.5.7

Proposition (Adapted basis Theorem) Let E be a finitely generated free A-module. For all submodules F of E there is basis for E adapted to F .

ADDENDUM. There is a basis $\left(e_{1}, \ldots e_{n}\right)$ for E , an integer $p \leqslant n$ and elements $a_{1}, \ldots, a_{p}$ of $\mathrm{A}^{*}$ such that $\left(a_{1} e_{1}, \ldots, a_{p} e_{p}\right)$ is a basis for F and $a_{1}\left\langle a_{2}\left\langle\cdots\left\langle a_{p}\right.\right.\right.$.
Proof By induction on the rank $p$ of F . If $p=0$, then there is nothing to show. Suppose that F is of rank $p \neq 0$ and that $f_{1} \in \mathrm{~F}$ has minimal content. Let $a_{1} \in \mathrm{~A}^{*}$ be a content of $f_{1}$ and $e_{1} \in \mathrm{E}$ such that $a_{1} e_{1}=f_{1}$. By 3.5.6, $\mathrm{E}=\mathrm{A} e_{1} \oplus \mathrm{E}_{1}$ and $\mathrm{F}=\mathrm{A} f_{1} \oplus \mathrm{~F}_{1}$, where $\mathrm{E}_{1}$ is a submodule of E and $\mathrm{F}_{1}=\mathrm{F} \cap \mathrm{E}_{1}$. By 3.5.1, the modules $\mathrm{E}_{1}$ and $\mathrm{F}_{1}$ are free. $\operatorname{As} \operatorname{rk}(\mathrm{F})=\operatorname{rk}\left(\mathrm{F}_{1}\right)+1, \mathrm{~F}_{1}$ is free of rank $p-1$. By the induction hypothesis, there is a basis $\left(e_{2}, \ldots, e_{n}\right)$ for $\mathrm{E}_{1}$ and elements $a_{2}, \ldots, a_{p}$ of $\mathrm{A}^{*}$ such that $\left(a_{2} e_{2}, \ldots, a_{p} e_{p}\right)$ is a basis for $\mathrm{F}_{1}$ and $a_{2}\left\langle\cdots\left\langle a_{p}\right.\right.$. Then $\left(e_{1}, \ldots, e_{n}\right)$ is a basis for E , $\left(a_{1} e_{1}, \ldots, a_{p} e_{p}\right)$ a basis for F , and, if $p>1, a_{1}\left\langle a_{2}\right.$ by 3.5.6, since $a_{2}$ is a content of $f_{2}=a_{2} e_{2} \in \mathrm{~F}$.

### 3.5.8

Theorem Every finitely generated A-module is the direct sum of cyclic submodules.

ADDENDUM. Every finitely generated A-module is isomorphic to a module $\mathrm{A}^{s} \oplus \bigoplus_{i \in \mathrm{I}} \mathrm{A} /\left(a_{i}\right)$, where $s \in \mathbb{N}$, I is a finite set and $a_{i} \in \mathrm{~A}^{*}$ for $i \in \mathrm{I}$. In such a module, the second direct factor is a torsion submodule.

Proof Every finitely generated A-module is isomorphic to the quotient of a finitely generated free A-module E by a submodule of E. The theorem and its addendum then follow from Proposition 3.5.7 and from 3.5.2.

Corollary 3.1 In a finitely generated A-module, the torsion submodule is a direct factor.

Corollary 3.2 All finitely generated torsion-free A-modules are free.
Corollary 3.3 All commutative finite groups are products of cyclic groups.

### 3.5.9

Two elements $x, y \in \mathrm{~A}$ are said to be relatively prime if the ideal they generate is A , i.e. if 1 is a g.c.d. of $x$ and $y$.

Proposition (Chinese Theorem for two elements) Let $x$ and $y$ be two relatively prime elements of A . Then the module $\mathrm{A} /(x y)$ is isomorphic to $\mathrm{A} /(x) \oplus \mathrm{A} /(y)$.

ADDENDUM. Let $\chi_{x}: \mathrm{A} \rightarrow \mathrm{A} /(x)$ and $\chi_{y}: \mathrm{A} \rightarrow \mathrm{A} /(y)$ be the canonical maps. The homomorphism $\Delta: z \mapsto\left(\chi_{x}(z), \chi_{y}(z)\right)$ from A to $\mathrm{A} /(x) \times \mathrm{A} /(y)$ factorizes naturally (3.3.7) to give an isomorphism from $\mathrm{A} /(x y)$ onto $\mathrm{A} /(x) \oplus \mathrm{A} /(y)$. A more general version will be considered in (5.1.5).

Proof The kernel of $\Delta$ is the ideal of common multiples of $x$ and $y$ and is generated by $x y$ since $x y$ is a 1.c.m. of $x$ and $y(3.2 .5)$.

We show that $\Delta$ is surjective. By Bezout's formula (3.2.11, Corollary), there exist $u, v \in \mathrm{~A}$ such that $u x+v y=1$. Let: $(\alpha, \beta)=\left(\chi_{x}(a), \chi_{y}(b)\right) \in \mathrm{A} /(x) \times \mathrm{A} /(y)$; set $c=v y a+u x b$. As $c=a+u(b-a) x=b+v(a-b) y$,

$$
\chi_{x}(c)=\alpha, \quad \chi_{y}(c)=\beta, \quad \Delta(c)=(\alpha, \beta) .
$$

Hence $\Delta$ is surjective, and so the proposition and its addendum follow.
Corollary Let $x \in \mathrm{~A}^{*}$. Write $x$ as $x=u \cdot p_{1}^{r_{1}} \cdots p_{k}^{r_{k}}$, where $u$ is invertible, all $p_{i}$ irreducible and $\chi\left(p_{i}\right)$ distinct in $\operatorname{Mon}^{*}(\mathrm{~A})$. Then the module $\mathrm{A} /(x)$ is isomorphic to $\mathrm{A} /\left(p_{1}^{r_{1}}\right) \oplus \cdots \oplus \mathrm{A} /\left(p_{k}^{r_{k}}\right)$.

### 3.5.10 Elementary Divisors

Proposition and Definition Let E be a finitely generated A-module.
(a) The module E is isomorphic to a module

$$
\mathrm{A}^{s} \oplus \bigoplus_{i \in \mathrm{I}} \mathrm{~A} /\left(p_{i}^{r_{i}}\right)
$$

where $s \in \mathbb{N}$, I is a finite set, and, for all $i \in \mathrm{I}, p_{i} \in \mathrm{~A}^{*}$ is irreducible and $r_{i}>0$.
(b) If $\mathrm{E} \approx \mathrm{A}^{s} \oplus \bigoplus_{i \in \mathrm{I}} \mathrm{A} /\left(p_{i}^{r_{i}}\right) \approx \mathrm{A}^{s^{\prime}} \oplus \bigoplus_{i \in \mathrm{I}^{\prime}} \mathrm{A} /\left(p_{i}^{\prime r_{i}^{\prime}}\right)$, then as the conditions of (a) hold for $s, \mathrm{I}, p_{i}, r_{i}$ and for $s^{\prime}, \mathrm{I}^{\prime}, p_{i}^{\prime}, r_{i}^{\prime}, s=s^{\prime}$ and there is a bijection $\tau$ from I onto $\mathrm{I}^{\prime}$ such that $\chi\left(p_{i}\right)=\chi\left(p_{\tau(i)}^{\prime}\right)$ and $r_{i}=r_{\tau(i)}^{\prime}$ for all $i \in \mathrm{I}$.
$\left(p_{i}^{r_{i}}\right)_{i \in \mathrm{I}}$ is the family of elementary divisors of E .
(a) follows from 3.5.8 and 3.5.9. The proof of (b) uses a lemma:

Lemma Let $\mathrm{E}=\mathrm{A}^{s} \oplus \bigoplus_{i \in \mathrm{I}} \mathrm{A} /\left(p_{i}^{r_{i}}\right)$, where all $p_{i}$ are irreducibles, $p \in \mathrm{~A}^{*}$ an irreducible element and $r>0$. Then $p^{r-1} \mathrm{E} / p^{r} \mathrm{E}$ is isomorphic to $(\mathrm{A} /(p))^{s+d(p, r)}$, $d(p, r)$ being the number of indices $i$ such that $\chi\left(p_{i}\right)=\chi(p)$ and $r_{i} \geqslant r$.
Proof of the Lemma The general case follows by passing to the direct sum in the following particular cases:
(i) If $\mathrm{E}=\mathrm{A}, p^{r-1} \mathrm{E} / p^{r} \mathrm{E}=\mathrm{A} p^{r-1} / \mathrm{A} p^{r} \approx \mathrm{~A} / \mathrm{A} p$.
(ii) If $\mathrm{E}=\mathrm{A} /\left(p^{\prime r^{\prime}}\right)$, where $\chi\left(p^{\prime}\right) \neq \chi(p)$, there exist $u, v \in \mathrm{~A}$ such that $u p+$ $v p^{\prime r^{\prime}}=1$. So the image of $p$ in the quotient ring $\mathrm{A} /\left(p^{\prime r^{\prime}}\right)$ is invertible, and $p^{r-1} \mathrm{E}=$ $p^{r} \mathrm{E}=\mathrm{E}$, and soù $p^{r-1} \mathrm{E} / p^{r} \mathrm{E}=0$.
(iii) If $\mathrm{E}=\mathrm{A} /\left(p^{r^{\prime}}\right)$ with $r^{\prime} \geqslant r$, then $p^{r} \mathrm{E}=\mathrm{A} p^{r} / \mathrm{A} p^{r^{\prime}}$ and $p^{r-1} \mathrm{E}=\mathrm{A} p^{r-1} / \mathrm{A} p^{r^{\prime}}$, and so $p^{r-1} \mathrm{E} / p^{r} \mathrm{E} \approx \mathrm{A} p^{r-1} / \mathrm{A} p^{r} \approx \mathrm{~A} / \mathrm{A} p$.
(iv) If $\mathrm{E}=\mathrm{A} /\left(p^{r^{\prime}}\right)$ with $r^{\prime}<r$, then $p^{r-1} \mathrm{E}=p^{r} \mathrm{E}=0$, and so $p^{r-1} \mathrm{E} / p^{r} \mathrm{E}=0$.

Proof of the Proposition (b) The quotient module of E by its torsion submodule is isomorphic to $\mathrm{A}^{s}$ and to $\mathrm{A}^{s^{\prime}}$. Thus $s=s^{\prime}$ by the uniqueness of dimension (3.4.5). Defining $d^{\prime}(p, r)$ as $d(p, r), s^{\prime}+d^{\prime}(p, r)=s+d(p, r)$ : it is the dimension of $p^{r-1} \mathrm{E} / p^{r} \mathrm{E}$ as a vector space over $\mathrm{A} /(p)$. Hence $d^{\prime}(p, r)=d(p, r)$.

The number $f(p, r)$ of indices $i \in \mathrm{I}$ such that $\chi\left(p_{i}\right)=\chi(p)$ and $r_{i}=r$ is $d(p, r)-d(p, r+1)$; The same holds for the number $f^{\prime}(p, r)$ defined likewise. Henceù $f(p, r)=f^{\prime}(p, r)$.
Remark In part (b) of the proposition, the submodules of E corresponding to $\mathrm{A} /\left(p_{i}^{r_{i}}\right)$ and to $\mathrm{A} /\left(p_{\tau(i)}^{r_{\tau(i)}}\right)$ are isomorphic, but may be distinct.

### 3.5.11 Primary Decomposition

Let E be a torsion A -module and $p \in \mathrm{~A}^{*}$ an irreducible element. The $p$-primary component of E , written $\mathrm{E}_{p}$, is the set of $x \in \mathrm{E}$ such that $(\exists r \in \mathbb{N}) p^{r} x=0$. It is
a submodule of E . Then $\mathrm{E}_{p}=\mathrm{E}_{p^{\prime}}$ if $\chi(p)=\chi\left(p^{\prime}\right)$. Hence we can talk of $\mathrm{E}_{p}$ for irreducible $p \in \operatorname{Mon}^{*}(\mathrm{~A})$. When $\mathrm{E}_{p}=\mathrm{E}, \mathrm{E}$ is said to be $p$-primary.

If F is a submodule of E , then $\mathrm{F}_{p}=\mathrm{F} \cap \mathrm{E}_{p}$. If E and $\mathrm{E}^{\prime}$ are torsion A-modules and $\varphi: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ is a homomorphism, then $\varphi\left(\mathrm{E}_{p}\right) \subset \mathrm{E}_{p}^{\prime}$.

Theorem Let E be a torsion A-module. Then $\mathrm{E}=\bigoplus_{p \in \mathrm{M}} \mathrm{E}_{p}$, where M is the set of irreducible elements of Mon*(A).

Proof Let $x \in \mathrm{E}$. The cyclic submodule $\mathrm{A} x$ of E is isomorphic to a $\mathrm{A} /(a)$, where $a \in$ $\mathrm{A}^{*}$ (3.3.6, Example). Write $a=u \cdot p_{1}^{r_{1}} \cdots p_{k}^{r_{k}}$ as in (3.5.9, Corollary); then A/ $\left.a\right) \approx$ $\mathrm{A} /\left(p_{1}^{r_{1}}\right) \oplus \cdots \oplus \mathrm{A} /\left(p_{k}^{r_{k}}\right)$. Set $e_{i}=\left(0, \ldots, \chi_{p_{i}^{r_{i}}}(1), \ldots, 0\right)$ in $\mathrm{A} /\left(p_{1}^{r_{1}}\right) \oplus \cdots \oplus$ $\mathrm{A} /\left(p_{k}^{r_{k}}\right)$. Let $\varphi$ be an isomorphism from this direct sum onto $\mathrm{A} x$ and set $x_{i}=\varphi\left(e_{i}\right)$. Then $p_{i}^{r_{i}} x_{i}=0$. So $x_{i} \in \mathrm{E}_{p_{i}}$, and $x$ is of the form $\sum \lambda_{i} x_{i}$. Hence E is generated by the submodules $\mathrm{E}_{p_{i}}$.

Let $p_{1}, \ldots, p_{k} \in \mathrm{~A}^{*}$ be irreducible elements for which $\chi\left(p_{1}\right), \ldots, \chi\left(p_{k}\right)$ are distinct, and let $x_{1} \in \mathrm{E}_{p_{1}}, \ldots, x_{k} \in \mathrm{E}_{p_{k}}$ be elements such that $\sum x_{i}=0$. We show that $x_{1}=\cdots=x_{k}=0$. For $i \in\{1, \ldots, k\}$, let $r_{i}$ be such that $p_{i}^{r_{i}} x_{i}=0$. By Bezout's formula (3.2.11, Corollary), there exist $u, v \in \mathrm{~A}$ such that $u \cdot p_{1}^{r_{1}}+v \cdot p_{2}^{r_{2}} \cdots p_{k}^{r_{k}}=1$. Then

$$
x_{1}=\left(1-u p_{1}^{r_{1}}\right) x_{1}=v \cdot p_{2}^{r_{2}} \cdots p_{k}^{r_{k}} x_{1}=-v \cdot p_{2}^{r_{2}} \cdots p_{k}^{r_{k}}\left(x_{2}+\cdots+x_{k}\right)=0
$$

Similarly, $x_{2}=\cdots=x_{k}=0$. Hence E is a direct sum of the submodules $\mathrm{E}_{p}$.
AdDEndum. For all submodules F of $\mathrm{E}, \mathrm{F}=\bigoplus_{p \in \mathrm{M}}\left(\mathrm{F} \cap \mathrm{E}_{p}\right)$.

## Exercises 3.5. (Modules over principal ideal domains)

1. (Invariant factors) Let A be a PID and E a finitely generated A-module.
(a) Show that E is isomorphic to a module $\mathrm{A} /\left(a_{1}\right) \oplus \cdots \oplus \mathrm{A} /\left(a_{n}\right)$ with $a_{1}$ ) $\left.a_{2}\right\rangle \cdots>a_{n}$, and non-invertible $a_{n}$.
(b) Show that the sequence $\left(a_{1}, \ldots, a_{n}\right)$ is uniquely determined up to multiplication of the elements $a_{i}$ by invertible elements (these $a_{i}$ are called invariant factors). Show that $a_{1}=0$ if and only if A is torsion-free.
(c) Find the elementary divisors and invariant factors of the $\mathbb{Z}$-module $\mathbb{Z} /(4) \oplus$ $\mathbb{Z} /(8) \oplus \mathbb{Z} /(6) \oplus \mathbb{Z} /(18) \oplus \mathbb{Z} /(7)$.
(d) Show that $n$ is the minimal cardinality of a generating family of M.
(e) How can the invariant factors be computed in terms of the elementary divisors and conversely?
2. (Multiplicative groups of modulo $n$ integers) - Let $n$ be an integer. Set $\mathrm{G}(n)$ to be the finite multiplicative group of invertible elements of the ring $\mathbb{Z} /(n)$. The aim is to express $\mathrm{G}(n)$ as a product of cyclic groups.
(a) Let G be a finite commutative group, written multiplicatively. Assume that, for every prime $p$, the number of elements of G of order $p$ is less than $p$. Show that G is cyclic.
(b) Let $p$ be a prime. Show that the group $\mathrm{G}(p)$ is cyclic (use the fact that a polynomial of degree $q$ has at most $q$ roots in the field $\mathbb{Z} /(p))$.
(c) Let $\mathrm{N}\left(p^{r}\right)$ denote the set of elements of $\mathbb{Z} /\left(p^{r}\right)$ with image 1 in $\mathbb{Z} /(p)$. Show that $\mathrm{N}\left(p^{r}\right)$ is a multiplicative group (every element in $\mathrm{N}\left(p^{r}\right)$ is of the form $1+u$ withù $u$ nilpotent). What is its order?
(d) Assume that $p \neq 2$. Show that $\left(1+a p^{k}\right)^{p} \equiv 1+a p^{k+1}\left(\bmod p^{k+2}\right)$. Which are the elements of order $p$ in $\mathrm{N}\left(p^{r}\right)$ ? Show that $\mathrm{N}\left(p^{r}\right)$ is cyclic. Is the group $\mathrm{G}\left(p^{r}\right)$ cyclic?
(e) For $r \geqslant 2$, let $\mathrm{N}^{\prime}\left(2^{r}\right)$ be the subgroup of $\mathrm{N}\left(2^{r}\right)=\mathrm{G}\left(2^{r}\right)$ consisting of the elements with image 1 in $\mathbb{Z} /(4)$. Show that $\mathrm{N}^{\prime}\left(2^{r}\right)$ is cyclic, and that it is a direct factor of $\mathrm{N}\left(2^{r}\right)$. Is the group $\mathrm{G}\left(2^{r}\right)$ cyclic?
(f) Show that, if $n=p_{1}^{r_{1}} \cdots p_{k}^{r_{k}}$, the group $\mathrm{G}(n)$ is isomorphic to the product group $\mathrm{G}\left(p_{1}^{r_{1}}\right) \times \cdots \times \mathrm{G}\left(p_{k}^{r_{k}}\right)$.
(g) Express $\mathrm{G}(4761)$ as a product of cyclic groups.
3. (Fitting ideal)—Let A be a PID and E a finitely generated torsion A-module.
(a) Show that E has a presentation $\mathrm{A}^{q} \xrightarrow{f} \mathrm{~A}^{p} \rightarrow \mathrm{E} \rightarrow 0$, with $f$ injective. Show that, for such a presentation, $q=p$ necessarily holds, and that, up to multiplication by an invertible element, the determinant of $f$ only depends on E and not on the presentation. The ideal generated by $\operatorname{det}(f)$ is called the Fitting ideal of E .
(b) Compute $\operatorname{det}(f)$ for $\mathrm{E}=\mathrm{A} /\left(a_{1}\right) \oplus \cdots \oplus \mathrm{A} /\left(a_{n}\right)$. Give an interpretation for it when $\mathrm{A}=\mathbb{Z}$.
(c) Let $0 \rightarrow \mathrm{E} \rightarrow \mathrm{F} \rightarrow \mathrm{G} \rightarrow 0$ be a short exact sequence of finitely generated torsion A-modules. What is the relation between the generators of the Fitting ideals of $\mathrm{E}, \mathrm{F}$ and G ?
4. -Let $n$ be an integer and A a ring. Set $\mathrm{B}=\mathrm{A}[\mathrm{X}]$. A vector of $\mathrm{B}^{n}$ (resp. a matrix with entries in B) can be written $u_{0}+\mathrm{X} u_{1}+\cdots+\mathrm{X}^{d} u_{d}$, with $u_{i} \in \mathrm{~A}^{n}$ (resp. $u_{i}$ is a matrix with entries in A).
(A) Let M be a matrix with entries in A and consider the B -linear map $\Phi_{\mathrm{M}}=$ $\mathrm{M}-\mathrm{XI}: \mathrm{B}^{n} \rightarrow \mathrm{~B}^{n}$. Show that the A-submodule $\mathrm{A}^{n}$ of $\mathrm{B}^{n}$ is a complement of the image of $\Phi_{\mathrm{M}}$.
(B) Consider the B-module $\mathrm{E}_{\mathrm{M}}$ obtained by equipping $\mathrm{A}^{n}$ with the A-bilinear map $\mathrm{B} \times \mathrm{A}^{n} \rightarrow \mathrm{~A}^{n}$ defined by $\mathrm{X} . v=\mathrm{M}(v)$ for $v \in \mathrm{~A}^{n}$.
(a) Show that $\mathrm{E}_{\mathrm{M}}$ is not isomorphic to any submodule of $\mathrm{B}^{n}$.
(b) When A is an algebraically closed field, give an interpretation of the primary decomposition of $\mathrm{E}_{\mathrm{M}}$, and more particularly of Proposition 3.5.10.
(c) Show that there is an exact B-linear sequence (resolution of $\mathrm{E}_{\mathrm{M}}$ )

$$
0 \rightarrow \mathrm{~B}^{n} \rightarrow \mathrm{~B}^{n} \rightarrow \mathrm{E}_{\mathrm{M}}
$$

with $\Phi_{\mathrm{M}}$, and a map $\varepsilon: \mathrm{B}^{n} \rightarrow \mathrm{E}_{\mathrm{M}}$ which should be defined.
(C) Let M and $\mathrm{M}^{\prime}$ be two $n \times n$ matrices with entries in A . Show that the following conditions are equivalent:
(i) The B-modules $\mathrm{E}_{\mathrm{M}}$ and $\mathrm{E}_{\mathrm{M}^{\prime}}$ are isomorphic;
(ii) there are invertible matrices V and W with entries in B such that $\mathrm{M}^{\prime}-\mathrm{XI}=$ $\mathrm{W}(\mathrm{M}-\mathrm{XI}) \mathrm{V}^{-1}$;
(iii) there is a matrix U with entries in A such that $\mathrm{M}^{\prime}=\mathrm{UMU}^{-1}$.
(D) Show by an example (with $\mathrm{A}=\mathbb{R}, n=2$ ) that condition (iii) of (C) is not generally equivalent to the analogous condition (iii') where $\operatorname{det} \mathrm{U}=1$. What happens when A is an algebraically closed field?
5.-Let E be a finite $n$-dimensional vector space over a field K and $u$ an endomorphism of E . Endow E with the structure of a $\mathrm{K}[\mathrm{T}]$-module defined by $\mathrm{T} \cdot x=u(x)$ for $x \in \mathrm{E}$.
(a) Show that E is a finitely generated torsion $\mathrm{K}[\mathrm{T}]$-module.
(b) Show that the first invariant factor of the $\mathrm{K}[\mathrm{T}]$-module E is the minimal polynomial ${ }^{15}$ of $u$. What is the Fitting ideal of E ?
(c) Suppose that K is algebraically closed. Using 3.5 .10 recover the fact that $\mathrm{E}=$ $\mathrm{E}_{1} \oplus \cdots \oplus \mathrm{E}_{k}$, where all $\mathrm{E}_{i}$ are stable under $u$, and where, for all $i$, the endomorphism $u_{i}$ of $\mathrm{E}_{i}$ induced by $u$ is the sum of a homothety and a nilpotent endomorphism. ${ }^{16}$
(d) Show that from the $i$ th factor, the product of invariants factors is equal to the g.c.d. of minors of order $n-i$ of a matrix representing $u$.
6.-Let $E$ be the $\mathbb{Z}$-module $\mathbb{Z}^{\mathbb{N}}$ of all integers. The aim of this exercise is to determine the module $\mathrm{E}^{\top}=\operatorname{Hom}(\mathrm{E}, \mathbb{Z})$.

Let F be the submodule $\mathbb{Z}^{(\mathbb{N})}$ of E consisting of sequences with finite support, and $\left(e_{n}\right)_{n \in \mathbb{N}}$ the canonical basis for F defined by $e_{n}=\left(\delta_{n, p}\right)_{p \in \mathbb{N}}$. For all $u \in \mathrm{E}^{\top}$, define $\varphi(u) \in \mathrm{E}$ by $\varphi(u)=\left(u\left(e_{n}\right)\right)_{n \in \mathbb{N}}$. This gives a homomorphism from $\mathrm{E}^{\top}$ to E .

Parts A an B aim to find the kernel and image of $\varphi$. They are independent.

## A. Determination of $\operatorname{Ker} \varphi$

(a) Show that $\operatorname{Ker} \varphi=\mathrm{G}^{\top}$, where $\mathrm{G}=\mathrm{E} / \mathrm{F}$ (module of the "germs of sequences at infinity").
(b) Let $\mathrm{H}_{2}$ be the set of elements of G divisible by $2^{k}$ for all $k$. Show that $\mathrm{H}_{2}$ is a submodule of G. Provide a method for finding the nonzero elements of $\mathrm{H}_{2}$.
(c) Show that all linear forms $u: G \rightarrow \mathbb{Z}$ vanish on $\mathrm{H}_{2}$.
(d) Define $\mathrm{H}_{3}$ similarly. Determine $\mathrm{H}_{2}+\mathrm{H}_{3}$. Conclude.
B. Determination of $\operatorname{Im} \varphi$

For all $x=2^{a} q \in \mathbb{Z}$ with $q$ odd, set $|x|_{2}=2^{-a}$; let $|0|_{2}=0$.
(a) Show that $(x, y) \mapsto|y-x|_{2}$ is a distance on $\mathbb{Z}$. Is the space $\mathbb{Z}$ complete with respect to this distance?

Show that, if $x_{1}, \ldots, x_{n}$ are integers such that $\left|x_{i}\right|_{2}$ are mutually distinct, then $\left|\sum x_{i}\right|_{2}$ is the largest of the $\left|x_{i}\right|_{2}$.
(b) For $x=\left(x_{n}\right)_{n \in \mathbb{N}} \in \mathrm{E}$, set $\|x\|_{2}=\sup \left|x_{n}\right|_{2}$. Show that:

$$
\left(\forall u \in \mathrm{E}^{\top}\right)(\forall x \in \mathrm{E}) \quad|u(x)|_{2} \leqslant\|x\|_{2} .
$$

(c) Let $x=\left(x_{n}\right)_{n \in \mathbb{N}}$. What condition must the sequence $\left(\left\|x-\sum_{0}^{n} x_{k} e_{k}\right\|_{2}\right)_{n \in \mathbb{N}}$ satisfy so as to tend to 0 ?

[^12](d) Let $u \in \mathrm{E}^{\top}$ and set S to be the support of $\varphi(u)$, i.e. $\mathrm{S}=\left\{n \mid u\left(e_{n}\right) \neq 0\right\}$. Let $x \in \mathrm{E}$ be an element whose support is S and such that the maps
$$
s \mapsto\left|x_{s}\right|_{2} \quad \text { and } \quad s \mapsto\left|u\left(e_{s}\right)\right|_{2}\left|x_{s}\right|_{2}
$$
from $S$ to $\mathbb{R}_{+}$are strictly decreasing.
Let $A$ be the set of sequences having support in $S$ and which only take values 0 and 1 , i.e. $\mathrm{A}=\{0,1\}^{\mathrm{S}}$. For $\varepsilon \in \mathrm{A}$, set $\psi(\varepsilon)=u(\varepsilon x)$, where $\varepsilon x=\left(\varepsilon_{n} x_{n}\right)_{n \in \mathbb{N}}$. Determine $\left|\psi(\varepsilon)-\psi\left(\varepsilon^{\prime}\right)\right|_{2}$ in terms of $s_{0}=\inf \left\{s \mid \varepsilon_{s} \neq \varepsilon_{s}^{\prime}\right\}$. Deduce that $\psi: \mathrm{A} \rightarrow \mathbb{Z}$ is injective.
(e) Conclude by considering the cardinality of A.
C. Conclusion
(a) Describe $E^{\top}$.
(b) Is the module E free?
(c) Do these results generalize to infinite sets other than $\mathbb{N}$, to rings other than $\mathbb{Z}$ ?
7.-Let A be a (not necessarily principal) ring all of whose ideals are projective A-modules (for example a Dedekind ring). Let E be a free A-module, $\left(e_{i}\right)_{i \in \mathrm{I}}$ a basis for E , and F a submodule of E . Show that F is isomorphic to a module of the form $\bigoplus_{i \in \mathrm{I}} \mathrm{J}_{i}$, where $\mathrm{J}_{i}$ is an ideal of A for each $i$ (use the Proof of Theorem 3.5.1).
8.-Let G be a countable $\mathbb{Z}$-module whose points are separated by linear forms (i.e. $(\forall x \neq 0)(\exists \varphi: \mathrm{G} \rightarrow \mathbb{Z}$ linear $) \varphi(x) \neq 0)$. The aim is to show that G is free. This exercise uses some of the results of 3.8.
(a) Show that $G$ is torsion-free. Set $V=\mathbb{Q} \otimes G$. Show that $G$ is embedded in V , and that V is a finite dimensional or countable vector space. Show that $\mathrm{G}^{\top}=$ $\operatorname{Hom}(\mathrm{G} ; \mathbb{Z})$ is embedded in $\mathrm{V}^{*}=\operatorname{Hom}(\mathrm{V} ; \mathbb{Q})$.
(b) Suppose $V$ is finite dimensional. Show that there is a family of elements of $G^{\top}$ forming a basis for $\mathrm{V}^{*}$, and that such a family separates the points of G . Deduce that G is embedded in a finitely generated $\mathbb{Z}$-module, and that it itself is finitely generated.
(c) Still supposing V to be finite dimensional, let $\mathrm{V}^{\prime}$ be a hyperplane in V . Set $\mathrm{G}^{\prime}=\mathrm{G} \cap \mathrm{V}^{\prime}$. Show that any basis for $\mathrm{G}^{\prime}$ can be extended to a basis for G .
(d) Suppose V is of countably infinite dimension. Let $\left(e_{1}, \ldots, e_{n}, \ldots\right)$ be a basis for V , and $\mathrm{V}_{n}$ the vector subspace generated by $\left(e_{1}, \ldots, e_{n}\right)$, and set $\mathrm{G}_{n}=\mathrm{V}_{n} \cap \mathrm{G}$. Construct by induction a sequence $\left(x_{n}\right)$ of elements of G such that $\left(x_{1}, \ldots, x_{n}\right)$ is a $\mathbb{Z}$-base for $\mathrm{G}_{n}$ for all $n$. Show that $\left(x_{1}, \ldots, x_{n}, \ldots\right)$ is then a basis for G .
(e) Deduce that every countable subgroup of $\mathbb{Z}^{\mathbb{N}}$ is free.
9.-Let $G \subset \mathbb{Z}^{\mathbb{N}}$ be a countable subgroup that is not finitely generated. Show that G is not a direct factor (it may be shown that there is element $\xi \neq 0$ in $\mathbb{Z}^{\mathbb{N}} / \mathrm{G}$ with infinitely many divisors).
10.-The aim is to find an example of a torsion-free countable $\mathbb{Z}$-module G all whose nonzero elements are only divisible by finitely many scalars, but which is not free.

Let $p$ be a prime and A the ring $\mathbb{Z}\left[\frac{1}{p}\right]=\left\{\frac{m}{p^{r}}\right\}_{m \in \mathbb{Z}, r \in \mathbb{N}}$; Denote by $\widehat{\mathbb{Z}}_{p}$ the ring of $p$-adic integers (see 2.5 , Exercise 7) and by $\widehat{\mathbb{Q}}_{p}$ its field of fractions.

Let $\alpha \in \widehat{\mathbb{Z}}_{p}-\mathbb{Q}$, and G the subgroup of $\mathrm{A}^{2}$ consisting of pairs $(x, y)$ such that $|y-\alpha x|_{p} \leqslant 1$ (absolute $p$-adic value).
(a) Show that the projection $(x, y) \mapsto x$ from G to A is surjective and that its kernel can be identified with $\mathbb{Z}$. Deduce that $G$ is countable, but not finitely generated.
(b) Show that G is embedded in $\mathbb{Q}^{2}$. Deduce that $G$ is not free.
(c) Show that the nonzero elements of G are only divisible by finitely many powers of $p$, and only have finitely many divisors.

### 3.6 Noetherian Rings

### 3.6.1

Theorem and Definition Let A be a ring. The following properties are equivalent:
(i) every ideal of A is finitely generated;
(ii) the set of ideals of A , ordered by inclusion, is Noetherian (i.e. every increasing sequence of ideals of A is stationary);
(iii) every submodule of a finitely generated A-module is finitely generated;
(iv) For any finitely generated A-module E, the ordered set of submodules of E is Noetherian.

If these conditions hold, the ring A is said to be Noetherian.
Implications (iii) $\Rightarrow$ (i) and (iv) $\Rightarrow$ (ii) are immediate (take $\mathrm{E}=\mathrm{A}$ ). The equivalences (i) $\Leftrightarrow$ (ii) and (iii) $\Leftrightarrow$ (iv) follow from the next lemma:

Lemma Let E be an $\mathrm{A}-m o d u l e$. The following conditions are equivalent:
( $i_{\mathrm{E}}$ ) every submodule of E is finitely generated;
(ii $\mathrm{E}_{\mathrm{E}}$ ) the ordered set of submodules of E is Noetherian.
Proof of the Lemma $\left(\mathrm{i}_{\mathrm{E}}\right) \Rightarrow$ (iii $)$. Let $\left(\mathrm{F}_{n}\right)$ be an increasing sequence of submodules of E . Set $\mathrm{F}=\bigcup \mathrm{F}_{n}$. Let $\left(x_{i}\right)_{i \in \mathrm{I}}$ be a finite generating family of F . For each $i$, there exists $n_{i}$ such that $x_{i} \in \mathrm{~F}_{n_{i}}$. Thus, $\mathrm{F}_{n}=\mathrm{F}$ for $n=\sup \left(n_{i}\right)$, and so the sequence $\left(\mathrm{F}_{n}\right)$ is stationary.
$\left(\mathrm{ii}_{\mathrm{E}}\right) \Rightarrow\left(\mathrm{i}_{\mathrm{E}}\right)$. Let F be a submodule of E . The set X of finitely generated submodules of $F$ is not empty since $\{0\} \in X$. So it contains a maximal element $G$ (1.6.1), and $\mathrm{G}=\mathrm{F}$. Indeed were $x$ in $\mathrm{G}-\mathrm{F}$, then G would be strictly contained in the submodule $G+A x$. Hence $F$ is finitely generated.

Proof of the Theorem (i) $\Rightarrow$ (iii) remains to be shown. We show by induction on $n$ that every submodule of an A-module E generated by $n$ elements is finitely generated. This is obvious for $n=0$. It holds for $n=1$. Indeed, a cyclic A-module is isomorphic to a quotient module $A / I$, and submodules of $A / I$ are of the form $J / I$, where $J$ is an ideal of $A$, and is finitely generated if $J$ is a finitely generated ideal. Let $E$ be an A-module generated by $x_{1}, \ldots, x_{n}$ with $n \geqslant 1$, and F a submodule of E . Let $\mathrm{E}^{\prime}$ be the submodule of E generated by $\left(x_{1}, \ldots, x_{n-1}\right), \mathrm{E}^{\prime \prime}=\mathrm{E} / \mathrm{E}^{\prime}, \mathrm{F}^{\prime}=\mathrm{F} \cap \mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime \prime}=\mathrm{F} / \mathrm{F}^{\prime}$.

The module $\mathrm{F}^{\prime \prime}$ can be identified with a submodule of $\mathrm{E}^{\prime \prime}$, and so is finitely generated, and by the induction hypothesis, $\mathrm{F}^{\prime}$ is finitely generated. This gives an exact sequence $0 \rightarrow \mathrm{~F}^{\prime} \rightarrow \mathrm{F} \rightarrow \mathrm{F}^{\prime \prime} \rightarrow 0$, and so F is finitely generated by (3.3.9, Corollary).

### 3.6.2 Examples

(1) Every PID is Noetherian.
(2) Every quotient ring of a Noetherian ring is Noetherian, Indeed, if $B=A / I$, the ideals of B correspond to ideals of A containing I .
(3) Let $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism. If A is Noetherian, and if $\varphi$ turns $B$ into a finitely generated A-module (see 3.3.1, Example 3), then the ring B is Noetherian. Indeed, every ideal in B is finitely generated over A, perforce over B. This will be generalized (3.7.7, Corollary 3.5).
(4) If K is a field, the ring of polynomials $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ in $n$ variables will be seen (3.7.7) to be Noetherian. The same holds for the ring $\mathrm{K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$ of formal series, and for the ring of convergent series $\mathbb{C}\left\{\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right\}$.
(5) Every ring of fractions of a Noetherian ring is Noetherian. In particular every localization of a Noetherian ring is Noetherian.
(6) The ring of analytic functions on a compact $\mathbb{R}$-analytic manifold is Noetherian.

### 3.6.3 Counterexamples

(1) Let K be a field and I an infinite set. The ring $\mathrm{K}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ is not Noetherian: the ideal of polynomials with constant terms is not finitely generated.
(2) If $E$ is an infinite dimensional Banach space, the ring of germs of analytic functions on E is not Noetherian: if $\left(\mathrm{F}_{n}\right)$ is a strictly decreasing sequence of closed vector subspaces of finite codimension in E , then we get a strictly increasing sequence $\left(\mathrm{I}_{n}\right)$ of ideals, where $\mathrm{I}_{n}$ is the ideal of germs of analytic functions vanishing on $\mathrm{F}_{n}$ in the neighbourhood of 0 .
(3) A subring of a Noetherian ring is not necessarily Noetherian. Indeed, a nonNoetherian integral domain is a subring of its field of fractions.
(4) The ring $C([0,1])$ of continuous functions on $[0,1]$ is not Noetherian: the ideal $\mathfrak{m}$ of functions $f \in \mathrm{C}([0,1])$ such that $f(0)=0$ is not finitely generated. Indeed, let $f_{1}, \ldots, f_{n} \in \mathfrak{m}$, and X be the set of $x \in[0,1]$ such that $f_{1}(x)=\cdots=f_{n}(x)=0$. If $\mathrm{X} \neq\{0\}$, the function $x \mapsto x$ is not in the ideal generated by $\left(f_{1}, \ldots, f_{n}\right)$; if $\mathrm{X}=\{0\}$, every function in the ideal generated by $\left(f_{1}, \ldots, f_{n}\right)$ is $\mathrm{O}(h)$, where $h=\sup \left|f_{i}\right|$, and the function $\sqrt{h}$ is not in this ideal.

The ring of germs at 0 of continuous functions on $\mathbb{R}$ can likewise be shown not be Noetherian.
(5) The ring $\mathrm{E}_{0}$ of germs at 0 of $\mathrm{C}^{\infty}$-functions on $\mathbb{R}$ is not Noetherian: the sequence of ideals $\left(f_{n}\right)$, where $f_{n}(x)=\frac{1}{x^{n}} e^{-1 / x^{2}}$, is strictly increasing.
(6)The ring $\mathcal{O}(\mathbb{R})$ of analytic functions on $\mathbb{R}$ is not Noetherian: the sequence of ideals $\left(\mathrm{I}_{n}\right)$, where $\mathrm{I}_{n}$ is the ideal of functions $f \in \mathcal{O}(\mathbb{R})$ such that $f(p)=0$ for all integers $p \geqslant n$, is strictly increasing. Indeed the function $x \mapsto \frac{\sin \pi x}{x-n}$ is in $\mathrm{I}_{n+1}$ but not in $\mathrm{I}_{n}$.
(7) An infinite dimensional Banach algebra over $\mathbb{R}$ or $\mathbb{C}$ is never a Noetherian ring (3.6, Exercise 3). However, there are Noetherian and even principal infinite dimensional Banach algebras over ultrametric fields (3.6, Exercise 4).

### 3.6.4 Minimal Prime Ideals

Proposition Let A be a ring.
(a) Every prime ideal contains a minimal prime ideal..
(b) The intersection of minimal prime ideals is the nilradical of A .
(c) If A is Noetherian, it only has finitely many minimal prime ideals.

Proof (a) The set of prime ideals of A is codirect (i.e. direct for the reverse inclusion order). Indeed, if $\left(\mathfrak{p}_{i}\right)$ is a totally ordered family of prime ideals, then $\mathfrak{p}=\bigcap \mathfrak{p}_{i}$ is prime since $\mathrm{A} / \mathfrak{p}$ can be identified with a subring of $\lim \mathrm{A} / \mathfrak{p}_{i}$, and so is integral.
(More fundamentally, if $x \notin \mathfrak{p}$ and $y \notin \mathfrak{p}$, there exists $i$ such that $x \notin \mathfrak{p}_{i}$ and $y \notin \mathfrak{p}_{i}$, and so $x y \notin \mathfrak{p}_{i}$. Hence $x y \notin \mathfrak{p}$.) By Zorn's theorem, every prime ideal contains a minimal prime ideal.
(b) By (a), the intersection of minimal prime ideals is equal to the intersection of all prime ideals.
(c) The ordered set E of ideals in A is Noetherian. For any ideal I of A, let $\mathscr{M}$ (I) be the set of minimal prime ideals containing I and let us show by Noetherian induction (1.6.3) that, for all I, the set $\mathscr{M}(\mathrm{I})$ is finite. Let $\mathrm{I} \in \mathrm{E}$ be such that, for any ideal J strictly containing I , the set $\mathscr{M}(\mathrm{J})$ is finite. If I is prime, then $\mathscr{M}(\mathrm{I})$ has at most one element. If I is not prime, let $x, y \notin \mathrm{I}$ such that $x y \in \mathrm{I}$, and set $\mathrm{J}^{\prime}=\mathrm{I}+\mathrm{A} x$ and $\mathrm{J}^{\prime \prime}=\mathrm{I}+\mathrm{A} y$. Every prime ideal containing I contains $x y$, and hence contains $\mathrm{J}^{\prime}$ or $\mathrm{J}^{\prime \prime}$; thus $\mathscr{M}(\mathrm{I})=\mathscr{M}\left(\mathrm{J}^{\prime}\right) \cup \mathscr{M}\left(\mathrm{J}^{\prime \prime}\right)$ is finite.

### 3.6.5 Associated Prime Ideals

Definition Let A be a ring and M an A-module. The annihilator of an element of M is called an associated prime ideal of M .

The set of all associated prime ideals associated of $M$ is denoted by Ass(M). For $x \in \mathrm{M}$, write $\operatorname{Ann}(x)$ for the annihilator of $x$.

Proposition Let A be a Noetherian ring and M a nonzero A-module. Then Ass(M) is not empty.

The set of ideals $\operatorname{Ann}(x)$ with $x \in \mathrm{M}, x \neq 0$, ordered by inclusion, is Noetherian. Hence it has a maximal element $\operatorname{Ann}\left(x_{0}\right)$. We show by contradiction that $\operatorname{Ann}\left(x_{0}\right)$ is prime. Consider $a$ and $b$ not in $\operatorname{Ann}\left(x_{0}\right)$ such that $a b \in \operatorname{Ann}\left(x_{0}\right)$. Then $\operatorname{Ann}\left(x_{0}\right) \subset$ $\operatorname{Ann}\left(b x_{0}\right), a \in \operatorname{Ann}\left(b x_{0}\right)$ and $a \notin \operatorname{Ann}\left(x_{0}\right)$; so $\operatorname{Ann}\left(b x_{0}\right)$ strictly contains $\operatorname{Ann}\left(x_{0}\right)$, contradicting maximality.

### 3.6.6

Proposition Let A be a ring and $0 \rightarrow \mathrm{E} \xrightarrow{u} \mathrm{~F} \xrightarrow{v} \mathrm{G} \rightarrow 0$ an exact sequence of A modules. Then $\operatorname{Ass}(\mathrm{F}) \subset \operatorname{Ass}(\mathrm{E}) \cup \operatorname{Ass}(\mathrm{G})$.

Proof Let $\mathfrak{p}=\operatorname{Ann}(x) \in \operatorname{Ass}(\mathrm{F})$. We show that $\mathfrak{p} \in \operatorname{Ass}(\mathrm{E})$ or $\mathfrak{p} \in \operatorname{Ass}(\mathrm{G})$. If $x \in$ $u(\mathrm{E})$, then $\mathfrak{p} \in \operatorname{Ass}(\mathrm{E}) ;$ Suppose that $x \notin u(\mathrm{E}) ;$ so $v(x) \neq 0$. Then $\mathfrak{p} \subset \operatorname{Ann}(v(x))$. If $\mathfrak{p}=\operatorname{Ann}(v(x))$, then $\mathfrak{p} \in \operatorname{Ass}(\mathrm{G})$. Suppose that $\mathfrak{p} \neq \operatorname{Ann}(v(x))$. Hence there exists $a \in \mathrm{~A}$ such that $a \in \operatorname{Ann}(v(x))$ and $a \notin \mathfrak{p}$, i.e. $a v(x)=0$ and $a x \neq 0$. We show that $\mathfrak{p}=\operatorname{Ann}(a x)$. It is obvious that $\mathfrak{p}=\operatorname{Ann}(x) \subset \operatorname{Ann}(a x)$; we prove the converse inclusion. Let $b \in \mathrm{~A}$ such that $b a x=0$, then $b a \in \mathfrak{p}$ and $a \notin \mathfrak{p}$, and so $b \in \mathfrak{p}$; hence $\operatorname{Ann}(a x) \subset \mathfrak{p}$.

### 3.6.7

Theorem Let A be a Noetherian ring and M a finitely generated A-module. Then Ass(M) is finite.

Proof We first prove the theorem when M is cyclic, i.e. $\mathrm{M}=\mathrm{A} / \mathrm{I}$ for some ideal I , by Noetherian induction on I (1.6.3). If I is prime then, for all $x \neq 0$ in $M$, the annihilator of $x$ in $\mathrm{A} / \mathrm{I}$ is the zero set and so the annihilator of $x$ in A is I and $\operatorname{Ass}(\mathrm{M})=\{\mathrm{I}\}$ contains a unique element, and so is finite. Suppose that $I$ is not prime and that, for all ideals I' strictly containing I, the set $\operatorname{Ass}\left(\mathrm{A} / \mathrm{I}^{\prime}\right)$ is finite. Let $a, b \in \mathrm{~A}-\mathrm{I}$ be such that $a b \in \mathrm{I}$; set $\mathrm{I}^{\prime}=\mathrm{I}+\mathrm{A} a$ and let $\mathrm{I}^{\prime \prime}$ be the ideal of $c \in \mathrm{~A}$ such that $a c \in \mathrm{I}$. As I $\mp \mathrm{I}^{\prime}$, $\mathrm{I} q \mathrm{I}^{\prime \prime}$ since $b \in \mathrm{I}^{\prime \prime}$. There is an exact sequence $0 \rightarrow \mathrm{~A} / \mathrm{I}^{\prime \prime} \xrightarrow{u} \mathrm{~A} / \mathrm{I} \xrightarrow{v} \mathrm{~A} / \mathrm{I}^{\prime} \rightarrow 0$, where $u$ is multiplication by $a$. Therefore, $\operatorname{Ass}(\mathrm{A} / \mathrm{I})$ is finite by Proposition 3.6.6.

Next suppose that $x_{1}, \ldots, x_{r}$ generate M . The proof is by induction on $r$. Let $\mathrm{M}^{\prime}$ be the submodule of M generated by $x_{1}, \ldots, x_{r-1}$. By the induction hypothesis, the set $\operatorname{Ass}\left(M^{\prime}\right)$ is finite and $\operatorname{Ass}\left(M / M^{\prime}\right)$ is finite since $M / M^{\prime}$ is cyclic. Hence $\operatorname{Ass}(M)$ is finite by Proposition 3.6.6.

### 3.6.8 Remark

Let A be a Noetherian ring, M a finitely generated A -module, I the annihilator of M , i.e. $\mathrm{I}=\bigcap_{x \in \mathrm{M}} \operatorname{Ann}(x)$. The inverse image of every minimal prime ideal of $\mathrm{A} / \mathrm{I}$ is an element of Ass(M).

## Exercises 3.6. (Noetherian rings)

1.-Let U be a non-empty open subset of $\mathbb{C}$. The aim is to show that the ring $\mathcal{O}(\mathrm{U})$ of holomorphic functions on $U$ is not Noetherian.
(a) Show that there is a holomorphic function on $U$ with infinitely many simple roots. (Study the case where U contains $\mathbb{R}_{+}$and is a neighbourhood of zero not containing $\mathbb{R}_{+}$.)
(b) Let $f \in \mathcal{O}(\mathrm{U})$ admit infinitely many simple roots $a_{n}$. Denote by $\mathrm{I}_{n}$ the set $h \in \mathcal{O}(\mathrm{U})$ such that $h\left(a_{k}\right)=0$ for all $k \geqslant n$. Show that $\left(\mathrm{I}_{n}\right)$ is a strictly increasing sequence of ideals.
2.-For every open subset U of $\mathbb{C}^{2}$, let $\mathcal{O}(\mathrm{U})$ be the ring of holomorphic functions on U and for every compact subset K of $\mathbb{C}^{2}, \mathcal{O}(\mathrm{~K})=\underset{\longrightarrow}{\lim } \mathcal{O}(\mathrm{U})$ for open neighbourhoods $U$ in $K$. Denote by $D$ the closed unit disc of $\mathbb{C}$.
$\mathcal{O}(\mathrm{D} \times \mathrm{D})$ may be shown to be Noetherian. ${ }^{17}$ The aim is to show that there are convex compact subsets $K_{1}$ and $K_{2}$ in $\mathbb{C}$ such that $\mathcal{O}\left(K_{1} \times K_{2}\right)$ is not Noetherian.

Let $\mathrm{K}_{1} \subset \mathrm{D}$ be a convex compact subset whose intersection with the boundary of D consists of 1 and of the points $a_{n}=e^{\frac{1}{2 n} i}$ and $\mathrm{K}_{2}$ be the image of $\mathrm{K}_{1}$ under conjugation $z \mapsto \bar{z}$. Set $\mathrm{K}=\mathrm{K}_{1} \times \mathrm{K}_{2}$.
(a) Describe the intersection of K with $\mathrm{H}=\left\{(x, y) \in \mathbb{C}^{2} \mid x y=1\right\}$.
(b) Set

$$
b_{n}=e^{i / 2 n+1}, \quad u_{n}(z)=-i \frac{z-b_{n}}{z+b_{n}}, \quad v_{n}(z)=i \frac{z-\overline{b_{n}}}{z+\overline{b_{n}}}
$$

Show that $u_{n}(z)\left(\right.$ resp. $\left.v_{n}(z)\right)$ has an imaginary part $\geqslant 0$ (resp. $\left.\leqslant 0\right)$ for $z \in \mathrm{~K}_{1}$ (resp. $z \in \mathrm{~K}_{2}$ ).

Show that $u_{n}\left(a_{i}\right)=v_{n}\left(\overline{a_{i}}\right) \in \mathbb{R}$; what is the sign of this number?
(c) Define $h_{n} \in \mathcal{O}\left(\mathrm{~K}_{1} \times \mathrm{K}_{2}\right)$ by $h_{n}(x, y)=\frac{1}{2 i \pi}\left(\log v_{n}(y)-\log u_{n}(x)\right)$, log being chosen so that that the imaginary parts are in $[0,2 \pi]$. Calculate $h_{n}\left(a_{k}, \overline{a_{k}}\right)$. Conclude as in the previous exercise.
3.-Let A a commutative Banach algebra over $\mathbb{R}$. The aim is to show that, if A is a Noetherian ring, then it is finite dimensional as a vector space over $\mathbb{R}$.

Assume the following results:
Banach Theorem. Let E and F be two Banach spaces, $f: \mathrm{E} \rightarrow \mathrm{F}$ a continuous bijective linear map. Then $f^{-1}$ is continuous.

Gelfand-Mazur Theorem. Let A be a Banach algebra over the field $\mathbb{R}$. Then A is isomorphic to $\mathbb{R}$ or $\mathbb{C}(4.4$, Exercise 5$)$.
(a) Let E and F be Banach spaces, $f: \mathrm{E} \rightarrow \mathrm{F}$ an injective closed morphism. Show that there exists $c>0$ such that every $f^{\prime}: \mathrm{E} \rightarrow \mathrm{F}$ with $\left\|f^{\prime}-f\right\| \leqslant c$ is injective and

[^13]closed. Show that if $f$ is not surjective, then this is also the case for all $f^{\prime}$ sufficiently near $f$ (extend $f$ to a closed injective morphism $g: \mathrm{E} \oplus \mathbb{R} \rightarrow \mathrm{F}$ ).
(b) Let E and F be Banach spaces and $f: \mathrm{E} \rightarrow \mathrm{F}$ a surjective morphism. Show that any morphism sufficiently near $f$ is surjective.
(c) Let A be a Banach algebra, E a Banach A-module and F a submodule of E . Show that if the closure $\bar{F}$ of $F$ is a finitely generated $A$-module, then $F$ is closed (find elements of F approaching the generators of $\overline{\mathrm{F}}$; these elements generate $\overline{\mathrm{F}}$ ). Deduce that if A is a Noetherian ring, then every ideal is closed.
(d) Show that every integral Noetherian Banach algebra is a field (consider the map $a \mapsto a \mathrm{I}$ from $\mathrm{A}-\{0\}$ to $\mathrm{L}(\mathrm{A}, \mathrm{A})$ and show that its image is contained in the set of automorphisms of $A$ regarded as a vector space by noting that $A-\{0\}$ is connected if $\mathrm{A} \neq \mathbb{R}$ ).
(e) Prove the result when A is reduced (use Proposition 3.6.4).
(f) Let A be a Noetherian ring and $\mathfrak{N}$ its nilradical. Show that there exists $r$ such that $\mathfrak{N}^{r}=0$. Show that the quotients $\mathfrak{N}^{k} / \mathfrak{N}^{k+1}$ are finitely generated $\mathrm{A} / \mathfrak{N}$-modules. Deduce that, if A is a $\mathbb{R}$-algebra such that $\mathrm{A} / \mathfrak{N}$ is a finite dimensional vector space, then A is finite dimensional. Conclude.
4.-Let K be a complete ultrametric valued field (i.e. equipped with a map $x \mapsto|x|$ from K to $\mathbb{R}_{+}$satisfying $|x y|=|x||y|,|1|=1,|x+y| \leqslant \sup (|x|,|y|)$, and complete with respect to the distance $d(x, y)=|x-y|)$. Denote by $\mathrm{K}\{\mathrm{X}\}$ the subring of $\mathrm{K}[[\mathrm{X}]]$ consisting of the formal series $\sum a_{n} \mathrm{X}^{n}$ such that the sequence $\left(\left|a_{n}\right|\right)$ tends to 0 .
(a) For $f=\sum a_{n} \mathrm{X}^{n} \in \mathrm{~K}\{\mathrm{X}\}$, set $\|f\|=\sup \left(\left|a_{n}\right|\right)$. Show that this defines a norm on $\mathrm{K}\{\mathrm{X}\}$ with respect to which this vector space is complete, and that multiplication is continuous.

The aim is to prove that $\mathrm{K}\{\mathrm{X}\}$ is a PID.
(b) Let $d \in \mathbb{N}$, and $\mathrm{E}_{d}$ the vector space of polynomials of degree $<d$ with coefficients in K. Equip this space with the norm $\|f\|=\sup \left(\left|a_{i}\right|\right)$ for $f=\sum a_{i} \mathrm{X}^{i}$, and the space $\mathrm{E}_{d} \oplus \mathrm{~K}\{\mathrm{X}\}$ with the norm defined by $\|(\mathrm{R}, f)\|=\sup (\|\mathrm{R}\|,\|f\|)$. Let $\mathrm{P} \in$ $\mathrm{K}[\mathrm{X}]$ be a monic polynomial of degree $d$. Show that the map $\varphi:(\mathrm{R}, f) \mapsto \mathrm{P} f+\mathrm{R}$ is an isometry from $\mathrm{E}_{d} \oplus \mathrm{~K}\{\mathrm{X}\}$ onto $\mathrm{K}\{\mathrm{X}\}$ first show that $\varphi$ is isometric, then that its image is dense, noting that the initial space is complete).
(c) Let $h=\sum a_{n} \mathrm{X}^{n}$ be a nonzero element of $\mathrm{K}\{\mathrm{X}\}$, and $d$ the largest integer $i$ such that $\left|a_{i}\right|=\|h\|$. Show that $\psi:(\mathrm{R}, f) \mapsto h f+\mathrm{R}$ is an isomorphism from $\mathrm{E}_{d} \oplus \mathrm{~K}\{\mathrm{X}\}$ onto $\mathrm{K}\{\mathrm{X}\}$ (reduce to the case $\left|a_{d}\right|=1$, where $\psi=\varphi+u$ for some isometry $\varphi$ and $u$ with $\|u\|<1$ ).
(d) Show that $K\{X\}$ is Euclidian.
5.-Let E be a finite dimensional vector space over $\mathbb{C}$, and A the ring $\mathcal{O}_{\mathrm{E}}$ of germs at the origin of holomorphic functions on E . Denote by $\mathfrak{m}$ the maximal ideal of A consisting of $f$ such that $f(0)=0$. For any linear form $u$ on E , let $u^{*}$ be the homomorphism $f \mapsto f \circ u$ from $\mathbb{C}\{\mathrm{X}\}=\mathcal{O}_{\mathrm{C}}$ to A , and for all A-modules $\mathrm{M}, u^{*} \mathrm{M}$ the $\mathbb{C}\{X\}$-module obtained by restricting scalars from A to $\mathbb{C}\{\mathrm{X}\}$ along $u^{*}$. Assume that the ring A is Noetherian.

Let M be a finitely generated A-module, N the submodule of M consisting of elements $x$ such that $(\exists r) \mathfrak{m}^{r} x=0$. Show that there is a linear form $u$ on E such that the torsion submodule of $u^{*} \mathrm{M}$ is N (reduce to the case $\mathrm{N}=0$, and note that then for all $\mathfrak{p} \in \operatorname{Ass}(M)$, the set of $u \in \mathrm{E}^{\top}$ in $\mathfrak{p}$ is a strict vector subspace, and that a finite union of such subspaces cannot be the whole of $\mathrm{E}^{\top}$ ).
6. (Fractional ideals)—Let A be an integral domain and K its field of fractions. The fractional ideal of A (or by abuse of language K ) is an A-submodule of K contained in a cyclic A-submodule of K .
(a) Let I be a fractional ideal of A and $f: \mathrm{I} \rightarrow \mathrm{K}$ an A-linear map. Show that $f$ is of the form: $x \mapsto a x$, with $a \in \mathrm{~K}$. In particular $f$ is injective or null, and its image is a fractional ideal.
(b) Show that all fractional ideals are finitely generated projective A-modules.
(c) Let I be a fractional ideal. It is said to be invertible if there is a fractional ideal $I^{\prime}$ such that $I \cdot I^{\prime}=A$. Show that the following conditions are equivalent:
(i) I is invertible;
(ii) there are $x_{1}, \ldots, x_{r} \in \mathrm{I}$ and $y_{1}, \ldots, y_{r} \in \mathrm{~K}$ such that $\mathrm{I} \cdot y_{1} \subset \mathrm{~A}, \ldots, \mathrm{I} \cdot y_{r} \subset \mathrm{~A}$, and $x_{1} y_{1}+\cdots+x_{r} y_{r}=1$;
(iii) I is a projective A-module.
(d) Let $I_{1}$ and $I_{2}$ be two ideals of $A$. $I_{1}$ is said to divide $I_{2}$ if there is an ideal $J$ of A such that $\mathrm{I}_{2}=\mathrm{J} \cdot \mathrm{I}_{1}$. Show that, if $\mathrm{I}_{1}$ divides $\mathrm{I}_{2}$, then $\mathrm{I}_{1} \supset \mathrm{I}_{2}$. Give an example for which the converse does not hold. Show that, if $I_{1}$ is invertible, $I_{1}$ divides $I_{2}$ if and only if $\mathrm{I}_{1} \supset \mathrm{I}_{2}$.
7. (Dedekind rings)-This exercise follows on the previous one.

A Dedekind ring is an integral domain A in which every ideal is a projective A-module.
(a) Show that all Dedekind rings are Noetherian.
(b) Let A be a Dedekind ring. Show that the multiplicative monoid of nonzero ideals in A is a free commutative monoid.
(c) Show that in a Dedekind ring, every nonzero prime ideal is maximal.
(d) Show that all PIDs are Dedekind. The ring of analytic functions of period $2 \pi$ on $\mathbb{R}$ (see 3.1, Exercise 12) is Dedekind, but not principal. All Dedekind UFDs are PIDs (3.2, Exercise 6).
(e) Let A be a Dedekind ring. For any multiplicatively closed subset $S$ in A, the ring of fractions $\mathrm{S}^{-1} \mathrm{~A}$ is Dedekind.

An integrally closed Noetherian ring in which every nonzero prime ideal is maximal can be shown to be a Dedekind ring.

### 3.7 Polynomial Algebras

Throughout this section, A denotes a (commutative unital) ring.

### 3.7.1 Monoid Algebra

Let $M$ be a multiplicative monoid. Consider the free A-module $A^{(M)}$ on a set $M$, and let $\left(e_{m}\right)_{m \in \mathrm{M}}$ be its canonical basis.
Proposition and Definition There is a unique bilinear multiplication on $\mathrm{A}^{(\mathrm{M})}$ such that $\left(\forall\left(m, m^{\prime}\right) \in \mathrm{M}^{2}\right) e_{m} \cdot e_{m^{\prime}}=e_{m m^{\prime}}$. With this multiplication, $\mathrm{A}^{(\mathrm{M})}$ is a unital associative commutative algebra if M is commutative. It is the monoid algebra M .

Proof Such a multiplication is necessarily given by:

$$
\left(\sum a_{m} e_{m}\right) \cdot\left(\sum b_{m^{\prime}} e_{m^{\prime}}\right)=\sum a_{m} b_{m^{\prime}} e_{m m^{\prime}}
$$

i.e. $\left(a_{m}\right)_{m \in \mathrm{M}} \cdot\left(b_{m^{\prime}}\right)_{m^{\prime} \in \mathrm{M}}=\left(c_{m^{\prime \prime}}\right)_{m^{\prime \prime} \in \mathrm{M}}$, where $c_{m^{\prime \prime}}=\sum_{m m^{\prime}=m^{\prime \prime}} a_{m} b_{m^{\prime}}$. The multiplication is defined by this formula is easily seen to satisfy the required conditions.

UnIVERSAL PROPERTY. For all unital associative A-algebras B and all homomorphisms $f$ from M to the multiplicative monoid underlying B , there is a unique unital algebra homomorphism: $\bar{f}: \mathrm{A}^{(\mathrm{M})} \rightarrow \mathrm{B}$ such that $(\forall m \in \mathrm{M}) \bar{f}\left(e_{m}\right)=f(m)$.
Proof It suffices to check that the linear map $\bar{f}: \mathrm{A}^{(\mathrm{M})} \rightarrow \mathrm{B}$ defined by $\bar{f}\left(e_{m}\right)=$ $f(m)$ (universal property 3.4.6) is a multiplicative homomorphism. This follows from an immediate computation.

Remark The universal property can also be stated as follows: the functor from the category of monoids (resp. commutative monoids) to the category of unital associative A -algebras (resp. unital associative commutative) assigning to a monoid M its algebra $\mathrm{A}^{(\mathrm{M})}$ is the left adjoint of the forgetful functor assigning to an algebra its underlying multiplicative monoid.

### 3.7.2 Polynomial Algebras

Definition Let $S$ be a set and $L=\mathbb{N}^{(S)}$ the free commutative monoid on $S$ (3.2.4). The polynomial algebra on S with coefficients in A , written $\mathrm{A}[\mathrm{S}]$, is the algebra $\mathrm{A}^{(\mathrm{L})}$ of the monoid L .

For $s \in \mathrm{~S}$, set $e_{s}$ to be the basis element of L corresponding to $s$ and $\varepsilon_{s}$ the element of the canonical basis for $\mathrm{A}[\mathrm{S}]=\mathrm{A}^{(\mathrm{L})}$ corresponding to $e_{s}$.
UnIVERSAL PROPERTY. For any unital associative commutative A-algebra B and any family $\left(x_{s}\right)_{s \in \mathrm{~S}}$ of elements in B , there is a unique unital algebra homomorphism $f$ from $\mathrm{A}[\mathrm{S}]$ to B satisfying $(\forall s \in \mathrm{~S}) f\left(\varepsilon_{s}\right)=x_{s}$.

In other words, the functor $\mathrm{S} \mapsto \mathrm{A}[\mathrm{S}]$ from $\mathscr{E}$ ns to the category $\mathscr{A}$ of unital associative commutative A-algebras is the left adjoint of the forgetful functor from $\mathscr{A}$ to $\mathscr{E n}$.

Proof Let $\mathscr{M}$ be the category of commutative monoids. The forgetful functor $\omega$ from $\mathscr{A}$ to $\mathscr{E} n s$ is the composite of the forgetful functor $\omega_{1}: \mathscr{A} \rightarrow \mathscr{M}$ assigning to an algebra its underlying multiplicative monoid, and of the forgetful functor $\omega_{2}$ : $\mathscr{M} \rightarrow \mathscr{E} n s$. The left adjoint of the functor $\omega_{2}$ is the functor $\lambda_{2}: S \mapsto \mathbb{N}^{(S)}$, and that of $\omega_{1}$ is $\lambda_{1}: \mathbf{M} \mapsto \mathrm{A}^{(\mathrm{M})}$; so the left adjoint of $\omega=\omega_{2} \circ \omega_{1}$ is $\lambda_{1} \circ \lambda_{2}$.

### 3.7.3 Notation

The elements of $\mathrm{A}[\mathrm{S}]$ will be denoted in one of the following manner:
1st method. For $s \in \mathrm{~S}$, identify $s$ with $\varepsilon_{s}$. The canonical basis element of $\mathrm{A}[\mathrm{S}]$ corresponding to $\left(n_{s}\right)_{s \in \mathrm{~S}}$ is then the monomial $\prod_{s \in \mathrm{~S}} s^{n_{s}}$ and all elements of $\mathrm{A}[\mathrm{S}]$ are A-linear combinations of finitely many monomials.

If, for example, $S=\{X, Y, Z\}$, then write $A[X, Y, Z]$ for $A[S]$, and so for instance $a \mathrm{X}^{3}+b \mathrm{Y}^{2} \mathrm{Z}^{4} \in \mathrm{~A}[\mathrm{X}, \mathrm{Y}, \mathrm{Z}]$ if $a, b \in \mathrm{~A}$.

2nd method. For $s \in \mathrm{~S}$, let $\mathrm{X}_{s}$ be the element $\varepsilon_{s}$ and denote $\mathrm{A}[\mathrm{S}]$ by $\mathrm{A}\left[\left(\mathrm{X}_{s}\right)_{s \in \mathrm{~S}}\right]$.
Then, for example, $a \mathrm{X}_{1}^{3}+b \mathrm{X}_{2}^{2} \mathrm{X}_{3}^{4} \in \mathrm{~A}\left[\left(\mathrm{X}_{s}\right)_{s \in \mathrm{~S}}\right]$. For $n=\left(n_{s}\right)_{s \in \mathrm{~S}} \in \mathbb{N}^{(\mathrm{S})}$, the canonical basis element of $\mathrm{A}[\mathrm{S}]$ corresponding to $n$ is often written $\mathrm{X}^{n}$. Then, $\mathrm{X}^{n}=\prod_{s \in \mathrm{~S}} \mathrm{X}_{s}^{n_{s}}$, and all elements of $\mathrm{A}\left[\left(\mathrm{X}_{s}\right)_{s \in \mathrm{~S}}\right]=\mathrm{A}[\mathrm{S}]$ can be uniquely written as $\sum_{n \in \mathbb{N}^{(S)}} a_{n} \mathrm{X}^{n}$.

The second method becomes necessary to avoid confusion when $S$ is a subset of $\mathbb{N}$, or of a set equipped with a composition law.

Let B be a unital associative commutative A-algebra and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of B . Let $\mathrm{A}\left[\left(x_{i}\right)\right]_{\mathrm{B}}$ denote the unital subalgebra of B generated by $\left(x_{i}\right)$ : it is the image of $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ under the homomorphism $f: \mathrm{A}\left[\left(\mathrm{X}_{i}\right)\right] \rightarrow \mathrm{B}$ defined by $f\left(\mathrm{X}_{i}\right)=x_{i}$. Similarly, if $x, y, z \in \mathrm{~B}$, let $\mathrm{A}[x, y, z]_{\mathrm{B}}$ denote the unital subalgebra of B generated by $x, y$ and $z$; it is the image of $\mathrm{A}[\mathrm{X}, \mathrm{Y}, \mathrm{Z}]$ under the homomorphism $f: \mathrm{A}[\mathrm{X}, \mathrm{Y}, \mathrm{Z}] \rightarrow \mathrm{B}$ defined by $f(\mathrm{X})=x, f(\mathrm{Y})=y, f(\mathrm{Z})=z$. When there is no possible confusion, write $\mathrm{A}\left[\left(x_{i}\right)\right]$ or $\mathrm{A}[x, y, z]$ instead of $\mathrm{A}\left[\left(x_{i}\right)\right]_{\mathrm{B}}$ or $\mathrm{A}[x, y, z]_{\mathrm{B}}$, even when $f$ is not injective, so that the polynomial algebra over I or over the set $\{x, y, z\}$ cannot be identified with its image in B .

For example $\mathbb{Z}[\sqrt{3}]$ can denote the unital subalgebra of $\mathbb{R}$ generated by $\sqrt{3}$, which is identified with $\mathbb{Z}[\mathrm{X}] /\left(\mathrm{X}^{2}-3\right)$. To avoid confusion, the following convention is followed: When the intention is to consider the polynomial algebra over $S$, the elements of S are called "indeterminates" and denoted by capital letters; when S is a subset of an algebra $B$ and $A[S]$ denotes the subalgebra of $B$ generated by $S$, the elements of $S$ are denoted by small Latin or Greek letters, or else by the expressions they are determined by in B.

A final point: if M is a monoid, algebraists often write $\mathrm{A}[\mathrm{M}]$ for the monoid algebra $A^{(M)}$. For example, if $G$ is a group, $\mathbb{Z}[G]$ generally denotes the group algebra $\mathbb{Z}^{(\mathrm{G})}$.

### 3.7.4

Proposition (a) If $\mathrm{S}^{\prime}$ is a subset of S , the algebra $\mathrm{A}\left[\mathrm{S}^{\prime}\right]$ is a subalgebra of $\mathrm{A}[\mathrm{S}]$.
(b) If $\mathrm{S}=\mathrm{S}^{\prime} \cup \mathrm{S}^{\prime \prime}$, with $\mathrm{S}^{\prime} \cap \mathrm{S}^{\prime \prime}=\varnothing$, then the algebra $\mathrm{A}[\mathrm{S}]$ is the same as ( $\left.\mathrm{A}\left[\mathrm{S}^{\prime}\right]\right)\left[\mathrm{S}^{\prime \prime}\right]$.
(c) If $\mathrm{S}=\underset{\longrightarrow}{\lim } \mathrm{S}_{i}$, then $\mathrm{A}[\mathrm{S}]=\underset{\longrightarrow}{\lim } \mathrm{A}\left[\mathrm{S}_{i}\right]$.
(d) Let I be an ideal of A and $\mathrm{J} \overrightarrow{\text { the }}$ ideal of $\mathrm{A}[\mathrm{S}]$ generated by I . Then the algebra (A/I)[S] can be identified with $(\mathrm{A}[\mathrm{S}]) / \mathrm{J}$.
Proof (a) $\mathbb{N}^{\left(S^{\prime}\right)} \subset \mathbb{N}^{(S)}$, and (a) follows.
(b) $\mathbb{N}^{(S)}=\mathbb{N}^{\left(S^{\prime}\right)} \times \mathbb{N}^{\left(S^{\prime \prime}\right)}$. Every element of $\mathrm{A}[\mathrm{S}]$ can be uniquely written as

$$
\begin{aligned}
\sum_{n \in \mathbb{N}^{(S)}} a_{n} \mathrm{X}^{n} & =\sum_{\left(n^{\prime}, n^{\prime \prime}\right) \in \mathbb{\mathbb { N } ^ { ( S ^ { \prime } ) } \times \mathbb { N } ( S ^ { \prime \prime } )}} a_{n^{\prime}, n^{\prime \prime}} \mathrm{X}^{n^{\prime}} \mathrm{X}^{n^{\prime \prime}} \\
& =\sum_{n^{\prime \prime} \in \mathbb{N}^{\left(S^{\prime \prime}\right)}}\left(\sum_{n^{\prime} \in \mathbb{N}^{\left(S^{\prime}\right)}} a_{n^{\prime}, n^{\prime \prime}} \mathrm{X}^{n^{\prime}}\right) \mathrm{X}^{n^{\prime \prime}}
\end{aligned}
$$

where $\left(a_{n}\right)_{n \in \mathbb{N}^{(s)}}$ is a family with finite support in A , and so uniquely as $\sum_{n^{\prime \prime} \in \mathbb{N}\left(\mathrm{s}^{\prime \prime}\right)}$ $b_{n^{\prime \prime}} \mathrm{X}^{n^{\prime \prime}}$, where $\left(b_{n^{\prime \prime}}\right)_{n^{\prime \prime} \in \mathbb{N}^{\left(s^{\prime \prime}\right)}}$ is a family with finite support in $\mathrm{A}\left[\mathrm{S}^{\prime}\right]$.
(c) follows from the universal property and from (2.7.5, Corollary).
(d) For any A-algebra $B$, set $F(B)=B^{S}$ if $I B=0$ and $F(B)=\varnothing$ if $I B \neq 0$. The algebras $(\mathrm{A} / \mathrm{I})[\mathrm{S}]$ and $(\mathrm{A}[\mathrm{S}]) / \mathrm{J}$ represent the functor F .

### 3.7.5

Proposition Suppose that A is integral. For all sets S , the ring $\mathrm{A}[\mathrm{S}]$ is integral.
Proof (a) First suppose that S has a unique element. Let P and $\mathrm{P}^{\prime}$ be two nonzero polynomials of respective degrees $d$ and $d^{\prime}$, and leading terms $a \mathrm{X}^{d}$ and $a^{\prime} \mathrm{X}^{d^{\prime}}$. The leading term of $\mathrm{PP}^{\prime}$ is $a a^{\prime} \mathrm{X}^{d+d^{\prime}}$; as it is not zero, $\mathrm{PP}^{\prime} \neq 0$.
(b) Next suppose that S is finite: the proof follows by induction from the previous case since $\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]=\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]\left[\mathrm{X}_{n}\right]$.
(c) General case: it follows from the previous one, together with (3.7.4, c) and (3.1, Exercise 4).

### 3.7.6 Degree of Polynomials

For $n=\left(n_{s}\right)_{s \in \mathrm{~S}} \in \mathbb{N}^{(\mathrm{S})}$, set $|n|=\sum_{s \in \mathrm{~S}} n_{s}$.
Let $\mathrm{P}=\sum a_{n} \mathrm{X}^{n} \in \mathrm{~A}[\mathrm{~S}]$. The degree of P , written $d(\mathrm{P})$ is the upper bound (in $\overline{\mathbb{Z}}$ ) of $|n|$ for $n$ such that $a_{n} \neq 0$. The polynomial P is said to be homogeneous of degree $d$ if $a_{n}=0$ for all $n$ such that $|n| \neq d$.

Remark The degree of the polynomial 0 is $-\infty$, although 0 is homogeneous of degree $d$ for all $d$. For any nonzero polynomial $\mathrm{P}, d(\mathrm{P}) \in \mathbb{N}$.

Polynomials of degree $\leqslant d$ (resp. homogeneous of degree $d$ ) form a free Amodule. If S is a finite set with $m$ elements, the dimension of this module is the binomial coefficient $\binom{m+d}{d}$ (resp. $\binom{m+d-1}{d}$ ).

Indeed $\left(n_{1}, \ldots, n_{m}\right) \mapsto\left(n_{1}+1, n_{1}+n_{2}+2, \ldots, n_{1}+\cdots+n_{m}+m\right)$ is a bijection from the set of $n \in \mathbb{N}^{m}$ such that $|n| \leqslant d$ (resp. $|n|=d$ ) onto the set of strictly increasing sequences with $m$ elements in $\{1, \ldots, d+m\}$ (resp. whose last element is $d+m$ ).

If P and Q are homogeneous of respective degrees $p$ and $q$, then PQ is homogeneous of degree $p+q$.

Proposition Let $\mathrm{P}, \mathrm{Q} \in \mathrm{A}[\mathrm{S}]$. Then $d(\mathrm{PQ}) \leqslant d(\mathrm{P})+d(\mathrm{Q})$. If A is an integral domain, then $d(\mathrm{PQ})=d(\mathrm{P})+d(\mathrm{Q})$.

Proof The inequality is immediate. We prove the equality when A is integral. We may assume that P and Q are nonzero. Write $\mathrm{P}=\mathrm{P}_{0}+\mathrm{P}_{1}$, where $\mathrm{P}_{0}$ is homogeneous of degree $d(\mathrm{P})$ and $\mathrm{P}_{1}$ of degree $<d(\mathrm{P}),{ }^{18}$ and likewise write Q as $\mathrm{Q}_{0}+\mathrm{Q}_{1}$. Then $\mathrm{PQ}=\mathrm{P}_{0} \mathrm{Q}_{0}+\mathrm{R}$, where $\mathrm{R}=\mathrm{P}_{0} \mathrm{Q}_{1}+\mathrm{P}_{1} \mathrm{Q}_{0}+\mathrm{P}_{1} \mathrm{Q}_{1}$ and $\mathrm{P}_{0} \mathrm{Q}_{0}$ are homogeneous of nonzero degree $d(\mathrm{P})+d(\mathrm{Q})$ by Proposition 3.7.5 and R is of degree $<d(\mathrm{P})+d(\mathrm{Q})$. Hence the degree of PQ is $d(\mathrm{P})+d(\mathrm{Q})$.

### 3.7.7

Theorem (Hilbert ${ }^{19}$ ) If A is a Noetherian ring, then so is $\mathrm{A}[\mathrm{X}]$.

Proof Let I be an ideal of $\mathrm{A}[\mathrm{X}]$; we show that I is finitely generated. For any $\mathrm{P} \in \mathrm{A}[\mathrm{X}]$ with $\mathrm{P} \neq 0$, let $d(\mathrm{P})$ be the degree of P and $\alpha(\mathrm{P})$ the leading coefficient of P (i.e. the coefficient of $\left.\mathrm{X}^{d(\mathrm{P})}\right)$.

Let J be the subset of A consisting of 0 and of $\alpha(\mathrm{P})$ for $\mathrm{P} \in \mathrm{I}-\{0\}$. The set J is an ideal in A. Indeed, let $a=\alpha(\mathrm{P})$ and $b=\alpha(\mathrm{Q})$ be elements of J , and suppose that $d(\mathrm{P}) \leqslant d(\mathrm{Q})$; if $a+b \neq 0$, then $a+b=\alpha\left(\mathrm{X}^{d(\mathrm{Q})-d(\mathrm{P})} \mathrm{P}+\mathrm{Q}\right) \in \mathrm{J}$; for $\lambda \in \mathrm{A}$, if $\lambda a \neq 0$, then $\lambda a=\alpha(\lambda \mathrm{P})$.

As A is Noetherian, J is finitely generated. Let $a_{1}, \ldots, a_{k}$ be its generators and $\mathrm{P}_{1}, \ldots, \mathrm{P}_{k}$ polynomials such that $a_{i}=\alpha\left(\mathrm{P}_{i}\right)$. All $\mathrm{P}_{i}$ may be assumed to be of degree $d$ to within multiplication by powers of X . Let L be the submodule of $\mathrm{A}[\mathrm{X}]$ consisting of polynomials of degree $<d$ and set $\mathrm{M}=\mathrm{L} \cap \mathrm{I}$. The module L is isomorphic to $\mathrm{A}^{d}$, and so M is finitely generated; let $\left(\mathrm{R}_{1}, \ldots, \mathrm{R}_{q}\right)$ generate M as an A-module.

[^14]We show that the polynomials $\mathrm{P}_{1}, \ldots, \mathrm{P}_{k}, \mathrm{R}_{1}, \ldots, \mathrm{R}_{q}$ generate the ideal I. We do this by showing by induction on $d(\mathrm{P})$ that all polynomials $\mathrm{P} \in \mathrm{I}$ are linear combinations of these polynomials, with coefficients in $\mathrm{A}[\mathrm{X}]$. If $d(\mathrm{P})<d$, the polynomial P is the A-linear combination of $\mathrm{R}_{1}, \ldots, \mathrm{R}_{q}$.

If $d(\mathrm{P}) \geqslant d$, the scalar $\alpha(\mathrm{P})$ is of the form $\sum \lambda_{i} a_{i}$, and $\mathrm{P}-\mathrm{X}^{d(\mathrm{P})-d} \sum \lambda_{i} a_{i}$ is of degree $<d(\mathrm{P})$, and so ...

Corollary 3.4 If A is a Noetherian ring then, for all $n$, so is $\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$.
This follows from the theorem by induction on $n$ since

$$
\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]=\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]\left[\mathrm{X}_{n}\right]
$$

Corollary 3.5 If A is a Noetherian ring, then so are all unital associative commutative A-algebras.

Indeed, such an algebra is isomorphic to a quotient of $\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ for some $n$.
Corollary 3.6 Every ring is the union of an increasing directed family of Noetherian subrings.

Indeed, the subring generated by a finite subset of $A$ is a finitely generated $\mathbb{Z}$-algebra, and so is Noetherian.

### 3.7.8 Polynomials with Coefficients in a UFD

Throughout this subsection, A is assumed to be a UFD, K denotes the field of fractions of A and for irreducible element $m \in \mathrm{~A}, \chi_{m}$ the canonical homomorphism from $\mathrm{A}[\mathrm{X}]$ to $(\mathrm{A} / m)[\mathrm{X}]$.

Definition A polynomial $\mathrm{P}=\sum a_{i} \mathrm{X}^{i} \in \mathrm{~A}[\mathrm{X}]$ is said to be primitive if 1 is a g.c.d. of the $a_{i}$.

P is primitive if and only if for all irreducible elements $m \in \mathrm{~A}, \chi_{m}(\mathrm{P}) \neq 0$.
Lemma 3.1 Every nonzero polynomial $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ can be written $\mathrm{P}=\lambda \mathrm{P}_{1}$ for some primitive polynomial $\mathrm{P}_{1} \in \mathrm{~A}[\mathrm{X}]$ and some nonzero $\lambda \in \mathrm{K}$.
Proof $\mathrm{P}=\frac{1}{d} \mathrm{P}_{0}$, for someù $d \in \mathrm{~A}$ with $d \neq 0$ and some $\mathrm{P}_{0} \in \mathrm{~A}[\mathrm{X}]$. Let $a$ be the g.c.d. of the coefficients of $\mathrm{P}_{0}$. Then $\mathrm{P}_{0}=a \mathrm{P}_{1}$, where $\mathrm{P}_{1} \in \mathrm{~A}[\mathrm{X}]$ is primitive, and $\mathrm{P}=\lambda \mathrm{P}_{1}$ with $\lambda=\frac{a}{d}$.
Lemma 3.2 Let P and Q be polynomials in $\mathrm{A}[\mathrm{X}]$. Assume P is primitive and that P divides Q in $\mathrm{K}[\mathrm{X}]$. Then P divides Q in $\mathrm{A}[\mathrm{X}]$.

Proof Suppose that $\mathrm{Q} \neq 0$, and $\mathrm{Q}=\mathrm{PR}$ with $\mathrm{R} \in \mathrm{K}[\mathrm{X}]$. There exist $a, b \in \mathrm{~A}$ with g.c.d. $(a, b)=1$ such that $\mathrm{R}=\frac{a}{b} \mathrm{R}_{1}$ for some primitive $\mathrm{R}_{1} \in \mathrm{~A}[\mathrm{X}]$. Then $b \mathrm{Q}=a \mathrm{PR}_{1}$. If $m \in \mathrm{~A}$ is irreducible element dividing $b$, then $\chi_{m}(a) \neq 0, \chi_{m}(\mathrm{P}) \neq 0$ and $\chi_{m}\left(\mathrm{R}_{1}\right) \neq 0$; the ring $\mathrm{A} / m$ is an integral domain (3.2.8, (ii)), and hence so is $(\mathrm{A} / m)[\mathrm{X}](3.7 .5)$, and $\chi_{m}(b \mathrm{Q})=\chi_{m}\left(a \mathrm{PR}_{1}\right) \neq 0$, which contradicts $\chi_{m}(b)=0$. Hence $b$ is invertible in A.

## Lemma 3.3 Let $\mathrm{P} \in \mathrm{A}[\mathrm{X}]$ be a nonzero polynomial.

(a) If P is of degree 0 , then P is irreducible in $\mathrm{A}[\mathrm{X}]$ if and only if it is irreducible in A .
(b) If P is of degree $>0, \mathrm{P}$ is irreducible in $\mathrm{A}[\mathrm{X}]$ if and only if it is primitive and irreducible in $\mathrm{K}[\mathrm{X}]$.

Proof (a) If $\mathrm{P}=\mathrm{QR}$, with Q and R non-invertible elements of $\mathrm{A}[\mathrm{X}]$, the polynomials Q and R are of degree 0 , and so are non-invertible elements of A . If $\mathrm{P}=a b$, where $a, b \in \mathrm{~A}$ are non-invertible, then $a, b$ remain non-invertible in $\mathrm{A}[\mathrm{X}]$.
(b) Suppose that P is irreducible in $\mathrm{A}[\mathrm{X}]$. Then P is primitive for otherwise $\mathrm{P}=a \mathrm{P}_{1}$, with $a \in \mathrm{~A}$ non-invertible in A and so non-invertible in $\mathrm{A}[\mathrm{X}]$, and $\mathrm{P}_{1}$ of degree $>0$, hence non-invertible in $\mathrm{A}[\mathrm{X}]$. We show that P is irreducible in $\mathrm{K}[\mathrm{X}]$. If $\mathrm{P}=\mathrm{QR}$ for non-invertible $\mathrm{Q}, \mathrm{R} \in \mathrm{K}[\mathrm{X}]$, and so of degree $>0$ then, by Lemma 3.1, Q may be assumed to be a primitive polynomial of $\mathrm{A}[\mathrm{X}]$, and thus $\mathrm{R} \in \mathrm{A}[\mathrm{X}]$ by Lemma 3.2.

Suppose that $P$ is primitive and irreducible in $K[X]$. Then there are no $Q, R \in A[X]$ of degree $>0$ such that $P=Q R$ since $Q$ and $R$ would then be non-invertible in $K[X]$. There are also no non-invertible Q of degree 0 in $\mathrm{A}[\mathrm{X}]$ such that $\mathrm{P}=\mathrm{QR}$ since Q would then be non-invertible element in A and P would not be primitive. Therefore P is irreducible in $\mathrm{A}[\mathrm{X}]$.

### 3.7.9

Theorem If A is a UFD then so is $\mathrm{A}[\mathrm{S}]$ for all sets S .

Proof We keep the notation of 3.2.7. By 3.2.8, it amounts to showing that A[S] satisfies the following two conditions:
(i) every decreasing sequence of elements of $\mathrm{Mon}^{*}(\mathrm{~A}[\mathrm{~S}])$ is stationary;
(ii) for every irreducible element $\mathrm{P} \in \mathrm{A}[\mathrm{S}]$, the ideal generated by P is prime.
(a) Condition (i). Let $\left(\mathrm{P}_{k}\right)_{k \in \mathbb{N}}$ be a sequence of elements of $\mathrm{A}[\mathrm{S}]$, with $\mathrm{P}_{k}=$ $\sum a_{k, n} \mathrm{X}^{n}$ such that $\left(\chi\left(\mathrm{P}_{k}\right)\right)_{k \in \mathbb{N}}$ is a decreasing sequence of elements of Mon* $(\mathrm{A}[\mathrm{S}])$. The degrees $d\left(\mathrm{P}_{k}\right)$ form a decreasing sequence of elements of $\mathbb{N}$, hence constant from some rank $k_{0}$. For $k \geqslant k_{0}, \mathrm{P}_{k}=c_{k} \mathrm{P}_{k+1}$, where $c_{k} \in \mathrm{~A}^{*}$. Let $n \in \mathbb{N}^{(\mathrm{S})}$ be such that $a_{k_{0}, n} \neq 0$. For $k>k_{0}, a_{k, n}=c_{k} a_{k+1, n}$ and $\left(\chi\left(a_{k, n}\right)\right)_{k \geqslant k_{0}}$ is a decreasing sequence of
elements of Mon*A. This sequence is therefore stationary, and all $c_{k}$ are invertible beyond a certain rank, proving (i).
(b) Condition / (ii). Suppose first that S has only one element X. Let P be an irreducible element of $\mathrm{A}[\mathrm{X}]$ and K the field of fractions of A . By Lemma 3.3, P is primitive, and by Lemma 3.2, $\mathrm{P} \cdot \mathrm{A}[\mathrm{X}]=\mathrm{A}[\mathrm{X}] \cap \mathrm{P} \cdot \mathrm{K}[\mathrm{X}]$. So $\mathrm{A}[\mathrm{X}] / \mathrm{P} \cdot \mathrm{A}[\mathrm{X}]$ can be identified with a subring of $\mathrm{K}[\mathrm{X}] / \mathrm{P} \cdot \mathrm{K}[\mathrm{X}]$. By Lemma 3.3, P is irreducible in $\mathrm{K}[\mathrm{X}]$; as $\mathrm{K}[\mathrm{X}]$ is a PID, and thus a UFD, $\mathrm{K}[\mathrm{X}] / \mathrm{P} \cdot \mathrm{K}[\mathrm{X}]$ is an integral domain, and hence so is $\mathrm{A}[\mathrm{X}] / \mathrm{P} \cdot \mathrm{A}[\mathrm{X}]$.
(c) Suppose next that S is finite. The theorem follows by induction from the previous case.
(d) General case. Let P be an irreducible element of $\mathrm{A}[\mathrm{S}]$. There is a finite subset $S^{\prime}$ of $S$ such that $P \in A\left[S^{\prime}\right]$. Then $P$ is irreducible in $A\left[S^{\prime}\right]$, and so $A\left[S^{\prime}\right] / P \cdot A\left[S^{\prime}\right]$ is an integral domain. Set. Then so is $\mathrm{A}[\mathrm{S}] / \mathrm{P} \cdot \mathrm{A}[\mathrm{S}]=\left(\mathrm{A}\left[\mathrm{S}^{\prime}\right] / \mathrm{P} \cdot \mathrm{A}\left[\mathrm{S}^{\prime}\right]\right)\left[\mathrm{S}^{\prime \prime}\right]$, where $S^{\prime \prime}=S-S^{\prime}$, since $A[S]=\left(A\left[S^{\prime}\right]\right)\left[S^{\prime \prime}\right]$ (3.7.5).

### 3.7.10 Substitution, Polynomial Functions

Let B be a unital associative commutative A-algebra, $\left(x_{i}\right)_{i \in \mathrm{I}}$, a family of elements in B and $f: \mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right] \rightarrow \mathrm{B}$ the homomorphism defined by $f\left(\mathrm{X}_{i}\right)=x_{i}$.

The image of $f$ is the unital subalgebra of B generated by $\left(x_{i}\right)_{i \in \mathrm{I}}$. If $f$ is injective, then the family $\left(x_{i}\right)_{i \in \mathrm{I}}$ is said to be algebraically free, or that the elements $x_{i}$ are algebraically independent. We say that $x \in \mathrm{~B}$ is transcendental over A if the family consisting of $x$ is algebraically free over A. In general, the kernel of $f$ is an ideal of $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ called the ideal of algebraic relations between the $x_{i}$.

For $\mathrm{P} \in \mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right], f(\mathrm{P})$ is the element of B obtained by substituting the $x_{i}$ for the indeterminates $\mathrm{X}_{i}$ in P , and it is written $\mathrm{P}_{\mathrm{B}}\left(\left(x_{i}\right)_{i \in \mathrm{I}}\right)$ or simply $\mathrm{P}\left(\left(x_{i}\right)_{i \in \mathrm{I}}\right) .{ }^{20}$

The map $\mathrm{P}_{\mathrm{B}}$ from $\mathrm{B}^{\mathrm{I}}$ to B thus defined is the polynomial function defined by P on $\mathrm{B}^{\mathrm{I}}$. The map $\mathrm{P} \mapsto \mathrm{P}_{\mathrm{B}}$ is a unital A -algebra homomorphism from $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)\right]$ to the algebra of maps from $\mathrm{B}^{\mathrm{I}}$ to B ; it associates to $\mathrm{X}_{i}$ the projection of index $i$ from $\mathrm{B}^{\mathrm{I}}$ to B.

Proposition If B an infinite integral domain and the homomorphism $i: \mathrm{A} \rightarrow \mathrm{B}$ defining the algebra structure of B is injective, then the homomorphism $\mathrm{P} \mapsto \mathrm{P}_{\mathrm{B}}$ is injective.
Proof (a) Suppose that I has only one element. Let K be the field of fractions of B and let $\mathrm{P} \in \mathrm{A}[\mathrm{X}] \subset \mathrm{K}[\mathrm{X}]$. If P is of degree $d$, there are at most $d$ roots ${ }^{21}$ in K ; hence $\mathrm{P}(x)=0$ for all $x \in \mathrm{~B}$ only if $\mathrm{P}=0$.
(b) Suppose next that I is finite. Assume $\mathrm{I}=\{1, \ldots, n\}$ and use induction on $n$. Let $\mathrm{P}=\sum a_{k} \mathrm{X}_{n}^{k} \in \mathrm{~A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]\left[\mathrm{X}_{n}\right]=\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$, where $a_{k} \in \mathrm{~A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]$,

[^15]and assume the map $\mathrm{P}_{\mathrm{B}}: \mathrm{B}^{n} \rightarrow \mathrm{~B}$ is zero. Let $\left(x_{1}, \ldots, x_{n-1}\right) \in \mathrm{B}^{n-1}$, and set $\mathrm{Q}=$ $\mathrm{P}\left(x_{1}, \ldots, x_{n-1}, \mathrm{X}_{n}\right)=\sum a_{k}\left(x_{1}, \ldots, x_{n-1}\right) \mathrm{X}_{n}^{k} \in \mathrm{~B}\left[\mathrm{X}_{n}\right]$. By assumption, $\left(\forall x_{n} \in \mathrm{~B}\right)$ $\mathrm{Q}\left(x_{n}\right)=\mathrm{P}\left(x_{1}, \ldots, x_{n-1}, x_{n}\right)=0$, and so $\mathrm{Q}=0$ by (a), and $(\forall k) a_{k}\left(x_{1}, \ldots, x_{n-1}\right)=0$. As this holds for all $x_{1}, \ldots, x_{n-1} \in \mathrm{~B},(\forall k) a_{k}=0$ by the induction hypothesis. Thus $\mathrm{P}=0$.
(c) General case. Let $(0 \neq) \mathrm{P} \in \mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$. There is a finite subset J of I such that $\mathrm{P} \in \mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{~J}}\right]$. Then the map $\mathrm{P}_{\mathrm{B}}: \mathrm{B}^{\mathrm{I}} \rightarrow \mathrm{B}$ is the composite of the surjective projection $B^{I} \rightarrow B^{J}$ and of the nonzero map $P_{B}: B^{J} \rightarrow B$, and so in not zero.

Remarks (1) If A is an integral domain, then we may take $\mathrm{B}=\mathrm{A}$ in the proposition. For all integral domains A , there is an infinite integral domain B , for example $\mathrm{A}[\mathrm{X}]$, with a subring with which A can be identified.
(2) The assumption of the infinity of B cannot be removed, nor that of the integrality of $B$ (see Exercise 15).
(3) A map $\varphi$ from $\mathrm{A}^{p}$ to $\mathrm{A}^{q}$ is said to be polynomial if the maps $\varphi_{1}, \ldots, \varphi_{q}$ from $\mathrm{A}^{p}$ to A such that $\varphi=\left(\varphi_{1}, \ldots, \varphi_{q}\right)$ are polynomials. Let M and N be finitely generated free A-modules. A map $\varphi$ from M to N is said to be polynomial with respect to the bases $e=\left(e_{1}, \ldots, e_{p}\right)$ of M and $f=\left(f_{1}, \ldots, f_{q}\right)$ of N if the expression of $\varphi$ with respect to the bases $e$ and $f$ (i.e. the map from $\mathrm{A}^{p}$ to $\mathrm{A}^{q}$ induced by $\varphi$ by identifying $\mathrm{A}^{p}$ and $\mathrm{A}^{q}$ with M and N via $e$ and $f$ respectively) is polynomial. If a map $\varphi$ from M to N is polynomial with respect to one pair of bases, it remains so with respect to any pair of bases. Hence it may be called a polynomial map from M to N .

The composite of two polynomial maps is polynomial.
If $A=\mathbb{R}$ or $\mathbb{C}$, equip $M$ with the topology obtained by transferring that of $A^{p}$ via a basis (this topology, called the canonical topology is independent of the choice of basis; It is the only Hausdorff topology compatible with the vector space structure); do the same for N . Every polynomial map from M to N is continuous (and even $\mathrm{C}^{\infty}$, or better still: analytic). Two polynomial maps from M to N coinciding on a nonempty open subset of $M$ are equal: this follows from the Taylor formula for polynomials. In particular, if $\varphi$ is a nonzero polynomial map from M to N , then $\varphi^{-1}(0)$ is closed and with empty interior in M.

### 3.7.11 Resultants

For $p \in \mathbb{N}$, denote by $\mathrm{A}[\mathrm{X}]_{p}$ the submodule of $\mathrm{A}[\mathrm{X}]$ consisting of polynomials of degree $\leqslant p$. It is a free module of rank $p+1$ as $\left(1, \mathrm{X}, \ldots, \mathrm{X}^{p}\right)$ is a basis.

Definition Let $\mathrm{P} \in \mathrm{A}[\mathrm{X}]_{p}$ and $\mathrm{Q} \in \mathrm{A}[\mathrm{X}]_{q}$. The resultant of P and Q for the degrees $p$ and $q$, written $\operatorname{Res}_{p, q}(\mathrm{P}, \mathrm{Q})$, is the determinant of the matrix representing the $\operatorname{map}(\mathrm{U}, \mathrm{V}) \mapsto \mathrm{U} . \mathrm{P}+\mathrm{V} . \mathrm{Q}$ from $\mathrm{A}[\mathrm{X}]_{q-1} \oplus \mathrm{~A}[\mathrm{X}]_{p-1}$ to $\mathrm{A}[\mathrm{X}]_{p+q-1}$, with respect to bases $\left((1,0), \ldots,\left(\mathrm{X}^{q-1}, 0\right),(0,1), \ldots,\left(0, \mathrm{X}^{p-1}\right)\right)$ and $\left(1, \ldots, \mathrm{X}^{p+q-1}\right)$.

If $p$ and $q$ are the degrees of P and Q respectively, we simply write $\operatorname{Res}(\mathrm{P}, \mathrm{Q})$, and call this the resultant of P and Q .

For $\mathrm{P}=a_{p} \mathrm{X}^{p}+\cdots+a_{0}$ and $\mathrm{Q}=b_{q} \mathrm{X}^{q}+\cdots+b_{0}$, the resultant $\operatorname{Res}_{p, q}(\mathrm{P}, \mathrm{Q})$ is therefore the determinant of the Sylvester matrix:

$$
\left(\begin{array}{cccccccc}
a_{0} & 0 & \cdots & 0 & b_{0} & 0 & \cdots & 0 \\
a_{1} & a_{0} & \ddots & 0 & b_{1} & b_{0} & \ddots & \vdots \\
\vdots & & \ddots & 0 & \vdots & & \ddots & 0 \\
& & & a_{0} & & & & \\
& & & & b_{q} & & & b_{0} \\
a_{p} & & & & 0 & b_{q} & & \\
0 & \ddots & & \vdots & & \ddots & \ddots & \vdots \\
\vdots & \ddots & \ddots & & \vdots & & \ddots & \\
0 & \cdots & 0 & a_{p} & 0 & \cdots & \cdots & 0
\end{array}\right)
$$

It is in A , and $\operatorname{Res}_{q, p}(\mathrm{Q}, \mathrm{P})=-1^{p q} \cdot \operatorname{Res}_{p, q}(\mathrm{P}, \mathrm{Q})$.
If $\varphi: \mathrm{A} \rightarrow \mathrm{A}^{\prime}$ is a ring homomorphism, defining $\varphi_{*}: \mathrm{A}[\mathrm{X}] \rightarrow \mathrm{A}^{\prime}[\mathrm{X}]$ by $\left(\varphi_{*}\right)_{\mid \mathrm{A}}=$ $\varphi$ and $\varphi_{*}(\mathrm{X})=\mathrm{X}$ gives

$$
\operatorname{Res}_{p, q}\left(\varphi_{*}(\mathrm{P}), \varphi_{*}(\mathrm{Q})\right)=\varphi\left(\operatorname{Res}_{p, q}(\mathrm{P}, \mathrm{Q})\right)
$$

Here, two elements $a$ and $b$ of A will be said to be relatively prime if the ideal they generate is the whole of A (even if A is not principal).

Proposition Let $\mathrm{P}, \mathrm{Q} \in \mathrm{A}[\mathrm{X}]$ be polynomials of respective degrees $p$ and $q$ with relatively prime leading coefficients $a_{p}, b_{q} \in \mathrm{~A}$. Then, P and Q are relatively prime in $\mathrm{A}[\mathrm{X}]$ if and only if the resultant $\operatorname{Res}(\mathrm{P}, \mathrm{Q})$ is invertible in A .

Proof (a) Suppose that $\operatorname{Res}(\mathrm{P}, \mathrm{Q})$ is invertible. The map $(\mathrm{U}, \mathrm{V}) \mapsto \mathrm{U} . \mathrm{P}+\mathrm{V} . \mathrm{Q}$ from $\mathrm{A}[\mathrm{X}]_{q-1} \oplus \mathrm{~A}[\mathrm{X}]_{p-1}$ to $\mathrm{A}[\mathrm{X}]_{p+q-1}$ is bijective. In particular there exist U and V such that $\mathrm{U} . \mathrm{P}+\mathrm{V} \cdot \mathrm{Q}=1$.
(b) Conversely, suppose that P and Q are relatively prime. Let $\mathrm{R} \in \mathrm{A}[\mathrm{X}]_{p+q-1}$, and $h$ be minimal of the form $\mathrm{R}=\mathrm{U} . \mathrm{P}+\mathrm{V} . \mathrm{Q}$ with $\mathrm{U} \in \mathrm{A}[\mathrm{X}]_{h-p}$ and $\mathrm{V} \in \mathrm{A}[\mathrm{X}]_{h-q}$; we show that $h<p+q$.

Assume $h \geqslant p+q$; let $\alpha$ and $\beta$ be the coefficients of degree $h-p$ and $h-q$ of U and V respectively. Since $\mathrm{R}=\mathrm{U} . \mathrm{P}+\mathrm{V} . \mathrm{Q}$ is of degree $<h, \alpha \cdot a_{p}+\beta . b_{q}=0$. As $a_{p}$ and $b_{q}$ are relatively prime, $\mu \cdot a_{p}+\nu \cdot b_{q}=1$. It follows that $\alpha=\lambda \cdot b_{q}$ and $\beta=-\lambda \cdot a_{p}$ with $\lambda=\alpha \cdot \nu-\beta \cdot \mu \cdot \operatorname{Set} \tilde{\mathrm{U}}=\mathrm{U}-\lambda . \mathrm{Q} \cdot \mathrm{X}^{d}$ and $\tilde{\mathrm{V}}=\mathrm{V}+\lambda . \mathrm{P} \cdot \mathrm{X}^{d}$, where $d=h-(p+q)$. Then $\tilde{\mathrm{U}} \cdot \mathrm{P}+\tilde{\mathrm{V}} \cdot \mathrm{Q}=\mathrm{U} \cdot \mathrm{P}+\mathrm{V} \cdot \mathrm{Q}=\mathrm{R}$, with $d^{o}(\tilde{\mathrm{U}})<h-p$ and $d^{o}(\tilde{\mathrm{~V}})<h-q$, contradicting the minimality of $h$.

### 3.7.12 Discriminants

Let $\quad \mathrm{P}=a_{0}+a_{1} \mathrm{X}+\cdots+a_{p} \mathrm{X}^{p} \in \mathrm{~A}[\mathrm{X}]$, and consider its derivative $\mathrm{P}^{\prime}=a_{1}+\cdots+p a_{p} \mathrm{X}^{p-1}$. The resultant $\operatorname{Res}_{p, p-1}$ can be written $a_{p} . \Delta$, where $\Delta$ is the determinant of

$$
\left(\begin{array}{ccccccc}
a_{0} & \cdots & 0 & a_{1} & 0 & \cdots & 0 \\
& \ddots & \vdots & 2 a_{2} & \ddots & \ddots & \vdots \\
\vdots & & a_{0} & \vdots & & \ddots & 0 \\
& & & p a_{p} & & & a_{1} \\
a_{p} & & \vdots & 0 & \ddots & & \vdots \\
\vdots & \ddots & & \vdots & \ddots & p a_{p} & \\
0 & \cdots & 1 & 0 & \cdots & 0 & p
\end{array}\right)
$$

This determinant is called the discriminant of P and is written discr $(\mathrm{P})$.
Examples For a second degree polynomial $\mathrm{P}=a \mathrm{X}^{2}+b \mathrm{X}+c$,

$$
\operatorname{discr}(\mathrm{P})=\operatorname{det}\left(\begin{array}{ccc}
c & b & 0 \\
b & 2 a & b \\
1 & 0 & 2
\end{array}\right)=4 a c-b^{2}
$$

For $\mathrm{P}=\mathrm{X}^{3}+p \mathrm{X}+q$,

$$
\operatorname{discr}(\mathrm{P})=\operatorname{det}\left(\begin{array}{lllll}
q & 0 & p & 0 & 0 \\
p & q & 0 & p & 0 \\
0 & p & 3 & 0 & p \\
1 & 0 & 0 & 3 & 0 \\
0 & 1 & 0 & 0 & 3
\end{array}\right)=\operatorname{det}\left(\begin{array}{ccc}
p-3 q & 0 \\
0 & -2 p & -3 q \\
3 & 0 & -2 p
\end{array}\right)=4 p^{3}+27 q^{2}
$$

### 3.7.13 Algebraic Sets

Let K be a field and $\left(\mathrm{P}_{\lambda}\right)_{\lambda \in \Lambda}$ a family of polynomials in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$. Denote by $\mathrm{V}\left(\left(\mathrm{P}_{\lambda}\right)\right)$ the set of $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathrm{K}^{n}$ such that $(\forall \lambda \in \Lambda) \mathrm{P}_{\lambda}(x)=0$. The subsets of $\mathrm{K}^{n}$ thus obtained are called algebraic sets. The set $\mathrm{V}\left(\left(\mathrm{P}_{\lambda}\right)\right)$ is defined by equations $P_{\lambda}$, and does not depend on the ideal I generated by $P_{\lambda}$. It will also be written $V(I)$. Since the ring $K\left[X_{1}, \ldots, X_{n}\right]$ is Noetherian, every algebraic set can be defined by finitely many equations. Let $\sqrt{\mathrm{I}}$ be the ideal consisting of elements $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ with some power in I . It is also the inverse image in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ of the nilradical of $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] / \mathrm{I}$. Hence it is the intersection of prime ideals containing I. The following hold:
(1) $\mathrm{V}(\sqrt{\mathrm{I}})=\mathrm{V}(\mathrm{I})$.

Let I and J be ideals in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$, then
(2) $\mathrm{I} \supset \mathrm{J} \Rightarrow \mathrm{V}(\mathrm{I}) \subset \mathrm{V}(\mathrm{J})$;
(3) $\mathrm{V}(\mathrm{I}+\mathrm{J})=\mathrm{V}(\mathrm{I}) \cap \mathrm{V}(\mathrm{J})$;
(4) $\mathrm{V}(\mathrm{I} \cap \mathrm{J})=\mathrm{V}(\mathrm{I} \cdot \mathrm{J})=\mathrm{V}(\mathrm{I}) \cup \mathrm{V}(\mathrm{J})$.

If $\left(\mathrm{I}_{\lambda}\right)$ is a family of ideals, then
(5) $\mathrm{V}\left(\sum \mathrm{I}_{\lambda}\right)=\bigcap_{\lambda} \mathrm{V}\left(\mathrm{I}_{\lambda}\right)$.

Thanks to formulas (4) and (5), there is a unique topology on $\mathrm{K}^{n}$ with respect to which the closed subsets are the algebraic sets. It is called the Zariski topology. It is Hausdorff (except when K is finite or $n=0$ ), but all points are closed.

If K is algebraically closed, then it may be shown that $\mathrm{V}(\mathrm{I})=\mathrm{V}(\mathrm{J})$ if and only if $\sqrt{\mathrm{I}}=\sqrt{\mathrm{J}}$. This result is known as the Hilbert's Nullstellensatz (Exercises 13, 14).

Let $k$ be a subfield of K . A subset of $\mathrm{K}^{n}$ is a $k$-algebraic set in $\mathrm{K}^{n}$ if it is defined by equations in $k\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$.

### 3.7.14 Principle of Extension of Identities

SCholium. Let $\mathscr{R}$ be a relation involving a ring (resp. a unital associative commutative $\mathbb{Q}$-algebra), stable under taking subrings (resp. subalgebras), quotients and directed unions. If $\mathscr{R}$ holds for $\mathbb{R}$ or $\mathbb{C}$, then it holds for all rings (resp. all $\mathbb{Q}$-algebras).

More precisely:
Proposition Let $\mathscr{A}$ be the category of commutative rings (resp. commutative, associative, unital $\mathbb{Q}$-algebras), and $\mathscr{A}^{\prime}$ a complete subcategory of $\mathscr{A}$. Suppose that
(i) if A and B are objects of $\mathscr{A}$ for which there exists an injective homomorphism $\mathrm{A} \rightarrow \mathrm{B}$, then $\mathrm{B} \in \mathscr{A}^{\prime} \Rightarrow \mathrm{A} \in \mathscr{A}^{\prime}$;
(ii) if A and B are objects of $\mathscr{A}$ for which there exists a surjective homomorphism $\mathrm{A} \rightarrow \mathrm{B}$, then $\mathrm{A} \in \mathscr{A}^{\prime} \Rightarrow \mathrm{B} \in \mathscr{A}^{\prime} ;$
(iii) if an object A de $\mathscr{A}$ is a directed union of objets of $\mathscr{A}^{\prime}$, then A is an object of $\mathscr{A}^{\prime}$.
Then, $\mathscr{A}^{\prime}=\mathscr{A}$ if $\mathbb{R} \in \mathscr{A}^{\prime}$, or if $\mathbb{C} \in \mathscr{A}^{\prime}$.
Lemma There is an algebraically free sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ in $\mathbb{R}$ over $\mathbb{Q}$.
Proof Let us construct such a sequence by induction. Let $\left(x_{0}, \ldots, x_{n-1}\right)$ be a finite sequence in $\mathbb{R}$ algebraically free over $\mathbb{Q}$, and set $\mathrm{A}_{n-1}=\mathbb{Q}\left[x_{0}, \ldots, x_{n-1}\right]$. The set $\mathrm{A}_{n-1}$ is countable. Let $x_{n} \in \mathbb{R}$. Then, $\left(x_{0}, \ldots, x_{n}\right)$ is algebraically free over $\mathbb{Q}$ if and only if $x_{n}$ is transcendental over $\mathrm{A}_{n-1}$. Elements of $\mathbb{R}$ that are not transcendental over $\mathrm{A}_{n-1}$ are algebraic over $\mathrm{A}_{n-1}$ i.e. roots of a nonzero polynomial $\mathrm{P} \in \mathrm{A}_{n-1}[\mathrm{X}]$. The set $\mathrm{A}_{n-1}[\mathrm{X}]-\{0\}$ is countable and every polynomial in it has only finitely many roots. So the set of algebraic elements over $\mathrm{A}_{n-1}$ is a countable subset of $\mathbb{R}$. As $\mathbb{R}$ is not countable, there is a transcendental real $x_{n}$ over $\mathrm{A}_{n-1}$, enabling us to continue the construction by induction.

Proof of the Proposition For all $n$, there are injective homomorphisms $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots\right.$, $\left.\mathrm{X}_{n}\right] \rightarrow \mathbb{Q}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] \rightarrow \mathbb{R} \rightarrow \mathbb{C}$, and so, if $\mathbb{R}$ or $\mathbb{C}$ is an object of $\mathscr{A}^{\prime}, \mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ (resp. $\mathbb{Q}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ ) is an object of $\mathscr{A}^{\prime}$. By (ii), any finitely generated $\mathbb{Z}$-algebra (resp. $\mathbb{Q}$-algebra) of $\mathscr{A}$ is an object of $\mathscr{A}^{\prime}$. Since any objet of $\mathscr{A}$ is an increasing directed union of finitely generated subobjects, by (iii) every object of $\mathscr{A}$ is an object of $\mathscr{A}^{\prime}$.

### 3.7.15 Examples

(1) EXPONENTIAL AND LOGARITHM.

Proposition Let A be a $\mathbb{Q}$-algebra. For a nilpotent $u \in \mathrm{~A}$, set $\exp (u)=\sum_{0}^{\infty} \frac{1}{n!} u^{n}$ and $\log (1+u)=\sum_{1}^{\infty}(-1)^{n-1} \frac{1}{n} u^{n}$. Then $\log (\exp (u))=u$ and $\exp (\log (1+u))=$ $1+u$.

Scholium 3.7.14 cannot be directed applied to this proposition, but it can to the next lemma from which the latter readily follows.

Lemma Let A be a $\mathbb{Q}$-algebra. Define the polynomials $\lambda_{p}, \varepsilon_{p} \in \mathrm{~A}[\mathrm{X}]$ by $\lambda_{p}(\mathrm{X})=$ $\sum_{1}^{p}(-1)^{n-1} \frac{1}{n} \mathrm{X}^{n}$ and $\varepsilon_{p}(\mathrm{X})=\sum_{0}^{p} \frac{1}{n!} \mathrm{X}^{n}$. Then $\lambda_{p}\left(\varepsilon_{p}(\mathrm{X})\right) \equiv \mathrm{X}\left(\bmod \mathrm{X}^{p+1}\right)$ and $\varepsilon_{p}\left(\lambda_{p}(\mathrm{X})\right) \stackrel{n}{\equiv} \mathrm{X}\left(\bmod \mathrm{X}^{p+1}\right)$.

For $\mathrm{A}=\mathbb{R}$, this lemma follows by considering the Taylor expansion of degree $p$ of exp and log. It readily follows that the category of unital associative commutative $\mathbb{Q}$-algebras for which this lemma holds satisfies conditions (i), (ii) and (iii) of Proposition 3.7.14.
(2) Hamilton- - Cayley Theorem.

Let $\mathscr{M}_{n}(\mathrm{~A})$ be the algebra of $n \times n$ square matrices with coefficients in A . For $\mathrm{M} \in$ $\mathscr{M}_{n}(\mathrm{~A})$, consider $\mathrm{XI}_{n}-\mathrm{M} \in \mathscr{M}_{n}(\mathrm{~A}[\mathrm{X}])$, where $\mathrm{I}_{n}$ is the unital element of $\mathscr{M}_{n}(\mathrm{~A})$, and define the characteristic polynomial $\mathrm{P}_{\mathrm{M}} \in \mathrm{A}[\mathrm{X}]$ by $\mathrm{P}_{\mathrm{M}}(\mathrm{X})=\operatorname{det}\left(\mathrm{XI}_{n}-\mathrm{M}\right)$. In the unital subalgebra of $\mathscr{M}_{n}(\mathrm{~A})$ generated by M , which is commutative, consider $\mathrm{P}_{\mathrm{M}}(\mathrm{M})$.

Theorem For all $\mathrm{M} \in \mathscr{M}_{n}(\mathrm{~A}), \mathrm{P}_{\mathrm{M}}(\mathrm{M})=0$.
Proof (a) Suppose first that $\mathrm{A}=\mathbb{C}$ and that $\mathrm{P}_{\mathrm{M}}$ has $n$ distinct roots. In this case, M is diagonalizable, and M may be assumed to be diagonal; let $\lambda_{1}, \ldots, \lambda_{n}$ be its diagonal entries. Then $\mathrm{P}_{\mathrm{M}}(\mathrm{X})=\left(\mathrm{X}-\lambda_{1}\right) \cdots\left(\mathrm{X}-\lambda_{n}\right)$, and $\mathrm{P}_{\mathrm{M}}(\mathrm{M})$ is diagonal with diagonal entries $\mathrm{P}_{\mathrm{M}}\left(\lambda_{1}\right), \ldots, \mathrm{P}_{\mathrm{M}}\left(\lambda_{n}\right)$ all zero, and so $\mathrm{P}_{\mathrm{M}}(\mathrm{M})=0$.
(b) Suppose next that $\mathrm{A}=\mathbb{C}$, and that M is arbitrary. For $\mathrm{M} \in \mathscr{M}_{n}(\mathbb{C})$, let $\Delta(\mathrm{M})$ be the discriminant of $\mathrm{P}_{\mathrm{M}}$. The map $\Delta: \mathscr{M}_{n}(\mathbb{C}) \rightarrow \mathbb{C}$ is polynomial because it is the composite of the map $\mathrm{M} \mapsto \mathrm{P}_{\mathrm{M}}$ from $\mathscr{M}_{n}(\mathbb{C})$ to the vector space $\mathrm{S}_{n+1}$ of polynomials of degree $\leqslant n$ and of the map $\mathrm{P} \mapsto$ discr P from $\mathrm{S}_{n+1}$ to $\mathbb{C}$, both of which are polynomial. It is nonzero because, if M is diagonal with distinct diagonal entries,
then $\Delta(\mathrm{M}) \neq 0$. Therefore the set W of $\mathrm{M} \in \mathscr{M}_{n}(\mathbb{C})$ such that $\Delta(\mathrm{M}) \neq 0$ is open and dense in $\mathscr{M}_{n}(\mathbb{C})$. The map $\mathrm{M} \mapsto \mathrm{P}_{\mathrm{M}}(\mathrm{M})$ from $\mathscr{M}_{n}(\mathbb{C})$ to itself is polynomial, and vanishes on W , and so is null.
(c) General case. The rings A for which the theorem holds form a subcategory $\mathscr{A}^{\prime}$ of $\mathscr{A} \boldsymbol{n} n$ satisfying conditions (i), (ii) and (iii) of Proposition 3.7.14, and $\mathbb{C} \in \mathscr{A}^{\prime}$; hence $\mathscr{A}^{\prime}=\mathscr{A} n n$.

## Exercises 3.7. (Polynomials)

1.-What are the prime ideals of $\mathbb{C}[\mathrm{X}]$ ? of $\mathbb{R}[\mathrm{X}]$ ?
2.-Show that in $\mathbb{R}[X, Y]$, the ideal of polymials vanishing on the circle $S^{1}=$ $\left\{(x, y) \in \mathbb{R}^{2} \mid x^{2}+y^{2}=1\right\}$ is principal and prime.
3.-(a) Show that the ideal generated by $\mathrm{Y}^{2}-\mathrm{X}^{2}-1$ in $\mathbb{R}[\mathrm{X}, \mathrm{Y}]$ is prime, but that the ideal generated by the same function in the ring of analytic functions on $\mathbb{R}^{2}$ is not prime.
(b) Give an example of a polynomial $\mathrm{P} \in \mathbb{R}[\mathrm{X}, \mathrm{Y}]$ generating a prime ideal in $\mathbb{R}[\mathrm{X}, \mathrm{Y}]$, but an ideal that is not prime in $\mathbb{C}[\mathrm{X}, \mathrm{Y}]$. Show that the converse is not possible.
4.-Let K be a field. For every ideal I in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$, define $\mathrm{V}(\mathrm{I})$ as in 3.7.13, and for every subset V of $\mathrm{K}^{n}$, set $\mathscr{J}(\mathrm{V})$ to be the ideal in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ consisting of polynomials vanishing on V .
(a) Show that $\mathscr{F}(\mathrm{V} \cup \mathrm{W})=\mathscr{J}(\mathrm{V}) \cap \mathscr{J}(\mathrm{W})$ and that $\mathscr{F}(\mathrm{V} \cap \mathrm{W}) \supset \mathscr{J}(\mathrm{V})+\mathscr{J}(\mathrm{W})$. Give an example where inclusion is strict.
(b) Show that $\mathrm{V}(\mathscr{J}(\mathrm{V}(\mathrm{I})))=\mathrm{V}(\mathrm{I})$ and $\mathscr{J}(\mathrm{V}(\mathscr{f}(\mathrm{V})))=\mathscr{J}(\mathrm{V})$.
5.-Let $\mathrm{I}=\left\{f \in \mathbb{R}[\mathrm{X}, \mathrm{Y}] \mid(\forall x \in \mathbb{R}) f(x, 0)=0\right.$ and $\left.f_{\mathrm{Y}}^{\prime}(0,0)=0\right\}$. Show that I is an ideal and give a family of generators. Is the ideal I principal? Describe V(I) and $\mathscr{J}(\mathrm{V}(\mathrm{I}))$.
6.-Describe $\mathrm{V}(\mathrm{I})$ and $\mathscr{J}(\mathrm{V}(\mathrm{I}))$ for the following ideals of $\mathbb{R}[\mathrm{X}, \mathrm{Y}]:\left(\mathrm{X}^{2}-\mathrm{Y}^{2}\right)$; $\left(\mathrm{X}^{2}+\mathrm{Y}^{2}\right) ;\left(\mathrm{X}^{2}+\mathrm{Y}^{2}-1\right) ;\left(\mathrm{X}^{2}+\mathrm{Y}^{2}+1\right) ;\left(\mathrm{X}\left(\mathrm{X}^{2}+\mathrm{Y}^{2}+1\right)\right) ;\left(\mathrm{X}^{2} \mathrm{Y}, \mathrm{X}^{3}\right) ;\left(\mathrm{X}^{2}-\right.$ $Y^{2}, X^{3}+Y^{3}$.

Which among these ideals can be defined by vanishing conditions on polynomials and their derivatives?
7.-(a) Show that in the space of one variable polynomials of degree $d$ with real coefficients, the set of polynomials with $d$ real distinct roots is open.
(b) Let $\mathrm{P} \in \mathbb{R}[\mathrm{X}, \mathrm{Y}]$ be a polynomial of degree $d$ in which the coefficient of $\mathrm{Y}^{d}$ is nonzero. Assume there exists $x \in \mathbb{R}$ such that the polynomial $\mathrm{P}(x, \mathrm{Y}) \in \mathbb{R}[\mathrm{Y}]$ has $d$ distinct real roots. Show that every polynomial vanishing on $\mathrm{V}(\mathrm{P})$ is a multiple of P (divide by P in $(\mathrm{A}[\mathrm{X}])[\mathrm{Y}])$.
(c) In (b) replace the assumption that the coefficient of $\mathrm{Y}^{d}$ is nonzero with: P is primitive in $(\mathrm{A}[\mathrm{X}])[\mathrm{Y}]$, i.e. $\mathrm{P}(\mathrm{X}, \mathrm{Y})=\sum a_{i}(\mathrm{X}) \mathrm{Y}^{i}$, where the polynomials $a_{i}$ do not have any common divisors (use 3.7.8, Lemma 3.2).
(d) Generalize to more than two variables.
8. (Continuity of roots)-(a) Let A be the Banach algebra of series $\sum a_{n} \mathrm{X}^{n} \in \mathbb{C}[[\mathrm{X}]]$ such that $\sum\left|a_{n}\right|<\infty$. Let $\mathrm{P} \in \mathbb{C}[\mathrm{X}]$ be a polynomial of degree $d$. Show that if P has
$k$ roots with absolute value $<1$ and $d-k$ roots with absolute value $>1$ (counting multiplicities), every element $f \in \mathrm{~A}$ can be uniquely written as $f=\mathrm{P} g+\mathrm{R}$, with $g \in \mathrm{~A}$ and $\mathrm{R} \in \mathbb{C}[\mathrm{X}]$ of degree $<k$. Converse.
(b) Let $\mathrm{P}_{0} \in \mathbb{C}[\mathrm{X}]$ be a polynomial of degree $d$ having 0 as a root of order $k$, and $c$ be the smallest absolute value of the nonzero roots. Show that, for all $\varepsilon>0$, there exists $\alpha>0$ such that any polynomial P of degree $d$ satisfying $\left\|\mathrm{P}-\mathrm{P}_{0}\right\|<\alpha$ has $k$ roots of absolute value $<\varepsilon$ and $d-k$ roots of absolute value $>c-\varepsilon$. (Use the fact that, if E and F are Banach spaces, then $\operatorname{Isom}(\mathrm{E}, \mathrm{F})$ is open in $\mathrm{L}(\mathrm{E}, \mathrm{F})$.)
(c) Let $\mathrm{P}_{0} \in \mathbb{C}[\mathrm{X}]$ be a polynomial of degree $d, x_{1}, \ldots, x_{d}$ its roots. Show that, for all $\varepsilon>0$, there exists $\alpha>0$ such that for any polynomial P of degree $d$ satisfying $\left\|\mathrm{P}-\mathrm{P}_{0}\right\|<\alpha$, its roots $y_{1}, \ldots, y_{d}$ can be classified so that $\left|y_{1}-x_{1}\right|<\varepsilon, \ldots$, $\left|y_{d}-x_{d}\right|<\varepsilon$.
(d) Let $d^{\prime} \geqslant d$. Show that, for all $\varepsilon>0$ and all $\mathrm{M} \in \mathbb{R}_{+}$, there exists $\alpha>0$ such that, for any polynomial P of degree $d^{\prime \prime} \leqslant d^{\prime}$ satisfying $\left\|\mathrm{P}-\mathrm{P}_{0}\right\|<\alpha$, its roots $y_{1}, \ldots, y_{d}$ can be classified so that $\left|y_{1}-x_{1}\right|<\varepsilon, \ldots,\left|y_{d}-x_{d}\right|<\varepsilon,\left|y_{d+1}\right|>\mathrm{M}, \ldots$, $\left|y_{d^{\prime \prime}}\right|>M$.
(e) Extend these results to an arbitrary algebraically closed valued field K (consider the completion $\widehat{\mathrm{K}}$ of K and define A as in (a). $\widehat{\mathrm{K}}$ can be shown to be algebraically closed, but it is not necessary in this exercise).
9.-Consider the rings

$$
\mathrm{A}=\mathbb{R}[\mathrm{X}, \mathrm{Y}] /\left(\mathrm{X}^{2}+\mathrm{Y}^{2}-(1), \quad \mathrm{B}=\mathbb{C}[\mathrm{X}, \mathrm{Y}] /\left(\mathrm{X}^{2}+\mathrm{Y}^{2}-1\right)\right.
$$

(a) Show that $B$ is isomorphic to $\mathbb{C}\left[Z, \frac{1}{Z}\right]$ (subring of $\mathbb{C}(Z)$ ), by setting $Z=$ $\mathrm{X}+i \mathrm{Y}$. What is the image of A under this isomorphism?
(b) Show that every maximal ideal of B is of the form $\mathfrak{n}_{x, y}=\{f \in \mathrm{~B} \mid f(x, y)=$ $0\}$, where $(x, y)$ is a point of $\mathbb{C}^{2}$ such that $x^{2}+y^{2}=1$ (the reader may use the injection from $\mathbb{C}[\mathrm{X}]$ to B). Show that all nonzero ideals in B are of the form $\mathfrak{n}_{x_{1}, y_{1}}^{k_{1}} \cdots \mathfrak{n}_{x_{r}, y_{r}}^{k_{r}}$.
(c) For $(x, y) \in \mathbb{C}^{2}, x^{2}+y^{2}=1$, set $\mathfrak{m}_{x, y}=\mathfrak{n}_{x, y} \cap \mathrm{~A}$. What is the codimension of $\mathfrak{m}_{x, y}$ as a real vector space in A? Show that every nonzero ideal in A is of the form $\mathfrak{m}_{x_{1}, y_{1}}^{k_{1}} \cdots \mathfrak{m}_{x_{r}, y_{r}}^{k_{r}}$, and that such an ideal is principal if and only if the sum of the $k_{i}$ corresponding to real pairs is even.
(d) Show that A is not a PID, but that B is, and is even Euclidean.
(e) Consider the rings

$$
\mathrm{A}^{\prime}=\mathbb{R}[\mathrm{X}, \mathrm{Y}] /\left(\mathrm{X}^{2}+\mathrm{Y}^{2}+1\right), \quad \mathrm{B}^{\prime}=\mathbb{C}[\mathrm{X}, \mathrm{Y}] /\left(\mathrm{X}^{2}+\mathrm{Y}^{2}+1\right)
$$

Show that $\mathrm{B}^{\prime}$ is isomorphic to B , but that $\mathrm{A}^{\prime}$ is not isomorphic to $A$. Describe the ideals of $A^{\prime}$ and $B^{\prime}$. Show that $A^{\prime}$ is a PID.
(f) What are the invertible elements of $\mathrm{B}^{\prime}$ ? of $\mathrm{A}^{\prime}$ ? Show that $\mathrm{A}^{\prime}$ is not Euclidean (use 3.1, Exercise 14).
10.-For any subset V of $\mathbb{C}^{n}$, set

$$
\mathscr{J}_{\mathbb{R}}(\mathrm{V})=\mathbb{R}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] \cap \mathscr{J}(\mathrm{V})
$$

and for any ideal I of $\mathbb{R}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$, set $\mathrm{V}_{\mathbb{C}}(\mathrm{I})=\mathrm{V}(\mathbb{C} \cdot \mathrm{I})$. Show that for any algebraic subset V of $\mathbb{C}^{n}, \mathrm{~V}_{\mathbb{C}}\left(\mathscr{f}_{\mathbb{R}}(\mathrm{V})\right)=\mathrm{V} \cup \tau(\mathrm{V})$, where $\tau: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ is the conjugation $\left(z_{1}, \ldots, z_{n}\right) \mapsto\left(\overline{z_{1}}, \ldots, \overline{z_{n}}\right)$.
11.-The aim is to give another proof of the fact that the ring $K\left[X_{1}, \ldots, X_{n}\right]$ is Noetherian when $K$ is an infinite field.
(a) Let $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ be a polynomial of degree $d$. Show that there is an automorphism $\varphi$ of $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ such that the coefficient of $\mathrm{X}_{n}^{d}$ in $\varphi(\mathrm{P})$ is nonzero (take a point $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathrm{K}^{n}$ such that $\mathrm{P}_{d}\left(x_{1}, \ldots, x_{n}\right) \neq 0$, where $\mathrm{P}_{d}$ denotes the homogeneous part of degree $d$ of P , and effect a change of coordinates in $\mathrm{K}^{n}$ such that $x$ becomes the last basis vector).
(b) Let $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ be a polynomial of degree $d$, where the coefficient of $\mathrm{X}_{n}^{d}$ is nonzero. Show that $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] /(\mathrm{P})$ is a free module of rank $d$ over $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]$.
(c) Let A be a ring. Show that, if $\mathrm{A} /(f)$ is Noetherian for all nonzero $f \in \mathrm{~A}$, then A is Noetherian.
(d) Conclude by induction on $n$ (use 3.6.2, Example 3).
12.-(a) Let K be a field and $f \in \mathrm{~K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$ a formal series of degree $d$ (i.e. such that the nonzero term of lowest degree is of degree $d$ ). Suppose that the coefficient of $\mathrm{X}_{n}^{d}$ is nonzero. Show that every formal series $g \in \mathrm{~K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$ can be uniquely written as $g=f q+\mathrm{R}$, with $q \in \mathrm{~K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$ and $\mathrm{R} \in \mathrm{K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n-1}\right]\right]\left[\mathrm{X}_{n}\right]$ a polynomial of degree $<d$.
(b) Suppose that K is infinite. Show that as in the previous exercise $\mathrm{K}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]$ is Noetherian.
(c) Show that the ring $\mathbb{R}\left\{\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right\}$ is Noetherian.
13. (Nullstellensatz for an uncountable field)—Let K be an uncountable field (for example $\mathbb{C}$ ).
(a) Show that the field $K(X)$ (field of fractions of $K[X]$ ) is a $K$-vector space of uncountable dimension (as $a$ runs through K , the $\frac{1}{\mathrm{X}-a}$ are linearly independent). Deduce that if a finite dimensional K -algebras is a field then it is an algebraic extension of K.
(b) Suppose that K is algebraically closed. Show that all maximal ideals in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ are of the form

$$
\mathfrak{m}_{\left(x_{1}, \ldots, x_{n}\right)}=\left\{\mathrm{P} \in \mathrm{~K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] \mid \mathrm{P}\left(x_{1}, \ldots, x_{n}\right)=0\right\}
$$

with $\left(x_{1}, \ldots, x_{n}\right) \in \mathrm{K}^{n}$. Deduce that for an ideal I in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right], \mathrm{V}(\mathrm{I})=\varnothing \Rightarrow \mathrm{I}=$ $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$.
(c) Continue to assume that K is algebraically closed; let I be an ideal in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ and set $\mathrm{A}=\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] / \mathrm{I}$. Let $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ be a polynomial vanishing at every point of $\mathrm{V}(\mathrm{I})$. Let J denote the ideal in $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}, \mathrm{Z}\right]$ generated by I and $(1-\mathrm{PZ})$. Show that $\mathrm{V}(\mathrm{J})$ is empty in $\mathrm{K}^{n+1}$.

Deduce that the image of $1-\mathrm{PZ}$ in $\mathrm{A}[\mathrm{Z}]$ is invertible, and that the image of P in A is nilpotent.
(d) Conclude.
14. (Nullstellensatz, general case)- Let $k$ be a field. The aim is to prove that a finitely generated field extension of $k$ (as a $k$-algebra) is an algebraic extension of $k$.
(A) Let L be a field, A a subring of L and $\mathrm{K} \subset \mathrm{L}$ the field of fractions of A .
(a) Show that if K is a finitely generated A -module, then $\mathrm{K}=\mathrm{A}$.
(b) Show that if $L$ is a finitely generated A-module, then $K=A$ (since $K$ is a direct factor in L as A-module).
(c) Show that if L is a finitely generated A -algebra and an algebraic extension of K , then there exists nonzero $b \in \mathrm{~A}$ such that $\mathrm{K}=\mathrm{A}\left[b^{-1}\right]$.
(B) (a) Set $\mathrm{A}=k\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]$ and $\mathrm{K}=k\left(\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right)$. Show that, if $n>0$, there are no $b \in \mathrm{~A}$ such that $\mathrm{K}=\mathrm{A}\left[b^{-1}\right]$.
(b) Let L be a finitely generated extension de $k$ (as a $k$-algebra). There are algebraically independent elements $x_{1}, \ldots, x_{n}$ in L such that L is an algebraic extension of $\mathrm{K}=k\left(x_{1}, \ldots, x_{n}\right)$. Show that $n=0$ necessarily holds. Conclude and continue as in the previous exercise.
15.-(a) Let $\mathrm{K}=\left\{a_{1}, \ldots a_{q}\right\}$ be a finite field. Show that the polynomial $\left(\mathrm{X}-a_{1}\right) \ldots\left(\mathrm{X}-a_{q}\right) \in \mathrm{K}[\mathrm{X}]$ is nonzero, but that the map from K to K defined by it is null.
(b) Let I be an infinite set, K the field $\mathbb{Z} /(2)$ and B the K -algebra $\mathrm{K}^{\mathrm{I}}$. Show that the map $\mathrm{P} \mapsto \mathrm{P}_{\mathrm{B}}$ from $\mathrm{K}[\mathrm{X}]$ to $\mathrm{B}^{\mathrm{B}}$ is not injective, although B is infinite (consider $\left.X^{2}+X\right)$.

### 3.8 Tensor Products

### 3.8.1

If $\mathrm{E}, \mathrm{F}, \mathrm{G}$ are three A -modules, let $\mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})$ be the module of bilinear maps from $\mathrm{E} \times \mathrm{F}$ to G .

Theorem and Definition Let E and F be A-modules. Then the covariant functor $\mathrm{G} \mapsto \mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})$ from $\mathrm{A}-\mathscr{M} \circ$ d to $\mathscr{E n}$ ns is representable. The tensor product of E and F is a representative of this functor.

Comments In other words, a tensor product of E and of F is an A -module T with a bilinear map $\theta$ from $\mathrm{E} \times \mathrm{F}$ to T , satisfying the following property: for all A -modules G and all bilinear maps $f$ from $\mathrm{E} \times \mathrm{F}$ to G , there is a unique linear map $\bar{f}: \mathrm{T} \rightarrow \mathrm{G}$ such that $f=\bar{f} \circ \theta$.
Notation The module T is written $\mathrm{E} \otimes_{\mathrm{A}} \mathrm{F}$, or simply $\mathrm{E} \otimes \mathrm{F}$, and the map $\theta(x, y) \mapsto x \otimes y$.
Cynical remark The different representatives of the functor $\mathrm{G} \mapsto \mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})$ can be mutually identified 2.4.11, hence they can be written $\mathrm{E} \otimes \mathrm{F}$, and this can be called "the" tensor product of E and F.

We give two proofs of the theorem.
$I^{s t}$ proof Let L be the free A-module over the set $\mathrm{E} \times \mathrm{F}$, and $\left(e_{x, y}\right)_{(x, y) \in \mathrm{E} \times \mathrm{F}_{-}}$the canonical basis for L. For all maps $f$ from $\mathrm{E} \times \mathrm{F}$ to an A-module G , write $\bar{f}$ for the homomorphism from L to G defined by $\bar{f}\left(e_{x, y}\right)=f(x, y)$. The map $f \mapsto \bar{f}$ is a bijection from the set of maps from $\mathrm{E} \times \mathrm{F}$ to G onto $\operatorname{Hom}(\mathrm{L}, \mathrm{G})$. The map $f$ is bilinear if and only if

$$
\begin{aligned}
& (\forall a \in \mathrm{~A})(\forall b \in \mathrm{~A})(\forall x \in \mathrm{E})(\forall y \in \mathrm{E})(\forall z \in \mathrm{~F})(\forall t \in \mathrm{~F}) \\
& \quad \bar{f}\left(e_{a x+b y, z}\right)-a \bar{f}\left(e_{x, z}\right)-b \bar{f}\left(e_{y, z}\right)=\bar{f}\left(e_{x, a z+b t}\right)-a \bar{f}\left(e_{x, z}\right)-b \bar{f}\left(e_{x, t}\right)=0,
\end{aligned}
$$

or if and only if $\bar{f}$ vanishes on every submodule N of L generated by elements of the form

$$
e_{a x+b y, z}-a e_{x, z}-b e_{y, z} \text { ou } e_{x, a z+b t}-a e_{x, z}-b e_{x, t}
$$

Hence $\mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})$ can be identified with $\operatorname{Hom}(\mathrm{T}, \mathrm{G})$, where $\mathrm{T}=\mathrm{L} / \mathrm{N}$. More precisely, the functor B is represented by $(\mathrm{T}, \theta)$, where $\theta(x, y)$ is the class of $e_{x, y}$.
$2^{\text {nd }}$ proof (a) Suppose first that E is free. If $\mathrm{E}=\mathrm{A}^{(\mathrm{I})}$, then $\mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})=\operatorname{Hom}\left(\mathrm{A}^{(\mathrm{I})}\right.$, $\operatorname{Hom}(\mathrm{F}, \mathrm{G}))=(\operatorname{Hom}(\mathrm{F}, \mathrm{G}))^{\mathrm{I}}=\operatorname{Hom}\left(\mathrm{F}^{(\mathrm{I})}, \mathrm{G}\right)$, and the functor $\mathrm{G} \mapsto \mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{G})$ is represented by $\mathrm{F}^{(\mathrm{I})}$.
(b) General case. Let $\mathrm{A}^{(\mathrm{J})} \xrightarrow{\delta} \mathrm{A}^{(\mathrm{I})} \xrightarrow{\varepsilon} \mathrm{E} \rightarrow 0$ be a presentation of E . As B $(\mathrm{M}, \mathrm{F} ; \mathrm{G})=\operatorname{Hom}(\mathrm{M}, \operatorname{Hom}(\mathrm{F}, \mathrm{G}))$ for all modules M , there is an exact sequence

$$
\begin{aligned}
& 0 \rightarrow \mathrm{~B}(\mathrm{E}, \mathrm{~F} ; \mathrm{G}) \xrightarrow{\varepsilon^{*}} \mathrm{~B}\left(\mathrm{~A}^{(\mathrm{I})}, \mathrm{F} ; \mathrm{G}\right) \xrightarrow{\delta^{*}} \mathrm{~B}\left(\mathrm{~A}^{(\mathrm{J})}, \mathrm{F} ; \mathrm{G}\right) \\
& \operatorname{Hom}\left(\mathrm{F}^{(\mathrm{I})}, \mathrm{G}\right)
\end{aligned} \quad \operatorname{Hom}\left(\mathrm{F}^{(\mathrm{J})}, \mathrm{G}\right) .
$$

By the Yoneda lemma 2.4.3, the morphism $\delta^{*}: \operatorname{Hom}\left(\mathrm{F}^{(\mathrm{I})}, \mathrm{G}\right) \rightarrow \operatorname{Hom}\left(\mathrm{F}^{(\mathrm{J})}, \mathrm{G}\right)$, functorial at G , arises from a morphism $\delta_{*}: \mathrm{F}^{(\mathrm{J})} \rightarrow \mathrm{F}^{(\mathrm{I})}$. Let T be the cokernel of $\delta_{*}$. For all G, there is an exact sequence

$$
0 \rightarrow \operatorname{Hom}(\mathrm{~T}, \mathrm{G}) \rightarrow \operatorname{Hom}\left(\mathrm{F}^{(\mathrm{I})}, \mathrm{G}\right) \xrightarrow{\delta^{*}} \operatorname{Hom}\left(\mathrm{~F}^{(\mathrm{J})}, \mathrm{G}\right)
$$

$\operatorname{SoHom}(T, G)=B(E, F ; G)$, and $T$ represents the functor $G \mapsto B(E, F ; G)$.
Remark The construction given in the first proof is intrinsic. That given in the second one shows that, for example, if E and F are finitely generated (resp. have finite presentation), the same holds for $\mathrm{E} \otimes \mathrm{F}$.

### 3.8.2

The elements of $\mathrm{E} \otimes \mathrm{F}$ are called tensors. A tensor $x \otimes y$, with $x \in \mathrm{E}$ and $y \in \mathrm{~F}$ is called a simple tensor.

Proposition The module $\mathrm{E} \otimes \mathrm{F}$ is generated by the simple tensors.
Proof Let $\theta: \mathrm{E} \times \mathrm{F} \rightarrow \mathrm{E} \otimes \mathrm{F}$ be the map $(x, y) \mapsto x \otimes y$, T the submodule of $\mathrm{E} \otimes \mathrm{F}$ generated by the simple tensors, and $\chi: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{E} \otimes \mathrm{F} / \mathrm{T}$ the canonical homomorphism. Then $\chi \circ \theta=0=0 \circ \theta$, and so $\chi=0$ by uniqueness of the universal property, which in turn gives $\mathrm{E} \otimes \mathrm{F} / \mathrm{T}=0$.
(This also follows by considering the construction given in the First Proof of Theorem 3.8.1.)

The product of a simple tensor and of a scalar is a simple tensor: indeed $a(x \otimes y)=$ $(a x) \otimes y$. Hence all tensors are the sum of finitely many simple ones. The rank of a tensor $t \in \mathrm{E} \otimes \mathrm{F}$ is the smallest integer $r$ for which $t$ is the sum of $r$ simple tensors.

If $f$ is a bilinear map from $\mathrm{E} \times \mathrm{F}$ to a module G , let $x \otimes y \mapsto f(x, y)$ denote the linear map of $\mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{G}$ corresponding to it, i.e. the unique linear map $\bar{f}: \mathrm{E} \otimes \mathrm{F} \rightarrow$ G such that $(\forall x \in \mathrm{E})(\forall y \in \mathrm{~F}) \bar{f}(x \otimes y)=f(x, y)$. For example, if $\mathrm{M}, \mathrm{N}, \mathrm{P}$ are three A-modules, then it is the map $(f \otimes g) \mapsto g \circ f$ from $\operatorname{Hom}(\mathbf{M}, \mathbf{N}) \otimes \operatorname{Hom}(\mathbf{N}, \mathrm{P})$ to $\operatorname{Hom}(\mathrm{M}, \mathrm{P})$.

Remark To show that a property holds for every element of $\mathrm{E} \otimes \mathrm{F}$, checking it for all simple tensors is not always sufficient. For example, $f(x \otimes y) \neq 0$ for $x \neq 0$ and $y \neq 0$ does not necessarily imply that the map $f: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{G}$ is injective. (Example: multiplication $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C} \rightarrow \mathbb{C}$ or see 3.8, Exercise 3.)

### 3.8.3 Functoriality

Let $\mathrm{E}, \mathrm{E}^{\prime}, \mathrm{F}, \mathrm{F}^{\prime}$ be A-modules, $f: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ and $g: \mathrm{F} \rightarrow \mathrm{F}^{\prime}$ homomorphisms. There is a linear map $x \otimes y \mapsto f(x) \otimes g(y)$ from $\mathrm{E} \otimes \mathrm{F}$ to $\mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}$, written $f \otimes g f \otimes g$ of $\operatorname{Hom}\left(\mathrm{E}, \mathrm{E}^{\prime}\right) \otimes \operatorname{Hom}\left(\mathrm{F}, \mathrm{F}^{\prime}\right)$. However these two elements can be identified when the Kronecker homomorphism (see 3.9) is injective. corresponding to $(x, y) \mapsto f(x) \otimes$ $g(y)$ from $\mathrm{E} \times \mathrm{F}$ to $\mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}$. In particular, $1_{\mathrm{E}} \otimes 1_{\mathrm{F}}=1_{\mathrm{E} \otimes \mathrm{F}}$ and, for homomorphisms $f^{\prime}: \mathrm{E}^{\prime} \rightarrow \mathrm{E}^{\prime \prime}$ and $g^{\prime}: \mathrm{F}^{\prime} \rightarrow \mathrm{F}^{\prime \prime},\left(f^{\prime} \circ f\right) \otimes\left(g^{\prime} \circ g\right)=\left(f^{\prime} \otimes g^{\prime}\right) \circ(f \otimes g)$.

### 3.8.4 Symmetry

Let E and F be A -modules. The linear map $x \otimes y \mapsto y \otimes x$ from $\mathrm{E} \otimes \mathrm{F}$ to $\mathrm{F} \otimes \mathrm{E}$ is an isomorphism, called symmetry. For $\mathrm{F}=\mathrm{E}$, the symmetry $\sigma: \mathrm{E} \otimes \mathrm{E} \rightarrow \mathrm{E} \otimes \mathrm{E}$ is not the identity, but $\sigma \circ \sigma=1_{\mathrm{E} \otimes \mathrm{E}}$. In general $\mathrm{E} \otimes \mathrm{F}$ and $\mathrm{F} \otimes \mathrm{E}$ not identified under $\sigma$.

### 3.8.5 Identity Object

Let E be an A-module. The maps $a \otimes x \mapsto a x$ from $\mathrm{A} \otimes \mathrm{E}$ to E and $x \mapsto 1 \otimes x$ from E to $\mathrm{A} \otimes \mathrm{E}$ are converse isomorphisms which enable us to identify $\mathrm{A} \otimes \mathrm{E}$ with E . Similarly, $\mathrm{E}=\mathrm{E} \otimes \mathrm{A}$.

### 3.8.6 Associativity

For finitely many A-modules $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ and an A-module G , let $\mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}} ; \mathrm{G}\right)$ denote the module of multilinear maps from $\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$ to G .

Proposition and Definition (a) For every finite family $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ of A-modules, the covariant functor $\mathrm{G} \mapsto \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}} ; \mathrm{G}\right)$ is representable. If $(\mathrm{T}, \theta)$ is a representative of this functor, then set $\bigotimes_{i \in \mathrm{I}} \mathrm{E}_{i}=\mathrm{T}$ and, for $\left(x_{i}\right)_{i \in \mathrm{I}} \in \prod_{i \in \mathrm{I}} \mathrm{E}_{i}, \otimes_{i \in \mathrm{I}} x_{i}=\theta\left(\left(x_{i}\right)_{i \in \mathrm{I}}\right)$.
(b) If $\mathrm{I}=\mathrm{J} \cup \mathrm{K}$ with $\mathrm{J} \cap \mathrm{K}=\varnothing$, the linear map $\otimes_{i \in \mathrm{I}} x_{i} \mapsto\left(\otimes x_{i}\right) \otimes\left(\otimes x_{i}\right)$ from $\bigotimes_{i \in \mathrm{I}} \mathrm{E}_{i}$ to $\left(\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}\right) \otimes\left(\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}\right)$ is an isomorphism.

Proof Let $\mathrm{I}=\mathrm{J} \cup \mathrm{K}$ with $\mathrm{J} \cap \mathrm{K}=\varnothing$, and suppose that the functors $\mathrm{G} \mapsto \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{~J}} ; \mathrm{G}\right)$ and $\mathrm{G} \mapsto \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{~K}} ; \mathrm{G}\right)$ are representable. Show that the functor $\mathrm{G} \mapsto \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}} ; \mathrm{G}\right)$ is represented by $\left(\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}\right) \otimes\left(\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}\right)$ equipped with the multilinear map $\left(x_{i}\right)_{i \in \mathrm{I}} \mapsto\left(\otimes_{i \in \mathrm{~J}} x_{i}\right) \otimes\left(\otimes_{i \in \mathrm{~K}} x_{i}\right)$.

For all A-modules G,

$$
\begin{aligned}
\mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}} ; \mathrm{G}\right) & =\mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{~J}} ; \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{~K}} ; \mathrm{G}\right)\right) \\
& =\mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \mathrm{~J}} ; \operatorname{Hom}\left(\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}, \mathrm{G}\right)\right) \\
& =\operatorname{Hom}\left(\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}, \operatorname{Hom}\left(\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}, \mathrm{G}\right)\right) \\
& =\mathrm{B}\left(\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}, \bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i} ; \mathrm{G}\right)=\operatorname{Hom}\left(\left(\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}\right) \otimes\left(\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}\right), \mathrm{G}\right),
\end{aligned}
$$

proving the existence of $\bigotimes_{i \in \mathrm{I}} \mathrm{E}_{i}$ as well as (b) if $\bigotimes_{i \in \mathrm{~J}} \mathrm{E}_{i}$ and $\bigotimes_{i \in \mathrm{~K}} \mathrm{E}_{i}$ are assumed to exist.

For $\mathrm{I}=\varnothing, \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in \varnothing} ; \mathrm{G}\right)=\mathrm{G}=\operatorname{Hom}(\mathrm{A}, \mathrm{G})$, and so $\bigotimes_{i \in \varnothing} \mathrm{E}_{i}=\mathrm{A}$.
For $\mathrm{I}=\{a\}, \mathrm{B}\left(\left(\mathrm{E}_{i}\right)_{i \in\{a\}} ; \mathrm{G}\right)=\operatorname{Hom}\left(\mathrm{E}_{a}, \mathrm{G}\right)$, and so $\bigotimes_{i \in\{a\}} \mathrm{E}_{i}=\mathrm{E}_{a}$.
(a) now follows by induction on the number of elements of I .

Corollary For any A-modules E, F, G,

$$
\mathrm{E} \otimes(\mathrm{~F} \otimes \mathrm{G})=\mathrm{E} \otimes \mathrm{~F} \otimes \mathrm{G}=(\mathrm{E} \otimes \mathrm{~F}) \otimes \mathrm{G} .
$$

### 3.8.7 Right Exactness, Distributivity, Passing to Direct Limits

Proposition Let F be an A-module.
(a) Let $f: \mathrm{E}^{\prime} \rightarrow \mathrm{E}$ be a homomorphism of A -modules. Then

$$
\operatorname{Coker}\left(f \otimes 1_{\mathrm{F}}\right)=(\operatorname{Coker} f) \otimes \mathrm{F}
$$

(b) Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ be a family of A-modules. Then

$$
\left(\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}\right) \otimes \mathrm{F}=\bigoplus_{i \in \mathrm{I}}\left(\mathrm{E}_{i} \otimes \mathrm{~F}\right)
$$

(c) Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ a direct system of A-modules. Then

$$
\left(\lim _{\longrightarrow} \mathrm{E}_{i}\right) \otimes \mathrm{F}=\underset{\longrightarrow}{\lim }\left(\mathrm{E}_{i} \otimes \mathrm{~F}\right) .
$$

Proof For all A-modules G,

$$
\operatorname{Hom}(\mathrm{E} \otimes \mathrm{~F}, \mathrm{G})=\operatorname{Hom}(\mathrm{E}, \operatorname{Hom}(\mathrm{~F}, \mathrm{G}))
$$

In other words, the functor $\mathrm{F} \mapsto \mathrm{E} \otimes \mathrm{F}$ from $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$ to $\mathrm{A}-\mathscr{M} \sigma d$ has a right adjoint functor $\mathrm{G} \mapsto \operatorname{Hom}(\mathrm{F}, \mathrm{G})$. The proposition then follows from (2.7.3, Corollary).

Corollary (a) If $\mathrm{E}^{\prime} \rightarrow \mathrm{E} \rightarrow \mathrm{E}^{\prime \prime} \rightarrow 0$ is an exact sequence of A -modules, then $\mathrm{E}^{\prime} \otimes$ $\mathrm{F} \rightarrow \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{E}^{\prime \prime} \otimes \mathrm{F} \rightarrow 0$ is an exact sequence.
(b) Let $f: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ and $g: \mathrm{F} \rightarrow \mathrm{F}^{\prime}$ be surjective homomorphisms, then $f \otimes g:$ $\mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}$ is surjective.
(c) Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ and $\left(\mathrm{F}_{j}\right)_{j \in \mathrm{~J}}$ be two families of A-modules. Then

$$
\left(\bigoplus_{i \in \mathrm{I}} \mathrm{E}_{i}\right) \otimes\left(\bigoplus_{j \in \mathrm{~J}} \mathrm{~F}_{j}\right)=\bigoplus_{(i, j) \in \mathrm{I} \times \mathrm{J}} \mathrm{E}_{i} \otimes \mathrm{~F}_{j}
$$

(d) Let $\left(\mathrm{E}_{i}\right)_{i \in \mathrm{I}}$ and $\left(\mathrm{F}_{j}\right)_{j \in \mathrm{~J}}$ be direct systems of A-modules. Then

$$
\left(\underset{\longrightarrow}{\lim } \mathrm{E}_{i}\right) \otimes\left(\underset{\longrightarrow}{\lim } \mathrm{F}_{j}\right)=\underset{(i, \overrightarrow{j) \in \mathrm{I} \times \mathrm{J}}}{\lim } \mathrm{E}_{i} \otimes \mathrm{~F}_{j}
$$

Remarks (1) Let F be an A-module and $f: \mathrm{E}^{\prime} \rightarrow \mathrm{E}$ an injective homomorphism. The homomorphism $f \otimes 1_{\mathrm{F}}: \mathrm{E}^{\prime} \otimes \mathrm{F} \rightarrow \mathrm{E} \otimes \mathrm{F}$ is not in general injective. For example, suppose that A is integral and take $\mathrm{E}=\mathrm{E}^{\prime}=\mathrm{A}$. For all nonzero $a \in \mathrm{~A}, a \cdot 1_{\mathrm{E}}$ is injective, but if F has torsion, then there exists $a \neq 0$ such that $a \cdot 1_{\mathrm{E}} \otimes 1_{\mathrm{F}}=a \cdot 1_{\mathrm{F}}$ is not injective.

In particular, if $\mathrm{E}^{\prime}$ is a submodule of E , then in general $\mathrm{E}^{\prime} \otimes \mathrm{F}$ cannot be identified with a submodule of $\mathrm{E} \otimes \mathrm{F}$. For example, for $\mathrm{A}=\mathbb{Z}, \mathrm{E}=\mathbb{Q}, \mathrm{E}^{\prime}=\mathbb{Z}, \mathrm{F}=\mathbb{Z} /(n)$, $\mathrm{E} \otimes \mathrm{F}=0$ but $\mathrm{E}^{\prime} \otimes \mathrm{F} \neq 0$.

However, if $\mathrm{E}^{\prime}$ is a direct factor of E , the module $\mathrm{E}^{\prime} \otimes \mathrm{F}$ is a direct factor of $\mathrm{E} \otimes \mathrm{F}$. For example, if $f: \mathrm{E}^{\prime} \rightarrow \mathrm{E}$ and $g: \mathrm{F}^{\prime} \rightarrow \mathrm{F}$ are injective vector space homomorphisms over a field K , the map $f \otimes g: \mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime} \rightarrow \mathrm{E} \otimes \mathrm{F}$ is injective, since $f$ and $g$ identify $\mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime}$ with direct factors of E and F respectively.
(2) Let $f: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ and $g: \mathrm{F} \rightarrow \mathrm{F}^{\prime}$ be homomorphisms. Then the modules $\operatorname{Coker}\left(f \otimes g: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}\right)$ and (Coker $\left.f\right) \otimes($ Coker $g)$ are not in general isomorphic. For example, for $\mathrm{A}=\mathbb{R}, \mathrm{E}=\mathbb{R}^{n}, \mathrm{E}^{\prime}=\mathbb{R}^{n+1}, \mathrm{~F}=\mathbb{R}^{p}, \mathrm{~F}^{\prime}=\mathbb{R}^{p+1}, f$ and $g$ injective, the above modules are vector spaces of dimension $n+p+1$ and 1 respectively.

### 3.8.8 Tensor Product of Free Modules

Proposition Let E and F be free A -modules with respective bases $\left(e_{i}\right)_{i \in \mathrm{I}}$ and $\left(f_{j}\right)_{j \in \mathrm{~J}}$. Then $\mathrm{E} \otimes \mathrm{F}$ is a free module and $\left(e_{i} \otimes f_{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ is a basis for $\mathrm{E} \otimes \mathrm{F}$.

Proof Taking $\mathrm{E}_{i}=\mathrm{F}_{j}=\mathrm{A}$ in Corollary (3.8.7, c), $\mathrm{A}^{(\mathrm{I})} \otimes \mathrm{A}^{(\mathrm{J})}=\mathrm{A}^{(\mathrm{I} \times \mathrm{J})}$, and the proposition follows.

Corollary 3.7 Let E and F be A-modules, $\left(x_{i}\right)_{i \in \mathrm{I}}$ and $\left(y_{j}\right)_{j \in \mathrm{~J}}$ generators of E and F respectively. Then $\left(x_{i} \otimes y_{j}\right)_{(i, j) \in \mathrm{I} \times \mathrm{J}}$ generate $\mathrm{E} \otimes \mathrm{F}$.

Corollary 3.8 The tensor product of two projective modules is projective.
Proof Let E and F be two projective A-modules. There are A-modules $\mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime}$ such that $\mathrm{E} \oplus \mathrm{E}^{\prime}$ and $\mathrm{F} \oplus \mathrm{F}^{\prime}$ are free. Then, $\left(\mathrm{E} \oplus \mathrm{E}^{\prime}\right) \otimes\left(\mathrm{F} \oplus \mathrm{F}^{\prime}\right)=(\mathrm{E} \otimes \mathrm{F}) \oplus(\mathrm{E} \otimes$ $\left.\mathrm{F}^{\prime}\right) \oplus\left(\mathrm{E}^{\prime} \otimes \mathrm{F}\right) \oplus\left(\mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}\right)$ is free. So $\mathrm{E} \otimes \mathrm{F}$ is projective.

### 3.8.9 Tensor Product of Matrices

Proposition and Definition Let $\mathrm{E}, \mathrm{E}^{\prime}, \mathrm{F}, \mathrm{F}^{\prime}$ be finitely generated free A -modules, $\left(e_{i}\right)_{i \in \mathrm{I}}, \quad\left(e_{i^{\prime}}^{\prime}\right)_{i^{\prime} \in \mathrm{I}^{\prime}},\left(f_{j}\right)_{j \in \mathrm{~J}}, \quad\left(f_{j^{\prime}}^{\prime}\right)_{j^{\prime} \in \mathrm{J}^{\prime}}$ bases for $\mathrm{E}, \mathrm{E}^{\prime}, \mathrm{F}, \mathrm{F}^{\prime}$ respectively. Let $u: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ and $v: \mathrm{F} \rightarrow \mathrm{F}^{\prime}$ be homomorphisms, $\mathrm{M}=\left(m_{i}^{i^{\prime}}\right)$ and $\mathrm{N}=\left(n_{j}^{j^{\prime}}\right)$ their respective matrices with respect to the given bases. Then the matrix of the homomorphism $u \otimes v: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}$ with respect to the bases $\left(e_{i} \otimes f_{j}\right)$ and $\left(e_{i^{\prime}}^{\prime} \otimes f_{j^{\prime}}^{\prime}\right)$ is $\left(t_{i, j}^{i^{\prime}, j^{\prime}}\right)$, where $t_{i, j}^{i^{\prime}, j^{\prime}}=m_{i}^{i^{\prime}} n_{j}^{j^{\prime}}$. This matrix is called the tensor product of M and N , and is written $\mathrm{M} \otimes \mathrm{N}$.
Proof Since $u\left(e_{i}\right)=\sum m_{i}^{i^{\prime}} e_{i^{\prime}}^{\prime}$ and $v\left(f_{j}\right)=\sum n_{j}^{j^{\prime}} f_{j^{\prime}}^{\prime}, u \otimes v\left(e_{i} \otimes f_{j}\right)=u\left(e_{i}\right) \otimes$ $v\left(f_{j}\right)=\sum_{i^{\prime}, j^{\prime}} m_{i}^{i^{\prime}} n_{j}^{j^{\prime}} e_{i^{\prime}}^{\prime} \otimes f_{j^{\prime}}^{\prime}$.

Corollary Let K be a field, M and N matrices with entries in K . Then the rank of $\mathrm{M} \otimes \mathrm{N}$ is the product of the ranks of M and N .

Proof We may assume M and N are as follows to within base changes

$$
\left(\begin{array}{llllllll}
1 & 0 & \ldots & 0 & 0 & \ldots & 0 \\
0 & 1 & \ldots & 0 & 0 & \ldots & 0 \\
\vdots & & & \vdots & \vdots & & & \vdots \\
0 & \ldots & \ldots & 1 & 0 & \ldots & 0 \\
0 & \ldots & \ldots & 0 & 0 & \ldots & 0 \\
0 & \ldots & \ldots & 0 & 0 & \ldots & 0
\end{array}\right)
$$

### 3.8.10 Tensor Product of Cyclic Modules

Proposition Let I and J be two ideals of A . Then

$$
(\mathrm{A} / \mathrm{I}) \otimes(\mathrm{A} / \mathrm{J})=\mathrm{A} /(\mathrm{I}+\mathrm{J})
$$

Proof Let $e$ and $f$ be the images of 1 in A/I and A/J. The module (A/I) $\otimes(\mathrm{A} / \mathrm{J})$ is generated by $e \otimes f$ (3.8.8, Corollary 3.7), and the annihilator of $e \otimes f$ contains I and J , and so $\mathrm{I}+\mathrm{J}$. On the other hand, the map $(a, b) \mapsto a b$ defined by passing to the quotient is a bilinear map from $(\mathrm{A} / \mathrm{I}) \times(\mathrm{A} / \mathrm{J})$ to $\mathrm{A} /(\mathrm{I}+\mathrm{J})$, and hence a homomorphism of $(\mathrm{A} / \mathrm{I}) \otimes(\mathrm{A} / \mathrm{J})$ mapping $e \otimes f$ onto the image $g$ of 1 . Therefore the annihilator of $e \otimes f$ is contained in that of $g$, i.e. in $\mathrm{I}+\mathrm{J}$; and so is equal to it.

### 3.8.11 Examples of Tensor Products over $\mathbb{Z}$

Consider the following $\mathbb{Z}$-modules: $\mathbb{Z}^{r}, \mathbb{Z} /(n), \mathbb{Q}, \mathbb{Q} / \mathbb{Z}, \mathbb{Z}^{(\mathbb{N})}, \mathbb{Z}^{\mathbb{N}}$. Their mutual tensor products are given by the following table,

| $\otimes$ | $\mathbb{Z}^{r}$ | $\mathbb{Z} /(n)$ | $\mathbb{Q}$ | $\mathbb{Q} / \mathbb{Z}$ | $\mathbb{Z}^{(\mathbb{N})}$ | $\mathbb{Z}^{\mathbb{N}}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mathbb{Z}^{r^{\prime}}$ | $\mathbb{Z}^{r r^{\prime}}$ |  |  |  |  |  |
| $\mathbb{Z} /\left(n^{\prime}\right)$ | $\left(\mathbb{Z} /\left(n^{\prime}\right)\right)^{r}$ | $\mathbb{Z} /(d)$ <br> $d=$ g.c.d. $\left(n, n^{\prime}\right)$ |  |  |  |  |
| $\mathbb{Q}$ | $\mathbb{Q}^{r}$ | 0 | $\mathbb{Q}$ |  |  |  |
| $\mathbb{Q} / \mathbb{Z}$ | $(\mathbb{Q} / \mathbb{Z})^{r}$ | 0 | 0 | 0 |  |  |
| $\mathbb{Z}^{(\mathbb{N})}$ | $\mathbb{Z}^{(\mathbb{N} \times\{1, \ldots, r\})}$ | $(\mathbb{Z} /(n))^{(\mathbb{N})}$ | $\mathbb{Q}^{(\mathbb{N})}$ | $(\mathbb{Q} / \mathbb{Z})^{(\mathbb{N})}$ | $\mathbb{Z}^{(\mathbb{N} \times \mathbb{N})}$ |  |
| $\mathbb{Z}^{\mathbb{N}}$ | $\mathbb{Z}^{\mathbb{N} \times\{1, \ldots, r\}}$ | $(\mathbb{Z} /(n))^{\mathbb{N}}$ | E | F | G | H |

where $\mathrm{E} \subset \mathbb{Q}^{\mathbb{N}}$ is the set of sequences of the form $\left(\frac{p_{i}}{q_{i}}\right)_{i \in \mathbb{N}}$, the sequence $\left(q_{i}\right)$ being bounded; $F$ is the image of $E$ in $(\mathbb{Q} / \mathbb{Z})^{\mathbb{N}} ; G \subset E^{\mathbb{N} \times \mathbb{N}}$ is the set of double sequences whose support is contained in a "strip", i.e. in a set $\{1, \ldots, k\} \times \mathbb{N} ; H \subset \mathbb{Z}^{\mathbb{N} \times \mathbb{N}}$ is the set of double sequences $(n, p) \mapsto u_{n, p}$ such that there exists a finitely generated submodule M of $\mathbb{Z}^{\mathbb{N}}$ for which, for all $n \in \mathbb{N}$, the partial sequence $u_{n}: p \mapsto u_{n, p}$ is in M .

Proof (a) Column of $\mathbb{Z}^{r}$ : for all $\mathbb{Z}$-modules $\mathrm{M}, \mathbb{Z}^{r} \otimes \mathrm{M}=\mathrm{M}^{r}$ by distributivity (3.8.7, b).
(b) Row of $\mathbb{Z}^{(\mathbb{N})}$ : for all $\mathbb{Z}$-modules $\mathrm{M}, \mathbb{Z}^{(\mathbb{N})} \otimes \mathrm{M}=\mathrm{M}^{(\mathbb{N})}$ by distributivity. The tensor products $\mathbb{Z}^{(\mathbb{N})} \otimes \mathbb{Z}^{r}$ and $\mathbb{Z}^{(\mathbb{N})} \otimes \mathbb{Z}^{(\mathbb{N})}$ are given by Proposition 3.8.8.
(c) Column of $\mathbb{Z} /(n)$ : For all Z-modules $\mathrm{M}, \mathbb{Z} /(n) \otimes \mathrm{M}=\mathrm{M} / n \mathrm{M}$.

Indeed, the exact sequence $\mathbb{Z} \xrightarrow{n} \mathbb{Z} \rightarrow \mathbb{Z} /(n) \rightarrow 0$ gives

$$
\mathrm{M} \xrightarrow{n} \mathrm{M} \rightarrow \mathbb{Z} /(n) \otimes \mathrm{M} \rightarrow 0
$$

For $\mathrm{M}=\mathbb{Z} /\left(n^{\prime}\right)$, let $\chi: \mathbb{Z} \rightarrow \mathbb{Z} /\left(n^{\prime}\right)$ be the canonical homomorphism. Then $n \mathrm{M}=$ $\chi(n \mathbb{Z})=\chi\left(n \mathbb{Z}+n^{\prime} \mathbb{Z}\right)=\left(n \mathbb{Z}+n^{\prime} \mathbb{Z}\right) / n^{\prime} \mathbb{Z}=d \mathbb{Z} / n^{\prime} \mathbb{Z}$ where $d=$ g.c.d. $\left(n, n^{\prime}\right)$, and so $\mathrm{M} / n \mathrm{M}=\left(\mathbb{Z} /\left(n^{\prime}\right)\right) /\left(d \mathbb{Z} /\left(n^{\prime}\right)\right)=\mathbb{Z} /(d)$.

For $\mathbb{M}=\mathbb{Q}, \mathbb{Q} \otimes \mathbb{Z} /(n)=\mathbb{Q} / n \mathbb{Q}=0$ since $n \mathbb{Q}=\mathbb{Q}$. Similarly, $\mathbb{Q} / \mathbb{Z} \otimes \mathbb{Z} /(n)=$ 0 since $n(\mathbb{Q} / \mathbb{Z})=\mathbb{Q} / \mathbb{Z}$.

For $\mathbb{M}=\mathbb{Z}^{\mathbb{N}}, \mathbb{Z}^{\mathbb{N}} \otimes \mathbb{Z} /(n)=\mathbb{Z}^{\mathbb{N}} / n \mathbb{Z}^{\mathbb{N}}=(\mathbb{Z} /(n))^{\mathbb{N}}$.
(d) $\mathbb{Q} \otimes \mathbb{Q}=\mathbb{Q}$. For $q \in \mathbb{N}^{*}$, consider the submodule $\mathbb{Z} \frac{1}{q}$ de $\mathbb{Q}$. Then $\mathbb{Q}=$ $\bigcup_{q \in \mathbb{N}^{*}} \mathbb{Z} \frac{1}{q}=\underset{\longrightarrow}{\lim } \mathbb{Z} \frac{1}{q}$, the direct limit being indexed by $\mathbb{N}^{*}$ ordered by divisibility. For all $q \in \mathbb{N}^{*}$, the homomorphism $x \otimes y \xrightarrow{\alpha_{q}} x y$ from $\mathbb{Q} \otimes\left(\mathbb{Z} \frac{1}{q}\right)$ to $\mathbb{Q}$ is an isomorphism thanks to the diagram


These isomorphisms form a coherent system, i.e. for all $q$ and $q^{\prime}$ such that $q$ divides $q^{\prime}$, the following diagrams

are commutative.

Hence, passing to the direct limit, we get an isomorphism $\alpha: \mathbb{Q} \otimes \mathbb{Q} \rightarrow \mathbb{Q}$ characterized by $\alpha(x \otimes y)=x y$.
(e) The tensor product of the exact sequence $0 \rightarrow \mathbb{Z} \xrightarrow{\iota} \mathbb{Q} \rightarrow \mathbb{Q} / \mathbb{Z} \rightarrow 0$ with $\mathbb{Q}$ is the following exact sequence $\mathbb{Q} \rightarrow \mathbb{Q} \rightarrow \mathbb{Q} \otimes \mathbb{Q} / \mathbb{Z} \rightarrow 0$, where the homomorphism $1 \otimes \iota: \mathbb{Q} \rightarrow \mathbb{Q}$ is the identity since the isomorphisms identifying $\mathbb{Q} \otimes \mathbb{Z} \rightarrow \mathbb{Q}$ with $\mathbb{Q} \otimes \mathbb{Q} \rightarrow \mathbb{Q}$ are defined by $1 \otimes 1 \mapsto 1$.
(f) The homomorphism $x \otimes\left(u_{i}\right)_{i \in \mathbb{N}} \mapsto\left(x u_{i}\right)_{i \in \mathbb{N}}$ from $\mathbb{Q} \otimes \mathbb{Z}^{\mathbb{N}}$ to $\mathbb{Q}^{\mathbb{N}}$ induced for all $q \in \mathbb{N}^{*}$ is an isomorphism from $\mathbb{Z} \cdot \frac{1}{q} \otimes \mathbb{Z}^{\mathbb{N}}$ onto the submodule $\mathbb{Z}^{\mathbb{N}} \cdot \frac{1}{q}$ of $\mathbb{Q}^{\mathbb{N}}$ consisting of the sequences in $\mathbb{Z} \cdot \frac{1}{q}$. Then

$$
\mathbb{Q} \otimes \mathbb{Z}^{\mathbb{N}}=\lim _{\longrightarrow} \mathbb{Z}^{\mathbb{N}} \frac{1}{q}=\bigcup_{q \in \mathbb{N}^{*}} \mathbb{Z}^{\mathbb{N}} \frac{1}{q}=\mathrm{E}
$$

(g) Column of $\mathbb{Q} / \mathbb{Z}$. The module $\mathbb{Q} / \mathbb{Z} \otimes \mathbb{Q} / \mathbb{Z}$ is a quotient of $\mathbb{Q} \otimes \mathbb{Q} / \mathbb{Z}=0$ and so is zero.

The tensor product of the exact sequence $0 \rightarrow \mathbb{Z} \rightarrow \mathbb{Q} \rightarrow \mathbb{Q} / \mathbb{Z} \rightarrow 0$ with $\mathbb{Z}^{\mathbb{N}}$ is $\mathbb{Z}^{\mathbb{N}} \rightarrow \mathrm{E} \rightarrow \mathbb{Z}^{\mathbb{N}} \otimes \mathbb{Q} / \mathbb{Z} \rightarrow 0$. As $\mathbb{Z}^{\mathbb{N}} \rightarrow \mathrm{E}$ is injective, $\mathbb{Z}^{\mathbb{N}} \otimes \mathbb{Q} / \mathbb{Z}=\mathrm{E} / \mathbb{Z}^{\mathbb{N}}$.
(h) $\mathbb{Z}^{\mathbb{N}} \otimes \mathbb{Z}^{(\mathbb{N})}=\left(\mathbb{Z}^{\mathbb{N}}\right)^{(\mathbb{N})}=\mathrm{G}$.
(i) For every $\mathbb{Z}$-module $\mathbf{M}$, let $\alpha_{\mathrm{M}}$ denote the homomorphisms $\left(u_{i}\right) \otimes m \mapsto\left(u_{i} m\right)$ from $\mathbb{Z}^{\mathbb{N}} \otimes \mathbf{M}$ to $\mathbf{M}^{\mathbb{N}}$. These define a morphism of functors, and $\alpha_{M}$ is an isomorphism if $M$ is free and finitely generated. Every finitely generated submodule $M$ of $\mathbb{Z}^{\mathbb{N}}$ is free (3.5.8, Corollary 3.2) and for such a submodule, $\mathbf{M}^{\mathbb{N}}$ can be identified with a submodule of $\left(\mathbb{Z}^{\mathbb{N}}\right)^{\mathbb{N}}=\mathbb{Z}^{\mathbb{N} \times \mathbb{N}}$. As $\mathbb{Z}^{\mathbb{N}}$ is the direct limit of its finitely generated submodules $M, \mathbb{Z}^{\mathbb{N}} \otimes \mathbb{Z}^{\mathbb{N}}=\xrightarrow{\lim } \mathbb{Z}^{\mathbb{N}} \otimes M=\bigcup_{M} M^{\mathbb{N}}=H$.

### 3.8.12 Extension of Scalars

Let $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism. Consider B as an $\mathrm{A}-m o d u l e$ (3.3.1, Example 3). For all A-modules $\mathrm{E}, \mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}$ can be endowed with a structure of B-module by $b \cdot\left(b^{\prime} \otimes m\right)=b b^{\prime} \otimes m$. With this structure, $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}$, is said to be the B -module obtained by the extension of scalars of A to B along $\varphi$. This B-module is written $\varphi_{*}(\mathrm{E})$.

If $f: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ is an A-module homomorphism, then the map $f_{*}=1_{\mathrm{B}} \otimes_{\mathrm{A}} f:$ $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E} \rightarrow \mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}^{\prime}$ is B -linear. This defines a functor $\varphi_{*}: \mathrm{A}-\mathscr{M}_{\mathrm{od}} \rightarrow \mathrm{B}-\mathscr{M} \operatorname{od}$ which commutes with cokernels, direct sums, and direct limits.

For a B-module F , let $\varphi^{*}(\mathrm{~F})$ be the A-module obtained by restricting scalars from B to A along $\varphi$ (3.3.2). This gives a functor $\varphi^{*}: \mathrm{B}-\mathscr{M} \sigma d \rightarrow \mathrm{~A}-\mathscr{M} \sigma d$. The modules $\varphi^{*}(\mathrm{~F})$ and F have the same underlying additive group.

Proposition The extension of scalars functor $\varphi_{*}: \mathrm{A}-\mathscr{M} \sigma d \rightarrow \mathrm{~B}-\mathscr{M} \sigma d$ is the left adjoint of the restriction of scalars functor $\varphi^{*}$.

In other words, for all A-modules E and all B-modules $\mathrm{F}, \operatorname{Hom}_{\mathrm{B}}\left(\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}, \mathrm{F}\right)=$ $\operatorname{Hom}_{\mathrm{A}}(\mathrm{E}, \mathrm{F})$. This is the set of $\varphi$-linear maps from E to F .

Proof For an A-linear map $f: \mathrm{E} \rightarrow \mathrm{F}$, set $\bar{f}(b \otimes x)=b f(x)$ for all $b \in \mathrm{~B}$ and $x \in \mathrm{E}$. For $g: \mathrm{B} \otimes_{\mathrm{A}} \mathrm{E} \rightarrow \mathrm{F}$, set $\bar{g}(x)=g(1 \otimes x)$ for all $x \in \mathrm{E}$. The maps $f \mapsto \bar{f}$ and $g \mapsto \bar{g}$ are mutual inverses.

Corollary If $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ and $\psi: \mathrm{B} \rightarrow \mathrm{C}$ are ring homomorphisms then, for all $\mathrm{A}-$ modules $\mathrm{E},(\psi \circ \varphi)_{*}(\mathrm{E})=\psi_{*}\left(\varphi_{*}(\mathrm{E})\right)$.

In other words, $\mathrm{C} \otimes_{\mathrm{B}}\left(\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}\right)=\mathrm{C} \otimes_{\mathrm{A}} \mathrm{E}$.
Proof The functor $\psi_{*} \circ \varphi_{*}$ is a left adjoint of $\varphi^{*} \circ \psi^{*}=(\psi \circ \varphi)^{*}$, and so may be identified with $(\psi \circ \varphi)_{*}$.

Examples (1) If $E$ is a real vector space, then $\mathbb{C} \otimes_{\mathbb{R}} E$ is a complex vector space called the complexification of E whose $\mathbb{C}$-dimension is equal to the $\mathbb{R}$-dimension of E .
(2) Let E be an A-module and I an ideal in A ; the ( $\mathrm{A} / \mathrm{I}$ )-module $\mathrm{A} / \mathrm{I} \otimes \mathrm{E}$ can be identified with $\mathrm{E} / \mathrm{IE}$, where IE is the submodule generated by the elements $a x, a \in \mathrm{I}$ and $x \in \mathrm{E}$. Indeed, IE is the image of $\mathrm{I} \otimes \mathrm{E}$ in E .

### 3.8.13 Extension of Scalars to the Field of Fractions

Suppose that A is an integral domain, K its field of fractions and $\iota: \mathrm{A} \rightarrow \mathrm{K}$ the canonical injection.

Proposition For all A-modules E , the kernel of $\iota \otimes 1_{\mathrm{E}}: \mathrm{E} \rightarrow \mathrm{K} \otimes_{\mathrm{A}} \mathrm{E}$ is the torsion submodule of E .

Proof For $q \in \mathrm{~A}^{*}$, consider the submodule $\mathrm{A} \frac{1}{q}$ of K and let $\iota_{q}: \mathrm{A} \rightarrow \mathrm{A} \frac{1}{q}$ be the inclusion homomorphism. Tensoring the commutative diagram

with E gives

and so $\operatorname{Ker}\left(\iota_{q} \otimes 1_{\mathrm{E}}\right)=\{x \in \mathrm{E} \mid q x=0\}$.

Passing to the direct limit over the ordered set $\operatorname{Mon}^{*}(\mathrm{~A})$, we get $\operatorname{Ker}\left(\iota \otimes 1_{\mathrm{E}}\right)=$ $\bigcup_{q}\{x \in \mathrm{E} \mid q x=0\}$, and this set is the torsion submodule of E .

### 3.8.14 Rank of a Module

Let A be an integral domain and K its field of fractions.
The rank of a finitely generated A-module M , written rk M , is the dimension of the $K$-vector space $K \otimes_{A} M$. If $A$ is principal, $M$ is isomorphic to some module $\mathrm{A}^{r} \oplus \mathrm{~T}$, where T is a torsion module (3.5.8), and the rank of M is $r$.

If $0 \rightarrow \mathrm{M} \rightarrow \mathrm{N} \rightarrow \mathrm{P} \rightarrow 0$ is a short exact sequence of finitely generated A-modules, then $\mathrm{rkN}=\mathrm{rk} \mathrm{M}+\mathrm{rk}$ P. Indeed, K being the direct limit of free modules of rank 1, there is an exact sequence

$$
0 \rightarrow \mathrm{~K} \otimes_{\mathrm{A}} \mathrm{M} \rightarrow \mathrm{~K} \otimes_{\mathrm{A}} \mathrm{~N} \rightarrow \mathrm{~K} \otimes_{\mathrm{A}} \mathrm{P} \rightarrow 0
$$

If $f: \mathrm{M} \rightarrow \mathrm{N}$ is a finitely generated A -module homomorphism, let $\mathrm{rk} f$ denote the rank of $\operatorname{Im} f$. Then $\mathrm{rk} \mathrm{N}=\mathrm{rk} f+\operatorname{rk}(\operatorname{Coker} f)$ and, assuming $\operatorname{Ker} f$ is finitely generated (which is always the case if A is Noetherian), $\mathrm{rkM}=\operatorname{rk}(\operatorname{Ker} f)+\operatorname{rk} f$.

Corollary Let A be an integral domain. Every finitely generated torsion free A-module is isomorphic to a submodule of a finitely generated free module of the same rank.

Proof Let K be the field of fractions of A, E a finitely generated torsion free Amodule generated by the finite family $\left(x_{i}\right)$ in E , and $\left(e_{j}\right)$ a basis for the finite dimensional K-vector space $\mathrm{K} \otimes_{\mathrm{A}} \mathrm{E}$. Define the scalars $a_{i}^{j} \in \mathrm{~K}$ by $\iota\left(x_{i}\right)=\sum a_{i}^{j} e_{j}$, and let $d$ be their common denominator. The injection $\iota: \mathrm{E} \rightarrow \mathrm{K} \otimes_{\mathrm{A}} \mathrm{E}$ maps E into the free A-submodule $\mathrm{K} \otimes_{\mathrm{A}} \mathrm{E}$ generated by $\frac{1}{d} e_{j}$.

### 3.8.15 Tensor Products of Algebras

Let E and F be A -algebras. There is a unique A -algebra structure on the A -module $\mathrm{E} \otimes \mathrm{F}$ such that $(x \otimes y)\left(x^{\prime} \otimes y^{\prime}\right)=x x^{\prime} \otimes y y^{\prime}$. If E and F are associative (resp. unital, resp. commutative), the same holds for $\mathrm{E} \otimes \mathrm{F}$. If E and F are unital, the maps $\iota_{1}$ : $\mathrm{E} \rightarrow \mathrm{E} \otimes \mathrm{F}$ and $\iota_{2}: \mathrm{F} \rightarrow \mathrm{E} \otimes \mathrm{F}$ defined by $\iota_{1}(x)=x \otimes 1$ and $\iota_{2}(y)=1 \otimes y$ are algebra homomorphisms.

Let A- $\mathscr{A l} \mathscr{g}$ denote the category of associative unital commutative, A-algebras.
Proposition Let E and F be objects of $\mathrm{A}-\mathscr{A l}$ g. The algebra $\mathrm{E} \otimes \mathrm{F}$, equipped with $\iota_{1}$ and $\iota_{2}$, is the sum of E and F in the category $\mathrm{A}-\mathscr{A} \mathscr{\mathscr { q }}$.

Proof We show that, for any A-algebra C and all morphisms $f: \mathrm{E} \rightarrow \mathrm{C}$, $g: \mathrm{F} \rightarrow \mathrm{C}$ there is a unique morphism $h: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{C}$ such that, for all $x \in \mathrm{E}$ and $y \in \mathrm{~F}, h(x \otimes 1)=f(x)$ and $h(1 \otimes y)=g(y)$, or such that

$$
\begin{equation*}
h(x \otimes y)=f(x) g(y) \tag{*}
\end{equation*}
$$

since $x \otimes y=(x \otimes 1)(1 \otimes y)$. We know that there is a unique A-linear map $h: \mathrm{E} \otimes$ $\mathrm{F} \rightarrow \mathrm{C}$ satisfying $(*)$; moreover, for $x, x^{\prime} \in \mathrm{E}$ and $y, y^{\prime} \in \mathrm{F} h\left((x \otimes y)\left(x^{\prime} \otimes y^{\prime}\right)\right)=$ $h\left(x x^{\prime} \otimes y y^{\prime}\right)=f\left(x x^{\prime}\right) g\left(y y^{\prime}\right)=f(x) g(y) f\left(x^{\prime}\right) g\left(y^{\prime}\right)=h(x \otimes y) h\left(x^{\prime} \otimes y^{\prime}\right)$ since C is commutative.

Corollary Under the same assumptions, $\mathrm{E} \otimes_{\mathrm{A}} \mathrm{F}$ is an amalgamated sum of E and F over A in the category of commutative rings.

### 3.8.16

Let $\mathrm{E}, \mathrm{B}, \mathrm{C}$ be algebras, $f: \mathrm{B} \rightarrow \mathrm{E}$ and $g: \mathrm{C} \rightarrow \mathrm{E}$ algebra homomorphisms, then the image of the homomorphism $x \otimes y \mapsto f(x) g(y)$ of $\mathrm{B} \otimes \mathrm{C} \rightarrow \mathrm{E}$ is the subalgebra of E generated by $f(\mathrm{~B})$ and $g(\mathrm{C})$. In particular, if B and C are unital subalgebras of an algebra $E$, the subalgebra $B \cdot C$ of $E$ generated by $B$ and $C$ is identified with a quotient algebra of $\mathrm{B} \otimes \mathrm{C}$.

### 3.8.17

Let A and B be two rings, $f: \mathrm{A} \rightarrow \mathrm{B}$ a ring homomorphism and E an A -algebra. The $B$-module $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{E}$, equipped with the multiplication defined in 3.8.15, is a $B$-algebra, said to be the $B$-algebra obtained from $E$ by extension of scalars from $A$ to B along $f$. For any B -algebra $\mathrm{F}, \mathrm{F}$ can be regarded as an A -algebra by restriction of scalars, and

$$
\operatorname{Hom}_{\mathrm{B}-\mathscr{A l} g}\left(\mathrm{~B} \otimes_{\mathrm{A}} \mathrm{E}, \mathrm{~F}\right)=\operatorname{Hom}_{\mathrm{A}-\mathscr{A} \mathscr{g}}(\mathrm{E}, \mathrm{~F})
$$

### 3.8.18 Examples

(1) $\mathrm{A}[\mathrm{X}, \mathrm{Y}]=\mathrm{A}[\mathrm{X}] \otimes_{\mathrm{A}} \mathrm{A}[\mathrm{Y}]$.

Similarly $\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{p}, \mathrm{Y}_{1}, \ldots, \mathrm{Y}_{q}\right]=\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{p}\right] \otimes_{\mathrm{A}} \mathrm{A}\left[\mathrm{Y}_{1}, \ldots, \mathrm{Y}_{q}\right]$. More generally, if I is an arbitrary set and $\left(\mathrm{I}^{\prime}, \mathrm{I}^{\prime \prime}\right)$ a partition of I , then $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]=$ $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}^{\mathrm{I}}}\right] \otimes_{\mathrm{A}} \mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}^{\prime}}\right]$.
(2) Let $\mathrm{P}_{1}, \ldots, \mathrm{P}_{k} \in \mathrm{~A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{p}\right], \mathrm{Q}_{1}, \ldots, \mathrm{Q}_{l} \in \mathrm{~A}\left[\mathrm{Y}_{1}, \ldots, \mathrm{Y}_{q}\right]$. Then

$$
\begin{array}{r}
\left(\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{p}\right] /\left(\mathrm{P}_{1}, \ldots, \mathrm{P}_{k}\right)\right) \otimes_{\mathrm{A}}\left(\mathrm{~A}\left[\mathrm{Y}_{1}, \ldots, \mathrm{Y}_{q}\right] /\left(\mathrm{Q}_{1}, \ldots, \mathrm{Q}_{l}\right)\right) \\
\quad=\mathrm{A}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{p}, \mathrm{Y}_{1}, \ldots, \mathrm{Y}_{q}\right] /\left(\mathrm{P}_{1}, \ldots, \mathrm{P}_{k}, \mathrm{Q}_{1}, \ldots, \mathrm{Q}_{l}\right)
\end{array}
$$

(3) By extension of scalars from A to B along $f$, the algebra $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ gives $\mathrm{B}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ and $\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right] /\left(\left(\mathrm{P}_{j}\right)_{j \in \mathrm{~J}}\right)$ gives $\mathrm{B}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}} \mathrm{J} /\left(\left(f_{*} \mathrm{P}_{j}\right)_{j \in \mathrm{~J}}\right)\right.$.
(4) The algebra $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}$ is isomorphic to the product algebra $\mathbb{C} \times \mathbb{C}$. Indeed $\mathbb{C}=\mathbb{R}[\mathrm{X}] /\left(\mathrm{X}^{2}+1\right)$, and so $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}=\mathbb{C}[\mathrm{X}] /\left(\mathrm{X}^{2}+1\right)=\mathbb{C}[\mathrm{X}] /((\mathrm{X}-i)(\mathrm{X}+i))$. The substitutions $\mathrm{X} \mapsto i$ and $\mathrm{X} \mapsto-i$ give two homomorphisms from $\mathbb{C}[\mathrm{X}]$ to $\mathbb{C}$, and hence there is a homomorphism from $\mathbb{C}[X]$ to $\mathbb{C} \times \mathbb{C}$ whose kernel is the ideal $(\mathrm{X}-i) \cap(\mathrm{X}+i)=((\mathrm{X}-i)(\mathrm{X}+i))$. This in turn gives an injective homomorphism from $\mathbb{C}[\mathrm{X}] /((\mathrm{X}-i)(\mathrm{X}+i))$ to $\mathbb{C} \times \mathbb{C}$. This homomorphism is bijective since both these algebras have dimension 4 over $\mathbb{R}$. (We will resume this argument in the more general framework of Galois theory.)
(5) If $X$ is a topological space, $\mathbb{C} \otimes_{\mathbb{R}} C(X, \mathbb{R})=C(X, \mathbb{C})$.

## Exercises 3.8. (Tensor product)

1.-Let E and F be two A-modules, $\mathrm{L}_{1} \xrightarrow{\delta} \mathrm{~L}_{0} \xrightarrow{\varepsilon} \mathrm{E} \rightarrow 0$ and $\mathrm{M}_{1} \xrightarrow{\delta^{\prime}} \mathrm{M}_{0} \xrightarrow{\varepsilon^{\prime}} \mathrm{F} \rightarrow 0$ presentations of E and F respectively. Show that

$$
\left(\mathrm{L}_{1} \otimes \mathrm{M}_{0}\right) \oplus\left(\mathrm{L}_{0} \otimes \mathrm{M}_{1}\right) \xrightarrow{\left(\delta \otimes 1,1 \otimes \delta^{\prime}\right)} \mathrm{L}_{0} \otimes \mathrm{M}_{0} \xrightarrow{\varepsilon \otimes \varepsilon^{\prime}} \mathrm{E} \otimes \mathrm{~F} \rightarrow 0
$$

is a presentation of $\mathrm{E} \otimes \mathrm{F}$.
2.-Let K be a field and $\mathrm{E}_{1} \xrightarrow{u} \mathrm{E} \xrightarrow{f} \mathrm{~F} \xrightarrow{v} \mathrm{~F}_{1}$ and $\mathrm{E}_{1}^{\prime} \xrightarrow{u^{\prime}} \mathrm{E}^{\prime} \xrightarrow{f^{\prime}} \mathrm{F}^{\prime} \xrightarrow{v^{\prime}} \mathrm{F}_{1}^{\prime}$ two exact sequences of vector spaces. Show that $\left(\mathrm{E}_{1} \otimes \mathrm{E}^{\prime}\right) \oplus\left(\mathrm{E} \otimes \mathrm{E}_{1}^{\prime}\right) \xrightarrow{\left(u \otimes 1,1 \otimes u^{\prime}\right)}$ $\mathrm{E} \otimes \mathrm{E}^{\prime} \xrightarrow{f \otimes f^{\prime}} \mathrm{F} \otimes \mathrm{F}^{\prime} \xrightarrow{\left(\begin{array}{ll}v \otimes 1 \\ 1 \otimes & v^{\prime}\end{array}\right)}\left(\mathrm{F}_{1} \otimes \mathrm{~F}^{\prime}\right) \oplus\left(\mathrm{F} \otimes \mathrm{F}_{1}^{\prime}\right)$ is an exact sequence. In particular $\operatorname{Ker}\left(f \otimes f^{\prime}\right)=\left(\operatorname{Ker} f \otimes \mathrm{E}^{\prime}\right)+\left(\mathrm{E} \otimes \operatorname{Ker} f^{\prime}\right)$.
3.-Let E and F be K -vector spaces and $t \in \mathrm{E} \otimes \mathrm{F}$. Set $\mathrm{G}=\mathrm{E} \otimes \mathrm{F} / \mathrm{K} \cdot t$ and let $f: \mathrm{E} \otimes \mathrm{F} \rightarrow \mathrm{G}$ be the canonical map. Show that $(f(x \otimes y)=0) \Rightarrow(x \otimes y=0)$ if and only if $t$ is of rank $\geqslant 2$ and that the restriction of $f$ to the set of simple tensors is injective if and only if $t$ is of rank $\geqslant 3$.
4.-Let E and F be modules over an integral domain $\mathrm{A}, x \in \mathrm{E}$ and $y \in \mathrm{~F}$ non-torsion elements. Show that $x \otimes y \neq 0$ (extend the scalars to the field of fractions of A).
5.-Let M and N be square matrices with entries in a field K .
(a) Show that if M and N are diagonal (resp. triangular, resp. symmetric), then the same holds for $\mathrm{M} \otimes \mathrm{N}$.
(b) Show that if $M_{1}$ and $N_{1}$ are respectively similar to $M$ and $N, M_{1} \otimes N_{1}$ is similar to $\mathrm{M} \otimes \mathrm{N}$.
(c) Write $\operatorname{det}(M \otimes N)$ in terms of $\operatorname{det}(M)$ and $\operatorname{det}(N)$.
(d) Determine the eigenvalues of $\mathrm{M} \otimes \mathrm{N}$ with their multiplicities assuming known those of M and N .
6.-Let K be a field.
(a) Let X be a set and E a K -vector space. Show that the map $f \otimes t \mapsto$ $(x \mapsto f(x) \cdot t)$ from $\mathrm{K}^{\mathrm{X}} \otimes \mathrm{E}$ to $\mathrm{E}^{\mathrm{X}}$ is bijective if either X is finite or E is finite
dimensional, and is always injective (consider E as the increasing union of finite dimensional vector spaces). What is its image?
(b) Let X and Y be two sets. Show that the map $f \otimes g \mapsto h$ from $\mathrm{K}^{\mathrm{X}} \otimes \mathrm{K}^{\mathrm{Y}}$ to $\mathrm{K}^{\mathrm{X} \times \mathrm{Y}}$ defined by $h(x, y)=f(x) \cdot g(y)$ is injective.
(c) Let X and Y be topological spaces. Show that the map $\mathrm{C}(\mathrm{X}, \mathbb{R}) \otimes_{\mathbb{R}} \mathrm{C}(\mathrm{Y}, \mathbb{R})$ assigning to every vector $f \otimes g h$ defined by $h(x, y)=f(x) \cdot g(y)$ is injective. Can $\mathbb{R}$ be replaced by $\mathbb{C}$ ?
(d) Supposing X and Y to be compact and metrizable, show that this map is image-dense.
7.-Let $h: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism, M and N two B -modules. Show that there is an exact sequence:

$$
\mathrm{M} \otimes_{\mathrm{A}} \mathrm{~B} \otimes_{\mathrm{A}} \mathrm{~N} \xrightarrow{d} \mathrm{M} \otimes_{\mathrm{A}} \mathrm{~N} \rightarrow \mathrm{M} \otimes_{\mathrm{B}} \mathrm{~N} \rightarrow 0
$$

where $d(x \otimes b \otimes y)=b x \otimes y-x \otimes b y$.
8.-Let $K$ be a field. For every algebraic subset $X$ of $K^{n}$, let $A(X)$ be the algebra of functions on X induced by polynomials.
(a) Let $\mathrm{X} \subset \mathrm{K}^{m}$ and $\mathrm{Y} \subset \mathrm{K}^{n}$ be algebraic sets. Show that $\mathrm{A}(\mathrm{X} \times \mathrm{Y})=\mathrm{A}(\mathrm{X}) \otimes_{\mathrm{K}}$ A(Y).
(b) Let $p: \mathrm{K}^{n+r}=\mathrm{K}^{n} \times \mathrm{K}^{r} \rightarrow \mathrm{~K}^{n}$ and $q: \mathrm{K}^{n+s} \rightarrow \mathrm{~K}^{n}$ be the first projections, $\mathrm{S} \subset \mathrm{K}^{n}, \mathrm{X} \subset \mathrm{K}^{n+r}$ and $\mathrm{Y} \subset \mathrm{K}^{n+s}$ algebraic sets such that $p(\mathrm{X}) \subset \mathrm{S}, q(\mathrm{Y}) \subset \mathrm{S}$. Define $p^{*}: \mathrm{A}(\mathrm{S}) \rightarrow \mathrm{A}(\mathrm{X})$ by $p^{*}(f)=f \circ p$, and $q^{*}: \mathrm{A}(\mathrm{S}) \rightarrow \mathrm{A}(\mathrm{Y})$ likewise. Show that $\mathrm{A}\left(\mathrm{X} \times_{\mathrm{S}} \mathrm{Y}\right)=\mathrm{A}(\mathrm{X}) \otimes_{\mathrm{A}(\mathrm{S})} \mathrm{A}(\mathrm{Y})$.
9.-(a) Let X and Y be two compact metrizable spaces, $f: \mathrm{Y} \rightarrow \mathrm{X}$ a continuous map, and E a vector bundle over X (3.4, Exercise 9). Define a bundle $\mathrm{F}=f^{*}(\mathrm{E})$ over Y by equipping the fibre product en $\mathrm{Y} \times_{\mathrm{X}} \mathrm{E}$ with its first projection onto Y and, for all $y \in \mathrm{Y}$, with the structure obtained on $\mathrm{F}(y)$ by identifying it with $\mathrm{E}(f(y))$. Show that the module $\mathrm{S}(\mathrm{F})$ of continuous sections of F can be identified with a module obtained from $\mathrm{S}(\mathrm{E})$ by extension of scalars along $f^{*}: \mathrm{C}(\mathrm{X}) \rightarrow \mathrm{C}(\mathrm{Y})$.
(b) Let E and F be two vector bundles over X . Define a vector bundle $\mathrm{E} \otimes \mathrm{F}$ over X such that $(\mathrm{E} \otimes \mathrm{F})(x)=\mathrm{E}(x) \otimes \mathrm{F}(x)$ for all $x \in \mathrm{X}$, and show that $\mathrm{S}(\mathrm{E} \otimes \mathrm{F})=$ $S(E) \otimes_{C(X)} S(F)$.
10.-Let $K$ be a field, A the ring $K[X, Y]$ and $\mathfrak{m}$ the maximal ideal in A consisting of polynomials without constant terms. The aim is to determine the module $\mathfrak{m} \otimes_{\mathrm{A}} \mathfrak{m}$.
(a) Show that there is an exact sequence $0 \longrightarrow A \xrightarrow{\binom{-Y}{X}} A^{2} \xrightarrow{(X, Y)} \mathfrak{m} \longrightarrow 0$.
(b) Show that this sequence induces an exact sequence

$$
0 \rightarrow \mathrm{~A} \rightarrow \mathfrak{m} \oplus \mathfrak{m} \rightarrow \mathfrak{m}^{2} \rightarrow 0
$$

where $\mathfrak{m}^{2}$ denotes the ideal of A generated by the monomials of degree 2 .
(c) Deduce an exact sequence $0 \rightarrow \mathrm{~K} \rightarrow \mathfrak{m} \otimes \mathfrak{m} \rightarrow \mathfrak{m}^{2} \rightarrow 0$, where K is endowed with the A-module structure given by its identification with $\mathrm{A} / \mathfrak{m}$.
(d) Show that the element $x \otimes y-y \otimes x$ of $\mathfrak{m} \otimes \mathfrak{m}$ is nonzero and generates the torsion submodule.
(e) Show that the submodule of $\mathfrak{m} \otimes \mathfrak{m}$ generated by $x \otimes x, x \otimes y$ and $y \otimes y$ is isomorphic to $\mathfrak{m}^{2}$, and finally that $\mathfrak{m} \otimes \mathfrak{m}$ is isomorphic to $\mathfrak{m}^{2} \oplus K$.
11.-In what follows, A and B denote two unitary commutative rings, and $\alpha: \mathrm{A} \rightarrow \mathrm{B}$ a unitary ring homomorphism. The A -module $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{B}$ being equipped with its canonical ring structure (multiplication being given by: $\left.(x \otimes y)\left(x^{\prime} \otimes y^{\prime}\right)=x x^{\prime} \otimes y y^{\prime}\right)$, let $p: \mathrm{B} \otimes_{\mathrm{A}} \mathrm{B} \rightarrow \mathrm{B}$ be the homomorphism given by: $p(x \otimes y)=x y$, and $j_{1}, j_{2}: \mathrm{B} \rightarrow$ $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{B}$ the homomorphisms defined by: $j_{1}(x)=x \otimes 1, j_{2}(x)=1 \otimes x$.
(a) Show that the kernel I of $p$ is generated by the elements $1 \otimes x-x \otimes 1(x \in \mathrm{~B})$.
(b) For any B-module M, let $\mathscr{D}_{\mathrm{A}}(\mathrm{B}, \mathrm{M})$ be the set of "A-derivation from B to M", defined by:

$$
\mathscr{D}_{\mathrm{A}}(\mathrm{~B}, \mathrm{M})=\left\{\mathrm{D} \in \operatorname{Hom}_{\mathrm{A}}-\mathscr{M}^{\boldsymbol{o d}}(\mathrm{B}, \mathrm{M}) \mid \forall x, y \in \mathrm{~B}, \mathrm{D}(x y)=x \mathrm{D}(y)+y \mathrm{D}(x)\right\} .
$$

$(\alpha)$ Show that $\mathscr{D}_{\mathrm{A}}(\mathrm{B}, \mathrm{M})$ is a B-submodule of $\operatorname{Hom}_{\mathrm{A}} \mathscr{M} \boldsymbol{\operatorname { O d }}(\mathrm{B}, \mathrm{M})$.
( $\beta$ ) Show that $\mathrm{D} \circ \alpha=0$, for any A-derivation D from B to M .
(c) The ring $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{B}$ being equipped with its B -algebra structure defined by $j_{1}$ (i.e. B acts on $\mathrm{B} \otimes_{\mathrm{A}} \mathrm{B}$ by: $\left.b \cdot(x \otimes y)=j_{1}(b)(x \otimes y)\right)$, consider the isomorphism of B-modules:

$$
v: \operatorname{Hom}_{\mathrm{A}-\mathscr{M} \boldsymbol{\sim}}(\mathrm{B}, \mathrm{M}) \rightarrow \operatorname{Hom}_{\mathrm{B}-\mathscr{M} \boldsymbol{\mathscr { d }}}\left(\mathrm{B} \otimes_{\mathrm{A}} \mathrm{~B}, \mathrm{M}\right)
$$

given by: $v(f)(x \otimes y)=x f(y)$. Show that $v\left(\mathscr{D}_{\mathrm{A}}(\mathrm{B}, \mathrm{M})\right)$ is the B -submodule of $\operatorname{Hom}_{\mathrm{B}}-\mathscr{M} \underset{\operatorname{d}}{ }\left(\mathrm{B} \otimes_{\mathrm{A}} \mathrm{B}, \mathrm{M}\right)$ whose elements vanish on $j_{1}(\mathrm{~B})$ and on $\mathrm{I}^{2}$.
(d) Deduce that there exists a B -module $\Omega_{\mathrm{B} / \mathrm{A}}$ (called the module of A-differentials of B ) and an A-derivation $d_{\mathrm{B} / \mathrm{A}}: \mathrm{B} \rightarrow \Omega_{\mathrm{B} / \mathrm{A}}$ such that, for all A -modules M and A-derivations $\mathrm{D}: \mathrm{B} \rightarrow \mathrm{M}$, there is a unique homomorphism of B -modules $u$ : $\Omega_{\mathrm{B} / \mathrm{A}} \rightarrow \mathrm{M}$ such that: $\mathrm{D}=u \circ d_{\mathrm{B} / \mathrm{A}}$.
(e) Describe $\Omega_{\mathrm{B} / \mathrm{A}}$ in the following cases:
( $\alpha$ ) $\mathrm{A}=\mathbb{R}, \mathrm{B}=\mathbb{C}, \alpha$ is the canonical injection;
$(\beta) \mathrm{A}=\mathbb{Z}, \mathrm{B}=\mathbb{Q}, \alpha$ is the canonical injection;
$(\gamma) \mathrm{A}=\mathrm{F}_{p}, \mathrm{~B}=\mathrm{F}_{q}, \alpha$ is the canonical injection $\left(\mathrm{F}_{p}\right.$ and $\mathrm{F}_{q}$ are respectively the fields of order $p$ and $q=p^{f}, p$ being prime and $\left.f \in \mathbb{N}^{*}\right)$;
( $\delta$ ) $\mathrm{A}=k\left[\mathrm{X}^{p}\right], \mathrm{B}=k[\mathrm{X}]$, where $p$ is prime and $k$ a field, $\alpha$ being the canonical injection;
( $\varepsilon$ ) $\mathrm{B}=\mathrm{A}^{n}, n \geqslant 1, \alpha$ being defined by: $\alpha(a)=(a, a, \ldots, a)$;
( $\zeta$ ) $\alpha: \mathrm{A} \rightarrow \mathrm{B}$ is surjective;
$(\eta) \mathrm{B}$ is an A-algebra of polynomials with indeterminates in an arbitrary set L , the set $\alpha: \mathrm{A} \rightarrow \mathrm{A}\left[\mathrm{X}_{e}\right]_{e \in \mathrm{~L}}$ being defined by $\alpha(a)=a$.
12.-Let E and F be real Euclidean spaces.
(a) Show that there is a unique symmetric bilinear form on $\mathrm{E} \otimes \mathrm{F}$ satisfying $\left\langle x \otimes y, x^{\prime} \otimes y^{\prime}\right\rangle=\left\langle x, x^{\prime}\right\rangle \cdot\left\langle y, y^{\prime}\right\rangle$.
(b) Let $\left(e_{1}, \ldots, e_{n}\right),\left(f_{1}, \ldots, f_{p}\right)$ be orthonormal bases for E and F respectively. Show that the elements $e_{i} \otimes f_{j}$ form an orthonormal basis for $\mathrm{E} \otimes \mathrm{F}$. Deduce that $\mathrm{E} \otimes \mathrm{F}$ is Euclidian.
(c) Generalize these results to Hermitian spaces over $\mathbb{C}$.
13.-Let $K$ be a field. For every vector space $E$, denote by $P(E)$ the projective space of $E$, i.e. the set of vector subspaces of rank 1 of $E$ and $\psi_{E}$ the map from $E-\{0\}$ to $\mathrm{P}(\mathrm{E})$ assigning to every vector the line it generates.
(a) Let E and F be two vector spaces. Show that the map $(\mathrm{L}, \mathrm{M}) \mapsto \mathrm{L} \otimes \mathrm{M}$ from $P(E) \times P(F)$ to $P(E \otimes F)$ is injective.
(b) Let $x_{0} \in \mathrm{E}$ and $y_{0} \in \mathrm{~F}$. Determine the tangent linear map $\mathrm{T}_{x_{0}, y_{0}} \theta$ to $\theta$ : $(x, y) \mapsto x \otimes y$ from $\mathrm{E} \times \mathrm{F}$ to $\mathrm{E} \otimes \mathrm{F}$ at $\left(x_{0}, y_{0}\right)$ (i.e. the degree 1 part of the expansion of $\left.(u, v) \mapsto\left(x_{0}+u\right) \otimes\left(y_{0}+v\right)\right)$.
(c) Suppose $x_{0} \neq 0$ and $y_{0} \neq 0$. Show that, by passing to quotients, $\mathrm{T}_{x_{0}, y_{0}} \theta$ gives an injective map from $\mathrm{E} / \mathrm{K} x_{0} \times \mathrm{F} / \mathrm{K} y_{0}$ to $\mathrm{E} \otimes \mathrm{F} / \mathrm{K}\left(x_{0} \otimes y_{0}\right)$.
(d) Suppose that $K=\mathbb{R}$ or $\mathbb{C}$ and that E and F are finite dimensional. Show that $\mathrm{P}(\mathrm{E}) \times \mathrm{P}(\mathrm{F})$ can be identified with a closed submanifold of $\mathrm{P}(\mathrm{E} \otimes \mathrm{F})$.

### 3.9 The Kronecker Homomorphism

### 3.9.1

Let $\mathrm{E}, \mathrm{F}, \mathrm{E}^{\prime}, \quad \mathrm{F}^{\prime}$ be A-modules. The map $\left(f, f^{\prime}\right) \mapsto f \otimes f^{\prime}$ from $\operatorname{Hom}(\mathrm{E}, \mathrm{F}) \times \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)$ to $\operatorname{Hom}\left(\mathrm{E} \otimes \mathrm{E}^{\prime}, \mathrm{F} \otimes \mathrm{F}^{\prime}\right)$ is bilinear. Hence there is a homomorphism

$$
\varphi_{\mathrm{F}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}: f \otimes f^{\prime} \mapsto f \otimes f^{\prime}
$$

from $\operatorname{Hom}(\mathrm{E}, \mathrm{F}) \otimes \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)$ to $\operatorname{Hom}\left(\mathrm{E} \otimes \mathrm{E}^{\prime}, \mathrm{F} \otimes \mathrm{F}^{\prime}\right)$, called the Kronecker homomorphism.

In particular, this gives morphisms $\alpha_{\mathrm{F}}^{\mathrm{E}}=\varphi_{\mathrm{A}, \mathrm{F}}^{\mathrm{E}, \mathrm{A}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ and $\beta^{\mathrm{E}, \mathrm{F}}=\varphi_{\mathrm{A}, \mathrm{A}}^{\mathrm{E}, \mathrm{F}}: \mathrm{E}^{\top} \otimes \mathrm{F}^{\top} \rightarrow(\mathrm{E} \otimes \mathrm{F})^{\top}$, where $\mathrm{E}^{\top}=\operatorname{Hom}(\mathrm{E}, \mathrm{A})$.

### 3.9.2 Functoriality

Let $u: \mathrm{E}_{1} \rightarrow \mathrm{E}, u^{\prime}: \mathrm{E}_{1}^{\prime} \rightarrow \mathrm{E}^{\prime}, v: \mathrm{F} \rightarrow \mathrm{F}_{1}$, and $v^{\prime}: \mathrm{F}^{\prime} \rightarrow \mathrm{F}_{1}^{\prime}$ be homomorphisms of A-modules. These give a homomorphism $v_{*} u^{*}=u^{*} v_{*}: f \mapsto v \circ f \circ u$ from $\operatorname{Hom}(\mathrm{E}, \mathrm{F})$ to $\operatorname{Hom}\left(\mathrm{E}_{1}, \mathrm{~F}_{1}\right)$. The commutativity of the following diagram follows readily:


This says that $\varphi$ is a morphism from the functor $\left(E, E^{\prime}, F, F^{\prime}\right) \mapsto \operatorname{Hom}(E, F) \otimes$ $\operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)$ to the functor $\left(\mathrm{E}, \mathrm{E}^{\prime}, \mathrm{F}, \mathrm{F}^{\prime}\right) \mapsto \operatorname{Hom}\left(\mathrm{E} \otimes \mathrm{E}^{\prime}, \mathrm{F} \otimes \mathrm{F}^{\prime}\right)$, both these functors being defined on $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}^{0} \times \mathrm{A}-\mathscr{M} \sigma d^{0} \times \mathrm{A}-\mathscr{M} \sigma \boldsymbol{d} \times \mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$, with values in A- $\mathscr{M} \sigma d$.

### 3.9.3

For homomorphisms $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime}$ and $g: \mathrm{N} \rightarrow \mathrm{N}^{\prime}$, let $f \oplus g$ be the homomorphism $(x, y) \mapsto(f(x), g(y))$ from $\mathrm{M} \oplus \mathrm{N}$ to $\mathrm{M}^{\prime} \oplus \mathrm{N}^{\prime}$.

Proposition If $\mathrm{E}=\mathrm{E}_{1} \oplus \mathrm{E}_{2}$, then $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}=\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}_{1}, \mathrm{E}^{\prime}} \oplus \varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}_{2}, \mathrm{E}^{\prime}}$. Similarly, if $\mathrm{F}=\mathrm{F}_{1} \oplus \mathrm{~F}_{2}$, then $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}=\varphi_{\mathrm{F}_{1}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}} \oplus \varphi_{\mathrm{F}_{2}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$

Proof Apply functoriality to $\iota_{1}: \mathrm{E}_{1} \rightarrow \mathrm{E}$ and $\iota_{2}: \mathrm{E}_{2} \rightarrow \mathrm{E}$ (resp. to $\pi_{1}: \mathrm{E} \rightarrow \mathrm{E}_{1}$ and $\pi_{2}: \mathrm{E} \rightarrow \mathrm{E}_{2}$ ).
Corollary If $\mathrm{E}=\mathrm{E}_{1} \oplus \mathrm{E}_{2}$, then $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$ is an isomorphism if and only if so are $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}_{1}, \mathrm{E}^{\prime}}$ and $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}_{2}, \mathrm{E}^{\prime}}$. Similarly, if $\mathrm{F}=\mathrm{F}_{1} \oplus \mathrm{~F}_{2}$, then $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$ is an isomorphism if and only if so $\operatorname{are} \varphi_{\mathrm{F}_{1}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$ and $\varphi_{\mathrm{F}_{2}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$.

This follows from the fact that $f \oplus g$ is an isomorphism if and only if so are $f$ and $g$.

### 3.9.4 Specific Cases When $\varphi$ is an Isomorphism

Proposition If E and $\mathrm{E}^{\prime}$ (resp. E and F ) are finitely generated and projective (i.e. isomorphic to direct factors of a finitely generated free module), $\varphi_{\mathrm{F}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}$ is an isomorphism.

Proof It is immediate if $\mathrm{E}=\mathrm{E}^{\prime}=\mathrm{A}$, or if $\mathrm{E}=\mathrm{F}=\mathrm{A}$. The case where E and $\mathrm{E}^{\prime}$ (resp. F and $\mathrm{F}^{\prime}$ ) are finitely generated and free follows by iteration from (3.9.3, Corollary). The case where ù E and $\mathrm{E}^{\prime}$ (resp. E and F) are finitely generated and projective then follows by applying the converse part of the same corollary.

Corollary If either E or F is finitely generated and projective, then $\alpha_{\mathrm{F}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is an isomorphism.

### 3.9.5 Specific Cases When $\alpha$ is Bijective

Proposition If E is free and F has finite presentation, then $\alpha_{\mathrm{F}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is an isomorphism.

Proof Let $\mathrm{L}_{1} \rightarrow \mathrm{~L}_{0} \rightarrow \mathrm{~F} \rightarrow 0$ be a presentation of F , where $\mathrm{L}_{0}$ and $\mathrm{L}_{1}$ are are finitely generated and free. By Corollary 3.9.4, the morphisms $\alpha_{\mathrm{L}_{0}}^{\mathrm{E}}$ and $\alpha_{\mathrm{L}_{1}}^{\mathrm{E}}$ are bijective. As the diagram

where the rows are exact, is commutative, $\mathrm{E}^{\top} \otimes \mathrm{F}$ and $\operatorname{Hom}(\mathrm{E}, \mathrm{F})$ can be identified with the cokernel of the same homomorphism.

### 3.9.6 The Noetherian Case

Proposition If A is Noetherian and E is free, then $\alpha_{\mathrm{F}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is injective. Its image is the set of $f: \mathrm{E} \rightarrow \mathrm{F}$ for which $f(\mathrm{E})$ is contained in a finitely generated submodule of F .
Proof Let $\left(\mathrm{F}_{i}\right)$ be an increasing directed family of finitely generated submodules of F such that $\mathrm{F}=\bigcup \mathrm{F}_{i}$. For all $i$, the module $\mathrm{F}_{i}$ has a finite presentation since A is Noetherian, and so $\alpha_{\mathrm{F}_{i}}^{\mathrm{E}}$ is bijective. As $\mathrm{F}=\xrightarrow{\lim } \mathrm{F}_{i}$, passing to the direct limit, gives an isomorphism from $\mathrm{E}^{\top} \otimes \mathrm{F}=\underset{\rightarrow}{\lim } \mathrm{E}^{\top} \otimes \mathrm{F}_{i}$ onto $\xrightarrow{\lim } \operatorname{Hom}\left(\mathrm{E}, \mathrm{F}_{i}\right)$. This direct limit is just the union of submodules $\overrightarrow{\operatorname{Hom}}\left(\mathrm{E}, \mathrm{F}_{i}\right)$ of $\operatorname{Hom}(\mathrm{E}, \mathrm{F})$.
Corollary 3.9 If A is Noetherian and E est free, then $\beta^{\mathrm{E}, \mathrm{F}}: \mathrm{E}^{\top} \otimes \mathrm{F}^{\top} \rightarrow(\mathrm{E} \otimes \mathrm{F})^{\top}$ is injective.

Indeed, $\beta^{\mathrm{E}, \mathrm{F}}=\alpha_{\mathrm{F}^{\top}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F}^{\top} \rightarrow \operatorname{Hom}\left(\mathrm{E}, \mathrm{F}^{\top}\right)=\mathrm{B}(\mathrm{E}, \mathrm{F} ; \mathrm{A})=(\mathrm{E} \otimes \mathrm{F})^{\top}$.
Corollary 3.10 If A is Noetherian, E F free, then

$$
\varphi_{\mathrm{F}, \mathrm{~F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}: \operatorname{Hom}(\mathrm{E}, \mathrm{~F}) \otimes \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right) \rightarrow \operatorname{Hom}\left(\mathrm{E} \otimes \mathrm{E}^{\prime}, \mathrm{F} \otimes \mathrm{~F}^{\prime}\right)
$$

is injective.
Proof If $\mathrm{F}=\mathrm{A}$, the homomorphism $\varphi_{\mathrm{A}, \mathrm{F}^{\prime}}^{\mathrm{E}, \mathrm{E}^{\prime}}=\alpha_{\mathrm{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)}^{\mathrm{E}}: \quad \operatorname{Hom}(\mathrm{E}, \mathrm{A}) \otimes \operatorname{Hom}$ $\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right) \rightarrow \operatorname{Hom}\left(\mathrm{E}, \operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)\right)=\mathrm{B}\left(\mathrm{E}, \mathrm{E}^{\prime} ; \mathrm{F}^{\prime}\right)=\operatorname{Hom}\left(\mathrm{E} \otimes \mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)$ is injective by Corollary 1 . The result follows by passing to the direct sum when $F$ is free and finitely generated, then by passing to the direct limit when $F$ is free. Hence it is the increasing directed union of finitely generated free modules.

### 3.9.7 The Principal Ideal Domain Case

Proposition If A is principal, $\alpha_{\mathrm{F}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is injective for all E and F .
Proof First suppose that F is finitely generated. By 3.5.8, it reduces to the case when F is cyclic. If $\mathrm{F}=\mathrm{A}$, then the morphism $\alpha_{\mathrm{F}}^{\mathrm{E}}$ is the identity of $\mathrm{E}^{\top}$. If $\mathrm{F}=\mathrm{A} /(a)$ with $a \neq$ 0 , the exact sequence $0 \rightarrow \mathrm{E}^{\top} \xrightarrow{a \cdot 1} \mathrm{E}^{\top} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ and $\mathrm{E}^{\top} \xrightarrow{a \cdot 1} \mathrm{E}^{\top} \rightarrow \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow 0$ associated to the exact sequence $0 \rightarrow \mathrm{~A} \xrightarrow{a \cdot 1} \mathrm{~A} \rightarrow \mathrm{~F} \rightarrow 0$ show that $\alpha$ identifies $\mathrm{E}^{\top} \otimes \mathrm{F}$ with a submodule de $\operatorname{Hom}(E, F)$.

The general case follows as in Proposition 3.9.6.

### 3.9.8 Image of $\alpha$

Proposition The image of $\alpha_{\mathrm{F}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{F} \rightarrow \operatorname{Hom}(\mathrm{E}, \mathrm{F})$ is the set of homomorphisms $f: \mathrm{E} \rightarrow \mathrm{F}$ factorizing through a finitely generated free $\mathrm{A}-m o d u l e$.
Proof If $t=\sum_{1}^{r} \xi_{i} \otimes y_{i} \in \mathrm{E}^{\top} \otimes \mathrm{F}$, then $\alpha(t)=g \circ f$, where $f: \mathrm{E} \rightarrow \mathrm{A}^{r}$ and $g$ : $\mathrm{A}^{r} \rightarrow \mathrm{~F}$ are defined by $f(x)=\left(\xi_{1}(x), \ldots, \xi_{r}(x)\right)$ and $g\left(a_{1}, \ldots, a_{r}\right)=a_{1} y_{1}+\cdots+a_{r} y_{r}$. Conversely, if $h=g \circ f$ with $f: \mathrm{E} \rightarrow \mathrm{A}^{r}$ and $g: \mathrm{A}^{r} \rightarrow \mathrm{~F}$, then $h=\sum \xi_{i} \otimes y_{i}$, where the $\xi_{i}$ are the coordinates of $f$ and the $y_{i}$ are the images under $g$ of the elements of a canonical basis.
Corollary Let E be an A-module. The following conditions are equivalent:
(i) $1_{\mathrm{E}}$ belongs to the image of $\alpha_{\mathrm{E}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{E} \rightarrow \operatorname{End}(\mathrm{E})$;
(ii) $\alpha_{\mathrm{E}}^{\mathrm{E}}$ is bijective;
(iii) E is finitely generated and projective.

Proof (i) $\Rightarrow$ (iii) By the above proposition and Proposition 3.4.8, (iii) $\Rightarrow$ (ii) by (3.9.4, Corollary), and (ii) $\Rightarrow$ (i) is then immediate.

### 3.9.9 Contraction and Composition

Proposition Let E, F, G be three A-modules. The diagram

where $k$ is the contraction $\xi \otimes y \otimes \eta \otimes z \mapsto \eta(y) \cdot \xi \otimes z$ and $c$ the composition $(f \otimes g) \mapsto g \circ f$, is commutative.

Proof The image of $\xi \otimes y \otimes \eta \otimes z$ under the two composite maps in Hom(E, G) is the map $x \mapsto \xi(x) \cdot \eta(y) \cdot z$.

### 3.9.10 Trace

Let E be an A-module. Consider the Kronecker homomorphism $\alpha_{\mathrm{E}}=\alpha_{\mathrm{E}}^{\mathrm{E}}: \mathrm{E}^{\top} \otimes$ $\mathrm{E} \rightarrow \operatorname{End}(\mathrm{E})$ and the contraction $\kappa_{\mathrm{E}}: \mathrm{E}^{\top} \otimes \mathrm{E} \rightarrow \mathrm{A}$ defined by $\kappa_{\mathrm{E}}(\xi \otimes x)=\xi(x)$. Suppose that $\operatorname{Ker} \alpha_{\mathrm{E}} \subset \operatorname{Ker} \kappa_{\mathrm{E}}$. Then, for all $f: \mathrm{E} \rightarrow \mathrm{E}$ in the image of $\alpha_{\mathrm{E}}$, define $\operatorname{Tr}(f) \in \mathrm{A}$ by $\operatorname{Tr}(f)=\kappa(t)$, where $t \in \alpha_{\mathrm{E}}^{-1}(f)$, a scalar independent of the choice of $t$. The scalar $\operatorname{Tr}(f)$ is called the trace of $f$. The map $\operatorname{Tr}: \operatorname{Im} \alpha_{\mathrm{E}} \rightarrow \mathrm{A}$ is A-linear.

Proposition Suppose that E is finitely generated and free and let $\left(e_{i}\right)$ be a basis for E. For all endomorphisms $f$ of E , the trace of $f$ is the sum $\sum a_{i}^{i}$ of the diagonal entries of the matrix representing $f$.

Proof As $\alpha_{\mathrm{E}}$ is bijective, $\operatorname{Tr}(f)$ is defined for all $f \in \operatorname{End}(\mathrm{E})$. Let $\left(e_{i}^{\prime}\right)$ be the dual basis $\mathrm{E}^{\top}$ of $\left(e_{i}\right)$. For all $x \in \mathrm{E}, x=\sum e_{i}^{\prime}(x) e_{i}$; so $f(x)=\sum e_{i}^{\prime}(x) f\left(e_{i}\right)$. Thus $f=$ $\alpha\left(\sum e_{i}^{\prime} \otimes f\left(e_{i}\right)\right)$ and $\operatorname{Tr}(f)=\sum e_{i}^{\prime}\left(f\left(e_{i}\right)\right)=\sum a_{i}^{i}$.

### 3.9.11

Proposition Let E and F be two $\mathrm{A}-m o d u l e s, f: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism in the image of $\alpha_{\mathrm{F}}^{\mathrm{E}}$ and $g: \mathrm{F} \rightarrow \mathrm{E}$ an arbitrary homomorphism. Assume that $\operatorname{Ker} \alpha_{\mathrm{E}} \subset$ $\operatorname{Ker} \kappa_{\mathrm{E}}$ and $\operatorname{Ker} \alpha_{\mathrm{F}} \subset \operatorname{Ker} \kappa_{\mathrm{F}}$. Then $\operatorname{Tr}(g \circ f)$ and $\operatorname{Tr}(f \circ g)$ are well defined, and $\operatorname{Tr}(g \circ f)=\operatorname{Tr}(f \circ g)$.

Proof Let $t=\sum x_{i}^{\prime} \otimes y_{i} \in \mathrm{E}^{\top} \otimes \mathrm{F}$ be such that $f=\alpha_{\mathrm{F}}^{\mathrm{E}}(t)$. Then $g \circ f=\alpha_{\mathrm{E}}\left(\sum x_{i}^{\prime} \otimes\right.$ $\left.g\left(y_{i}\right)\right)$ and $f \circ g=\alpha_{\mathrm{F}}\left(\sum\left(x_{i}^{\prime} \circ g\right) \otimes y_{i}\right)$, so that $\operatorname{Tr}(g \circ f)$ and $\operatorname{Tr}(f \circ g)$ are defined and equal to $\sum x_{i}^{\prime}\left(g\left(y_{i}\right)\right)$.

### 3.9.12 Trace and Extension of Scalars

Let $h: \mathrm{A} \rightarrow \mathrm{A}^{\prime}$ be a ring homomorphism, E and F A-modules, $\mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime}$ the $\mathrm{A}^{\prime}$-modules resulting from E and F by extension of scalars. Denote by $h_{*}$ the $x \mapsto 1 \otimes x$ from E to $\mathrm{E}^{\prime}$ (resp. from F to $\mathrm{F}^{\prime}$ ), and by $h_{*}$ the map $\xi \mapsto h \circ \xi$ from $\mathrm{E}^{\top}$ to $\operatorname{Hom}_{\mathrm{A}}\left(\mathrm{E}, \mathrm{A}^{\prime}\right)=\operatorname{Hom}_{\mathrm{A}^{\prime}}\left(\mathrm{E}^{\prime}, \mathrm{A}^{\prime}\right)=\mathrm{E}^{\prime \top}$. For all A-linear maps $f: \mathrm{E} \rightarrow \mathrm{F}$, let $f_{*}$ be the $\mathrm{A}^{\prime}$-linear map $1_{\mathrm{A}^{\prime}} \otimes f: \mathrm{E}^{\prime} \rightarrow \mathrm{F}^{\prime}$.

Proposition The following diagrams are commutative:
(a)

(b)


Proof (a) Let $\xi \in \mathrm{E}^{\top}$ and $y \in \mathrm{~F}$. Then $\alpha_{\mathrm{F}}^{\mathrm{E}}(\xi \otimes y)=(x \mapsto \xi(x) \cdot y)$, which gives a map $a^{\prime} \otimes x \mapsto a^{\prime} \otimes \xi(x) \cdot y$ in $\operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}^{\prime}\right)$. In other words, given the definition of the A-module structure on $\mathrm{A}^{\prime}, \xi \otimes y$ gives the element $h_{*}(\xi) \otimes(1 \otimes y)$ in $\mathrm{E}^{\top \top} \otimes \mathrm{F}^{\prime}$, whose image under the Kronecker homomorphism is the map $a^{\prime} \otimes x \mapsto a^{\prime} \cdot h(\xi(x))$. $1 \otimes y=a^{\prime} \otimes \xi(x) \cdot y$.
(b) For $\xi \in \mathrm{E}^{\top}$ and $x \in \mathrm{E}, \kappa \circ\left(h_{*} \otimes h_{*}\right)(\xi \otimes x)=\kappa\left(h_{*}(\xi) \otimes(1 \otimes x)\right)=h(\xi(x))$.

Corollary Suppose that $\operatorname{Ker} \alpha_{\mathrm{E}} \subset \operatorname{Ker} \kappa_{\mathrm{E}}$ and $\operatorname{Ker} \alpha_{\mathrm{E}^{\prime}} \subset \operatorname{Ker} \kappa_{\mathrm{E}^{\prime}}$ hold. Then, for $f \in \operatorname{Im}\left(\alpha_{\mathrm{E}}\right), \operatorname{Tr}\left(f_{*}\right)=h(\operatorname{Tr}(f))$.

### 3.9.13

Proposition Suppose that the ring A is reduced. Then, $\operatorname{Ker} \alpha_{\mathrm{E}} \subset \operatorname{Ker} \kappa_{\mathrm{E}}$ for all A-modules E .

Proof Let $t \in \mathrm{E}^{\top} \otimes \mathrm{E}$ be such that $\alpha_{\mathrm{E}}(t)=0$. We show that $\kappa_{\mathrm{E}}(t)=0$. By 3.1.7, all we need to show is that, for all homomorphisms $h$ from A to a field $\mathrm{A}^{\prime}, h\left(\kappa_{\mathrm{E}}(t)\right)=0$. With the notation of 3.9.12 and assuming $\mathrm{A}^{\prime}$ is a field, the homomorphism $\kappa_{\mathrm{E}^{\prime}}$ : $\mathrm{E}^{\prime \top} \otimes_{\mathrm{A}^{\prime}} \mathrm{E}^{\prime} \rightarrow \operatorname{End}\left(\mathrm{E}^{\prime}\right)$ is injective by Proposition 3.9.6, and so the image of $t$ under $h_{*} \otimes h_{*}: \mathrm{E}^{\top} \otimes \mathrm{E} \rightarrow \mathrm{E}^{\top \top} \otimes \mathrm{E}^{\prime}$ is zero, and $h\left(\kappa_{\mathrm{E}}(t)\right)=\kappa_{\mathrm{E}^{\prime}} \circ\left(h_{*} \otimes h_{*}\right)(t)=0$.

### 3.9.14 Computation of the Rank of a Tensor

Proposition Let K be a field, E and F vector spaces over K and $t \in \mathrm{E}^{\top} \otimes \mathrm{F}$. The rank (see 3.8.2) of t equals the rank (i.e. the dimension of the image) of the linear $\operatorname{map} \alpha_{\mathrm{F}}^{\mathrm{E}}(t): \mathrm{E} \rightarrow \mathrm{F}$.

Proof Let $t \in \mathrm{E}^{\top} \otimes \mathrm{F}$ and set $f=\alpha_{\mathrm{F}}^{\mathrm{E}}(t)$. If $t=\sum_{i} x_{i} \otimes y_{i}$, then the image of $f$ is contained in the vector subspace of F generated by the $y_{i}$, and so $\mathrm{rk}(f) \leqslant r$. Hence rk $f \leqslant \mathrm{rk}(t)$.

Let $d$ be the dimension of the image V of $f,\left(y_{1}, \ldots, y_{d}\right)$ a basis for V and define the linear forms $u_{1}, \ldots, u_{d}$ on E by $f(x)=\sum u_{i}(x) \cdot y_{i}$. Then $f=\alpha\left(\sum u_{i} \otimes y_{i}\right)$, and as $\alpha_{\mathrm{F}}^{\mathrm{E}}$ is injective (Proposition 3.9.6), $t=\sum u_{i} \otimes y_{i}$. Thus $\operatorname{rk}(t) \leqslant d=\operatorname{rk}(f)$.

Corollary Let E and F be vector spaces over $\mathrm{K}, \mathrm{E}^{\prime}$ a vector space of linear forms on E separating the points of E (i.e. such that $(\forall x \in \mathrm{E}, x \neq 0)\left(\exists f \in \mathrm{E}^{\prime}\right)$ $f(x) \neq 0)$. Let $t=\sum x_{i} \otimes y_{i} \in \mathrm{E} \otimes \mathrm{F}$. Then the rank of t equals the rank of the map $f \mapsto \sum f\left(x_{i}\right) \cdot y_{i}$ from $\mathrm{E}^{\prime}$ to F .

Proof The map $\delta: \mathrm{E} \rightarrow \mathrm{E}^{\prime \top}$ defined by $\delta(x)=(f \mapsto f(x))$ is injective, and identifies E with a direct factor of $\mathrm{E}^{\prime \top}$. As a result, the rank of $t$ equals the rank of $\left(\delta \otimes 1_{\mathrm{F}}\right)(t) \in \mathrm{E}^{\prime \top} \otimes \mathrm{F}$. The rank of the latter tensor equals the rank of its image in $\operatorname{Hom}\left(\mathrm{E}^{\prime}, \mathrm{F}\right)$ under $\alpha_{\mathrm{F}}^{\mathrm{E}^{\prime}}$, and this image is the map $f \mapsto \sum f\left(x_{i}\right) \cdot y_{i}$.

## Exercises 3.9. (Kronecker Homomorphism)

1. (Norm $\varepsilon$ )— Let E and F be normed vector spaces over $\mathbb{R}$ or $\mathbb{C}$. Denote by $\mathrm{E}^{\top}$ the space of continuous linear forms on $E$.
(a) Define injections from $E \otimes F$ to the space $L\left(E^{\top}, F\right)$ of continuous linear forms from $\mathrm{E}^{\top}$ to F , to $\mathrm{L}\left(\mathrm{F}^{\top}, \mathrm{E}\right)$ and to the space $\mathrm{B}\left(\mathrm{E}^{\top}, \mathrm{F}^{\top}\right)$ of continuous bilinear forms on $\mathrm{E}^{\top} \times \mathrm{F}^{\top}$. Show that the norms induced on $\mathrm{E} \otimes \mathrm{F}$ by these three injections agree (use the Hahn-Banach theorem). Let $x \mapsto\|x\|_{\varepsilon}$ be the norm thus defined on $\mathrm{E} \otimes \mathrm{F}$.
(b) Let $\left(\xi_{i}\right)$ be a family of elements of $\mathrm{E}^{\top}$ such that $(\forall x \in \mathrm{E})\|x\|=\sup \left|\xi_{i}(x)\right|$. Show that, for $t \in \mathrm{E} \otimes \mathrm{F},\|t\|_{\varepsilon}=\sup \left\|\left(\xi_{i} \otimes 1_{\mathrm{F}}\right)(t)\right\|$. If $\left(\eta_{j}\right)$ is a family of elements of $\mathrm{F}^{\top}$ satisfying the same property, then $\|t\|_{\varepsilon}=\sup _{i, j}\left|\left(\xi_{i} \otimes \eta_{j}\right)(t)\right|$.
(c) Let X be a compact metric space and suppose that F is complete. Show that the space $C(X, F)$ of continuous maps from $X$ to $F$ can be identified with the completion relative to the norm $\varepsilon$ of $\mathrm{C}(\mathrm{X}) \otimes \mathrm{F}$.
(d) Let X and Y be compact metric spaces. Show that $\mathrm{C}(\mathrm{X} \times \mathrm{Y})$ is identified with the completion with respect to the norm $\varepsilon$ of $\mathrm{C}(\mathrm{X}) \otimes \mathrm{C}(\mathrm{Y})$.
(e) Let $\mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime}$ be subspaces of E and F respectively, equipped with the induced norms. Show that the norm $\varepsilon$ on $\mathrm{E}^{\prime} \otimes \mathrm{F}^{\prime}$ is induced by that of $\mathrm{E} \otimes \mathrm{F}$. Assuming that $\mathrm{E}^{\prime}$ and $\mathrm{F}^{\prime}$ closed, is the norm $\varepsilon$ on $\left(\mathrm{E} / \mathrm{E}^{\prime}\right) \otimes\left(\mathrm{F} / \mathrm{F}^{\prime}\right)$ the quotient of the norm $\varepsilon$ on $\mathrm{E} \otimes \mathrm{F}$ ?
(f) Show that the seminorm $\pi$ defined in [2.4, Exercise 3] is a norm.
2. (Nuclear norm)—Let E and F be finite dimensional normed spaces over $\mathbb{R}$ or $\mathbb{C}$. For all $t \in \mathrm{E}^{\top} \otimes \mathrm{F}$, let $\|t\|_{\pi}$ denote the lower bound of $\sum\left\|\xi_{i}\right\|\left\|y_{i}\right\|$ for the finite families $\left(\left(\xi_{i}, y_{i}\right)\right)$ such that $t=\sum \xi_{i} \otimes y_{i}$. The nuclear norm, written $f \mapsto\|f\|_{\nu}$, is the norm obtained by transferring the norm $\pi$ on $\mathrm{L}(\mathrm{E}, \mathrm{F})$ via the Kronecker isomorphism.
(a) Show that for all $f \in \mathrm{~L}(\mathrm{E}, \mathrm{E}),\|\operatorname{Tr}(f)\| \leqslant\|f\|_{\nu}$.
(b) Show that $\left\|1_{\mathrm{E}}\right\|_{\nu}=\operatorname{dim}(\mathrm{E})$ (use (a) gives one of the inequalities; to prove the other one, first study the dimension 2 case by considering a parallelogram with maximal area among all those with sides of length 1).
(c) Let F be the space $\mathbb{R}^{3}$ equipped with the norm $x \mapsto \sup \left|x_{i}\right|$, and E the plane $\left\{x \in \mathbb{R}^{3} \mid x_{1}+x_{2}+x_{3}=0\right\}$ equipped with the induced norm. Draw a unit ball of E . Compute the nuclear norm of the canonical injection from E to F .
(d) Let $f \in \mathrm{~L}(\mathrm{E}, \mathrm{F})$ and $\mathrm{F}_{1}$ a subspace de F containing the image of $f$ (resp. N a subspace de E contained in the kernel of $f$ ). Is the nuclear norm of $f$ in $\mathrm{L}\left(\mathrm{E}, \mathrm{F}_{1}\right)$ (resp. in $\mathrm{L}(\mathrm{E} / \mathrm{N}, \mathrm{F})$ ) necessarily equal to the nuclear norm of $f$ in $\mathrm{L}(\mathrm{E}, \mathrm{F})$ ?
(e) Take $f: \mathbb{R}^{3} \rightarrow \mathbb{R}^{3}$ to be the map represented by the matrix

$$
\left(\begin{array}{lll}
0 & 1 & 1 \\
1 & 0 & 1 \\
1 & 1 & 0
\end{array}\right)
$$

the space $\mathbb{R}^{3}$ being equipped with the norm $x \mapsto \sup \left|x_{i}\right|$. Compare the sum of the absolute values of the eigenvalues with the nuclear norm.
(f) Show that, for $f \in \mathrm{~L}(\mathrm{E}, \mathrm{F})$, the nuclear norm of $f$ is the upper bound of $\operatorname{Tr}(g \circ f)$ for $g: \mathrm{F} \rightarrow \mathrm{E}$ such $\|g\| \leqslant 1$, i.e. $(\forall y \in \mathrm{D})\|g(y)\| \leqslant\|y\|$.
3.-Let A be a ring, E and F finitely generated projective A -modules. Show that the bilinear form $(f, g) \mapsto \operatorname{Tr}(g \circ f)$ enables the identification of each of the modules $\operatorname{Hom}(E, F)$ and $\operatorname{Hom}(F, E)$ with the dual of the other.
4.-Let V be the real vector space of $n \times n$ matrices. Denote by $u^{\top}$ the transpose of the matrix $u$, and by $\operatorname{Tr}(u)$ its trace. Set $p(u)=\operatorname{Tr}(u \cdot u)$ and $q(u)=\operatorname{Tr}\left(u^{\top} \cdot u\right)$.
(a) Check that $p$ and $q$ are quadratic forms on V .
(b) If these quadratic forms are decomposed as sums of squares with signs of linearly independent forms, then in the case of $q$, how many have a + sign and how many a - one? In the case of $p$ ? (Start with the case $n=2$.)
(c) Is there a nonzero matrix $u$ such that $p(u)=0$ ? such that $q(u)=0$ ?
5.-Let $M$ be a square matrix with real entries. Show that

$$
\frac{d}{d t} \operatorname{det}(1+t \mathbf{M})_{t=0}=\operatorname{Tr}(\mathbf{M})
$$

(the computation may be done in $\mathbb{C}$ by reducing to the case when M is triangular). Compute the derivative for arbitrary $t$.
6.-Let K be a field and E a finite dimensional K -vector space. Show that every linear form $\theta: \operatorname{End}(\mathrm{E}) \rightarrow \mathrm{K}$ such that $(\forall f, g) \theta(f \circ g)=\theta(g \circ f)$ can be written $\lambda$ Tr. Can K be replaced by a ring by taking E to be a finitely generated free A-module? a finitely generated projective module?
7.-(a) Let E and F be finite dimensional vector spaces over a field $\mathrm{K}, f$ : $\mathrm{E} \rightarrow \mathrm{F}$ and $g: \mathrm{F} \rightarrow \mathrm{E}$ two homomorphisms such that $u=g \circ f$ and $v=f \circ g$ are of finite rank. Show that the nonzero eigenvalues of $u$ are eigenvalues of $v$ with the
same multiplicity. Deduce that $\operatorname{Tr}(v)=\operatorname{Tr}(u)$. (If necessary, extend the scalars to the algebraic closure of K.)
(b) Let E and F be two A-modules, $f: \mathrm{E} \rightarrow \mathrm{F}$ and $g: \mathrm{F} \rightarrow \mathrm{E}$ homomorphisms such that the traces of $u=g \circ f$ and $v=f \circ g$ are well defined. Assume A is reduced. Show that $\operatorname{Tr}(u)=\operatorname{Tr}(v)$.
(c) In question (b), can the assumption that A is reduced be omitted?
8.-(a) Let $\mathrm{P} \in \mathbb{C}[Z]$ be nonzero and A the algebra $\mathbb{C}[Z] /(\mathrm{P})$. For $f \in \mathrm{~A}$, let $\operatorname{Tr}_{\mathrm{A}}: \mathbb{C}(f)$ be the trace of the endomorphism $\mu_{f}: g \mapsto f \cdot g$ of the vector space A. What are the eigenvalues of $\mu_{f}$ ? Show that if F is a representative of $f$, then $\operatorname{Tr}_{\mathrm{A}: \mathbb{C}}(f)$ is the sum of the values of F at the roots of P .
(b) Let $\mathrm{P} \in \mathbb{C}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}, \mathrm{Z}\right]$ be a monic polynomial in Z (i.e. in $\mathbb{C}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right][\mathrm{Z}]$ ). Let $\mathrm{F} \in \mathbb{C}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}, \mathrm{Z}\right]$, and for $\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{C}^{n} \operatorname{set} \mathrm{~S}\left(x_{1}, \ldots, x_{n}\right)=\sum_{i} \mathrm{~F}\left(x_{1}, \ldots\right.$, $\left.x_{n}, z_{i}\right)$, where the $z_{i}$ are the roots of $\mathrm{P}\left(x_{1}, \ldots, x_{n}, \mathrm{Z}\right) \in \mathbb{C}[\mathrm{Z}]$, counted with their multiplicities. Show that the function $S: \mathbb{C}^{n} \rightarrow \mathbb{C}$ thus defined is polynomial.
9.-Let A be a ring and B an A -algebra, finitely generated and projective as A module. For all $f \in \mathrm{~B}$, let $\operatorname{Tr}_{\mathrm{B}: \mathrm{A}}(f)$ be the trace of the endomorphism $g \mapsto f \cdot g$ of the A-module B.
(a) Compute $\operatorname{Tr}_{\mathbb{C}: \mathbb{R}}(z)$ for $z \in \mathbb{C}$.
(b) Show that, if $\mathrm{B}=\mathrm{A}[\mathrm{X}] /\left(\mathrm{X}^{d}-a\right)$, with $a \in \mathrm{~A}$, then $\operatorname{Tr}\left(x^{k}\right)=0$ for $0<k<$ $d$. Give an example where $\operatorname{Tr}_{\mathrm{B}: \mathrm{A}}: \mathrm{B} \rightarrow \mathrm{A}$ is zero, while $\mathrm{B} \neq 0$.
(c) Let C be a B -algebra, finitely generated and projective as B -module. Show that C is a finitely generated A -module and that, if $h \in \mathrm{C}$, then $\operatorname{Tr}_{\mathrm{C}: \mathrm{A}}(h)=\operatorname{Tr}_{\mathrm{B}: \mathrm{A}}\left(\operatorname{Tr}_{\mathrm{C}: \mathrm{B}}(h)\right)$. Give a slightly more general statement.
10.-Let K be a field, $\mathrm{A}=\mathrm{K}[\mathrm{X}], \mathrm{B}=\mathrm{K}[\mathrm{Y}]$, and $h$ a polynomial of degree $d>0$. consider the algebra homomorphism $h^{*}: \mathrm{A} \rightarrow \mathrm{B}$ defined by $h^{*}(\mathrm{X})=h(\mathrm{Y})$, whereù $h^{*}(u)=u \circ h$ for $u \in \mathrm{~A}$.

Equip B with the A-module structure defined by this homomorphism.
(a) Show that B is a free A-module of rank $d$ with $\left(1, \ldots, \mathrm{Y}^{d-1}\right)$ as basis.
(b) For $f \in \mathrm{~B}$, let $h_{!}(f)$ be the trace of the endomorphism $g \mapsto f g$ of the A-module B. Show that $h_{!}(f)$ is a polynomial of degree $\leqslant \frac{n}{d}$, where $n$ is the degree of $f$. Show that, for $u \in \mathrm{~A}, h_{!}\left(h^{*}(u)\right)=d u$.
(c) Suppose that K is algebraically closed. Show that, for $x \in \mathrm{~K}$ and $f \in \mathbf{B}$, the scalar $h_{!}(f)(x)$ is the sum of $f(y)$ for $y \in \mathrm{~K}$ such that $h(y)=x$ (counted with their multiplicities as roots of $h(\mathrm{Y})-x)$.
(d) Let $f \in \mathrm{~B}$. Show that there is a polynomial $g \in \mathrm{~A}$ such that, for all $x \in \mathrm{~K}$, the scalar $g(x)$ is the product of all $f(y)$ for $y \in \mathrm{~K}$ such that $h(y)=x$ (resp. $g(x)=$ $\sum_{i<j} f\left(y_{i}\right) f\left(y_{j}\right)$, whereù $y_{1}, \ldots, y_{d}$ are the roots of $\left.h(\mathrm{Y})-x\right)$.

### 3.10 Chain Complexes

### 3.10.1 Graded Module

Definition A graded A-module is an A-module E. equipped with a sequence $\left(\mathrm{E}_{n}\right)_{n \in \mathbb{Z}}$ of submodules such that $\mathrm{E} .=\bigoplus \mathrm{E}_{n}$.

Let E . be a graded A-module. Equipping the module E . with the sequence $\mathrm{E}_{n}^{\prime}$ defined by $\mathrm{E}_{n}^{\prime}=\mathrm{E}_{n-k}$ gives a graded A-module E . $(k)$ obtained by diminishing the degrees by $k$.

The A-module E. is called positively (resp. negatively) graded if $\mathrm{E}_{n}=0$ for all $n<0$ (resp. $>0$ ). It is said to be bounded below (resp. above) if there exists $n_{0} \in \mathbb{Z}$ such that $\mathrm{E}_{n}=0$ for $n<n_{0}$ (resp. $n>n_{0}$ ); If it is bounded both below and above, it is said to be bounded.

Let E. and F. be graded A-modules. A degree $k$ morphism from E. to F. is a homomorphism $f$ from the module E . to the module F . such that $f\left(\mathrm{E}_{n}\right)$ is contained in $\mathrm{F}_{n+k}$ for all $n$. A degree $k$ morphism induces a homomorphism $f_{n}$ from $\mathrm{E}_{n}$ to $\mathrm{F}_{n+k}$ for all $n$, and $f \mapsto\left(f_{n}\right)_{n \in \mathbb{Z}}$ is a bijection of the set $\operatorname{Hom}_{k}(\mathrm{E} ., \mathrm{F}$.) of degree $k$ morphisms from E. to F. onto $\prod_{n \in \mathbb{Z}} \operatorname{Hom}\left(\mathrm{E}_{n}, \mathrm{~F}_{n+k}\right)$. Clearly,

$$
\operatorname{Hom}_{k}(\mathrm{E} ., \mathrm{F} .)=\operatorname{Hom}_{0}(\mathrm{E} .(k), \mathrm{F} .)=\operatorname{Hom}_{0}(\mathrm{E} ., \mathrm{F} .(-k)) .
$$

### 3.10.2 Chain Complexes

Definition A descending (resp. ascending) chain complex of A-modules is a graded A-module equipped with a degree -1 (resp. +1 ) endomorphism $d$ such that $d \circ d=0$.

Descending chain complexes are usually written with lower indices and ascending ones with upper ones. A descending (resp. ascending) complex is said to be right bounded if it is bounded below (resp. above), left bounded if its bounded above (resp. below). If $\mathrm{E}^{*}$ is a chain complex, a descending complex can be obtained by setting $\mathrm{E}_{n}=\mathrm{E}^{-n}$. This remark enables us to reduce the study of ascending complexes to that of complexes and conversely. In what follows, we study ascending complexes.

Let ( $\mathrm{E} ., d$ ) be a descending chain complex. The endomorphism $d$ is called the differential of the complex. $\operatorname{Set} \mathrm{Z}_{n}(\mathrm{E})=.\operatorname{Ker} d_{n}$ and $\mathrm{B}_{n}(\mathrm{E})=.\operatorname{Im} d_{n+1}$; the elements de $\mathrm{Z}_{n}(\mathrm{E}$.$) are the degree n$ cycles, those of $\mathrm{B}_{n}(\mathrm{E}$.$) the degree n$ boundaries, and $\mathrm{B}_{n}(\mathrm{E}.) \subset \mathrm{Z}_{n}(\mathrm{E}).$. Set $\mathrm{H}_{n}(\mathrm{E})=.\mathrm{Z}_{n}(\mathrm{E}.) / \mathrm{B}_{n}(\mathrm{E}$.$) .$

Definition $\mathrm{H}_{n}(\mathrm{E}$.$) is said to be the n$-th homology module of E . .
We have the following exact sequences

$$
0 \rightarrow \mathrm{Z}_{n} \rightarrow \mathrm{E}_{n} \rightarrow \mathrm{~B}_{n-1} \rightarrow 0
$$

$$
0 \rightarrow \mathrm{~B}_{n} \rightarrow \mathrm{Z}_{n} \rightarrow \mathrm{H}_{n} \rightarrow 0
$$

The module E. is said to be acyclic in degree $n$ if $\mathrm{H}_{n}(\mathrm{E})=$.0 , in other words if $\mathrm{E}_{n+1} \xrightarrow{d} \mathrm{E}_{n} \xrightarrow{d} \mathrm{E}_{n-1}$ is an exact sequence.

### 3.10.3 Morphisms of Chain Complexes

Let E. and F. be complexes. A morphism of complexes from E. to F. is a degree 0 morphism $f$ from the graded module E . to the graded module F . such that $d \circ f=$ $f \circ d$, where $d$ denotes the differential of E . as well as that of F . If $f: \mathrm{E}$. $\rightarrow \mathrm{F}$. is a morphism of chain complexes, $f$ induces a homomorphism from $\mathrm{Z}_{n}\left(\mathrm{E}\right.$.) to $\mathrm{Z}_{n}(\mathrm{~F}$.) and from $\mathrm{B}_{n}(\mathrm{E}$.$) to \mathrm{B}_{n}(\mathrm{~F}$.). Passing to the quotient, this gives a homomorphism $f_{*}: \mathrm{H}_{n}(\mathrm{E}.) \rightarrow \mathrm{H}_{n}(\mathrm{~F}$.$) .$

More generally, a degree $k$ morphism of complexes from E. to F. is a graded degree $k$ morphism $f: \mathrm{E} . \rightarrow$ F. such that $d \circ f=(-1)^{k} f \circ d$. This gives a homomorphism $f_{*}: \mathrm{H}_{n}(\mathrm{E}.) \rightarrow \mathrm{H}_{n+k}(\mathrm{~F}$.$) for all n$. The morphisms of chain complexes are the degree 0 morphisms.

The chain complex E. $(k)$ is defined by equipping the graded module E. $(k)$ with the differential $(-1)^{k} d$. The degree $k$ morphisms of chain complexes from E . to F . are the (degree 0 ) morphisms of chain complexes from E. ( $k$ ) to F. .

### 3.10.4 Homotopic Morphisms

Let E . and F . be two chain complexes.
Definition The degree $k$ morphisms $f$ and $g$ of complexes from E. to F. are said to be homotopic if there is a graded degree $k+1$ morphism $s:$ E. $\rightarrow$ F. such that $g-f=d \circ s+(-1)^{k} s \circ d$.

Proposition If $f$ and $g$ are homotopic degree $k$ morphisms from E . to F ., then the homomorphisms $f_{*}$ and $g_{*}$ from $\mathrm{H}_{n}(\mathrm{E}$.$) to \mathrm{H}_{n+k}(\mathrm{~F}$.$) are equal.$
Proof For $\quad x \in \mathrm{Z}_{n}(\mathrm{E}),. \quad g(x)-f(x)=d(s(x))+(-1)^{k} s(d(x))=d(s(x)) \in$ $\mathrm{B}_{n+k}(\mathrm{~F}$.$) , and so, denoting the canonical homomorphism by \chi, \mathrm{Z}_{n} \rightarrow \mathrm{H}_{n}$, $g_{*}(\chi(x))-f_{*}(\chi(x))=\chi(d(s(x)))=0$.
Remark Set $\operatorname{Hom}_{k}\left(\mathrm{E}_{\mathrm{I}}, \mathrm{F}.\right)$ to be the module of graded degree $k$ morphisms from E. to F., and $\operatorname{Hom}\left(\mathrm{E}_{\mathrm{E}}, \mathrm{F}.\right)$ the graded module $\bigoplus_{k} \operatorname{Hom}_{k}$ (E., F.). The module underlying $\operatorname{Hom}(E .$, F.) is a submodule of the module of module homomorphisms from E. to F.: a module homomorphism $h: \mathrm{E} . \rightarrow$ F. can be represented by a matrix $\left(h_{n, p}\right)$, where $h_{n, p}$ is a homomorphism from $\mathrm{E}_{n}$ to $\mathrm{F}_{p} . h \in \operatorname{Hom}(\mathrm{E} ., \mathrm{F}$.$) if and only if the support$ of the family $\left(h_{n, p}\right)$ is contained in a 'strip 'parallel to the diagonal", i.e. in some set $\{(n, p) \mid a \leqslant p-n \leqslant b\}$.

For $f \in \operatorname{Hom}_{k}(\mathrm{E} ., \mathrm{F}$.$) , set \mathrm{D}(f)=d \circ f+(-1)^{k-1} f \circ d$. We check that $\mathrm{D}(\mathrm{D}(f))=0$, and $\operatorname{Hom}(\mathrm{E} ., \mathrm{F}$.) equipped with the differential D is a chain complex. Degree $k$ cycles of this complex are degree $k$ morphisms of complexes from E. to F., and the $k$-th homology module of $\operatorname{Hom}(E .$, F.) is the module of homotopy classes of degree $k$ morphisms of complexes from E. to F. .

### 3.10.5

Proposition Let E., F., G. be chain complexes.
(a) The homotopy relation is an equivalence relation between degree $k$ morphisms of complexes from E. to F .;
(b) let $f_{0}$ and $f_{1}$ be degree $p$ morphisms of complexes from E. to F ., $g_{0}$ and $g_{1}$ degree $q$ morphisms of complexes from F . to G .. If $f_{1}$ is homotopic to $f_{0}$ and $g_{1}$ homotopic to $g_{0}$, then $g_{1} \circ f_{1}$ is homotopic to $g_{0} \circ f_{0}$.

Proof (a) follows from the previous remark. We prove (b).

1. Suppose that $g_{0}=g_{1}=g$ : if $f_{1}-f_{0}=d \circ s+(-1)^{p} s \circ d$, then $g \circ f_{1}-g \circ$ $f_{0}=g \circ d \circ s+(-1)^{p} g \circ s \circ d=(-1)^{q} d \circ g \circ s+(-1)^{p} g \circ s \circ d=$ $d \circ s^{\prime}+(-1)^{p+q} s^{\prime} \circ d$, with $s^{\prime}=(-1)^{q} g \circ s$.
2. Next suppose that $f_{0}=f_{1}=f$ : if $g_{1}-g_{0}=d \circ s+(-1)^{q} s \circ d$, then $g_{1} \circ$ $f-g_{0} \circ f=d \circ s \circ f+(-1)^{q} s \circ d \circ f=d \circ s \circ f+(-1)^{p+q} s \circ f \circ d$.
3. The general case follows from the previous ones.

Definition A morphism of chain complexes $f:$ E. $\rightarrow$ F. is a homotopy equivalence if there is a morphism of chain complexes $g: \mathrm{F} . \rightarrow \mathrm{E}$. such that $g \circ f$ and $f \circ g$ are respectively homotopic to the identity of E . and F .

Two complexes E. and F. are said to be homotopy equivalent if there is a homotopy equivalence from E. to F. . The relation thus defined is an equivalence relation between complexes.

### 3.10.6 Connecting Homomorphism

Let $0 \rightarrow \mathrm{E} . \xrightarrow{f} \mathrm{~F} . \xrightarrow{g}$ G. $\rightarrow 0$ be a short exact sequence, where $f$ and $g$ are degree 0 morphisms of complexes. We next define for all $n$ a homomorphism $\delta: \mathrm{H}_{n}(\mathrm{G}.) \rightarrow$ $\mathrm{H}_{n-1}$ (E.).

Let $\gamma \in \mathrm{H}_{n}$ (G.). Choose $z \in \mathrm{Z}_{n}(\mathrm{G}$.$) such that \chi(z)=\gamma$, where $\chi: \mathrm{Z}_{n} \rightarrow \mathrm{H}_{n}$ is a canonical homomorphism. Choose $y \in \mathrm{~F}_{n}$ such that $g(y)=z$; as $d(y) \in \mathrm{F}_{n-1}$ and $g(d(y))=0$, there is a unique $x \in \mathrm{E}_{n-1}$ such that $f(x)=d(y)$.

Since $f(d(x))=d(d(y))=0, d(x)=0$ as $f$ is injective. Set $\eta=\chi(x) \in \mathrm{H}_{n-1}(\mathrm{E}$.$) .$

Proposition and Definition The element $\eta$ is independent of the choices made. The map $\delta: \gamma \mapsto \eta$ from $\mathrm{H}_{n}(\mathrm{G}$.$) to \mathrm{H}_{n-1}(\mathrm{E}$.$) thus defined is a homomorphism called the$ connecting homomorphism.

Proof (a) Independence from the choice of $y$ : let $y^{\prime} \in \mathrm{F}_{n}$ such that $g\left(y^{\prime}\right)=z$; then $y^{\prime}=y+f(u)$ for some $u \in \mathrm{E}_{n}$. Then $d\left(y^{\prime}\right)=d(y)+d(f(u))$, and $d\left(y^{\prime}\right)=f\left(x^{\prime}\right)$ with $x^{\prime}=x+d(u)$; so $\chi\left(x^{\prime}\right)=\chi(x)$.
(b) Independence from the choice of $z:$ let $z^{\prime} \in \mathrm{Z}_{n}(\mathrm{G}$.$) be such that \chi\left(z^{\prime}\right)=\chi(z)=$ $\gamma$, then $z^{\prime}$ is of the form $z+d(w)$, with $w \in \mathrm{G}_{n+1}$, and $w$ is of the form $g(v)$, with $v \in \mathrm{~F}_{n+1}$. For $y^{\prime}$ such that $g\left(y^{\prime}\right)=z^{\prime}$, we may choose $y^{\prime}=y+d(v)$. Then $d\left(y^{\prime}\right)=d(y)=f(x)$.
(c) Linearity of $\delta$ : follows immediately.

### 3.10.7 Diagramme du serpent

Théorème. Soit $0 \rightarrow \mathrm{E} . \xrightarrow{f} \mathrm{~F} \xrightarrow{g}$ G. $\rightarrow 0$ une suite exacte courte où $f$ et $g$ sont des morphismes de complexes de degré 0 . Alors

est une suite exacte.

Proof (a) Exactness in $\mathrm{H}_{n}(\mathrm{~F}):. g_{*} \circ f_{*}=(g \circ f)_{*}=0$. Let $\beta \in \mathrm{H}_{n}(\mathrm{~F}$.$) be such that$ $g_{*}(\beta)=0$. We show that $\beta=f_{*}(\alpha)$ for some $\alpha \in \mathrm{H}_{n}(\mathrm{E}$.$) . Let y \in \mathrm{Z}_{n}(\mathrm{~F})$ be such that $\chi(y)=\beta$. As $g_{*}(\beta)=0, g(y)=d(w)$ for some $w \in \mathrm{G}_{n+1}$ and $w=g(v)$ for some $v \in \mathrm{~F}_{n+1}$ since $g$ is surjective. Set $y^{\prime}=y-d(v)$. As $\chi\left(y^{\prime}\right)=\chi(y)=\beta$, and $g\left(y^{\prime}\right)=0, y^{\prime}=f(x)$ for some $x \in \mathrm{E}_{n}$. As $f(d(x))=d\left(y^{\prime}\right)=0, d(x)=0$ since $f$ is injective, and $x \in \mathrm{Z}_{n}(\mathrm{E}$.$) . Setting \chi(x)=\alpha, f(\alpha)=\beta$.
(b) $\delta \circ g_{*}=0$ : if $\gamma=g_{*}(\beta)$ with $\beta \in \mathrm{H}_{n}(\mathrm{~F}$.$) , choose y_{1} \in \mathrm{Z}_{n}(\mathrm{~F}$.) such that $\chi\left(y_{1}\right)=\beta$. With the notation of 3.10 .6 , we may take $z=g\left(y_{1}\right)$ and $y=y_{1}$. Then $d(y)=0$, and so $\delta(\gamma)=0$.
(c) Exactness in $\mathrm{H}_{n}(\mathrm{G}$.$) : let y \in \mathrm{H}_{n}(\mathrm{G}$.) be such that $\delta(\gamma)=0$. We show that $\gamma=g_{*}(\beta)$ for some $\beta \in \mathrm{H}_{n}(\mathrm{~F}$.$) . With the notation of 3.10.6, x=d(u)$ for some $u \in \mathrm{E}_{n}$. Set $y^{\prime}=y-f(u)$; as $d\left(y^{\prime}\right)=0, y^{\prime} \in \mathrm{Z}_{n}\left(\mathrm{~F}\right.$.), and $g\left(y^{\prime}\right)=g(y)=z$; so $\gamma=g_{*}\left(\chi\left(y^{\prime}\right)\right)$.
(d) $f_{*} \circ \delta=0$ : with the notation of 3.10.6,

$$
f_{*}(\delta(\gamma))=f_{*}(\chi(x))=\chi(f(x))=\chi(d(y))=0
$$

(e) Exactness in $\mathrm{H}_{n-1}(\mathrm{E}$.$) : let \alpha \in \mathrm{H}_{n-1}(\mathrm{E}$.$) be such that f_{*}(\alpha)=0$. Choose $x \in$ $\mathrm{Z}_{n-1}(\mathrm{E}$.$) such that \chi(x)=\alpha$. As $f(x) \in \mathrm{B}_{n-1}(\mathrm{~F}),. f(x)=d(y)$ for some $y \in \mathrm{~F}_{n}$. Set $g(y)=z$. As $d(z)=g(d(y))=g(f(x))=0, z \in \mathrm{Z}_{n}(\mathrm{G}$. $)$, and $\alpha=\delta(\gamma)$ with $\gamma=\chi(z)$.

### 3.10.8 Mapping Cylinder

Let E . and F . be two chain complexes and $f: \mathrm{E} . \rightarrow \mathrm{F}$. a morphism of chain complexes. Set $\mathrm{M}_{n}=\mathrm{F}_{n} \oplus \mathrm{E}_{n-1}$ and let $d: \mathrm{M}_{n} \rightarrow \mathrm{M}_{n-1}$ be the homomorphism defined by the matrix

$$
\left(\begin{array}{cc}
d & f \\
0 & -d
\end{array}\right) .
$$

We check that $d \circ d=0: \mathrm{M}_{n} \rightarrow \mathrm{M}_{n-2}$. The complex $\mathrm{M} .=\bigoplus \mathrm{M}_{n}$ thus constructed is called the mapping cylinder of $f$.

Proposition With the above notation,
(a) there is a short exact sequence of complexes

$$
0 \rightarrow \text { F. } \xrightarrow{u} \text { M. } \xrightarrow{v} \text { E. (1) } \rightarrow 0,
$$

where $u_{n}: \mathrm{F}_{n} \rightarrow \mathrm{M}_{n}$ and $v_{n}: \mathrm{M}_{n} \rightarrow \mathrm{E}_{n-1}$ are the canonical injection and projection;
(b) the connecting homomorphism $\delta: \mathrm{H}_{n-1}(\mathrm{E}.) \rightarrow \mathrm{H}_{n-1}(\mathrm{~F}$.$) of this exact sequence$ is $f_{*}$.
Proof (a) is immediate; we show (b). Let $\gamma \in \mathrm{H}_{n-1}$ (E.) and choose $z \in \mathrm{Z}_{n-1}$ (E.) such that $\chi(z)=\gamma$. For $y \in \mathbf{M}_{n}$ such that $v(y)=z$ we may choose $y=\binom{0}{z}$. Then $d(y)=\binom{f(z)}{0}=u(x)$ with $x=f(z)$, and by construction $\delta(\gamma)=\chi(x)=$ $\chi(f(z))=f_{*}(\gamma)$.
Corollary 3.11 There is an exact sequence

$$
\cdots \rightarrow \mathrm{H}_{n}(\mathrm{E} .) \xrightarrow{f_{*}} \mathrm{H}_{n}(\text { F. }) \xrightarrow{u_{*}} \mathrm{H}_{n}(\text { M. }) \xrightarrow{v_{*}} \mathrm{H}_{n-1}(\mathrm{E} .) \rightarrow \cdots
$$

Corollary 3.12 Let $a, b \in \mathbb{Z}$ be such that $a \leqslant b . \mathrm{H}_{n}(\mathrm{M})=$.0 for $a \leqslant n \leqslant b$ if and only if $f_{*}: \mathrm{H}_{n}(\mathrm{E}.) \rightarrow \mathrm{H}_{n}(\mathrm{~F}$.) is an isomorphism for $a \leqslant n<b$, surjective for $n=b$ and injective for $n=a-1$.

### 3.10.9 Resolutions

Definition Let M be an A-module. A left (right) resolution of M is a descending positively graded chain complex L. (resp. ascending $L^{\bullet}$ ), acyclic in degrees $\neq 0$ and equipped with an isomorphism $\iota: \mathrm{H}_{0}(\mathrm{~L}.) \rightarrow \mathrm{M}$ (resp. $\iota: \mathrm{M} \rightarrow \mathrm{H}_{0}\left(\mathrm{~L}^{\circ}\right)$ ).

A left (resp. right) resolution of $M$ gives an exact sequence

$$
\begin{gathered}
\cdots \rightarrow \mathrm{L}_{n} \rightarrow \mathrm{~L}_{n-1} \rightarrow \cdots \rightarrow \mathrm{~L}_{1} \rightarrow \mathrm{~L}_{0} \rightarrow \mathrm{M} \rightarrow 0 \\
(\text { resp. } 0
\end{gathered} \rightarrow \mathrm{M} \rightarrow \mathrm{~L}^{0} \rightarrow \mathrm{~L}^{1} \rightarrow \cdots \rightarrow \mathrm{~L}^{n} \rightarrow \mathrm{~L}^{n+1} \rightarrow \cdots \text { ) } . ~ \$
$$

We say that L . is a free resolution (resp. finite free) of M if L . is a left resolution of M such that $\mathrm{L}_{n}$ is free (resp. finite free) for all $n$. Projective and flat resolutions are defined likewise; however an injective resolution is a right resolution such that $\mathrm{L}^{n}$ is injective for all $n$.

Proposition (a) All A-modules admit a free resolution.
(b) If A is a Noetherian ring, then all finitely generated A-modules admit a finite free resolution.

Proof (a) (resp. (b)): by (3.4.6, Remark 2), there is a pair $\left(\mathrm{L}_{0}, \varepsilon\right)$ where $\mathrm{L}_{0}$ is a free (resp. finite free) module libre and $\varepsilon: \mathrm{L}_{0} \rightarrow \mathrm{M}$ a surjective morphism, and a sequence of pairs $\left(\left(\mathrm{L}_{n}, d_{n}\right)\right)_{n \geqslant 1}$ where all $\mathrm{L}_{n}$ are free (resp. finite free) modules can be inductively constructed and so can homomorphisms $d_{n}: \mathrm{L}_{n} \rightarrow \mathrm{~L}_{n-1}$ such that $\operatorname{Im} d_{1}=\operatorname{Ker} \varepsilon$ and $\operatorname{Im} d_{n}=\operatorname{Ker} d_{n-1}$ for $n>1\left(\operatorname{Ker} d_{n-1}\right.$ is a finitely generated submodule of $\mathrm{L}_{n-1}$ in case (b)).

Let L . and $\mathrm{L}^{\prime}$. be two resolutions of M . A resolution morphism from L. to L . is a morphism of chain complexes $h: \mathrm{L} . \rightarrow$ L. such that $\iota^{\prime} \circ h_{*}=\iota$. More generally, let $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime}$ be an A-module homomorphism, L . a resolution of M and $\mathrm{L}_{\text {. }}$ a resolution of $\mathrm{M}^{\prime}$. A $f$-morphism from L . to $\mathrm{L}^{\prime}$. is a morphism of chain complexes $h: \mathrm{L} . \rightarrow \mathrm{L}^{\prime}$ for which the diagram

commutes.

### 3.10.10 Projective Resolutions

Let M be an A-module. A projective resolution of M is a left resolution L . of M such that $\mathrm{L}_{n}$ is a projective module for all $n$, or, equivalently, such that the module $\bigoplus \mathrm{L}_{n}$ is projective.

A free resolution is projective, and so all modules admit a projective resolution.

Theorem Let M and $\mathrm{M}^{\prime}$ be two A -modules, L . and $\mathrm{L}^{\prime}$. projective resolutions of M and $\mathrm{M}^{\prime}$ respectively, and $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime}$ a homomorphism. There exists a $f$-morphism $h: \mathrm{L} . \rightarrow \mathrm{L}^{\prime}$, unique up to homotopy.

This theorem is a particular case of the following proposition:
Proposition Let L. and E. be positively graded complexes. Suppose that L. is projective and E . acyclic in degrees $\neq 0$. Then, for all homomorphisms $f: \mathrm{H}_{0}(\mathrm{~L}.) \rightarrow$ $\mathrm{H}_{0}(\mathrm{E}$.$) , there is a morphism of chain complexes h$, unique up to homotopy, such that $f=h_{*}$.

Proof (a) Existence. Construct by induction a sequence $\left(h_{n}\right)_{n \in \mathbb{N}}$ of homomorphisms $h_{n}: \mathrm{L}_{n} \rightarrow \mathrm{E}_{n}$ for which the diagrams

commute. This is possible since $L_{0}$ being projective and $\varepsilon: \mathrm{E}_{0} \rightarrow \mathrm{H}_{0}(\mathrm{E}$. ) being surjective, we can complete the diagram

to obtain a commutative one with $h_{0}: \mathrm{L}_{0} \rightarrow \mathrm{E}_{0}$; suppose that $h_{n}$ has been defined such that $d_{n} \circ h_{n}=h_{n-1} \circ d_{n}$. Then $h_{n}$ maps $\mathrm{Z}_{n}(\mathrm{~L}$.$) to \mathrm{Z}_{n}(\mathrm{E}$.$) and, as \mathrm{L}_{n+1}$ is projective and $d_{n+1}: \mathrm{E}_{n+1} \rightarrow \mathrm{Z}_{n}(\mathrm{E}$.$) surjective, we can complete the diagram$

to obtain a commutative one with $h_{n+1}: \mathrm{L}_{n+1} \rightarrow \mathrm{E}_{n+1}$.
(b) Uniqueness up to homotopy. By difference, it suffices to show that, if $h$ : L. $\rightarrow$ E. is a morphism of chain complexes such that $h_{*}=0: H(L.) \rightarrow \mathrm{H}_{0}(\mathrm{E}$.$) ,$ and by induction there is a sequence $\left(s_{n}\right)_{n \in \mathbb{N}}$ of homomorphisms $s_{n}: \mathrm{L}_{n} \rightarrow \mathrm{E}_{n+1}$ such that $h_{0}=d_{1} \circ s_{0}$ and $h_{n}=d_{n+1} \circ s_{n}+s_{n-1} \circ d_{n}$ for $n \geqslant 1$. Since $h_{*}=0$, the homomorphism $h_{0}$ maps $\mathrm{L}_{0}$ to $\mathrm{B}_{0}(\mathrm{E}$.), and the diagram

where $\mathrm{L}_{0}$ is projective and $d_{1}$ surjective, can be completed to obtain a commutative diagram with $s_{0}: \mathrm{L}_{0} \rightarrow \mathrm{E}_{1}$. Suppose that $s_{n}$ has been defined such that $h_{n}=d_{n+1} \circ$ $s_{n}+s_{n-1} \circ d_{n}$, and set $s_{n+1}$ such that $h_{n+1}=d_{n+2} \circ s_{n+1}+s_{n} \circ d_{n+1}$, i.e. $d_{n+2} \circ$ $s_{n+1}=h_{n+1}-s_{n} \circ d_{n+1}$. As $d_{n+1} \circ s_{n} \circ d_{n+1}=0, h_{n+1}-s_{n} \circ d_{n+1}$ maps $\mathrm{L}_{n+1}$ to $\mathrm{Z}_{n+1}$ (E.) $=\operatorname{Im} d_{n+2}$ and this homomorphism can be written $d_{n+2} \circ s_{n+1}$.

Corollary Any two projective resolutions of an A-module M are homotopy equivalent.

### 3.10.11 Injective Modules

Definition An A-module E is said to be injective if, for all injective A-module homomorphisms $f: \mathrm{M}^{\prime} \rightarrow \mathrm{M}$ and all homomorphisms $g: \mathrm{M}^{\prime} \rightarrow \mathrm{E}$, there is a homomorphism $h: \mathrm{M} \rightarrow \mathrm{E}$ making the diagram

commutative.
In other words, E is injective if and only if the contravariant functor $\mathrm{M} \mapsto$ $\operatorname{Hom}(\mathrm{M}, \mathrm{E})$ from the category A- $\mathbb{M}$ od to itself transforms injective homomorphisms into surjective ones.

Proposition Let E be an A-module. The following conditions are equivalent:
(i) the module E is injective;
(ii) for all ideals I of A , and all homomorphisms $f: \mathrm{I} \rightarrow \mathrm{E}$, there exists $x \in \mathrm{E}$ such that, for all $a \in \mathrm{I}, f(a)=a x$.

Proof (i) $\Rightarrow$ (ii): follows readily by applying the definition of the canonical injection $\mathrm{I} \rightarrow \mathrm{A}$ and by noting that $\operatorname{Hom}(\mathrm{A}, \mathrm{E})=\mathrm{E}$.
(ii) $\Rightarrow$ (i): let $f: \mathrm{M}^{\prime} \rightarrow \mathrm{M}$ be an injective homomorphism and $g: \mathrm{M}^{\prime} \rightarrow \mathrm{E}$ a homomorphism. Let $\Omega$ be the set of submodules N of M equipped with a homomorphism $h: \mathrm{N} \rightarrow$ E. Define an order on $\Omega$ by $(\mathrm{N}, h) \leqslant\left(\mathrm{N}^{\prime}, h^{\prime}\right)$ if and only if $\mathrm{N} \subset \mathrm{N}^{\prime}$ and $\left.h^{\prime}\right|_{\mathrm{N}}=h$. Equipped with this order, $\Omega$ is inductive. Indeed, if $\left(\left(\mathrm{N}_{\lambda}, h_{\lambda}\right)\right)_{\lambda \in \Lambda}$ is a totally ordered family of elements of $\Omega$, then the pair ( $\mathrm{N}, h$ ), where $\mathrm{N}=\bigcup_{\lambda \in \Lambda} \mathrm{N}_{\lambda}$ and $h$ is obtained by gluing together all $h_{\lambda}$, is an upper bound for the family. The set
of upper bounds for $\left(\mathrm{M}^{\prime}, g\right)$ in $\Omega$ is also inductive. Hence, by Zorn's theorem, there is a maximal upper bound $(\mathrm{N}, h)$ for $\left(\mathrm{M}^{\prime}, g\right)$ in $\Omega$. Assuming that (ii) holds, we show by contradiction that, for such an element, $\mathrm{N}=\mathrm{M}$. Let $x \in \mathrm{M}-\mathrm{N}$, and I be the ideal consisting of $a \in \mathrm{~A}$ such that $a x \in \mathrm{~N}$. There exists $y \in \mathrm{E}$ such that, for all $a \in \mathrm{I}$, $h(a x)=a y$. The homomorphism $\bar{h}: \mathrm{N} \oplus \mathrm{A} \rightarrow \mathrm{E}$ defined by $\bar{h}(n, a)=h(n)+a y$ can be factorized as

where $\mathrm{N}^{\prime}=\mathrm{N}+\mathrm{A} x$ and $\chi(n, a)=n+a x$. The element $\left(\mathrm{N}^{\prime}, h^{\prime}\right)$ is a strict upper bound for ( $\mathrm{N}, h$ ), which is impossible.

Definition Suppose that A is an integral domain and E an A-module. An element $x \in \mathrm{E}$ is said to be divisible if for all $a \in \mathrm{~A}-\{0\}$, there exists $y \in \mathrm{E}$ such that $a y=x$. The module E is said to be divisible if all elements of E are divisible.

The above proposition has the following corollary:
Corollary Suppose that A is an integral domain.
(a) All injective A-modules are divisible.
(b) If A is a PID, all divisible modules are injective.

Proof Condition (ii) of the proposition always holds for the 0 ideal. As E is divisible it also holds for every nonzero principal ideal.
(a) If the module E is injective, (ii) holds for all ideals, and hence in particular for all principal ideals.
(b) If A is a PID, and (ii) holds for all principal ideals, it holds for all ideals.

### 3.10.12

The results of this subsection will be needed in the following one.
Let $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ be a ring homomorphism, and E an A-module. Set $\widetilde{\mathrm{E}}=\mathrm{Hom}_{\mathrm{A}}$ (B, E). Define a B-module structure on $\widetilde{\mathrm{E}}$ by setting $b f(x)=f(x b)$ for $b, x \in \mathrm{~B}$ and $f \in \widetilde{\mathrm{E}}$.

Proposition For all B-modules M,

$$
\operatorname{Hom}_{\mathrm{B}}(\mathrm{M}, \widetilde{\mathrm{E}})=\operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E})
$$

In other words, the functor $\mathrm{E} \mapsto \widetilde{\mathrm{E}}$ is the right adjoint of the scalar restriction functor.
More precisely, define $\alpha: \widetilde{\mathrm{E}} \rightarrow \mathrm{E}$ by $\alpha(f)=f(1)$. The map $\alpha$ is A-linear, and so $\alpha_{*}: \operatorname{Hom}_{\mathrm{B}}(\mathrm{M}, \widetilde{\mathrm{E}}) \rightarrow \operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E})$. For A-linear maps $h: \mathrm{M} \rightarrow \mathrm{E}$, define $\tilde{h}: \mathrm{M} \rightarrow$ $\widetilde{\mathrm{E}}$ by $\tilde{h}(x)(b)=h(b x)$, and define $\beta: \operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E}) \rightarrow \operatorname{Hom}_{\mathrm{B}}(\mathrm{M}, \widetilde{\mathrm{E}})$ by $\beta(h)=\tilde{h}$. The maps $\alpha_{*}$ and $\beta$ are mutually inverses.

Proof For $h \in \operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E}), \alpha_{*}(\beta(h))(x)=\alpha_{*}(\tilde{h})(x)=\tilde{h}(x)(1)=h(1 \cdot x)=h(x)$. For $u \in \operatorname{Hom}_{\mathrm{B}}(\mathrm{M}, \widetilde{\mathrm{E}})$, set $v=\alpha_{*}(u)$, i.e. $v(x)=u(x)(1)$. As $\tilde{v}(x)(b)=v(b x)=$ $u(b x)(1)=b u(x)(1)=u(x)(1 \cdot b)=u(x)(b), \tilde{v}=u$, i.e. $\beta \circ \alpha_{*}=1$.

Corollary If E is an injective A-module, then $\widetilde{\mathrm{E}}$ is an injective B -module.
Proof The functor $\mathrm{M} \mapsto \operatorname{Hom}_{\mathrm{B}}(\mathrm{M}, \widetilde{\mathrm{E}})=\operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E})$ is the composition of the scalar restriction functor $\mathrm{M} \mapsto \mathrm{M}$ from $\mathrm{B}-\mathscr{M} \sigma \boldsymbol{d}$ to $\mathrm{A}-\mathscr{M} \sigma \boldsymbol{d}$, which transforms injections into injections, and of the functor $\mathrm{M} \mapsto \operatorname{Hom}_{\mathrm{A}}(\mathrm{M}, \mathrm{E})$ which transforms injections into surjections.

### 3.10.13

Theorem Any module is isomorphic to a submodule of an injective module.

Proof (a) Suppose first that the ring A is principal. Let K be the field of fractions of $\mathrm{A}, \mathrm{M}$ an A-module, $\mathrm{A}^{(\mathrm{J})} \xrightarrow{u} \mathrm{~A}^{(\mathrm{I})} \rightarrow \mathrm{M}$ a presentation of M , and $\iota: \mathrm{A}^{(\mathrm{I})} \rightarrow \mathrm{K}^{(\mathrm{I})}$ the canonical injection. Set $\mathrm{E}=\operatorname{Coker}(\iota \circ u)$. The module E is divisible since it is the quotient of $\mathrm{K}^{(\mathrm{I})}$ which is divisible, and so E is injective, and M is isomorphic to a submodule of E .
(b) Now let A be arbitrary. Let M be an A-module. As $\mathbb{Z}$ is a PID, there is a $\mathbb{Z}$-linear injection $f$ from M to an injective $\mathbb{Z}$-module E . Set $\widetilde{\mathrm{M}}=\operatorname{Hom}_{\mathbb{Z}}(\mathrm{A}, \mathrm{M})$ and $\widetilde{\mathrm{E}}=\operatorname{Hom}_{\mathbb{Z}}(\mathrm{A}, \mathrm{E})$, and let $h: \mathrm{M} \rightarrow \widetilde{\mathrm{M}}$ be the map defined by $h(x)(a)=a x$. The maps $h$ and $f_{*}: \widetilde{\mathrm{M}} \rightarrow \widetilde{\mathrm{E}}$ are A-linear injections, and so $f_{*} h: \mathrm{M} \rightarrow \widetilde{\mathrm{E}}$ is an A-linear injection and by (3.10.12, corollary), E is injective.

### 3.10.14 Injective Resolutions

Let M be an A-module. An injective resolution of M is a right resolution $\mathrm{E}^{*}$ of M such that $\mathrm{E}^{n}$ is an injective module for all $n$.

Remark This does not imply the injectivity of the module $\mathrm{E}^{*}$ : a direct sum of injective modules is generally not injective. As can be checked, this is however the case when the ring A is Noetherian.

Theorem (a) All A-modules admit an injective resolution.
(b) Let M be an A-module. Any two injective resolutions of M are homotopy equivalent.
(c) Let M and $\mathrm{M}^{\prime}$ be two $\mathrm{A}-m o d u l e s, \mathrm{E}^{*}$ and $\mathrm{E}^{\circ}$ injective resolutions of M and $\mathrm{M}^{\prime}$ respectively, and $f: \mathrm{M} \rightarrow \mathrm{M}^{\prime}$ a homomorphism. Then there is a $f$-morphism $h: \mathrm{E}^{*} \rightarrow \mathrm{E}^{\prime}$. It is unique up to homotopy.

Proof It is similar to those of (3.10.9, Proposition, a), (3.10.10, Corollary and Theorem), by reversing the direction of the arrows.

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## Chapter 4 <br> Coverings

### 4.1 Spaces over B

### 4.1.1

Definition Let B be a topological space. A space over B is a topological space X together with a continuous map $\pi: \mathrm{X} \rightarrow \mathrm{B}$ called projection.

If $X$ is a space over $B$, a continuous section or simply a section of $X$ is a continuous map $\sigma: \mathrm{B} \rightarrow \mathrm{X}$ such that $\pi \circ \sigma=1_{\mathrm{B}}$. We say that $\sigma$ passes through a point $x \in \mathrm{X}$ if $\sigma(\pi(x))=x$.

For all $b \in \mathrm{~B}$, a fiber of X at $b$, written $\mathrm{X}(b)$, is the subspace $\pi^{-1}(b)$ of X .
Let X and Y be spaces over $\mathrm{B}, \pi$ and $\pi^{\prime}$ their projections. A B-morphism, or simply a morphism from X to Y is a continuous map $f: \mathrm{X} \rightarrow \mathrm{Y}$ for which the diagram

commutes. If $f: \mathrm{X} \rightarrow \mathrm{Y}$ is a B-morphism, the map $f$ induces a map $f_{b}: \mathrm{X}(b) \rightarrow$ $\mathrm{Y}(b)$ for all $b \in \mathrm{~B}$. Define a category B-Top by taking its objects to be the spaces over B and its morphisms the B-morphisms.

Let X and Y be spaces over $\mathrm{B}, \pi$ and $\pi^{\prime}$ their projections. The fiber product of X and Y , written $\mathrm{X} \times_{\mathrm{B}} \mathrm{Y}$, is the subspace of the topological space $\mathrm{X} \times \mathrm{Y}$ consisting of pairs $(x, y)$ such that $\pi(x)=\pi^{\prime}(y)$. The space $\mathrm{X} \times_{\mathrm{B}} \mathrm{Y}$ equipped with the map $(x, y) \mapsto \pi(x)$ is the product of X and Y in $\mathrm{B}-\mathscr{T} a p$. For all $b \in \mathrm{~B},(\mathrm{X} \times \mathrm{B} \mathrm{Y})(b)=$ $\mathrm{X}(b) \times \mathrm{Y}(b)$.

The sum of X and Y in $\mathrm{B}-\mathscr{T} a p$ is the disjoint union $\mathrm{X} \sqcup \mathrm{Y}$ equipped with the projection inducing $\pi$ on X and $\pi^{\prime}$ on Y. So

$$
(\mathrm{X} \sqcup \mathrm{Y})(b)=\mathrm{X}(b) \sqcup \mathrm{Y}(b)
$$

### 4.1.2 Change of Basis

Let B and $\mathrm{B}^{\prime}$ be topological spaces, X a space over B and $h: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ a continuous map. The topological space $\mathrm{X}^{\prime}=\mathrm{B}^{\prime} \times_{\mathrm{B}} \mathrm{X}$, equipped with the first projection $\mathrm{X}^{\prime} \rightarrow$ $\mathrm{B}^{\prime}$, is a space over $\mathrm{B}^{\prime}$ written $h^{*}(\mathrm{X})$ and called the space over $\mathrm{B}^{\prime}$ obtained from X by a change of basis from B to $\mathrm{B}^{\prime}$ along $h$. The fibre of $h^{*}(\mathrm{X})$ at a point $b^{\prime} \in \mathrm{B}^{\prime}$ can be identified with the fibre of X at $h\left(b^{\prime}\right)$.

Let A be a subspace of B , and $\iota: \mathrm{A} \rightarrow \mathrm{B}$ the canonical injection. The space $\iota^{*}(\mathrm{X})$ can be identified with $\pi^{-1}(\mathrm{~A})$, equipped with the projection onto A induced by $\pi$. This space is written $\left.\mathrm{X}\right|_{\mathrm{A}}$.

### 4.1.3 Hausdorff Spaces over B

A space X over B is said to be Hausdorff over B if the diagonal $\Delta_{\mathrm{X}}=\{(x, x) \in$ $\mathrm{X} \times \mathrm{X}\}$ is closed in $\mathrm{X} \times_{\mathrm{B}} \mathrm{X}$, i.e. if for all pairs $(x, y)$ of distinct elements lying in the same fibre of $X$, there are two neighbourhoods of $x$ and $y$ respectively, open in X and disjoint.

Remark If X is a Hausdorff space over B , then the fibres of X are Hausdorff.
If B is Hausdorff, the conditions " X is Hausdorff over B " and " X is Hausdorff" are equivalent.

### 4.1.4 Etale Spaces over B

Let X and Y be topological spaces. A continuous map $f: \mathrm{X} \rightarrow \mathrm{Y}$ is a local homeomorphism if, for all $x \in \mathrm{X}$, there is an open neighbourhood U of $x$ such that $f$ induces a homeomorphism from $U$ onto an open subset of $Y$.

A space $X$ over $B$ is said to be etale over $B$ if the projection is a local homeomorphism.

Example (Implicit function theorem) Let X and Y be $\mathscr{C}^{1}$-manifolds, $f: \mathrm{X} \rightarrow \mathrm{Y}$ a $\mathscr{C}^{1}$-map. If, for all $x \in \mathrm{X}$, the tangent linear map

$$
\mathrm{T}_{x} f: \mathrm{T}_{x} \mathrm{X} \rightarrow \mathrm{~T}_{f(x)} \mathrm{Y}
$$

is an isomorphism, then $f$ is a local homeomorphism.

### 4.1.5 Proper Spaces

The related concept of a proper space over B corresponds to that of a quasi-compact space.

Proposition and Definition Let X be a space over B. The following conditions are equivalent:
(i) Every filter $\mathscr{F}$ on X whose image $\pi_{*} \mathscr{F}$ in B converges to a point $b \in \mathrm{~B}$, has a cluster point $x$ such that $\pi(x)=b$;
(ii) for all $b \in \mathrm{~B}$, and open subsets $\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ of X satisfying

$$
\bigcup_{i \in \mathrm{I}} \mathrm{U}_{i} \supset \mathrm{X}(b),
$$

there is a finite subset J of I and a neighbourhood V of b such that

$$
\bigcup_{i \in \mathrm{~J}} \mathrm{U}_{i} \supset \pi^{-1}(\mathrm{~V}) ;
$$

(iii) The fibres of X are quasi-compact and the map $\pi$ is closed.

If these conditions hold, then X is proper over B .
Proof (i) $\Rightarrow$ (iii). Let $b \in \mathrm{~B}, \iota$ the injection $\mathrm{X}(b) \rightarrow \mathrm{X}$, and $\mathscr{G}$ a filter on $\mathrm{X}(b)$. Set $\mathscr{F}=\iota_{*} \mathscr{G}$, the filter $\pi_{*} \mathscr{F}$ is defined by $b$, and so $\mathscr{F}$ admits a cluster point $x \in \mathrm{X}(b)$. Then $x$ is a cluster point of $\mathscr{G}$. Since this holds for all filters $\mathscr{G}$, the space $\mathrm{X}(b)$ is quasi-compact.

Let Y be a closed subset of X and $b \in \mathrm{~B}$ a cluster point of $\pi(\mathrm{Y})$. Set $\mathscr{V}_{b}$ to be the filter of neighbourhoods of $b$ in B. For all U in $\mathscr{V}_{b}, \pi^{-1}(\mathrm{U}) \cap \mathrm{Y}$ generated a filter $\mathscr{F}$ on X. The filter $\pi_{*} \mathscr{F}$ converges to $b$, and so $\mathscr{F}$ admits a cluster point $x \in \mathrm{X}(b)$. Then $x$ is a cluster point of Y , and so $x \in \mathrm{Y}$ and $b \in \pi(\mathrm{Y})$. Therefore, $\pi(\mathrm{Y})$ is closed.
(iii) $\Rightarrow$ (ii). Let $b \in \mathrm{~B}$, and $\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ open subsets in X whose union cover $\mathrm{X}(b)$. As $\mathrm{X}(b)$ is quasi-compact, there is a finite subset J of I such that $\bigcup_{i \in \mathrm{~J}} \mathrm{U}_{i} \supset \mathrm{X}(b)$. Then $\mathrm{Y}=\mathrm{X}-\bigcup_{i \in \mathrm{~J}} \mathrm{U}_{i}$ is closed in X and does meet $\mathrm{X}(b)$, the set $\pi(\mathrm{Y})$ is closed in B and does not contain $b$, and $\mathrm{V}=\mathrm{B}-\pi(\mathrm{Y})$ is an open neighbourhood of $b$ with the desired property.
(ii) $\Rightarrow$ (i). Let $\mathscr{F}$ be a filter on X such that $\pi_{*} \mathscr{F}$ converges to $b \in \mathrm{~B}$.

Suppose that, for all $x \in \pi^{-1}(b)$ the filter $\mathscr{F}$ is incompatible with the filter of neighbourhoods of $x$, i.e., there exists an open neighbourhood $\mathrm{U}_{x}$ of $x$ and a subset $\mathrm{A}_{x} \in \mathscr{F}$ such that $\mathrm{A}_{x} \cap \mathrm{U}_{x}=\varnothing$. The subsets $\mathrm{U}_{x}$ form an open cover of $\pi^{-1}(b)$. Let J be a finite subset of X and V a neighbourhood of $b$ such that $\bigcup_{x \in \mathrm{~J}} \mathrm{U}_{x} \supset \pi^{-1}(\mathrm{~V})$. Set $\mathrm{A}=\bigcap_{x \in \mathrm{~J}} \mathrm{~A}_{x}$. Then $\pi(\mathrm{A}) \in \pi_{*} \mathscr{F}$ and $\pi(\mathrm{A}) \cap \mathrm{V}=\varnothing$; so $\pi_{*} \mathscr{F}$ does not converge to $b$, contradicting assumption.

Corollary If X is proper over B then for all $b \in \mathrm{~B}$, all neighbourhoods $\mathrm{X}(b)$ in X contain a subset $\pi^{-1}(\mathrm{~V})$, where V is a neighbourhood of $b$ in B .

Condition (ii) for I reduces to one element.

### 4.1.6

Proposition Let X be a Hausdorff space over B , where B is assumed to be locally compact. Then X is proper over B if and only if the inverse image of every compact subset of B is compact. If this is the case, then X is locally compact.

Proof Suppose first that X is proper over B. We first show that if B is compact, then so is X. Let $\mathscr{F}$ be a filter on X. The filter $\pi_{*} \mathscr{F}$ admits a cluster point in B since B is compact. Let $\mathscr{G}$ be a filter on B finer than $\pi_{*} \mathscr{F}$ and convergent. Then the filter $\mathscr{F} \wedge \pi^{*} \mathscr{G}$ admits a cluster point in X since $\pi$ is proper, and so $\mathscr{F}$ admits a cluster point.

In the general case, for any compact subset $K$ of $B$, the space $\pi^{-1}(\mathrm{~K})$ is proper over K , and hence compact.

We now prove the converse. Let $b \in \mathrm{~B}$ and K a compact neighbourhood of $b$. Every filter $\mathscr{F}$ on X such that $\pi_{*} \mathscr{F}$ converges to $b$ induces a compact filter on $\pi^{-1}(\mathrm{~K})$, and thus admits a cluster point. The projection of this value must be $b$ since B is Hausdorff. So $\pi$ is proper.

X is locally compact. Let $x \in \mathrm{X}$ and set $b=\pi(x)$. Let K be a compact neighbourhood of $b$. Then $\pi^{-1}(\mathrm{~K})$ is a compact neighbourhood of $x$.

### 4.1.7 Fibres

Let X be a space over B . A trivialization of X is an isomorphism $\tau: \mathrm{X} \rightarrow \mathrm{B} \times \mathrm{F}$ of spaces over B , where F is a topological space, the product $\mathrm{B} \times \mathrm{F}$ being equipped with its projection onto the first factor. We say that X is a trivial fibre bundle if it admits a trivialization. It is a locally trivial fibre bundle or simply a fibre bundle if for all $b \in \mathrm{~B}$, there is a neighbourhood U of $b$ in B such that $\left.\mathrm{X}\right|_{\mathrm{U}}$ is a trivial fibre bundle, or equivalently, if there exists an open cover $\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ of B such that, for all $i,\left.\mathrm{X}\right|_{\mathrm{U}_{i}}$ is a trivial fibre bundle. It is a fibre bundle with fibre F if all the fibres are homeomorphic to F .

Remark Let X be a fibre bundle over B. If B is connected, then the fibres of X are mutually homeomorphic. Indeed, the equivalences classes of the equivalence relation " $\mathrm{X}(b)$ is homeomorphic to $\mathrm{X}\left(b^{\prime}\right)$ " between elements $b, b^{\prime} \in \mathrm{B}$ are open; hence there is a unique class.

### 4.1.8 Example. The Möbius Strip

Consider the set X of pairs $(\mathrm{L}, x)$ where $\mathrm{L} \subset \mathbb{R}^{2}$ is a 1-dimensional vector subspace and $x \in \mathrm{~L}$. Let $\pi$ be the projection $(\mathrm{L}, x) \mapsto \mathrm{L}$ of X onto the set B of 1-dimensional vector subspaces. Identify $B$ with $\mathbb{R} / \pi \cdot \mathbb{Z}$, and $X$ with the quotient of $\mathbb{R}^{2}$ by the equivalence relation identifying $(\theta, \rho)$ with $\left(\theta+k \pi,(-1)^{k} \rho\right)$ for $k \in \mathbb{Z}$. Equip X and B with the topologies defined by these identifications, and X with the projection $\pi$. Then X is a fibre space with fibre $\mathbb{R}$ over B .

This fibre space is not trivial. Indeed a continuous section $\sigma$ of X is given by a continuous function $h: \theta \mapsto \rho$ from $\mathbb{R}$ to $\mathbb{R}$ satisfying $h(\theta+k \pi)=(-1)^{k} h(\theta)$. Given two continuous sections $\sigma_{1}$ and $\sigma_{2}$ represented by two functions $h_{1}$ and $h_{2}, g=$ $h_{2}-h_{1}$ vanishes at some point: indeed, if $g(\theta)>0$ the $g(\theta+\pi)<0$ and conversely. Hence there exists $b \in \mathrm{~B}$ such that $\sigma_{2}(b)=\sigma_{1}(b)$.

The set X can be identified with the set of affine lines in $\mathbb{R}^{2}$ by assigning to ( $\mathrm{L}, x$ ) the line D perpendicular to L at $x$.

This set can in turn be identified with the set of 1-dimensional vector spaces that are not vertical in $\mathbb{R}^{3}$, by assigning to every affine line D in the horizontal plane H with side 1 the 1-dimensional vector space $\Delta$ perpendicular to the plane containing O and D. In other words, X may be considered a projective plane $\mathbf{P}^{2} \mathbb{R}$ without a point. All these identifications are compatibles with the natural topologies.

On the other hand, the space X can be topologically embedded in $\mathbb{R}^{3}$ by

$$
(\theta, \rho) \mapsto\left(\begin{array}{c}
\cos 2 \theta(1+u(\rho) \cos \theta) \\
\sin 2 \theta(1+u(\rho) \cos \theta) \\
u(\rho) \sin \theta
\end{array}\right)
$$

where $u(\rho)=\frac{1}{\pi} \operatorname{Arctg} \rho$ varies from $-\frac{1}{2}$ to $\frac{1}{2}$ as $\rho$ varies from $-\infty$ to $+\infty$.
The closure M of the image of X in $\mathbb{R}^{3}$ is the image of $\mathbb{R} \times\left[-\frac{1}{2}, \frac{1}{2}\right]$ under

$$
(\theta, s) \mapsto\left(\begin{array}{c}
\cos 2 \theta(1+s \cos \theta) \\
\sin 2 \theta(1+s \cos \theta) \\
s \sin \theta
\end{array}\right)
$$

It is the Möbius strip, which may be regarded as the (non trivial) fibre bundle over the circle $S^{1}$ with fibre $\left[-\frac{1}{2}, \frac{1}{2}\right]$.


### 4.1.9 Fibre Bundle over a Segment

Theorem Every fibre bundle over a compact interval $[a, b]$ of $\mathbb{R}$ is trivial.
The proof of this theorem will be given in 4.1.11.
Remark This result holds for any interval in $\mathbb{R}$ (4.1, Exercise 12), not just a compact one.

### 4.1.10 Partition Subordinate to a Cover

Lemma and Definition Let X be a compact metric space and $\mathscr{U}=\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ an open cover of X . There exists $\mu>0$ such that every ball of radius $<\mu$ is contained in some $\mathrm{U}_{i}$. Every such $\mu$ is called a mesh of $\mathscr{U}$.
Proof The distance $\rho_{i}(x)$ from $x$ to $\mathrm{X}-\mathrm{U}_{i}$ (and $\infty$ if $\mathrm{U}_{i}=\mathrm{X}$ ) for $i \in \mathrm{I}$ and $x \in \mathrm{X}$ is a Lipschitz function with constant 1 or identically $+\infty$. So is $\rho(x)=\sup _{i}\left(\rho_{i}(x)\right)$; it is $>0$ at all points since $\left(\mathrm{U}_{i}\right)$ is a cover of X . A lower bound $\mu>0$ of the function $\rho$ is a mesh of $\mathscr{U}$.
Proposition and Definition Let $[a, b]$ be a compact interval in $\mathbb{R}$ and $\mathscr{U}=\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ an open cover of $[a, b]$. Then there is an increasing finite sequence $\left(t_{0}, \ldots, t_{n}\right)$ with $t_{0}=a, t_{n}=b$ such that

$$
(\forall k \in\{1, \ldots, n\})(\exists i \in \mathrm{I})\left[t_{k-1}, t_{k}\right] \subset \mathrm{U}_{i}
$$

The intervals $\left[t_{k-1}, t_{k}\right]$ are said to form a partition subordinate to $\mathscr{U}$.

Proof Choose $n$ such that $\frac{b-a}{n}$ are less than a mesh of $\mathscr{U}$ and set $t_{k}=a+k \frac{b-a}{n}$.

### 4.1.11 Gluing Trivializations

Lemma 4.1 Let E and F be two topological spaces, $f$ a map from E to F and $\left(\mathrm{A}_{k}\right)$ a finite closed cover of E . If $\left.f\right|_{\mathrm{A}_{k}}$ continuous for all $k$, then $f$ is continuous.

Proof of Lemma 4.1. Let $G$ be a closed subset of F . For all $k, f^{-1}(\mathrm{G}) \cap \mathrm{A}_{k}=$ $\left(\left.f\right|_{\mathrm{A}_{k}}\right)^{-1}(\mathrm{G})$ is closed in $\mathrm{A}_{k}$, hence in E , and $f^{-1}(\mathrm{G})=\bigcup_{k}\left(f^{-1}(\mathrm{G}) \cap \mathrm{A}_{k}\right)$ is closed.

Lemma 4.2 Let $a, b, c \in \mathbb{R}$ with $a<b<c$, and X a space over $[a, c]$. If $\left.\mathrm{X}\right|_{[a, b]}$ and $\left.\mathrm{X}\right|_{[b, c]}$ are trivial fibre bundles, then so is X over $[a, c]$.

Proof Let $\tau_{1}:\left.\mathrm{X}\right|_{[a, b]} \rightarrow[a, b] \times \mathrm{F}_{1}$ and $\tau_{2}:\left.\mathrm{X}\right|_{[b, c]} \rightarrow[b, c] \times \mathrm{F}_{2}$ be trivializations. Define a homeomorphism $\gamma: \mathrm{F}_{1} \rightarrow \mathrm{~F}_{2}$ by $(b, \gamma(y))=\tau_{2}\left(\tau_{1}^{-1}(b, y)\right)$. Setting $\tau(x)=(\operatorname{Id} \times \gamma)\left(\tau_{1}(x)\right)$ if $\left.x \in \mathrm{X}\right|_{[a, b]}$ and $\tau_{2}(x)$ if $\left.x \in \mathrm{X}\right|_{[b, c]}$ gives a bijection $\mathrm{X} \rightarrow[a, c] \times \mathrm{F}_{2}$. Applying Lemma 4.1 to $\tau$ and $\tau^{-1}$ shows that $\tau$ is a homeomorphism, hence a trivialization.

Proof of Theorem 4.1.9. Let X be a fibre bundle over $[a, b]$, and $\mathscr{U}=\left(\mathrm{U}_{i}\right)$ an open cover of $[a, b]$ such that $\left.X\right|_{\mathrm{U}_{i}}$ is trivial for all $i$. Let $\left(t_{0}, \ldots, t_{n}\right)$ be such that $\left[t_{k-1}, t_{k}\right]$ form a partition subordinate to $\mathscr{U}$. Applying Lemma 4.2, induction on $k$ shows that $\left.\mathrm{X}\right|_{\left[a, t_{k}\right]}$ is trivial for $k=1, \ldots, n$. For $k=n, \mathrm{X}$ is a trivial fibre bundle.

## Exercises 4.1. (Spaces over B)

1.-Let B be a topological space.
(a) Show that if X is a Hausdorff (resp. quasi-compact) space, then $\mathrm{B} \times \mathrm{X}$ is Hausdorff (resp. proper) over B.
(b) Let X be a space over B, Hausdorff (resp. etale, resp. proper) over B, $\mathrm{B}^{\prime}$ a topological space and $f: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ a continuous map. Show that $f^{*}(\mathrm{X})$ is Hausdorff (resp. etale, resp. proper) over $\mathrm{B}^{\prime}$.
(c) Let X and Y be spaces Hausdorff (resp. etale, resp. proper) over B. Show that $\mathrm{X} \times_{\mathrm{B}} \mathrm{Y}$ is Hausdorff (resp. etale, resp. proper) over B.
(d) Let $X$ be a space over B and $\left(U_{i}\right)$ an open cover of B. Show that, if $\left.X\right|_{\mathrm{U}_{i}}$ is Hausdorff (resp. etale, resp. proper) over $\mathrm{U}_{i}$ for all $i$, then the space X is Hausdorff (resp. etale, resp. proper) over B. Does the same hold for a locally finite closed cover?
2.-Let $X$ be a space over B.
(a) If X is proper over B , The fibre $\mathrm{X}(b)$ is quasi-compact for all $b \in \mathrm{~B}$.
(b) If X is quasi-compact, then X is proper over B .
(c) If B is quasi-compact, then X is proper over B if and only if X is quasi-compact.
(d) If $X$ is proper over $B$ then for any quasi-compact $K \subset B$, the space $\pi^{-1}(K)$ is quasi-compact.
3.-Let $X$ be a space proper over B.
(a) Every closed subset of X is proper over B .
(b) Let Y be a space over B and $f: \mathrm{X} \rightarrow \mathrm{Y}$ a B-morphism. Then $f(\mathrm{X})$ is proper over $\mathbf{B}$ and, if Y is Hausdorff over $\mathrm{B}, f(\mathrm{X})$ is closed in Y . If $f$ is bijective and Y Hausdorff over $\mathbf{B}$, then the map $f$ is a homeomorphism.
(c) Let $f: \mathrm{X} \rightarrow \mathbb{R}$ be a continuous map. If the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is surjective, then the function $g: \mathrm{B} \rightarrow \mathbb{R}$ defined by $g(b)=\sup _{x \in \pi^{-1}(b)} f(x)$ is upper semicontinuous. If moreover the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is an open map, then $g$ is continuous.
4.-A map $f: \mathrm{X} \rightarrow \mathrm{Y}$ is said to be proper (resp. Hausdorff) if X , equipped with $f$, is proper (resp. Hausdorff) over Y. Let $f: \mathrm{X} \rightarrow \mathrm{Y}$ and $g: \mathrm{Y} \rightarrow \mathrm{Z}$ be continuous maps.
(a) If $f$ and $g$ are proper, then so is $g \circ f$.
(b) If $g \circ f$ is proper and $g$ is Hausdorff, then $f$ is proper.
5.-Let B be a topological space. If K is a compact set, then all closed subsets of $\mathrm{B} \times \mathrm{K}$ are proper over B. Conversely, let X be a proper space over B. Show that if X is completely regular (2.6, Exercise 2), then there is a compact set K and a homeomorphism over B from $X$ onto a closed subset of $B \times K$ (take the Čech compactification of X [2.6, Exercise 2]).
6.-Let B be a topological space, X a locally compact space and Y a subspace of $\mathrm{B} \times \mathrm{X}$. Show that the following conditions are equivalent:
(i) Y is proper over B ;
(ii) Y is closed in $\mathrm{B} \times \mathrm{X}$ and, for all $b \in \mathrm{~B}$, there is a compact subset $\mathrm{K} \subset \mathrm{X}$ and a neighbourhood V of $b$ in B such that $\left.\mathrm{Y}\right|_{\mathrm{V}} \subset \mathrm{V} \times \mathrm{K}$.
7.-Let $X$ be a space over $B$, and $\pi$ its projection. Assume that $X$ and $B$ are metrizable (or simply that X is metrizable and that all points of B admit a countable fundamental system of neighbourhoods). The aim is to show that the following conditions are equivalent:
(i) X is proper over B ;
(ii) every sequence $\left(x_{n}\right)$ in X such that the sequence $\left(\pi\left(x_{n}\right)\right)$ converges to a point $b \in \mathrm{~B}$ admits an a cluster point in $\mathrm{X}(b)$.
(a) Show that (i) $\Rightarrow$ (ii).
(b) Choose a distance on X . Show that the following conditions are equivalent:
$(\alpha)$ every Cauchy filter of $\mathscr{F}$ over X such that $\pi_{*} \mathscr{F}$ converges to a point $b \in \mathrm{~B}$ converges to a point of $\mathrm{X}(b)$;
$(\beta)$ every Cauchy sequence $\left(x_{n}\right)$ in X such that the sequence $\left(\pi\left(x_{n}\right)\right)$ converges to a point $b \in \mathrm{~B}$ converges to a point of $\mathrm{X}(b)$.

If these conditions hold, then X will be said to be B -complete.
(c) X will be said to be pre-proper over B if, for all $b \in \mathrm{~B}$ and all $\varepsilon>0$, there is a neighbourhood V of $b$ such that $\pi^{-1}(\mathrm{~V})$ has a finite cover consisting of sets having diameter $\leqslant \varepsilon$. Show that if X is pre-proper over B , then every ultrafilter $\mathscr{U}$ over X such that $\pi_{*} \mathscr{U}$ converges in B is a Cauchy filter. Deduce that, if X is pre-proper over B and B -complete, then X is proper over B .
(d) Show that if X is not pre-proper over B , then there is a sequence $\left(x_{n}\right)$ in X and $a>0$ such that $d\left(x_{p}, x_{q}\right) \geqslant a$ for $p \neq q$, and the sequence $\left(\pi\left(x_{n}\right)\right)$ converges B. Conclude.
8.-Let X and Y be topological spaces and $f: \mathrm{X} \rightarrow \mathrm{Y}$ a continuous map. The map $f$ is said to be closed if, for any closed subset A of X , the image $f(\mathrm{~A})$ is closed in Y.
(a) Show that $f$ is closed if and only if for all $y \in \mathrm{Y}$, every neighbourhood of $f^{-1}(y)$ contains a set $f^{-1}(\mathrm{~V})$, where V is a neighbourhood of $y$ in Y .
(b) Show that every proper map is closed. Give an example of a non-proper closed map and an example of a non-proper map with compact fibres.
(c) Let X be a space over B . Show that X is proper over B if and only if for every space $\mathrm{B}^{\prime}$ and continuous map $h: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$, the projection of $h^{*}(\mathrm{X})$ onto $\mathrm{B}^{\prime}$ is closed. (To show $X$ is proper over $B$, the previous exercise may be used and $B^{\prime}$ taken to be $\mathbb{N} \cup\{\infty\}$ if X and B are metrizable; in the general case, $\mathrm{B}^{\prime}$ may be taken to be $X \sqcup\{\infty\}$, equipped with a topology inducing the discrete topology on $X$ and with a carefully chosen filter of neighbourhoods of $\infty$.)
9.-Let X be a topological space, $\left(\mathrm{A}_{i}\right)$ subspaces of X , and $f: \bigsqcup \mathrm{A}_{i} \rightarrow \mathrm{X}$ the map inducing the canonical injection on each $\mathrm{A}_{i}$. Show that, $f$ is proper if and only if $\left(\mathrm{A}_{i}\right)$ is a locally finite family of closed sets. When is $f$ etale? open?
10. (Descent)—Let B and $\mathrm{B}^{\prime}$ be topological spaces and $f: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ a continuous map. Let X be a space over B and set $\mathrm{X}^{\prime}=f^{*}(\mathrm{X})$. Assume that $f$ is open and surjective, or else proper and surjective. Show that, if $X^{\prime}$ is proper (resp. etale, resp. Hausdorff) over $\mathrm{B}^{\prime}$, then the space X is proper (resp. etale, resp. Hausdorff) over B. Show that, if $\pi^{\prime}: \mathrm{X}^{\prime} \rightarrow \mathrm{B}^{\prime}$ is an open map, then the same holds for $\pi: \mathrm{X} \rightarrow \mathrm{B}$.
11. (Classical fibrations)—Consider $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ ) equipped with the Euclidean (resp. Hermitian) norm. Define the following spaces:

- $\mathrm{S}^{n-1}$ (resp. $\mathrm{S}^{2 n-1}$ ) is the unit sphere of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ );
- $\mathrm{O}_{n}$ (resp. $\mathrm{U}_{n}$ ) is the space of linear isometries of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ );
- $\mathrm{SO}_{n}\left(\right.$ resp. $\left.\mathrm{SU}_{n}\right)$ is the subspace of $\mathrm{O}_{n}\left(\right.$ resp. $\left.\mathrm{U}_{n}\right)$ of maps having determinant +1 ;
- for $p \leqslant n$, the Stiefel manifold $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)\left(\right.$ resp. $\left.\mathrm{V}_{p}\left(\mathbb{C}^{n}\right)\right)$ is the space of orthonormal families $\left(x_{1}, \ldots, x_{p}\right)$ of vectors of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ ), the Grassmannian $\mathrm{G}_{p}\left(\mathbb{R}^{n}\right)$ (resp. $\mathrm{G}_{p}\left(\mathbb{C}^{n}\right)$ ) is the set of $p$-dimensional subspaces of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ );
- for integers $p_{1}, p_{2}, p_{3}$ such that $p_{1} \leqslant p_{2} \leqslant p_{3} \leqslant n, \mathrm{D}_{p_{1}, p_{2}, p_{3}}$ is the set of families $\left(\mathrm{F}_{1}, \mathrm{~F}_{2}, \mathrm{~F}_{3}\right)$ of subspaces of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ ), where $\mathrm{F}_{i}$ has dimension $p_{i}$, such that

$$
\mathrm{F}_{1} \subset \mathrm{~F}_{2} \subset \mathrm{~F}_{3} .
$$

Consider the surjective map $\psi: \mathrm{V}_{p}\left(\mathbb{R}^{n}\right) \rightarrow \mathrm{G}_{p}\left(\mathbb{R}^{n}\right)$ which assigns to every family the subspace it generates. Equip $\mathrm{G}_{p}\left(\mathbb{R}^{n}\right)$ with the quotient topology of $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)$ and $\mathrm{D}_{p_{1}, p_{2}, p_{3}}\left(\mathbb{R}^{n}\right)$ with the topology induced by that of $\mathrm{G}_{p_{1}} \times \mathrm{G}_{p_{2}} \times \mathrm{G}_{p_{3}}$.

Show that the map $\psi$ is a fibration. What is its fibre? Same questions for the following maps:

- $\left(x_{1}, \ldots, x_{p_{3}}\right) \mapsto\left(\mathrm{F}_{1}, \mathrm{~F}_{2}, \mathrm{~F}_{3}\right)$ from $\mathrm{V}_{p_{3}}$ to $\mathrm{D}_{p_{1} p_{2} p_{3}}$, where $\mathrm{F}_{i}$ is the subspace generated by $x_{1}, \ldots, x_{p_{i}}$;
- $\left(\mathrm{F}_{1}, \mathrm{~F}_{2}, \mathrm{~F}_{3}\right) \mapsto \mathrm{F}_{i}$ from $\mathrm{D}_{p_{1}, p_{2}, p_{3}}$ to $\mathrm{G}_{p_{i}}$ for $i=1,2,3$;
- for $q \leqslant p,\left(x_{1}, \ldots, x_{p}\right) \mapsto\left(x_{1}, \ldots, x_{q}\right)$ from $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)\left(\right.$ resp. $\left.\mathrm{V}_{p}\left(\mathbb{C}^{n}\right)\right)$ onto $\mathrm{V}_{q}\left(\mathbb{R}^{n}\right)$ $\left(\operatorname{resp} . \mathrm{V}_{q}\left(\mathbb{C}^{n}\right)\right)$.
This in particular gives a fibration of $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)$ over $\mathrm{S}^{n-1}$ with fibre $\mathrm{V}_{p-1}\left(\mathbb{R}^{n-1}\right)$, of $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)$ over $\mathrm{V}_{p-1}\left(\mathbb{R}^{n}\right)$ with fibre $\mathrm{S}^{n-p}$, of $\mathrm{O}_{n}$ over $\mathrm{S}^{n-1}$ with fibre $\mathrm{O}_{n-1}$, of $\mathrm{U}_{n}$ over $\mathrm{S}^{2 n-1}$ with fibre $\mathrm{U}_{n-1}$, of $\mathrm{S}^{2 n-1}$ over $\mathrm{P}_{n-1}(\mathbb{C})=\mathrm{G}_{1}\left(\mathbb{C}^{n}\right)$ with fibre $\mathrm{U}_{1} \approx \mathrm{~S}^{1}$. More particularly, since $P_{1}(\mathbb{C}) \approx S^{2}$ (Riemann sphere), this gives a fibration of $S^{3}$ over $S^{2}$ with fibre $S^{1}$, called the Hopf fibration. Show that $S^{3}$ is not homeomorphic to $S^{2} \times S^{1}$.
12.-(a) Let $\mathscr{U}=\left(U_{i}\right)$ be an open cover of $\mathbb{R}$. Show that there is an increasing family $\left(t_{k}\right)_{k \in \mathbb{Z}}$ of points of $\mathbb{R}$ such that for all $k$, the interval $\left[t_{k-1}, t_{k}\right]$ is contained in some $U_{i}$ and such that the union of these intervals is $\mathbb{R}$.
(b) Show that a fibre bundle over a not necessarily compact interval in $\mathbb{R}$ is trivial.


### 4.2 Locally Connected Spaces

### 4.2.1

Definition A topological space X is said to be locally connected if all points of X admit a fundamental system of connected neighbourhoods.

This definition is equivalent to the following condition:
(CCOO) Every connected component of open subsets of X is open in X .
Indeed, suppose that X is locally connected. Let U be an open subset of X and $\mathrm{U}^{\prime}$ a connected component of U . For all $x \in \mathrm{U}^{\prime}$, the set U contains a connected neighbourhood V of $x$. Then $\mathrm{U}^{\prime} \supset \mathrm{V}$, and so $\mathrm{U}^{\prime}$ is a neighbourhood of $x$, proving (CCOO). Conversely, suppose that condition (CCOO) holds. Then, for all $x \in \mathrm{X}$ and all open neighbourhoods U of X , the connected component of U containing $x$ is open, and so is a neighbourhood of $x$ contained in $U$.

### 4.2.2 Examples

The space $\mathbb{R}^{n}$ is locally connected. Every topological manifold, i.e. every space locally homeomorphic to $\mathbb{R}^{n}$, is locally connected.

A finite union of locally connected closed sets is locally connected. More generally:

Proposition Let X be a topological space and Y a proper space over X . If Y is locally connected and the projection $\pi: \mathrm{Y} \rightarrow \mathrm{X}$ is surjective, then the space X is locally connected.

Proof Let $x \in \mathrm{X}$ and U be a neighbourhood of $x$ in X ; set $\mathrm{F}=\pi^{-1}(x)$. The intersection of the union W of connected components of $\pi^{-1}(\mathrm{U})$ and of F is a neighbourhood of F . By 4.1.5, corollary, it contains a set $\pi^{-1}(\mathrm{~V})$, where V is a neighbourhood of $x$. Then $\pi(\mathrm{W})$ contains V ; so it is a connected neighbourhood of $x$ in U .

### 4.2.3 Counterexamples

A totally discontinuous space is locally connected only if it is discrete. Thus $\mathbb{Q}$ is not locally connected. Non-finite profinite sets are not locally connected. A compact space is locally connected only if it has finitely many connected components.

In $\mathbb{T}^{2}=\mathbb{R}^{2} / \mathbb{Z}^{2}$, the image of a line with irrational slope in $\mathbb{R}^{2}$ is an arc-connected set, but is not locally connected.

Let $\mathbf{m}=\left(m_{n}\right)$ be a sequence of numbers tending to infinity, such that $m_{n}$ divides $m_{n+1}$ for all $n$. The solenoid $\mathrm{S}_{\mathbf{m}}=\lim \mathbb{R} / \mathbb{Z} . m_{n}$ is a connected compact space, but is neither locally connected, nor arc-connected.

We next give three examples of connected compact subsets of $\mathbb{R}^{2}$, but not locally connected.

Examples (1) Comb Space. It is the set

$$
\mathrm{K}_{1}=([0,1] \times\{0\}) \cup(\mathrm{A} \times[0,1]),
$$

where $\mathrm{A}=\{0\} \cup\left\{\frac{1}{n}\right\}_{n \in \mathbb{N} *}$.
(2) Topologists' sine curve. It is the set

$$
\left.\left.\mathrm{K}_{2}=\left\{(x, y) \in \mathbb{R}^{2} \mid x \in\right] 0,1\right], y=\sin \frac{1}{x}\right\} \cup\{0\} \times[-1,1] .
$$

(3) Spiral around circle. It is the set

$$
\mathrm{K}_{3}=\left\{\left(1+\frac{1}{t}\right) \cdot e^{i t}\right\}_{t \in[1, \infty[ } \cup S^{1} \subset \mathbb{C},
$$

where $S^{1}$ is the unit circle.
Note that $K_{1}$ is arc-connected, whereas $K_{2}$ and $K_{3}$ are not.

### 4.2.4

Theorem Every connected and locally connected compact metrizable space is arcconnected.

To prove this theorem, we introduce the notion of uniformly locally connected metric space. Then the theorem follows from Propositions 4.2.6 and 4.2.8 given below.

### 4.2.5

Definition Let X be a metric space and $\mathrm{B}(x, r)$ the closed ball in X centered at $x$ with radius $r$. Then X is uniformly locally connected if for all $\varepsilon>0$, there exists $\alpha>0$ such that, for all $x \in \mathrm{X}, \mathrm{X}$ contains a connected set L with $\mathrm{B}(x, \alpha) \subset \mathrm{L} \subset \mathrm{B}(x, \varepsilon)$.

Or equivalently, there is a function $h:[0, a[\rightarrow \mathbb{R}$, continuous at 0 with $h(0)=0$ and $a>0$, such that, for all $x, x^{\prime} \in \mathrm{X}$ with $d\left(x, x^{\prime}\right)<a$, there is a connected set L having diameter $\leqslant h\left(d\left(x, x^{\prime}\right)\right)$ containing $x$ and $x^{\prime}$. Such a function is called a modulus of local connectivity (an allusion to uniform continuity and to the notion of the modulus of continuity).

### 4.2.6

Proposition Every locally connected metric space is uniformly locally connected.
Proof Let $\varepsilon>0$. All $x \in \mathrm{X}$ have a connected neighbourhood L contained in $\mathrm{B}\left(x, \frac{\varepsilon}{2}\right)$, and hence having diameter $\leqslant \varepsilon$; For such L , there is open subset $\mathrm{U} \leqslant \mathrm{L}$ containing $x$. Hence there is an open cover $\mathscr{U}=\left(U_{i}\right)$ of X such that each $\mathrm{U}_{i}$ is contained in a connected set having diameter $\leqslant \varepsilon$. Let $\mu$ be a mesh of $\mathscr{U}$ (4.1.10) and $\alpha<\mu$. Then

$$
(\forall x \in \mathrm{X})(\exists i) \mathrm{B}(x, \alpha) \subset \mathrm{U}_{i} \subset \mathrm{~L}_{i} \subset \mathrm{~B}(x, \varepsilon) .
$$

### 4.2.7 Embeddings in a Banach Space

The above Proposition 4.2.8 is easier to prove for a closed subset of a Banach space. The next lemma shows that this does not lessen its generality.

Lemma For every metric space X , there is an isometric embedding $\iota$ of X in a Banach space E .

Proof Suppose first that E is the vector space of bounded continuous functions on X , equipped with the norm $\|f\|=\sup _{x \in \mathrm{X}}|f(x)|$. If X has finite diameter, then we may take $\iota(x)=h_{x}$, where $h_{x}\left(x^{\prime}\right)=d\left(x, x^{\prime}\right)$. In the general case, assuming $\mathrm{X} \neq \varnothing$, $\iota(x)=h_{x}-h_{x_{0}}$ for some $x_{0} \in \mathrm{X}$. In all cases, this gives an isometric embedding of X in E .

### 4.2.8

Proposition Every complete uniformly locally connected metric space is arcconnected.

Proof Thanks to Lemma 4.2.7, X may be assumed to be in some Banach space E, in which it is closed since it is complete. Let $x, y \in X$. Set $t_{0,0}=0, t_{0,1}=1, x_{0,0}=x$, $x_{0,1}=y$.

Let $\left(\varepsilon_{n}\right)$ be a sequence of reals $>0$ such that $\sum \varepsilon_{n}<\infty$, and for all $n$ let $\alpha_{n}>0$ such that any two points at a distance $\leqslant \alpha_{n}$ are in the same connected set having diameter $\leqslant \varepsilon_{n}$. For all $n$, we inductively construct a finite sequence $\left(t_{n, 0}, \ldots, t_{n, k(n)}\right)$ in $\mathrm{I}=[0,1]$ with $0=t_{n, 0}<t_{n, 1}<\cdots<t_{n, k(n)}=1$, and a finite sequence $\left(x_{n, 0}, \ldots, x_{n, k(n)}\right)$ in X in such a way that, setting $\sigma_{n}=\left\{t_{n, i}\right\}_{0 \leqslant i \leqslant k(n)}$ :

- $\left\|x_{n, i}-x_{n, i-1}\right\| \leqslant \alpha_{n}$ for $i=1, \ldots, k(n), n \geqslant 1$.
- $\sigma_{n+1} \supset \sigma_{n}$
- $x_{n+1, j}=x_{n, i}$ if $t_{n+1, j}=t_{n, i}$
- $\left\|x_{n+1, j}-x_{n, i-1}\right\| \leqslant \varepsilon_{n}$ and $\left\|x_{n+1, j}-x_{n, i}\right\| \leqslant \varepsilon_{n}$ if $t_{n, i-1} \leqslant t_{n+1, j} \leqslant t_{n, i}$.
$t_{n, i}$ and $x_{n, i}$ being defined for $0 \leqslant i \leqslant k(n)$, the points $x_{n, i-1}$ and $x_{n, i}$ are in a connected set $\mathrm{L}_{n, i}$ having diameter $\leqslant \varepsilon_{n}$. There is a finite sequence $\left(x_{n+1, i, s}\right)_{0 \leqslant s \leqslant k(n+1, i)}$ with $x_{n+1, i, 0}=x_{n, i-1}, x_{n+1, i, k(n+1, i)}=x_{n, i},\left\|x_{n+1, i, s}-x_{n+1, i, s-1}\right\| \leqslant \alpha_{n}$. Choose an arbitrary sequence $\left(t_{n+1, i, s}\right)$ with

$$
t_{n, i-1}=t_{n+1, i, 0}<t_{n+1, i, 1}<\cdots<t_{n+1, i, k(n+1, i)}=t_{n, i} .
$$

Set

$$
\begin{gathered}
\sigma_{n+1, i}=\left\{t_{n+1, i, s}\right\}_{0 \leqslant s \leqslant k(n+1, i)}, \\
\sigma_{n+1}=\bigcup \sigma_{n+1, i}=\left\{t_{n+1,0}, \ldots, t_{n+1, k(n+1)}\right\}
\end{gathered}
$$

with $0=t_{n+1,0}<\cdots<t_{n+1, k(n+1)}=1, k(n+1)=\sum_{i} k(n+1, i)$. This ends the inductive construction.

Then define an affine function $\gamma_{n}: \mathrm{I} \rightarrow \mathrm{E}$ over $\left[t_{n, i-1}, t_{n, i}\right]$ by $\gamma_{n}\left(t_{n, i}\right)=x_{n, i}$. So $\gamma_{n}(0)=x, \gamma_{n}(1)=y,\left\|\gamma_{n+1}(t)-\gamma_{n}(t)\right\| \leqslant \varepsilon_{n}$ and $(\forall t \in \mathrm{I}) d\left(\gamma_{n}(t), \mathrm{X}\right) \leqslant \alpha_{n}$.

Therefore the functions $\gamma_{n}$ form a Cauchy sequence in the space $\mathrm{C}(\mathrm{I}, \mathrm{E})$, converging to a path in X from $x$ to $y$.

## Exercises 4.2. (Locally connected spaces)

1.-Let X be a topological space. Show that there is a surjective continuous map $h$ : $[0,1] \rightarrow X$ if and only if X is compact, metrizable, connected and locally connected.
2.-A space X is locally connected (resp. openly locally connected) at a point $x_{0}$ if $x_{0}$ admits a fundamental system of connected (resp. open connected) neighbourhoods. The aim here is to construct an example of a locally connected space at a point but not openly locally connected at that point.

Let $\mathrm{I}_{p, q}$ be the segment in $\mathbb{R}^{2}$ having as endpoints $\left(\frac{1}{p}, 0\right)$ and $\left(\frac{1}{p+1}, \frac{1}{(p+1) p q}\right)$. Set

$$
\mathrm{X}=[0,1] \times\{0\} \cup \bigcup_{p \geqslant 1, q \geqslant 1} \mathrm{I}_{p, q}
$$

and $x_{0}=(0,0)$.
(a) Show that X is locally connected at $x_{0}$.
(b) Show that $[0,1] \times\{0\}$ is in every open connected subset of $X$ containing $x_{0}$.
(c) Show that a locally connected space (i.e. locally connected at all points) is openly locally connected at all points.
3.-(a) Given two metric spaces X and Y , a map $f: \mathrm{X} \rightarrow \mathrm{Y}$ and an increasing continuous function $h:\left[0, a\left[\rightarrow \mathbb{R}_{+}\right.\right.$, with $h(0)=0$ and $a>0, f$ is said to admit $h$ as a modulus of continuity if $d\left(f(x), f\left(x^{\prime}\right)\right) \leqslant h\left(d\left(x, x^{\prime}\right)\right)$ for all pairs $\left(x, x^{\prime}\right)$ such that $d\left(x, x^{\prime}\right)<a$.

Show that a map admits a modulus of continuity if and only if it is uniformly continuous.
(b) Given a metric space X and a function $h$ as above, X is said to admit $h$ as a modulus of local connectivity if, for all $\left(x, x^{\prime}\right)$ such that $d\left(x, x^{\prime}\right)<a$, there is a connected subset L of X having diameter $\leqslant h\left(d\left(x, x^{\prime}\right)\right)$ and containing $x$ and $x^{\prime}$.

Show that a metric space admits a modulus of local connectivity if and only if it is uniformly locally connected.
4.-(a) Let X be a connected compact metric space. Assuming that the diameter D of X and a modulus of local connectivity of X are known, is it possible to determine a modulus of continuity of a path from $x$ to $x^{\prime}$ in X for $r \leqslant \mathrm{D}$, for all $\left(x, x^{\prime}\right)$ such that $d\left(x, x^{\prime}\right)=r$ ?
(b) Let X be a metric space. Assume there is a surjective continuous map $g$ : $[0,1] \rightarrow \mathrm{X}$. Then $g$ is uniformly continuous and X compact and locally connected, and so uniformly locally connected.

If a modulus of continuity of $g$ is known, is it possible to find a modulus of local connectivity of X?
5.-Show that none of the three spaces given in 4.2 .3 as examples of compact non locally connected subsets of $\mathbb{R}^{2}$ are homeomorphic to subspaces of the other two.

Is there a surjective continuous map from one of these compact subsets to any of the others?
6.-Show that the solenoid $S_{m}$ defined in 4.2.3 can be embedded in $\mathbb{R}^{3}$, but not in $\mathbb{R}^{2}$.

### 4.3 Coverings

### 4.3.1

Definition Let B be a topological space. A covering of B is a fibre bundle with discrete fibres over B.

In other words, a covering of B is a space X together with a map $\pi: \mathrm{X} \rightarrow \mathrm{B}$, satisfying the following condition:

For all $b \in \mathrm{~B}$, there is a neighbourhood U of $b$ in B , a discrete space F and a homeomorphism $\varphi$ from $\pi^{-1}(\mathrm{U})$ onto $\mathrm{U} \times \mathrm{F}$ such that the diagram

commutates.
If $B$ is locally connected, then the condition can be restated as follows:
There is an open connected cover $\left(\mathrm{U}_{i}\right)_{i \in \mathrm{I}}$ of B such that, for all $i$ and all connected components V of $\pi^{-1}\left(\mathrm{U}_{i}\right)$, the projection $\pi$ induces a homeomorphism from V onto U .

A covering X of B is trivial if it is isomorphic to a covering $\mathrm{B} \times \mathrm{F}$, where F is a discrete space; an isomorphism $\varphi: \mathrm{X} \rightarrow \mathrm{B} \times \mathrm{F}$ is called a trivialization of X . Let $f: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ be a continuous map and X a covering of B . Then $\left(\mathrm{B}^{\prime}, f\right)$ trivializes X if the covering $\mathrm{X}^{\prime}=f^{*}(\mathrm{X})$ of $\mathrm{B}^{\prime}$ is trivial. In particular, if A is a subspace of $\mathrm{B}, \mathrm{A}$ is said to trivialize X if $\left.\mathrm{X}\right|_{\mathrm{A}}$ is a trivial covering of A .

The coverings of B form a full subcategory $\mathscr{C} \bullet \boldsymbol{v}_{\mathrm{B}}$ of the category of spaces over B.

### 4.3.2

A covering $X$ of $B$ is etale and Hausdorff over $B$. Indeed these conditions are local in B and if F is discrete, then $\mathrm{U} \times \mathrm{F}$ is etale and Hausdorff over U .

Therefore, if $s$ and $s^{\prime}$ are two continuous sections of X , the set of $b \in \mathrm{~B}$ such that $s(b)=s^{\prime}(b)$ is clopen ${ }^{1}$ in B.

The disjoint union and the fibre product of two coverings of $B$ are coverings of $B$.
If X is a covering and $f: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ a continuous map, then the space $f^{*}(\mathrm{X})$ obtained from X by a change of basis from B to $\mathrm{B}^{\prime}$ along $f$ is a covering of $\mathrm{B}^{\prime}$.

Let X be a covering of B and $b \in \mathrm{~B}$. The degree of X in $b$, written $\operatorname{deg}_{b}(\mathrm{X})$, is the cardinal of the fibre $\mathrm{X}(b) .^{2}$ The function $b \mapsto \operatorname{deg}_{b}(\mathrm{X})$ is locally constant: indeed, if $\mathrm{X}_{\mathrm{U}} \approx \mathrm{U} \times \mathrm{F}$, then $\operatorname{deg}_{b}(\mathrm{X})=$ Card F for all $b \in \mathrm{U}$. The image of X in B is clopen.

If $\mathrm{B} \neq \varnothing$ and if $\operatorname{deg}_{b}(\mathrm{X})$ does not depend on $b$ (for example if B is connected), then $\operatorname{deg}(\mathrm{X})$ or $\operatorname{deg}_{\mathrm{B}}(\mathrm{X})$ denote its value.

The covering X is finite of B if, for all $b \in \mathrm{~B}, \operatorname{deg}_{b}(\mathrm{X})$ is finite.

### 4.3.3

Proposition Suppose that B is connected and the space X etale and Hausdorff over B. X is a trivial covering of B if and only if, for all $x \in \mathrm{X}$, there is a continuous section $s: \mathrm{B} \rightarrow \mathrm{X}$ such that $s(\pi(x))=x$.
Proof Suppose that such continuous sections exist. Let $\Gamma$ be the set of sections $\mathrm{B} \rightarrow \mathrm{X}$, equipped with the discrete topology, and define the map $\varepsilon: \mathrm{B} \times \Gamma \rightarrow \mathrm{X}$ by $\varepsilon(b, s)=s(b)$. The map $\varepsilon$ is obviously continuous, it is etale since $\mathrm{B} \times \Gamma$ and X are etale over B , it is surjective by assumption. It is injective. Indeed, if $s(b)=s^{\prime}\left(b^{\prime}\right)$, then $b=b^{\prime}=\pi(s(b))$, and the set of points where $s$ and $s^{\prime}$ agree is closed since X is Hausdorff over $B$, closed since $\pi$ is locally injective, nonempty since it contains $b$, and so is equal to B since B is connected. Hence $\varepsilon$, being etale and bijective, is a homeomorphism. The converse is obvious.

### 4.3.4 Examples

(1) Let $X$ be clopen in a space $B$. Then $X$, equipped with inclusion, is a covering of B.
(2) Let $S^{1}=\{z \in \mathbb{C},|z|=1\}$. Then $\mathbb{R}$ together with the map $\varphi: t \mapsto e^{i t}$ is a covering of $S^{1}$. Indeed $\varphi$ induces a homeomorphism from $\mathbb{R} / 2 \pi \mathbb{Z}$ to $S^{1}$, and $\mathbb{R}$ equipped with the quotient map is a covering of $\mathbb{R} / 2 \pi \mathbb{Z}$ as will be seen in (4.3.14).

[^16]Similarly, $\mathbb{C}$ equipped with $z \mapsto e^{z}$ is a covering of $\mathbb{C}^{*}$.
(3) The space $\mathrm{S}^{1}$ together with $z \mapsto z^{d}$ is a $d$-fold covering of $\mathrm{S}^{1}$. Similarly, $\mathbb{C}^{*}$ together with $z \mapsto z^{d}$ is a covering of $\mathbb{C}^{*}$.
(4) The action of $\mathbb{R}^{*}$ identifies the set $\mathbf{P}^{n} \mathbb{R}$ of 1-dimensional vector subspaces of $\mathbb{R}^{n+1}$ with the quotient of $\mathbb{R}^{n+1}-\{0\}$, while that of $\{+1,-1\}$ identifies it with the unit sphere $\mathbf{S}^{n}$ of $\mathbb{R}^{n+1}$. Equip $\mathbf{P}^{n} \mathbb{R}$ with the quotient topology (both quotient topologies agree). The space $\mathbf{S}^{n}$, equipped with the quotient map, is a 2-fold covering of $\mathbf{P}^{n} \mathbb{R}$.

This covering is non trivial for $n>1$, since $\mathbf{S}^{n}$ is connected.
(5) More generally, let $\mathbf{G}_{p}\left(\mathbb{R}^{n}\right)$ be the set of $p$-dimensional vector subspaces of $\mathbb{R}^{n}$ (Grassmannian), and $\widetilde{\mathbf{G}}_{p}\left(\mathbb{R}^{n}\right)$ the set of oriented $p$-dimensional vector subspaces. Equip $\mathbf{G}_{p}\left(\mathbb{R}^{n}\right)$ (resp. $\left.\widetilde{\mathbf{G}}_{p}\left(\mathbb{R}^{n}\right)\right)$ with the topology obtained by considering it as the quotient of the Stiefel manifold $\mathbf{V}_{p}\left(\mathbb{R}^{n}\right) \subset \mathbb{R}^{n p}$ consisting of the orthonormal families $\left(x_{1}, \ldots, x_{p}\right)$ in $\mathbb{R}^{n}$ with respect to the action of the group $\mathrm{O}_{p}\left(\right.$ resp. $\left.\mathrm{SO}_{p}\right)(4.1$, Exercise 11). Then $\widetilde{\mathbf{G}}_{p}\left(\mathbb{R}^{n}\right)$ is a 2 -fold covering of $\mathbf{G}_{p}\left(\mathbb{R}^{n}\right)$, non trivial if $1 \leqslant p \leqslant n-1$.
(6) Let B be a topological space and $b \mapsto \mathrm{P}_{b}$ a continuous map from B to the vector space of polynomials of $\mathbb{C}[\mathbb{Z}]$ of degree $\leqslant d$. Assume that, for all $b \in \mathrm{~B}$, the polynomial $\mathrm{P}_{b}$ admits $d$ distinct roots in $\mathbb{C}$. Then the subspace X of $\mathrm{B} \times \mathbb{C}$ consisting of pairs $(b, z)$ such that $\mathrm{P}_{b}(z)=0$, equipped with the projection $(b, z) \mapsto b$, is a $d$ fold covering of B .

Indeed, let E be the vector space of polynomials of degree $\leqslant d, \mathrm{H}$ the subspace of $\mathrm{E} \times \mathbb{C}$ consisting of pairs $(\mathrm{P}, z)$ such that $\mathrm{P}(z)=0$, and W the set of $\mathrm{P} \in \mathrm{E}$ having $d$ distinct roots. Changing bases, it becomes a matter of showing that $\left.\mathrm{H}\right|_{\mathrm{W}}$ is a $d$ fold covering of W . Let $\mathrm{P}_{0} \in \mathrm{~W}$; the roots of $\mathrm{P}_{0}$ are simple and, by the implicit function theorem, for each root $z_{i}$ of $\mathrm{P}_{0}$, there is neighbourhood $\mathrm{U}_{i}$ of $\mathrm{P}_{0}$ in E and a neighbourhood $\mathrm{V}_{i}$ of $z_{i}$ in $\mathbb{C}$ such that $\mathrm{G}_{i}=\mathrm{H} \cap\left(\mathrm{U}_{i} \times \mathrm{V}_{i}\right)$ is the graph of a continuous map from $\mathrm{U}_{i}$ to $\mathrm{V}_{i}$. The sets $\mathrm{V}_{i}$ may be assumed to be disjoint, and the sets $\mathrm{U}_{i}$ to be equal to a neighbourhood U of $\mathrm{P}_{0}$. Then $\left.\mathrm{H}\right|_{\mathrm{U}}$ contains the union of the disjoint sets $\mathrm{G}_{i}$. Since every $\mathrm{P} \in \mathrm{U}$ has at least $d$ roots, $\left.\mathrm{H}\right|_{\mathrm{U}}=\bigcup \mathrm{G}_{i}$, and so $\left.\mathrm{H}\right|_{\mathrm{U}}$ is a trivial $d$-fold covering of U . It follows that W is open in E and that $\left.\mathrm{H}\right|_{\mathrm{W}}$ is a $d$-fold covering of W .

### 4.3.5

Proposition Let B be a topological space. Let X and Y be coverings of $\mathrm{B}, f$ and $g$ B -morphisms from Y to X . The set of $y \in \mathrm{Y}$ such that $f(y)=g(y)$ is clopen in Y . In particular, if Y is connected and if there exists $y \in \mathrm{Y}$ such that $f(y)=g(y)$, then $f=g$.

Proof The set of points where $f$ and $g$ agree is closed since X is Hausdorff over B, and it is opens since the projection $\mathrm{X} \mapsto \mathrm{B}$ is locally injective.

Corollary With the same notation, if B is connected and there exists $b \in \mathrm{~B}$ such that $f_{b}=g_{b}$, then $f=g$.

Proof Each connected component of Y contains a point of $\mathrm{Y}(b)$.

### 4.3.6

Proposition Let B be a locally connected space, X and Y be coverings of B , and $f: \mathrm{Y} \rightarrow \mathrm{X}$ a morphism. Then Y together with $f$ is a covering of X .

Lemma Assume B is connected, and that X and Y are trivial coverings of B and $f: \mathrm{Y} \rightarrow \mathrm{X}$ a covering map. Then Y together with $f$ is a covering of X , trivial on each connected component of X .

The proposition follows from the lemma by covering B with open and connected sets trivializing both X and Y .

Proof of Lemma We may assume that $\mathrm{X}=\mathrm{B} \times \mathrm{F}$ and $\mathrm{Y}=\mathrm{B} \times \mathrm{G}$, with F and G are discrete. Then $f$ is of the form $(b, u) \mapsto(b, \alpha(b, u))$, where $\alpha: \mathrm{B} \times \mathrm{G} \rightarrow \mathrm{F}$ is a continuous map. Since B is connected and F discrete, $\alpha(b, u)$ does not depend on $b$, and so $\alpha(b, u)=\varphi(u)$, where $\varphi$ is a map from G to F . Then, for $t \in \mathrm{~F}$, $f^{-1}(\mathrm{~B} \times\{t\})=\mathrm{B} \times \varphi^{-1}(t)$, and this space is a trivial covering of $\mathrm{B} \times\{t\}$.

Corollary 4.1 Let $f: \mathrm{Y} \rightarrow \mathrm{X}$ be a covering map. Then $f(\mathrm{Y})$ is clopen in X . In particular, if X is connected and $\mathrm{Y} \neq \varnothing$, then $f$ is surjective.

Corollary 4.2 Let $f: \mathrm{Y} \rightarrow \mathrm{X}$ be a covering map. If B is connected and there exists $b \in \mathrm{~B}$ such that $f_{b}: \mathrm{Y}(b) \rightarrow \mathrm{X}(b)$ is bijective, then $f$ is an isomorphism.

### 4.3.7

Let B be a locally connected space.
Proposition Assume that B is connected; let X be a covering of B . If there exists $b \in \mathrm{~B}$ such that there is a continuous section passing through all $x \in \mathrm{X}(b)$, then X is trivial.

Proof Let $\Gamma$ be the set of continuous sections of X equipped with the discrete topology. The map $\varepsilon: \mathrm{B} \times \Gamma \rightarrow \mathrm{X}$ defined by $\varepsilon(b, s)=s(b)$ is a covering map and ( $\mathrm{B} \times \Gamma, \varepsilon$ ) is a covering of X (4.3.6). We show that $\varepsilon$ is a isomorphism: for this it suffices that, for every connected component Y of $\mathrm{X}, \operatorname{deg}_{\mathrm{Y}}(\mathrm{B} \times \Gamma, \varepsilon)=1$. Let Y be a connected component of X . It is a covering of B , and its projection $\Pi(\mathrm{Y})$ in B is clopen and so equal to B ; in particular $\mathrm{Y} \cap \mathrm{X}_{b} \neq \varnothing$ : for $y \in \mathrm{Y} \cap \mathrm{X}_{b}$, $\operatorname{deg}_{y}(\mathrm{~B} \times \Gamma, \varepsilon)=\operatorname{deg}_{\mathrm{Y}}(\mathrm{B} \times \Gamma, \varepsilon)>0$. If $s$ and $s^{\prime}$ are two continuous sections X passing through $y$, then $s=s^{\prime}$ by 4.3.5, and so $\operatorname{deg}_{\mathrm{Y}}(\mathrm{B} \times \Gamma, \varepsilon)=1$.

### 4.3.8 Subcoverings, Quotient Coverings

Proposition Let B be a locally connected space and X a covering of B .
(a) A subset Y of X is a covering of B if and only if Y is clopen in X .
(b) If R be an equivalence relation on X , then $\mathrm{X} / \mathrm{R}$ is a covering of B if and only if the graph of R is clopen in $\mathrm{X} \times_{\mathrm{B}} \mathrm{X}$.

## Proof

(a) Suppose that $\mathrm{Y} \subseteq \mathrm{X}$ is a covering of B . Let U be a connected open subset of B trivializing X and Y . Let W be a connected component of $\left.\mathrm{Y}\right|_{\mathrm{U}}$. Then W is contained in a connected component V of $\left.\mathrm{X}\right|_{\mathrm{U}}$. As $\left.\pi\right|_{\mathrm{W}}: \mathrm{W} \rightarrow \mathrm{U}$ and $\left.\pi\right|_{\mathrm{V}}: \mathrm{V} \rightarrow \mathrm{U}$ are bijective, $\mathrm{W}=\mathrm{V}$. In other words, $\left.\mathrm{Y}\right|_{\mathrm{U}}$ is the union of the connected components of $\left.X\right|_{U}$, and the same holds for $\left.(X-Y)\right|_{U}$. Hence $\left.Y\right|_{U}$ and $\left.(X-Y)\right|_{U}$ are open in $\left.X\right|_{U}$. Since B can be covered with connected open subsets trivializing X and $\mathrm{Y}, \mathrm{Y}$ and $(\mathrm{X}-\mathrm{Y})$ must be open.
Conversely, let Y be clopen in $\mathrm{X}, \mathrm{U}$ an open connected subset in B trivializing X , and $\varphi:\left.\mathrm{X}\right|_{\mathrm{U}} \rightarrow \mathrm{U} \times \mathrm{F}$ a trivialization. The sets $\varphi^{-1}(\mathrm{U} \times\{t\})$ for $t \in \mathrm{~F}$ are the connected components of $\left.\mathrm{X}\right|_{\mathrm{U}}$. The space $\left.\mathrm{Y}\right|_{\mathrm{U}}$ is clopen in $\left.\mathrm{X}\right|_{\mathrm{U}}$, and so $\left.Y\right|_{U}=\varphi^{-1}(U \times G)$, for some subset $G$ of $F$. Hence $\left.Y\right|_{U}$ is a trivial covering of U , and since all points of B admit a connected open neighbourhood trivializing X , the space Y is a covering of B .
(b) Suppose $R$ is as stated and that $X / R$ is a covering of $B$. Let $\chi: X \rightarrow X / R$ be the canonical morphism. The graph of $R$ is the set of points of $X \times_{B} X$, where $\chi \circ \operatorname{pr}_{1}$ and $\chi \circ \mathrm{pr}_{2}: \mathrm{X} \times_{\mathrm{B}} \mathrm{X} \rightarrow \mathrm{X} / \mathrm{R}$ agree. It is therefore clopen by Proposition 4.3.5. Conversely, let U be a connected open subset in B trivializing X . Then $\mathrm{X} /\left.\mathrm{R}\right|_{\mathrm{U}}=$ $\left(\left.X\right|_{U}\right) / R$, and the topologies agree since $\left.X\right|_{U}$ is a saturated set for $R$ open in $X$. Let $\varphi:\left.\mathrm{X}\right|_{\mathrm{U}} \rightarrow \mathrm{U} \times \mathrm{F}$ be a trivialization and let $\mathrm{R}^{\prime}$ be the equivalence relation on $\mathrm{U} \times \mathrm{F}$ obtained by transferring via $\varphi$ the relation induced by R onto $\left.\mathrm{X}\right|_{\mathrm{U}}$. The graph of $\mathrm{R}^{\prime}$ is clopen in

$$
(\mathrm{U} \times \mathrm{F}) \times_{\mathrm{U}}(\mathrm{U} \times \mathrm{F})=\mathrm{U} \times \mathrm{F} \times \mathrm{F},
$$

and so can be written as $U \times G$, where $G$ is a subset of $F \times F$. The set $G$ is the graph of an equivalence relation on F since, if $x \in \mathrm{U}$, then

$$
(x, t) \sim_{\mathrm{R}^{\prime}}\left(x, t^{\prime}\right) \Longleftrightarrow\left(t, t^{\prime}\right) \in \mathrm{G} .
$$

The canonical bijection $(\mathrm{U} \times \mathrm{F}) / \mathrm{R}^{\prime} \rightarrow \mathrm{U} \times(\mathrm{F} / \mathrm{G})$ is a homeomorphism since the canonical map $\mathrm{U} \times \mathrm{F} \rightarrow \mathrm{U} \times(\mathrm{F} / \mathrm{G})$ is open. Hence $\left.(\mathrm{X} / \mathrm{R})\right|_{\mathrm{U}}$ is isomorphic to $U \times(F / G)$, and so is a trivial covering of $U$. Since $B$ can be covered by a family of nonempty connected open sets trivializing $X$, the space $X / R$ is a covering of $B$.

### 4.3.9 Examples

For $f \in \mathbb{C}[z]$ considered a map $\mathbb{C} \rightarrow \mathbb{C}$, let $f^{k}$ be $f \circ \cdots \circ f$, where $f$ appears $k$ times. Then $z$ is said to be a periodic point with period $k$ if $f^{k}(z)=z$ and $z_{i}=$ $f^{i}(z) \neq z$ for $0<i<k$. In this case, the set $\left\{z_{0}, \ldots, z_{k-1}\right\}$ is a $k$-cycle of $f$.

Fix $d \geqslant 2$ and $k \geqslant 1$. Let E denote the vector space of polynomials of degree $\leqslant d$, B the set of $f \in \mathrm{E}$ for which the equation $f^{k}(z)-z=0$ has $d^{k}$ distinct roots and set $\mathrm{X}=\left\{(f, z) \in \mathrm{B} \times \mathbb{C} \mid f^{k}(z)=z\right\}$. In other words, X is the set of $(f, z) \in \mathrm{B} \times \mathbb{C}$ such that $z$ is periodic point of $f$ with period dividing $k$.

The set B is open and dense in E since it is the set of $f$ of degree $d$ for which the discriminant of $f^{k}(z)-z$ is nonzero. Equipped with the projection $\pi:(f, z) \mapsto f$, the space X is a covering of B (4.3.4, Example 5).

The set Y of $(f, z)$, for $z$ with period $k$, is clopen in X . Indeed, it is clearly open; its complement is the set of $(f, z)$, where $z$ is periodic with period dividing some $k^{\prime}$ strictly dividing $k$. Hence it is the union of subcoverings, and so also open.

The set Z of $(f, \zeta)$, where $f \in \mathrm{~B}$ and $\zeta$ is a $k$-cycle of $f$ may be considered a covering of B . Indeed Z is the quotient of Y by the equivalence relation identifying $(f, z)$ and $\left(f, f^{i}(z)\right)$ for $0 \leqslant i \leqslant k-1$. The graph of this relation is the union of the images of the morphisms $(f, z) \mapsto\left(f, z, f^{i}(z)\right)$ from Y to $\mathrm{Y} \times_{\mathrm{B}} \mathrm{Y}$. Hence it is clopen and Proposition 4.3.8. can be applied.

### 4.3.10

Definition A space $X$ is simply connected if $X \neq \varnothing$ and every covering of $X$ is trivial. ${ }^{3}$

All simply connected spaces are connected. Indeed, if Y is clopen in X and different from $\varnothing$ and X , then Y is a non trivial covering of X .

The space X is called locally simply connected if all points of X admit a fundamental system of simply connected neighbourhoods.

Examples (1) All intervals of $\mathbb{R}$ are simply connected. For compact intervals this follows from (4.1.9); otherwise see 4.1, Exercise 12, or 4.3.17, Corollary 4.2.
(2) More generally, in a topological vector space, every nonempty convex subset is simply connected (4.3.17, Corollary 4.2).

[^17]
### 4.3.11

If $X$ is a covering of $B$ and $Y$ a covering of $X$, then the space $Y$, equipped with the composite function onto B, is not necessarily a covering of B (4.3, Exercise 7). This holds in most cases met with in practice. The next proposition provides a criterion for this to be so.

If Y is a space over B and $b \in \mathrm{~B}, \mathrm{Y}$ is said to be a $b$-neighbourhood covering if there is neighbourhood V of $b$ in B such that $\left.\mathrm{Y}\right|_{\mathrm{V}}$ is a covering of V .

Proposition Let B be a locally connected space, X a covering of B and Y a covering of X . Let $b \in \mathrm{~B}$ and U a connected neighbourhood of $b$ trivializing X . Let $\left(\mathrm{U}_{t}\right)_{t \in \mathrm{~F}}$ denote the family of connected components of $\left.\mathrm{X}\right|_{\mathrm{U}}$, and for open subsets V of U , set $\mathrm{V}_{t}=\left.\mathrm{U}_{t} \cap \mathrm{X}\right|_{\mathrm{V}}$. Then, Y is a b-neighbourhood covering of B if and only if there is an open neighbourhood V of $b$ in U such that, for all $t \in \mathrm{~F}$, the open subset $\mathrm{V}_{t}$ trivializes Y .

Proof Suppose that $\mathrm{V}_{t}$ trivializes Y for all $t \in \mathrm{~F}$. Then $\left.\mathrm{Y}\right|_{\mathrm{V}}$ is the disjoint union of the trivial coverings $f^{-1}\left(\mathrm{~V}_{t}\right)$ of V , and so is a trivial covering of V . The converse follows by (4.3.1).

Corollary 4.3 All coverings of a finite cover of B are coverings of B .
Corollary 4.4 If B is locally simply connected, all coverings of a covering of B are coverings of B .

### 4.3.12

Proposition A space X over B is a finite covering of B if and only if the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is etale, proper and Hausdorff.

Proof Suppose the projection is etale, proper and Hausdorff, and let $b \in \mathrm{~B}$. The fibre $\mathrm{X}(b)$ is discrete since the projection is etale; it is compact since it is is proper and Hausdorff, and hence finite. Let $\mathrm{X}(b)=\left\{x_{1}, \ldots, x_{d}\right\}$. There are open neighbourhoods $\mathrm{U}_{1}, \ldots, \mathrm{U}_{d}$ of $x_{1}, \ldots, x_{d}$ such that $\pi$ induces respectively a homeomorphism from $\mathrm{U}_{i}$ onto the open subsets $\mathrm{V}_{i}$ of B . Since $\pi$ is Hausdorff, the sets $\mathrm{U}_{i}$ may be assumed to be disjoint. Set

$$
\mathrm{V}=\bigcap \mathrm{V}_{i}-\pi\left(\mathrm{X}-\bigcup \mathrm{U}_{i}\right)
$$

Then $\pi^{-1}(\mathrm{~V})=\mathrm{U}_{1}^{\prime} \cup \cdots \cup \mathrm{U}_{d}^{\prime}$, where $\mathrm{U}_{i}^{\prime}=\mathrm{U}_{i} \cap \pi^{-1}(\mathrm{~V})$, and $\left.\mathrm{X}\right|_{\mathrm{V}}$ is isomorphic to $\mathrm{V} \times\{1, \ldots, d\}$. Conversely, "etale", "proper" and "Hausdorff" are local properties of B, and hence obviously hold for a finite covering.

### 4.3.13

Theorem Let X be a topological space and G a group acting on X such that, for all $g \in \mathrm{G}$, the map $x \mapsto g x$ is continuous. We assume the following:
(L) For all $x \in \mathrm{X}$, there is a neighbourhood U of $x$ such that, for all $g \neq e$ in G , $g \mathrm{U} \cap \mathrm{U}=\varnothing$.
Then X is a covering of $\mathrm{X} / \mathrm{G}$.
Proof Let U be open in X and $\forall g \neq e, \quad g \mathrm{U} \cap \mathrm{U}=\varnothing$. For $g \neq g^{\prime}$,

$$
g \mathbf{U} \cap g^{\prime} \mathbf{U}=\varnothing, \text { as } g^{\prime}=g h
$$

with $h \neq e$ and $g \mathrm{U} \cap g^{\prime} \mathrm{U}=g(\mathrm{U} \cap h \mathrm{U})=\varnothing$. So GU is a disjoint union of all $g \mathrm{U}$ for $g \in \mathrm{G}$. The canonical map $\chi: \mathrm{X} \rightarrow \mathrm{X} / \mathrm{G}$ is open since the saturation of all open subsets of $X$, i.e. the union of the orbits, is open. In particular $V=\chi(U)$ is open in $\mathrm{X} / \mathrm{G}$. The assumption $g \mathrm{U} \cap \mathrm{U}=\varnothing$ implies that $\chi$ induces a bijection from U onto V . This bijection is continuous and open, and hence a homeomorphism. For all $g \in \mathrm{G}$, the restriction of $\chi$ to $g \mathbf{U}$ is the composition of $g^{-1}: g \mathrm{U} \rightarrow \mathbf{U}$ and $\chi: \mathrm{U} \rightarrow \mathrm{V}$. So it is a homeomorphism from $g \mathrm{U}$ onto V . This gives a partition of $\chi^{-1}(\mathrm{~V})=\mathrm{GU}$ into open subsets homeomorphic to V under a projection.

Remark The assumptions of the theorem hold if X is a Hausdorff space and G a finite group acting continuously and freely on X, i.e.

$$
(\forall g \in \mathrm{G}), \quad x \mapsto g x \text { is continuous and }(\forall x \in \mathrm{X})(\forall g \neq e) \quad g \cdot x \neq x
$$

### 4.3.14

Corollary Let G be a topological group and H a discrete subgroup of G . Then G is a covering of $\mathrm{G} / \mathrm{H}$.

Proof There is a neighbourhood W of $e$ in G such that $\mathrm{W} \cap \mathrm{H}=\{e\}$, and a neighbourhood U of $e$ in G such that $\mathrm{U}^{-1} \cdot \mathrm{U} \subset \mathrm{W}$. Then $\mathrm{U} h \cap \mathrm{U}=\varnothing$ for $h \neq e$ in H , and for all $x \in \mathrm{G}, x \mathrm{U} h \cap x \mathrm{U}=\varnothing$. Hence the theorem can be applied to the right action of H on G .

### 4.3.15

Theorem Let B be a topological space, $\rho$ the map $(t, b) \mapsto(0, b)$ from $[0,1] \times \mathrm{B}$ to $[0,1] \times \mathrm{B}$, and X a covering of $[0,1] \times \mathrm{B}$. Then there is a unique covering isomorphism $\varphi: \mathrm{X} \rightarrow \rho^{*}(\mathrm{X})$ inducing the identity over $\{0\} \times \mathrm{B}$.

Proof Suppose first that X is trivial. We may then assume that $\mathrm{X}=[0,1] \times \mathrm{B} \times \mathrm{F}$, where F is discrete. So $\rho^{*}(\mathrm{X})=\mathrm{X}$, and $1_{\mathrm{X}}$ is the desired isomorphism. If $\varphi$ is another isomorphism with the required property, then $\varphi$ is of the form

$$
(t, b, x) \mapsto(t, b, u(t, b, x))
$$

where $u:[0,1] \times \mathrm{B} \times \mathrm{F} \rightarrow \mathrm{F}$ is continuous. For fixed $b$ and $x, t \mapsto u(t, b, x)$ is a constant continuous map from $[0,1]$ to F since $[0,1]$ is connected and F is discrete, so that $u(t, b, x)=u(0, b, x)=x$. Hence $\varphi=1_{\mathrm{X}}$.

Lemma 4.3 Let $\alpha, \beta, \gamma \in \mathbb{R}$ be such that $\alpha \leqslant \beta \leqslant \gamma$, and X a covering of $[\alpha, \gamma] \times \mathrm{B}$ such that $\mathrm{X}_{1}=\left.\mathrm{X}\right|_{[\alpha, \beta] \times \mathrm{B}}$ and $\mathrm{X}_{2}=\left.\mathrm{X}\right|_{[\beta, \gamma] \times \mathrm{B}}$ are trivial. Then X is trivial.

Proof of Lemma 4.3 Let $\tau_{1}: \mathrm{X}_{1} \rightarrow[\alpha, \beta] \times \mathrm{B} \times \mathrm{F}$ and $\tau_{2}: \mathrm{X}_{2} \rightarrow[\beta, \gamma] \times \mathrm{B} \times \mathrm{G}$ be trivializations. Let $\theta$ the homeomorphism from $\mathrm{B} \times \mathrm{F}$ onto $\mathrm{B} \times \mathrm{G}$ induced by $\tau_{2} \circ \tau_{1}^{-1}$ over $\{\beta\} \times \mathrm{B}$, and

$$
\theta_{*}=1_{[\beta, \gamma]} \times \theta:[\beta, \gamma] \times \mathrm{B} \times \mathrm{F} \rightarrow[\beta, \gamma] \times \mathrm{B} \times \mathrm{G}
$$

Then $\tau_{2}^{\prime}=\theta_{*}^{-1} \circ \tau_{2}: \mathrm{X}_{2} \rightarrow[\beta, \gamma] \times \mathrm{B} \times \mathrm{F}$ is a trivialization agreeing with $\tau_{1}$ over $\{\beta\} \times \mathrm{B}$. The trivializations $\tau_{1}$ and $\tau_{2}^{\prime}$ can be glued together to give a continuous bijection $\tau: \mathrm{X} \rightarrow[\alpha, \gamma] \times \mathrm{B} \times \mathrm{F}$, which is continuous. So is its inverse by (4.1.11, Lemma 4.4). Hence $\tau$ is a trivialization of X .

Lemma 4.4 Keeping the notation of the theorem, for all $b \in \mathrm{~B}$, there is a neighbourhood U of $b$ in B such that $\left.\mathrm{X}\right|_{[0,1] \times \mathrm{U}}$ is trivial.

Proof of Lemma 4.4 For $b \in \mathrm{~B},[0,1] \times\{b\}$ can be covered by a finite family of open subsets $\mathrm{J}_{i} \times \mathrm{U}_{i}$ of $[0,1] \times \mathrm{B}$, where $\mathrm{J}_{i}$ is open in $[0,1]$ and $\mathrm{U}_{i}$ a neighbourhood of $b$. Set $\mathrm{U}=\bigcap \mathrm{U}_{i}$. By Proposition 4.1.10, there is an increasing sequence $\left(t_{0}, \ldots, t_{n}\right)$ with

$$
t_{0}=0, t_{n}=1 \text { and }(\forall k \in\{1, \ldots, n\})(\exists i)\left[t_{k-1}, t_{k}\right] \subset \mathbf{J}_{i}
$$

Then $\mathrm{X}_{\left[t_{k-1}, t_{k}\right] \times \mathrm{U}}$ is trivial for all $k \in\{1, \ldots, n\}$. Using Lemma 4.3, it follows by induction that $\left.\mathrm{X}\right|_{\left[0, t_{k}\right] \times \mathrm{U}}$ is trivial; for $k=n$, this means that $\left.\mathrm{X}\right|_{[0,1] \times \mathrm{U}}$ is trivial.

Proof of Theorem (general case) By Lemma 4.4, there is an open cover $\left(\mathrm{U}_{i}\right)$ of B such that $\mathrm{X}_{[0,1] \times \mathrm{U}_{i}}$ is trivial for all $i$. Then, for all $i$, there is an isomorphism $\varphi_{i}$ : $\left.\mathrm{X}\right|_{[0,1] \times \mathrm{U}_{i}} \rightarrow \rho^{*}\left(\left.\mathrm{X}\right|_{[0,1] \times \mathrm{U}_{i}}\right)$ inducing the identity over $\{0\} \times \mathrm{U}_{i}$, and by uniqueness in the trivial case, $\varphi_{i}$ and $\varphi_{j}$ agree over $[0,1] \times \mathrm{U}_{i} \cap \mathrm{U}_{j}$. The maps $\varphi_{i}$ can be glued together to give an isomorphism $\varphi$ with the desired property. Its uniqueness follows from that of the maps $\varphi_{i}$.

### 4.3.16

Let X and Y be topological spaces. Two continuous maps $f_{0}$ and $f_{1}$ from X to Y are said to be homotopic if there is a continuous map $\mathrm{F}:[0,1] \times \mathrm{X} \rightarrow \mathrm{Y}$ such that $\mathrm{F}(0, x)=f_{0}(x)$ and $\mathrm{F}(1, x)=f_{1}(x)$ for all $x \in \mathrm{X}$. Then the relation $\simeq$ defined by $f_{0} \simeq f_{1}$ is an equivalence relation on the set $\mathrm{C}(\mathrm{X}, \mathrm{Y})$ of continuous maps from X to Y (for transitivity, use 4.1.11, Lemma 4.3). The topological spaces and homotopy classes of continuous maps form the quotient category $\mathscr{H}$ amatap of $\mathscr{T} a p$.

A continuous map $f: \mathrm{X} \rightarrow \mathrm{Y}$ is a called a homotopy equivalence if the homotopy class of $f$ is an isomorphism in $\mathscr{H}$ amotap, i.e. if there is a continuous map $g: \mathrm{Y} \rightarrow \mathrm{X}$ such that $g \circ f \simeq 1_{\mathrm{X}}$ and $f \circ g \simeq 1_{\mathrm{Y}}$. If there is a homotopy equivalence $\mathrm{X} \rightarrow \mathrm{Y} \mathrm{X}$ and Y are said to be homotopy equivalent, or to have the same homotopy type.

A space X is said to be contractible if it is homotopy equivalent to a point, i.e. if $X \neq \varnothing$ and $1_{X}$ is homotopic to a constant map.

Let X be a topological space and A a subset of X . A retraction from X onto A is a map $r: \mathrm{X} \rightarrow$ A such that $r(x)=x$ for $x \in \mathrm{~A}$. If there is a continuous retraction from X onto $\mathrm{A}, \mathrm{A}$ is said to be a retract of X (this requires A to be closed if X is Hausdorff).

A continuous retraction $r: \mathrm{X} \rightarrow \mathrm{A}$ is called a deformation retraction if $\iota \circ r$ : $\mathrm{X} \rightarrow \mathrm{X}$, where $\iota: \mathrm{A} \rightarrow \mathrm{X}$ is the canonical injection, is homotopic to $1_{\mathrm{X}}$. If there is a deformation retraction $\mathrm{X} \rightarrow \mathrm{A}, \mathrm{A}$ is a deformation retract of X . This implies that X and A are homotopy equivalent.

### 4.3.17

Theorem 4.3.15 admits the following corollaries:

Corollary 4.5 Let B and $\mathrm{B}^{\prime}$ be topological spaces, X a covering of $\mathrm{B}, f$ and $g$ continuous maps from $\mathrm{B}^{\prime}$ to B . If $f$ and $g$ are homotopic, then the coverings $f^{*}(\mathrm{X})$ and $g^{*}(\mathrm{X})$ of $\mathrm{B}^{\prime}$ are isomorphic.

Proof Let $h:[0,1] \times \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ be a homotopy from $f$ to $g, r$ the injection $b^{\prime} \mapsto\left(1, b^{\prime}\right)$ from $\mathrm{B}^{\prime}$ to $[0,1] \times \mathrm{B}^{\prime}$ and $\rho$ the map $\left(t, b^{\prime}\right) \mapsto\left(0, b^{\prime}\right)$ from $[0,1] \times \mathrm{B}^{\prime}$ to itself. Then

$$
f=h \circ \rho \circ r \text { and } g=h \circ r, \text { and } \rho^{*}\left(h^{*}(\mathrm{X})\right) \approx h^{*}(\mathrm{X})
$$

by the theorem; thus $f^{*}(\mathrm{X}) \approx g^{*}(\mathrm{X})$.
Corollary 4.6 Every contractible space is simply connected.

Remark In Theorem 4.3.15, assuming the paracompactness of B, the existence (but not uniqueness) of $\varphi$ generalizes to fibre spaces with arbitrary fibre. The proof is different. Assuming the paracompactness of $\mathrm{B}^{\prime}$, Corollary 4.5 generalizes likewise. All fibre spaces over a contractible paracompact space are trivial.

### 4.3.18 Coverings of Products

Proposition Let X and Y be topological spaces. Assume Y is connected and locally connected. Let E be a covering of $\mathrm{X} \times \mathrm{Y}$. If $\left.\mathrm{E}\right|_{\{x\} \times \mathrm{Y}}$ is trivial for all $x \in \mathrm{X}$ and if there exists $y_{0} \in \mathrm{Y}$ such that $\mathrm{E}_{\mathrm{X} \times\left\{y_{0}\right\}}$ is trivial, then E is trivial.

Lemma Let X and Y be topological spaces. Assume Y is connected and locally connected. Let E be a covering of $\mathrm{X} \times \mathrm{Y}$, and $\sigma$ a section of E , which is not assumed to be continuous. If $\left.\sigma\right|_{\{x\} \times \mathrm{Y}}$ is continuous for all $x \in \mathrm{X}$ and if there exists $y_{0} \in \mathrm{Y}$ such that $\left.\sigma\right|_{\mathrm{X} \times\left\{y_{0}\right\}}$ is continuous, then $\sigma$ is continuous.

Proof of Lemma Let $\tau:\left.\mathrm{E}\right|_{\mathrm{U} \times \mathrm{V}} \rightarrow \mathrm{U} \times \mathrm{V} \times \mathrm{F}$ be the trivialization with open $\mathrm{U} \subset$ X and connected, open $\mathrm{V} \subset \mathrm{Y}$. For $(x, y) \in \mathrm{U} \times \mathrm{V}$, we may write $\tau(\sigma(x, y))=$ $(x, y, s(x, y))$. This defines a function $s: \mathrm{U} \times \mathrm{V} \rightarrow \mathrm{F}$ continuous at $y$, and hence independent from $y$ since V is connected and F discrete. Hence, if $\sigma$ is continuous (resp. non continuous) at $(x, y) \in \mathrm{U} \times \mathrm{V}$, then it is continuous (resp. non continuous) at $\left(x, y^{\prime}\right)$ for all $y^{\prime} \in \mathrm{V}$. If $x \mapsto \sigma(x, y)$ is continuous on U for some $y \in \mathrm{~V}$, then $\sigma$ is continuous on $\mathrm{U} \times \mathrm{V}$.

It follows that, for each $x \in \mathrm{X}$, the set of $y$ such that $\sigma$ is continuous at $(x, y)$ is clopen in Y , and contains $y_{0}$. As Y is connected, this set is the whole of Y .

Proof of Proposition Let $\tau_{0}:\left.\mathrm{E}\right|_{\mathrm{X} \times\left\{y_{0}\right\}} \rightarrow \mathrm{X} \times\left\{y_{0}\right\} \times \mathrm{F}$ be a trivialization. For $t \in \mathrm{~F}$ and $x \in \mathrm{X}$, there is a unique continuous section $\sigma_{t, x}$ of $\left.\mathrm{E}\right|_{\{x\} \times \mathrm{Y}}$ passing through $\tau_{0}^{-1}\left(x, y_{0}, t\right)$. For each $t$, the section $\sigma_{t}:(x, y) \mapsto \sigma_{t, x}(y)$ satisfies the assumptions of the lemma, and so is continuous. The map $\alpha:(x, y, t) \mapsto \sigma_{t}(x, y)$ is a covering map.

The morphism $\alpha$ is bijective. Indeed, for $z \in \mathrm{E}(x, y)$, there is a unique section $\sigma_{z}$ of $\left.\mathrm{E}\right|_{\{x\} \times \mathrm{Y}}$ passing through $z$. Then $\tau\left(\sigma_{z}\left(x, y_{0}\right)\right)$ is of the form $\left(x, y_{0}, t_{z}\right)$ and $\alpha^{-1}(z)$ is reduced to $\left(x, y, t_{z}\right)$.

Hence, $\alpha$ is a covering isomorphism $\mathrm{X} \times \mathrm{Y} \times \mathrm{F} \rightarrow \mathrm{E}$.

### 4.3.19 Product of Simply Connected Spaces

Corollary Let X and Y be two simply connected spaces, at least one of which is locally connected. Then $\mathrm{X} \times \mathrm{Y}$ is simply connected.

## Exercises 4.3. (Coverings)

1.-Show that through a change of basis the covering $\widetilde{\mathrm{G}}_{p}\left(\mathbb{R}^{n}\right) \rightarrow \mathrm{G}_{p}\left(\mathbb{R}^{n}\right)$ (4.3.4, Example 5) is induced by the covering of $S^{N} \rightarrow P^{N}(\mathbb{R})$ for sufficiently large $N$.
2.-Let X and Y be topological spaces, and $f: \mathrm{Y} \rightarrow \mathrm{X}$ a proper map. Consider the transition functor $f^{*}$ from the category $\mathscr{C} \bullet \boldsymbol{v}_{\mathrm{X}}$ of coverings of X to $\mathscr{C} \boldsymbol{\imath} \boldsymbol{v}_{\mathrm{Y}}$.
(a) Show that, if $f$ is surjective with connected fibres, then the functor $f^{*}$ is fully faithful.
(b) Deduce that, if $f$ is surjective with connected fibres and Y simply connected, X is simply connected.
(c) Show that, if the fibres of $f$ are simply connected, then the functor $f^{*}$ is an equivalence of categories. Deduce that, with this assumption, if X is simply connected, then so is Y.
3.-Let $X$ be a topological space, $A$ and $B$ two open (or two closed) subsets of $X$ such that $\mathrm{A} \cup \mathrm{B}=\mathrm{X}$ and $\mathrm{A} \cap \mathrm{B} \neq 0$ is connected.
(a) Show that any covering E of X such that $\left.\mathrm{E}\right|_{\mathrm{A}}$ and $\left.\mathrm{E}\right|_{\mathrm{B}}$ are trivial is trivial.
(b) Deduce that if A and B are simply connected, so is X.
(c) Show that the sphere $S^{n}$ is simply connected for $n \geqslant 2$.
4.-Show that the Stiefel manifolds $\mathrm{V}_{p}\left(\mathbb{R}^{n}\right)$ (4.1, Exercise 11) are simply connected for $p \leqslant n-2$. Deduce that so are the oriented Grassmannians $\widetilde{\mathrm{G}}_{p}\left(\mathbb{R}^{n}\right)$ (4.3, Exercise 1).
5.-Let $\Phi$ be a functor from the category of sets and bijections into itself, B a locally connected space and X a covering of B . Define the structure of a B-covering on the set $\Phi_{\mathrm{B}}(\mathrm{X})=\bigsqcup_{b \in \mathrm{~B}} \Phi(\mathrm{X}(b))$. Study the properties of this action.
6.-Set $\mathrm{C}_{n}$ to be the circle of diameter $\left[(0,0),\left(\frac{1}{n}, 0\right)\right]$ in $\mathbb{R}^{2}$, and

$$
\mathrm{B}=\bigcup_{n \in \mathbb{N}^{*}} \mathrm{C}_{n}
$$

(a) Show that B is compact.
(b) Set $\mathrm{B}_{+}=\mathrm{B} \cap\left(\mathbb{R} \times \mathbb{R}_{+}\right)$and $\mathrm{B}_{-}=\mathrm{B} \cap\left(\mathbb{R} \times \mathbb{R}_{-}\right)$. Show that $\mathrm{B}_{+}$and $\mathrm{B}_{-}$are simply connected (they are contractible).
(c) Find a covering $\mathrm{E}_{n}$ of B such that $\left.\mathrm{E}_{n}\right|_{\mathrm{C}_{n}}$ is not trivial.
(d) Show that the disjoint union E of all $\mathrm{E}_{n}$ is a covering of $\mathrm{B} \times \mathbb{N}^{*}$, which is itself a trivial covering of $B$, but is not a covering of $B$.

Show that nonetheless $\left.\mathrm{E}\right|_{\mathrm{B}_{+}}$and $\left.\mathrm{E}\right|_{\mathrm{B}_{-}}$are trivial coverings.
7.-Let B and $\mathrm{B}^{\prime}$ be topological spaces. Describe the maps $f: \mathrm{B}^{\prime} \rightarrow \mathrm{B}$ for which a space X over B is a covering of B when $f^{*} \mathrm{X}$ is a covering of $\mathrm{B}^{\prime}$.
8.-Consider the non locally connected space $\overline{\mathbb{N}}=\mathbb{N} \cup\{\infty\}$.
(a) Show that the subspace $\overline{\mathbb{N}} \times \mathbb{N}-\{(\infty, 0)\}$ of $\overline{\mathbb{N}} \times \mathbb{N}$ is also a trivial covering of $\overline{\mathbb{N}}$.
(b) Find a clopen subset in $\overline{\mathbb{N}} \times \mathbb{N}$ which is not a covering of $\overline{\mathbb{N}}$.
9.-Let B be a not necessarily locally connected topological space. Define $\mathscr{R} i g$ to be the category of what are called rigid coverings of B as follows. An object of $\mathscr{R i g}$ is a covering X of B together with sets $\mathrm{U}_{i}$ forming an open cover of B and with trivializations $\left.\tau_{i} \mathrm{X}\right|_{\mathrm{U}_{i}} \rightarrow \mathrm{U}_{i} \times \mathrm{F}_{i}, i \in \mathrm{I}$, satisfying the condition below. Over $\mathrm{U}_{i, j}=\mathrm{U}_{i} \cap \mathrm{U}_{j}$, the map $\tau_{j} \circ \tau_{i}^{-1}$ is of the form $(b, y) \mapsto\left(b, \gamma_{i, j}(b)(y)\right)$, where $\gamma_{i, j}$ is a map from $\mathrm{U}_{i, j}$ to the set $\mathrm{S}_{i, j}=\operatorname{Isom\mathscr {E}\boldsymbol {ns}}\left(\mathrm{F}_{i}, \mathrm{~F}_{j}\right)$. The required condition is:
(lc) For all $(i, j)$, the map $\gamma_{i, j}: \mathrm{U}_{i, j} \rightarrow \mathrm{~S}_{i, j}$ is locally constant.
(a) Note that this condition is automatically satisfied if B is locally connected. Give an example with B compact but not locally connected where it is is not satisfied.

Given two objects $\mathbf{X}=\left(\mathrm{X},\left(\mathrm{U}_{i}, \tau_{i}\right)_{i \in \mathrm{I}}\right)$ and $\mathbf{X}^{\prime}=\left(\mathrm{X}^{\prime},\left(\mathrm{U}_{j}, \tau_{j}\right)_{j \in \mathrm{I}^{\prime}}\right)$ of $\mathscr{R} i g$, define $\operatorname{Hom}_{\mathscr{R i} \mathfrak{g}}\left(\mathbf{X}, \mathbf{X}^{\prime}\right)$ as the set of covering morphisms $h: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ such that, for all $(i, j) \in \mathrm{I} \times \mathrm{I}^{\prime}$ the map $\eta_{i, j}: \mathrm{U}_{i} \cap \mathrm{U}_{j}^{\prime} \rightarrow \operatorname{Isom}_{\mathscr{E} n s}\left(\mathrm{~F}_{i}, \mathrm{~F}_{j}^{\prime}\right)$ expressing $h$ is locally constant.
(b) Show that, if B is locally connected, then the forgetful functor $\mathscr{R} i g \rightarrow \mathscr{C} \not \boldsymbol{v}_{B}$ is an equivalence of categories. Give examples where it is not fully faithfully, or essentially surjective.
10.-Let $\Omega$ be the set of $\left(z_{1}, \ldots, z_{d}\right) \in \mathbb{C}^{d}$, with distinct $z_{i}$.
(a) Show that $\Omega$ is connected.
(b) Denote by $\mathscr{P}_{d}$ the set of monic polynomials $z^{d}+a_{d-1} z^{d-1}+\cdots+a_{0}$ of degree $d$, and $\mathscr{B}$ the set of $\mathrm{P} \in \mathscr{P}_{d}$ with $d$ distinct roots. Define $\pi: \Omega \rightarrow \mathscr{B}$ by $\pi\left(z_{1}, \ldots, z_{d}\right)=\prod\left(z-z_{i}\right)$.

Show that $\Omega$ together with $\pi$ is a $d!$-fold covering of $\mathscr{B}$ and is not trivial.
(c) Denote by X the set of $(\mathrm{P}, z)$ such that $\mathrm{P} \in \mathscr{B}$ and $\mathrm{P}(z)=0$. Define $\varpi: \mathrm{X} \rightarrow$ $\mathscr{B}$ by $\varpi(\mathrm{P}, z)=\mathrm{P}$ and $\omega: \Omega \rightarrow \mathrm{X}$ by $\omega\left(z_{1}, \ldots, z_{d}\right)=\left(\pi\left(z_{1}, \ldots, z_{d}\right), z_{1}\right)$. Show that $(\Omega, \omega)$ is a non trivial covering of X and that $(\mathrm{X}, \varpi)$ is a non trivial covering of $\mathscr{B}$.
11. (A proof of d'Alembert's theorem.) -Let $f \in \mathbb{C}[z]$ be a polynomial of degree $d$ and F a primitive of $f$ (i.e. a polynomial such that $\mathrm{F}^{\prime}=f$ ).
(a) Show that $\mathrm{F}: \mathbb{C} \rightarrow \mathbb{C}$ is a proper map.
(b) Assume that $f$ does not vanish at any point. Show that F is a $d+1$-fold covering. Given that $\mathbb{C}$ is simply connected, show that $d=0$.

### 4.4 Universal Coverings

### 4.4.1

Let ( $\mathrm{B}, b_{0}$ ) be a pointed space, i.e. a topological B with a distinguished point $b_{0} \in \mathrm{~B}$ called a basepoint. A pointed covering of $\left(\mathrm{B}, b_{0}\right)$ is a covering E of B with a point $t_{0} \in \mathrm{E}\left(b_{0}\right)$.

Definition Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space. If the functor $\mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$ from the category of coverings of B to the category of sets is representable, then a pair ( $\mathrm{E}, t_{0}$ ) representing this functor is called a pointed universal covering of ( $\mathrm{B}, b_{0}$ ).

In other words, $\left(\mathrm{E}, t_{0}\right)$ is a pointed universal covering of $\left(\mathrm{B}, b_{0}\right)$ if and only if for all coverings X of B and all $x \in \mathrm{X}\left(b_{0}\right)$, there is a unique B -morphism $f: \mathrm{E} \rightarrow \mathrm{X}$ such that $f\left(t_{0}\right)=x$.

A universal pointed covering of ( $\mathrm{B}, b_{0}$ ) is an initial object (2, Exercise 2 ) in the category of pointed coverings of ( $\mathrm{X}, x_{0}$ ).

Two universal pointed coverings of ( $\mathrm{B}, b_{0}$ ) are uniquely isomorphic.

### 4.4.2

Let ( $\mathrm{B}, b_{0}$ ) be a pointed locally connected space.
Proposition A pointed covering ( $\mathrm{E}, t_{0}$ ) of $\left(\mathrm{B}, b_{0}\right)$ is universal if and only if E is connected and trivializes all coverings of B .

## Proof

(a) If $\left(\mathrm{E}, t_{0}\right)$ is universal, then E is connected.

Let F be clopen in E and $t_{0} \in \mathrm{~F}$. Define $f$ and $g$ from E to $\mathrm{B} \times\{0,1\}$ by $f(t)=$ $(\pi(t), 0)$ for all $t \in \mathrm{E}$ and $g(t)=(\pi(t), 1)$ for $t \notin \mathrm{~F}$ and $g(t)=(\pi(t), 0)$ for $t \in \mathrm{~F}$. Then $f\left(t_{0}\right)=g\left(t_{0}\right)$; so $f=g$ by uniqueness of the universal property; hence $\mathrm{F}=\mathrm{E}$. Therefore E is connected.
(b) If $\left(\mathrm{E}, t_{0}\right)$ is universal, then E trivializes every covering of B .

Let $X$ be a covering of $B$. We show that $E \times_{B} X$ is trivial. For this it suffices (4.3.7) to show that for $x \in \mathrm{X}\left(b_{0}\right)$, there is a continuous section $\mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ passing through $\left(t_{0}, x\right)$. By the universal property of E , there is a B-morphism $f: \mathrm{E} \rightarrow \mathrm{X}$ such that $f\left(t_{0}\right)=x$. Then $t \mapsto(t, f(t))$ is the desired section.
(c) Converse.

Let X be a covering of B and $x \in \mathrm{X}\left(b_{0}\right)$. There is a continuous section $\mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ passing through $\left(t_{0}, x\right)$ since E trivializes X ; it is unique since E is connected. A unique B -morphism $f: \mathrm{E} \rightarrow \mathrm{X}$ such that $f\left(t_{0}\right)=x$ corresponds to it.

Remark The assumption that B is locally connected was only needed to prove "necessity". "Sufficiency" holds without this assumption.

### 4.4.3

Let ( $\mathrm{B}, b_{0}$ ) be a pointed locally connected space.
Corollary Let $\left(\mathrm{E}, t_{0}\right)$ be a universal pointed covering of $\left(\mathrm{B}, b_{0}\right), b \in \mathrm{~B}$ and $t \in \mathrm{E}(b)$. Then ( $\mathrm{E}, t$ ) is a universal pointed covering of $(\mathrm{B}, b)$.

Indeed, the characterization given in Proposition 4.4.2 does not involve basepoints. (Not that the assumption $t \in \mathrm{E}(b)$ only implies that $b$ and $b_{0}$ belong to the same connected component.)

A covering E of B is said to be a universal covering of B if there are $b \in \mathrm{~B}$ and $t \in \mathrm{E}(b)$ such that $(\mathrm{E}, t)$ is a universal pointed covering of $(\mathrm{B}, b)$. This holds for all $(b, t)$ with $t \in \mathrm{E}(b)$.

If E is a universal covering of B , the image $\mathrm{B}^{\prime}$ of E in B is a connected component of B , and E a universal covering of $\mathrm{B}^{\prime}$.

If B is connected, then two universal coverings of B are isomorphic, but not uniquely. Worse, there is no natural way of choosing an isomorphism, and so the universal coverings of B cannot be mutually identified, nor is it possible to speak of "the universal covering of B".

### 4.4.4

Proposition Every simply connected covering is universal.
Proof Let E be a simply connected covering of B. Then E is connected and every covering of E is trivial, and so E trivializes every covering of B , and E is universal by 4.4.2.

Remark Let B be a connected space; if B admits a simply connected covering, then every universal covering of $B$ is simply connected. There are spaces $B$ with universal coverings but without any simply connected one (4.4, Exercise 2).

### 4.4.5

Theorem Every connected, simply connected space has a simply connected covering.

The proof requires the notion of product covering of an infinite family of coverings. For an example of locally connected space without any universal covering, see 4.4, Exercise 3.

### 4.4.6 Product Covering

Let B be a totally connected space and $\left(X_{\lambda}\right)_{\lambda \in \Lambda}$ a family of coverings of $B$. If $\Lambda$ is finite, then the fibre product $X$ of $X_{\lambda}$ is a covering of $B$. But this is not always the case when $\Lambda$ is infinite; then X may not in general be locally connected nor discrete since its fibres are infinite products of discrete spaces.

Proposition If there is an open cover of B trivializing all $\mathrm{X}_{\lambda}$, then the set X can be equipped with a unique topology $\mathscr{T}$, finer than the product topology, turning X into a cover of B . The cover thus obtained is a product of $\mathrm{X}_{\lambda}$ in the category of coverings of B .

## Proof

(a) Suppose first that B is connected and all $\mathrm{X}_{\lambda}$ are trivial. For all $\lambda \in \Lambda$, assume $\mathrm{X}_{\lambda}=\mathrm{B} \times \mathrm{F}_{\lambda}$. Set $\mathrm{F}=\prod_{\lambda} \mathrm{F}_{\lambda}$. Let $\mathscr{T}_{0}$ be the product topology on $\mathrm{B} \times \prod_{\lambda} \mathrm{F}_{\lambda}$ (which is not generally locally connected), and $\mathscr{T}$ the product topology of the topology on B and of the discrete topology on F . The topology $\mathscr{T}$ is finer than $\mathscr{T}_{0}$ and turns X into a covering of B.
We show that $(\mathrm{X}, \mathscr{T})$ is a product of $\mathrm{X}_{\lambda}$ in $\mathscr{C} \bullet \boldsymbol{v}_{\mathrm{B}}$. For this, it suffices to show that, if Y is a covering of B and for all $\lambda, f_{\lambda}$ is a morphism from Y to $\mathrm{X}_{\lambda}$, then the map $f$ from Y to X defined by $f_{\lambda}$ is continuous from Y to $(\mathrm{X}, \mathscr{T})$. If necessary by restricting to an open connected subset of B trivializing Y , we may assume that $\mathrm{Y}=\mathrm{B} \times \mathrm{G}$, where G is discrete. Then, B being connected and $\mathrm{F}_{\lambda}$ discrete, $f_{\lambda}:(b, g) \mapsto\left(b, h_{\lambda}(g)\right)$ for some map $h_{\lambda}$ from G to $\mathrm{F}_{\lambda}$. The map $h: \mathrm{G} \rightarrow \mathrm{F}$ defined by $h_{\lambda}$ is continuous since G is discrete. Hence $f:(b, g) \mapsto(b, h(g))$ is continuous, which shows that X is the product of the $\mathrm{X}_{\lambda}$ in $\mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}}$.
Let $\mathscr{T}^{\prime}$ be another topology finer than $\mathscr{T}_{0}$ and turning X into a covering of B. The projection $\left(\mathrm{X}, \mathscr{T}^{\prime}\right) \rightarrow \mathrm{X}_{\lambda}$ is continuous for all $\lambda$, and so the identity $\left(\mathrm{X}, \mathscr{T}^{\prime}\right) \rightarrow(\mathrm{X}, \mathscr{T})$ is continuous. It is a bijective covering morphism, hence an isomorphism, and $\mathscr{T}^{\prime}=\mathscr{T}$. This proves the proposition in case (a).
(b) In the general case, let ( $\mathrm{U}_{i}$ ) be open connected sets covering B and trivializing all $\mathrm{X}_{\lambda}$. For all $i$, let $\mathscr{T}_{i}$ be the topology on the product covering $\mathrm{X}_{i}=\left.\mathrm{X}\right|_{\mathrm{U}_{i}}$ of $\left.\mathrm{X}_{\lambda}\right|_{\mathrm{U}_{i}}$. By uniqueness in (a), for all $(i, j)$ the topologies $\mathscr{T}_{i}$ and $\mathscr{T}_{j}$ agree over each connected component of $\mathrm{U}_{i, j}$. The following lemma shows that $\mathscr{T}_{i}$ and $\mathscr{T}_{j}$ agree over $\mathrm{U}_{i, j}$, then that the $\mathscr{T}_{i}$ can be glued together to give the desired topology $\mathscr{T}$.

Gluing Lemma Let B be a topological space and X a set with a map $\pi: \mathrm{X} \rightarrow \mathrm{B}$. Let $\left(\mathrm{U}_{i}\right)$ be an open cover of B . Set $\mathrm{X}_{i}=\pi^{-1}\left(\mathrm{U}_{i}\right)$ and for each $i$ let $\mathscr{T}_{i}$ be a topology on X making $\pi_{i}=\left.\pi\right|_{\mathrm{X}_{i}}$ continuous. Assume that $\mathscr{T}_{i}$ and $\mathscr{T}_{j}$ induce the same topology on
$\mathrm{X}_{i} \cap \mathrm{X}_{j}$ for all $i$ and $j$. Then there is a unique topology $\mathscr{T}$ on X making $\pi$ continuous and inducing $\mathscr{T}_{i}$ on $\mathrm{X}_{i}$ for all $i$.

Proof Let $\mathscr{T}$ be the topology on X whose open sets are the subsets V of X such that for all $i, \mathrm{~V} \cap \mathrm{X}_{i}$ is open with respect to $\mathscr{T}_{i}$. Then $\pi$ is continuous with respect to $\mathscr{T}$. We show that $\left.\mathscr{T}\right|_{\mathrm{U}_{i}}=\mathscr{T}_{i}$ for all $i$. It readily follows that $\mathscr{T}_{i}$ is finer than $\left.\mathscr{T}\right|_{\mathrm{U}_{i}}$. Let W be an open subset in $\mathrm{X}_{i}$ for $\mathscr{T}_{i}$. For all $j, \mathrm{~W} \cap \mathrm{X}_{j}$ is open in $\mathrm{X}_{i j}=\mathrm{X}_{i} \cap \mathrm{X}_{j}$ with respect to $\mathscr{T}_{i}$, hence with respect to $\mathscr{T}_{j}$. Since $\pi_{j}$ is continuous with respect to $\mathscr{T}_{j}, \mathrm{X}_{i j}$ is open in $\mathrm{X}_{j}$ with respect to $\mathscr{T}_{j}$, and so $\mathrm{W} \cap \mathrm{X}_{j}$ is open with respect to $\mathscr{T}_{j}$. Therefore W is open with respect to $\mathscr{T}$, and $\mathscr{T}_{i}=\left.\mathscr{T}\right|_{\mathrm{U}_{i}}$. If $\mathscr{T}^{\prime}$ is any other topology making $\pi$ continuous and inducing the $\mathscr{T}_{i}$, it readily follows that $\mathscr{T}^{\prime}$ is coarser than $\mathscr{T}$. Let V be an open subset of W with respect to $\mathscr{T}$. Then $\mathrm{V} \cap \mathrm{X}_{i}$ is open in $\mathrm{X}_{i}$ with respect to $\mathscr{T}^{\prime}$, and so $\mathrm{V} \cap \mathrm{X}_{i}$ is open in X with respect to $\mathscr{T}^{\prime}$ and $\mathrm{V}=\bigcup\left(\mathrm{V} \cap \mathrm{X}_{i}\right)$ is open with respect to $\mathscr{T}^{\prime}$. Hence $\mathscr{T}=\mathscr{T}^{\prime}$.

### 4.4.7

Cardinality Lemma Let B be a topological space, $\Omega$ an open basis for B having infinite cardinality $\mathbb{W}$, and X a connected etale Hausdorff space over B . Then $(\forall b \in \mathrm{~B})$ $\operatorname{Card} \mathbf{X}(b) \leqslant \boldsymbol{w}$.

Proof The open sets in $\Omega$ may be assumed to be connected without changing $ש$. A finite sequence $\left(\mathrm{U}_{0}, \ldots, \mathrm{U}_{n}\right)$ of open connected sets in $\Omega$ such that $\mathrm{U}_{i-1} \cap \mathrm{U}_{i} \neq \varnothing$ for $i=1, \ldots, n$ will be called a chain in B. If $\left(\mathrm{U}_{0}, \ldots, \mathrm{U}_{n}\right)$ is a chain in B, a finite sequence $\left(\mathrm{V}_{0}, \ldots, \mathrm{~V}_{n}\right)$ of open sets in X will be called a chain in X over $\left(\mathrm{U}_{0}, \ldots, \mathrm{U}_{n}\right)$ if $\pi$ induces a homeomorphism from $\mathrm{V}_{i}$ onto $\mathrm{U}_{i}$ for $i \in\{0, \ldots, n\}$ and $\pi\left(\mathrm{V}_{i-1} \cap \mathrm{~V}_{i}\right)=\mathrm{U}_{i-1} \cap \mathrm{U}_{i}$ for $i \in\{1, \ldots, n\}$. If $\left(\mathrm{V}_{0}, \ldots, \mathrm{~V}_{n}\right)$ and $\left(\mathrm{V}_{0}^{\prime}, \ldots, \mathrm{V}_{n}^{\prime}\right)$ are two chains over the same chain $\left(\mathrm{U}_{0}, \ldots, \mathrm{U}_{n}\right)$ of B , and if $\mathrm{V}_{0} \cap \mathrm{~V}_{0}^{\prime} \neq \varnothing$, then these two chains are equal. Indeed, writing $s_{i}$ for the section $\mathrm{U}_{i} \rightarrow \mathrm{~V}_{i}$ and $s_{i}^{\prime}$ for the section $\mathrm{U}_{i} \rightarrow \mathrm{~V}_{i}^{\prime}$, induction shows that $s_{i}$ and $s_{i}^{\prime}$ agree at least at one point of $\mathrm{U}_{i}$, and so are equal.

If there is a chain $\mathscr{U}$ in B such that $\mathscr{V}$ is a chain over $\mathscr{U}, \mathscr{V}$ will be said to be a chain in X. For $x, x^{\prime} \in \mathrm{X}$, let $\mathscr{C}_{x}$ (resp. $\left.\mathscr{C}_{x, x^{\prime}}\right)$ be the set of chains $\left(\mathrm{V}_{0}, \ldots, \mathrm{~V}_{n}\right)$ of X such that $x \in \mathrm{~V}_{0}$ (resp. $x \in \mathrm{~V}_{0}$ and $x^{\prime} \in \mathrm{V}_{n}$ ). The cardinality of the set of chains of B is $\leqslant \boldsymbol{ש}$, and over each chain of B there is at most one chain in $\mathscr{C}_{x}$, and so $\operatorname{Card} \mathscr{C}_{x} \leqslant \boldsymbol{ש}$.

The relation $\mathscr{C}_{x, x^{\prime}} \neq \varnothing$, where $x, x^{\prime} \in \mathrm{X}$ is an equivalence relation. Indeed, if

$$
\left(\mathrm{V}_{0}, \ldots, \mathrm{~V}_{n}\right) \in \mathscr{C}_{x, x^{\prime}} \text { and }\left(\mathrm{V}_{0}^{\prime}, \ldots, \mathrm{V}_{n}^{\prime}\right) \in \mathscr{C}_{x^{\prime}, x^{\prime \prime}}
$$

writing W for the connected component of $x^{\prime}$ in $\mathrm{V}_{n} \cap \mathrm{~V}_{0}^{\prime}$,

$$
\left(\mathrm{V}_{0}, \ldots, \mathrm{~V}_{n}, \mathrm{~W}, \mathrm{~V}_{0}^{\prime}, \ldots, \mathrm{V}_{n}^{\prime}\right) \in \mathscr{C}_{x^{\prime}, x^{\prime \prime}}
$$

reflexivity and symmetry are immediate. The equivalence classes are open. Hence, since X is connected, all points are equivalent. Let $b \in \mathrm{~B}$ and $x_{0} \in \mathrm{X}$. As the restriction of $\pi$ to the last open set of a chain in X is injective, for $x \in \mathrm{X}(b), \mathscr{C}_{x_{0}, x}$ form a disjoint family of nonempty subsets of $\boldsymbol{\mathscr { C }}_{x_{0}}$, and so $\operatorname{Card} \mathrm{X}(b) \leqslant \boldsymbol{ש}$.

Proof of Theorem 4.4.5 Let $(\mathrm{B}, b)$ be a pointed space. If ( $\mathrm{X}, x$ ) and (Y,y) are pointed connected coverings of $(\mathrm{B}, b),(\mathrm{Y}, y)$ will be said to dominate $(\mathrm{X}, x)$ if there is a pointed covering morphism from $(\mathrm{Y}, y)$ to $(\mathrm{X}, x)$. Such a morphism is unique by (4.3.5). Hence, if ( $\mathrm{X}, x$ ) and ( $\mathrm{Y}, y$ ) mutually dominate one other, then they are isomorphic.

If $(\mathrm{E}, t)$ is a pointed connected covering of $(\mathrm{B}, b)$ dominating all others, then it is a universal covering by Definition (4.4.1). If $B$ is moreover locally simply connected, then E is simply connected. Indeed, if $(\mathrm{Y}, y)$ is a pointed connected covering of $(\mathrm{E}, t)$, then Y is a covering of B by (4.3.11, Corollary 4.6), and so ( $\mathrm{E}, t$ ) dominates ( $\mathrm{Y}, y$ ) and the projection $\mathrm{Y} \rightarrow \mathrm{E}$ is an isomorphism. Hence every connected covering of E is trivial, and every sum of connected coverings is trivial.

Let B be a connected, locally simply connected space. Choose a point $b \in \mathrm{~B}$, an open basis $\Omega$ for B and a set $\Xi$ with cardinality $\geqslant \mathrm{Card} \mathrm{B} \times \operatorname{Card} \Omega$. By the cardinality lemma, every covering of $B$ has cardinality $\leqslant \operatorname{Card} \Xi$, and so is isomorphic to a covering whose underlying set is a subset of $\Xi$. Let $\left(\mathrm{X}_{\lambda}, x_{\lambda}\right)_{\lambda \in \Lambda}$ be a family of pointed coverings of $(\mathrm{B}, b)$ containing every connected pointed coverings of $(\mathrm{B}, b)$ whose underlying set is in $\Xi$. A simply connected open cover of $B$ trivializes all $\mathrm{X}_{\lambda}$. Let X be the product covering of $\mathrm{X}_{\lambda}$, set $x=\left(x_{\lambda}\right)$ and let E be the connected component of $x$ in X . The pointed covering ( $\mathrm{E}, x$ ) is connected and dominates all $\mathrm{X}_{\lambda}$, hence all connected pointed coverings of $(\mathrm{B}, b)$. Thus it is simply connected.

## Exercises 4.4. (Universal coverings)

1. (Non-functoriality of the universal covering.) -(a) Let E be the universal covering of $\mathbf{S}^{1}=\{z \in \mathbb{C}| | z \mid=1\}$, and $f$ the $z \mapsto-z$ from $S^{1}$ to itself. Show that there is a $f$-morphism $g: \mathrm{E} \rightarrow \mathrm{E}$, but none satisfying $g \circ g=1_{\mathrm{E}}$.
(b) Let $\mathscr{C}$ be the category of locally simply connected topological spaces. Show that there is no functor $\mathrm{R}: \mathscr{C} \rightarrow \mathscr{C}$ and nor any functorial morphism $\pi: \mathrm{R} \rightarrow 1_{\mathscr{C}}$ such that, for all $X \in \mathscr{C}$, the space $R(X)$ with $\pi_{X}: R(X) \rightarrow X$ is a universal covering of X.
(c) What do you think of the following argument: "Consider the Hopf fibre bundle $\mathrm{H}=\left(\mathrm{S}^{3}, \pi: \mathrm{S}^{3} \rightarrow \mathrm{~S}^{2}\right)$ with fibre $\mathrm{S}^{1}$ (4.1, Exercise 11). Replacing each fibre by its universal covering gives a fibre bundle with fibre $\mathbb{R}$. However, every fibre bundle with fibre $\mathbb{R}$ has a continuous section; the image of such a section is a continuous section of H , and so H has a continuous section"?
2.-Let $\mathrm{C}_{n}$ be the circle of diameter $\left[(0,0),\left(\frac{1}{2 n+1}, 0\right)\right]$ for $n \in \mathbb{Z}$ in $\mathbb{R}^{2}$, and set $\mathrm{M}=$ $\bigcup \mathrm{C}_{n}$. Let $f$ be a homeomorphism from M onto itself inducing a homeomorphism from $\mathrm{C}_{n}$ onto $\mathrm{C}_{n+1}$ for all $n$. Let X be the quotient space of $\mathbb{R} \times \mathrm{M}$ by the equivalence relation identifying $(t, m)$ to $(t+1, f(m))$ for $t \in \mathbb{R}$ and $m \in \mathrm{M}$.
(a) Show that $X$ is a fibre bundle on $S^{1}=\mathbb{R} / \mathbb{Z}$ with fibre M. Show that every covering of X is trivial on each fibre of this fibration.
(b) Show that X has a non simply connected universal covering.
3.-Show that the space $B$ of 4.3 , Exercise 6 , does not have a universal covering.
4.-(A) Let $B$ be a topological space. Consider the following properties:
(i) there is an open cover $\left(U_{i}\right)$ of B such that every cover of B induces a trivial covering of each $U_{i}$;
(ii) Every covering of a covering of B is a covering of B .

Do any one of these properties imply the other?
(B) Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed connected space.
(a) Let $\mathscr{U}=\left(\mathrm{U}_{i}\right)$ be a connected open cover of B , and $\mathscr{\mathscr { C }} \boldsymbol{\boldsymbol { v } _ { \mathrm { B } , \mathrm { U } }}$ the category of coverings of B that are trivial over each $\mathrm{U}_{i}$. Show that B has a universal covering in $\mathscr{C} \bullet \boldsymbol{v}_{B, U}$, i.e. that the functor $\mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$ from $\mathscr{C} \bullet \boldsymbol{v}_{\mathrm{B}, \mathrm{U}}$ to $\mathscr{E} n \boldsymbol{n}$ is representable.
(b) Show that B has a universal covering if and only if B satisfies condition (i) of part A, and that this covering is simply connected if B satisfies (ii).
5. (Another proof of d'Alembert's theorem.) -Let A be a commutative Banach algebra over $\mathbb{R}$ or $\mathbb{C}$, and let $G$ be the group of invertible elements of $A$. Define the map $\exp : \mathrm{A} \rightarrow \mathrm{A}$ by

$$
\exp (x)=\sum \frac{1}{n!} x^{n}
$$

(a) Show that A with $\exp$ is a universal covering of G.
(b) Show that if G is simply connected, then the map exp is injective.
(c) Deduce that a commutative complex Banach algebra is a field only if it has dimension 1 (Gelfand-Mazur theorem).
(d) Deduce that $\mathbb{C}$ is algebraically closed (d'Alembert's theorem).
(e) Extend the result of (c) to a not necessarily commutative algebra over $\mathbb{C}$ (consider a maximal commutative subalgebra).

For other proofs of d'Alembert's theorem, see (5.2.4), (4.6, Exercise 4) and (6.1, Exercise 3).

### 4.5 Galois Coverings

### 4.5.1

Theorem and Definition Let B be a connected space and E a connected covering of B . The following conditions are equivalent:
(i) E trivializes itself (i.e. $\mathrm{E} \times{ }_{\mathrm{B}} \mathrm{E}$ with $\mathrm{pr}_{1}: \mathrm{E} \times_{\mathrm{B}} \mathrm{E} \rightarrow \mathrm{E}$ is a trivial covering of E);
(ii) the group $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$ of automorphisms of E acts transitively on each fibre of E .

If these conditions hold, then E is said to be a Galois covering of B , and the group $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$ is called the Galois group of E .

Proof By (4.3.3), condition (i) is equivalent to the following:
(i') There is a continuous section $\mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{E}$ passing through all points of $\mathrm{E} \times_{\mathrm{B}} \mathrm{E}$. Condition ( $\mathrm{i}^{\prime}$ ) is equivalent to:
( $\mathrm{i}^{\prime \prime}$ ) For all $\left(x, x^{\prime}\right) \in \mathrm{E} \times_{\mathrm{B}} \mathrm{E}$, there is a B-morphism $f: \mathrm{E} \rightarrow \mathrm{E}$ such that $f(x)=x^{\prime}$.
(ii) $\Rightarrow\left(\mathrm{i}^{\prime \prime}\right)$ is obvious; we show the converse. If ( $\mathrm{i}^{\prime \prime}$ ) holds, for $x$ and $x^{\prime}$ in the same fibre, there are B -morphisms $f$ and $f^{\prime}$ from E to itself, such that $f(x)=x^{\prime}$ and $f^{\prime}\left(x^{\prime}\right)=x$. Then $f^{\prime} \circ f$ and $f \circ f^{\prime}$ have a fixed point, and so agree with $1_{\mathrm{E}}$ by (4.3.5).

Remark If $E$ is a Galois covering of $B$, then by (4.3.5) the Galois group acts simply transitively on the fibres.

### 4.5.2 Examples

(1) Under the assumptions of Theorem 4.3.13, if $X$ is connected, then $X$ is a Galois covering of $X / G$. Indeed, every element of $G$ defines an automorphism of X as a covering map of $\mathrm{X} / \mathrm{G}$, and by definition G acts transitively on the fibres of $\mathrm{X} \rightarrow \mathrm{X} / \mathrm{G}$. This applies to Example (4) of (4.3.4): $\mathrm{S}^{n} \rightarrow \mathbf{P}^{n} \mathbb{R}$.
(2) Similarly, under the assumptions of Corollary 4.3.14, if G is connected, then G is a Galois covering of $\mathrm{G} / \mathrm{H}$. This applies to Examples (2) and (3) of (4.3.4).
(3) Let B be a locally connected space, E a Galois covering of B, X a connected cover of B and $f: \mathrm{E} \rightarrow \mathrm{X}$ aB-morphism. Then E together with $f$ is a Galois covering of X . Indeed if $t, t^{\prime} \in \mathrm{E}$ such that $f(t)=f\left(t^{\prime}\right)$, then there exist $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$ such that $g(t)=t^{\prime}$. As $f$ and $f \circ g$ are two morphisms de E in X agreeing at $t, f=f \circ g$ by (4.3.5), and $g \in \operatorname{Aut}_{X}(\mathrm{E})$. The Galois group Aut ${ }_{X}(\mathrm{E})$ is the subgroup of $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$ consisting of all $g$ such that $f \circ g=f$ (4.5.7).
(4) Let $B$ be a connected space and $E$ a universal covering of $B$. Then $E$ is a Galois covering of $B$. Indeed, $E$ trivializes all coverings of $B$, in particular itself.
(5) If $X$ and $X^{\prime}$ are Galois coverings of $B$ and $B^{\prime}$ with Galois groups $G$ and $G^{\prime}$ respectively, then $X \times X^{\prime}$ is a Galois covering of $B \times B^{\prime}$, with group $G \times G^{\prime}$. Indeed $\mathrm{G} \times \mathrm{G}^{\prime}$ may be checked to act simply transitively on the fibres of $\mathrm{X} \times \mathrm{X}^{\prime}$.

For an example of a non Galois covering, see 4.5, Exercise 3.

### 4.5.3 The Functor $S$

Let $B$ be a connected and locally connected space, E a Galois covering of B and $\mathrm{G}=\operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$ its Galois group. Let G - $\mathscr{E n s}$ denote the category of sets on which G acts and $\mathscr{D}$ the category of coverings of B trivialized by E. For every object X of $\mathscr{D}$, let $\mathrm{S}(\mathrm{X})$ be the set $\operatorname{Hom}_{\mathrm{B}}(\mathrm{E}, \mathrm{X})$ on which G acts by $(g, f) \mapsto f \circ g^{-1}$. This defines a functor $S$ from $\mathscr{D}$ to G-EXns.

For $X \in \mathscr{D}, \operatorname{deg}_{B}(X)=\operatorname{Card} S(X)$. Indeed, $\operatorname{Card} S(X)$ is equal to the number of sections $\mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ equal to $\operatorname{deg}_{\mathrm{E}}\left(\mathrm{E} \times_{\mathrm{B}} \mathrm{X}\right)$ since $\mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ is a trivial covering of $E$ which is connected and

$$
\operatorname{deg}_{E}\left(E \times{ }_{B} X\right)=\operatorname{deg}_{B}(X)
$$

since the fibres are the same. In particular, $\mathrm{S}(\mathrm{X})=\varnothing \Leftrightarrow X=\varnothing$.
Theorem The functor $\mathrm{S}: \mathscr{D} \rightarrow$ G- $\mathscr{E n s}$ is an equivalence of categories.

This theorem amounts to the above Propositions 4.5.5 and 4.5.6.

### 4.5.4

Proposition Let X be an object of $\mathscr{D}$; equip $\mathrm{S}(\mathrm{X})$ with the discrete topology. The map $\varepsilon_{\mathrm{X}}: \mathrm{E} \times \mathrm{S}(\mathrm{X}) \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ defined by $\varepsilon_{\mathrm{X}}(t, s)=(t, s(t))$ a covering isomorphism of E compatible with the actions $\perp$ and $\cdot$ of G on $\mathrm{E} \times \mathrm{S}(\mathrm{X})$ and $\mathrm{E} \times{ }_{\mathrm{B}} \mathrm{X}$ respectively defined by

$$
g \perp(t, s)=\left(g(t), s \circ g^{-1}\right) \quad g \cdot(t, x)=(g(t), x)
$$

Proof Clearly, $\varepsilon_{\mathrm{X}}$ is a covering morphism of E; in particular it is open. We show it is bijective. If $(t, s(t))=\left(t^{\prime}, s^{\prime}\left(t^{\prime}\right)\right)$, then $t=t^{\prime}$, and $s=s^{\prime}$ by (4.3.5); so $\varepsilon_{\mathrm{X}}$ is injective. For all $(t, x) \in \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$, there is a section $\sigma: \mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ passing through $(t, x)$ since $\mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ is a trivial covering of E ; it corresponds to a B -morphism $s: \mathrm{E} \rightarrow \mathrm{X}$ such that $s(t)=x$. Hence $\varepsilon_{\mathrm{X}}$ is surjective, and it is a homeomorphism, and

$$
\begin{aligned}
\varepsilon_{\mathrm{X}}(g \perp(t, s))=\varepsilon_{\mathrm{X}}\left(g(t), s \circ g^{-1}\right) & =\left(g(t), s\left(g^{-1}(g(t))\right)\right) \\
& =(g(t), s(t))=g \cdot \varepsilon_{\mathrm{X}}(t, s) .
\end{aligned}
$$

### 4.5.5

## Proposition The functor S is fully faithful.

In other words, if X and Y are coverings of B trivialized by E , then the map $f \mapsto f_{*}$ from $\operatorname{Hom}_{\mathrm{B}}(\mathrm{X}, \mathrm{Y})$ to $\operatorname{Hom}_{\mathrm{G}}(\mathrm{S}(\mathrm{X}), \mathrm{S}(\mathrm{Y}))$ is bijective. Here, $\operatorname{Hom}_{\mathrm{B}}(\mathrm{X}, \mathrm{Y})$ denotes the set of covering morphisms from $X$ to $Y$, and $\operatorname{Hom}_{G}(S(X), S(Y))$ the set of maps $\varphi$ from $S(X)$ to $S(Y)$ compatible with the action of $G$, i.e. such that

$$
\varphi(g \cdot s)=g \cdot \varphi(s) \text { for } g \in \mathrm{G} \text { and } s \in \mathrm{~S}(\mathrm{X})
$$

and $f_{*}$ is defined by $f_{*}(s)=f \circ s$.
Lemma Let X be a covering of B . For the action of G on $\mathrm{E} \times{ }_{\mathrm{B}} \mathrm{X}$ given by the action of G on $\mathrm{E},\left(\mathrm{E} \times{ }_{\mathrm{B}} \mathrm{X}\right) / \mathrm{G}=\mathrm{X}$.

Proof of Lemma The fibres of $E \times_{B} X \rightarrow X$ are the same as the fibres of $E \rightarrow B$, and so are not empty and G acts transitively on them. Hence the set $\left(E \times_{B} X\right) / G$ can be identified with $X$. The topology of $X$ agrees with the quotient topology of $\mathrm{E} \times_{\mathrm{B}} \mathrm{X}$, since the map $\mathrm{E} \times_{\mathrm{B}} \mathrm{X} \rightarrow \mathrm{X}$ defines a covering, and so is open.

Proof of Proposition Let $\varphi: \mathrm{S}(\mathrm{X}) \rightarrow \mathrm{S}(\mathrm{Y})$ be a morphism in G-Éns. There is mor$\operatorname{phism} \varphi_{*}: \mathrm{E} \times \mathrm{S}(\mathrm{X}) \rightarrow \mathrm{E} \times \mathrm{S}(\mathrm{Y})$ defined by $\varphi_{*}(t, s)=(t, \varphi(s))$ corresponding to it. The morphism $\varphi_{*}$ is compatible with the action $\perp$ of G .

As $\varepsilon_{\mathrm{X}}$ and $\varepsilon_{\mathrm{Y}}$ are isomorphisms, there is a unique E-morphism $\varphi_{!}$from $\mathrm{E} \times_{\mathrm{B}} \mathrm{X}$ to $\mathrm{E} \times_{\mathrm{B}} \mathrm{Y}$ making the diagram

commutative. This morphism is compatible with the actions of G. It follows that there is a continuous map $f$ from $\mathrm{X}=\left(\mathrm{E} \times_{\mathrm{B}} \mathrm{X}\right) / \mathrm{G}$ to $\mathrm{Y}=\left(\mathrm{E} \times_{\mathrm{B}} \mathrm{Y}\right) / \mathrm{G}$, and $f$ is a B-morphism. We show that $f_{*}: \mathrm{S}(\mathrm{X}) \rightarrow \mathrm{S}(\mathrm{Y})$ is equal to $\varphi$, i.e. that $f(s(t))=(\varphi(s))(t)$ for $s \in \mathrm{~S}(\mathrm{X})$ and $t \in \mathrm{E}$. This follows from the commutativity of the diagram


Indeed $f_{*}=\varphi$, which shows that the map $f \mapsto f_{*}$ is surjective. We show that, if $\varphi=u_{*}$, for some B-morphism $u: \mathrm{X} \rightarrow \mathrm{Y}$, then $f=u$. The injectivity of $u \mapsto u_{*}$ will follow. For all $s \in \mathrm{~S}(\mathrm{X}), f \circ s=u \circ s$, and so $f$ and $u$ agree on $s(\mathrm{E})$. For all $s \in \mathrm{~S}(\mathrm{X})$, there is a section $\tilde{s}: \mathrm{E} \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$, and $\mathrm{E} \times_{\mathrm{B}} \mathrm{X}=\bigcup_{s \in \mathrm{~S}(\mathrm{X})} \tilde{s}(\mathrm{E})$. Hence $\mathrm{X}=\bigcup_{s \in \mathrm{~S}(\mathrm{X})} s(\mathrm{E})$, and so $f=u$.

### 4.5.6

Proposition The functor S is essentially surjective.
In other words, for all sets $A$ on which $G$ acts, there is a covering $X \in \mathscr{D}$ such that $S(X)$ is isomorphic to $A$ in $G-\mathscr{E} n s$.

Proof (a) Assume $\mathrm{A}=\mathrm{G} / \mathrm{H}$ for some subgroup H of G . Set $\mathrm{X}=\mathrm{E} / \mathrm{H}$ and let $\xi$ denote the canonical map $\mathrm{E} \rightarrow \mathrm{E} / \mathrm{H}$. We show that X is an object of $\mathscr{D}$. The $\operatorname{map} \varepsilon_{\mathrm{E}}$ (4.5.4) is a homeomorphism from $\mathrm{E} \times \mathrm{G}$ onto $\mathrm{E} \times{ }_{\mathrm{B}} \mathrm{E}$. The graph of the equivalence relation induced on E by H is the image under $\varepsilon_{\mathrm{E}}$ of $\mathrm{E} \times \mathrm{H}$ which is clopen in $\mathrm{E} \times \mathrm{G}$; so this graph is clopen in $\mathrm{E} \times_{\mathrm{B}} \mathrm{E}$. Hence by (4.3.8, (b)) X is a covering of B . On the other hand, by passing to the quotient, $\varepsilon_{\mathrm{E}}$ gives a morphism $\varepsilon^{\prime}: \mathrm{E} \times(\mathrm{H} \backslash \mathrm{G}) \rightarrow \mathrm{E} \times_{\mathrm{B}} \mathrm{X}$, where $\mathrm{H} \backslash \mathrm{G}$ is the set of classes $\mathrm{H} g$. This morphism is bijective, and so is an isomorphism, and $\mathrm{X} \in \mathscr{D}$.

We show that the map $g \mapsto g \cdot \xi=\xi \circ g^{-1}$ from G to $\mathrm{S}(\mathrm{X})$ is surjective. Choose $t_{0} \in \mathrm{E}$. For $f \in \mathrm{~S}(\mathrm{X})$, let $t \in \mathrm{E}$ such that $\xi(t)=f\left(t_{0}\right)$. There exists $g \in \mathrm{G}$ such that $g(t)=t_{0}$, and the morphisms $f$ and $\xi \circ g^{-1}$ agree at $t_{0}$, and so are equal (4.3.5, Corollary).

Therefore the group G acts transitively on $\mathrm{S}(\mathrm{X})$. We show that the stabilizer of $\xi$ is $H$. $A \approx S(X)$ since

$$
\begin{aligned}
\xi \circ h^{-1}=\xi & \Longleftrightarrow \xi=\xi \circ h \Longleftrightarrow \xi\left(h\left(t_{0}\right)\right)=\xi\left(t_{0}\right) \\
& \Longleftrightarrow \exists h^{\prime} \in \mathrm{H} h\left(t_{0}\right)=h^{\prime}\left(t_{0}\right) \Longleftrightarrow h \in \mathrm{H} .
\end{aligned}
$$

(b) If $\left(X_{i}\right)$ is a family of objects of $\mathscr{D}$, then $S\left(\bigsqcup X_{i}\right)=\bigsqcup S\left(X_{i}\right)$. Indeed, for all morphisms from E to $\bigsqcup \mathrm{X}_{i}$, the image of E is in some $\mathrm{X}_{i}$ since E is connected.
(c) Every $A \in G$ - $\mathscr{E n s}$ is of the form $\bigsqcup \mathrm{A}_{i}$, with $\mathrm{A}_{i} \approx G / \mathrm{H}_{i}$, where $H_{i}$ is a subgroup of G. By (a), for each $i, A_{i}$ is isomorphic to some $S\left(X_{i}\right)$ and (b) $A \approx S(X)$, where $\mathrm{X}=\bigsqcup \mathrm{X}_{i}$.

### 4.5.7

Proposition Let B be a locally connected space, X and E connected coverings of B and $f: \mathrm{E} \rightarrow \mathrm{X}$ a morphism. Assume E is Galois and let G be the Galois group $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$. Then,
(a) E trivializes X over B ;
(b) ( $\mathrm{E}, f$ ) is a Galois covering of X ;
(c) X is a Galois covering of B if and only if the group $\mathrm{H}=\operatorname{Aut}_{\mathrm{X}}(\mathrm{E}, f)$ consisting of $g \in \mathrm{G}$ such that $f \circ g=f$ is a normal subgroup of G .

Proof (a) By (4.3.1) it suffices to show that for $b \in \mathrm{~B}, t \in \mathrm{E}(b)$ and $x \in \mathrm{X}(b)$, there is a morphism $h: \mathrm{E} \rightarrow \mathrm{X}$ such that $h(t)=x$. By (4.3.6, Corollary 4.5), $f(\mathrm{E})=\mathrm{X}$, and so there exists $t^{\prime} \in \mathrm{E}(b)$ such that $f\left(t^{\prime}\right)=x$. As E is Galois, there exists $g \in \mathrm{G}$ such that $g(t)=t^{\prime}$. Then $h=f \circ g$ satisfies $h(t)=x$.
(b) Let $t, t^{\prime} \in f^{-1}(x)$. As E is Galois, there exists $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$ such that $g(t)=t^{\prime}$. Then $f(g(t))=f(t)$, and by (4.3.5) $f \circ g=f$. In other words, $g \in \operatorname{Aut}_{\mathrm{X}}(\mathrm{E}, f)$. Hence ( $\mathrm{E}, f$ ) is a Galois covering of X .

To prove (c), a lemma is needed.
Lemma X is a Galois covering of B if and only if $\mathrm{Aut}_{\mathrm{B}}(\mathrm{X})$ acts transitively on $S(X)=\operatorname{Hom}_{B}(E, X)$.

Proof of Lemma Suppose that X is Galois, and let $f$ and $f^{\prime}$ be two morphisms $\mathrm{E} \rightarrow \mathrm{X}$. Choose $t \in \mathrm{E}$ and set $x=f(t)$ and $x^{\prime}=f^{\prime}(t)$. There exists $h \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{X})$ such that $h(x)=x^{\prime}$. Then $h \circ f$ and $f^{\prime}$ agree at $t$, and so $h \circ f=f^{\prime}$ by (4.3.5).

Conversely, let $x$ and $x^{\prime}$ be in the same fibre $\mathrm{X}(b)$. Choose $t \in \mathrm{E}(b)$; there are morphisms $f$ and $f^{\prime}$ from E to X such that $f(t)=x, f^{\prime}(t)=x^{\prime}$. If $\operatorname{Aut}_{\mathrm{B}}(\mathrm{X})$ acts transitively on $\mathrm{S}(\mathrm{X})$, there exists $h \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{X})$ such that $h \circ f=f^{\prime}$. Then $h(x)=x^{\prime}$.

End of the Proof of the Proposition (c) The G-set $\mathrm{S}(\mathrm{X})$ is isomorphic to $\mathrm{G} / \mathrm{H}$ under $g \mapsto f \circ g^{-1}$. The stabilizer of $f \circ g^{-1}$ is $g \mathrm{H} g^{-1}$; if the automorphisms of $\mathrm{S}(\mathrm{X})$ act transitively, these stabilizers are all equal, and H is normal.

Conversely, if H is normal, then the right translations in G act transitively on G/H under G-automorphisms, and so X is Galois.

## Exercises 4.5. (Galois coverings)

1.-Let E be a 2-fold covering of B . Show that the map $\tau: \mathrm{E} \rightarrow \mathrm{E}$ assigning to each point the other point of the same fibre is an automorphism of E . Deduce that, if B and E are connected, then E is a Galois covering.
2.-Let $B$ be a connected space and $E$ a covering of $B$ such that $A u t_{B} E$ acts transitively on the fibres of E . Show that each connected component of E is a Galois covering of B.
3.-Let $X$ and $Y$ be curves having double points represented in the figure below and $\pi: \mathrm{Y} \rightarrow \mathrm{X}$ a continuous map sending bijectively each (open) arc $\mathrm{A}_{i}$ onto A and each $\mathrm{B}_{i}$ onto B by preserving the direction of the arrows.

(a) Show that Y together with $\pi$ is a covering of X . Give its degree.
(b) Determine Aut ${ }_{X} Y$. Is the covering Y Galois?
(c) Construct a 2 -fold covering $\tilde{Y}$ of $Y$ such that $\tilde{Y}$ is a Galois covering of X .
(d) Construct a 2-fold covering $\widetilde{X}$ of X such that $\tilde{\mathrm{Y}}$ is a 3-fold Galois covering of $\widetilde{X}$.
(e) Show (using the results of 4.7) that if B is a connected and locally connected space and Z a connected 3-fold non Galois covering of B , then there is a continuous map $f: \mathrm{X} \rightarrow \mathrm{B}$ such that $f^{*}(\mathrm{Z})$ is isomorphic to Y .
4.-Let $B$ be a locally connected space and $X$ a covering of $B$. Assume $X$ has degree 3 and is not Galois. Show that there is a 2-fold covering E of X which is a Galois covering of B , and that there is 2-fold covering Y of B such that E is a 3-fold Galois covering of Y. How many quotient coverings of E are there?

### 4.6 Fundamental Groups

### 4.6.1

Definition Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space, $\Phi$ the functor $\mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$ from the category $\mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}}$ of coverings of B to $\mathscr{E} \mathrm{ns}$. The fundamental group of B with respect to $b_{0}$, written $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$, is the automorphism group of the functor $\Phi$.

Hence an element $\gamma \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ consists in being given for every covering X of B a permutation $\gamma_{\mathrm{X}}$ of $\mathrm{X}\left(b_{0}\right)$, subject to the condition that, for all covering morphisms $f: \mathrm{X} \rightarrow \mathrm{Y}$, the following diagram is commutative.


Remark (1) If B is locally connected, then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)=\pi_{1}\left(\mathrm{~B}_{0}, b_{0}\right)$ where $\mathrm{B}_{0}$ is the connected component of B containing $b_{0}$.
(2) If $\mathscr{U}$ is an open cover of B , let $\mathscr{C} \mathscr{a} \boldsymbol{v}_{\mathrm{B}, \mathscr{U}}$ denote the category of coverings of B trivialized by the open subsets of $\mathscr{U}$. We may define $\pi_{1}\left(\mathrm{~B}, \mathscr{U}, b_{0}\right)$ as the automorphism group of the restriction $\Phi_{\mathscr{U}}$ of $\Phi$ to $\mathscr{C} \mathscr{\vartheta} \boldsymbol{v}_{\mathrm{B}, \mathscr{U}}$. Then

$$
\pi_{1}\left(\mathrm{~B}, b_{0}\right)=\underset{\longleftarrow}{\lim } \pi_{1}\left(\mathrm{~B}, \mathscr{U}, b_{0}\right)
$$

If $B$ is locally simply connected, then the inverse limit is reached by all simply connected open covers, in other words $\mathscr{C} \boldsymbol{a} \boldsymbol{v}_{\mathrm{B}}=\mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}, \mathscr{U}}$ and $\pi_{1}\left(\mathrm{~B}, b_{0}\right)=$ $\pi_{1}\left(\mathrm{~B}, \mathscr{U}, b_{0}\right)$ if $\mathscr{U}$ is such a cover. In other cases, the pro-object of the category of groups defined by $\mathrm{d} \pi_{1}\left(\mathrm{~B}, \mathscr{U}, b_{0}\right)$ is not all that interesting.

### 4.6.2

Proposition Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space and E a universal covering of B such that $\mathrm{E}\left(b_{0}\right) \neq \varnothing$. Then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ acts simply transitively on $\mathrm{E}\left(b_{0}\right)$.
Proof Let $t, u \in \mathrm{E}\left(b_{0}\right)$. We first show that there is at most one element $\gamma \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ such that $\gamma_{\mathrm{E}}(t)=u$. If $\gamma$ has this property, then for all coverings X of B and all $x \in$ $\mathrm{X}\left(b_{0}\right)$, there is a morphism $f: \mathrm{E} \rightarrow \mathrm{X}$ such that $f(t)=x$ and $\gamma_{\mathrm{X}}(x)=\gamma_{\mathrm{X}}(f(t))=$ $f\left(\gamma_{\mathrm{E}}(t)\right)=f(u)$, which defines $\gamma_{\mathrm{X}}(x)$ for all X and $x \in \mathrm{X}\left(b_{0}\right)$.

We next show existence. Define $\gamma_{u t} \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ by setting for all coverings X of B and all $x \in \mathrm{X}\left(b_{0}\right),\left(\gamma_{u t}\right)_{\mathrm{X}}(x)=f(u)$ where $f: \mathrm{E} \rightarrow \mathrm{X}$ is the unique morphism such that $f(t)=x$. We check that $\gamma_{u t}$ is a functorial morphism from $\Phi$ to $\Phi$. Then $\gamma_{t t}=1$ and $\gamma_{v u} \circ \gamma_{u t}=\gamma_{v t}$; so $\gamma_{u t}$ is an isomorphism, and we also have $\gamma_{u t}(t)=u$.

### 4.6.3 Comparison of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and Aut $\mathrm{B}_{\mathrm{B}}(\mathrm{E})$

Let E be a universal covering of B such that $\mathrm{E}\left(b_{0}\right) \neq \varnothing$. The group $\operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$ also acts simply transitively on $\mathrm{E}\left(b_{0}\right)$ since E is Galois (4.5.1, Remark). By definition
of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$, for $\gamma \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$, the permutations $\sigma_{\gamma}=\gamma_{\mathrm{E}}$ and $\rho_{g}=g_{\mid \mathrm{E}\left(b_{0}\right)}$ of $\mathrm{E}\left(b_{0}\right)$, respectively defined by $\gamma$ and $g$, commute.

Let $t_{0} \in \mathrm{E}\left(b_{0}\right)$. The inverse of the bijection $\theta \mapsto \theta \cdot t_{0}$ identifies $\mathrm{E}\left(b_{0}\right)$ with $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and the action $\sigma_{\gamma}$ becomes the left translation by $\gamma$. The action $\rho_{g}$ becomes the right translation by some $\alpha_{t_{0}}(g) \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$. Indeed, every permutation of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ commuting with left translations is a right translation.

Proposition The map $\alpha_{t_{0}}: \operatorname{Aut}_{\mathrm{B}}(\mathrm{E}) \rightarrow \pi_{1}\left(\mathrm{E}, b_{0}\right)$ is an anti-isomorphism.
This means that $\alpha_{t_{0}}$ is bijective and that $\alpha_{t_{0}}\left(g g^{\prime}\right)=\alpha_{t_{0}}\left(g^{\prime}\right) \cdot \alpha_{t_{0}}(g)$.
Proof The map $\alpha_{t_{0}}$ is bijective since $\operatorname{Aut}_{\mathrm{B}}(\mathrm{E})$ acts simply transitively on $\mathrm{E}\left(b_{0}\right)$, and

$$
x \cdot \alpha_{t_{0}}\left(g g^{\prime}\right)=g g^{\prime}(x)=g\left(g^{\prime}(x)\right)=g\left(x \cdot \alpha_{t_{0}}\left(g^{\prime}\right)\right)=x \cdot \alpha_{t_{0}}\left(g^{\prime}\right) \cdot \alpha_{t_{0}}(g)
$$

Corollary The groups $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$ are isomorphic.
Remark The anti-isomorphism $\alpha_{t_{0}}$ depends on the choice of $t_{0}$. Indeed, if $t_{1}$ is another point of $\mathrm{E}\left(b_{0}\right)$, then let $\beta \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ be the unique element such that $t_{1}=$ $\beta \cdot t_{0}$; so $\alpha_{t_{1}}$ is the composition of $\alpha_{t_{0}}$ and of an inner automorphism of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ defined by $\beta$. Indeed,

$$
g\left(\theta \cdot t_{1}\right)=\theta \cdot \alpha_{t_{1}}(g) \cdot t_{1} \text { for } \theta \in \pi_{1}\left(\mathrm{~B}, b_{0}\right) \text { and } g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{E}),
$$

where $g\left(\theta \cdot \beta \cdot t_{0}\right)=\theta \cdot \alpha_{t_{1}}(g) \cdot \beta \cdot t_{0}$, but $g\left(\theta \cdot \beta \cdot t_{0}\right)=\theta \cdot \beta \cdot \alpha_{t_{0}}(g) \cdot t_{0}$, where

$$
\alpha_{t_{1}}(g)=\beta \cdot \alpha_{t_{0}}(g) \cdot \beta^{-1}
$$

Hence, except when $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is commutative, $\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$ should not be identified with $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$.

### 4.6.4 Functoriality

Let $\left(\mathrm{B}, b_{0}\right)$ and $\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ be two pointed spaces and $f:\left(\mathrm{B}, b_{0}\right) \rightarrow\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ a continuous pointed map, i.e. a continuous map $f: \mathrm{B} \rightarrow \mathrm{B}^{\prime}$ such that $f\left(b_{0}\right)=b_{0}^{\prime}$. Let $\gamma \in \pi_{1}\left(\mathbf{B}, b_{0}\right)$. For every covering $\mathbf{X}^{\prime}$ of $\mathbf{B}^{\prime}, f^{*}\left(\mathbf{X}^{\prime}\right)$ is a covering of $\mathbf{B}$, and $\gamma_{f^{*}\left(X^{\prime}\right)}$ is a permutation of $f^{*}\left(\mathrm{X}^{\prime}\right)\left(b_{0}\right)=\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)$. When $\mathrm{X}^{\prime}$ varies, these $\gamma_{f^{*}\left(\mathrm{X}^{\prime}\right)}$ form an automorphism of the functor $\mathrm{X}^{\prime} \mapsto \mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)$, i.e. an element of $\pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$, which we will denote $f_{*}(\gamma)$.

The map $f_{*}: \pi_{1}\left(\mathrm{~B}, b_{0}\right) \rightarrow \pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ thus defined is a group homomorphism. If $g: \mathrm{B}^{\prime} \rightarrow \mathrm{B}^{\prime \prime}$ is a continuous map such that $g\left(b_{0}^{\prime}\right)=b_{0}^{\prime \prime}$, then $(g \circ f)_{*}=g_{*} \circ f_{*}$. In other words, $\pi_{1}$ is a functor from the category of pointed spaces to the category of groups.

Let $\mathrm{B}_{0}$ be the connected component of $b_{0}$ in B and $\iota: \mathrm{B}_{0} \rightarrow \mathrm{~B}$ the canonical injection. Then $\iota_{*}: \pi_{1}\left(\mathrm{~B}_{0}, b_{0}\right) \rightarrow \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is an isomorphism. Indeed, apart from the functor $\iota^{*}: \mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}} \rightarrow \mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}_{0}}$, there is a functor $\iota_{*}: \mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}_{0}} \rightarrow \mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}}$ assigning to all coverings $X$ of $B_{0}$ the space $X$ over $B$. These functors can be checked to define mutual inverses between $\pi_{1}\left(\mathrm{~B}_{0}, b_{0}\right)$ and $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$.

### 4.6.5

Proposition Let $f:\left(\mathrm{B}, b_{0}\right) \rightarrow\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ be a pointed space morphism, X and $\mathrm{X}^{\prime}$ respective coverings of B and $\mathrm{B}^{\prime}$, and $h: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ an $f$-morphism. Then the map $h: \mathrm{X}\left(b_{0}\right) \rightarrow \mathrm{X}^{\prime}\left(b_{0}\right)$ is an $f_{*}$-morphism in the following sense: for $g \in \pi_{1}\left(\mathrm{~B}_{0}, b_{0}\right)$ and $x \in \mathrm{X}\left(b_{0}\right), h(g \cdot x)=f_{*}(g) \cdot h(x)$.
Proof Let $\tilde{h}: \mathrm{X} \rightarrow \mathrm{B} \times_{B^{\prime}} \mathrm{X}^{\prime}$ be the B-morphism $x \mapsto(p(x), h(x))$, where $p$ is the projection $\mathrm{X} \rightarrow \mathrm{B}$.

The map $\tilde{h}: \mathrm{X}\left(b_{0}\right) \rightarrow\left(\mathrm{B} \times{ }_{\mathrm{B}^{\prime}} \mathrm{X}^{\prime}\right)\left(b_{0}\right)=\left\{b_{0}\right\} \times \mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)=\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)$ is a $\pi_{1}\left(\mathrm{~B}, b_{0}\right)-$ morphism, and by definition of $f_{*}$, the identity of $\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)$. The proposition follows.

### 4.6.6 Basepoint-Change

Proposition Let B be a connected space admitting a universal covering, and $b_{0}, b_{1} \in$ B. Then the groups $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\pi_{1}\left(\mathrm{~B}, b_{1}\right)$ are isomorphic.

Proof If E is a universal covering, then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\pi_{1}\left(\mathrm{~B}, b_{1}\right)$ are isomorphic to Aut $_{\mathrm{B}}$ (E) by (4.6.3, Corollary).

Remark (1) The assumption that B admits a universal covering can be weakened, but not removed. All that is needed is the existence of a simply connected space J and of a continuous map $g: \mathrm{J} \rightarrow \mathrm{B}$ whose image contains $b_{0}$ and $b_{1}$. For example, it is sufficient for $b_{0}$ and $b_{1}$ to be in the same arc-connected component of B.

For an example of a compact, connected non locally connected space B with two points $b_{0}$ and $b_{1}$ such that $\pi_{1}\left(\mathrm{~B}, b_{0}\right) \not \approx \pi_{1}\left(\mathrm{~B}, b_{1}\right)$, see (4.6, Exercise $\left.15, \mathrm{~B}\right)$. We do not know if there is such an example when $B$ is locally connected.
(2) There is no natural isomorphism from $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ onto $\pi_{1}\left(\mathrm{~B}, b_{1}\right)$, and hence these cannot be identified, and "the fundamental group of B " is not well defined even when B is connected and admits a universal covering. More precisely, let $\mathscr{C}$ be the category of connected spaces admitting a universal covering and $\mathscr{C}^{\bullet}$ the category of pointed spaces $\left(\mathrm{B}, b_{0}\right)$ such that $\mathrm{B} \in \mathscr{C}$. We defined a functor $\pi_{1}: \mathscr{C}^{*} \rightarrow \mathscr{G} \boldsymbol{r}$. There is no functor $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{E}_{\boldsymbol{\mu}}$ such that $\pi_{1}$ is isomorphic to the composition of the forgetful functor $\mathscr{C}^{\bullet} \rightarrow \mathscr{C}$ and of F (cf. 4.6, Exercise 8).

### 4.6.7

Let ( $\mathrm{B}, b_{0}$ ) and ( $\mathrm{B}^{\prime}, b_{0}^{\prime}$ ) be pointed topological spaces, $f$ and $g$ pointed continuous maps from ( $\mathrm{B}, b_{0}$ ) to $\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$. A pointed homotopy from $f$ to $g$ is a continuous map $h:[0,1] \times \mathrm{B} \rightarrow \mathrm{B}^{\prime}$ such that

$$
(\forall t \in[0,1]) h\left(t, b_{0}\right)=b_{0}^{\prime},(\forall b \in \mathrm{~B}) h(0, b)=f(b) \text { and } h(1, b)=g(b) .
$$

Theorem If there is a pointed homotopy from $f$ to $g$, then the homomorphisms $f_{*}$ and $g_{*}$ from $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ to $\pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ agree.

Proof Set I $=[0,1]$, let $\iota_{0}$ and $\iota_{1}$ be the injections $b \mapsto(0, b), b \mapsto(1, b)$ from B to $\mathrm{I} \times \mathrm{B}$, and $\rho$ the map $(t, b) \mapsto(0, b)$ from $\mathrm{I} \times \mathrm{B}$ to itself. Let $h$ be a pointed homotopy from $f$ to $g$ and $\mathrm{X}^{\prime}$ a covering of $\mathrm{B}^{\prime}$. Set $\mathrm{X}=h^{*}\left(\mathrm{X}^{\prime}\right), \mathrm{X}_{0}=\iota_{0}^{*}(\mathrm{X})=f^{*}\left(\mathrm{X}^{\prime}\right)$ and $\mathrm{X}_{1}=\iota_{1}^{*}(\mathrm{X})=g^{*}\left(\mathrm{X}^{\prime}\right)$. Then $\mathrm{X}_{0}\left(b_{0}\right)=\mathrm{X}_{1}\left(b_{0}\right)=\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)=\mathrm{X}\left(t, b_{0}\right)$ for all $t \in \mathrm{I}$, and

$$
\left.(h \circ \rho)\right|_{\mathrm{I} \times\left\{b_{0}\right\}}=\left.h\right|_{\mathrm{I} \times\left\{b_{0}\right\}}, \text { and so }\left.\rho^{*} \mathrm{X}\right|_{\mathrm{I} \times\left\{b_{0}\right\}}=\left.\mathrm{X}\right|_{\mathrm{I} \times\left\{b_{0}\right\}} .
$$

By Theorem 4.3.15, there is an isomorphism $\varphi: \mathrm{X} \rightarrow \rho^{*}(\mathrm{X})$ agreeing with $1_{\mathrm{X}}$ over $\{0\} \times$ B. By (4.3.5, Corollary), $\varphi$ agrees with $1_{\mathrm{X}}$ on $\mathrm{I} \times\left\{b_{0}\right\}$, in particular $\varphi_{\left(1, b_{0}\right)}=$ $1_{\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)}$. Then $\alpha=\iota_{1}^{*}(\varphi)$ is an isomorphism from $\mathrm{X}_{1}$ onto $\mathrm{X}_{0}$ such that $\alpha_{b_{0}}=1_{\mathrm{X}^{\prime}\left(b_{0}^{\prime}\right)}$.

Let $\gamma \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$. Since $\gamma$ is a functorial morphism, there is a commutative diagram

which can also be written as

and $\gamma_{\mathrm{X}_{0}}=\gamma_{\mathrm{X}_{1}}$. As this holds for all $\mathbf{X}^{\prime}, f_{*}(\gamma)=g_{*}(\gamma)$. As this holds for all $\gamma$, $f_{*}=g_{*}$.

### 4.6.8

Theorem Let B be a connected, locally connected space admitting a universal covering and $b_{0} \in \mathrm{~B}$. Then the functor $\Phi: \mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$ from $\mathscr{C}$ ou $\boldsymbol{u}_{\mathrm{B}}$ to $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ - $\mathscr{E n s}$ is an equivalence of categories.

Proof Let $\left(\mathrm{E}, t_{0}\right)$ be a pointed universal covering of $\left(\mathrm{B}, b_{0}\right)$. Set $\mathrm{G}=\mathrm{Aut}_{\mathrm{B}}(\mathrm{E})$. The isomorphism $g \mapsto \alpha_{t_{0}}\left(g^{-1}\right)$ (4.6.3) identifies G with $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$. Since E is a Galois covering of B and trivializes every covering X of B , by Theorem 4.5.3 the functor $S: X \mapsto \operatorname{Hom}(E, X)$ from $\mathscr{C} a v_{B}$ to G- $\mathscr{E} n s$ is an equivalence of categories. We show that the functors S and $\Phi$ are isomorphic. By the universal property of $\left(\mathrm{E}, t_{0}\right)$, for all coverings X of B , the map $\delta_{t_{0}}: f \mapsto f\left(t_{0}\right)$ from $\operatorname{Hom}(\mathrm{E}, \mathrm{X})$ onto $\mathrm{X}\left(b_{0}\right)$ is bijective. Hence it is a G-set isomorphism. Functoriality is obvious and so $\delta_{t_{0}}$ is an isomorphism of functors.

### 4.6.9 Dictionary

Let B be a connected, locally connected space admitting a universal covering, $b_{0} \in \mathrm{~B}$ and X a covering of B .
(a) If X is a disjoint union (resp. product covering) of all $\mathrm{X}_{\lambda}$, then the $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$-set $\mathrm{X}\left(b_{0}\right)$ is the sum (resp. product) of $\mathrm{X}_{\lambda}\left(b_{0}\right)$.
(b) Assuming $\mathrm{X} \neq \varnothing, \mathrm{X}$ is connected if and only if $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ acts transitively on $X\left(b_{0}\right)$. Indeed, $X$ est connected if and only if $X=X_{1} \sqcup X_{2}$ for some $X_{1} \neq \varnothing$ and $\mathrm{X}_{2} \neq \varnothing$, and a $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$-set F is non-homogenous if $\mathrm{F}=\mathrm{F}_{1} \sqcup \mathrm{~F}_{2}$ for some $\mathrm{F}_{1} \neq \varnothing$ and $\mathrm{F}_{2} \neq \varnothing$.

More generally, let $\left(\mathrm{X}_{\lambda}\right)$ be the connected components of X . Then $\mathrm{X}_{\lambda}\left(b_{0}\right)$ are the orbits in $\mathrm{X}\left(b_{0}\right)$ of the action of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$.
(c) Assigning to every pointed covering ( $\mathrm{X}, x_{0}$ ) of ( $\mathrm{B}, b_{0}$ ) the stabilizer of $x_{0}$ gives a bijection between the isomorphic classes of connected pointed coverings of ( $\mathrm{B}, b_{0}$ ) and the subgroups of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$. Indeed, for every group G this gives a bijection between the isomorphic classes of homogenous pointed G-sets and the subgroups of G mapping ( $\mathrm{F}, x_{0}$ ) onto the stabilizer of $x_{0}$; the inverse is given by $\mathrm{H} \mapsto \mathrm{G} / \mathrm{H}$.
(d) Let $\left(\mathrm{X}, x_{0}\right)$ and $\left(\mathrm{X}^{\prime}, x_{0}^{\prime}\right)$ be two connected pointed coverings of $\left(\mathrm{B}, b_{0}\right)$. The coverings X and $\mathrm{X}^{\prime}$ are isomorphic if and only the stabilizers S and $\mathrm{S}^{\prime}$ of $x_{0}$ and $x_{0}^{\prime}$ are conjugate. Indeed, for $g \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$, the stabilizer of $g x_{0}$ is $g \mathrm{~S} g^{-1}$; the coverings X and $\mathrm{X}^{\prime}$ are isomorphic if an only if there exists $g \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ such that all pointed coverings $\left(\mathrm{X}, g x_{0}\right)$ and $\left(\mathrm{X}^{\prime}, x_{0}^{\prime}\right)$ are isomorphic, i.e. $g \mathrm{~S}^{-1}=\mathrm{S}^{\prime}$.
(e) Let $\left(\mathrm{X}, x_{0}\right)$ be a connected pointed covering of $\left(\mathrm{B}, b_{0}\right)$. The covering X is Galois if and only if the stabilizer of $x_{0}$ is normal. Indeed, X is Galois if and only if, for all $g \in \mathrm{G}$, the pointed coverings of ( $\mathrm{X}, x_{0}$ ) and ( $\mathrm{X}, g x_{0}$ ) are isomorphic.

### 4.6.10

Proposition Let B be a connected space admitting a simply connected covering, $b_{0} \in \mathrm{~B},\left(\mathrm{X}, x_{0}\right)$ a pointed covering of $\left(\mathrm{B}, b_{0}\right)$ and $p: \mathrm{X} \rightarrow \mathrm{B}$ the projection. Then $p_{*}: \pi_{1}\left(\mathrm{X}, x_{0}\right) \rightarrow \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is injective and its image is the stabilizer of $x_{0}$.

Proof Let E be a simply connected pointed covering of B and $t_{0} \in \mathrm{E}\left(b_{0}\right)$. The pointed covering ( $\mathrm{E}, t_{0}$ ) is universal by (4.4.4), and so there is a morphism $f:\left(\mathrm{E}, t_{0}\right) \rightarrow$ ( $\mathrm{X}, x_{0}$ ). The space E together with $f$ is a covering of X by (4.3.6), which is universal by (4.4.4).

Consider the diagram

where $\iota$ is the canonical injection, $u(g)=g \cdot t_{0}, v\left(g^{\prime}\right)=g^{\prime} \cdot t_{0}$ and $w\left(g^{\prime}\right)=g^{\prime} \cdot x_{0}$. This diagram is commutative: by applying (4.6.5) to the $p$-morphism $1_{\mathrm{E}}$, this holds for the upper square, and is obvious for the lower triangle.

The maps $u$ and $v$ are bijective by (4.6.2), $\iota$ is injective and its image is $f^{-1}\left(x_{0}\right)$. Hence $p_{*}$ is injective and its image is $w^{-1}\left(x_{0}\right)$.

### 4.6.11

Corollary Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space admitting a simply connected pointed covering.
(a) For all subgroups H of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$, there is a connected pointed covering ( $\mathrm{X}, x_{0}$ ) of $\left(\mathrm{B}, b_{0}\right)$ such that the image of $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ in $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is H , and it is unique up to isomorphism.
(b) Assume that B is connected and let $\left(\mathrm{X}, x_{0}\right)$ be a connected pointed covering of $\left(\mathrm{B}, b_{0}\right)$. Then X is Galois if and only if the image of $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ is a normal subgroup of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$.

If X is Galois, then $\mathrm{Aut}_{\mathrm{B}}(\mathrm{X})$ is isomorphic to the quotient group

$$
\pi_{1}\left(\mathrm{~B}, b_{0}\right) / \pi_{1}\left(\mathrm{X}, x_{0}\right)
$$

### 4.6.12 Fundamental Group of a Product

Proposition Let $\mathrm{B}\left(\right.$ resp. $\left.\mathrm{B}^{\prime}\right)$ be a space admitting a connected pointed covering E (resp. $\mathrm{E}^{\prime}$ ), $b_{0} \in \mathrm{~B}$ and $b_{0}^{\prime} \in \mathrm{B}^{\prime}$ such that $\mathrm{E}\left(b_{0}\right)$ and $\mathrm{E}^{\prime}\left(b_{0}^{\prime}\right)$ are nonempty. Assume that B or $\mathrm{B}^{\prime}$ is locally connected. Then the homomorphism

$$
\pi_{1}\left(\mathrm{~B} \times \mathrm{B}^{\prime},\left(b_{0}, b_{0}^{\prime}\right)\right) \rightarrow \pi_{1}\left(\mathrm{~B}, b_{0}\right) \times \pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)
$$

given by the projections is an isomorphism.
Proof The groups $\mathrm{G}=\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\mathrm{G}^{\prime}=\pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ act simply transitively on $\mathrm{E}\left(b_{0}\right)$ and $\mathrm{E}^{\prime}\left(b_{0}^{\prime}\right)$ respectively, and so $\mathrm{G} \times \mathrm{G}^{\prime}$ acts simply transitively on $\mathrm{E} \times \mathrm{E}^{\prime}\left(b_{0}, b_{0}^{\prime}\right)=\mathrm{E}\left(b_{0}\right) \times \mathrm{E}^{\prime}\left(b_{0}^{\prime}\right)$. If B (resp. $\left.\mathrm{B}^{\prime}\right)$ is locally connected, then so is E (resp. $\mathrm{E}^{\prime}$ ). Thanks to (4.3.19, Corollary), the space $\mathrm{E} \times \mathrm{E}^{\prime}$ is simply connected, and so is a universal covering of $\mathrm{B} \times \mathrm{B}^{\prime}$. It follows that $\pi_{1}\left(\mathrm{~B} \times \mathrm{B}^{\prime},\left(b_{0}, b_{0}^{\prime}\right)\right)$ can be identified with $\mathrm{G} \times \mathrm{G}^{\prime}$.

## Exercises 4.6. (Fundamental group)

1.-Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space admitting a pointed universal covering. Show that every endomorphism of the functor $\Phi$ of (4.6.1) is an automorphism. Does this result still hold without the assumption that $\left(\mathrm{B}, b_{0}\right)$ admits a pointed universal covering?
2.-Let $X$ be a topological space and $G$ a group acting on $X$ in such a way that assumption (L) of Theorem 4.3.13 is satisfied. Set $\mathrm{B}=\mathrm{X} / \mathrm{G}$. Let $x_{0} \in \mathrm{X}$ and $b_{0}$ be its image in $B$.
(a) Show that if X is simply connected, then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is isomorphic to G .
(b) Show that if X is connected and admits a simply connected covering, then X is a Galois covering of B and G is isomorphic to the quotient group of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ by the image of $\pi_{1}\left(\mathrm{X}, x_{0}\right)$.
(c) Show that if X is locally simply connected, then each connected component of X is a Galois covering of a connected component of B , and the quotient of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ by the image of $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ is isomorphic to the stabilizer of the connected component of $b_{0}$ in B.
3. (Degree of a map from $S^{1}$ to $S^{1}$.) -Set $S^{1}=\{z \in \mathbb{C}| | z \mid=1\}$ and let $\varepsilon$ be the map $t \mapsto e^{2 i \pi t}$ from $\mathbb{R}$ to $\mathrm{S}^{1}$.
(a) Show that $\pi_{1}\left(\mathrm{~S}^{1}, 1\right)$ is isomorphic to $\mathbb{Z}$. Since this group is commutative, $\pi_{1}\left(\mathrm{~S}^{1}\right)$ is well defined.
(b) Let $f: \mathrm{S}^{1} \rightarrow \mathrm{~S}^{1}$ be a continuous map, define $\mathrm{d}^{\circ}(f) \in \mathbb{Z}$ by $f_{*}(\gamma)=\mathrm{d}^{\circ}(f) \cdot \gamma$ for $\gamma \in \pi_{1}\left(\mathrm{~S}^{1}\right)$; this number is called the degree of $f$. Show that there is a continuous $\operatorname{map} \tilde{f}: \mathbb{R} \rightarrow \mathbb{R}$ such that $\varepsilon \circ \tilde{f}=f \circ \varepsilon$ and that for any $\tilde{f}$ with this property and all $t \in \mathbb{R}$,

$$
\tilde{f}(t+1)=\tilde{f}(t)+\mathrm{d}^{\circ}(f) .
$$

(c) Show that any two continuous maps $f$ and $g$ from $\mathrm{S}^{1}$ to $\mathrm{S}^{1}$ are homotopic if and only if $\mathrm{d}^{\circ}(f)=\mathrm{d}^{\circ}(g)$.
(d) Let $f: \mathrm{S}^{1} \rightarrow \mathrm{~S}^{1}$ be a differentiable map (i.e. continuous and such that $\tilde{f}$ is differentiable). For $z \in \mathrm{~S}^{1}$, let $f^{\prime}(z)$ be the common value of all $\tilde{f}^{\prime}(t)$ for $t \in \varepsilon^{-1}(z)$. Let $a \in \mathrm{~S}^{1}$; assume that $f^{\prime}(z) \neq 0$ for all $z \in f^{-1}(a)$. Show that the set $f^{-1}(a)$ is finite and that

$$
\mathrm{d}^{\circ}(f)=\operatorname{Card}\left\{z \in f^{-1}(a) \mid f^{\prime}(z)>0\right\}-\operatorname{Card}\left\{z \in f^{-1}(a) \mid f^{\prime}(z)<0\right\}
$$

(e) Let $f$ be a continuous map from $\mathrm{S}^{1}$ to $\mathbb{C}^{*}=\mathbb{C}-\{0\}$; set $\mathrm{d}^{\circ}(f)=\mathrm{d}^{\circ}\left(\frac{f}{|f|}\right)$. Assuming that $f$ is infinitely differentiable, express $\mathrm{d}^{\circ}(f)$ as an integral.
4. (A proof of d'Alembert's theorem.) -(a) Calculate the degree of the map $z \mapsto z^{d}$ from $S^{1}$ to $S^{1}$.
(b) Show that two sufficiently near (relative to uniform convergence) maps from $S^{1}$ to $S^{1}$ are homotopic.
(c) Let P be a complex polynomial of degree $d$, and for $r \in \mathbb{R}_{+}$let $f_{r}$ be the map $z \mapsto \mathrm{P}(r z)$ from $\mathrm{S}^{1}$ to $\mathbb{C}$. Show that, for sufficiently large $r, f_{r}$ is a map from $\mathrm{S}^{1}$ to $\mathbb{C}^{*}$ of degree $d$.
(d) Show that, if P has no roots $z$ such that $r \leqslant|z| \leqslant r^{\prime}$, then the maps $f_{r}$ and $f_{r^{\prime}}$ from $\mathrm{S}^{1}$ to $\mathbb{C}^{*}$ have the same degree. What is the degree of $f_{0}$ (assuming $\left.\mathrm{P}(0) \neq 0\right)$ ? Show that P has at least one root if $d>0$.
5.-Let $G$ be a finite group. The aim is to show that there is a compact manifold whose fundamental group is isomorphic to G.
(a) Show that, for sufficiently large $n, \mathrm{G}$ is isomorphic to a subgroup of the group $\mathrm{O}_{n}$ of isometries of Euclidean $\mathbb{R}^{n}$.
(b) Show that, for $p \geqslant n+2$, the space $\mathrm{V}_{n}^{p}$ of isometric embeddings of $\mathbb{R}^{n}$ in $\mathbb{R}^{p}$ is a simply connected compact manifold (4.1, Exercise 11 and 4.3, Exercise 4).
(c) Make G act freely on $\mathrm{V}_{n}^{p}$. Conclude.
6.-Let $G$ be an arbitrary group. The aim is to show that there is a connected space

B admitting a contractible universal covering and a point $b_{0} \in \mathrm{~B}$ such that $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is isomorphic to G .
(a) Let E be the space of maps $f:[0,1] \rightarrow \mathrm{G}$ such that there is an increasing finite sequence $\left(t_{0}, \ldots, t_{n}\right)$ with $t_{0}=0, t_{n}=1, f(0)=0$ and $f$ constant on $\left.] t_{i-1}, t_{i}\right]$ for $i \in\{1, \ldots, n\}$. For $f, g \in \mathrm{E}$, let $d(f, g)$ be the sum of the lengths of intervals on
which they differ. Show that E becomes a metric group and that G can be identified to a discrete subgroup of E .
(b) For $t \in[0,1]$ and $f \in \mathrm{E}$, define $f_{t}$ by $f_{t}(s)=f(s)$ if $s \leqslant t$ and $f_{t}(s)=e$ if $s>t$. Show that the map $(t, f) \mapsto f_{t}$ from [0, 1] $\times \mathrm{E}$ to E is continuous. Deduce that $E$ is contractible.
(c) Conclude.
7.-Let K be a connected space admitting a universal contractible covering $\mathrm{E}, k_{0} \in \mathrm{~K}$ and $e_{0} \in \mathrm{E}\left(k_{0}\right)$. The aim is to show that, for any metrizable compact pointed space $\left(\mathrm{B}, b_{0}\right)$ admitting a universal covering, and for any homomorphism $\varphi: \pi_{1}\left(\mathrm{~B}, b_{0}\right) \rightarrow$ $\pi_{1}\left(\mathrm{~K}, k_{0}\right)$, there is a pointed continuous map $f:\left(\mathrm{B}, b_{0}\right) \rightarrow\left(\mathrm{K}, k_{0}\right)$ such that $f_{*}=\varphi$, and that it is unique up to pointed homotopy.
(a) Let $G$ be a group. A principal G-covering of $B$ is a covering of $B$ equipped with a continuous $G$-action inducing a simply transitive action of $G$ on each fibre. Show that every principal G-cover admitting a section is trivial.
(b) Set $\mathrm{G}=\pi_{1}\left(\mathrm{~K}, k_{0}\right)$, and consider E as a principal G-covering. Prove that the problem amounts to showing that, for any principal pointed G-covering ( $\mathrm{X}, x_{0}$ ) of $\left(\mathrm{B}, b_{0}\right)$, there is, up to pointed homotopy, a unique pointed continuous map $f$ : $\left(\mathrm{B}, b_{0}\right) \rightarrow\left(\mathrm{K}, k_{0}\right)$ such that all principal pointed G-coverings $\left(\mathrm{X}, x_{0}\right)$ and $f^{*}\left(\mathrm{E}, e_{0}\right)$ of $\left(\mathrm{B}, b_{0}\right)$ are isomorphic.
(c) Let F be a contractible space, B a compact metrizable space, A a closed subset of $\mathrm{B}, \mathrm{U}$ a neighbourhood of A and $f: \mathrm{U} \rightarrow \mathrm{F}$ a continuous map. Show that there is a continuous map $g: \mathrm{B} \rightarrow \mathrm{F}$ agreeing with $f$ in the neighbourhood of A (if F is a convex containing 0 in a normed space, multiply $f$ by a function with support in U ).
(d) Let B be a compact metrizable space, A a closed subset of B, U a neighbourhood of A and X a trivial principal G-covering of B. Show that, for all continuous $f: \mathrm{U} \rightarrow \mathrm{K}$ such that $f^{*}(\mathrm{E})$ is isomorphic to $\left.\mathrm{X}\right|_{\mathrm{U}}$, there is a continuous map $g: \mathrm{B} \rightarrow \mathrm{K}$ such that $g^{*}(\mathrm{E})$ is isomorphic to X . More precisely, if $\alpha:\left.\mathrm{X}\right|_{\mathrm{U}} \rightarrow f^{*}(\mathrm{E})$ is an isomorphism, then $g$ may be chosen so that there is an isomorphism $\beta: \mathrm{X} \rightarrow g^{*}(\mathrm{E})$ agreeing with $\alpha$ over a neighbourhood of A .
(e) Prove the result without assuming that X is trivial (take a finite open cover $\left(\mathrm{U}_{i}\right)_{1 \leqslant i \leqslant n}$ of B trivializing X , a shrinking $\left(\mathrm{V}_{i}\right)$, i.e. an open cover such that $\overline{\mathrm{V}_{i}} \subset \mathrm{U}_{i}$, a shrinking $\left(\mathrm{W}_{i}\right)$ of $\left(\mathrm{V}_{i}\right)$, set $\mathrm{A}_{i}=\mathrm{A} \cup \overline{\mathrm{W}_{1}} \cup \cdots \cup \overline{\mathrm{~W}_{i}}$, and proceed by induction by applying (d) to $\left(\overline{\mathrm{V}_{i}}, \overline{\mathrm{~V}_{i}} \cap \mathrm{~A}_{i-1}\right)$ ).
(f) Conclude: for existence, take $\mathrm{A}=\left\{b_{0}\right\}$; for uniqueness, consider

$$
\left([0,1] \times \mathrm{B},\{0,1\} \times \mathrm{B} \cup[0,1] \times\left\{b_{0}\right\}\right) .
$$

(g) Can the assumption that B is compact and metrizable be replaced by a weaker assumption? Can the result be formulated in terms of the representability of some functor?
8. (Non-functoriality of $\pi_{1}$ without basepoint.) -With the notation of (4.6.6, Remark 2), assume there is a functor $\mathrm{F}: \mathscr{C} \rightarrow \mathscr{G}_{\boldsymbol{r}}$ such that the diagram

commutes.
(a) Show that, for any space B, the maps $\iota_{0 *}$ and $\iota_{1 *}$ from $\mathrm{F}(\mathrm{B})$ to $\mathrm{F}([0,1] \times \mathrm{B})$ agree. Deduce that, if $f$ and $g$ are two pointed continuous maps from ( $\mathrm{B}, b_{0}$ ) to ( $\mathrm{B}^{\prime}, b_{0}^{\prime}$ ) such there is a non-pointed homotopy from $f$ to $g$, then the maps $f_{*}$ and $g_{*}$ from $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ to $\pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ agree.
(b) Let $f$ and $g$ are two pointed continuous maps from ( $\mathrm{B}, b_{0}$ ) to ( $\mathrm{B}^{\prime}, b_{0}^{\prime}$ ), and assume there is non-pointed homotopy $h$ from $f$ to $g$. Show that $g_{*}: \pi_{1}\left(\mathbf{B}, b_{0}\right) \rightarrow$ $\pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ is the composition of $f_{*}$ with the inner automorphism defined by the element $\gamma \in \pi_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ corresponding to the loop $t \mapsto h\left(t, b_{0}\right)$ (see 4.9.9). Deduce a contradiction.
(c) Can there be a functor $\mathrm{G}: \mathscr{C} \rightarrow \mathscr{E} \boldsymbol{\mu}$ such that the group $\mathrm{G}(\mathrm{B})$ is isomorphic to $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ for all $\left(\mathrm{B}, b_{0}\right) \in \mathscr{C}^{*}$ ?
9.-In Proposition 4.6.10, instead of a simply connected covering, can B be assumed to admit a universal covering? (study the example of 4.3, Exercise 2).
10. Let U be the set of all $(p, q) \in \mathbb{C}^{2}$ such that $4 p^{3}+27 q^{2} \neq 0$. We remind the reader that, for $(p, q) \in \mathrm{U}$, the equation $\mathrm{X}^{3}+p \mathrm{X}+q=0$ has three distinct roots. Set $\mathrm{V}=\left\{(p, q, x) \in \mathrm{U} \times \mathbb{C} \mid x^{3}+p x+q=0\right\}$ and $\mathrm{E}=\left\{(x, y, z) \in \mathbb{C}^{3} \mid\right.$ $x+y+z=0, x \neq y \neq z \neq x\}$. Equip V with the projection $(p, q, x) \mapsto(p, q)$ from V to U and E with the map $(x, y, z) \mapsto(p, q, x)$ from E to V , where $p$ and $q$ are defined by

$$
\mathrm{X}^{3}+p \mathrm{X}+q=(\mathrm{X}-x)(\mathrm{X}-y)(\mathrm{X}-z)
$$

(a) Show that $E$ is a covering of $V$, that $V$ is a covering of $U$ and that $E$ is a covering of $U$.

What are the degrees of these coverings?
(b) Show that E is connected. Deduce that V and U are connected.
(c) Let $a=(-1,0) \in \mathrm{U}$. Show that the action of $\pi_{1}(\mathrm{U}, a)$ on the fibre $\mathrm{V}_{a}$ defines a surjective homomorphism from $\pi_{1}(\mathrm{U}, a)$ onto $\mathfrak{S}_{3}$.
(d) Which of the coverings $\mathrm{E} \rightarrow \mathrm{V}, \mathrm{V} \rightarrow \mathrm{U}, \mathrm{E} \rightarrow \mathrm{U}$, are Galois?
(e) Let $\mathrm{S}^{3}$ be the unit sphere of $\mathbb{C}^{2}$, i.e. $\mathrm{S}^{3}=\{(p, q) \mid p \bar{p}+q \bar{q}=1\}$, and set $\Gamma=\left\{(p, q) \in \mathrm{S}^{3} \mid 4 p^{3}+27 q^{2}=0\right\}$. Show that $\mathrm{S}^{3}-\Gamma$ has the same homotopy type as U. Deduce that $\pi_{1}\left(\mathrm{~S}^{3}-\Gamma, a\right)$ is not commutative.
(f) Let $\Gamma_{0}=\{(\lambda, 0)| | \lambda \mid=1\}$. Show that $\Gamma$ is homeomorphic to $\Gamma_{0}$, but that there is no homeomorphism $f: \mathrm{S}^{3} \rightarrow \mathbf{S}^{3}$ such that $f(\Gamma)=\Gamma_{0}$.
(g) The space $S^{3}-\{a\}$ is homeomorphic to $\mathbb{R}^{3}$. Represent with an iron wire the image of $\Gamma$ under a homeomorphism from $S^{3}-\{a\}$ onto $\mathbb{R}^{3}$.
11.-Let G be a topological group and $e$ its identity element. Assume that G admits a universal covering.
(a) Show that $\pi_{1}(\mathrm{G}, e)$ is commutative (2.5, Exercise 2, (d)).
(b) Let $(\mathrm{H}, e)$ be a connected pointed covering of $(\mathrm{G}, e)$. Show that there is a unique topological group law on H having $e$ as its identity, such that the projection $\mathrm{H} \rightarrow \mathrm{G}$ is a homomorphism. (The case of a universal covering may be first considered.)
(c) Can the space $\mathbb{C}-\{+1,-1\}$ be endowed with a composition law making it a topological group?
12.-Let $\left(\mathrm{X}, x_{0}\right)$ and $\left(\mathrm{Y}, y_{0}\right)$ be two pointed spaces and let $f:\left(\mathrm{X}, x_{0}\right) \rightarrow\left(\mathrm{Y}, y_{0}\right)$ be a pointed continuous map. Let $\left(\mathrm{E}, e_{0}\right)$ be a pointed universal covering of $\left(\mathrm{Y}, y_{0}\right)$. Show that the image of $\pi_{1}\left(f^{*} \mathrm{E},\left(x_{0}, e_{0}\right)\right)$ in $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ is the kernel of the homomor$\operatorname{phism} f_{*}: \pi_{1}\left(\mathrm{X}, x_{0}\right) \rightarrow \pi_{1}\left(\mathrm{Y}, y_{0}\right)$.
13.-Let V be a manifold and $\omega$ a real-valued closed differential form of degree 1 on V . A subset in $\mathrm{V} \times \mathbb{R}$ will be called a $\omega$-slice if it is the graph of a function $f: \mathrm{U} \rightarrow \mathbb{R}$ such that $d f=\left.\omega\right|_{\mathrm{U}}$, where U is open in V . The slices form a basis with respect to the topology $\mathscr{T}$ on $\mathrm{V} \times \mathbb{R}$.
(a) Show that $\mathrm{V} \times \mathbb{R}$, equipped with the topology $\mathscr{T}$ and the projection $p r_{1}$ onto V , is a covering of V .
(b) Let $b \in \mathrm{~V}$. Show that there is a homomorphism $\varphi: \pi_{1}(\mathrm{~V}, b) \rightarrow \mathbb{R}$ such that the action of $\pi_{1}(\mathrm{~V}, b)$ on the fibre of this covering, namely on $\mathbb{R}$ equipped with the discrete topology, is given by $\gamma \cdot x=x+\varphi(\gamma)$.
(c) Let $\gamma \in \pi_{1}(\mathrm{~V}, b)$ be represented by a differentiable loop C . Show that

$$
\varphi(\gamma)=\int_{\mathrm{C}} \omega
$$

(d) Show that the following conditions are equivalent:
(i) $\omega$ is exact, i.e. $(\exists f) \omega=d f$;
(ii) $\varphi=0$;
(iii) the covering $(\mathrm{V} \times \mathbb{R}, \mathscr{T})$ is trivial.
(e) Assume V is connected and let E be a connected component of $(\mathrm{V} \times \mathbb{R}, \mathscr{T})$ and $\pi: \mathrm{E} \rightarrow \mathrm{V}$ the projection. Show that E is a Galois covering of V and that the differential form $\pi^{*} \omega$ on E is exact.
(f) Show that $\omega \mapsto \varphi$ is a linear map from the space of closed differential forms on V to the vector space $\operatorname{Hom}\left(\pi_{1}(\mathrm{~V}, b), \mathbb{R}\right)$. What is its kernel?
(g) Generalize the results to complex valued differential forms.
14.-(a) Let G be a group and H a subgroup of G . Show that $\mathrm{Aut}_{\mathrm{G}} \mathrm{G} / \mathrm{H}$ is isomorphic to $\mathrm{N} / \mathrm{H}$, where N is the normalizer of H in G .
(b) Let ( $\mathrm{B}, b_{0}$ ) be a connected, locally simply connected pointed space, and let $\left(\mathrm{X}, x_{0}\right)$ be a connected pointed covering of $\left(\mathrm{B}, b_{0}\right)$. Describe Aut ${ }_{\mathrm{B}} \mathrm{X}$.
15.-(A) Let B be a topological space, $b_{0}, b_{1} \in \mathrm{~B}$. Assume that $b_{0}$ and $b_{1}$ are in the same arc-connected component, or more generally, that there is a simply connected space J and a continuous map $g: \mathrm{J} \rightarrow \mathrm{B}$ such that $g(\mathrm{~J}) \supset\left\{b_{0}, b_{1}\right\}$. Let $\Phi_{0}$ (resp. $\left.\Phi_{1}\right)$ denote the functor $\mathscr{C} \boldsymbol{a} \boldsymbol{v}_{\mathrm{B}} \rightarrow \mathscr{E} n s$ assigning the fibre $\mathrm{X}\left(b_{0}\right)$ (resp. $\mathrm{X}\left(b_{1}\right)$ ) to the covering X of B. Show that the functors $\Phi_{0}$ and $\Phi_{1}$ are isomorphic (they are both isomorphic to the functor $\operatorname{Sec} \circ g^{*}$ assigning to $\mathbf{X}$ the set of sections de $g^{*}(\mathbf{X})$ ).

Deduce that the groups $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\pi_{1}\left(\mathrm{~B}, b_{1}\right)$ are isomorphic.
(B) Let $p$ be a (not necessarily prime) integer $\geqslant 2$, and take B to be the solenoid $\widehat{\mathbb{R}}=$ $\lim _{\leftrightarrows} \mathbb{R} / p^{n} \mathbb{Z}$ with its group structure, the projections $\pi_{n}: \widehat{\mathbb{R}} \rightarrow \mathbb{R} / p^{n} \mathbb{Z}$, in particular $\pi_{0}: \widehat{\mathbb{R}} \rightarrow \mathbb{R} / \mathbb{Z}$, and the injection $\iota: \mathbb{R} \rightarrow \widehat{\mathbb{R}}$. Let $b_{0}, b_{1} \in \mathrm{~B}$.

The aim is to prove that if $b_{0}$ and $b_{1}$ are not in the same arc-connected component of B, then the functors $\Phi_{0}$ and $\Phi_{1}$ are not isomorphic.
(a) Show that $b_{0}$ and $b_{1}$ are in the same arc-connected component of $\mathrm{B}=\widehat{\mathbb{R}}$ if and only if $b_{1}-b_{0} \in \iota(\mathbb{R})$.
(b) Let $\widehat{\mathbb{Z}}$ be the subgroup $\lim _{\leftarrow} \mathbb{Z} / p^{n} . \mathbb{Z}$ of $\widehat{\mathbb{R}}$, and by $\psi:(s, t) \mapsto s+\iota(t)$, consider $\mathrm{E}=\widehat{\mathbb{Z}} \times \mathbb{R}$ as a space over $\widehat{\mathbb{R}}$.
Show that E is a covering of $\widehat{\mathbb{R}}$ isomorphic to $\pi_{0}^{*}(\mathbb{R})$, where $\mathbb{R}$ is considered the universal covering of $\mathbb{R} / \mathbb{Z}$.
For all $n$, consider $\mathrm{E}_{n}=p^{n} \mathbb{Z} \times \mathbb{R} \subset \mathrm{E}$. Show that $\mathrm{E}_{n}$, equipped with the induced projection is a covering of $\widehat{\mathbb{R}}$, isomorphic to $\pi_{n}^{*}(\mathbb{R})$.
(c) Let $b \in \widehat{\mathbb{R}}$. Show that $\bigcap \mathrm{E}_{n}(b)=\{0\}$ if $b=0$, and is empty if $b \notin \iota(\mathbb{R})$. Conclude.
(C) Keep $B=\widehat{\mathbb{R}}$.
(a) Show that, for all $b \in \mathrm{~B}$, the fundamental group $\pi_{1}(\mathrm{~B}, b)$ acts trivially on $\mathrm{E}(b)$, although the fiber bundle E is not trivial.
(b) Show that there is some $n$ such that every covering of B is induced by a covering of $\mathbb{R} / p^{n} . \mathbb{Z}$ by change of basis by $\pi_{n}$. Show that for all $b, \pi_{1}(\mathrm{~B}, b)=\{0\}$ although B is connected and not simply connected.
(D) Let $\mathrm{B}^{\prime}$ be a connected space consisting of the union of two closed subsets of $B$ and $\Gamma$, with $B=\widehat{\mathbb{R}}, \Gamma$ homeomorphic to $S^{1}$ and $B \cap \Gamma=\{0\}$.

Let X be a covering of $\mathrm{B}^{\prime}$ inducing a universal covering over $\Gamma$ and a trivial covering over B . Denote by $\pi_{\mathrm{X}}$ the projection $\mathrm{X} \rightarrow \mathrm{B}^{\prime}$. Consider X as the real line $\mathbb{R}$ with a copy $\mathrm{B}_{k}$ of B at each point $k \in \mathbb{Z}$; for each $k$, there exists $b_{k} \in \mathrm{~B}_{k}$ over $b$.
(a) Show that $\pi_{1}\left(\mathrm{~B}^{\prime}, 0\right)$ acts non trivially on $\mathrm{X}(0)$. Deduce that $\pi_{1}\left(\mathrm{~B}^{\prime}, 0\right) \neq\{0\}$.
(b) The aim is to show that, for all $b \in \mathrm{~B}$, the group $\pi_{1}\left(\mathrm{~B}^{\prime}, b\right)$ acts trivially on $\mathrm{X}(b)$ if $b \notin \iota(\mathbb{R})$.

Let $\mathscr{C} \bullet v_{\mathrm{B}^{\prime}, \mathrm{X}}$ be the category of coverings Y of $\mathrm{B}^{\prime}$ together with a morphism $\varphi: \mathrm{Y} \rightarrow \mathrm{X}$. If Y is an object of $\mathscr{C} \bullet \boldsymbol{v}_{\mathrm{B}^{\prime}, \mathrm{X}}$, then the fibre $\mathrm{Y}(b)$ is the disjoint union of $\mathrm{Y}\left(b_{k}\right)=\varphi^{-1}\left(b_{k}\right)$. Let $\alpha \in \pi_{1}\left(\mathrm{~B}^{\prime}, b\right)$. Then $\alpha$ acts on $\mathrm{X}(b)$ and on $\mathrm{Y}(b)$, and $\alpha . \mathrm{Y}\left(b_{k}\right)=\mathrm{Y}\left(\alpha \cdot b_{k}\right)$.

Assuming that $\alpha$ acts non trivially on $\mathrm{X}(b)$, let $k$ be such that $\alpha\left(b_{k}\right)=b_{k^{\prime}}$ with $k^{\prime} \neq k$. Restricting to the objects of $\mathscr{C} a \boldsymbol{v}_{\mathrm{B}^{\prime}, \mathrm{X}}$ inducing a trivial covering over $\mathrm{B}_{k}$ (but an arbitrary one over $\mathrm{B}_{k^{\prime}}$, and starting with $\alpha$, define an isomorphism between the functors $\Phi_{0}$ and $\Phi_{b}: \mathscr{C} \not \boldsymbol{v}_{\mathrm{B}} \rightarrow \mathscr{E} n s$. Conclude.

### 4.7 Van Kampen's Theorem

The aim of this section is to find the fundamental group of a space B when $\mathrm{B}=$ $\mathrm{U}_{1} \cup \mathrm{U}_{2}$ and the fundamental groups of $\mathrm{U}_{1}, \mathrm{U}_{2}$ and $\mathrm{U}_{1} \cap \mathrm{U}_{2}$ are known.

### 4.7.1 Categorical Preliminaries

Let G and H be two groups and $f: \mathrm{H} \rightarrow \mathrm{G}$ a homomorphism. For all G -sets X , let $f^{*}(\mathrm{X})$ be the H -set obtained by making H act on X by $h \cdot x=f(h) x$. This gives a functor $f^{*}: \mathrm{G}-\mathscr{E} n s \rightarrow \mathrm{H}-\mathscr{E} n s$.

Proposition If $f^{*}$ is an equivalence of categories, then $f$ is an isomorphism.

## Proof

(a) $f^{*}$ essentially surjective $\Rightarrow f$ injective. For any set X , all actions of H on X stems from an action of G , and so the elements of $\operatorname{Ker} f$ act trivially. Taking $\mathrm{X}=\mathrm{H}$ and making H act by left translation, $\operatorname{Ker} f=\{e\}$ follows.
(b) $f^{*}$ fully faithful $\Rightarrow f$ surjective. For all G-sets X and Y , every H-morphism from X to Y is a G-morphism. Take $\mathrm{X}=\{x\}$ and $\mathrm{Y}=\mathrm{G} / f(\mathrm{H})$. The map $\varphi: \mathrm{X} \rightarrow \mathrm{Y}$ such that $\varphi(\mathrm{X})$ is the class of the identity $e$ is a H-morphism, hence a G-morphism, and $e$ is fixed by G. As G acts transitively on $\mathrm{G} / f(\mathrm{H})$,

$$
\mathrm{G} / f(\mathrm{H})=\{e\}, \text { and } f(\mathrm{H})=\mathrm{G}
$$

### 4.7.2

If $\mathscr{C}_{1}, \mathscr{C}_{2}, \mathscr{D}$ are categories, $\mathrm{R}_{1}: \mathscr{C}_{1} \rightarrow \mathscr{D}$ and $\mathrm{R}_{2}: \mathscr{C}_{2} \rightarrow \mathscr{D}$ functors, define a category $\mathscr{C}=\mathscr{C}_{1} \times \mathscr{D} \mathscr{C}_{2}$ as follows: an object of $\mathscr{C}$ is a triplet $\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right)$ where $\mathrm{X}_{1} \in \mathscr{C}_{1}, \mathrm{X}_{2} \in \mathscr{C}_{2}$, and $f: \mathrm{R}_{1}\left(\mathrm{X}_{1}\right) \rightarrow \mathrm{R}_{2}\left(\mathrm{X}_{2}\right)$ is an isomorphism of $\mathscr{D}$. A morphism from $\mathscr{C}$ to $\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right)$ in $\left(\mathrm{Y}_{1}, \mathrm{Y}_{2}, g\right)$ is a pair $\left(h_{1}, h_{2}\right)$, with $h_{1}: \mathrm{X}_{1} \rightarrow \mathrm{Y}_{1}$ a morphism in $\mathscr{C}_{1}, h_{2}: \mathrm{X}_{2} \rightarrow \mathrm{Y}_{2}$ a morphism in $\mathscr{\mathscr { C }}_{2}$, and with commuting diagram


If

is a commutative diagram of functors, assign to it a functor $\Phi: \mathscr{C}_{1} \times \mathscr{D} \mathscr{C}_{2} \rightarrow$ $\mathscr{C}_{1}^{\prime} \times \mathscr{D}^{\prime} \mathscr{C}_{2}^{\prime}$. If $\Phi_{1}, \Phi_{2}$ are $\psi$ equivalences of categories, then so is $\Phi$.

### 4.7.3

Let $B$ be a locally connected space, $\mathrm{U}_{1}, \mathrm{U}_{2}$ subspaces of B , both open or both closed. Assume that $\mathrm{B}=\mathrm{U}_{1} \cup \mathrm{U}_{2}$ and set $\mathrm{V}=\mathrm{U}_{1} \cap \mathrm{U}_{2}$ (if both $\mathrm{U}_{i}$ are open, then they are locally connected and so is $V$, if they are closed, then $U_{1}, U_{2}$, and $V$ have to be assumed to be locally connected).

Proposition The functor $\alpha: \mathscr{C} \bullet \boldsymbol{v}_{\mathrm{B}} \rightarrow \mathscr{C} \bullet v_{\mathrm{U}_{1}} \times \mathscr{C} \bullet \boldsymbol{v}_{\mathrm{V}} \mathscr{C} \bullet \boldsymbol{v}_{\mathrm{U}_{2}}$ defined by $\alpha(\mathrm{X})=$ $\left(\left.\mathrm{X}\right|_{\mathrm{U}_{1}},\left.\mathrm{X}\right|_{\mathrm{U}_{2}},\left.1_{\mathrm{X}}\right|_{\mathrm{V}}\right)$ is an equivalence of categories.

Proof The functor $\alpha$ is fully faithful. Indeed, if X and Y are coverings of B and if $f_{1}:\left.\left.X\right|_{\mathrm{U}_{1}} \rightarrow \mathrm{Y}\right|_{\mathrm{U}_{1}}, f_{2}:\left.\left.\mathrm{X}\right|_{\mathrm{U}_{2}} \rightarrow \mathrm{Y}\right|_{\mathrm{U}_{2}}$ are morphisms agreeing over V , then gluing together $f_{1}$ and $f_{2}$ gives a morphism $f$ from X to Y . (When $\mathrm{U}_{1}$ and $\mathrm{U}_{2}$ are closed, the continuity of $f$ follows from (4.1.11, Lemma 4.3).)

The functor $\alpha$ is essentially surjective. Let $\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right) \in \mathscr{C} \boldsymbol{v}_{\mathrm{U}_{1}} \times \mathscr{C} \boldsymbol{a} \boldsymbol{u}_{\mathrm{V}} \mathscr{C} \not \boldsymbol{v}_{\mathrm{U}_{2}}$. Denote by $X$ the space over $B$ consisting of the quotient of the disjoint union $X_{1} \sqcup \mathrm{X}_{2}$
by the equivalence relation identifying $x$ with $f(x)$ for $\left.x \in \mathrm{X}_{1}\right|_{\mathrm{V}}$. The canonical injection $\iota_{1}: \mathrm{X}_{1} \rightarrow \mathrm{X}$ is a homeomorphism from $\mathrm{X}_{1}$ onto $\left.\mathrm{X}\right|_{\mathrm{U}_{1}}$. Indeed the canonical injection $X_{1} \rightarrow X_{1} \sqcup X_{2}$ is clopen, and the canonical map $X_{1} \sqcup X_{2} \rightarrow X$ is open if both $U_{i}$ are open, and closed if they are both closed. Similarly for $\iota_{2}$.

We show that $X$ is a covering of $B$. This is obvious if both $U_{i}$ are open. Assume they are closed and let $b \in \mathrm{~B}$. If $b \in \mathrm{~B}-\mathrm{V}, b$ is an interior point of one of the $\mathrm{U}_{i}$, and X is a covering over a neighbourhood of $b$. Assume $b \in \mathrm{~V}$. Since V is locally connected, there is a neighbourhood S of $b$ in B such that $\mathrm{T}=\mathrm{S} \cap \mathrm{V}$ is connected, and such that $\left.\mathrm{X}_{1}\right|_{\mathrm{s}_{1}}$ and $\left.\mathrm{X}_{2}\right|_{\mathrm{s}_{2}}$ are trivial, with $\mathrm{S}_{i}=\mathrm{S} \cap \mathrm{U}_{i}$. Let $\tau_{i}:\left.\mathrm{X}_{i}\right|_{\mathrm{s}_{i}} \rightarrow \mathrm{~S}_{i} \times \mathrm{F}_{i}$ be trivializations. We may assume that $\mathrm{F}_{1}=\mathrm{F}_{2}=\mathrm{F}$ and that the diagram

commutes within to identification of $\mathrm{F}_{1}$ with $\mathrm{F}_{2}$ by $\tau_{2, b} \circ f_{b} \circ \tau_{1, b}^{-1}$. Then the diagram

commutes. Indeed $\tau_{2} \circ f$ and $\tau_{1}$ are morphisms from $\left.\mathrm{X}_{1}\right|_{\mathrm{T}}$ to $\mathrm{T} \times \mathrm{F}$ agreeing over $b$, hence at least at one point in each of the connected components (4.3.5). So gluing together $\tau_{1}$ and $\tau_{2}$ gives a trivialization $\tau: \mathrm{X}_{\mathrm{S}} \rightarrow \mathrm{S} \times \mathrm{F}$; the continuity of $\tau$ and of its inverse follows from (4.1.11, Lemma 4.3).

Then $\mathrm{X} \in \mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{B}}$ and $\alpha(\mathrm{X}) \approx\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right)$.

### 4.7.4 Free Groups on Sets

Let X be a set. The free product $\mathrm{L}(\mathrm{X})=\star_{x \in \mathrm{X}} \mathrm{L}_{x}$ (2.5.2, Example 3), where, for all $x \in \mathrm{X}, \mathrm{L}_{x}=\mathbb{Z}$, is called the free group on X . For all $x \in \mathrm{X}$, set $e_{x}=\iota_{x}(1) \in \mathrm{L}_{x}$. The group $\mathrm{L}(\mathrm{X})$, equipped with the family $\left(e_{x}\right)_{x \in \mathrm{X}}$, represents the functor $\mathrm{G} \rightarrow \mathrm{G}^{\mathrm{X}}$, in other words:

UNIVERSAL PROPERTY. For every group G and family $\left(g_{x}\right)_{x \in \mathrm{X}}$ of elements in G , there is a unique homomorphism $\varphi: \mathrm{L}(\mathrm{X}) \rightarrow \mathrm{G}$ such that, for all $x \in \mathrm{X}, \varphi\left(e_{x}\right)=g_{x}$. Indeed

$$
\operatorname{Hom}_{\mathscr{G}_{\boldsymbol{r}}}(\mathrm{L}(\mathrm{X}) ; \mathrm{G})=\prod_{x \in \mathrm{X}} \operatorname{Hom}_{\mathscr{G}_{\mu}}\left(\mathrm{L}_{x} ; \mathrm{G}\right)=\mathrm{G}^{\mathrm{X}}
$$

The homomorphism $\varphi$ is surjective if and only if $\left(g_{x}\right)_{x \in X}$ generates G. If $\varphi$ is bijective (resp. injective), $\left(g_{x}\right)_{x \in \mathrm{X}}$ is said to be a basis (resp. a free family), and G a free group if it has a basis.

### 4.7.5

We keep the notation of 4.7.3. Let $b_{0} \in \mathrm{~V},\left(\mathrm{~V}_{\lambda}\right)_{\lambda \in \Lambda}$ the family of connected components of V , with $\mathrm{V}_{0}$ the connected component of $b_{0}$, and $\Lambda^{*}=\Lambda-\{0\}$. Assume that there is a universal covering of $B$.

Theorem If $\mathrm{U}_{1}$ and $\mathrm{U}_{2}$ are simply connected, then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is isomorphic to the free group over $\Lambda^{*}$.

More precisely, let X be a covering of B . Set $\mathrm{F}=\mathrm{X}\left(b_{0}\right)$ and let $\tau_{1}:\left.\mathrm{X}\right|_{\mathrm{U}_{1}} \rightarrow$ $\mathrm{U}_{1} \times \mathrm{F}$ and $\tau_{2}:\left.\mathrm{X}\right|_{\mathrm{U}_{2}} \rightarrow \mathrm{U}_{2} \times \mathrm{F}$ be the trivializations inducing the identity over $b_{0}$. Over $\mathrm{V}, \tau_{2} \circ \tau_{1}^{-1}$ is an automorphism of $\mathrm{V} \times \mathrm{F}$ defining a permutation $g_{b}^{\mathrm{X}}$ of F for all $b \in \mathrm{~V}$. This permutation only depends on the connected component of $b$ in V , and we get a family $\left(g_{\lambda}^{\mathrm{X}}\right)_{\lambda \in \Lambda}$ of permutations of F . For all $\lambda \in \Lambda$, as X varies, $g_{\lambda}^{\mathrm{X}}$ gives an automorphism $g_{\lambda}$ of the functor $\mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$ and hence an element of $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$. Obviously $g_{0}=e$. We prove that:

Proposition Under the assumptions of the theorem, $\left(g_{\lambda}\right)_{\lambda \in \Lambda^{*}}$ is a basis for $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ as a (not necessarily commutative) group.

Proof Consider the diagram of functors

where L is the free group over $\Lambda^{*}, \gamma: \mathrm{L} \rightarrow \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ the homomorphism defined by $\gamma\left(e_{\lambda}\right)=g_{\lambda}, \mathscr{C}$ the category of pairs $\left(\mathrm{F},\left(f_{\lambda}\right)_{\lambda \in \Lambda^{*}}\right)$, where F is a set and $f_{\lambda}$ a permutation of F for all $\lambda, \mathscr{D}$ the category of triplets $\left(\mathrm{F}_{1}, \mathrm{~F}_{2},\left(f_{\lambda}\right)_{\lambda \in \mathrm{V}}\right)$, where $\mathrm{F}_{1}$ and $\mathrm{F}_{2}$ are sets and $f_{\lambda}: \mathrm{F}_{1} \rightarrow \mathrm{~F}_{2}$ a bijection for $\lambda$; the functor $\delta$ assigns to a L-set F the set F equipped with the permutations defined by $e_{\lambda}$, the functor $\varepsilon$ assigns to $\left(\mathrm{F},\left(f_{\lambda}\right)\right.$ ) the object $\left(\mathrm{F}, \mathrm{F},\left(f_{\lambda}\right)\right.$ ), with $f_{0}=1_{\mathrm{F}}$. The functor $\alpha$ is the one defined in 4.7.3, and the functor $\beta$ is defined as follows:

$$
\beta\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right)=\left(\mathrm{F}_{1}, \mathrm{~F}_{2},\left(f_{\lambda}\right)\right), \text { where } \mathrm{F}_{1}=\mathrm{X}_{1}\left(b_{0}\right), \mathrm{F}_{2}=\mathrm{X}_{2}\left(b_{0}\right),
$$

$f_{\lambda}$ is the map $x \mapsto \tau_{2} \circ \tau_{1}^{-1}(b, x)$ for $b \in \mathrm{~V}_{\lambda}$, where $\tau_{1}: \mathrm{X}_{1} \rightarrow \mathrm{U}_{1} \times \mathrm{F}_{1}$ and $\tau_{2}$ : $\mathrm{X}_{2} \rightarrow \mathrm{U}_{2} \times \mathrm{F}_{2}$ are the trivializations inducing the identity over $b_{0}$.

This diagram is easily seen to be commutative. By Proposition (4.7.3), the functor $\alpha$ is an equivalence of categories. That $\delta$ and $\varepsilon$ are equivalences follows readily, and it is easily seen that so is $\beta$. As a result, $\gamma^{*} \circ \Phi$ is an equivalence of categories.

The functor $\Phi$ is an equivalence by (4.6.8), and hence so is $\gamma^{*}$, and by Proposition (4.7.1) $\gamma$ is an isomorphism.

### 4.7.6 Amalgamated Sums

Let $\mathrm{G}_{1}, \mathrm{G}_{2}$, H be groups, $u_{1}: \mathrm{H} \rightarrow \mathrm{G}_{1}$ and $u_{2}: \mathrm{H} \rightarrow \mathrm{G}_{2}$ homomorphisms. The sum of $\mathrm{G}_{1}$ and $\mathrm{G}_{2}$ amalgamated by H , written $\mathrm{G}_{1} \stackrel{H}{\star} \mathrm{G}_{2}$, is the quotient $\mathrm{Q}=\left(\mathrm{G}_{1} \star \mathrm{G}_{2}\right) / \mathrm{N}$, where $G_{1} \star G_{2}$ is the summation group of $G_{1}$ and $G_{2}$ and $N$ the normal subgroup generated by $u_{1}(h) u_{2}\left(h^{-1}\right), h \in \mathrm{H}$. Equip Q with the homomorphisms $j_{1}: \mathrm{G}_{1} \rightarrow \mathrm{Q}$ and $j_{2}: \mathrm{G}_{2} \rightarrow \mathrm{Q}$ composed of the canonical injections $\mathrm{G}_{1} \rightarrow \mathrm{G}_{1} \star \mathrm{G}_{2}$ and $\mathrm{G}_{2} \rightarrow$ $\mathrm{G}_{1} \star \mathrm{G}_{2}$ and the canonical quotient map.

UNIVERSAL PROPERTY.
(a) The diagram

is commutative.
(b) Let $\Gamma$ be a group, $f_{1}$ and $f_{2}$ homomorphisms such that the diagram

commutes. Then there is a unique homomorphism $g: \mathrm{Q} \rightarrow \Gamma$ such that $g \circ j_{1}=f_{1}$ and $g \circ j_{2}=f_{2}$.

## Proof

(a) For all $h \in \mathrm{H}, u_{1}(h) u_{2}\left(h^{-1}\right) \in \mathrm{N}$, and so $j_{1} \circ u_{1}=j_{2} \circ u_{2}$.
(b) Define $g: \mathrm{G}_{1} \star \mathrm{G}_{2} \rightarrow \Gamma$ by $g(x)=f_{1}(x)$ if $x \in \mathrm{G}_{1}$ and $g(x)=f_{2}(x)$ if $x \in \mathrm{G}_{2}$. For $h \in \mathrm{H}, g\left(u_{1}(h) u_{2}\left(h^{-1}\right)\right)=f_{1} \circ u_{1}(h) f_{2} \circ u_{2}\left(h^{-1}\right)=e$, where $e$ is the identity element of $\Gamma$. Hence there is a homomorphism $f$ from Q to $\Gamma$. As $\mathrm{G}_{1}$ and $\mathrm{G}_{2}$
generate $\mathrm{G}_{1} \star \mathrm{G}_{2}$, their images $j_{1}\left(\mathrm{G}_{1}\right)$ and $j_{2}\left(\mathrm{G}_{2}\right)$ generate Q and the homomorphism $f: \mathrm{Q} \rightarrow \Gamma$ thus constructed is the only one with the desired property.

The universal property is equivalent to the statement that $G_{1} \stackrel{H}{\star} G_{2}$, equipped with $\left(j_{1}, j_{2}\right)$ represents the functor

$$
\Gamma \mapsto \operatorname{Hom}\left(\mathrm{G}_{1}, \Gamma\right) \times \begin{aligned}
& \operatorname{Hom}(\mathrm{H}, \Gamma) \\
& \operatorname{Hom}\left(\mathrm{G}_{2}, \Gamma\right)
\end{aligned}
$$

### 4.7.7

We keep the notation of 4.7.3. Let $b_{0} \in \mathrm{~B}$. Assume that there are universal pointed coverings of $\left(\mathrm{B}, b_{0}\right),\left(\mathrm{U}_{1}, b_{0}\right),\left(\mathrm{U}_{2}, b_{0}\right)$ and of $\left(\mathrm{V}, b_{0}\right)$.

Theorem (Van Kampen) If V is connected, then $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is the sum of $\pi_{1}\left(\mathrm{U}_{1}, b_{0}\right)$ and $\pi_{1}\left(\mathrm{U}_{2}, b_{0}\right)$ amalgamated by $\pi_{1}\left(\mathrm{~V}, b_{0}\right)$.

More precisely, set $\mathrm{G}=\pi_{1}\left(\mathrm{~B}, b_{0}\right), \mathrm{G}_{i}=\pi_{1}\left(\mathrm{U}_{i}, b_{0}\right)$ and $\mathrm{H}=\pi_{1}\left(\mathrm{~V}, b_{0}\right)$. The canonical inclusions define a commutative diagram.


The theorem states that the homomorphism $v: \mathrm{G}_{1} \stackrel{\mathrm{H}}{\star} \mathrm{G}_{2} \rightarrow \mathrm{G}$ defined by $v_{1}$ and $v_{2}$ is an isomorphism.

Proof Let $\mathrm{U}_{i}^{\prime}$ be the connected component of $b_{0}$ in $\mathrm{U}_{i}$ and $\mathrm{B}^{\prime}$ the connected component of $b_{0}$ in B . Then $\mathrm{U}_{1}^{\prime} \cup \mathrm{U}_{2}^{\prime}=\mathrm{B}^{\prime}$ and $\mathrm{U}_{1}^{\prime} \cap \mathrm{U}_{2}^{\prime}=\mathrm{V}$ and the fundamental groups remain the same when $\mathrm{U}_{i}$ is replaced by $\mathrm{U}_{i}^{\prime}$ and B by $\mathrm{B}^{\prime}$. Hence B and both $\mathrm{U}_{i}$ may be assumed to be connected.

There is a commutative diagram of functors

where $\quad \theta(\mathrm{X})=\mathrm{X}\left(b_{0}\right), \quad \Phi\left(\mathrm{X}_{1}, \mathrm{X}_{2}, f\right)=\left(\mathrm{X}_{1}\left(b_{0}\right), \mathrm{X}_{2}\left(b_{0}\right), f_{b_{0}}\right), \quad \alpha(\mathrm{X})=\left(\mathrm{X}_{\mathrm{U}_{1}}\right.$, $\left.\mathrm{X}_{\mathrm{U}_{2}}, \mathrm{I}\right), \beta(\mathrm{F})=(\mathrm{F}, \mathrm{F}, \mathrm{I}), \gamma(\mathrm{F})=(\mathrm{F}, \mathrm{F}, \mathrm{I})$. By Theorem (4.6.8), in this diagram, $\theta$ is an equivalence of categories, and so is $\Phi$. Clearly, this is also the case of $\gamma$. By Proposition 4.7.3, $\alpha$ is an equivalence of categories. Hence so is $v^{*}$, and by Proposition 4.7.1 $v$ is an isomorphism.

## Exercises 4.7. (Van Kampen's theorem)

1.-Let $\mathrm{G}_{1}, \mathrm{G}_{2}$, H be groups, $u_{1}: \mathrm{H} \rightarrow \mathrm{G}_{1}$ and $u_{2}: \mathrm{H} \rightarrow \mathrm{G}_{2}$ homomorphisms. Show that the amalgamated sum $G=G_{1} \stackrel{H}{\star} G_{2}$ is generated by the union of the images of $\mathrm{G}_{1}$ and $\mathrm{G}_{2}$.
2.-Let $S^{n}$ be the unit sphere in $\mathbb{R}^{n+1}, S_{+}^{n}$ the upper hemisphere, i.e.

$$
S_{+}^{n}=\left\{x=\left(x_{1}, \ldots, x_{n+1}\right) \in S^{n} \mid x_{n+1} \geqslant 0\right\}
$$

and $S_{-}^{n}$ the lower one. Let X be a fibre bundle over $\mathrm{S}^{n}$. Assume that the fibres $\mathrm{X}_{+}$and $X_{-}$induced on $S_{+}^{n}$ and $S_{-}^{n}$ are trivial (this condition can be shown to be necessarily satisfied). Set $s_{0}=(1,0, \ldots, 0), \mathrm{F}=\mathrm{X}\left(s_{0}\right)$ and let $x_{0} \in \mathrm{~F}$.

The aim is to find $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ assuming known $\pi_{1}\left(\mathrm{~F}, x_{0}\right)$.
(a) Find $\pi_{1}\left(\mathrm{X}_{+}, x_{0}\right)$ and $\pi_{1}\left(\mathrm{X}_{-}, x_{0}\right)$ as well as $\pi_{1}\left(\mathrm{X}_{0}, x_{0}\right)$, where

$$
\mathrm{X}_{0}=\mathrm{X}_{+} \cap \mathrm{X}_{-}
$$

(consider the cases $n=2$ and $n \geqslant 3$ separately).
(b) Find $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ when $n \geqslant 3$.
c) Show that, if $n=2$, then the group $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ can be identified with the quotient of $\pi_{1}\left(\mathrm{~F}, x_{0}\right)$ by a cyclic subgroup contained in the centre.
3.-Let $\mathrm{SO}_{n}$ be the group of direct linear isometries of $\mathbb{R}^{n}$. The aim is to find $\pi_{1}\left(\mathrm{SO}_{n}\right)$.
(a) Show that $\mathrm{SO}_{2}$ is homeomorphic to $\mathrm{S}^{1}$.
(b) Let $\mathbf{H}$ denote the field of quaternions: it is a 4-dimensional vector space over $\mathbb{R}$ with basis $(1, i, j, k)$, and bilinear multiplication defined by

$$
i^{2}=j^{2}=k^{2}=-1, i j=-j i=k, j k=-k j=i, k i=-i k=j
$$

This multiplication is associative, but not commutative. For

$$
x=r+a i+b j+c k, \text { set }\|x\|=\sqrt{r^{2}+a^{2}+b^{2}+c^{2}} ;
$$

Then $\|x y\|=\|x\| \cdot\|y\|$, and as a result, $\mathbf{H}$ is a skew field. Let G be the group of the quaternions having norm 1 . Show that for all $x \in \mathrm{G}$, the inner automorphism of $\mathbf{H}$ defined by $x$ induces an isometry of the vector subspace of $\mathbf{H}$ generated by $(i, j, k)$. Deduce a homomorphism from G to $\mathrm{SO}_{3}$. Show that it is surjective. Give its kernel.
(c) Show that $\pi_{1}\left(\mathrm{SO}_{3}\right)=\mathbb{Z} /(2)$.
(d) Show that, for all $n \geqslant 3, \pi_{1}\left(\mathrm{SO}_{n}\right)=\mathbb{Z} /(2)$ (use Exercise 2). Deduce that for all $n \geqslant 2$, there is a topological group $\operatorname{Spin}_{n}$, unique up to isomorphism, which is a connected 2-fold covering of $\mathrm{SO}_{n}$.
4.-Show that in Theorem 4.7.7, the assumption that there is a universal covering of ( $\mathrm{B}, b_{0}$ ) is a consequence of the other assumptions.
5.-With the notation of 4.7.6, assume that $u_{1}$ and $u_{2}$ are injective. Show that $j_{1}$ and $j_{2}$ are injective.
6.-Let G and H be two groups.
(a) Show that, if the categories G- $\mathscr{E} n s$ and $\mathrm{H}-\mathscr{E} n s$ are equivalent, then the groups G and H are isomorphic.
(b) Are all equivalences from $\mathrm{H}-\mathscr{E} n s$ to $\mathrm{G}-\mathscr{E} n s$ of the form $f^{*}$, where $f: \mathrm{G} \rightarrow \mathrm{H}$ is a homomorphism?
7.-Let E be the set of pairs $(x, y)$ of orthonormal vectors of $\mathbb{R}^{3}$, i.e. with $\|x\|=$ $\|y\|=1,(x \cdot y)=0$. Equip E with the projection $\pi:(x, y) \mapsto x$ from E onto $\mathrm{S}^{2}$.
(a) Show that $\pi$ is a fibration with fibre $\mathrm{S}^{1}$.
(b) Show that E is homeomorphic to $\mathrm{SO}_{3}$.
(c) Assuming there is a continuous section $\sigma: \mathrm{S}^{2} \rightarrow \mathrm{E}$ of $\pi$, show that E is homeomorphic to $S^{2} \times S^{1}$. With this assumption, determine $\pi_{1}(\mathrm{E})$. Compare with Exercise 3 , (c). Conclude.
(d) A continuous tangent vector field on $S^{2}$ is a continuous map $\xi: S^{2} \rightarrow \mathbb{R}^{3}$ such that, for all $x \in \mathrm{~S}^{2}, \xi(x)$ is tangent to $\mathrm{S}^{2}$ at $x$, i.e. orthogonal to $x$. Show that there is a point of $S^{2}$ at which the continuous tangent vector field vanishes.
(e) A continuous normal tangent vector field on $S^{2}$ is a continuous map $\nu: \mathrm{S}^{2} \rightarrow \mathbb{R}^{3}$ such that, for all $x \in \mathrm{~S}^{2}, \nu(x)$ is normal to $\mathrm{S}^{2}$ at $x$, i.e. proportional to $x$. Consider the space $\mathscr{T}$ of continuous tangent vector fields and the space $\mathscr{N}$ of continuous normal tangent vector fields as modules over the ring A of continuous functions $\mathrm{S}^{2} \rightarrow \mathbb{R}$. Show that $\mathscr{N}$ is free of rank 1 and that $\mathscr{N} \oplus \mathscr{T}$ is free of rank 3. Deduce that $\mathscr{T}$ is projective. Show that $\mathscr{T}$ is not free.

### 4.8 Graphs. Subgroups of Free Groups

### 4.8.1 Graphs

A simple graph is a pair $(\mathrm{S}, \mathrm{A})$, where S is a set and A a subset of $\mathrm{S} \times \mathrm{S}$ defining a reflexive and symmetric relation. In this chapter, we will simply use the term graph. In chapter 7 we will give a broader definition. The points of $S$ are called vertices
of the graph. Two vertices $s$ and $s^{\prime}$ are said to be connected if $\left(s, s^{\prime}\right) \in \mathrm{A}$. If (S, A) is a graph, then its geometric realization is defined as follows: consider the vector space $\mathbb{R}^{(S)}$ equipped with the norm $x \mapsto \sum\left|x_{s}\right|$, and embed S in $\mathbb{R}^{(\mathrm{S})}$ by identifying $s$ with the element of the canonical basis having index $s$. For $s, t \in \mathrm{~S}$, let $[s, t]$ be the segment with endpoints $s$ and $t$, i.e. the set of $(1-\theta) s+\theta t$ for $\theta \in[0,1]$. The geometric realization of $(\mathrm{S}, \mathrm{A})$ is the space

$$
|\mathrm{S}, \mathrm{~A}|=\bigcup_{(s, t) \in \mathrm{A}}[s, t]
$$

equipped with the distance induced by the norm.
The segments $[s, t]$ with $(s, t) \in \mathrm{A}$ are called the edges. If $s, t, s^{\prime}, t^{\prime}$ are distinct points,

$$
[s, t]=\{s\} \text { if } t=s,[s, t]=[t, s],[s, t] \cap\left[s, t^{\prime}\right]=\{s\} \text { if } t \neq t^{\prime}
$$

and

$$
[s, t] \cap\left[s^{\prime}, t^{\prime}\right]=\varnothing
$$

More precisely, if

$$
[s, t] \cap\left[s^{\prime}, t^{\prime}\right]=\varnothing,
$$

then

$$
(\forall x \in[s, t])\left(\forall x^{\prime} \in\left[s^{\prime}, t^{\prime}\right]\right)\left\|x-x^{\prime}\right\|=2
$$

For $s \in \mathrm{~S}$, the union of the edges of the graph having $s$ as an endpoint is called the star of $s$. It is a contractible neighbourhood of $s$ in $|\mathrm{S}, \mathrm{A}|$, hence simply connected (4.3.17, Corollary 4.6).

Let ( $\mathrm{S}, \mathrm{A}$ ) and ( $\mathrm{S}^{\prime}, \mathrm{A}^{\prime}$ ) be two graphs. A morphism from ( $\mathrm{S}, \mathrm{A}$ ) to $\left(\mathrm{S}^{\prime}, \mathrm{A}^{\prime}\right)$ is a map $f$ from S to $\mathrm{S}^{\prime}$ such that $f \times f: \mathrm{S} \times \mathrm{S} \rightarrow \mathrm{S}^{\prime} \times \mathrm{S}^{\prime}$ maps A to $\mathrm{A}^{\prime}$. Then there is continuous map $f_{*}$ from $|\mathrm{S}, \mathrm{A}|$ to $\left|\mathrm{S}^{\prime}, \mathrm{A}^{\prime}\right|$ corresponding to $f$. For $\mathrm{S} \subset \mathrm{S}^{\prime}$ and $\mathrm{A} \subset \mathrm{A}^{\prime}$, $(S, A)$ is said to be a subgraph of $\left(S^{\prime}, A^{\prime}\right)$; then $|S, A|$ is a subspace of $\left|S^{\prime}, A^{\prime}\right|$ with the induced metric. Let ( $\mathrm{S}, \mathrm{A})$ and (T, B) be two graphs, $\left(\mathrm{S}^{\prime}, \mathrm{A}^{\prime}\right)$ a subgraph of ( $\mathrm{S}, \mathrm{A}$ ) and $f:(\mathrm{T}, \mathrm{B}) \rightarrow(\mathrm{S}, \mathrm{A})$ a morphism. Set $\left(\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right)=f^{-1}\left(\mathrm{~S}^{\prime}, \mathrm{A}^{\prime}\right)$; then

$$
\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|=f_{*}^{-1}\left(\left|\mathrm{~S}^{\prime}, \mathrm{A}^{\prime}\right|\right)
$$

If $G$ is a graph, then its geometric realization $|\mathrm{G}|$ is a locally connected space, and every connected component of $|\mathrm{G}|$ is the geometric realization of a subgraph of G .

### 4.8.2

Proposition and Definition Let (S, A) and (T, B) be two graphs and $\rho:(\mathrm{T}, \mathrm{B}) \rightarrow$ $(\mathrm{S}, \mathrm{A})$ a morphism. $|\mathrm{T}, \mathrm{B}|$ together with $\rho_{*}$ is a covering of $|\mathrm{S}, \mathrm{A}|$ if and only if

$$
\left(\forall\left(x, x^{\prime}\right) \in \mathrm{A}\right)(\forall y \in \mathrm{~T}, \rho(y)=x)\left(\exists y^{\prime} \in \mathrm{T}\right) \rho\left(y^{\prime}\right)=x^{\prime} \text { et }\left(y, y^{\prime}\right) \in \mathrm{B}
$$

If these conditions hold, then (T, B) together with $\rho$ is said to be a covering of (S, A).
Proof Suppose it is a covering. The connected component of $y$ in $\rho_{*}^{-1}\left(\left[x, x^{\prime}\right]\right)$ is of the form $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$, where $\left(\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right)$ is a subgraph of $(\mathrm{T}, \mathrm{B})$, and $\rho_{*}$ induces a homeomorphism from $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$ onto $\left[x, x^{\prime}\right]$. In particular, $\rho$ injects $\mathrm{T}^{\prime}$ into $\left\{x, x^{\prime}\right\}$, and so $\mathrm{T}^{\prime}$ has at most two points. If $x \neq x^{\prime}, \mathrm{T}^{\prime}$ has two points since $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$ is homeomorphic to $[0,1]$; these are $y$ and a point $y^{\prime}$ with the desired property. If $x=x^{\prime}$, then $y^{\prime}=y$.

Conversely, let $\mathrm{T}_{x}$ be the inverse image of $x$ in T. Two distinct vertices $y_{1}$ and $y_{2}$ of $\mathrm{T}_{x}$ are never connected, for otherwise $\left(y_{1}, y_{1}\right)$ and $\left(y_{2}, y_{2}\right)$ would be two coverings distinct from $(x, x)$.

The union of all $\left[x, x^{\prime}\left[\right.\right.$ with $x^{\prime}$ is connected $x$ is called the open star of $x$ and written $\mathrm{U}_{x}$. Then $\rho^{-1}\left(\mathrm{U}_{x}\right)=\bigcup_{y \in \mathrm{~T}_{x}} \mathrm{U}_{y}$. Indeed, let $z \in|\mathrm{~T}, \mathrm{~B}|$ be such that $\rho(z) \in\left[x, x^{\prime}[\right.$. If $z$ is a vertex, then $\rho(z)=x$; if $z \in] y, y^{\prime}\left[\right.$, one of the points $y, y^{\prime}$, say $y$, projects onto $x$, and $z \in \mathrm{U}_{y}$ and $y \in \mathrm{~T}_{x}$.

For all $y \in \mathrm{~T}_{x}$, the assumption implies that $\rho$ induces an isomorphism from the star of $y$ onto the star of $x$. The same holds for open stars.

Let $y_{1}$ and $y_{2}$ be two distinct points of $\mathrm{T}_{x}$. Then the stars of $y_{1}$ and $y_{2}$ are disjoint. Indeed, $y_{1}$ and $y_{2}$ are not connected, and if there was a vertex $z$ connected to $y_{1}$ and to $y_{2}$, then $\left(z, y_{1}\right)$ and $\left(z, y_{2}\right)$ would be two coverings of $(\rho(z), x)$. Moreover, the distance of these stars is 2 .

Hence $\rho^{-1}\left(\mathrm{U}_{x}\right)$ is a trivial covering of $\mathrm{U}_{x}$.

### 4.8.3

Proposition Let (S, A) be a graph. The functor $((\mathrm{T}, \mathrm{B}), \rho) \mapsto\left(|\mathrm{T}, \mathrm{B}|, \rho_{*}\right)$ defines an equivalence from the category of coverings of $(\mathrm{S}, \mathrm{A})$ onto the category of coverings of $|\mathrm{S}, \mathrm{A}|$.

## Proof

(a) The functor is essentially surjective. Set $\mathrm{X}=|\mathrm{S}, \mathrm{A}|$ and let Y be a covering of $X$. Denote by $p$ the projection of $Y$ onto $X$, and set $S^{\prime}=p^{-1}(S)$. Define $A^{\prime} \subset S^{\prime} \times S^{\prime}$ by $\left(s^{\prime}, t^{\prime}\right) \in \mathrm{A}^{\prime} \Leftrightarrow(s, t) \in \mathrm{A}$ and $s^{\prime}$ and $t^{\prime}$ are in the same connected component of $p^{-1}\left[s, t\left[\right.\right.$, where $s=p\left(s^{\prime}\right)$ and $t=p\left(t^{\prime}\right)$. Set $\mathrm{X}^{\prime}=\left|\mathrm{S}^{\prime}, \mathrm{A}^{\prime}\right|$. The map $p$ induces a map from $S^{\prime}$ to $S$. It is a morphism from $\left(S^{\prime}, A^{\prime}\right)$ to $(S, A)$, hence there is a map $p_{*}: \mathrm{X}^{\prime} \rightarrow \mathrm{X}$.

Define $\varphi: \mathrm{X}^{\prime} \rightarrow \mathrm{Y}$ as follows: for $\left(s^{\prime}, t^{\prime}\right) \in \mathrm{A}^{\prime}$, set $s=p\left(s^{\prime}\right)$ and $t=p\left(t^{\prime}\right)$, then $\left.\mathrm{Y}\right|_{[s, t]}$ is a trivial covering since $[s, t]$ is simply connected, and so there is a unique continuous section $\sigma$ of Y over $[s, t]$ such that $\sigma(s)=s^{\prime}$. Then $\sigma(t)=t^{\prime}$ by definition of $\mathrm{A}^{\prime}$. Set $\varphi_{s^{\prime}, t^{\prime}}=\left.\sigma \circ p_{*}\right|_{\left[s^{\prime}, t^{\prime}\right]}$. Gluing together all $\varphi_{s^{\prime}, t^{\prime}}$ gives a map $\varphi$ commuting with the projections onto X .

Show that $\varphi$ is a homeomorphism. As the interiors of the stars covering $X$, it suffices to show that for every star V in $\mathrm{X}, \varphi_{\mathrm{V}}:\left.\left.\mathrm{X}^{\prime}\right|_{\mathrm{V}} \rightarrow \mathrm{Y}\right|_{\mathrm{V}}$ is a homeomorphism. As V is contractible, hence simply connected, $\left.\mathrm{Y}\right|_{\mathrm{V}}$ is a trivial covering. Thus the question reduces to the case of a trivial covering.

It may then be assumed that $\mathrm{Y}=\mathrm{X} \times \mathrm{F}$ for some discrete set F . Equip Y with the distance defined by

$$
d((x, u),(y, u))=d(x, y) \text { and } d((x, u),(y, v))=2 \text { if } u \neq v
$$

As can be checked, $\varphi$ is an isometry, and Y is isomorphic to $\left|\mathrm{S}^{\prime}, \mathrm{A}^{\prime}\right|$.
(b) The functor is fully faithful. Let ( $\mathrm{T}, \mathrm{B}$ ) and ( $\mathrm{T}^{\prime}, \mathrm{B}^{\prime}$ ) be coverings of ( $\mathrm{S}, \mathrm{A}$ ) and $f:|\mathrm{T}, \mathrm{B}| \rightarrow\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$ a covering morphism of $|\mathrm{S}, \mathrm{A}|$. We show that there a unique morphism $g:(\mathrm{T}, \mathrm{B}) \rightarrow\left(\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right)$ inducing $f$. For $s \in \mathrm{~S}$, two distinct points of $p^{-1}(s)$ are never connected since $p_{*}^{-1}(s)$ is discrete. Hence $p_{*}^{-1}(\mathrm{~S})=\mathrm{T}$. Similarly $p_{*}^{\prime-1}(\mathrm{~S})=\mathrm{T}^{\prime}$, and $f(\mathrm{~T}) \subset \mathrm{T}^{\prime}$. Let $g: \mathrm{T} \rightarrow \mathrm{T}^{\prime}$ be the map induced by $f$.

Let $t_{1}$ and $t_{2}$ be two vertices of T ; set $s_{1}=p\left(t_{1}\right)$ and $s_{2}=p\left(t_{2}\right)$. Then, $t_{1}$ and $t_{2}$ are connected if and only if $s_{1}$ and $s_{2}$ are connected and $t_{1}$ and $t_{2}$ are in the same connected component in $p_{*}^{-1}\left(\left[s_{1}, s_{2}\right]\right)$. Hence, if $t_{1}$ and $t_{2}$ are connected, so are $f\left(t_{1}\right)$ and $f\left(t_{2}\right)$; hence $g$ is a morphism from ( $\mathrm{T}, \mathrm{B}$ ) to ( $\mathrm{T}^{\prime}, \mathrm{B}^{\prime}$ ). The covering (topological) morphisms $g_{*}$ and $f$ agree on the vertices, and hence at least at one point in each connected component of $|\mathrm{T}, \mathrm{B}|$; hence they are equal. Uniqueness is obvious.

### 4.8.4

Definition A tree is a graph whose geometric realization is simply connected.
Proposition Let $(\mathrm{S}, \mathrm{A})$ be a graph and $\left(\mathrm{S}_{i}, \mathrm{~A}_{i}\right)_{i \in \mathrm{I}}$ a nonempty increasing directed family of subgraphs. Assume that $(\mathrm{S}, \mathrm{A})$ is the union of $\left(\mathrm{S}_{i}, \mathrm{~A}_{i}\right)$, i.e. $\mathrm{S}=\bigcup \mathrm{S}_{i}$ and $\mathrm{A}=\bigcup \mathrm{A}_{i}$, and that $\left(\mathrm{S}_{i}, \mathrm{~A}_{i}\right)$ is a tree for all $i$. Then $(\mathrm{S}, \mathrm{A})$ is a tree.

Proof Let $s \in \mathrm{~S}$. If necessary replacing I by a cofinal set, we may assume that $s \in \mathrm{~S}_{i}$ for all $i$. Then $|\mathrm{S}, \mathrm{A}|$ is the union of the connected $\left|\mathrm{S}_{i}, \mathrm{~A}_{i}\right|$ containing $s$, and so is connected.

Let Y be a covering of $|\mathrm{S}, \mathrm{A}|$ and $t \in \mathrm{Y}(s)$. We may assume $\mathrm{Y}=|\mathrm{T}, \mathrm{B}|$ for some covering ( $\mathrm{T}, \mathrm{B}$ ) of $(\mathrm{S}, \mathrm{A})$. Then $t \in \mathrm{~T}$. For all $i$, there is a unique continuous section of $|\mathrm{T}, \mathrm{B}|$ over $\left|\mathrm{S}_{i}, \mathrm{~A}_{i}\right|$ passing through $t$. By (4.8.3), this section stems from a unique morphism $\varphi_{i}:\left(\mathrm{S}_{i}, \mathrm{~A}_{i}\right) \rightarrow(\mathrm{T}, \mathrm{B})$. By uniqueness, gluing together these
morphisms gives a morphism $\varphi:(\mathrm{S}, \mathrm{A}) \rightarrow(\mathrm{T}, \mathrm{B})$ and $\varphi_{*}:|\mathrm{S}, \mathrm{A}| \rightarrow|\mathrm{T}, \mathrm{B}|$ is a continuous section passing through $t$. By (4.3.11), Y is a trivial covering.

### 4.8.5

A graph is connected if so is its geometric realization.
Proposition Every connected graph (S, A) has a spanning tree, i.e. a subgraph ( $\mathrm{T}, \mathrm{B}$ ) with $\mathrm{T}=\mathrm{S}$.

Proof By the previous proposition, the set of trees in (S, A) is inductive. Hence, by Zorn's theorem, there is a maximal tree ( $\mathrm{T}, \mathrm{B}$ ). We show that $\mathrm{T}=\mathrm{S}$. If $\mathrm{S}-\mathrm{T} \neq \varnothing$, there exists $\left(t_{0}, t_{0}^{\prime}\right) \in \mathrm{A}$ with $t_{0} \in \mathrm{~T}$ and $t_{0}^{\prime} \in \mathrm{S}-\mathrm{T}$; otherwise ( $\mathrm{S}, \mathrm{A}$ ) would not be connected. Set

$$
\mathrm{T}^{\prime}=\mathrm{T} \cup\left\{t_{0}^{\prime}\right\} \text { and } \mathrm{B}^{\prime}=\mathrm{B} \cup\left\{\left(t_{0}, t_{0}^{\prime}\right),\left(t_{0}^{\prime}, t_{0}\right),\left(t_{0}^{\prime}, t_{0}^{\prime}\right)\right\}
$$

We show that $\left(\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right)$ is a tree, which will contradict maximality.
For $\theta \in[0,1]$ and $t \in \mathrm{~T}^{\prime}$, set

$$
h(\theta, t)=t \text { if } t \neq t_{0}^{\prime} \text { and } h\left(\theta, t_{0}^{\prime}\right)=(1-\theta) t_{0}+\theta t_{0}^{\prime} .
$$

For each value of $\theta$, extend $h$ to a continuous map from $[0,1] \times \mathbb{R}^{\left(\mathrm{T}^{\prime}\right)}$ to $\mathbb{R}^{\left(\mathrm{T}^{\prime}\right)}$ by linearity. This map can be checked to be continuous and to induce a homotopy (4.3.16) between a retraction from $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$ onto $|\mathrm{T}, \mathrm{B}|$ and the identity of $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$. Therefore $|T, B|$ is a deformation retract of $\left|\mathrm{T}^{\prime}, \mathrm{B}^{\prime}\right|$, of the same homotopy type (4.3.16) as $|\mathrm{T}, \mathrm{B}|$, and so is simply connected.

### 4.8.6

Theorem Let $(\mathrm{S}, \mathrm{A})$ be a graph and $s_{0} \in \mathrm{~S}$. Then $\pi_{1}\left(|\mathrm{~S}, \mathrm{~A}|, s_{0}\right)$ is a free group.
Replacing (S, A) by the connected component of $s_{0}$ if necessary, which leaves $\pi_{1}$ unchanged, ( $\mathrm{S}, \mathrm{A}$ ) can be assumed to be connected. Then, by the previous proposition, there is a tree $(\mathrm{S}, \mathrm{B})$. As B contains all $(s, s), \mathrm{A}-\mathrm{B}=\mathrm{C} \sqcup \mathrm{C}^{\prime}$, where C and $\mathrm{C}^{\prime}$ are mutually symmetric sets, i.e. $\mathrm{C}^{\prime}=\left\{\left(x^{\prime}, x\right) \mid\left(x, x^{\prime}\right) \in \mathrm{C}\right\}$. The theorem then follows from the more precise following proposition:

Proposition The group $\pi_{1}\left(|\mathrm{~S}, \mathrm{~A}|, s_{0}\right)$ has a basis indexed by C .
Proof Set

$$
\left.\mathrm{U}_{1}=|\mathrm{S}, \mathrm{~A}|-\bigcup_{\left(x, x^{\prime}\right) \in \mathrm{C}}\right] \frac{x+x^{\prime}}{2}, x^{\prime}\left[=|\mathrm{S}, \mathrm{~B}| \cup \bigcup_{\left(x, x^{\prime}\right) \in \mathrm{C}}\left[x, \frac{x+x^{\prime}}{2}\right]\right.
$$

and define $\mathrm{U}_{2}$ likewise by replacing C by $\mathrm{C}^{\prime}$. The sets $\mathrm{U}_{1}$ and $\mathrm{U}_{2}$ are closed in $|\mathrm{S}, \mathrm{A}|$ since $] \frac{x+x^{\prime}}{2}, x^{\prime}\left[\right.$ is open for all $\left(x, x^{\prime}\right)$, and their union is $|\mathrm{S}, \mathrm{A}|$.

The space $\mathrm{U}_{1}$ is of the same homotopy as $|\mathrm{S}, \mathrm{B}|$. Indeed, define a retraction $\rho: \mathrm{U}_{1} \rightarrow|\mathrm{~S}, \mathrm{~B}|$ by $\rho(x)=x$ for $x \in|\mathrm{~S}, \mathrm{~B}|$ and $\rho(x)=s$ for

$$
x \in\left[s, \frac{s+s^{\prime}}{2}\right], \quad\left(s, s^{\prime}\right) \in \mathrm{C} .
$$

This retraction is continuous at each vertex since $\|\rho(x)-s\| \leqslant\|x-s\|$ for all vertices $s$ and all $x \in \mathrm{U}_{1}$ such that $\|x-s\|<1$; it is continuous at each point that is not a vertex since in the neighbourhood of such a point it agrees with a constant map, the identity say. The map $(\theta, x) \mapsto(1-\theta) \rho(x)+\theta x$ from $[0,1] \times \mathrm{U}_{1}$ to $\mathrm{U}_{1}$ is a homotopy between $\rho$ and $1_{\mathrm{U}_{1}}$, and so $\mathrm{U}_{1}$ deformation retracts onto $|\mathrm{S}, \mathrm{B}|$. In particular, $\mathrm{U}_{1}$ is simply connected. So is $\mathrm{U}_{2}$.

The space $\mathrm{U}_{1} \cap \mathrm{U}_{2}$ consists of the connected subset $|\mathrm{S}, \mathrm{B}|$ and of the isolated points $\frac{x+x^{\prime}}{2}$ for $\left(x, x^{\prime}\right) \in \mathrm{C}$. By (4.7.5, Theorem) the group $\pi_{1}\left(|\mathrm{~S}, \mathrm{~A}|, s_{0}\right)$ has a basis indexed by C .

### 4.8.7

Proposition Every free group is isomorphic to a group of the form $\pi_{1}\left(|\mathrm{~S}, \mathrm{~A}|, s_{0}\right)$, where $(\mathrm{S}, \mathrm{A})$ is a connected graph and $s_{0} \in \mathrm{~S}$.

Proof Let $\Lambda$ be a set and L the free group on $\Lambda$ (4.7.4). Define a graph (S, A) by setting $S=\{0\} \cup\{1,2\} \times \Lambda$, and connecting 0 to all points and $(1, \lambda)$ to $(2, \lambda)$ for all $\lambda \in \Lambda$. Connecting only 0 to all points gives a tree (S, B). In the construction of 4.8.6, we may take $\mathrm{C}=\{((1, \lambda),(2, \lambda))\}_{\lambda \in \Lambda}$, which can be identified with $\Lambda$, and so $\pi_{1}(|\mathrm{~S}, \mathrm{~A}|, 0)$ is isomorphic to L .

### 4.8.8

Lemma Let $(\mathrm{S}, \mathrm{A})$ be a connected graph. Then the space $|\mathrm{S}, \mathrm{A}|$ has a simply connected covering.

Proof The space $|\mathrm{S}, \mathrm{A}|$ is locally contractible. Indeed, if $s$ is a vertex, then every ball centered at $s$ with radius $<2$ is contractible, and a point which is not a vertex has a fundamental system of neighbourhoods consisting of segments. Hence $|\mathrm{S}, \mathrm{A}|$ must necessarily be locally connected, and the lemma follows by (4.4.5).

### 4.8.9

Theorem All subgroups of a free group are free.

Proof Let L be a free group and H a subgroup of L. By (4.8.7), there is a connected graph ( $\mathrm{S}, \mathrm{A}$ ) and $s_{0} \in \mathrm{~S}$ such that $\pi_{1}\left(|\mathrm{~S}, \mathrm{~A}|, s_{0}\right)$ is isomorphic to L. Identify this fundamental group with L by an isomorphism. By (4.6.11, (a)), there is connected pointed covering ( $\mathrm{X}, x_{0}$ ) of $\left(|\mathrm{S}, \mathrm{A}|, s_{0}\right)$ such that $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ is isomorphic to H . By (4.8.3), the space X is homeomorphic to the geometric realization of a graph, and so $\mathrm{H}=\pi_{1}\left(\mathrm{X}, x_{0}\right)$ is free by (4.8.6).

Remark We have also proved (3.5.1) that a subgroup of a commutative group is free, i.e. of a free $\mathbb{Z}$-module is a free $\mathbb{Z}$-module. These two results only appear to be similar. For example, a submodule of a free $\mathbb{Z}$-module of rank $n$ is free of rank $\leqslant n$, whereas a free group on a set of two elements contains subgroups that are not finitely generated (4.8, Exercises 3 and 4).

## Exercises 4.8. (Graphs)

1.-Let (S, A) be a graph. Consider the topology $\mathscr{T}_{0}$ on $|\mathrm{S}, \mathrm{A}|$ defined as follows: a subset U in $|\mathrm{S}, \mathrm{A}|$ is open with respect to $\mathscr{T}_{0}$ if and only if $\mathrm{U} \cap[s, t]$ is open in $[s, t]$ for all $(s, t) \in \mathrm{A}$.
(a) Show that the topology $\mathscr{T}_{0}$ is finer than the topology $\mathscr{T}$ defined in (4.8.1), and that these topologies agree if and only if $(\mathrm{S}, \mathrm{A})$ is locally finite, i.e. if each vertex is only connected to finitely many vertices.
(b) Let $s_{0} \in \mathrm{~S}$. Show that the identity map

$$
i:(|\mathrm{S}, \mathrm{~A}|, \mathscr{T}) \rightarrow(|\mathrm{S}, \mathrm{~A}|, \mathscr{T})
$$

induces an isomorphism on the fundamental group with basepoint $s_{0}$. Is this map a homotopy equivalence?
2.-Let $G$ be a free group. Show that any two bases of $G$ are equipotent (consider the Abelianization of G , i.e. the quotient de G by the normal subgroup generated by the commutators).
3.-Consider the lemniscate $\mathrm{B}=\left\{z| | 1-z^{2} \mid=1\right\}$ in $\mathbb{C}$.
(a) Draw B. Show that it is homeomorphic to the geometric realization of a graph and give the latter's Cartesian diagram.
(b) Consider the covering X over B defined by the function $\log (1-z)$, i.e. $\mathrm{X}=$ $\left\{(z, x) \in \mathbb{C}^{2} \mid z \in \mathrm{~B}, e^{x}=1-z\right\}$ equipped with the projection $(z, x) \mapsto z$. Describe X.
(c) Describe the pointed covering Q of B whose fundamental group can be identified with the commutator group of $\pi_{1}(\mathrm{~B}, 0)$.
(d) Consider the polynomials $\mathrm{P}_{1}=\frac{3 \mathrm{X}-\mathrm{X}^{3}}{2}$ and $\mathrm{P}_{2}=1-4 \mathrm{X}^{2}+2 \mathrm{X}^{4}$. Show that $\mathrm{P}_{1}^{-1}(\mathrm{~B})$ and $\mathrm{P}_{2}^{-1}(\mathrm{~B})$ are coverings of B . Describe these coverings.
(e) Which of the above coverings are Galois?
4.-Let $\mathrm{L}(a, b)$ be a free group on a set $\{a, b\}$ with two elements.
(a) Let $f: \mathrm{L}(a, b) \rightarrow \mathbb{Z}$ be the homomorphism defined by $f(a)=1, f(b)=0$. Show that $\left(a^{n} \cdot b \cdot a^{-n}\right)_{n \in \mathbb{Z}}$ is a basis for the kernel of $f$.
(b) Let $g: \mathrm{L}(a, b) \rightarrow \mathbb{Z}^{2}$ be defined by $g(a)=(1,0), g(b)=(0,1)$. Show that the kernel of $g$ is the commutator group of $\mathrm{L}(a, b)$ with basis $\left(a^{p} \cdot b^{q} \cdot a \cdot b \cdot a^{-1}\right.$. $\left.b^{-1} \cdot b^{-q} \cdot a^{-p}\right)_{p, q \in \mathbb{Z}}$, or $\left(a^{p} \cdot b^{q} \cdot a \cdot b^{-q} \cdot a^{-p-1}\right)_{p, q \in \mathbb{Z}}$.
5.-Show that any countable graph is of homotopy equivalent to a graph all of whose vertices are connected to at most three others.
6.-Show that all trees are contractible.
7.-Let $U \subset \mathbb{R}^{2}$ be an open set $\mathbb{R}^{2}-K$ for some compact not necessarily locally connected set K , and let $x_{0} \in \mathrm{U}$. The aims is to show that $\pi_{1}\left(\mathrm{U}, x_{0}\right)$ is a free group with countable basis.
(a) Let $\mathrm{P}=[a, b] \times[c, d]$ be a rectangle in $\mathbb{R}^{2}$ and K a compact subset in the interior of P . Show that there is a sequence $\left(\mathrm{Q}_{n}\right)$ of rectangles such that
(i) $\mathrm{P}-\mathrm{K}=\bigcup \mathrm{Q}_{n}$;
(ii) $\stackrel{\circ}{\mathrm{Q}}_{i} \cap \stackrel{\circ}{\mathrm{Q}}_{j}=\varnothing$ for $j \neq i$;
(iii) each $a \in \mathrm{P}-\mathrm{K}$ has a neighbourhood meeting only finitely many $\mathrm{Q}_{i}$.

Show that, if $\mathrm{K} \neq \varnothing$, then the rectangles $\mathrm{Q}_{i}$ can be ordered so that

$$
(\forall n) \quad \partial \mathbf{Q}_{n} \not \subset \bigcup_{0}^{n-1} \mathrm{Q}_{i}
$$

where $\partial \mathbf{Q}_{n}$ denotes the boundary of $\mathbf{Q}_{n}$.
(b) Let $\left(\mathrm{Q}_{n}\right)$ be a sequence of rectangles satisfying the conditions of (a), and for all $n$, let $k(n)$ be the smallest $k$ such that $\mathrm{Q}_{k} \cap \mathrm{Q}_{n} \neq \varnothing$. Show that $k(n)$ tends to infinity as $n$ tends to infinity.
(c) Let $\mathrm{A} \subset \mathbb{R}^{2}$ be a finite union of rectangles $\mathrm{Q}_{i}$ satisfying condition (ii) of (a),
 where I is the set of indices $i$, and

$$
\mathrm{T} \subset\left\{(i, j) \mid \mathrm{Q}_{i} \cap \mathrm{Q}_{j} \neq \varnothing\right\}
$$

Assume that there is a homeomorphism embedding the geometric realization of $\Gamma$ in A under a such that the image of each vertex $i$ is contained in the interior of the rectangle $\mathrm{Q}_{i}$, and the image of each edge meets $\mathrm{Q}_{i}$ along a line segment. Assume that A deformation retracts onto $|\Gamma|$; let $\rho$ be a retraction and $h$ a homotopy from $1_{\mathrm{A}}$ to $\iota \circ \rho$. Show that there is a graph $\Gamma^{\prime}$ containing $\Gamma$ and such that the same situation holds for $\mathrm{A}^{\prime}$. Show that $\rho^{\prime}$ and $h^{\prime}$ can then be chosen in such a way that $h^{\prime}$ agrees with $h$ except on $\mathrm{Q}_{i}$ such that $\mathrm{Q}_{i} \cap \mathrm{Q} \neq \varnothing$.
(d) Keeping the notation of (a), show that $\mathrm{P}-\mathrm{K}$ deformation retracts onto the geometric realization of a countable graph. Deduce that the same holds for $\mathbb{R}^{2}-K$. Conclude.
(e) Show that, if U is connected and if K only has finitely many connected components, then $\pi_{1}\left(\mathrm{U}, x_{0}\right)$ has a basis with the same number of elements as the number of connected components of $K$. Does this still hold if $U$ is not connected? What if $K$ has infinitely many connected components? (Take $K$ to be the Cantor set in $\mathbb{R} \subset \mathbb{R}^{2}$.)
(f) If U is an arbitrary open subset of $\mathbb{R}^{2}$ and $x_{0} \in \mathrm{U}$, is the group $\pi_{1}\left(\mathrm{U}, x_{0}\right)$ free? 8.-Let $K$ be a compact subset of $\mathbb{R}^{2}$.
(a) Show that the category $\mathscr{C} \bullet v_{K}$ can be identified with $\lim _{\longrightarrow} \mathscr{C} \bullet v_{U}$, the inductive limit being that of neighbourhoods U of K .
(b) Show that K is simply connected if and only if its complement is connected.

### 4.9 Loops

### 4.9.1 Paths

Let B be a topological space, and $b, b^{\prime} \in \mathrm{B}$. A path from $b$ to $b^{\prime}$ in B is a continuous map $c$ from an interval $[0, a] \subset \mathbb{R}_{+}$to B such that $c(0)=b$ and $c(a)=b^{\prime}$. The length of the path is defined to be $a$. Let $c_{0}$ and $c_{1}$ be two paths from $b$ to $b^{\prime}$. A homotopy from $c_{0}$ to $c_{1}$ is a family $\left(h_{s}\right)_{s \in[0,1]}$ of paths from $b$ to $b^{\prime}$, such that $h_{0}=c_{0}$, $h_{1}=c_{1}$, such that the length $a_{s}$ of $h_{s}$ depends continuously on $s$ and that the map $(s, t) \mapsto h_{s}(t)$ from $\left\{(s, t) \mid s \in[0,1]\right.$ and $\left.t \in\left[0, a_{s}\right]\right\}$ to B is continuous. If there is a homotopy from $c_{0}$ to $c_{1}$, they are said to be homotopic and we write $c_{0} \simeq c_{1}$. The homotopy relation between paths from $b$ to $b^{\prime}$ is an equivalence relation. The set of homotopy classes of paths from $b$ to $b^{\prime}$ is written $\mathrm{S}\left(b, b^{\prime}\right)$.

### 4.9.2 Standard Paths

A path of length 1 is said to be standard.

If $c_{0}$ and $c_{1}$ are two homotopic standard paths, then there is a standard homotopy from $c_{0}$ to $c_{1}$, i.e. a homotopy $\left(h_{s}\right)$ such that $h_{s}$ is standard for all $s$, or equivalently a continuous map $h:[0,1] \times[0,1] \rightarrow \mathrm{B}$ such that $h(s, 0)=b, h(s, 1)=b^{\prime}$, $h(0, t)=c_{0}(t)$ and $h(1, t)=c_{1}(t)$ for $s, t \in[0,1]$. Indeed, if $\left(h_{s}\right)$ is a homotopy, $\bar{h}_{s}(t)=h_{s}\left(a_{s} t\right)$ defines a standard homotopy $\left(\bar{h}_{s}\right)$.

Moreover, every path is homotopic to a standard path. Indeed, if $c$ is a path of length $a$, define a homotopy $\left(h_{s}\right)$ from $c$ to a standard path by $h_{s}(t)=c\left(\frac{a t}{s+(1-s) a}\right)$; to check continuity, note that, for $a \neq 0$, the denominator does not vanish, and for $a=0$, the function $h$ is constant.

Hence, $\mathrm{S}\left(b, b^{\prime}\right)$ can be defined as the set of standard homotopy classes of standard paths from $b$ to $b^{\prime}$.

### 4.9.3 Juxtaposition

Let $b, b^{\prime}, b^{\prime \prime} \in \mathrm{B}, c$ a path from $b$ to $b^{\prime}$ of length $a$ and $c^{\prime}$ a path from $b^{\prime}$ to $b^{\prime \prime}$ of length $a^{\prime}$. Define a path $c^{\prime \prime}$ from $b$ to $b^{\prime \prime}$ of length $a+a^{\prime}$ by

$$
c^{\prime \prime}(t)=c(t) \text { for } t \in[0, a] \text { and } c^{\prime \prime}(t)=c^{\prime}(t-a) \text { for } t \in\left[a, a+a^{\prime}\right]
$$

We say that $c^{\prime \prime}$ is the juxtaposition of $c$ and $c^{\prime}$, and is written $c^{\prime} \cdot c$.
Juxtaposition is associative in the following sense: Let $b, b^{\prime}, b^{\prime \prime}, b^{\prime \prime \prime} \in \mathrm{B}, c$ a path from $b$ to $b^{\prime}, c^{\prime}$ from $b^{\prime}$ to $b^{\prime \prime}$ and $c^{\prime \prime}$ from $b^{\prime \prime}$ to $b^{\prime \prime \prime}$; then

$$
\left(c^{\prime \prime} \cdot c^{\prime}\right) \cdot c=c^{\prime \prime} \cdot\left(c^{\prime} \cdot c\right)
$$

For all $b \in \mathrm{~B}$, the path $e_{b}$ of zero length defined by $e_{b}(0)=b$ is the identity in the following sense: $c \cdot e_{b}=c$ for all paths $c$ with initial point $b$ and $e_{b} \cdot c=c$ for all paths $c$ with endpoint $b$.

Juxtaposition is compatible with the homotopy relation: if $c_{0}$ and $c_{1}$ are two homotopic paths from $b$ to $b^{\prime}, c_{0}^{\prime}$ and $c_{1}^{\prime}$ two homotopic paths from $b^{\prime}$ to $b^{\prime \prime}$, then the paths $c_{0}^{\prime} \cdot c_{0}$ and $c_{1}^{\prime} \cdot c_{1}$ are homotopic. Thus there is a well-defined map $\left(\gamma, \gamma^{\prime}\right) \mapsto$ $\gamma^{\prime} \cdot \gamma$, also called juxtaposition from $\mathbf{S}\left(b, b^{\prime}\right) \times \mathbf{S}\left(b^{\prime}, b^{\prime \prime}\right)$ to $\mathbf{S}\left(b, b^{\prime \prime}\right)$. This law is associative is the above sense, and for all $b$ the class $\varepsilon_{b}$ of $e_{b}$ is the identity. Moreover, for all $\gamma \in \mathbf{S}\left(b, b^{\prime}\right)$, there is an inverse $\gamma^{-1} \in \mathbf{S}\left(b^{\prime}, b\right)$ : let $c$ be a representative of $\gamma$ of length $a$, define $c^{*}$ by $c^{*}(t)=c(a-t), c^{*}$ is a path from $b^{\prime}$ to $b$; setting $h_{s}=$ $\left.\left(\left.c\right|_{[0, s a]}\right)^{*} \cdot c\right|_{[0, s a]}$, this gives a homotopy $\left(h_{s}\right)_{s \in[0,1]}$ from $e_{b}$ to $c^{*} \cdot c$ and a similar homotopy from $e_{b^{\prime}}$ to $c \cdot c^{*}$; so $\gamma^{-1} \cdot \gamma=\varepsilon_{b}$ and $\gamma \cdot \gamma^{-1}=\varepsilon_{b^{\prime}}$.

### 4.9.4 Lifting of Paths

Proposition Let X be a covering of $\mathrm{B}, b, b^{\prime} \in \mathrm{B}, x \in \mathrm{X}(b)$ and $c$ a path from $b$ to $b^{\prime}$. Then there is a unique path $\tilde{c}$ in X with initial point $x$ and such that $\pi \circ \tilde{c}=c$.

Proof Let $a$ be the length of $c$. The path $c$ is a map from $[0, a]$ to B and $c^{*} \mathrm{X}$ is a covering of $[0, a]$, which is trivial since $[0, a]$ is simply connected. Hence there is a unique section of $c^{*} \mathrm{X}$ passing through $(0, x)$, and a unique lifting $\tilde{c}$ corresponding to it with the desired property.

### 4.9.5

Proposition and Definition Let X be a covering of $\mathrm{B}, \gamma \in \mathrm{S}\left(b, b^{\prime}\right)$ and $x \in \mathrm{X}(b)$. Choose a representative $c$ of $\gamma$ and let $\tilde{c}$ be the unique lifting of $c$ with initial point $x$. Then the endpoint $x^{\prime}$ of $\tilde{c}$ in $\mathrm{X}\left(b^{\prime}\right)$ does not depend on the choice of $c$. It is denoted by $\gamma \cdot \mathrm{x} x$, or simply $\gamma \cdot x$.

Proof Let $c_{0}$ and $c_{1}$ be two homotopic paths from $b$ to $b^{\prime}$ with respective lengths $a_{0}$ and $a_{1}$, and let ( $h_{s}$ ) be a homotopy from $c_{0}$ to $c_{1}$. This homotopy defines a continuous map $h: \Delta \rightarrow \mathrm{B}$, where $\Delta=\left\{(s, t) \mid s \in[0,1]\right.$ and $\left.t \in\left[0, a_{s}\right]\right\}$. The space $\Delta$ deformation retracts onto $[0,1]$, so is simply connected, and $h^{*} \mathrm{X}$ is a trivial lifting of $\Delta$. Hence there is a continuous section of $h^{*} \mathrm{X}$ passing through $(0,0, x)$, and thus a continuous map

$$
\tilde{h}: \Delta \rightarrow \mathrm{X} \text { such that } \tilde{h}(0,0)=x \text { and } \pi \circ \tilde{h}=h .
$$

Since $t \mapsto \tilde{h}(0, t)$ is a lifting of $c_{0}$ with initial point $x, \tilde{h}(0, t)=\tilde{c}_{0}(t)$; and $\tilde{h}(s, 0)=$ $x$ for $s \in[0,1]$ since $s \mapsto \tilde{h}(s, 0)$ is a continuous map from the connected interval $[0,1]$ to the discrete set $\mathrm{X}(b)$. In particular $\tilde{h}(1,0)=x$. Since $t \mapsto \tilde{h}(1, t)$ is a lifting of $c_{1}$ with initial point $x, \tilde{h}(1, t)=\tilde{c}_{1}(t)$. The map $s \mapsto \tilde{h}\left(s, a_{s}\right)$ is continuous from the connected interval $[0,1]$ to the discrete set $\mathrm{X}\left(b^{\prime}\right)$, and so is constant and

$$
\tilde{c}_{0}\left(a_{0}\right)=\tilde{h}\left(0, a_{0}\right)=\tilde{h}\left(1, a_{1}\right)=\tilde{c}_{1}\left(a_{1}\right) .
$$

Remark Let $x \in \mathbf{X}(b), \gamma \in \mathbf{S}\left(b, b^{\prime}\right)$ and $\gamma^{\prime} \in \mathbf{S}\left(b^{\prime}, b^{\prime \prime}\right)$. Then

$$
\left(\gamma^{\prime} \cdot \gamma\right) \cdot x=\gamma^{\prime} \cdot(\gamma \cdot x)
$$

This follows readily from the definitions.

### 4.9.6 Poincaré Group

Let $\left(\mathrm{B}, b_{0}\right)$ be a pointed space. A loop in $\left(\mathrm{B}, b_{0}\right)$ is a path in B from $b_{0}$ to $b_{0}$. The set $\mathrm{S}\left(b_{0}, b_{0}\right)$ of homotopy classes of loops is a group with respect to the juxtaposition law called the singular fundamental group or the Poincaré group of ( $\mathrm{B}, b_{0}$ ). Il will be denoted it by $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$.

If $b_{0}, b_{1} \in \mathrm{~B}$ are in the same arc-connected component, then the groups $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ and $\widetilde{\pi}_{1}\left(\mathbf{B}, b_{1}\right)$ are isomorphic. Indeed, for all $\beta \in \mathbf{S}\left(b_{0}, b_{1}\right)$, the map $\gamma \mapsto \beta \cdot \gamma \cdot \beta^{-1}$ is an isomorphism from $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ onto $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{1}\right)$.

B is said to be simply arc-connected if it is arc-connected and there exists $b_{0} \in \mathrm{~B}$ such that $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)=\{e\}$; if this is the case then, for all $b \in \mathrm{~B}, \widetilde{\pi}_{1}(\mathrm{~B}, b)=\{e\} . \mathrm{B}$ is said to be locally simply arc-connected if all points of $B$ have a fundamental system of simply arc-connected neighbourhoods.

### 4.9.7 Functoriality

Let $f:\left(\mathrm{B}, b_{0}\right) \rightarrow\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ be a pointed space morphism. For all loops $c$ in $\left(\mathrm{B}, b_{0}\right)$, $f \circ c$ is a loop in ( $\mathrm{B}^{\prime}, b_{0}^{\prime}$ ) and the map $c \mapsto f \circ c$ is compatible with the homotopy relation and with juxtaposition. Hence $f_{*}: \widetilde{\pi_{1}}\left(\mathrm{~B}, b_{0}\right) \rightarrow \widetilde{\pi}_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)$ is a group homomorphism.

Together with $f \mapsto f_{*}$, the map $\widetilde{\pi}_{1}$ is a functor from the category of pointed est spaces to the category of groups. This functor commutes with products.

### 4.9.8 Poincaré Group of a Product

Proposition Let $\left(\mathrm{B}, b_{0}\right)$ and $\left(\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ be pointed topological spaces. The homomorphism

$$
\tilde{\pi}_{1}\left(\mathrm{~B} \times \mathrm{B}^{\prime},\left(b_{0}, b_{0}^{\prime}\right)\right) \rightarrow \widetilde{\pi_{1}}\left(\mathrm{~B}, b_{0}\right) \times \widetilde{\pi}_{1}\left(\mathrm{~B}^{\prime}, b_{0}^{\prime}\right)
$$

defined by the projections is an isomorphism.
Proof A standard loop in $\left(\mathbf{B} \times \mathbf{B}^{\prime},\left(b_{0}, b_{0}^{\prime}\right)\right)$ is given by a pair $\left(\gamma, \gamma^{\prime}\right)$, where $\gamma$ (resp. $\gamma^{\prime}$ ) is a standard loop in ( $\mathrm{B}, b_{0}$ ) (resp. ( $\left.\mathrm{B}^{\prime}, b_{0}^{\prime}\right)$ ); similarly a standard homotopy in $\left(\mathrm{B} \times \mathrm{B}^{\prime},\left(b_{0}, b_{0}^{\prime}\right)\right)$ is given by a pair of standard homotopies. The proposition follows.

### 4.9.9 The Poincaré Group and the Fundamental Group

For any covering $\mathbf{X}$ of $\mathbf{B}$, the map $(\gamma, x) \mapsto \gamma \cdot x$ is an action of the group $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ on $\mathrm{X}\left(b_{0}\right)$. In particular, for all $\gamma \in \widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$, the map $\gamma_{\mathrm{X}}: x \mapsto \gamma \cdot x$ is a permutation of $\mathrm{X}\left(b_{0}\right)$. If $f: \mathrm{X} \rightarrow \mathrm{X}^{\prime}$ is a covering morphism of B , then the map $f_{b_{0}}: \mathrm{X}\left(b_{0}\right) \rightarrow \mathrm{X}^{\prime}\left(b_{0}\right)$ is compatible with the actions of $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$, i.e. $f(\gamma \cdot x)=\gamma \cdot f(x)$ for $x \in \mathrm{X}\left(b_{0}\right)$ and $\gamma \in \widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$. Hence the maps $\gamma_{\mathrm{X}}$ form an automorphism of the functor $\mathrm{X} \mapsto \mathrm{X}\left(b_{0}\right)$, and so is an element $\gamma$. $\in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$. The $\operatorname{map} \varepsilon: \gamma \mapsto \gamma$. from $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ to $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is a group homomorphism.

We aim to prove the following result:
Theorem Let B be a locally simply arc-connected space and $b_{0} \in \mathrm{~B}$. Then $\varepsilon$ : $\tilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right) \rightarrow \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ is an isomorphism.

This theorem will be proved in 4.9.14. Surjectivity is easy (4.9.10). To prove injectivity, the existence of a covering X of $\left(\mathrm{B}, b_{0}\right)$ such that $\gamma$ acts non trivially on $\mathrm{X}\left(b_{0}\right)$, for nonzero elements $\gamma \in \widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ needs to be shown. The universal covering satisfies this property, but this will be more clearly seen by considering a different construction of the universal covering from that given in section 4.1.

### 4.9.10 Surjectivity

Proposition Let $\left(\mathrm{B}, b_{0}\right)$ be a locally simply arc-connected pointed space having a universal covering. Then $\varepsilon$ is surjective.

Proof Let $\gamma \in \pi_{1}\left(\mathrm{~B}, b_{0}\right)$ and (E, $t$ ) be a universal pointed covering of ( $\mathrm{B}, b_{0}$ ); set $u=\gamma_{\mathrm{E}}(t)$. Since B is locally simply arc-connected, so is E , and as E is connected, it is arc-connected. Hence there is a path $c$ from $t$ to $u$ in E , and $p \circ c$ is a loop in ( $\mathrm{B}, b_{0}$ ), where $p$ is the projection $\mathrm{E} \rightarrow \mathrm{B}$. Writing $\alpha$ for the class of $p \circ c$ in $\tilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right), \alpha \cdot t=u$, and so $\varepsilon(\alpha)(t)=u=\gamma(t)$; thus $\varepsilon(\alpha)=\gamma$ since $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ acts simply transitively on $\mathrm{E}\left(b_{0}\right)$ by (4.6.2).

### 4.9.11

Proposition All locally arc-connected, simply arc-connected spaces are simply connected.

Proof Since all coverings of B are sums of connected coverings, it suffices to shows that every connected covering is trivial. Let X be a connected covering of $\mathrm{B}, b \in \mathrm{~B}$, $x, x^{\prime} \in \mathrm{X}(b)$. Since X is locally arc-connected and connected, it is arc-connected, and so there is a path $c$ from $x$ to $x^{\prime}$, and $p \circ c$ is a loop in $(\mathrm{B}, b)$. Writing $\alpha$ for the
class of $p \circ c$ in $\widetilde{\pi}_{1}(\mathrm{~B}, b), \alpha \cdot x=x^{\prime}$; however $\alpha=e$, so $x=x^{\prime}$. Hence X has degree 1 , and so is trivial.

Corollary All connected, locally simply arc-connected spaces have a simply connected covering.

Indeed by the previous proposition, such a space is locally simply connected, and so the result follows from Theorem 4.4.5.

### 4.9.12 Construction of Universal Coverings as Spaces of Paths

Let ( $\mathrm{B}, b_{0}$ ) be a locally simply arc-connected pointed space. By the previous corollary it has a universal covering. To prove injectivity in Theorem 4.9.9, we give a new construction of this universal covering.

For all $b \in \mathrm{~B}$, let $\mathrm{S}(b)=\mathrm{S}\left(b_{0}, b\right)$ be the set of homotopy classes of paths from $b_{0}$ to $b$. Set $\mathrm{S}=\bigcup_{b \in \mathrm{~B}} \mathrm{~S}(b)$, and let $p$ be the map from S to B sending $\mathrm{S}(b)$ to $b$.

For all simply arc-connected open subset U of B , define an equivalence relation $\sim_{\mathrm{U}}$ on $\left.\mathrm{S}\right|_{\mathrm{U}}$ as follows: let $\gamma,\left.\gamma^{\prime} \in \mathrm{S}\right|_{\mathrm{U}}, b=p(\gamma)$ and $b^{\prime}=p\left(\gamma^{\prime}\right)$. All paths in U from $b$ to $b^{\prime}$ are homotopic since U is simply arc-connected; let $\beta \in \mathrm{S}\left(b, b^{\prime}\right)$ be the unique class containing the paths from $b$ to $b^{\prime}$ in U . We write $\gamma \sim_{\mathrm{U}} \gamma^{\prime}$ if and only if $\beta \cdot \gamma=\gamma^{\prime}$. Set $\mathrm{F}_{\mathrm{U}}=\left.\mathrm{S}\right|_{\mathrm{U}} / \sim_{\mathrm{U}}$ and let $\chi$ denote the canonical map $\left.\mathrm{S}\right|_{\mathrm{U}} \rightarrow \mathrm{F}_{\mathrm{U}}$.

Lemma (a) The map $\tau_{\mathrm{U}}:\left.\mathrm{S}\right|_{\mathrm{U}} \rightarrow \mathrm{U} \times \mathrm{F}_{\mathrm{U}}$ defined by $\tau_{\mathrm{U}}(\gamma)=(p(\gamma), \chi(\gamma))$ is bijective.
(b) If $\mathrm{V} \subset \mathrm{U}$ is simply arc-connected and open, then the map $\iota: \mathrm{F}_{\mathrm{V}} \rightarrow \mathrm{F}_{\mathrm{U}}$ induced by the canonical injection $\mathrm{S}_{\mathrm{V}} \rightarrow \mathrm{S}_{\mathrm{U}}$ is bijective.

## Proof

(a) Injectivity: let $\gamma,\left.\gamma^{\prime} \in \mathrm{S}\right|_{\mathrm{U}}$ such that $\tau_{\mathrm{U}}(\gamma)=\tau_{\mathrm{U}}\left(\gamma^{\prime}\right)$. Then $p(\gamma)=p\left(\gamma^{\prime}\right)$ and $\gamma^{\prime}=\beta \cdot \gamma$ with $\beta \in \widetilde{\pi}_{1}(\mathrm{U}, p(\gamma))$. Since U is simply connected, $\beta=e$ and $\gamma=\gamma^{\prime}$.

Surjectivity: let $(b, f) \in \mathrm{U} \times \mathrm{F}_{\mathrm{U}}$, and $\gamma$ be a representative of $f$. Set $b^{\prime}=p(\gamma)$. Let $\beta^{\prime} \in \mathbf{S}\left(b, b^{\prime}\right)$ be the class containing the paths from $b$ to $b^{\prime}$ in U . Then $\tau_{\mathrm{U}}\left(\beta^{\prime} \cdot \gamma\right)=$ $(b, f)$.
(b) Let $b \in \mathrm{~V}$. There is a commutative diagram


By (a), $\tau_{\mathrm{U}}$ and $\tau_{\mathrm{V}}$ are bijective, and hence so is $\iota$.

Topology on S For every simply arc-connected U, equip $F_{U}$ with the discrete topology and $\mathrm{U} \times \mathrm{F}_{\mathrm{U}}$ with the product topology. Let $\mathscr{T}_{\mathrm{U}}$ be the topology on $\left.\mathrm{S}\right|_{\mathrm{U}}$ obtained by transferring that of $\mathrm{U} \times \mathrm{F}_{\mathrm{U}}$ by $\tau_{\mathrm{U}}$. If V is an open simply arc-connected subset of U , the lemma implies that $\mathscr{T}_{\mathrm{V}}$ is induced by $\mathscr{T}_{\mathrm{U}}$. If U and $\mathrm{U}^{\prime}$ are two simply arcconnected sets, then $\mathrm{U} \cap \mathrm{U}^{\prime}$ can be covered by simply arc-connected $\mathrm{V}_{\lambda}$, and for all $\lambda, \mathscr{T}_{\mathrm{U}}$ and $\mathscr{T}_{\mathrm{U}^{\prime}}$ induce the same topology $\mathscr{T}_{\mathrm{V}_{\lambda}}$ on $\left.\mathrm{S}\right|_{\mathrm{V}_{\lambda}}$. Hence, by the topology gluing lemma (4.4.6, Gluing lemma) $\mathscr{T}_{\mathrm{U}}$ and $\mathscr{T}_{\mathrm{U}^{\prime}}$ induce the same topology on $\left.\mathrm{S}\right|_{\mathrm{UnU}^{\prime}}$. This lemma also tells us that there is a unique topology $\mathscr{T}$ on S inducing $\mathscr{T}_{\mathrm{U}}$ on $\left.\mathrm{S}\right|_{\mathrm{U}}$ for any simply arc-connected cover U .

Equipped with the topology $\mathscr{T}$ and the projection $p, \mathrm{~S}$ is a covering of B . We will see in subsection 4.9.15 that this covering is universal.

### 4.9.13

Proposition With the notation of 4.9.12, the action of $\tilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ on $\mathrm{S}\left(b_{0}\right)$ is simply transitive.

Proof As $\mathrm{S}\left(b_{0}\right)=\tilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$, we show that $\tilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$ acts on itself by left translations. More precisely, for $\gamma \in \widetilde{\pi_{1}}\left(\mathrm{~B}, b_{0}\right)$ and $\alpha \in \mathbf{S}\left(b_{0}\right)$, we show that $\gamma \cdot \mathrm{s} \alpha$ defined in (4.9.5) is equal to $\gamma \cdot \alpha$, defined by juxtaposition. More generally, we show that for $b \in \mathbf{B}, b^{\prime} \in \mathbf{B}, \alpha \in \mathbf{S}(b)$ and $\gamma \in \mathbf{S}\left(b, b^{\prime}\right), \gamma \cdot \mathrm{s} \alpha=\gamma \cdot \alpha$.

Suppose first that $\gamma$ is represented by a path $c$ of length $a$ contained in a simply arc-connected open subset $\mathbf{U}$ of $\mathbf{B}$. Define $\tilde{c}:[0, a] \rightarrow \mathbf{S}$ by $\tilde{c}(t)=\gamma_{t} \cdot \alpha$, where $\gamma_{t}$ is the homotopy class of $\left.c\right|_{[0, t]}$. Then $p \circ \tilde{c}=c$. We check that $\tilde{c}$ is continuous. For this, it suffices to check that $\tau_{\mathrm{U}} \circ \tilde{c}$ is continuous. Now, $\tau_{\mathrm{U}}=(p, \chi), p \circ \tilde{c}=c$ is continuous, and $\chi \circ \tilde{c}$ is constant. Indeed $\tilde{c}(t) \underset{\mathrm{U}}{\sim} \alpha$ since $\gamma_{t}$ is the class in $\mathrm{S}(b, c(t))$ of paths contained in U . Hence, $\tilde{c}$ is the lifting of $c$ in S with initial point $\alpha$, and $\gamma \cdot \mathrm{s} \alpha=c(\alpha)=\gamma \cdot \alpha$.

By the subdivision lemma (Proposition 4.1.10), all paths can be obtained by juxtaposition of paths contained in a simply arc-connected open set. So the result holds for all $\gamma$.

### 4.9.14 Proof of Theorem 4.9.9

Let ( $\mathrm{B}, b_{0}$ ) be a locally simply arc-connected pointed space. By (4.9.11, Corollary), the connected component of $b_{0}$ has a universal covering; so ( $\mathrm{B}, b_{0}$ ) has a universal pointed cover. Moreover B is locally arc-connected, and so $\varepsilon$ is surjective by (4.9.10, Proposition).

Let $\gamma$ be an identity element in $\widetilde{\pi}_{1}\left(\mathrm{~B}, b_{0}\right)$. By (4.9.13, Proposition), $\gamma$ acts non trivially on $\mathrm{S}\left(b_{0}\right)$; so $\varepsilon(\gamma) \neq e$, and $\varepsilon$ is injective.

### 4.9.15

Corollary With the notation of (4.9.12), S is a universal covering of B .
Proof If necessary replacing B by the connected component of $b_{0}$, B can be assumed to be connected. By (4.9.11, Corollary), B has a universal covering E. The fundamental group $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ acts simply transitively on $\mathrm{E}\left(b_{0}\right)$. By Theorem 4.9.9 and Proposition 4.9.13, $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ also acts simply transitively on $\mathrm{S}\left(b_{0}\right)$. The $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ sets $\mathrm{S}\left(b_{0}\right)$ and $\mathrm{E}\left(b_{0}\right)$ are therefore isomorphic. By Theorem 4.6.8, it follows that the coverings $S$ are $E$ isomorphic.

## Exercises 4.9. (Loops)

1.-Let $\mathrm{C}_{n}$ be the circle in $\mathbb{R}^{2}$ centered at $\left(\frac{1}{2^{n}}, 0\right)$ and passing through O . Set

$$
\mathrm{X}=\bigcup \mathrm{C}_{n}, \quad \mathrm{X}_{n}=\mathrm{C}_{0} \cup \cdots \cup \mathrm{C}_{n} .
$$

(a) For all $n$, let $f_{n}: \mathrm{X}_{n+1} \rightarrow \mathrm{X}_{n}$ be the map inducing the identity on $\mathrm{X}_{n}$ and sending $\mathrm{C}_{n+1}$ onto O . Show that X can be identified with the projective limit of the system defined by all $\mathrm{X}_{n}$ and $f_{n}$. Show that the category $\mathscr{C} \boldsymbol{\bullet} \boldsymbol{v}_{\mathrm{X}}$ can be identified with $\underset{\longrightarrow}{\lim } \mathscr{C} \boldsymbol{v}_{\mathrm{X}_{n}}$. Deduce that $\pi_{1}(\mathrm{X}, 0)=\lim \pi_{1}\left(\mathrm{X}_{n}, 0\right)$.
(b) Show that the map $\varepsilon: \widetilde{\pi}_{1}(X, 0) \rightarrow \pi_{1}(X, 0)$ is injective, but not surjective.
(c) Let $Y$ be the cone in $\mathbb{R}^{3}$, with basis $X$ and vertex $(0,0,1)$. Set $Z$ to be the union of Y and of its symmetric $\mathrm{Y}^{\prime}$ with respect to 0 (cf. Zisman [1], 3.3.5.3, p. 129). Show that Z is simply connected but not simply connected by loops, i.e. $\widetilde{\pi}_{1}(\mathrm{Z}, 0) \neq 0$. In particular $\varepsilon_{\mathrm{Z}}$ is surjective, but not injective.
(d) In a topological space, two compact subsets A and B whose intersection is arc-connected form a Van Kampen pair for $\pi_{1}$ (resp. for $\left.\widetilde{\pi}_{1}\right)$ if $\pi_{1}\left(\mathrm{~A} \cup \mathrm{~B}, x_{0}\right)$ is the sum of $\pi_{1}\left(\mathrm{~A}, x_{0}\right)$ and $\pi_{1}\left(\mathrm{~B}, x_{0}\right)$ amalgamated by $\pi_{1}\left(\mathrm{~A} \cap \mathrm{~B}, x_{0}\right)$ for $x_{0} \in \mathrm{~A} \cap \mathrm{~B}$ (resp. same for $\widetilde{\pi}_{1}$ ). Among the following pairs, which are Van Kampen for $\pi_{1}$ ? for $\widetilde{\pi}_{1}$ ?

- $\left(X, X^{\prime}\right)$, where $\mathrm{X}^{\prime}$ is the symmetric of X with respect to 0 ;
- ( $\mathrm{Y}, \mathrm{Y}^{\prime}$ );
- $\left(\mathrm{X}_{1}, \mathrm{X}_{1}^{\prime}\right)$, where $\mathrm{X}_{1}$ is the union of the segment $[0,(1,0)]$ and of the translation of X by the vector $(1,0)$, and where $\mathrm{X}_{1}^{\prime}$ is the symmetric of $\mathrm{X}_{1}$ with respect to 0 .
2.-Let $X$ be a topological space, $\mathscr{U}=\left(U_{i}\right)_{i \in \mathrm{I}}$ a simply arc-connected open cover of X , and $\mathscr{U}^{\prime}=\left(\mathrm{U}_{i}^{\prime}\right)_{i \in \mathrm{I}}$ a shrinking of $\mathscr{U}$ (i.e. an open cover such that $\overline{\mathrm{U}}_{i}^{\prime} \subset \mathrm{U}_{i}$ for all $i)$. For every 2-element subset $\{i, j\}$ of I , let $\mathrm{C}_{i j}$ be the set of connected components of $\mathrm{U}_{i} \cap \mathrm{U}_{j}$ containing at least one point of $\mathrm{U}_{i}^{\prime} \cap \mathrm{U}_{j}^{\prime}$.

Let $\Gamma$ be a graph with $\mathrm{I} \cup \bigcup \mathrm{C}_{i j}$ as its set of vertices in which a vertex $c \in \mathrm{C}_{i j}$ is connected to the vertices $i$ and $j$.
(a) Show that $\Gamma$ is connected if and only if X is connected.
(b) For all $i \in \mathrm{I}$ let $x_{i} \in \mathrm{U}_{i}$ for all $c \in \mathrm{C}_{i j}$; let $x_{c} \in c \cap \mathrm{U}_{i}^{\prime} \cap \mathrm{U}_{j}, \gamma_{c, i}$ and $\gamma_{c, j}$ paths respectively from $x_{i}$ and $x_{j}$ to $x_{c}$. Define a continuous map $\varphi:|\Gamma| \rightarrow \mathrm{X}$. Show that

$$
\varphi_{*}: \pi_{1}(\Gamma, e) \rightarrow \pi_{1}(\mathrm{X}, \varphi(e))
$$

is surjective for $e \in|\Gamma|$.
(c) Show that for all locally simply arc-connected compact spaces X and all $x \in \mathrm{X}$, the group $\pi_{1}(\mathrm{X}, x)$ is finitely generated.
(d) Show that the result continues to hold when X is compact, locally simply connected (not arc-connected) [there is no map $\varphi$, but $\pi_{1}(\mathrm{X}, x)$ is isomorphic to a quotient of $\left.\pi_{1}(\Gamma, e)\right]$.
3. (Partial blowup of $\mathbb{R}^{2}$ at all points of the axis $\mathrm{O} y$.) $-\operatorname{Set} \mathrm{U}=\mathbb{R}^{*} \times \mathbb{R}=$ $\mathbb{R}^{2}-(\{0\} \times \mathbb{R})$ and $\left.\mathrm{E}=\mathbb{R} \times\right]-1,1[$. Denote by X the disjoint union $\mathrm{U} \sqcup \mathrm{E}$.

For $\left(y_{0}, p\right) \in \mathrm{E}$ and $\varepsilon>0$, set
$\mathrm{V}_{\varepsilon}\left(y_{0}, p\right)=\left\{(x, y)\left|0<|x|<\varepsilon,\left|\frac{y-y_{0}}{x}-p\right|<\varepsilon\right\} \sqcup\left(\left\{y_{0}\right\} \times\right] p-\varepsilon, p+\varepsilon[)\right.$.
(a) Show that there is a unique topology on $X$ inducing the topology of $\mathbb{R}^{2}$ on $U$, that of $\left.\mathbb{R}_{\text {discr }} \times\right]-1,1\left[\right.$ on $E$, where $\mathbb{R}_{\text {discr }}$ is the space obtained by equipping $\mathbb{R}$ with the discrete topology, such that U is open and $\mathrm{V}_{\varepsilon}\left(y_{0}, p\right)$ form a fundamental system of neigbourhoods of $\left(y_{0}, p\right) \in \mathrm{E}$. Equip X with this topology.
(b) Show that X is a topological surface. Define an $\mathbb{R}$-analytic manifold structure on X.
(c) Show that the space X is Hausdorff, arc-connected, locally compact, admitting a countable dense set but that it is not the countable union of compact sets. Show that X is neither metrizable, nor normal.
(d) Show that for $b_{0} \in \mathrm{X}$, the Poincaré group $\widetilde{\pi}_{1}\left(\mathrm{X}, b_{0}\right)$ can be identified with the fundamental group $\pi_{1}\left(\mathrm{X}, b_{0}\right)$ and is not countable.
(e) Can X be embedded in the transfinite line (chap. 1, §6, Exercise 1)?

## Reference

1. M. Zisman, Topologie algébrique éléémentaire (Collection U, Armand Colin, 1972)

## Chapter 5 <br> Galois Theory

Unless stated otherwise, all algebras considered in this chapter are assumed to be unital associative commutative, and the homomorphisms to be unital.

The letter K will denote a field.

### 5.1 Extensions

### 5.1.1 Finite Algebras

Let A be a K -algebra. Denote by $\iota_{\mathrm{A}}$ the unique homomorphism from K to A , defined by $\iota(x)=x \cdot 1$. If $\mathrm{A} \neq 0$, then $\iota_{\mathrm{A}}$ is an injection, called canonical.

A is said to be a finite algebra on K if A is a finite dimensional K -vector space. The degree of $A$, written $(A: K)$ or $\operatorname{deg}_{K} A$, is the dimension of $A$ as a $K$-vector space.

### 5.1.2

Let A be a K-algebra and $x \in \mathrm{~A}$.
Proposition and Definition The following conditions are equivalent:
(i) the subalgebra $\mathrm{K}[x]_{\mathrm{A}}$ of A generated by $x$ is finite over K ;
(ii) there is a nonzero polynomial $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ such that $\mathrm{P}(x)=0$;
(iii) the powers of $x$ are linearly dependent.

If these conditions hold, then $x$ is said to be algebraic over K.
Proof (ii) is clearly equivalent to (iii). If the powers of $x$ are linearly independent, then the subalgebra $\mathrm{K}[x]_{\mathrm{A}}$ of A is not finite, and so (i) implies (iii). (ii) $\Rightarrow$ (i) since
the algebra homomorphism $\chi: \mathrm{K}[\mathrm{X}] \rightarrow \mathrm{A}$ defined by $\chi(\mathrm{X})=x$ vanishes at P and so factorizes through the finite dimensional K -vector space $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$.

The polynomials $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ such that $\mathrm{P}(x)=0$ are sometimes called the equations of $x$. They form an ideal with a single monic generator (i.e. whose leading coefficient is 1 ), called the minimal polynomial of $x$ in A over K. Its degree is the degree of $x$ over K , as well as that of the subalgebra of A generated by $x$.

### 5.1.3

Proposition Let A be a K-algebra. The algebraic elements of A over K form a subalgebra of A .

Proof Let $x, y \in \mathrm{~A}$ be algebraic over K. The subalgebra generated by $(x+y)$ (resp. $x y)$ is contained in the algebra $\mathrm{K}[x, y]_{\mathrm{A}}$ isomorphic to a quotient of $\mathrm{K}[x]_{\mathrm{A}} \otimes \mathrm{K}[y]_{\mathrm{A}}$ (3.8.16). Thus $x+y$ and $x y$ are algebraic over K.

Remark This proposition can be proved directly by taking condition (ii) de 5.1.2. as the definition of algebraic elements, but this is harder. Note that the above proof does enable us to find an equation of $x+y$ (resp. $x y$ ) from those of $x$ and $y$.

### 5.1.4

Let A be a K-algebra.
Definition We say that A is an extension of K if A is a field with respect to its ring structure.

A is said to be a trivial extension of K if $\iota_{\mathrm{A}}$ is an isomorphism. All K -algebras of degree 1 are trivial extensions of $K$.

If A is a extension field of K , i.e. a field containing K as a subfield, then A equipped with a natural structure of K -algebra is an extension of K ; all extensions are isomorphic to an extension of this type.

If A is of the form $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$ then A is a finite extension of K if and only if P is irreducible (3.2.10, Lemma).

Proposition If A is finite and integral, then A is an extension of K .
Proof Let $a(\neq 0) \in \mathrm{A}$. The map $x \mapsto a x$ from A to A is injective and hence surjective since A is a finite dimensional vector space. Hence there exists $x \in \mathrm{~A}$ such that $a x=1$, and so A is a field.

### 5.1.5

Chinese Remainder Theorem Let A be a ring, $\left(\mathrm{B}_{i}\right)_{i \in \mathrm{I}}$ a finite family of rings and for all $i \in \mathrm{I}, \varphi_{i}: \mathrm{A} \rightarrow \mathrm{B}_{i}$ a ring homomorphism. Assume that, for each $i \in \mathrm{I}, \varphi_{i}$ is surjective, and that for all $(i, j) \in \mathrm{I} \times \mathrm{I}$ with $i \neq j, \operatorname{Ker} \varphi_{i}+\operatorname{Ker} \varphi_{j}=\mathrm{A}$. Then the map $\Phi: \mathrm{A} \rightarrow \prod_{i} \mathrm{~B}_{i}$ defined by $\Phi(x)=\left(\varphi_{i}(x)\right)_{i \in \mathrm{I}}$ is surjective.

Proof For all $(i, j) \in \mathrm{I} \times \mathrm{I}$ with $i \neq j$, there exists $u_{i j} \in \mathrm{~A}$ such that $\varphi_{i}\left(u_{i j}\right)=1$ and $\varphi_{j}\left(u_{i j}\right)=0$. Indeed, for $u \in \operatorname{Ker} \varphi_{j}$ and $v \in \operatorname{Ker} \varphi_{i}$ such that $u+v=1, \varphi_{j}(u)=0$ and $\varphi_{i}(u)=\varphi_{i}(1-v)=1$. Take $u_{i j}=u$.

For all $i \in \mathrm{I}$, set $u_{i}=\prod_{j \neq i} u_{i j}$; then $\varphi_{i}\left(u_{i}\right)=1$ and $\varphi_{j}\left(u_{i}\right)=0$ for all $j \neq i$.
Let $b=\left(b_{i}\right)_{i \in \mathrm{I}} \in \prod_{i} \mathrm{~B}_{i}$ and for all $i \in \mathrm{I}$, let $x_{i} \in \mathrm{~A}$ be such that $\varphi_{i}\left(x_{i}\right)=b_{i}$. Setting $x=\sum_{i \in \mathrm{I}} u_{i} x_{i}, \Phi(x)=b$,

Remark The condition $\operatorname{Ker} \varphi_{i}+\operatorname{Ker} \varphi_{j}=\mathrm{A}$ holds in particular if $\operatorname{Ker} \varphi_{i}$ and $\operatorname{Ker} \varphi_{j}$ are distinct maximal ideals. See also (3.5.9).

### 5.1.6

Corollary 5.1 Let A be a finite K-algebra.
(1) The set X of maximal ideals of A is finite.
(2) The map $\Phi: \mathrm{A} \rightarrow \prod_{\mathfrak{m} \in \mathrm{X}} \mathrm{A} / \mathfrak{m}$ is surjective and its kernel consists of the nilpotent ideals of A .

Proof (1) Let Y be a finite subset of X . By the Chinese remainder theorem, $\mathrm{A} \rightarrow \prod_{\mathfrak{m} \in \mathrm{Y}} \mathrm{A} / \mathfrak{m}$ is surjective, and so card $\mathrm{Y} \leqslant(\mathrm{A}: K)$. Hence card $\mathrm{X} \leqslant(\mathrm{A}: \mathrm{K})$.
(2) By the Chinese remainder theorem and by (1), $\Phi$ is surjective; by 5.1.4, all prime ideals are maximal and by 3.1.6, the intersection of all prime ideals of A is the set of nilpotent elements of A.

An algebra with no nonzero nilpotent elements is said to be reduced (see 3.1.7).
Corollary 5.2 Every reduced finite K-algebra is isomorphic to a finite product of finite extensions of K .

## 5.1 .7

Dedekind Theorem Let A be a K-algebra and L an extension of K . Consider the L -vector space $\mathscr{L}_{\mathrm{K}}(\mathrm{A}, \mathrm{L})$ of K -linear maps from A to L . In this vector space, the K -algebra homomorphisms from A to L are linearly independent.

Proof By (3.8.17), $\operatorname{Hom}_{\mathrm{K}-\mathscr{A l g}_{g}(\mathrm{~A}, \mathrm{~L}) \text { can be identified with the set }}$ $\operatorname{Hom}_{\mathrm{L}-\mathscr{A l g}}\left(\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}, \mathrm{L}\right)$. Let $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ be a finite family of elements of $\operatorname{Hom}_{\mathrm{K}-\mathscr{A} \ell_{g}}(\mathrm{~A}, \mathrm{~L})$ and $\left(\tilde{\varphi}_{i}\right)_{i \in \mathrm{I}}$ the family of corresponding elements in $\operatorname{Hom}_{\mathrm{L}-\mathscr{A l g}}\left(\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}, \mathrm{L}\right)$. For each $i \in \mathrm{I}, \tilde{\varphi}_{i}$ is surjective; as L is a field, the kernel $\mathfrak{m}_{i}$ of $\tilde{\varphi}_{i}$ is a maximal ideal and $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}=\mathfrak{m}_{i} \oplus \mathrm{~L} \cdot 1$. If $\varphi_{i} \neq \varphi_{j}$ for $i \neq j$, then moreover $\mathfrak{m}_{i} \neq \mathfrak{m}_{j}$ since $\tilde{\varphi}_{i}(1)=\tilde{\varphi}_{j}(1)=1$ 。

By the Chinese remainder theorem, the map $\Phi: x \mapsto\left(\tilde{\varphi}_{i}(x)\right)_{i \in \mathrm{I}}$ from $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ to $L^{\mathrm{I}}$ is surjective.

Assuming there is a relation $\sum a_{i} \varphi_{i}=0, \sum a_{i} \tilde{\varphi}_{i}=0$ and the image of $\Phi$ is contained in the hyperplane $\left\{\left(t_{i}\right) \in \mathrm{L}^{\mathrm{I}} \mid \sum a_{i} t_{i}=0\right\}$, which contradicts the surjectivity of $\Phi$.

### 5.1.8 Transcendental Elements

Let A be an extension of K . An element $x \in \mathrm{~A}$ is said to be transcendental over K if it is not algebraic, in other words if the homomorphism $\varphi: \mathrm{K}[\mathrm{X}] \rightarrow \mathrm{A}$ defined by $\varphi(\mathrm{X})=x$ is injective.

A family $\left(x_{i}\right)_{i \in \mathrm{I}}$ of elements of A is said to be algebraically free if the homomorphism $\phi: \mathrm{K}\left[\mathrm{X}_{i}\right]_{i \in \mathrm{I}} \rightarrow$ A defined by $\phi\left(\mathrm{X}_{i}\right)=x_{i}$ is injective, in other words if there is a homomorphism $\Phi: K\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}} \rightarrow \mathrm{A}$ such that $\Phi\left(\mathrm{X}_{i}\right)=x_{i}\left(\right.$ where $\mathrm{K}\left(\mathrm{X}_{i}\right)$ denotes the field of fractions of $\mathrm{K}\left[\mathrm{X}_{i}\right]$ ). A maximal algebraically free family is a transcendental basis of A over K. If $\left(x_{i}\right)_{i=1, \ldots, k}$ is a transcendental basis of A, then the field A equipped with $\Phi$ is an algebraic extension of $K\left(X_{1}, \ldots, X_{k}\right)$.

Proposition and Definition All transcendental bases of A over K have the same cardinality. This cardinal is called the transcendental degree of A over K .

Lemma (Exchange Lemma) Let $\left(x_{1}, \ldots, x_{k}\right)$ be a transcendental basis of A over K and $y \in \mathrm{~A}$ a transcendental element over K . Then there exists $i_{0} \in\{1, \ldots, k\}$ such that the family $\left(x_{i}^{\prime}\right)$ defined by $x_{i_{0}}^{\prime}=y, x_{i}^{\prime}=x_{i}$ for $i \neq i_{0}$ is a transcendental basis of A over K.

Proof By maximality of $\left(x_{1}, \ldots, x_{k}\right)$, there is a nonzero polynomial $\mathrm{P}_{0} \in \mathrm{~K}\left[\mathrm{X}_{1}, \ldots\right.$, $\left.\mathrm{X}_{k}, \mathrm{Y}\right]$ such that $\mathrm{P}\left(x_{1}, \ldots, x_{k}, y\right)=0$, and $\mathrm{P} \notin \mathrm{K}[\mathrm{Y}]$ since $y$ is transcendental. Choose $i_{0}$ such that the degree of P with respect to $\mathrm{X}_{i_{0}}$ is nonzero, and define $x_{i}^{\prime}$ as in the lemma. Then $x_{i_{0}}$ is algebraic over $\mathrm{K}\left(x_{1}^{\prime}, \ldots, x_{k}^{\prime}\right)$, and so are all elements of A . The family $\left(x_{i}\right)_{i \neq i_{0}}$ is algebraically free, and $y$ is transcendental over $\mathrm{K}\left(x_{i}\right)_{i \neq i_{0}}$; otherwise $x_{i_{0}}$ would be algebraic over this field. Hence $\left(x_{i}^{\prime}\right)$ is algebraically free, and thus a transcendental basis.

Proof of the Proposition The proof is by induction on $k$.
( $\mathrm{E}_{k}$ ) For any extension A of a field K , if A has a transcendental basis with $k$ elements, all other transcendental bases of A also have $k$ elements.

Statement $\left(\mathrm{E}_{0}\right)$ is immediate: if the empty family is a transcendental basis, then A is algebraic over K , and every transcendental basis is empty. Let $k>0$, and assumed that $\left(\mathrm{E}_{k-1}\right)$ holds. We prove $\left(\mathrm{E}_{k}\right)$. Let A be an extension of K with two transcendental bases $\left(x_{i}\right)_{i \in \mathrm{I}}$ and $\left(y_{j}\right)_{j \in \mathrm{~J}}$ with card $\mathrm{I}=k$. Then $\mathrm{J} \neq \varnothing$; choose $j_{0} \in \mathrm{~J}$. The element $y_{j_{0}}$ is transcendental over K , and it follows from the lemma that there exists $i_{0}$ such that $\left(x_{i}\right)_{i \neq i_{0}}$ is a transcendental basis of A over $\mathrm{K}\left(y_{j_{0}}\right) . \operatorname{But}\left(y_{j}\right)_{j \neq j_{0}}$ is also a transcendental basis of A over $\mathrm{K}\left(y_{j_{0}}\right)$. By the induction hypothesis, $\operatorname{card}\left(\mathrm{J}-\left\{j_{0}\right\}\right)=k-1$, and so $\operatorname{card} \mathrm{J}=k$.

This proves the proposition when A has a finite transcendental basis. Assume that A has two transcendental bases $\left(x_{i}\right)_{i \in \mathrm{I}}$ and $\left(y_{j}\right)_{j \in \mathrm{~J}}$ with I and J infinite. For each $j \in \mathrm{~J}$, it is possible to find a finite subset $\mathrm{I}_{j}$ of I such that $y_{j}$ is algebraic over $\mathrm{K}\left(x_{i}\right)_{i \in \mathrm{I}_{j}}$. The extension A is algebraic over $\mathrm{K}\left(x_{i}\right)_{i \in \mathrm{I}^{\prime}}$, where $\mathrm{I}^{\prime}=\bigcup \mathrm{I}_{j}$, and so $\mathrm{I}^{\prime}=\mathrm{I}$. Let $\mathrm{S} \subset \mathrm{I} \times \mathrm{J}$ be the set of $(i, j)$ such that $i \in \mathrm{I}_{j}$. As J is an infinite set, card $\mathrm{S}=\operatorname{card} \mathrm{J}$. The projection of $S$ onto I is I. So card $S \geqslant$ card I, and card $J \geqslant$ card I. Similarly $\operatorname{card} \mathrm{I} \geqslant \operatorname{card} \mathrm{J}$, and thusù card $\mathrm{J}=\operatorname{card} \mathrm{I}$,

## Exercises 5.1 (Extensions)

1.-Find equations over $\mathbb{Q}$ for the following numbers:

$$
\begin{gathered}
\sqrt{2}+\sqrt{3}, \sqrt{2}+\sqrt[3]{2},\left(1+\frac{\sqrt{5}}{2}\right)\left(1+\frac{i \sqrt{3}}{2}\right) \\
x+\sqrt{2} \quad \text { where } \quad x^{3}+x+1=0
\end{gathered}
$$

2.-Show that, if $x$ has equation P , then

$$
\sum \frac{(-1)^{r}}{r!s!} d^{\frac{r+s}{2}} \mathbf{P}^{(r)} \mathbf{P}^{(s)}
$$

is an equation for $x+\sqrt{d}$, the sum being over all pairs $(r, s)$ of integers $0 \leqslant r \leqslant$ $\mathrm{d}^{\circ}(\mathrm{P}), 0 \leqslant s \leqslant \mathrm{~d}^{\circ}(\mathrm{P})$ with $r+s$ even, and $\mathrm{P}^{(r)}$ denoting the $r$-th derivative of P . Interpret this formula when the field K has nonzero characteristic.
3.-(a) Let $B$ be an integral ring, $P$ and $Q$ polynomials in $B[X]$ of respective degrees $p$ and $q$, and $\Omega$ the closure of the field of fractions of B . Show that the following conditions are equivalent:
(i) P and Q have a common root in $\Omega$;
(ii) there are nonzero polynomials $u$ with degree $\leqslant q-1$ and $v$ with degree $\leqslant$ $p-1$ such that $u \mathrm{P}+v \mathrm{Q}=0$.

Let $\mathrm{B}[\mathrm{X}]_{m}$ be the B -module of polynomials of degrees $\leqslant m$ (isomorphic to $\mathrm{B}^{m+1}$ ), and $\operatorname{Res}(\mathrm{P}, \mathrm{Q})$ the resultant of P and Q (3.7.11).
(b) For $\mathrm{P}, \mathrm{Q} \in \mathrm{K}[\mathrm{X}, \mathrm{Y}, \mathrm{Z}]$, define the resultant $\operatorname{Res}_{\mathrm{Z}}(\mathrm{P}, \mathrm{Q}) \in \mathrm{K}[\mathrm{X}, \mathrm{Y}]$ of P and Q as elements of $(\mathrm{K}[\mathrm{X}, \mathrm{Y}])[\mathrm{Z}]$.

Let A be a K-algebra, $x, y \in \mathrm{~A}$ algebraic elements over K having respectively P (of degree $p$ ) and Q (of degree $q$ ) as equations. Show that

$$
\begin{aligned}
& \mathrm{R}_{1}=\operatorname{Res}_{\mathrm{X}}(\mathrm{P}(\mathrm{X}), \mathrm{Q}(\mathrm{Z}-\mathrm{X})) \\
& \mathrm{R}_{2}=\operatorname{Res}_{\mathrm{X}}\left(\mathrm{P}(\mathrm{X}), \mathrm{X}^{q} \mathrm{Q}\left(\frac{\mathrm{Z}}{\mathrm{X}}\right)\right)
\end{aligned}
$$

are respectively equations for $x+y$ and $x y$ over K .
4.-(a) Draw the curve in $\mathbb{R}^{2}$ with parametric representation:

$$
\left\{\begin{array}{l}
x=\cos 2 t \\
y=\cos 3 t
\end{array}\right.
$$

Show that this curve is a subset of an algebraic curve and give its equation.
(b) Let $\mathbb{R}(\cos 2 \mathrm{~T}, \cos 3 \mathrm{~T})$ be the field of fractions of the ring of functions generated by the real constants and the functions $t \mapsto \cos 2 t$ and $t \mapsto \cos 3 t$. Show that this field is a finite extension of $\mathbb{R}(\cos 2 \mathrm{~T})$. Give an equation for $\cos 3 \mathrm{~T}$.
5.-Let $K$ be a field and $M$ a commutative monoid. Consider the vector space $K^{M}$ of maps from M to K. Show that in this vector space the homomorphisms from M to the multiplicative monoid K are linearly independent (consider the algebra $\mathrm{K}^{(\mathrm{M})}$ of M ). Does this result generalize to non-commutative monoids?
6.-Let A be a K -algebra, B and C subalgebras of A . If the canonical map $\mathrm{B} \otimes_{\mathrm{K}} \mathrm{C} \rightarrow \mathrm{A}$ is injective, then B and C are said to be linearly disjoint (over K ). Let $L$ be an extension of $K$ in $B$. Show that $B$ and $C$ are linearly disjoint over $K$ if and only if this is also the case of L and C over B as well as that of B and the algebra $\mathrm{L} \cdot \mathrm{C}$ generated by L and C over L .
7.-A non-trivial ring is said to be local if it has a unique maximal ideal. The aim is to show that every finite K -algebra is a unique product of local K -algebras. Let A be a finite K-algebra and $\left(\mathfrak{m}_{i}\right)_{i \in \mathrm{I}}$ the family of maximal ideals. Set $\mathfrak{N}=\bigcap \mathfrak{m}_{i}$.
(a) Show that $\mathfrak{N}^{r}=0$ for some $r$. Deduce that the product ideal of $\mathfrak{m}_{i}^{r}$ is 0 .
(b) Show that the annihilator of $x \in \mathfrak{m}_{i}^{r}$ is not contained in $\mathfrak{m}_{i}$. Deduce that $\cap \mathfrak{m}_{i}^{r}=$ 0 .
(c) Show that $\mathfrak{m}_{i}^{r}+\mathfrak{m}_{j}^{r}=\mathrm{A}$ for $j \neq i$. Deduce that the canonical map $\mathrm{A} \rightarrow \prod \mathrm{A} / \mathfrak{m}_{i}^{r}$ is an isomorphism. Show that $\mathrm{A} / \mathfrak{m}_{i}^{r}$ is local.
(d) Let $f: \mathrm{A} \rightarrow \prod_{j \in \mathrm{~J}} \mathrm{~B}_{j}$ be an isomorphism, where $\mathrm{B}_{j}$ is a local algebra for all $j$. Show that there is a bijection $s: \mathrm{I} \rightarrow \mathrm{J}$ and an isomorphism from $\mathrm{A} / \mathfrak{m}_{i}^{r}$ onto $\mathrm{B}_{s(i)}$ for all $i$.
(e) Let $\left(k_{i}\right)$ be integers. Show that the product ideal of $\mathfrak{m}_{i}^{k_{i}}$ is equal to their intersection.
8.-Let A and B be finite algebras over $\mathrm{K}, f: \mathrm{A} \rightarrow \mathrm{B}$ an injective homomorphism, and $\left(\mathfrak{m}_{i}\right)_{i \in \mathrm{I}}$ and $\left(\mathfrak{n}_{j}\right)_{j \in \mathrm{~J}}$ the respective families of maximal ideals of A and B. Show that there is a surjective map $f^{*}: \mathrm{J} \rightarrow \mathrm{I}$ and, for all $j$, a homomorphism

$$
f_{j}: \mathrm{A} / \mathfrak{m}_{f^{*}(j)} \rightarrow \mathrm{B} / \mathfrak{n}_{j}
$$

making the diagram

where $f_{*}$ is defined by all $f_{j}$, commutative. Show that $f^{*}$ and all $f_{j}$ are uniquely determined.
9.-Let $\mathrm{P}= \pm 1+\prod_{i=0}^{k}\left(\mathrm{X}-a_{i}\right)$ be a polynomial in $\mathbb{Z}[\mathrm{X}]$, where $a_{i} \in \mathbb{Z}$ take at least three distinct values for $i \neq 0, a_{0}=0$, and $\left|a_{i}\right| \geqslant 2$.
(a) Show that P has a root in $\mathbb{C}$ with absolute value $<1$, while the other roots have absolute value $>1$.
(b) Show that P is irreducible in $\mathbb{Q}[\mathrm{X}]$.
(c) Can $\mathbb{Z}$ be replaced by $\mathbb{Z}[i]$ ? By what other subring of $\mathbb{C}$ can it replaced?
10. (Liouville numbers)-(a) Let P be a polynomial of degree $d$ with integral coefficients and $x \in \mathbb{R}$ a root of P. Let $\left(\frac{a_{n}}{b_{n}}\right)\left(a_{n} \in \mathbb{Z}, b_{n} \in \mathbb{N}^{*}\right)$ a sequence of rational numbers tending to $x$ while remaining distinct from $x$. Show the existence of a constant $c$ such that

$$
\left|\frac{a_{n}}{b_{n}}-x\right| \geq c\left(\frac{1}{b_{n}}\right)^{d}
$$

(note that $\left.b_{n}^{d} \mathrm{P}\left(\frac{a_{n}}{b_{n}}\right) \in \mathbb{Z}\right)$.
(b) Show that if $\left(u_{n}\right)$ is a sequence of integers equal to 1 or 2 , then the real number $\sum \frac{u_{n}}{10^{n i}}$ is transcendental, i.e. is not algebraic over $\mathbb{Q}$. Show that such numbers form a set having the cardinality of the continuum, and homeomorphic to the Cantor set (2.8, Exercise 2).
(c) Is the criterion obtained in (a), applied to partial sums of the series $\sum \frac{1}{n!}$, sufficient to show that $e$ is transcendental?
11. (Transcendence of e)—(a) For $\mathrm{P} \in \mathbb{R}[\mathrm{X}]$, set

$$
\mathrm{L}(\mathrm{P})=\int_{0}^{\infty} e^{-x} \mathrm{P}(x) d x
$$

Show that if P is of degree $\leqslant n$, then $\mathrm{L}(\mathrm{P})=\mathrm{P}(0)+\mathrm{P}^{\prime}(0)+\cdots+\mathrm{P}^{(n)}(0)$, where $\mathrm{P}^{(k)}$ denotes the $k$-th derivative of P .
(b) Let $n$ be a fixed integer and for $p \in \mathbb{N}$, set

$$
\mathrm{P}_{p}(x)=\frac{1}{(p-1)!} x^{p-1}(x-1)^{p} \cdots(x-n)^{p}
$$

Show that $\mathrm{M}_{p}=\mathrm{L}\left(\mathrm{P}_{p}\right)$ is an integer congruent to $\pm(n!)^{p}(\bmod p)$.
(c) Let $\tau_{-k} \mathrm{P}$ denote the polynomial $x \mapsto \mathrm{P}(x+k)$. Set

$$
\mathrm{M}_{p, k}=\mathrm{L}\left(\tau_{-k} \mathrm{P}_{p}\right)=e^{k} \int_{k}^{\infty} e^{-x} \mathrm{P}_{p}(x) d x
$$

Show that for $p \leqslant k \leqslant n, \mathbf{M}_{p, k}$ is an integral multiple of $p$.
(d) Suppose that $a_{n} e^{n}+a_{n-1} e^{n-1}+\cdots+a_{0}=0$. Then

$$
\sum a_{k}\left(\mathbf{M}_{p, k}+\epsilon_{p, k}\right)=0
$$

where $\epsilon_{p, k}=\frac{e^{k}}{(p-1)!} \int_{0}^{k} \mathrm{P}_{p}(x) d x$. Show that $\epsilon_{p, k}$ tends to 0 as $p$ tends to infinity.
(e) Suppose that $a_{n} e^{n}+\cdots+a_{0}=0$ with $a_{i} \in \mathbb{Z}$. Show that $\sum a_{k} \mathbf{M}_{p, k}=0$ for sufficiently large $p$.

Assuming $a_{0} \neq 0$, find a contradiction by taking $p$ to be a sufficiently large prime.

### 5.2 Algebraic Extensions

### 5.2.1

Definition An extension $A$ of $K$ is algebraic if every element of $A$ is algebraic over K or, equivalently, if as an algebra, A is generated by algebraic elements over K .

Every finite extension is algebraic.
Proposition Every subalgebra of an algebraic extension is an extension.
Proof Let A be an algebraic extension of K and B a subalgebra of A . The extension A is the union of finite sub-extensions $\left(\mathrm{A}_{i}\right), i \in \mathrm{I}$. For all $i \in \mathrm{I}, \mathrm{B} \cap \mathrm{A}_{i}$ is a finite integral algebra and hence an extension (5.1.4, Proposition) and $B=\bigcup\left(B \cap A_{i}\right)$ is an extension.

### 5.2.2

Proposition Let A be an algebraic extension of K. Every endomorphism of A is an automorphism.

Proof Let $f: \mathrm{A} \rightarrow \mathrm{A}$ be an extension endomorphism. Then $f$ is injective since A is a field. For all $P \in K[X]$, let $R_{P}$ be the set of roots of $P$ in $A$. The endomorphism $f$ induces an injective map from $\mathrm{R}_{\mathrm{P}}$ in $\mathrm{R}_{\mathrm{P}}$, which is therefore surjective since $\mathrm{R}_{\mathrm{P}}$ is finite. As $\mathrm{A}=\bigcup_{\mathrm{P} \in \mathrm{K}[\mathrm{X}]} \mathrm{R}_{\mathrm{P}}, f$ is surjective.

### 5.2.3

Proposition and Definition The following conditions are equivalent:
(i) all polynomials $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ of degree $>0$ have at least one root in K ;
(ii) all polynomials can be uniquely written as $\alpha\left(\mathrm{X}-a_{1}\right) \ldots\left(\mathrm{X}-a_{d}\right)$ with $\alpha, a_{1}$, $\ldots, a_{d} \in \mathrm{~K}$;
(iii) every algebraic extension over K is trivial.

If these conditions hold, then K is said to be algebraically closed.
Proof (ii) $\Rightarrow$ (i) is immediate. We show (i) $\Rightarrow$ (iii).
Let A be an extension of $\mathrm{K}, x \in \mathrm{~A}$. Embed K in A by $\iota_{\mathrm{A}}$ and let P be the minimal polynomial of $x, a$ a root of P . Then $\mathrm{P}=(\mathrm{X}-a) \mathrm{Q}$ with $\mathrm{d}^{\circ} \mathrm{Q}<\mathrm{d}^{\circ} \mathrm{P}, \mathrm{P}(x)=0$ and $\mathrm{Q}(x) \neq 0$. Hence $x-a=0$ since A is integral and $x=a \in \mathrm{~K}$. Thus $\mathrm{A}=\mathrm{K}$.
(iii) $\Rightarrow$ (ii). As $K[X]$ is factorial, it suffices to show that all irreducible polynomials have degree 1. Let P be an irreducible polynomial. The K -algebra $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$ is finite, and integral since $\mathrm{K}[\mathrm{X}]$ is factorial. So it is an algebraic extension of K (5.1.4, Proposition). It is trivial by assumption, so that P has degree 1 ,

### 5.2.4 Example. d'Alembert's Theorem

The field $\mathbb{C}$ of complex numbers is algebraically closed.
Proof Let $\mathrm{P} \in \mathbb{C}[Z]$ be a polynomial of degree $\geqslant 1$ and suppose that P has no roots in $\mathbb{C}$. Then the function $f=\frac{1}{\mathrm{P}}$ is holomorphic on $\mathbb{C}$, and Cauchy's formula applies:

$$
f(x)=\frac{1}{2 \pi} \int_{0}^{2 \pi} f\left(x+r e^{i \theta}\right) d \theta
$$

As $r$ tends to infinity, the second term tends to 0 ; so $f(x)=0$, and $f$ vanishes everywhere, which is absurd.

For other proofs, see 4.3, Exercise 10, 4.4, Exercise 5, 4.6, Exercise 4 and 6.1, Exercise 3 (b).

### 5.2.5

Definition An algebraic closure of K is an algebraic extension of K that is an algebraically closed field.

Theorem Every field has an algebraic closure.
Cardinality Lemma For any algebraic extension A of K,

$$
\operatorname{CardA} \leqslant \operatorname{CardA}(\mathrm{K}[\mathrm{X}] \times \mathbb{N})
$$

Indeed, set $\Gamma=\{(\mathrm{P}, x) \mid \mathrm{P} \in \mathrm{K}[\mathrm{X}]-\{0\}, x \in \mathrm{~A}$ and $\mathrm{P}(x)=0\}$. The map $(\mathrm{P}, x) \mapsto \mathrm{P}$ from $\Gamma$ to $\mathrm{K}[\mathrm{X}]$ has finite fibres since every polynomial has only finitely many roots. So

$$
\operatorname{card} \Gamma \leqslant \operatorname{card}(\mathrm{K}[\mathrm{X}] \times \mathbb{N})
$$

The map $(\mathrm{P}, x) \mapsto x$ from $\Gamma$ to A is surjective since A is an algebraic extension of $K$. Hence card $A \leqslant \operatorname{card} \Gamma$, proving the lemma.

Proof of the Theorem Let U be a set containing K and such that

$$
\operatorname{card} \mathrm{U}>\operatorname{card}(\mathrm{K}[\mathrm{X}] \times \mathbb{N})
$$

Let $\mathscr{E}$ be the set of algebraic extensions of K whose underlying set is in U . Define the following order relation on $\mathscr{E}$ : " $\mathrm{A} \leqslant \mathrm{B}$ " if A is a sub-extension of B , i.e. the underlying set of A is in the underlying set of B and the laws of A are induced by those of B.

The ordered set $\mathscr{E}$ is inductive. Let $\Omega$ be a maximal element. We show by contradiction that $\Omega$ is algebraically closed.

Let W be a non trivial algebraic extension of $\Omega$. By the cardinality lemma, $\operatorname{card} \mathrm{W} \leqslant \operatorname{card}(\mathrm{K}[\mathrm{X}] \times \mathbb{N})$. Writing $\iota$ for the canonical injection from $\Omega$ to W ,

$$
\operatorname{card}(\mathrm{W}-\iota(\Omega)) \leqslant \operatorname{card} \mathrm{W} \leqslant \operatorname{card}(\mathrm{~K}[\mathrm{X}] \times \mathbb{N})<\operatorname{card} \mathrm{U}=\operatorname{card}(\mathrm{U}-\Omega)
$$

Let $\varphi$ be a bijection from $\mathrm{W}-\iota(\Omega)$ onto a subset F of $\mathrm{U}-\Omega$, and $\Omega^{\prime}=\Omega \cup \mathrm{F}$. Define $\Phi: \mathrm{W} \rightarrow \Omega^{\prime}$ by $\Phi(x)=\iota^{-1}(x)$ if $x \in \iota(\Omega)$ and $\Phi(x)=\varphi(x)$ if $x \in \mathrm{~W}-\iota(\Omega)$. The map $\Phi$ is a bijection. Use $\Phi$ to transfer the extension structure of W to $\Omega^{\prime}$. Then $\Omega^{\prime}$ becomes a non trivial algebraic extension of $\Omega$ containing $\Omega$, which contradicts the maximality of $\Omega$. Hence $\Omega$ is algebraically closed.

### 5.2.6

Theorem Let A be an algebraic extension of $\mathrm{K}, \mathrm{B}$ a sub-extension of A and $\Omega$ an algebraic closure of K. Every homomorphism from B to $\Omega$ extends to a homomorphism from A to $\Omega$.

Proof Let $\mathscr{E}$ be the set of pairs (C, $h$ ) where C is a subalgebra of A and $h: \mathrm{C} \rightarrow \Omega$ a homomorphism. Equip $\mathscr{E}$ with the order defined by: " $(\mathrm{C}, h) \leqslant\left(\mathrm{C}^{\prime}, h^{\prime}\right)$ " if and only if $\mathrm{C} \subset \mathrm{C}^{\prime}$ and $\left.h^{\prime}\right|_{\mathrm{C}}=h$. The ordered set $\mathscr{E}$ is inductive. Let $f: \mathrm{B} \rightarrow \Omega$ be a homomorphism and $(\overline{\mathrm{C}}, \bar{h})$ a maximal element in $\mathscr{E}$ and an upper bound of ( $\mathrm{B}, f$ ).

We show by contradiction that $\overline{\mathrm{C}}=\mathrm{A}$. Let $x \in \mathrm{~A}-\overline{\mathrm{C}}$ and $\mathrm{P}=\sum a_{i} \mathrm{X}^{i}$ be the minimal polynomial of $x$ over $\overline{\mathrm{C}}$. The polynomial $\sum \bar{h}\left(a_{i}\right) \mathrm{X}^{i} \in \Omega[\mathrm{X}]$ has at least one root $\xi$ since $\Omega$ is algebraically closed. The homomorphism $g: \overline{\mathrm{C}}[x]_{\mathrm{A}} \rightarrow \Omega$ defined by $g(x)=\xi$ extends $h$ and $\left(\overline{\mathrm{C}}[x]_{\mathrm{A}}, g\right)$ is a strict upper bound of the maximal element $(\overline{\mathrm{C}}, \bar{h})$, giving a contradiction.

### 5.2.7

Corollary 5.3 Let A be an algebraic extension of K and $\Omega$ an algebraic closure of K . Then there is an embedding of A into $\Omega$.

This follows from the above Theorem 5.2 .6 applied to $(\mathrm{K}, \iota)$, where $\iota: \mathrm{K} \rightarrow \Omega$ is the canonical injection.

Corollary 5.4 Let $\Omega$ be an algebraic closure of $\mathrm{K}, \mathrm{A}$ and B sub-extensions of $\Omega$. Every homomorphism from A to B extends to an automorphism of $\Omega$. In particular, every automorphism of A extends to an automorphism of $\Omega$.

Indeed, by Theorem 5.2.6, every homomorphism from A to B extends to an endomorphism of $\Omega$ that is bijective by 5.2.2.

Corollary 5.3 and Definition Let $\Omega$ be an algebraic closure of K , and $x, x^{\prime} \in \Omega$. The following conditions are equivalent:
(i) $x$ and $x^{\prime}$ have the same minimal polynomial;
(ii) there is an automorphism $g: \Omega \rightarrow \Omega$ such that $g(x)=x^{\prime}$.

If these conditions hold, then $x$ and $x^{\prime}$ are said to be conjugate.
Proof (i) $\Rightarrow$ (ii). Sub-extensions $\mathrm{K}[x]_{\Omega}$ and $\mathrm{K}\left[x^{\prime}\right]_{\Omega}$ of $\Omega$ are isomorphic to $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$, where P is the minimal polynomial of both $x$ and $x^{\prime}$.

The isomorphism $f: x \mapsto x^{\prime}$ from $\mathrm{K}[x]_{\Omega}$ to $\mathrm{K}\left[x^{\prime}\right]_{\Omega}$ extends to an automorphism of $\Omega$ (5.2.7. Corollary 5.2).
(ii) $\Rightarrow$ (i). Every equation of $x$ is an equation of $x^{\prime}$, proving (i).

Corollary 5.4 All algebraic closures of K are isomorphic.

Indeed, let $\Omega$ and $\Omega^{\prime}$ be two algebraic closures of $\mathrm{K}, f: \Omega \rightarrow \Omega^{\prime}$ and $g: \Omega^{\prime} \rightarrow \Omega$ homomorphisms (these exist by Corollary 1). The endomorphisms $g \circ f$ and $f \circ g$ of $\Omega$ and $\Omega^{\prime}$ respectively are automorphisms by 5.2 .2 , and so $f$ and $g$ are isomorphisms.

Remark Given two algebraic closures of a field, there is no method to choose an isomorphism from the one onto the other. Hence algebraic closures of K should not be identified; the notion of "the algebraic closure" of K is not well-defined. This is similar to universal coverings in the context of topology (5.2, Exercise 2).

### 5.2.8 Characteristic of a Field and the Frobenius Endomorphism

We remind the reader that the characteristic of a field K is the positive generator $p$ of the kernel of the homomorphism $\epsilon: n \mapsto n \cdot 1$ from $\mathbb{Z}$ to K . If $\epsilon$ injective, then $p=0$ and the subfield of K generated by 1 is isomorphic to $\mathbb{Q}$. If $\epsilon$ is not injective, the ring $\mathbb{Z} /(p)$, isomorphic to a subring of K , is integral and so $p$ is prime and $\mathbb{Z} /(p)$ is isomorphic to a subfield of K ; if $p \neq 0, p$ is also the order of 1 in the additive group K.

For $p$ prime, the field $\mathbb{Z} /(p)$ is written $\mathbb{F}_{p}$.
Theorem and Definition Let K be a field of characteristic $p \neq 0$. Then the map $\mathrm{F}: x \mapsto x^{p}$ is an endomorphism of K , called the Frobenius endomorphism of K .

The theorem follows from the binomial formula and from the fact that the binomial coefficients $\binom{p}{k}=\frac{p!}{k!(p-k)!}$ are multiples of $p$ for $0<k<p$.

### 5.2.9 Finite Fields

If K is a finite field, its characteristic is a prime $p \neq 0$ and K is a finite extension of $\mathbb{F}_{p}$. If $r$ is the degree of K over $\mathbb{F}_{p}$, then

$$
\operatorname{card} \mathrm{K}=p^{r}
$$

Proposition Let $p$ be a prime, and $q=p^{r}(r \geqslant 1)$ a power of $p$. There is a field with $q$ elements, and all such fields are isomorphic.

Proof Let $\Omega$ be an algebraic closure of $\mathbb{F}_{p}$. The endomorphism $\mathrm{F}^{r}: x \mapsto x^{q}$ de $\Omega$ is an automorphism. Hence the elements $x \in \Omega$ left invariant by $\mathrm{F}^{r}$ form a subfield $\mathbb{F}_{q}$ of $\Omega$ whose elements are the roots of the polynomial $\mathrm{P}=\mathrm{X}^{q}-\mathrm{X}$. The derivative of P is

$$
q \mathrm{X}^{q-1}-1=p p^{r-1} \mathrm{X}^{q-1}-1=-1 \neq 0
$$

The polynomial P has $q$ simple roots which are therefore distinct and $\mathbb{F}_{q}$ has $q$ elements.

Let K be a field with $q$ elements, the multiplicative group $\mathrm{K}^{*}$ has $q-1$ elements and for all $x \in \mathrm{~K}^{*}, x^{q-1}=1$; so $x^{q}=x$ for all $x \in \mathrm{~K}$. As K is a finite extension of $\mathbb{F}_{p}$, it may be embedded in $\Omega$, and thus is isomorphic to a subfield of $\mathbb{F}_{q}$, and is isomorphic to $\mathbb{F}_{q}$ since card $\mathrm{K}=\operatorname{card} \mathbb{F}_{q}$.

Remark The field $\mathbb{F}_{q}$ constructed here depends on the choice of the algebraic closure $\Omega$ of $\mathbb{F}_{p}$. We do not give any intrinsic definition of $\mathbb{F}_{q}$ for $q=p^{\nu}, \nu>1$. There is no natural isomorphism between two fields with $q$ elements.

### 5.2.10

Theorem If K is a finite field, then the multiplicative group $\mathrm{K}^{*}$ is cyclic.
We prove the following more general result:
Proposition Let K be a field. Every finite subgroup G of the multiplicative group $\mathrm{K}^{*}$ is cyclic.

Proof The multiplicative group G is isomorphic to an additive group $\bigoplus \mathbb{Z} / p_{i}^{r_{i}}$, where all $p_{i}$ are primes (3.5.10, Proposition). We show that all $p_{i}$ are distinct. Suppose that for $i \neq j p_{i}=p_{j}=p$, then the group G would contain a subgroup of the form $(\mathbb{Z} /(p))^{2}$, and hence at least $p^{2}$ elements $x$ such that $x^{p}=1$. This is impossible since the polynomial $\mathrm{X}^{p}-1$ has at most $p$ roots. By (3.5.9, Corollary), G is isomorphic to $\mathbb{Z} / \prod p_{i}^{r_{i}}$, and so is cyclic.

Remark When $\mathrm{K}=\mathbb{C}$, the finite subgroups can be explicitly described. The map $t \mapsto e^{2 i \pi t}$ defines an isomorphism from $\mathbb{Q} / \mathbb{Z}$ onto the group of elements of finite order of $\mathbb{C}^{*}$. Every subgroup of $\mathbb{Q} / \mathbb{Z}$ is generated by the image of a number $\frac{1}{q}$, and so every finite subgroup of $\mathbb{C}^{*}$ is of the form $\left\{e^{2 i \pi \frac{p}{q}}\right\}_{0 \leqslant p \leqslant q-1}$.

## Exercises 5.2. (Algebraic extensions)

1.-Find a generator of the group $(\mathbb{Z} / p \mathbb{Z})^{*}$ for each prime $p<100$. For which primes is 2 a generator?
2. (Non-functoriality of the algebraic closure)-(a) Consider the complex fields $\mathrm{K}=\mathbb{Q}(i), \mathrm{L}=\mathrm{K}(\sqrt{2}), \mathrm{M}=\mathrm{K}(\sqrt[4]{2})$ and the field $\Omega$ of algebraic numbers over $\mathbb{Q}$. Let $f \in \operatorname{Aut}_{\mathrm{K}}(\mathrm{L})$ be defined by $f(\sqrt{2})=-\sqrt{2}$. Determine $\mathrm{Aut}_{\mathrm{K}}(\mathrm{M})$. Show that there is an automorphism $g$ of M extending $f$, but none satisfying $g \circ g=1_{\mathrm{M}}$.
(b) Show that every automorphism $\Omega$ preserves M. Show that there is no automorphism $h$ of $\Omega$ inducing $f$ and satisfying $h \circ h=1_{\Omega}$.
(c) Let $\mathscr{C}$ be the category of fields. Show that there is no functor $\mathrm{S}: \mathscr{C} \rightarrow \mathscr{C}$ and no functorial morphism $\iota: 1_{\mathscr{C}} \rightarrow \mathrm{S}$ such that, for all $\mathrm{K} \in \mathscr{C}$, the field $\mathrm{S}(\mathrm{K})$, considered an extension of K by $\iota_{\mathrm{K}}: \mathrm{K} \rightarrow \mathrm{S}(\mathrm{K})$, is an algebraic closure of K . (Note the similarity with 4.4, Exercise 1.)
3.-Let $p$ be a prime and $\Omega$ an algebraic closure of the field $\mathbb{F}_{p}=\mathbb{Z} /(p)$.
(a) Show that if $p$ does not divide $n$, then the equation $x^{n}-1=0$ has $n$ distinct roots in $\Omega$.
(b) Let $r$ be an integer $>0$. Show that the set of fixed points of $\mathrm{F}^{r}$, where F is the Frobenius automorphism of $\Omega$, is a subfield $\mathbb{F}_{p^{r}}$ of $\Omega$ with $p^{r}$ elements.
(c) Show that $\mathbb{F}_{p^{r}}$ is the unique subfield of $\Omega$ with $p^{r}$ elements. Deduce that every finite subfield of $\Omega$ is of the form $\mathbb{F}_{p^{r}}$, and that every finite field of characteristic $p$ is isomorphic to some $\mathbb{F}_{p^{r}}$.
(d) Show that the group of automorphisms from $\mathbb{F}_{p^{r}}$ onto $\mathbb{F}_{p}$ is generated by the Frobenius automorphism, and is isomorphic to $\mathbb{Z} /(r)$.
(e) Show that $\mathbb{F}_{p^{r}}$ is contained in $\mathbb{F}_{p^{s}}$ if and only if $r$ divides $s$. Are there infinite subfields in $\Omega$ other than $\Omega$ ? Show that the group of automorphisms from $\Omega$ onto $\mathbb{F}_{p}$ is isomorphic to the profinite completion $\widehat{\mathbb{Z}}$ of $\mathbb{Z}$ (see 2.9.5).
(f) Show that the multiplicative group of $\Omega$ is isomorphic to the direct sum of the $l$-primary components $\mathrm{M}_{l}$ of $\mathbb{Q} / \mathbb{Z}$ for some prime $l \neq p\left(\mathrm{M}_{l}\right.$ is the subgroup of $\mathbb{Q} / \mathbb{Z}$ consisting of the elements whose order is a power of $l$ ).
(g) Show that the ring of endomorphism of $\Omega^{*}$ is isomorphic to $\prod_{l \neq p} \widehat{\mathbb{Z}}_{l}$ (3.3, Exercise 9) and that this gives an injective homomorphism from the additive group $\widehat{\mathbb{Z}}$ to the multiplicative group of invertible elements of this ring.
4. (Another proof of the existence of an algebraic closure)—If $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ is an infinite family of algebras, set $\bigotimes_{i \in \mathrm{I}} \mathrm{A}_{i}=\underset{ }{\lim } \bigotimes_{j \in \mathrm{~J}} \mathrm{~A}_{j}$, the direct limit being taken over the finite subsets J of I . Let $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ be the family of extensions $\mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right] / \mathfrak{p}$ of $K$, where $\mathfrak{p}$ is a maximal ideal of $K\left[X_{1}, \ldots, X_{n}\right]$, and let $\mathfrak{M}$ be a maximal ideal of $\bigotimes_{i \in \mathrm{I}} \mathrm{A}_{i}$. Show that $\bigotimes_{i \in \mathrm{I}} \mathrm{A}_{i} / \mathcal{M}$ is an algebraic closure of K .
5.-Let A be a finite K -algebra, B a subalgebra of A and $\Omega$ an algebraic closure of K. Show that every homomorphism from $B$ to $\Omega$ extends to a homomorphism from A to $\Omega$ (use 5.1 , Exercise 8 ). Does this result generalize to infinite algebras A, but all of whose elements are algebraic over K ?
6.-(a) Let A and B be two integral rings and $f: \mathrm{A} \rightarrow \mathrm{B}$ an injective homomorphism, $M$ the field of fractions of $A$. Assume that $B$, considered an A-algebra by using $f$, is a finitely generated A-module. Show that $\mathrm{M} \otimes_{\mathrm{A}} \mathrm{B}$ can be identified with the field of fractions of B.
(b) Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}, \mathrm{Y}]$ be an irreducible polynomial not in $\mathrm{K}[\mathrm{X}]$. Show that the field of fractions of $\mathrm{K}[\mathrm{X}, \mathrm{Y}] /(\mathrm{P})$ is a natural finite extension of $\mathrm{K}(\mathrm{X})$.

### 5.3 Diagonal Algebras

### 5.3.1 Algebra of Functions on a Finite Set

Let X be a a finite set with $d$ elements. The algebra $\mathrm{K}^{\mathrm{X}}$ of K -valued functions defined on X is a finite K -algebra over K of degree $d$ isomorphic to $\mathrm{K}^{d}$ equipped with the multiplication

$$
\left(\left(x_{1}, \ldots, x_{d}\right),\left(y_{1}, \ldots, y_{d}\right)\right) \mapsto\left(x_{1} y_{1}, \ldots, x_{d} y_{d}\right)
$$

For all $x \in \mathrm{X}$, the map $\delta_{x}: f \mapsto f(x)$ from $\mathrm{K}^{\mathrm{X}}$ to K is an algebra homomorphism, and its kernel $\mathfrak{m}_{x}$ is a maximal ideal.

For all subsets Y of X , the elements $f \in \mathrm{~K}^{\mathrm{X}}$ such that $\left.f\right|_{\mathrm{Y}}=0$ form an ideal $\mathscr{I}_{\mathrm{Y}}$ and the algebra $\mathrm{K}^{\mathrm{X}} / \mathscr{F}_{\mathrm{Y}}$ can be identified with $\mathrm{K}^{\mathrm{Y}}$. If Y and Z are subsets of X , then

$$
\begin{aligned}
& \mathscr{I}_{\mathrm{Y} \cap \mathrm{Z}}=\mathscr{I}_{\mathrm{Y}}+\mathscr{I}_{\mathrm{Z}} \\
& \mathscr{J}_{\mathrm{Y} \mathrm{Z}}=\mathscr{J}_{\mathrm{Y}} \cap \mathscr{I}_{\mathrm{Z}}=\mathscr{J}_{\mathrm{Y}} \cdot \mathscr{I}_{\mathrm{Z}} \\
& \mathrm{Y} \subset \mathrm{Z} \Longleftrightarrow \mathscr{J}_{\mathrm{Y}} \supset \mathscr{I}_{\mathrm{Z}}
\end{aligned}
$$

Proposition With the above notations, the map $\mathrm{Y} \mapsto \mathscr{J}_{\mathrm{Y}}$ from $\mathfrak{P}(\mathrm{X})$ to the set of ideals of $\mathrm{K}^{\mathrm{X}}$ is a bijection.

Proof The map is clearly injective; we show that it is surjective. Let $\mathscr{I}$ be an ideal of $\mathrm{K}^{\mathrm{X}}$ and Y the set of $x \in \mathrm{X}$ such that $f(x)=0$ for all $f \in \mathscr{F}$. For all $x \notin \mathrm{Y}$, there exists $u_{x} \in \mathscr{J}$ such that $u_{x}(x) \neq 0$. Let $e_{x}$ be the element of $\mathrm{K}^{\mathrm{X}}$ defined by $e_{x}(x)=1$ and $e_{x}\left(x^{\prime}\right)=0$ for all $x^{\prime} \neq x$. Then

$$
e_{x}=\frac{e_{x}}{u_{x}(x)} u_{x} \in \mathscr{I}
$$

Any $f \in \mathscr{J}_{\mathrm{Y}}$ can be written

$$
f=\sum_{x \in \mathrm{X}} f(x) e_{x}=\sum_{x \in \mathrm{X}-\mathrm{Y}} f(x) e_{x} \in \mathscr{I} .
$$

Hence $\mathscr{J}_{\mathrm{Y}} \subset \mathscr{J}$; as $\mathscr{J} \subset \mathscr{J}_{\mathrm{Y}}, \mathscr{J}_{\mathrm{Y}}=\mathscr{I}$.
Corollary The map $x \mapsto \mathfrak{m}_{x}$ is a bijection from X onto the set of maximal ideals of $\mathrm{K}^{\mathrm{X}}$ and $\mathrm{x} \mapsto \delta_{x}$ is a bijection from X onto $\operatorname{Hom}_{\mathrm{K}-\operatorname{seg}_{g}\left(\mathrm{~K}^{\mathrm{X}}, \mathrm{K}\right) \text {. }}$

Indeed, the first map is bijective because singletons are the minimal nonempty subsets. Clearly, $x \mapsto \delta_{x}$ is injective; let $f: \mathrm{K}^{\mathrm{X}} \rightarrow \mathrm{K}$ be a homomorphism, $\operatorname{Ker} f$ a maximal ideal and so of the form $\mathfrak{m}_{x}$, the maps $f$ and $\delta_{x}$ agree on $\mathfrak{m}_{x}$ and on $\mathrm{K} \cdot 1$, hence on $\mathrm{K}^{\mathrm{X}}$.

### 5.3.2 The Gelfand Transform

Let A be a K-algebra, and $\mathrm{X}=\operatorname{Hom}_{\mathrm{K}-\operatorname{seg}}(\mathrm{A}, \mathrm{K})$. For all $a \in \mathrm{~A}$, let $\widehat{a}$ be the map $\xi \mapsto \xi(a)$ from X to K. For $a \in \mathrm{~A}, b \in \mathrm{~A}, \lambda \in \mathrm{~K}$,

$$
\widehat{a+b}=\widehat{a}+\widehat{b}, \quad \widehat{a b}=\widehat{a} \widehat{b}, \quad \widehat{\lambda a}=\lambda \widehat{a}, \quad \text { and } \widehat{1}=1
$$

In other words the map $a \mapsto \widehat{a}$ from A to $\mathrm{K}^{\mathrm{X}}$ is an algebra homomorphism. This homomorphism is called the Gelfand transform of A (see 2.6.7).

If A is a finite algebra, then by the Chinese remainder theorem, the Gelfand transform of A is surjective (5.1.5, Theorem and Remark).
Remark Associating to every element $a$ of a set A a function $\widehat{a}$ on a set X of functions on A by $\widehat{a}(\xi)=\xi(a)$ often gives interesting mathematical results, for instance the Gelfand transform, biduality, Čech compactification, Dirac measures.

### 5.3.3

Proposition and Definition Let A be a finite K-algebra of degree d. Set $\mathrm{X}=\operatorname{Hom}_{\mathrm{K}-\operatorname{seg}_{g}}(\mathrm{~A}, \mathrm{~K})$. The following conditions are equivalent:
(i) A is isomorphic to an algebra of functions on a finite set;
(ii) the Gelfand transform $\gamma: \mathrm{A} \rightarrow \mathrm{K}^{\mathrm{X}}$ is an algebra isomorphism;
(iii) for all $a \in \mathrm{~A}, a \neq 0$, there exists $\xi \in \mathrm{X}$ such that $\xi(a) \neq 0$;
(iv) $\operatorname{card} \mathrm{X}=d$;
(v) $\operatorname{card} \mathrm{X} \geqslant d$.

## If these conditions hold, then A is diagonal.

Proof (ii) $\Rightarrow$ (iv) and (i) is obvious.
(i) $\Rightarrow$ (iii) since if $\mathrm{A}=\mathrm{K}^{\mathrm{Y}}$, then for all $a \neq 0$ there exists $x \in \mathrm{Y}$ such that $a(x) \neq 0$; take $\xi=\delta_{x}$.
(iii) $\Rightarrow$ (v): $\gamma$ is injective, proving (v).
(v) $\Rightarrow$ (ii) and (iv) since $\gamma$ is surjective by the Chinese remainder theorem.
(iv) $\Rightarrow$ (v) is obvious.

### 5.3.4 Example. (The Fourier Transform for $\mathbb{Z} /(n)$.

Let $G$ be the group $\mathbb{Z} /(n)$. Equip the $\mathbb{C}$-vector space $\mathbb{C}^{G}$ with the convolution law defined on the canonical basis by

$$
e_{x} * e_{y}=e_{x+y}, \text { which gives } f * g=h \text { with } h(x)=\sum_{y \in \mathrm{G}} f(y) g(x-y)
$$

for $f, g \in \mathbb{C}^{\mathrm{G}}$. Let A denote the algebra $\left(\mathbb{C}^{\mathrm{G}}, *\right)$. We show that A is diagonal.
For all maps $\varphi$ from $G$ to $\mathbb{C}$, there is a unique $\mathbb{C}$-linear map $\xi: \mathrm{A} \rightarrow \mathbb{C}$ such that $\xi\left(e_{x}\right)=\varphi(x)$. It is an algebra homomorphism if and only if $\xi\left(e_{0}\right)=1$ and $\xi\left(e_{x} * e_{y}\right)=\xi\left(e_{x}\right) \xi\left(e_{y}\right)$, i.e. $\varphi(0)=1$ and $\varphi(x+y)=\varphi(x) \varphi(y)$ for $x, y \in \mathrm{G}$. The set $X=\operatorname{Hom}_{\mathbb{C}-\mathscr{A} g}(A, \mathbb{C})$ can thus be identified with $\operatorname{Hom}_{\mathcal{G}_{c}}(G, \mathbb{C}-\{0\})$. For any $n$-th root of unity $\theta$, the map $x \mapsto \theta^{x}$ is a group homomorphism from G to $\mathbb{C}-\{0\}$. This give $n$ distinct homomorphisms from $G$ to $\mathbb{C}-\{0\}$. Hence condition (v) of 5.3.3 holds and A is diagonal.

We express the Gelfand transform of A explicitly. For all $k \in \mathrm{G}$, set $\theta_{k}=e^{\frac{2 i \pi k}{n}}$ and let $\varphi_{k}$ be the homomorphism $p \mapsto \theta_{k}{ }^{p}$ from G to $\mathbb{C}^{*}-\{0\}$. The homomorphism $\xi_{k}$ from A to $\mathbb{C}$ corresponding to $\varphi_{k}$ is defined by

$$
\xi_{k}(f)=\sum_{p \in \mathrm{G}} f(p) \varphi_{k}(p)=\sum_{p \in \mathrm{G}} f(p) e^{\frac{2 i \pi k p}{n}}
$$

The set $X$ contains $\xi_{0}, \ldots, \xi_{n-1}$, and so $X=\left\{\xi_{0}, \ldots, \xi_{n-1}\right\}$ since card $X=n$. Hence X can be identified with G . Thus for all $f \in \mathrm{~A}$,

$$
\hat{f}(k)=\sum_{p \in \mathrm{G}} f(p) e^{\frac{2 i \pi k p}{n}}
$$

and $f \mapsto \hat{f}$ is an isomorphism from the algebra $\left(\mathbb{C}^{\mathrm{G}}, *\right)$ onto the algebra $\left(\mathbb{C}^{\mathrm{G}}, \cdot\right)$, where - denotes function multiplication.

### 5.3.5

Proposition (a) All subalgebras of a diagonal algebra are diagonal.
(b) All quotient algebras of a diagonal algebra are diagonal.
(c) The product algebra of two diagonal algebras is diagonal.
(d) The tensor product algebra of two diagonal algebras is diagonal.
(e) Let A be an algebra, B and C diagonal subalgebras of A . Then the subalgebra $\mathrm{B} \cdot \mathrm{C}$ of A generated by B and C is diagonal.

Proof (a) Condition (iii) of (5.3.3) holds for subalgebras. (b) follows from (5.3.1), (c) and (d) from $\mathrm{K}^{\mathrm{X}} \times \mathrm{K}^{\mathrm{Y}}=\mathrm{K}^{\mathrm{X} \cup \mathrm{Y}}$ and $\mathrm{K}^{\mathrm{X}} \otimes \mathrm{K}^{\mathrm{Y}}=\mathrm{K}^{\mathrm{X} \times \mathrm{Y}}$. Finally (e) follows as $\mathrm{B} \cdot \mathrm{C}$ is isomorphic to a quotient of $\mathrm{B} \otimes \mathrm{C}$.

### 5.3.6

Proposition Let Setf be the category of finite sets, and $\mathscr{D}$ the category of diagonal algebras over K . Then $\mathscr{D}$ is anti-equivalent to $\mathcal{S e t f}$.

Proof Let F be the contravariant functor $\mathrm{X} \mapsto \mathrm{K}^{\mathrm{X}}$ from $\operatorname{set} \boldsymbol{f}$ to $\mathscr{D}$ and G the contravariant functor $\mathrm{A} \mapsto \operatorname{Hom}_{\mathrm{K}-\mathscr{A L}_{\mathscr{g}}}(\mathrm{A}, \mathrm{K})$ from $\mathscr{D}$ to $\operatorname{Setf}$. For all $\mathrm{A} \in \mathscr{D}$, the Gelfand transform $\gamma_{\mathrm{A}}: \mathrm{A} \rightarrow \mathrm{K}^{\mathrm{G}(\mathrm{A})}=\mathrm{F}(\mathrm{G}(\mathrm{A}))$ is an isomorphism, and these $\gamma_{\mathrm{A}}$ form an isomorphism from the identity functor $1_{\mathscr{D}}$ to $\mathrm{F} \circ \mathrm{G}$. For all finite sets X , the $\operatorname{map} \beta_{\mathrm{X}}: x \mapsto \delta_{x}$ from X to $\operatorname{Hom}_{\mathrm{K}-\mathscr{A g}}\left(\mathrm{K}^{\mathrm{X}}, \mathrm{K}\right)=\mathrm{G}(\mathrm{F}(\mathrm{X}))$ is bijective, and these $\beta_{\mathrm{X}}$ form an isomorphism from the identity functor of $\operatorname{Set} \mathcal{f}$ to $\mathrm{G} \circ \mathrm{F}$. The functors $F$ and $G$ are therefore anti-equivalent.

Addenda. Let X and Y be objects of $\boldsymbol{\operatorname { s e t }} \boldsymbol{\mathcal { F }}$ and $f$ a map from X to Y . The homomorphism $f^{*}: \mathrm{K}^{\mathrm{Y}} \rightarrow \mathrm{K}^{\mathrm{X}}$ is surjective (resp. injective) if and only if $f$ is injective (resp. surjective). The subsets (resp. quotient sets) of X are in bijective correspondence with the quotient algebras (resp. with subalgebras) of $\mathrm{K}^{\mathrm{X}}$.

In particular a diagonal algebra only has finitely many subalgebras.

## Exercises 5.3. (Diagonal algebras)

1.-A Boolean ring is a ring A all of whose elements are idempotents (i.e. $x^{2}=x$ for all $x$ ).
(a) Show that all finite Boolean rings are $\mathbb{Z} /(2)$-diagonal algebras.
(b) Let X be a set. Show that $\mathfrak{P}(\mathrm{X})$, equipped with the laws $\Delta$ and $\cap$, where

$$
\mathrm{Y} \Delta \mathrm{Z}=(\mathrm{Y} \cup \mathrm{Z})-(\mathrm{Y} \cap \mathrm{Z})
$$

is a Boolean ring. Show that every Boolean ring is embedded in a ring of this type, and that every finite Boolean ring is of this type with X finite.
2.-(a) Show that all reduced algebras over an algebraically closed field are diagonal.
(b) Show that an algebra over $\mathbb{F}_{p^{r}}$, where $x^{p^{r}}=x$ for all $x$, is diagonal.
3.-Let $G$ be a finite commutative group, and $A$ the convolution $\mathbb{C}^{\mathrm{G}}$.
(a) Show that the algebra A is diagonal.
(b) Set $\widehat{\mathrm{G}}=\operatorname{Hom}_{\operatorname{Ar}_{g}}(\mathrm{~A}, \mathbb{C})$. Show that $\widehat{\mathrm{G}}$ can be identified with $\operatorname{Hom}_{\mathscr{G}_{\boldsymbol{r}}}\left(\mathrm{G}, \mathbb{C}^{*}\right)$ or with $\operatorname{Hom}_{\mathscr{E}_{\boldsymbol{r}}}(\mathrm{G}, \mathbb{Q} / \mathbb{Z})$.
(c) The diagonal map $\mathrm{G} \rightarrow \mathrm{G} \times \mathrm{G}$ defines a homomorphism $\Delta: \mathrm{A} \rightarrow \mathrm{A} \otimes \mathrm{A}$, giving a composition law on $\widehat{G}$. Show that this law agrees with the law given by the identification with $\operatorname{Hom}_{\mathscr{E}_{\mu}}\left(G, \mathbb{C}^{*}\right)$.
(d) Show that $\widehat{\widehat{G}}$ can be identified with $G$.
(e) Show that, in the previous questions, $\mathbb{C}$ can be replaced by an algebraically closed field of characteristic $p$ as long as $p$ does not divide the order of G. For what values of $r$ can $\mathbb{F}_{p^{r}}$ be taken?
(f) For $f, g \in \mathbb{C}^{\mathrm{G}}, \widehat{(f \cdot g)}=\frac{1}{n}(\hat{f} * \hat{g})$, where $n$ is the order of G. Is there a $\mathbb{C}$-linear bijection $f \mapsto f^{\prime}$ from $\mathbb{C}^{\mathrm{G}}$ onto some $\mathbb{C}^{\mathrm{G}^{\prime}}$, where $\mathrm{G}^{\prime}$ is a group such that $(f * g)^{\prime}=f^{\prime} \cdot g^{\prime}$ and $(f \cdot g)^{\prime}=f^{\prime} * g^{\prime}$ for $f, g \in \mathbb{C}^{\mathrm{G}}$ ?

### 5.4 Etale Algebras

### 5.4.1

Definition Let A be a finite K-algebra. An extension L of K is said to diagonalize A if $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ is a diagonal L-algebra.

L diagonalizes A if and only if

$$
\operatorname{card} \operatorname{Hom}_{\mathrm{K}-\operatorname{seg}_{g}}(\mathrm{~A}, \mathrm{~L})=\operatorname{deg}_{\mathrm{K}} \mathrm{~A}
$$

### 5.4.2

Proposition Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ be a polynomial of degree $d$, and L an extension of K .
L diagonalizes $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$, if and only if P has d distinct roots in L .
Proof The algebra $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{K}[\mathrm{X}] /(\mathrm{P})$ is isomorphic to $\mathrm{L}[\mathrm{X}] /(\mathrm{P})$ and the degree of $\mathrm{L}[\mathrm{X}] /(\mathrm{P})$ over L is $d$. The map $\varphi \mapsto \varphi(x)$ is a bijection from $\operatorname{Hom}_{\mathrm{L}-\Delta t g}(\mathrm{~L}[\mathrm{X}] /(\mathrm{P}), \mathrm{L})$ onto the set of roots of P in L .

There are $d$ distinct homomorphisms from $\mathrm{L}[\mathrm{X}] /(\mathrm{P})$ to L if and only if P has $d$ distinct roots.

### 5.4.3

Proposition and Definition Let $\Omega$ be an algebraic closure of K and A a finite K -algebra. The following conditions are equivalent:
(i) there is an extension L of K diagonalizing A ;
(ii) there is a finite extension L of K diagonalizing A ;
(iii) $\Omega$ diagonalizes A .

If these conditions hold, then A is said to be an etale algebra.

Proof (ii) $\Rightarrow$ (iii). Indeed, L embeds in $\Omega$ and $\Omega \otimes_{\mathrm{K}} \mathrm{A}=\Omega \otimes_{\mathrm{L}}\left(\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}\right)$.
(iii) $\Rightarrow$ (i) is obvious.
(i) $\Rightarrow$ (ii). Let $d$ be the degree of A on K . There are $d$ distinct homomorphisms $\xi_{1}, \ldots, \xi_{d}$ from A to L . The subalgebra $\mathrm{L}^{\prime}$ of L generated by $\xi_{1}(\mathrm{~A}), \ldots, \xi_{d}(\mathrm{~A})$ is a finite extension of K contained in L , and $\xi_{1}, \ldots, \xi_{d}$ are $d$ homomorphisms from A to $L^{\prime}$. So L' diagonalizes A.

### 5.4.4

Proposition and Definition Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$. The following conditions are equivalent:
(i) the algebra $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$ is etale;
(ii) the roots of P in $\Omega$ are distinct;
(iii) g.c.d. $\left(\mathrm{P}, \mathrm{P}^{\prime}\right)=1$;
(iv) The discriminant of P is nonzero.

If these conditions hold, then the polynomial P called separable.
(i) $\Leftrightarrow$ (ii) is an immediate consequence of (5.4.2).
(ii) $\Leftrightarrow$ (iii) follows from the characterization of simple roots.
(iii) $\Leftrightarrow$ (iv) follows from (3.7.12) and (3.7.11, Proposition).

Let A be a K-algebra. An algebraic element $x \in A$ over $K$ is said to be separable if its minimal polynomial is separable, in other words if the subalgebra $\mathrm{K}[x]_{\mathrm{A}}$ generated by $x$ is etale. If A is algebraic over K , then A is said to be separable if all elements of A are separable.

### 5.4.5

Proposition (a) All subalgebras of an etale algebra are etale.
(b) All quotient algebras of an etale algebra are etale.
(c) The product algebra of two etale algebras is etale.
(d) The tensor product algebra of two etale algebras is etale.
(e) Let A be an algebra, B and C etale subalgebras of A , then the subalgebra $\mathrm{B} \cdot \mathrm{C}$ of A generated by B and C is etale.

Proof Let L be an extension of K. By (5.3.5), if L diagonalizes an algebra A , then L diagonalizes every subalgebra of A , and every quotient algebra of A . If L diagonalizes algebras $B$ and $C$, then it diagonalizes the product algebra $B \times C$ and the tensor product algebra $\mathrm{B} \otimes \mathrm{C}$.

If L diagonalizes the subalgebras B and C of an algebra A , then it also diagonalizes B.C.

Corollary 5.5 An algebra is etale if and only if it is finite and separable.
Corollary 5.2 and Definition Let A be a K-algebra. The algebraically separable elements over K of A form a subalgebra of A , written $\mathrm{A}_{\text {sep }}$ and called the separable closure of K in A .

### 5.4.6

Proposition Let L be an etale extension of K . An etale algebra A over L is etale over K.

Proof Let $\Omega$ be an algebraic closure of L . Then $\Omega \otimes_{\mathrm{K}} \mathrm{L} \approx \Omega^{d}$ and $\Omega \otimes_{\mathrm{L}} \mathrm{A} \approx \Omega^{e}$ with $d \in \mathbb{N}^{*}$ and $e \in \mathbb{N}$, and so

$$
\Omega \otimes_{\mathrm{K}} \mathrm{~A} \approx\left(\Omega \otimes_{\mathrm{K}} \mathrm{~L}\right) \otimes_{\mathrm{L}} \mathrm{~A} \approx \Omega^{d} \otimes_{\mathrm{L}} \mathrm{~A} \approx\left(\Omega \otimes_{\mathrm{L}} \mathrm{~A}\right)^{d} \approx \Omega^{e d}
$$

### 5.4.7

Let A be a finite cyclic extension of K , then there is an irreducible polynomial P such that $\mathrm{A}=\mathrm{K}[\mathrm{X}] /(\mathrm{P})$.

Proposition The following conditions are equivalent:
(i) A is not etale;
(ii) $\mathrm{P}^{\prime}=0$;
(iii) K has characteristic $p \neq 0$ and $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}^{p}\right]$.

Proof (ii) $\Leftrightarrow$ (i): The g.c.d. of $\left(\mathrm{P}, \mathrm{P}^{\prime}\right)$ is a divisor of P and so is 1 or P . If $\mathrm{P}^{\prime} \neq 0$, then $\operatorname{deg} \mathrm{P}^{\prime}<\operatorname{deg} \mathrm{P}$. Hence P does not divide $\mathrm{P}^{\prime}$ and g.c.d. $\left(\mathrm{P}, \mathrm{P}^{\prime}\right)=1$.

If $\mathrm{P}^{\prime}=0$, g.c.d. $\left(\mathrm{P}, \mathrm{P}^{\prime}\right)=\mathrm{P} \neq 1$ since $\mathrm{A} \neq 0$, and (i) $\Leftrightarrow$ (ii) thanks to (5.4.4).
(ii) $\Leftrightarrow$ (iii). If the characteristic of K is 0 then the derivative of a polynomial of degree $d>0$ is a nonzero polynomial of degree $d-1$.

If the characteristic of K is $p \neq 0$ then, for $\mathrm{P}=\sum a_{k} \mathrm{X}^{k} \mathrm{P}^{\prime}=\sum k a_{k} \mathrm{X}^{k-1}$. Hence $\mathrm{P}^{\prime}=0$ if and only if $a_{k}=0$ for all $k$ that is not a multiple of $p$, in other words if $\mathrm{P} \in \mathrm{K}\left[\mathrm{X}^{p}\right]$,

### 5.4.8

Theorem and Definition Let p be the characteristic of K. The following conditions are equivalent:
(i) Every finite extension of K is etale.
(ii) Either $p=0$, or $p \neq 0$ and the Frobenius endomorphism of K is surjective.

If these conditions hold, then K is said to be a perfect field.
Proof Condition (i) is equivalent to
(i') Every finite cyclic extension of K is etale.
Indeed (i) $\Rightarrow$ ( $\mathrm{i}^{\prime}$ ) is obvious, and ( $\left.\mathrm{i}^{\prime}\right) \Rightarrow$ (i) thanks to (5.4.5, (e) since every finite extension is generated by cyclic extensions.
(ii) $\Rightarrow\left(\mathrm{i}^{\prime}\right)$. Let $\mathrm{E}=\mathrm{K}[\mathrm{X}] /(\mathrm{P})$ be a finite cyclic extension. Since P is irreducible, if $K$ is of characteristic 0 then $E$ is etale (5.4.7, (iii) $\Rightarrow$ (i)).

Suppose that K is of characteristic $p \neq 0$. Then $\mathrm{P} \notin \mathrm{K}\left[\mathrm{X}^{p}\right]$. Indeed, if $\left.\mathrm{P} \in \mathrm{K}^{2} \mathrm{X}^{p}\right]$, then $\mathrm{P}=\sum_{k} a_{k} \mathrm{X}^{p k}$, where for all $k, a_{k}=b_{k}^{p}$ and $b_{k} \in \mathrm{~K}$, and so $\mathrm{P}=\mathrm{Q}^{p}$ for some $\mathrm{Q}=\sum b_{k} \mathrm{X}^{k}$. Hence P is reducible.
(i') $\Rightarrow$ (ii). Suppose that $p \neq 0$ and that F is not surjective. Let $a \in \mathrm{~K}, a \notin \mathrm{~F}(\mathrm{~K})$, and $\Omega$ be an algebraic closure of K and $\alpha \in \Omega$ such that $\alpha^{p}=a$. Then $\mathrm{X}^{p}-a=$ $(\mathrm{X}-\alpha)^{p}$. The minimal polynomial P of $\alpha$ is monic and divides $\mathrm{X}^{p}-a$, and so $\mathrm{P}=(\mathrm{X}-\alpha)^{d}$, where $1<d \leqslant p$, since $\alpha \notin \mathrm{K}$. Therefore $\alpha$ is a multiple root of P and the algebra $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$ is a non-etale finite cyclic extension (in fact, $d=p$ necessarily holds).

Example Every finite field is perfect. Indeed, as $\mathrm{F}: \mathrm{K} \rightarrow \mathrm{K}$ is injective, it is necessarily surjective if K is finite.

Counterexample If the characteristic of K is $p \neq 0$, then the field $\mathrm{K}(\mathrm{X})$ is not perfect. Indeed, there is no $f \in \mathbf{K}(\mathbf{X})$ such that $f^{p}=\mathbf{X}$.

### 5.4.9

Proposition Let A be a finite K-algebra.
(a) If K is perfect, then A is etale if and only if A is reduced.
(b) Let A be generated by $x_{1}, \ldots, x_{r}$ with respectively $\mathrm{P}_{1}, \ldots, \mathrm{P}_{r}$ as minimal polynomials. Then A is etale if and only if $(\forall i)$ g.c.d. $\left(\mathrm{P}_{i}, \mathrm{P}_{i}^{\prime}\right)=1$.
(c) With the notation of $(\mathrm{b})$, if K is of characteristic $p \neq 0$ and moreover A is an extension, then A is etale if and only if $(\forall i) \mathrm{P}_{i} \notin \mathrm{~K}\left[\mathrm{X}^{p}\right]$.

Proof (a) Let $L$ be an extension of $K$ diagonalizing A; then $A=K \otimes_{K} A$ embeds in $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}=\mathrm{L}^{d}$ for some $d$. The algebra $\mathrm{L}^{d}$ being a product of fields is reduced and hence so is A .

Conversely, A being finite and reduced is a finite product of finite extensions over a perfect field, and so is etale (5.4.5, c).
(b) and (c): A is etale if and only if $(\forall i) \mathrm{K}[\mathrm{X}] /\left(\mathrm{P}_{i}\right)$ is etale (5.4.5, a, e), in other words, if and only if $(\forall i)$ g.c.d. $\left(\mathrm{P}_{i}, \mathrm{P}_{i}^{\prime}\right)=1$ (5.4.4).

If A is an extension, all $\mathrm{P}_{i}$ are irreducible, $\mathrm{K}[\mathrm{X}] /\left(\mathrm{P}_{i}\right)$ are finite cyclic extensions, and by (5.4.7) A is etale if and only if $(\forall i) \mathrm{P}_{i} \notin \mathrm{~K}[\mathrm{X}]^{p}$.

Corollary The tensor product of two reduced finite algebras over a perfect field is a reduced finite algebra. In particular, the tensor product of two finite extensions is a reduced algebra.

Proof This corollary follows from (5.4.9, a) and (5.4.5, d).
Remark (1) In general, the tensor product of two extensions is not an extension. For example $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C} \approx \mathbb{C} \times \mathbb{C}$.
(2) The previous corollary has a converse. Over any non-perfect field $K$, there are finite extensions $A$ and $B$ such that $A \otimes_{K} B$ is not reduced (5.4, Exercise 3 ).

### 5.4.10 Example

Let K be a field of characteristic $p, \mathrm{G}$ a finite commutative group of order $n$ indivisible by $p$. Then $\mathrm{K}^{\mathrm{G}}$ equipped with convolution is an etale algebra (Maschke's theorem).

Indeed, for $x \in \mathrm{G}$, let $e_{x} \in \mathrm{~K}^{\mathrm{G}}$ be defined by $e_{x}(y)=0$ if $x \neq y$ and $e_{x}(x)=1$. If $x$ has order $k$, then $e_{x}$ is a root of $\mathrm{P}(\mathrm{X})=\mathrm{X}^{k}-1$. As $p$ does not divide $n$, neither does it divide $k$; $k$ is invertible in K and the ideal of $\mathrm{K}[\mathrm{X}]$ generated by $\left(\mathrm{X}^{k}-1\right)$ and $k X^{k-1}$ is the whole ring, and so $K[X] /\left(X^{k}-1\right)$ is etale. Since the elements $e_{x}$ generate $\mathrm{K}^{\mathrm{G}}$, the algebra $\left(\mathrm{K}^{\mathrm{G}}, *\right)$ must be etale by $(5.4 .9, \mathrm{~b})$.

### 5.4.11 Primitive Element Theorem

Theorem If K is infinite, every etale algebra over K is cyclic.
Lemma 5.1 An etale algebra has only finitely many subalgebras.
Let $A$ be an etale algebra and $L$ an extension of $K$ such that $L \otimes_{K} A$ is diagonal. For any subalgebra $B$ of $A, L \otimes_{K} B$ is a subalgebra of $L \otimes_{K} A$ and $B=A \cap\left(L \otimes_{K} B\right)$. So $B \mapsto L \otimes_{K} B$ is an injection from the set of subalgebras of $A$ to the set of subalgebras of $L \otimes_{K} A$, which is finite.

Lemma 5.2 Let K be an infinite field, E a finite dimensional vector space over K and let $\mathrm{F}_{i}$ be finitely many strict subspaces. Then $\bigcup \mathrm{F}_{i} \neq \mathrm{E}$.

Indeed, each $\mathrm{F}_{i}$ can be embedded in a hyperplane $\mathrm{H}_{i}$ defined by an equation $h_{i}=0$. We may assume that $\mathrm{E}=\mathrm{K}^{n}$. Every $h_{i} \in \mathrm{~K}\left[x_{1}, \ldots, x_{n}\right]$ is a homogeneous polynomial of degree 1 . So $h=\prod h_{i}$ is $\neq 0$. As K is infinite, the polynomial function $\mathrm{K}^{n} \rightarrow \mathrm{~K}$ represented by $h$ is nonzero. There exists $x \in \mathrm{~K}^{n}$ such that $h(x) \neq 0$ and hence $x \notin \prod \mathrm{H}_{i}$.

Proof of the Theorem Let A be an etale K-algebra. It is finite and has finitely many strict subalgebras $\mathrm{A}_{i}$ (Lemma 1). So $\bigcup \mathrm{A}_{i} \neq \mathrm{A}$. Let $x \in \mathrm{~A}-\bigcup \mathrm{A}_{i}$. The subalgebra generated by $x$ is not strict, hence must be A.

Definition A generator of A is called a primitive element of A .
Corollary (Primitive element Theorem) All etale extensions L of K are cyclic.
If K is infinite, this amounts to the theorem.
If $K$ is finite, then $L^{*}$ is cyclic and any of its generators generates $L$.

## Exercises 5.4. (Etale algebras)

1.-Let $E$ be an algebraic extension of $K$ and $\Omega$ an algebraic closure of $K$. Show that the following conditions are equivalent:
(i) E is separable;
(ii) All finite sub-extensions of E are etale;
(iii) All $\Omega$-algebra homomorphisms from $\Omega \otimes_{\mathrm{K}} \mathrm{E}$ to $\Omega$ separate the elements of $\Omega \otimes_{\mathrm{K}} \mathrm{E}$;
(iv) For all finite dimensional vector subspaces V of $\mathrm{E}, \operatorname{Hom}_{\mathrm{K}}(\mathrm{V}, \Omega)$ is generated as a vector space over $\Omega$ by $\left.f\right|_{\mathrm{V}}, f \in \operatorname{Hom}_{\mathrm{K}-\Delta \ell_{g}}(\mathrm{E}, \Omega)$.
2.-(a) Let A and B be K -algebras and $f: \mathrm{A} \rightarrow \mathrm{B}$ a surjective homomorphism such that $\mathrm{I}=\operatorname{Ker} f$ is a nilpotent ideal. Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ be a separable polynomial and $y \in \mathrm{~B}$ such that $\mathrm{P}(y)=0$. Show that there is a unique $x \in \mathrm{~A}$ such that $f(x)=y$ and $\mathrm{P}(x)=0$ ( $x$ may be constructed through successive approximations by noting that for all $x^{\prime} \in f^{-1}(y), \mathrm{P}^{\prime}\left(x^{\prime}\right)$ is invertible in A).

Deduce that $f$ induces a bijection from $\mathrm{A}_{\text {sep }}$ onto $\mathrm{B}_{\text {sep }}$.
(b) Let A and B be finite K -algebras and $f: \mathrm{A} \rightarrow \mathrm{B}$ a surjective homomorphism. Show that $f$ induces a surjective homomorphism

$$
f_{*}: \mathrm{A}_{\text {sep }} \rightarrow \mathrm{B}_{\text {sep }}
$$

(see [1], Exercise 7).
(c) Show that if $K$ is perfect then, for all K-algebras $A, A=A_{\text {sep }} \oplus \mathfrak{N}$, where $\mathfrak{N}$ is the nilradical of A .
3.-Let K be a non-perfect field of characteristic $p$, and $a \in \mathrm{~K}$ an element without any $p$-th root in K . Set $\mathrm{L}=\mathrm{K}[\mathrm{X}] /\left(\mathrm{X}^{p}-a\right)$. Show that L is an extension of K , but that the algebra $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{L}$ is not reduced.
4.-Let $k$ be an infinite field of characteristic $p$; set $\mathrm{K}=k(\mathrm{X}, \mathrm{Y})$ and $\mathrm{L}=k\left(\mathrm{X}^{1 / p}\right.$, $\left.\mathrm{Y}^{1 / p}\right)\left(\right.$ or, if preferred, $\mathrm{L}=k\left(\mathrm{X}^{\prime}, \mathrm{Y}^{\prime}\right)$ and $\left.\mathrm{K}=k\left(\mathrm{X}^{\prime p}, \mathrm{Y}^{\prime p}\right)\right)$.
(a) Show that L is a non cyclic extension of K of degree $p^{2}$.
(b) For all $a \in k$, let $\mathrm{M}_{a}$ be the subfield of L generated by K and $(\mathrm{Y}-a \mathrm{X})^{1 / p}$. Show that these subfields are mutually distinct.
5.-Let K be a field of characteristic $p \neq 0$ and $\mathrm{G}=\mathbb{Z} /\left(p^{r}\right)$. Show that the algebra of the group G is isomorphic to $\mathrm{K}[\mathrm{T}] /\left(\mathrm{T}^{p^{r}}\right)$.
6.-Let $k$ be a field, and $\mathrm{A}=k[\mathrm{X}], \mathrm{B}=k[\mathrm{X}, \mathrm{Y}] /(\mathrm{P})$, where $\mathrm{P} \in k[\mathrm{X}, \mathrm{Y}]$ is a polynomial of degree $d$ in Y ; set $\mathrm{K}=k(\mathrm{X})$ and $\mathrm{E}=\mathrm{K} \otimes_{\mathrm{A}} \mathrm{B}$.

Show that E is etale over K if and only if there exists $x$ in an algebraic closure $\Omega$ of $k$ such that $\mathrm{P}(x, \mathrm{Y}) \in \Omega[\mathrm{Y}]$ has degree $\geqslant d-1$ and distinct roots in $\Omega$.

### 5.5 Purely Inseparable Extensions

In this section, K is assumed to be of characteristic $p \neq 0$.

### 5.5.1

Proposition Let $\Omega$ be an algebraic closure of K and $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ an irreducible monic polynomial. The following conditions are equivalent:
(i) the polynomial P has only one (possibly multiple) root in $\Omega$;
(ii) there exists $r \in \mathbb{N}$ and $a \in \mathrm{~K}$ such that $\mathrm{P}=\mathrm{X}^{p^{r}}-a$.

Proof (ii) $\Rightarrow$ (i). The Frobenius endomorphism F : $x \mapsto x^{p}$ of $\Omega$ is injective since $\Omega$ is a field. Hence the same holds for $\mathrm{F}^{r}$.
(i) $\Rightarrow$ (ii). Let $r$ be the largest integer for which $\mathrm{P}(\mathrm{X})$ can be written as $\mathrm{Q}\left(\mathrm{X}^{p^{r}}\right)$. Then $\mathrm{Q}(\mathrm{Y}) \notin \mathrm{K}\left[\mathrm{Y}^{p}\right]$, for otherwise $\mathrm{P}(\mathrm{X})$ could be written as $\mathrm{Q}_{1}\left(\mathrm{X}^{r^{r+1}}\right)$. Moreover Q is irreducible since so is P . By (5.4.7) and (5.4.4), the roots of Q in $\Omega$ are distinct; as P has only one root, so does Q and thus Q is of degree 1 .

### 5.5.2

Definition Let L be an extension of K . An element $x \in \mathrm{~L}$ is said to be radical over K if there is an integer $r$ such that $x^{p^{r}} \in \mathrm{~K}$. If all elements of L are radical over K , L is a purely inseparable extension (or radical) of K .

### 5.5.3

Proposition Let L be an algebraic extension of K and $\Omega$ an algebraic closure of K . The following conditions are equivalent:
(i) L is purely inseparable;
(ii) there is no K -algebra homomorphism from L to $\Omega$.

Proof (i) $\Rightarrow$ (ii). Let $x \in \mathrm{~L}$. There exist $r \in \mathbb{N}$ and $a \in \mathrm{~K}$ such that $x^{p^{r}}=a$. If $f: \mathrm{L} \rightarrow \Omega$ is a homomorphism, then $f(x)$ is necessarily the unique root of $\mathrm{X}^{p^{r}}-a$ in $\Omega$.
not (i) $\Rightarrow$ not (ii). Suppose $x \in \mathrm{~L}$ is not radical. The minimal polynomial of $x$ has at least two distinct roots in $\Omega$, and thus there are two distinct homomorphisms from $\mathrm{K}[x]_{\mathrm{L}}$ to $\Omega$. Embedding in $\Omega$, which is possible by (5.2.6), gives two distinct homomorphisms from L to $\Omega$.

### 5.5.4

Proposition Every algebraic extension L of K is a purely inseparable extension of $\mathrm{L}_{\text {sep }}$.

Proof Let $x \in \mathrm{~L}, \mathrm{P}$ its minimal polynomial, $r$ the largest integer for which $\mathrm{P}(\mathrm{X})$ can be written as $\mathrm{Q}\left(\mathrm{X}^{p^{r}}\right)$. Then $\mathrm{Q} \notin \mathrm{K}\left[\mathrm{X}^{p}\right]$ and Q is irreducible, so that Q is separable, and $x^{p^{r}} \in \mathrm{~L}_{\text {sep }}$ since its minimal polynomial is separable.

Remark Let L be an extension of K . Then the set $\mathrm{L}_{\mathrm{rad}}$ of the radicals over K in L is a sub-extension of $L$. Indeed $L_{\text {rad }}$ is the increasing union of $\left(\mathrm{F}^{r}\right)^{-1}(\mathrm{~K})$, and so is a subfield, and $L_{\text {rad }} \supset K$. In general, even when the extension $L$ is algebraic, it is not separable over $L_{r a d}$. In other words, an algebraic extension of $K$ is always a purely inseparable extension of a separable extension of $K$, but not necessarily a separable extension of a purely inseparable extension of K (5.5, Exercise 3).

However, if $\Omega$ is an algebraic closure of K , then $\Omega$ is a separable extension of $\Omega_{\mathrm{rad}}$ since $\Omega_{\mathrm{rad}}$ is a perfect field.

### 5.5.5

Theorem (MacLane's criterion) Let A be an algebraic K -algebra over K , and $\Omega$ an algebraic closure of K .
(a) The following conditions are equivalent:
(i) A is separable;
(ii) $\Omega_{r a d} \otimes_{\mathrm{K}} \mathrm{A}$ is reduced;
(iii) $\Omega \otimes_{\mathrm{K}} \mathrm{A}$ is reduced.
(b) A is a separable extension of K if and only if $\Omega_{r a d} \otimes_{\mathrm{K}} \mathrm{A}$ is an extension of K .

Proof All the properties considered continue to hold for the direct limit. Hence A may be assumed to be finite over K.
(a) (i) $\Rightarrow \mathrm{A}$ is etale $\Rightarrow \Omega \otimes_{\mathrm{K}} \mathrm{A}$ is diagonal $\Rightarrow$ (iii) $\Rightarrow$ (ii). We show that (ii) $\Rightarrow$ (i): if (ii) holds, then A is reduced and hence is a product of extensions. So it suffices
to prove that these extensions are separable. We do so by showing that if L is a not a separable extension, then $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{L}$ is not reduced: let $x \in \mathrm{~L}-\mathrm{L}_{\text {sep }}$ and P its minimal polynomial. Then $\mathrm{P}(\mathrm{X})=\mathrm{Q}\left(\mathrm{X}^{p}\right)$ for some polynomial Q ; let $\mathrm{Q}_{1} \in \Omega_{\mathrm{rad}}[\mathrm{X}]$ be the polynomial all of whose coefficients are the roots of the $p$-th coefficients of Q . Then $\mathrm{Q}_{1}^{p}=\mathrm{P}$, and so $\left(\mathrm{Q}_{1}(x)\right)^{p}=0$. But $\mathrm{Q}_{1}(x) \neq 0$. Indeed, $\mathrm{K}[x]_{\mathrm{L}}=\mathrm{K}[\mathrm{X}] /(\mathrm{P})$; hence the subalgebra $\Omega_{\mathrm{rad}}[x]$ of $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{L}$ can be identified with $\Omega_{\mathrm{rad}}[\mathrm{X}] /(\mathrm{P})$, and the minimal polynomial of $x$ in $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{L}$ is P . However, $\mathrm{Q}_{1}$ is not a multiple of P since its degree is strictly smaller. Therefore $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{L}$ is not reduced.
(b) If $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{A}$ is an extension of K , then the algebra A is an extension of K ; it is separable by (a). Assume that A is a separable extension. Then $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{A}$ is a finite reduced algebra over $\Omega_{\mathrm{rad}}$ and can be written $\prod_{i \in \mathrm{I}} \mathrm{E}_{i}$, where all $\mathrm{E}_{i}$ are extensions of $\Omega_{\mathrm{rad}}$ (5.1.6, Corollary 1). These $\mathrm{E}_{i}$ can also be considered extensions of A. Embed A in $\Omega$. By (5.5.3), the set $\operatorname{Hom}_{\mathrm{A}-\mathscr{A l g}}\left(\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{A}, \Omega\right)=\operatorname{Hom}_{\mathrm{K}-\mathscr{A l g}}\left(\Omega_{\mathrm{rad}}, \Omega\right)$ consists of a single element. If there were distinct elements $i$ and $j$ in I , projecting respectively the product onto $\mathrm{E}_{i}$ and $\mathrm{E}_{j}$ and embedding $\mathrm{E}_{i}$ and $\mathrm{E}_{j}$ in $\Omega$ would give two distinct A-homomorphisms from $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{A}$ to $\Omega$. Hence, I has a unique element, and $\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{A}$ is an extension of K .

### 5.5.6 Characteristic 0 Conventions

If K is a field of characteristic 0 , the identity map of K is called the Frobenius endomorphism. Un element of an extension of K belonging to K is said to be radical. Only the trivial extension is purely inseparable. With these conventions, the results of 5.5.3, 5.5.4 and 5.5.5 hold in characteristic 0 .

## Exercises 5.5. (Purely inseparable extensions)

1.-(a) Let $\Omega$ and $\Omega^{\prime}$ be two algebraic closures of $K$. Show that there is a unique isomorphism from $\Omega_{\mathrm{rad}}$ to $\Omega_{\mathrm{rad}}^{\prime}$. (Thus $\Omega_{\mathrm{rad}}$ and $\Omega_{\mathrm{rad}}^{\prime}$ can be identified, and the purely inseparable closure of K is well-defined.)
(b) Show that every purely inseparable and perfect extension of K can be uniquely identified with $\Omega_{\mathrm{rad}}$.
(c) Consider the direct system

$$
\mathrm{K} \xrightarrow{\mathrm{~F}} \mathrm{~K} \xrightarrow{\mathrm{~F}} \mathrm{~K} \xrightarrow{\mathrm{~F}} \cdots \xrightarrow{\mathrm{~F}} \mathrm{~K} \xrightarrow{\mathrm{~F}} \cdots
$$

indexed by $\mathbb{N}$, where $F$ is the Frobenius endomorphism. Let $E$ be the direct limit of this system. Show that, equipped with $f_{0}^{\infty}: \mathrm{K} \rightarrow \mathrm{E}, \mathrm{E}$ is an extension of K isomorphic to $\Omega_{\mathrm{rad}}$.
(d) State a functorial property of the purely inseparable closure.
2.-(a) Show that if $K$ is perfect, then every finite extension of $K$ is a perfect field $\left(\operatorname{deg}_{\mathrm{K}}(\mathrm{F}(\mathrm{L}))=\operatorname{deg}_{\mathrm{F}(\mathrm{K})}(\mathrm{F}(\mathrm{L}))=\operatorname{deg}_{\mathrm{K}}(\mathrm{L})\right)$. Deduce that all algebraic extensions of K are perfect fields.
(b) Let K be a field and $\Omega$ an algebraic closure of K . Show that $\Omega=\Omega_{\mathrm{rad}} \otimes_{\mathrm{K}} \Omega_{\text {sep }}$ (the term on the right can be identified with a perfect sub-extension of $\Omega$ containing $\Omega_{\text {sep }}$ ).
3.-Let $k$ be a field of characteristic $p \neq 0,2$. Consider the subfields $\mathrm{M}=k\left(\mathrm{X}, \mathrm{Y}^{p}\right)$, $\mathrm{L}=k\left(\mathrm{X}^{p}, \mathrm{Y}^{p}\right)$ and

$$
\mathrm{K}=\{f \in \mathrm{~L} \mid f(\mathrm{Y}, \mathrm{X})=f(\mathrm{X}, \mathrm{Y})\}
$$

in $\mathrm{E}=k(\mathrm{X}, \mathrm{Y})$. (a) Show that L is a separable extension of K and that M is a purely inseparable extension of $L$. Determine $\mathrm{M}_{\text {sep }}$ relative to K .
(b) Let $f \in \mathrm{M}$ be such that $f(\mathrm{Y}, \mathrm{X})=f(\mathrm{X}, \mathrm{Y})$. Show that $f \in \mathrm{~K}$ (reduce to the case $f \in k[\mathrm{X}, \mathrm{Y}])$.
(c) Show that the purely inseparable closure of $K$ in $M$ is $K$. Deduce that $M$ is not a separable extension of a purely inseparable extension of $K$.
4.-Let $L$ an algebraic extension of the field $K$. Consider the canonical homomorphism $\varphi: \mathrm{L}_{\text {rad }} \otimes_{\mathrm{K}} \mathrm{L}_{\text {sep }} \rightarrow \mathrm{L}$ inducing the canonical injection on $\mathrm{L}_{\text {rad }}$ and $\mathrm{L}_{\text {sep }}$.
(a) Show that $\varphi$ is injective (see 5.5.5).
(b) Show that the image of $\varphi$ is the separable closure M of $\mathrm{L}_{\mathrm{rad}}$ in L (reduce to the case of finite L and prove the equalities

$$
\begin{aligned}
\left(\mathrm{M}: \mathrm{L}_{\mathrm{rad}}\right) & =\operatorname{card} \operatorname{Hom}_{\mathrm{L}_{\text {rad }}-\mathscr{A} \ell_{g}(\mathrm{M}, \Omega)}=\operatorname{card} \operatorname{Hom}_{\mathrm{K}-\mathscr{A} \ell_{g}(\mathrm{M}, \Omega)} \\
& \left.\left.=\operatorname{card} \operatorname{Hom}_{\mathrm{K}-\mathscr{A} \ell_{g}} \mathrm{~L}_{\text {sep }}, \Omega\right)=\left(\mathrm{L}_{\text {sep }}: \mathrm{K}\right)=\left(\mathrm{L}_{\mathrm{rad}} \otimes_{\mathrm{K}} \mathrm{~L}_{\text {sep }}: \mathrm{L}_{\mathrm{rad}}\right)\right) .
\end{aligned}
$$

In particular $L=L_{\text {rad }} \otimes_{K} L_{\text {sep }}$ if and only if $L$ is a separable extension of $L_{\text {rad }}$.
(c) Let $\Omega$ be an algebraic closure of $\mathrm{K}, x \in \Omega, \mathrm{P}$ its minimal polynomial, and L the sub-extension of $\Omega$ generated by the roots of P . Show that $\mathrm{L}=\mathrm{L}_{\text {rad }} \otimes_{\mathrm{K}} \mathrm{L}_{\text {sep }}$ (consider the symmetric functions of the roots of P , each counted only once).
5.-Let $A$ be a $K$-algebra and $B$ a subalgebra of $A$. We will say that $A$ is radical over B if $\Omega \otimes_{\mathrm{K}} \mathrm{A}$ is generated by $\Omega \otimes_{\mathrm{K}} \mathrm{B}$ and by nilpotent elements. ${ }^{1}$
(a) Show that for an extension $A$ of $K$, the definition is equivalent to that of (5.5.2).
(b) Show that every algebraic algebra A over K is radical over $\mathrm{A}_{\text {sep }}$ (use 5.4, Exercise 2).
(c) Generalize to the algebras of Proposition 5.5.3.

[^18]6.-Let A be a finite K-algebra.
(a) Let $\Omega$ be an algebraic closure of K. Show that all K-algebra homomorphisms from $\mathrm{A}_{\text {sep }}$ to $\Omega$ extend uniquely to A . Deduce that the cardinality of $\operatorname{Hom}_{\mathrm{K}-\operatorname{seg}_{g}(\mathrm{~A}, \Omega)}$ is equal to $\operatorname{deg}_{\mathrm{K}}\left(\mathrm{A}_{\text {sep }}\right)$.
(b) Let L be an extension of K . Show that the separable closure of L in $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ is $L \otimes_{K} A_{\text {sep }}$.
(c) If A and B are two finite K -algebras, show that
$$
\left(A \otimes_{K} B\right)_{\text {sep }}=A_{\text {sep }} \otimes_{K} B_{\text {sep }}
$$

Extend this result to algebraic K-algebras.
7.-Let $\Omega$ be an algebraic closure of K and $\Omega^{\prime} \subset \Omega$ a sub-extension. Assume that all $P \in K[X]$ have at least one root in $\Omega^{\prime}$.
(a) Show that every finite separable extension K embeds in $\Omega^{\prime}$ (use the primitive element theorem). Show that $\Omega_{\text {sep }}^{\prime}=\Omega_{\text {sep }}$.
(b) Prove that, for all $n$, every separable polynomial in $\mathrm{K}[\mathrm{X}]$ has at least one root in $\Omega_{n}^{\prime}=\mathrm{F}^{n}\left(\Omega^{\prime}\right) \cap \Omega_{\text {sep }}$. Show that $\Omega^{\prime}=\Omega$.
(c) Show that in 5.2, Exercise 4, the family $\left(\mathrm{A}_{i}\right)$ of extensions K can be taken to be of the form $\mathrm{K}[\mathrm{X}] /(\mathrm{P})$, where P is an irreducible polynomial.

### 5.6 Finite Galois Extensions

### 5.6.1

Definition A finite extension L of K is said to be Galois if L diagonalizes itself (5.4.1).

If $L$ is a Galois extension of $K$, the group $G=A u t_{K} L$ is called the Galois group of $L$ over $K$.

### 5.6.2

Proposition Let L be a Galois extension of K and A a sub-extension of L . Then L diagonalizes A over K , and as an A-algebra, L is a Galois extension of A .

Proof $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ is a subalgebra of $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{L}$, and $\mathrm{L} \otimes_{\mathrm{A}} \mathrm{L}$ is a quotient L-algebra of $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{L}$ since $(\lambda, \mu) \mapsto \lambda \otimes_{\mathrm{A}} \mu$ is A-bilinear and so necessarily K-bilinear, and thus defines a map $\epsilon: \lambda \otimes_{\mathrm{K}} \mu \mapsto \lambda \otimes_{\mathrm{A}} \mu$ whose image contains the simple tensors; hence
$\epsilon$ is surjective. Besides, $\epsilon$ is a L -algebra homomorphism, and so $\mathrm{L} \otimes_{K} A$ and $L \otimes_{K} L$ are diagonal L -algebras.

Remark A Galois extension of a Galois extension is not necessarily Galois (5.6, Exercise 3).

### 5.6.3

Let $L$ be a finite extension of $K$ and $\Omega$ an algebraic closure of $K$ containing L. Set

$$
d=\operatorname{deg}_{\mathrm{K}} \mathrm{~L}, \mathrm{G}=\operatorname{Aut}_{\mathrm{K}} \mathrm{~L}, \operatorname{Fix}_{\mathrm{G}} \mathrm{~L}=\{x \in \mathrm{~L} \mid(\forall g \in \mathrm{G}) g(x)=x\}
$$

Theorem The following conditions are equivalent:
(i) L is Galois;
(ii) $\operatorname{card} \mathrm{G}=d$;
(ii') $\operatorname{card} \mathrm{G} \geqslant d$;
(iii) L is etale and every automorphism of $\Omega$ preserves L ;
(iv) $\mathrm{Fix}_{\mathrm{G}} \mathrm{L}=\mathrm{K}$;
(v) for all $x \in \mathrm{~L}$, the roots of the minimal polynomial of $x$ in $\Omega$ are simple and belong to L .

Proof By 5.3.3, (i) $\Leftrightarrow$ (ii) $\Leftrightarrow$ (ii') since

$$
\mathrm{G}=\operatorname{Aut}_{\mathrm{K}} \mathrm{~L}=\operatorname{End}_{\mathrm{K}} \mathrm{~L}=\operatorname{Hom}_{\mathrm{L}}\left(\mathrm{~L} \otimes_{\mathrm{K}} \mathrm{~L}, \mathrm{~L}\right)
$$

(i) $\Rightarrow$ (iii). As L is etale, $\operatorname{card} \operatorname{Hom}_{\mathrm{K}-\mathscr{A t g}}(\mathrm{L}, \Omega)=d$. As $\mathrm{G} \subset \operatorname{Hom}_{\mathrm{K}-\mathscr{A t g}}(\mathrm{L}, \Omega)$ and (i) $\Rightarrow$ (ii), G $=\operatorname{Hom}_{\mathrm{K}-\mathscr{A l g}}(\mathrm{L}, \Omega$ ), proving (iii).
(iii) $\Rightarrow$ (ii). As L is etale, there are $d$ distinct homomorphisms from L to $\Omega$. By (5.2.6) and (5.2.2), these homomorphisms are induced by automorphisms of $\Omega$ and so are automorphisms of L, proving (ii).
(iii) $\Rightarrow$ (iv). All points of $K$ are fixed under the action of G. Assume $K \neq$ Fix $_{G} L$; let $x \in \mathrm{Fix}_{\mathrm{G}} \mathrm{L}-\mathrm{K}$ and P its minimal polynomial. As $x \notin \mathrm{~K}, \mathrm{~d}^{\circ} \mathrm{P} \geqslant 2$; and as L is etale, $x$ is a simple root of P . Let $y$ be another root of P in $\Omega$. The homomorphism $\varphi: \mathrm{K}[x] \rightarrow \Omega$ defined by $\varphi(x)=y$ extends to an automorphism of $\Omega$ which induces an automorphism of L not fixing $x$, giving a contradiction.
(iv) $\Rightarrow(\mathrm{v})$. Let $\alpha \in \mathrm{L}$ and $\left(a_{i}\right)_{i \in \mathrm{I}}$ such that $i \mapsto a_{i}$ is a bijection from I onto the set of conjugates of $\alpha$ in L . Set $\mathrm{P}(\mathrm{X})=\prod\left(\mathrm{X}-a_{i}\right)$. For $g \in \mathrm{G}$, define

$$
g_{*}: \mathrm{L}[\mathrm{X}] \rightarrow \mathrm{L}[\mathrm{X}] \text { by } g_{*}(\lambda)=g(\lambda) \text { for } \lambda \in \mathrm{L} \text { and } g_{*}(\mathrm{X})=\mathrm{X}
$$

Then

$$
g_{*}(\mathrm{P})=\prod\left(\mathrm{X}-g\left(a_{i}\right)\right)=\prod\left(\mathrm{X}-a_{i}\right)
$$

since $g$ induces a permutation of these $a_{i}$. Hence the coefficients of P are invariant under $G$, and $P \in K[X]$. The roots of $P$ are simple and in L. Hence the minimal polynomial of $\alpha$, which divides P , has the same property (in fact this polynomial may be shown to be P).
(v) $\Rightarrow$ (iii). Clearly, L is etale. For all automorphisms $\varphi$ of $\Omega$ and all $x \in \mathrm{~L}, \varphi(x)$ is a root of the minimal polynomial of $x$ and so $\varphi(x) \in \mathrm{L}$ by assumption. Hence $\varphi$ preserves L.

### 5.6.4

Corollary 5.6 Let L and $\mathrm{L}^{\prime}$ be Galois extensions of K contained in $\Omega$. Then the product extension L • L' is Galois.

This follows from (i) $\Leftrightarrow$ (iii) in Theorem 5.6.3.
Corollary 5.7 Let L and $\mathrm{L}^{\prime}$ be Galois extensions of K with $\mathrm{K} \subset \mathrm{L}^{\prime} \subset \mathrm{L} \subset \Omega$. Set $\mathrm{G}=\mathrm{Aut}_{\mathrm{K}} \mathrm{L}, \mathrm{G}^{\prime}=\mathrm{Aut}_{\mathrm{K}} \mathrm{L}^{\prime}$, and $\mathrm{H}=\mathrm{Aut}_{\mathrm{L}^{\prime}} \mathrm{L}$. Then H is a normal subgroup of G and $\mathrm{G}^{\prime}$ identified with $\mathrm{G} / \mathrm{H}$.

Proof Every element of G leaves $\mathrm{L}^{\prime}$ invariant, and so there is a restriction homomorphism $\rho: \mathrm{G} \rightarrow \mathrm{G}^{\prime}$. Every element of $\mathrm{G}^{\prime}$ extends to an automorphism of $\Omega$, leaving L invariant. Hence $\rho$ is surjective. The kernel of $\rho$ is obviously H.

### 5.6.5 Decomposition Field

Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ be a separable polynomial (5.4.4). The K -sub-extension E of $\Omega$ generated by the roots of P is Galois (Theorem 5.6.3, (iii) $\Rightarrow$ (i)). It is called the decomposition field of P .

Proposition With this notation, the group $\mathrm{G}=\mathrm{Aut}_{\mathrm{K}} \mathrm{E}$ acts transitively on the set X of roots of P in $\Omega$ if and only if P is irreducible.

Proof Suppose that $P$ is not irreducible. Then $P=P_{1} \cdot P_{2}$ for some $P_{1}$ and $P_{2}$ of degree $>0$. Then $\mathrm{X}=\mathrm{X}_{1} \cup \mathrm{X}_{2}$, where $\mathrm{X}_{i}=\mathrm{P}_{i}^{-1}(0)$. These two sets are disjoint since the roots of P are simple, and they are invariant under G . Hence G does not act transitively on X.

Conversely, let $\alpha$ and $\beta$ be two roots of P . If P is irreducible, it is the minimal polynomial of $\alpha$ as well as of $\beta$. Hence there is an automorphism $\sigma$ from $\mathrm{K}[\alpha]_{\Omega}$ onto $\mathrm{K}[\beta]_{\Omega}$ such that $\sigma(\alpha)=\beta$. Then $\sigma$ extends to an automorphism of $\Omega$ inducing an automorphism of E .

### 5.6.6

Proposition Let A be an etale K-algebra. There is a finite Galois extension diagonalizing A.

Proof Let $\Omega$ be an algebraic closure of K and $x_{1}, \ldots, x_{r}$ generators of A. For each $i$, let $\mathrm{P}_{i}$ be the minimal polynomial of $x_{i}$ and $\mathrm{L}_{i}$ the sub-extension of $\Omega$ generated by the roots of $\mathrm{P}_{i}$. For each $i, \mathrm{~L}_{i}$ diagonalizes the subalgebra $\mathrm{A}_{i}$ of A generated by $x_{i}$, and the sub-extension L of $\Omega$ generated by all $\mathrm{L}_{i}$ diagonalizes all $\mathrm{A}_{i}$, and hence diagonalizes A . The roots of $\mathrm{P}_{i}$ are simple since A is etale, and so L is etale; every automorphism of $\Omega$ permutes the roots of $\mathrm{P}_{i}$ for each $i$, and so preserves L , and L is Galois.

Corollary Let $\Omega$ be an algebraic closure of K and L an etale extension of K contained in $\Omega$. Then there is a finite Galois extension of K in $\Omega$ and containing L .

Proof Let $\mathrm{E} \subset \Omega$ be a finite Galois extension diagonalizing L. Then $\mathrm{L} \subset \mathrm{E}$ since the image of any homomorphism, in particular the canonical injection, from L to $\Omega$ is in E .

### 5.6.7

Proposition Let L be a field, H a finite group of automorphisms of L , and set $\mathrm{F}=$ $\mathrm{Fix}_{\mathrm{H}} \mathrm{L}$. Then L is a finite Galois extension of F and $\mathrm{Aut}_{\mathrm{F}}(\mathrm{L})=\mathrm{H}$.

Proof (a) Assume that L is finite over F . Set $\mathrm{H}^{\prime}=$ Aut $_{\mathrm{F}} \mathrm{L}$. Then $\mathrm{H}^{\prime} \supset \mathrm{H}$, and $\mathrm{F} \subset$ $\mathrm{Fix}_{H^{\prime}} \mathrm{L} \subset \mathrm{Fix}_{\mathrm{H}} \mathrm{L}=\mathrm{F}$; so $\mathrm{F}=\mathrm{Fix}_{\mathrm{H}^{\prime}} \mathrm{L}$, and by (5.6.3, (iv) $\Rightarrow$ (i)), L is Galois over F . In particular, L is etale, and hence cyclic by the primitive element Theorem (5.4.11, Corollary). Let $a \in \mathrm{~L}$ such that $\mathrm{L}=\mathrm{F}[a]$. The polynomial

$$
\mathrm{P}=\prod_{g \in \mathrm{H}}(\mathrm{X}-g(a))
$$

is in $\mathrm{F}[\mathrm{X}]$ since its coefficients are fixed by H , and the minimal polynomial Q of $a$ over F divides P (in fact $\mathrm{Q}=\mathrm{P}$ ). For all $f \in \mathrm{H}^{\prime}, f(a)$ is another root of Q , hence of P , and there exists $g \in \mathrm{H}$ such that $f(a)=g(a)$. The automorphisms $f$ of $g$ agree on F and at $a$, and so $f=g$ and $\mathrm{H}^{\prime}=\mathrm{H}$.
(b) We now prove that L is finite over F . First, L is algebraic over F . Indeed, for all $a \in \mathrm{~L}$, the polynomial $\mathrm{P}=\prod_{g \in \mathrm{H}}(\mathrm{X}-g(a))$ is in $\mathrm{F}[\mathrm{X}]$ and $\mathrm{P}(a)=0$. The extension L is the directed union of finite sub-extensions $\mathrm{L}_{i}$ over F , each preserved by H . For each $i$, let $\mathrm{H}_{i}$ be the image of H in $\operatorname{Aut}_{\mathrm{F}} \mathrm{L}_{i}$. Then $\mathrm{Fix}_{\mathrm{H}_{i}} \mathrm{~L}_{i}=\mathrm{Fix}_{\mathrm{H}} \mathrm{L}_{i}=\mathrm{F}$, and so $\mathrm{Aut}_{\mathrm{F}} \mathrm{L}_{i}=\mathrm{H}_{i}$ by (a). As $\mathrm{L}_{i}$ is Galois, $\operatorname{deg}_{\mathrm{F}} \mathrm{L}_{i}=\operatorname{card} \mathrm{H}_{i} \leqslant \operatorname{card} \mathrm{H}$; hence finally $\operatorname{deg}_{\mathrm{F}} \mathrm{L} \leqslant \operatorname{card} \mathrm{H}<\infty$.

### 5.7 Finite Galois Theory

### 5.7.1

Let L be a finite Galois extension of K , A a finite K -algebra diagonalized by L , $\mathrm{G}=\mathrm{Aut}_{\mathrm{K}} \mathrm{L}$, and $\mathrm{X}=\operatorname{Hom}_{\mathrm{K}-\Delta \ell_{g}}(\mathrm{~A}, \mathrm{~L})$. The set X can be identified with Hom $\mathrm{H}_{\mathrm{L}}-\Delta \ell_{g}$ ( $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}, \mathrm{L}$ ). Set $g \cdot \xi=g \circ \xi$ for all $g \in \mathrm{G}$ and $\xi \in \mathrm{X}$. Thus G acts on X.

The Gelfand transform of the L-algebra $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ is the map $\gamma: t \mapsto \hat{t}$ from $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$


For all $t$ of the form $\lambda \otimes a,(\widehat{\lambda \otimes a)}(\xi)=\xi(\lambda \otimes a)=\lambda \xi(a)$. As G acts on L, G can be made to act on $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}$ by $g_{*}(\lambda \otimes a)=g(\lambda) \otimes a$. The Gelfand transform is an isomorphism since $L \otimes_{K} A$ is diagonal, and the action of $G$ on $L \otimes_{K} A$ can be transferred on $\mathrm{L}^{\mathrm{X}}$ giving an action of G on $\mathrm{L}^{\mathrm{X}}$ which will be written $(g, f) \mapsto g \perp f$. This action is defined by the commutativity of the diagram:


### 5.7.2

Proposition The law $\perp: \mathrm{G} \times \mathrm{L}^{\mathrm{X}} \rightarrow \mathrm{L}^{\mathrm{X}}$ is given by the formula

$$
(g \perp f)(\xi)=g\left(f\left(g^{-1} \circ \xi\right)\right)
$$

Proof Let $\xi \in \mathrm{X}$ and $g \in \mathrm{G}$. The map $\varphi: f \mapsto(g \perp f)(\xi)$ from $\mathrm{L}^{\mathrm{X}}$ to L is a K -algebra homomorphism. Its kernel is a maximal ideal of $\mathrm{L}^{\mathrm{X}}$, and is of the form $\mathfrak{m}_{\eta}$, with $\eta \in \mathrm{X}$. The map $g_{*}$ in diagram (1) is $g$-linear; as $\gamma$ is L-linear, $f \mapsto g \perp f$ is $g$-linear, in other words, for $\lambda \in \mathrm{L}$ and $f \in \mathrm{~L}^{\mathrm{X}}, \varphi(\lambda f)=g(\lambda) \varphi(f)$. For all $f \in \mathrm{~L}^{\mathrm{X}}$, there exists $f_{1} \in \mathfrak{m}_{\eta}$ such that $f=f(\eta) \cdot 1+f_{1}$. Hence $\varphi(f)=g(f(\eta)) . \eta$ remains to be determined. For $\lambda \in \mathrm{L}$ and $a \in \mathrm{~A}$, by definition of $\perp$ :

$$
(g \perp(\widehat{\lambda \otimes a}))(\xi)=\widehat{(g(\lambda) \otimes a})(\xi)=g(\lambda) \cdot \xi(a)
$$

Besides,

$$
\begin{aligned}
& (g \perp(\widehat{\lambda \otimes a}))(\xi)=\varphi(\widehat{\lambda \otimes a})=g(\widehat{\lambda \otimes a})(\eta)=g(\lambda \cdot \eta(a))=g(\lambda) \cdot g(\eta(a)) \\
& \text { and so taking } \lambda=1, \xi(a)=g(\eta(a)) \text {. Therefore } \eta=g^{-1} \circ \xi
\end{aligned}
$$

### 5.7.3

Let L be a Galois extension of K and $\mathrm{G}=$ Aut $_{\mathrm{K}} \mathrm{L}$ its Galois group. Let G- $\operatorname{Setf}$ denote the category of finite sets on which G acts, and for X and Y objects of G - $\mathcal{S e t} \boldsymbol{\mathcal { P }}$, let $\operatorname{Hom}_{\mathrm{G}}(\mathrm{X}, \mathrm{Y})$ be the set of morphisms from X to Y , i.e. of maps $\varphi$ from X to Y such that, for all $x \in \mathrm{X}$ and $g \in \mathrm{G}, \varphi(g \cdot x)=g \cdot \varphi(x)$. Let $\mathscr{D}$ be the category of finite K-algebras diagonalized by L . For $\mathrm{A} \in \mathscr{D}$, let $\mathrm{S}(\mathrm{A})$ be the set $\operatorname{Hom}_{\mathrm{K} \text { - } \mathscr{A l g}}(\mathrm{A}, \mathrm{L})$, where G acts by $(g, f) \mapsto g \circ f$. This defines a contravariant functor S from $\mathscr{D}$ to $\mathrm{G}-\operatorname{Set} f$. For $f: \mathrm{A} \rightarrow \mathrm{B}$, the map $f^{*}: \mathrm{S}(\mathrm{B}) \rightarrow \mathrm{S}(\mathrm{A})$ is defined by

$$
f^{*}(\eta)=\eta \circ f
$$

The aim of this section is to prove the following result:
Theorem The functor $\mathrm{S}: \mathscr{D} \rightarrow$ G-Setf is an anti-equivalence of categories.

This theorem combines Propositions 5.7.4 and 5.7.5.

### 5.7.4

## Proposition The functor S is fully faithful.

In other words, for all algebras A and B of $\mathscr{D}$, the map $f \mapsto f^{*}$ from $\mathrm{Hom}_{\mathrm{K}-\mathscr{A} \ell g}$ $(A, B)$ to $\operatorname{Hom}_{G}(\mathrm{~S}(\mathrm{~B}), \mathrm{S}(\mathrm{A}))$ is bijective.

Lemma For all objects A of $\mathscr{D}, \operatorname{Fix}_{\mathrm{G}}\left(\mathrm{L} \otimes_{\mathrm{K}} \mathrm{A}\right)=\mathrm{A}$.
Proof of the Lemma The statement does not involve the algebra structure of A, only its vector space structure. Hence A can be identified with $\mathrm{K}^{d}$; the lemma is then a consequence of Theorem 5.6.3, (i) $\Rightarrow$ (iv).

Proof of the Proposition For any morphism $\varphi: \mathrm{S}(\mathrm{B}) \rightarrow \mathrm{S}(\mathrm{A})$ of G- $\operatorname{Set} \boldsymbol{f}$, there is an L-algebra morphism $\varphi^{*}: \mathrm{L}^{\mathrm{S}(\mathrm{A})} \rightarrow \mathrm{L}^{\mathrm{S}(\mathrm{B})}$ defined by $\varphi^{*}(h)=h \circ \varphi$ compatible with the action $\perp$ of G . In other words, $\varphi^{*}(g \perp h)=g \perp \varphi^{*}(h)$. Indeed, for $\xi \in \mathrm{S}(\mathrm{A})$,

$$
\begin{aligned}
\varphi^{*}(g \perp h)(\xi) & =(g \perp h)(\varphi(\xi))=g\left(h\left(g^{-1} \circ \varphi(\xi)\right)\right)=g\left(h\left(\varphi\left(g^{-1} \circ \xi\right)\right)\right) \\
& =(g \perp(h \circ \varphi))(\xi)=\left(g \perp \varphi^{*}(h)\right)(\xi)
\end{aligned}
$$

As the Gelfand transforms

$$
\gamma_{\mathrm{A}}: \mathrm{L} \otimes_{\mathrm{K}} \mathrm{~A} \rightarrow \mathrm{~L}^{\mathrm{S}(\mathrm{~A})} \text { and } \gamma_{\mathrm{B}}: \mathrm{L} \otimes_{\mathrm{K}} \mathrm{~B} \rightarrow \mathrm{~L}^{\mathrm{S}(\mathrm{~B})}
$$

are isomorphisms, there is a unique homomorphism

$$
\varphi^{\otimes}: \mathrm{L} \otimes_{\mathrm{K}} \mathrm{~A} \rightarrow \mathrm{~L} \otimes_{\mathrm{K}} \mathrm{~B}
$$

making the diagram

commutative. This homomorphism is compatible with the actions of G, and so $\varphi^{\otimes}$ induces a K -algebra homomorphism $f$ from

$$
\mathrm{A}=\operatorname{Fix}_{\mathrm{G}}\left(\mathrm{~L} \otimes_{\mathrm{K}} \mathrm{~A}\right) \text { to } \mathrm{B}=\operatorname{Fix}_{G}\left(\mathrm{~L} \otimes_{\mathrm{K}} \mathrm{~B}\right)
$$

We show that $f^{*}: \mathrm{S}(\mathrm{B}) \rightarrow \mathrm{S}(\mathrm{A})=\varphi$, i.e. that

$$
\left(\forall \eta \in \operatorname{Hom}_{\mathrm{K}-\operatorname{seg}}(\mathrm{B}, \mathrm{~L})\right) \quad \eta \circ f=\varphi(\eta)
$$

Let $\delta_{\eta}: \mathrm{L}^{\mathrm{S}(\mathrm{B})} \rightarrow \mathrm{L}$ be the homomorphism defined by $\delta_{\eta}(h)=h(\eta)$. By definition of the Gelfand transform, $\eta=\delta_{\eta} \circ \gamma_{\mathrm{B}} \circ \iota_{\mathrm{B}}$, where $\iota_{\mathrm{B}}: \mathrm{B} \rightarrow \mathrm{L} \otimes_{\mathrm{K}} \mathrm{B}$ is the canonical injection. Hence

$$
\begin{aligned}
\eta \circ f & =\delta_{\eta} \circ \gamma_{\mathrm{B}} \circ \iota_{\mathrm{B}} \circ f=\delta_{\eta} \circ \gamma_{\mathrm{B}} \circ \varphi^{\otimes} \circ \iota_{\mathrm{A}}=\delta_{\eta} \circ \varphi^{*} \circ \gamma_{\mathrm{A}} \circ \iota_{\mathrm{A}} \\
& =\delta_{\varphi(\eta)} \circ \gamma_{\mathrm{A}} \circ \iota_{\mathrm{A}}=\varphi(\eta)
\end{aligned}
$$

So indeed $f^{*}=\varphi$, which shows that the map $f \mapsto f^{*}$ is surjective.
We show that if $\varphi=u^{*}$ for some homomorphism $u$ from A to B , then $f=u$. It will follow that $u \mapsto u^{*}$ is injective. For all $\eta \in \mathrm{S}(\mathrm{B})$,

$$
\begin{aligned}
\delta_{\eta} \circ \gamma_{\mathrm{B}} \circ \iota_{\mathrm{B}} \circ u & =\eta \circ u=u^{*}(\eta)=\varphi(\eta)=\delta_{\varphi(\eta)} \circ \gamma_{\mathrm{A}} \circ \iota_{\mathrm{A}} \\
& =\delta_{\eta} \circ \varphi^{*} \circ \gamma_{\mathrm{A}} \circ \iota_{\mathrm{A}}=\delta_{\eta} \circ \gamma_{\mathrm{B}} \circ \iota_{\mathrm{B}} \circ f .
\end{aligned}
$$

As $\bigcap_{\eta \in \mathrm{S}(\mathrm{B})} \operatorname{Ker} \delta_{\eta}=0$ and $\gamma_{\mathrm{B}} \circ \iota_{\mathrm{B}}$ is injective, it follows that $u=f$.

### 5.7.5

Proposition The functor S is essentially surjective.
In other words, any finite G-set is isomorphic to a G-set of the form $S(A)$, where $\mathrm{A} \in \mathscr{D}$.

Proof (a) Let H be a subgroup of G , and set $\mathrm{F}=\mathrm{Fix}_{\mathrm{H}}(\mathrm{L})$. The restriction
 $\Omega$ of K; then any K-homomorphism of F to L extends to a homomorphism from L to $\Omega$ (Theorem 5.2.6), inducing an automorphism of L (Theorem 5.6.3, (i) $\Rightarrow$ (iii)). Let $f, g \in \mathrm{G}=\mathrm{Aut}_{\mathrm{K}}(\mathrm{L})$; then, $\rho(f)=\rho(g)$ if and only if $h \in \mathrm{G}$ defined by $g=f \circ h$ is in $\operatorname{Aut}_{\mathrm{F}}(\mathrm{L})$. Now, $\operatorname{Aut}_{\mathrm{F}}(\mathrm{L})=\mathrm{H}$ by (5.6.7), and so $\mathrm{S}(\mathrm{F})=\operatorname{Hom}_{\mathrm{K}-\mathscr{A l g}}(\mathrm{F}, \mathrm{L})$ can be identified with the set $\mathrm{G} / \mathrm{H}=\{f \mathrm{H}\}_{f \in \mathrm{G}}$, and this identification is compatible with the left actions of G.
(b) Let $A$ and $B$ be objects of $\mathscr{D}$. Then, $S(A \times B)=S(A) \sqcup S(B)$. Indeed, $L \otimes_{K}$ $\mathrm{A}=\mathrm{L}^{\mathrm{X}}$ with $\mathrm{X}=\mathrm{S}(\mathrm{A})$, and $\mathrm{L} \otimes_{\mathrm{K}} \mathrm{B}=\mathrm{L}^{\mathrm{Y}}$ with $\mathrm{Y}=\mathrm{S}(\mathrm{B})$, and

$$
\mathrm{L} \otimes_{\mathrm{K}}(\mathrm{~A} \times \mathrm{B})=(\mathrm{L} \otimes \mathrm{~A}) \times(\mathrm{L} \otimes \mathrm{~B})=\mathrm{L}^{\mathrm{X}\lrcorner \mathrm{Y}}
$$

So

$$
\mathrm{S}(\mathrm{~A} \times \mathrm{B})=\operatorname{Hom}_{\mathrm{L}-\Delta \ell_{g}}\left(\mathrm{~L}^{\mathrm{X}} \sqcup \mathrm{Y}, \mathrm{~L}\right)=\mathrm{X} \sqcup \mathrm{Y}
$$

(c) The proposition then follows from the fact that every finite G-set is of the form $\bigsqcup_{i \in \mathrm{I}} \mathrm{X}_{i}$, where I is finite, and where, for all $i \in \mathrm{I}, \mathrm{X}_{i}$ is isomorphic to a G-set of the form $\mathrm{G} / \mathrm{H}_{i}$, where $\mathrm{H}_{i}$ is a subgroup of G .

## Exercises 5.6 (Finite Galois extensions) and 5.7 (Finite Galois theory)

1.-Let $K$ be a field of characteristic $\neq 2$. Show that every extension of degree 2 of K is Galois and give its Galois group. What happens in characteristic 2?
2.-Let $K=\mathbb{C}(Z)$ and $L=K[X] /(P)$, where $P(X)=X^{3}-3 X+2 Z$.
(a) Show that L is an algebraic extension of K . Give its degree.
(b) Find the group of K-automorphisms of L. Is L a Galois extension of K?
(c) Let $\widetilde{\mathrm{L}}$ be a decomposition field of P (i.e. the field generated in an algebraic closure of $K$ by the roots of $P$ ). Show that $\widetilde{L}$ is a Galois extension of $K$. Give the degree of $\tilde{L}$ over $K$ and over $L$. Find the number of intermediate fields between $K$ and $\widetilde{\mathrm{L}}$.
3.-Let $\Omega$ be the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$. Set

$$
\varphi=\frac{1+\sqrt{5}}{2} \text { and } \varphi^{\prime}=\frac{1-\sqrt{5}}{2}
$$

and consider the field

$$
\mathrm{L}=\mathbb{Q}(\varphi)=\mathbb{Q}\left(\varphi^{\prime}\right)=\mathbb{Q}(\sqrt{5})
$$

(a) Show that fields $\mathrm{L}_{1}=\mathrm{L}(\sqrt{\varphi})$ and $\mathrm{L}_{2}=\mathrm{L}\left(\sqrt{\varphi^{\prime}}\right)$ of $\Omega$ are conjugate over $\mathbb{Q}$, but not over $L$. Show that $L_{1}$ is a Galois extension of $L$, which is a Galois extension of $\mathbb{Q}$, but that $L_{1}$ is not a Galois extension of $\mathbb{Q}$. Is this also true of $L_{2}$ ?
(b) Let M be the smallest Galois sub-extension of $\Omega$ containing $\mathrm{L}_{1}$ and $\mathrm{L}_{2}$. Give the degree of $M$ over $K$ ? Describe the Galois group $G=A u t_{K}(M)$. Is a normal subgroup of a normal subgroup of G necessarily a normal subgroup of G?
4.-Let K be of characteristic $p \geqslant 0$ and $\Omega$ an algebraic closure of K . Let the integer $n$ not be a multiple of $p$.
(a) Show that, for all $a \neq 0$, the polynomial $\mathrm{X}^{n}-a$ is separable.
(b) Show that the group $\mu_{n}$ of $n$-th roots ofèunity in $\Omega$ is isomorphic to $\mathbb{Z} /(n)$. In the following, assume that $\mu_{n} \subset \mathrm{~K}$.
(c) Let $a \in \mathrm{~K}$ be nonzero, and $\alpha$ a $n$-th root of $a$ in $\Omega$. Set $\mathrm{L}=\mathrm{K}(\alpha)$. Show that any conjugate of $\alpha$ can be written $\zeta \alpha$, where $\zeta \in \mu_{n}$. Deduce that L is Galois and that the group $\mathrm{G}=\operatorname{Aut}_{\mathrm{K}}(\mathrm{L})$ can be identified with a subgroup of $\mu_{n}$.
(d) Show that $\mathrm{X}^{n}-a \in \mathrm{~K}[\mathrm{X}]$ is irreducible if and only if there is no $m$-th root of $a$ of K , where $m \neq 1$ is a divisor of $n$. Show that then $\mathrm{G}=\mu_{n}\left(\prod_{g \in \mathrm{G}} g(\alpha)\right.$ may be considered).
5.-Let $E$ be the subfield of $\mathbb{C}$ generated by the roots of $X^{4}-5$. Give the Galois groups of E over

$$
\mathbb{Q}, \quad \mathbb{Q}(\sqrt{5}), \quad \mathbb{Q}(\sqrt{-5}), \quad \mathbb{Q}(i)
$$

6.-Let $\mathrm{P} \in \mathrm{K}[\mathrm{X}]$ a monic polynomial of degree $d, \Omega$ an algebraic closure of the field K and $x_{1}, \ldots, x_{d}$ the roots of P in $\Omega$. Set

$$
\Delta=\prod_{i<j}\left(x_{i}-x_{j}\right)^{2} .
$$

(a) Show that, up to sign, $\Delta$ is the discriminant of P (resulting from P and its derivative), and that in particular $\Delta \in \mathrm{K}$. Give the sign.
(b) Let E be the sub-extension of $\Omega$ generated by $x_{1}, \ldots, x_{d}$. Every $g \in \operatorname{Aut}_{K}(\mathrm{E})$ induces a permutation of $\left\{x_{1}, \ldots, x_{d}\right\}$, hence, assuming these elements are all distinct, i.e. $\Delta \neq 0$, an embedding of $\mathrm{G}=\mathrm{Aut}_{\mathrm{K}} \mathrm{E}$ in $\mathfrak{S}_{d}$. Show that $\Delta$ is a square in K if and only if G is contained in the alternating group $\mathfrak{A}_{d}$ (see 5.8.7). More generally, show that $\mathrm{G} \cap \mathfrak{A}_{d}=$ Aut $_{\mathrm{L}} \mathrm{E}$, where $\mathrm{L}=\mathrm{K}(\sqrt{\Delta})$.
(c) Suppose that $\mathrm{G}=\mathfrak{S}_{d}$ and $d \geqslant 5$. Show that L is the only strictly intermediate Galois extension between K and E . Study the cases $d=4, d=3$.
7.-Let $k$ be a non-perfect field of characteristic $p, a \in k, a \notin k^{p}$ and $\mathrm{K}=k(\mathrm{~T})$, where T is an indeterminate. Consider the algebraic extension $\mathrm{L}=\mathrm{K}(x)$ of K generated by a root of

$$
f(\mathrm{X})=\mathrm{X}^{p}-\mathrm{T}^{p-1} \mathrm{X}-a
$$

(a) Show that L is a Galois extension of K .
(b) Show that K is not perfect.
(c) Let $b \in k\left(b \neq a\right.$ and $\left.b \notin k^{p}\right)$, and $\mathrm{Q}=\mathrm{L}(y)$ the extension of L generated by a root $y$ of

$$
g(\mathrm{X})=\mathrm{X}^{p-1}-\mathrm{T}^{p-1} \mathrm{X}-b
$$

Study the Galois extension Q of K : Galois group, intermediate extensions.
8.-(A) Let K be an infinite field, E a Galois extension of K of degree $d$, G the Galois group $\mathrm{Aut}_{\mathrm{K}} \mathrm{E}$ and $g_{1}, \ldots, g_{d} \in \mathrm{G}$.
(a) Let $\left(e_{1}, \ldots, e_{d}\right)$ be a basis for the vector space E over K . Show that

$$
x \mapsto \operatorname{det}_{\left(e_{1}, \ldots, e_{d}\right)}\left(g_{1}(x), \ldots, g_{d}(x)\right)
$$

is a K-polynomial map that does not vanish everywhere (use 5.1.7). Deduce that there exists $x \in \mathrm{E}$ whose conjugates form a basis for the K -vector space E (such a basis is called a normal basis).
(b) Show that the elements $g_{i}(x)$ form a normal basis if and only if the determinant of the matrix

$$
\left(g_{i} \circ g_{j}(x)\right)_{\substack{1 \leqslant i \leqslant d \\ 1 \leqslant j \leqslant d}}
$$

with entries in E is nonzero.
(B) Let K be a finite field with $q=p^{r}$ elements, where $p$ is prime.
(a) Show that $\operatorname{Aut}_{\mathbb{F}_{p}} \mathrm{~K}$ is cyclic of order $r$, generated by the Frobenius automorphism.
(b) Show that, for any $r^{\prime}$ dividing $r$, there is a unique subfield $\mathrm{K}^{\prime}$ of K with $p^{r^{\prime}}$ elements. Show that $\mathrm{Aut}_{\mathrm{K}^{\prime}} \mathrm{K}$ is cyclic and give a generator.
(C) (a) Let E be a vector space of finite dimension $d$ over a field $k$ and $\varphi$ an endomorphism of E . The minimal polynomial of $\varphi$ is assumed to be of degree $d$. Show that there exists $x \in \mathrm{E}$ such that the smallest vector subspace of E preserved by $\varphi$ containing $x$ is E (the results on principal ring modules may be used).

Show that then $\left(x, \varphi(x), \ldots, \varphi^{d-1}(x)\right)$ is a basis for E .
(b) Let K be a field and E a finite Galois extension of E such that the Galois group $\mathrm{G}=\mathrm{Aut}_{\mathrm{K}} \mathrm{E}$ is cyclic of order $d$. If $\varphi$ is a generator of G , show that the minimal polynomial of $\varphi$ is $\mathrm{X}^{d}-1$ (5.1.7 may be used).
(c) Prove that all finite fields have a normal basis.
9.-Let $K$ be a field, $K^{\prime}$ a not necessarily algebraic extension of $K$, and $P \in K[X]$ an irreducible separable polynomial. Its coefficients may be taken to be in $\mathrm{K}^{\prime}$, but it is not necessarily irreducible in $\mathrm{K}^{\prime}[\mathrm{X}]$. Let $\Omega$ and $\Omega^{\prime}$ be respective algebraic closures of K and $\mathrm{K}^{\prime}, \mathrm{E}$ the sub-extension of $\Omega$ generated by the roots of P and $\mathrm{E}^{\prime}$ a sub-extension of $\Omega^{\prime}$ generated by a root $x^{\prime}$ of P and its conjugates.
(a) Set $\mathrm{L}=\mathrm{K}(x)$, where $x$ is a root of P , and $\mathrm{L}^{\prime}=\mathrm{K}^{\prime}\left(x^{\prime}\right)$. Show that $\mathrm{L}^{\prime}$ is isomorphic to a quotient of $K^{\prime} \otimes_{K} L$.
(b) Show that $G=A u t_{K} E$ can be identified with the group of permutations of

(c) Show that $\mathrm{G}^{\prime}=\operatorname{Aut}_{\mathrm{K}^{\prime}}\left(\mathrm{E}^{\prime}\right)$ is isomorphic to a subgroup of G .

### 5.8 Solvability

In this section, K is a field of characteristic $p$ with $p=0$ orprime, and $\Omega$ an algebraic closure of K .

### 5.8.1

Proposition Let the integer $n$ not be a multiple of $p$.
(a) The multiplicative group $\mu_{n}=\left\{u \in \Omega \mid u^{n}=1\right\}$ is isomorphic to the additive group $\mathbb{Z} /(n)$.
(b) The K-sub-extension $\mathrm{M}_{n}$ of $\Omega$ generated by $\mu_{n}$ is Galois.
(c) The Galois group of $\mathrm{Aut}_{\mathrm{K}} \mathrm{M}_{n}$ can be identified with the subgroup of the multiplicative group $\mathbb{Z} /(n)^{*}$ of invertible elements in the ring $\mathbb{Z} /(n)$. In particular, Aut $_{\mathrm{K}}\left(\mathrm{M}_{n}\right)$ is commutative.

Proof (a) By Proposition 5.2.10, the group $\mu_{n}$ is cyclic. The roots of $Z^{n}-1$ are simple, the derivative $n Z^{n-1}$ only vanishes at 0 (because of the assumption that $p$ does not divide $n$ ), and hence there are $n$ of them in $\Omega$ and $\operatorname{Card}\left(\mu_{n}\right)=n$.
(b) The extension $\mathrm{M}_{n}$ is etale (Proposition 5.4.9) and invariant under all automorphisms of $\Omega$. Hence it is Galois (Theorem 5.6.3).
(c) The map sending $m \in \mathbb{Z}$ to the endomorphism $x \mapsto m \cdot x$ of $\mathbb{Z} /(n)$ is a ring homomorphism from $\mathbb{Z}$ to $\operatorname{End}(\mathbb{Z} /(n))$, with respect to multiplication in the initial set and composition in the end one. Taking quotients, this gives an isomorphism $\mathbb{Z} /(n) \rightarrow \operatorname{End}(\mathbb{Z} /(n))$ assigning to the elements of $\mathbb{Z} /(n)^{*}$ the automorphisms of $\mathbb{Z} /(n)$.

Transfer using the isomorphism $\phi: \mathbb{Z} /(n) \rightarrow \mu_{n}$ shows that the map from $\mathbb{Z}$ to $\operatorname{End}\left(\mu_{n}\right)$ sending $m$ to the map $u \mapsto u^{m}$ gives a group isomorphism from $\mathbb{Z} /(n)$ onto $\operatorname{End}\left(\mu_{n}\right)$, which in turn induces a group isomorphism from $\mathbb{Z} /(n)^{*}$ (equipped with multiplication) onto $\operatorname{Aut}\left(\mu_{n}\right)$ (equipped with composition). Note that this isomorphism does not depend on the choice of $\phi$.

Any automorphism of the extension $\mathrm{M}_{n}$ of K induces a group automorphism $\mu_{n}$, and this action is injective since $\mathbf{M}_{n}$ is generated by $\mu_{n}$. Hence Aut $_{K}\left(\mathrm{M}_{n}\right)$ can be identified with a subgroup of $\operatorname{Aut}\left(\mu_{n}\right)$, which can itself be identified with $\mathbb{Z} /(n)^{*}$. No choice is involved in these identifications.

Remark For $\mathrm{K}=\mathbb{Q}$, the group $\operatorname{Aut}\left(\mathrm{M}_{n}\right)$ can be shown to be identified with the whole of $\mathbb{Z} /(n)^{*}$ (irreductibility of cyclotomic polynomials, 5.8, Exercise 2).

### 5.8.2

For $a \in \mathrm{~K}^{*}$ and $n \in \mathbb{N}^{*}$, let $\mathrm{R}_{\mathrm{K}}(a, n)$, or simply $\mathrm{R}(a, n)$, be the extension of K generated by the $n$-th roots of $a$ in $\Omega$, i.e. by the elements $x \in \Omega$ such that $x^{n}=a$. If $n$ is not a multiple of $p$, then the extension $\mathrm{R}(a, n)$ is Galois: Indeed it is the decomposition of the separable polynomial $\mathrm{X}^{n}-a$.

If $x$ is an $n$-th root of $a$, the other $n$-th roots of $a$ are $u . x$ for $u \in \mu_{n}$. Hence $\mathrm{M}_{n}=\mathrm{R}_{\mathrm{K}}(1, n) \subset \mathrm{R}_{\mathrm{K}}(a, n)$.

More generally, let $a_{1}, \ldots, a_{k} \in \mathrm{~K}$ and $n_{1}, \ldots, n_{k}$ integers $>0$. Write $\mathrm{R}_{\mathrm{K}}\left(a_{1}\right.$, $\left.n_{1} ; \ldots ; a_{k}, n_{k}\right)$ for the extension of K generated by $\mathrm{R}_{\mathrm{K}}\left(a_{1}, n_{1}\right), \ldots, \mathrm{R}_{\mathrm{K}}\left(a_{k}, n_{k}\right)$ in $\Omega$. This extension contains $\mathrm{M}_{n}$ for $n=$ 1.c.m. $\left(n_{1}, \ldots, n_{k}\right)$. If no $n_{i}$ is a multiple of $p$, this extension is Galois, being generated by Galois extensions.

Proposition Let $\mathrm{E}=\mathrm{R}_{k}\left(a_{1}, n_{1} ; \ldots ; a_{k}, n_{k}\right)$ and $n=1 . \mathrm{c} . \mathrm{m} .\left(n_{1}, \ldots, n_{k}\right)$. Suppose that $\mu_{n}$ is contained in K and that $n$ is not a multiple of $p$. Then the Galois group of $\mathrm{G}=\operatorname{Aut}_{\mathrm{K}}(\mathrm{E})$ can be identified with a subgroup of $\mu_{n_{1}} \times \cdots \times \mu_{n_{k}}$. In particular, G is commutative.

Proof Let $g \in \mathrm{G}$. For each $i \in\{1, \ldots, k\}$, let $x_{i}$ be a $n_{i}$-th root of $a_{i}$. Then $g\left(x_{i}\right)=$ $u_{i} \cdot x_{i}$ for some $u_{i} \in \mu_{n_{i}}$ which does not depend on the choice of $x_{i}$. For if $x_{i}^{\prime}$ is another $n_{i}$-th root of $a_{i}$, then $x_{i}^{\prime}=w_{i} \cdot x_{i}$ with $w_{i} \in \mu_{n_{i}}$, and $g\left(x_{i}^{\prime}\right)=w_{i} \cdot g\left(x_{i}\right)$ since $w_{i} \in \mathrm{~K}$. Assigning to each $g \in \mathrm{G}$ the family $\left(u_{1}, \ldots, u_{k}\right)$ defines a map $\varphi: \mathrm{G} \rightarrow$ $\mu_{n_{1}} \times \cdots \times \mu_{n_{k}}$ which is injective since E is generated by the elements $x_{i}$.

It is also a group homomorphism. For if $\varphi(g)=\left(u_{1}, \ldots, u_{k}\right)$ and $\varphi(h)=\left(v_{1}, \ldots\right.$, $v_{k}$, then $g \circ h\left(x_{i}\right)=u_{i} \cdot v_{i} \cdot x_{i}$.

### 5.8.3

Proposition Let the integer $n$ not be a multiple of $p$. Suppose $\mu_{n} \subset K$.
An extension E of K contained in $\Omega$ is of the form $\mathrm{R}_{\mathrm{K}}(a, n)$ if and only if E is Galois and $\mathrm{G}=\mathrm{A} u t_{\mathrm{K}}(\mathrm{E})$ cyclic of order dividing $n$.

Proof (a) Necessity: It is Proposition (5.8.2), with $k=1$.
(b) Sufficiency: Suppose G is cyclic of order $m$ dividing $n$; let $g$ be a generator of G. Consider E as a vector space over K and $g$ as a K-linear endomorphism of E.

The endomorphism $g$ satisfies $g^{m}-\mathrm{I}=0$, and the roots of $\mathrm{Z}^{m}-1$ are simple and in K. Hence $g$ is diagonalizable; in other words, E has a basis consisting of eigenvectors of $g$. The eigenvalues are in $\mu_{m}$. If $x$ and $y$ are eigenvectors of $g$ with respective eigenvalues $u$ and $v$, then

$$
g(x y)=g(x) \cdot g(y)=u x \cdot v y=u v \cdot x y
$$

and $x y$ are eigenvectors for the eigenvalue $u v$. Consequently, the eigenvalues of $g$ form a subgroup of $\mu_{m}$.

Moreover, $y \mapsto x y$ is a linear injection of the eigenspace corresponding to $v$ into the eigenspace corresponding to $u v$. As a result, the eigenspaces corresponding to the various eigenvalues all have the same dimension.

But the eigenspace corresponding to the eigenvalue 1 is $\operatorname{Fix}_{G}(\mathrm{E})=\mathrm{K}$, and so has dimension 1. As E is an $m$-dimensional vector space, there are $m$ distinct eigenvalues, and all elements of $\mu_{m}$ are eigenvalues.

Let $w$ be a generator of $\mu_{m}$ and $z$ an eigenvector of $g$ for $w$. Then $z^{i}, i=0, \ldots$, $m-1$, are the eigenvectors for all the eigenvalues of $g$. Hence they form a basis for E . In particular E is generated by $z$ as extension of K .

Set $b=z^{m}$ and $a=z^{n}$. The $m$-th roots of $b$ (resp. the $n$-th roots of $a$ ) are $u \cdot z$ with $u \in \mu_{m}$ (resp. $u \in \mu_{n}$ ). Hence they are in E and

$$
\mathrm{E}=\mathrm{R}(b, m)=\mathrm{R}(a, n)
$$

Remark If $\mathrm{K}=\mathbb{Q}$, G has order $n$, and E is then called a Kummer extension.

### 5.8.4 Solvable Extensions

Definition Let $L \subset \Omega$ be a finite extension of $K$. Then $L$ is said to be a solvable extension of $K$ if there is a finite sequence ( $\mathrm{L}_{0}, \ldots, \mathrm{~L}_{\mathrm{N}}$ ) of extensions with $\mathrm{K}=\mathrm{L}_{0} \subset$ $\mathrm{L}_{1} \subset \cdots \subset \mathrm{~L}_{\mathrm{N}} \subset \Omega, \mathrm{L} \subset \mathrm{L}_{\mathrm{N}}$ and $\mathrm{L}_{i+1}$ is of the form $\mathrm{R}_{\mathrm{L}_{i}}\left(a_{i, 1}, n_{i, 1} ; \ldots ; a_{i, k_{i}}, n_{i, k_{i}}\right)$ for $0 \leqslant i \leqslant \mathrm{~N}-1$.

This definition implies that a sub-extension of a solvable extension of K is solvable, an extension generated by finitely many solvable extensions is solvable, and a solvable extension of a solvable extension of K is a solvable extension of K .

In this definition, imposing $k_{i}=1$ for all $i$ does not restrict the notion of solvability,
A finite radical extension, i.e. a purely inseparable one, is solvable. An extension $L$ of $K$ is solvable if and only if so is $L_{\text {sep }}$ : indeed $L_{\text {sep }}$ is a sub-extension of $L$ and $L$ is a radical extension of $\mathrm{L}_{\text {sep }}$.

### 5.8.5 Separable Solvable Extensions

A separable extension $L$ of $K$ is solvable if and only if the Galois extension $E$ generated by $L$ is solvable. Indeed $L$ is a sub-extension of E , and E is generated by L and its conjugates under the action of $\operatorname{Aut}_{\mathrm{K}}(\Omega)$.

If $L$ is a separable solvable extension of $K$, then there is a sequence $\left(L_{0}, \ldots, L_{N}\right)$ satisfying the conditions of definition 5.8.4., where no $n_{i, j}$ is a multiple of $p$. This follows from the next lemma.

Lemma If $n=p^{r} m$ where $m$ is not a multiple of $p$, then for $a \in \mathrm{~K}^{*} \mathrm{R}(a, n)_{\text {sep }}=$ $\mathrm{R}(a, m)$.

Proof The extension $\mathrm{R}(a, m)$ is separable, and $\mathrm{R}(a, n)$ is a purely inseparable extension of $\mathrm{R}(a, m)$ since it is generated by the $p^{r}$-th roots of the $m$-th roots of $a$.

If $L$ is a separable solvable extension of $K$, there is a sequence $\left(L_{0}, \ldots, L_{N}\right)$ satisfying conditions of definition 5.8.4, with all $\mathrm{L}_{i}$ Galois. Indeed, by the previous reduction, all the extensions $L_{i}$ may be assumed to be separable. Then, each time $n$-th roots of an element are added, we also add those of all its conjugates over K .

### 5.8.6 Solvable Groups

Definition A group $G$ is said to be solvable if there is a decreasing finite sequence $\left(\mathrm{G}_{0}, \ldots, \mathrm{G}_{\mathrm{N}}\right)$ of subgroups of G with $\mathrm{G}_{0}=\mathrm{G}, \mathrm{G}_{\mathrm{N}}=\{e\}, \mathrm{G}_{i+1}$ normal in $\mathrm{G}_{i}$ and $\mathrm{G}_{i} / \mathrm{G}_{i+1}$ commutative for $0 \leqslant i \leqslant \mathrm{~N}-1$.

A subgroup of a solvable is solvable. A quotient group of a solvable group is solvable. Given a group G and a normal subgroup H , if H and $\mathrm{G} / \mathrm{H}$ are solvable, then so is G.

If $G$ is a finite solvable group, then there is a sequence $\left(G_{0}, \ldots, G_{N}\right)$ satisfying the conditions of the above definition with $\mathrm{G}_{i} / \mathrm{G}_{i+1}$ cyclic for all $i$. This follows from the fact that every finite commutative group is the direct sum of cyclic groups (Theorem 3.5.8, with $\mathrm{A}=\mathbb{Z}$ ).

### 5.8.7 The Groups $\mathfrak{S}_{\boldsymbol{n}}$ and $\mathfrak{A}_{\boldsymbol{n}}$

If the set X is finite, then let $\mathfrak{S}(\mathrm{X})$ be the permutation group on X (symmetric group), and $\mathfrak{A}(\mathrm{X})$ the kernel of the homomorphism $\epsilon: \mathfrak{S}(\mathrm{X}) \rightarrow\{+1,-1\}$ which assigns to each permutation its signature (3.3.18) (alternating group). For $\mathrm{Y} \subset \mathrm{X}$, embed $\mathfrak{S}(\mathrm{Y})$ in $\mathfrak{S}(\mathrm{X})$ by extending the permutations by the identity on $\mathrm{X}-\mathrm{Y}$. Then $\mathfrak{A}(\mathrm{Y})=\mathfrak{S}(\mathrm{Y}) \cap \mathfrak{A}(\mathrm{X})$. If X has $n$ elements, then $\mathfrak{S}(\mathrm{X})$ and $\mathfrak{A}(\mathrm{X})$ are respectively isomorphic to $\mathfrak{S}_{n}=\mathfrak{S}(\{1, \ldots, n\})$ and $\mathfrak{A}_{n}=\mathfrak{A}(\{1, \ldots, n\})$.

Let the set X be finite, $k \geqslant 2$ and $a_{1}, \ldots, a_{k}$ distinct points of X. Write $\left[a_{1}, \ldots, a_{n}\right]$ for the permutation $\zeta$ defined by $\zeta\left(a_{i}\right)=a_{i+1}$ for $1 \leqslant i \leqslant k-1, \zeta\left(a_{k}\right)=a_{1}$ and $\zeta(x)=x$ for $x \in \mathrm{X}-\left\{a_{1}, \ldots, a_{k}\right\}$. Such a permutation is called a $k$-cycle and
$\left\{a_{1}, \ldots, a_{k}\right\}$ is the support of $\zeta$. Cycles with disjoint supports commute. Every permutation on X can, up to order, be uniquely written as a product of cycles with disjoint supports.

If $\sigma=\zeta_{1} \ldots \zeta_{\nu}$ with disjoint supports, where $\zeta_{i}$ is a $k_{i}$-cycle and $k_{1} \geqslant \cdots \geqslant k_{\nu}$, then $\sigma$ will be said of type $\left(k_{1}, \ldots, k_{\nu}\right)$. Two permutations on X are conjugate in $\mathfrak{S}(\mathrm{X})$ if and only they are of the same type.

If $\sigma$ and $\sigma^{\prime}$ are of the same type $\left(k_{1}, \ldots, k_{n}\right)$ and if some $k_{i}$ is even, or any two $k_{i}$ are equal, or $k_{1}+\cdots+k_{\nu} \leqslant n-2$, then the permutations $\sigma$ and $\sigma^{\prime}$ are conjugate by an element of $\mathfrak{A}(\mathrm{X})$. In all three cases, there exists $\tau \in \mathfrak{S}(\mathrm{X})-\mathfrak{A}(\mathrm{X})$ commuting with $\zeta$. However, if $\sigma^{\prime}=\alpha \sigma \alpha^{-1}$, then $\sigma^{\prime}=\alpha^{\prime} \sigma \alpha^{\prime-1}$ when $\alpha^{\prime}=\alpha \tau$ with $\tau$ commuting with $\sigma$; if moreover $\epsilon(\tau)=-1$, one of $\alpha, \alpha^{\prime}$ is in $\mathfrak{A}(\mathrm{X})$.

### 5.8.8 Simplicity of $\mathfrak{A}_{\boldsymbol{n}}$ for $n \geqslant 5$

A group G is said to be simple if its only normal subgroups are G and $\{e\}$.
Proposition Let X be an $n$ element set with $n \geqslant 5$. Then the group $\mathfrak{A}(\mathrm{X})$ is simple.
This proposition follows from the next two lemmas.
Lemma 5.1 Let H be a nontrivial normal subgroup of $\mathfrak{A}(\mathrm{X})$. Then H contains a 3 -cycle or an element of type $(2,2)$.

Proof Let $\sigma \in \mathrm{H}$ be a nontrivial element, $\left(k_{1}, \ldots, k_{\nu}\right)$ the type of $\sigma$ and write $\sigma=$ $\zeta_{1} \cdots \zeta_{\nu}$, where $\zeta_{i}$ is a $k_{i}$-cycle and the supports are all disjoint.

If there is some $k_{i} \neq 3$, then there is a 3 element set $\Delta$ such that $\mathrm{Y}=\Delta \cup \sigma(\Delta)$ has 4 elements. Indeed if $k_{1}>3$ and $\zeta_{1}=\left[a_{1}, \ldots, a_{k_{1}}\right]$ then we may take $\Delta=\left\{a_{1}, a_{2}, a_{3}\right\}$; if $k_{\nu}=2$, then $\nu>1$ for otherwise $\epsilon(\sigma)=-1$ and we could take $\Delta=\operatorname{supp}\left(\zeta_{\nu}\right) \cup$ $\{x\}$, where $x$ is not a fixed element.

In this case, let $\alpha$ be a cycle of support $\Delta$. Then $\alpha^{\prime}=\sigma \alpha \sigma^{-1}$ is a cycle of support $\sigma(\Delta)$, and $\gamma=\alpha^{\prime} \alpha^{-1}$ is in $\mathfrak{A}(\mathrm{Y})$ and is nontrivial. Hence $\gamma$ is a 3-cycle or an element of type (2,2). But $\gamma$ can be written as $\sigma \sigma^{\prime-1}$, where $\sigma^{\prime}=\alpha \sigma \alpha^{-1} \in \mathrm{H}$, and so $\gamma \in \mathrm{H}$. This proves the lemma in this case.

It remains to prove the case where all $k_{i}=3$. If $\nu=1, \sigma$ itself is a 3-cycle. If $\nu \geqslant 2$, write $\zeta_{1}=\left[a_{1}, a_{2}, a_{3}\right], \zeta_{2}=\left[b_{1}, b_{2}, b_{3}\right]$, and set $\alpha=\left[a_{1}, b_{1}\right]\left[a_{2}, b_{2}\right]$. Then $\sigma^{\prime}=\alpha \sigma \alpha^{-1}=\left[b_{1}, b_{2}, a_{3}\right]\left[a_{1}, a_{2}, b_{3}\right] \zeta_{3} \cdots \zeta_{\nu}$, and $\gamma=\sigma \sigma^{\prime-1}=\left[a_{1}, b_{1}\right]\left[a_{3}, b_{3}\right] \in$ H.

Lemma 5.2 Let H be a normal subgroup of $\mathfrak{A}(\mathrm{X})$ with $\operatorname{Card}(\mathrm{X}) \geqslant 5$. If H contains a 3 -cycle or an element of type $(2,2)$, then $\mathrm{H}=\mathfrak{A}(\mathrm{X})$.

Proof First note that if there is an element of type (2,2), then there is also a 3-cycle. Indeed, if $\sigma=\left[a_{1}, a_{2}\right]\left[b_{1}, b_{2}\right] \in \mathrm{H}$, choose $x \in \mathrm{X}$ different from $a_{1}, a_{2}, b_{1}, b_{2}$, which is possible since $\operatorname{Card}(\mathrm{X}) \geqslant 5$. Set $\alpha=\left[a_{1}, a_{2}, x\right]$. Then

$$
\gamma=\sigma \alpha \sigma^{-1} \alpha^{-1}=\left[x, a_{2}, a_{1}\right] \in \mathrm{H}
$$

Hence this is the case where H contains a 3-cycle. Then all 3-cycles are in H since all of them are conjugate in $\mathfrak{A}(\mathrm{X})$. The proof then follows by a "jeu de taquin".

Take $\mathrm{X}=\{1, \ldots, n\}$, and let $\sigma \in \mathfrak{A}_{n}$. We show that $\sigma \in \mathrm{H}$. Construct by induction a sequence $\left(\sigma_{i}\right)_{0 \leqslant i \leqslant n-2}$ such that $\sigma_{i}$ fixes $1, \ldots, i$ with $\sigma_{0}=\sigma$ and $\sigma_{i}=h_{i} \sigma_{i-1}$, where $h_{i} \in \mathrm{H}$ for $1 \leqslant i \leqslant n-2$. If $\sigma_{i-1}(i)=i$, then take $h_{i}=\mathrm{I}$; otherwise take $h_{i}=\left[\sigma_{i-1}(i), i, x\right]$, where $x \in\{i+1, \ldots, n\}$ is different from $\sigma_{i-1}(i)$, which is possible. Then $\sigma_{n-2}$ fixes $1, \ldots, n-2$, and as $\epsilon\left(\sigma_{n-2}\right)=1, \sigma_{n-2}=\mathrm{I}$, and so $\sigma=$ $\left(h_{n-2} h_{n-3} \ldots h_{1}\right)^{-1} \in \mathrm{H}$.

### 5.8.9

Corollary For an $n$ element set X with $n \geqslant 5$, the only normal subgroups of $\mathfrak{S}(\mathrm{X})$ are $\{\mathrm{I}\}, \mathfrak{A}(\mathrm{X})$ and $\mathfrak{S}(\mathrm{X})$.

Proof Let H be a normal subgroup of $\mathfrak{S}(\mathrm{X})$. Then $\mathrm{H}^{\prime}=\mathrm{H} \cap \mathfrak{A}(\mathrm{X})$ is a normal subgroup of $\mathfrak{A}(X)$, and so is $\mathfrak{A}(X)$ or $\{I\}$.

If $H^{\prime}=\mathfrak{A}(X)$, i.e. if $H \supset \mathfrak{A}(X)$, then $H=\mathfrak{A}(X)$ or $H=\mathfrak{S}(X)$, since the only subgroups of $\mathfrak{S}(\mathrm{X}) / \mathfrak{A}(\mathrm{X})=\{+1,-1\}$ are $\{1\}$ and the whole group.

If $\mathrm{H}^{\prime}=\{\mathrm{I}\}$, then H has at most 2 elements. If H is nontrivial, then the nontrivial element must be left invariant by all inner automorphisms of $\mathfrak{S}(\mathrm{X})$, hence is central, which is impossible.

### 5.8.10

Theorem A finite Galois extension E of K is a solvable extension if and only if the Galois group $\mathrm{G}=\operatorname{Aut}_{\mathrm{K}}(\mathrm{E})$ is solvable and its order is not a multiple of $p$.

Proof (a) necessity: Suppose that E is solvable and that $\left(\mathrm{L}_{0}, \ldots, \mathrm{~L}_{\mathrm{N}}\right)$ satisfies the conditions of definition 5.8.4., where for all $i, \mathrm{~L}_{i}$ is a Galois extension of K and no $n_{i, j}$ is a multiple of $p$ for $i \in\{0, \ldots, \mathrm{~N}\}$ and $j \in\left\{1, \ldots, k_{i}\right\}$ (5.8.5). Set $\mathrm{K}^{\prime}=\mathrm{M}_{n}$, where $n$ is the l.c.m. of $n_{i, j}$.

For $i \in\{0, \ldots, \mathrm{~N}\}$, let $\mathrm{L}_{i}^{\prime}$ be the extension generated by $\mathrm{K}^{\prime}$ and $\mathrm{L}_{i}$ in $\Omega$. All $\mathrm{L}_{i}^{\prime}$ are Galois extensions of $\mathrm{K}^{\prime}$ (5.6.4, Corollary 1 ), and

$$
\mathrm{L}_{i+1}^{\prime}=\mathrm{R}_{\mathrm{L}_{i}^{\prime}}\left(a_{i, 1}, n_{i, 1} ; \ldots ; a_{i, k_{i}}, n_{i, k_{i}}\right)
$$

The Galois groups $\mathrm{G}_{i}^{\prime}$ of $\mathrm{Aut}_{\mathrm{L}_{i}^{\prime}}\left(\mathrm{L}_{\mathrm{N}}^{\prime}\right)$ are normal subgroups of $\mathrm{G}_{0}^{\prime}$ (5.6.4, Corollary 2$), \mathrm{G}_{\mathrm{N}}^{\prime}=\{e\}$ and by Proposition 5.8.2., for $i \in\{0, \ldots, \mathrm{~N}-1\}$, the quotient group $\mathrm{G}_{i}^{\prime} / \mathrm{G}_{i+1}^{\prime}=\operatorname{Aut}_{\mathrm{L}_{i}^{\prime}}\left(\mathrm{L}_{i+1}^{\prime}\right)$ is commutative. So the group $\mathrm{G}_{0}^{\prime}=\operatorname{Aut}_{\mathrm{K}^{\prime}}\left(\mathrm{L}_{\mathrm{N}}^{\prime}\right)$ is solvable.

By Proposition 5.8.1, (c), the group $\Gamma=\operatorname{Aut}_{K}\left(\mathrm{~K}^{\prime}\right)$ is commutative, and hence solvable. As $\Gamma=\mathrm{G}^{\prime \prime} / \mathrm{G}_{0}^{\prime}$, where $\mathrm{G}^{\prime \prime}=\operatorname{Aut}_{\mathrm{K}}\left(\mathrm{L}_{\mathrm{N}}^{\prime}\right)$ is solvable (5.8.6). The group $\mathrm{G}=$ $\mathrm{Aut}_{\mathrm{K}}(\mathrm{E})$ can be identified with a quotient of $\mathrm{G}^{\prime \prime}$ (5.6.4, Corollary 2), hence is also solvable.
(b) Sufficiency: Suppose that $G$ is solvable and that its order $n$ is not a multiple of $p$, and let $\left(\mathrm{G}_{0}, \ldots, \mathrm{G}_{\mathrm{N}}\right)$ be a sequence of subgroups satisfying the conditions of definition 5.8.6, with $\mathrm{G}_{i} / \mathrm{G}_{i+1}$ cyclic of order $n_{i}$. The order $n$ of G is the product of $n_{i}$, and so no $n_{i}$ is a multiple of $p$. Let $m$ be the l.c.m. of the integers $n_{i}, \mathrm{~K}^{\prime}=\mathrm{M}_{m}, \mathrm{E}^{\prime}$ the extension generated by $\mathrm{K}^{\prime}$ and E in $\Omega$, and $\mathrm{G}^{\prime}$ the Galois group of Aut $\mathrm{K}^{\prime}\left(\mathrm{E}^{\prime}\right)$. Setting $\mathrm{G}_{i}^{\prime}=\mathrm{G}^{\prime} \cap \mathrm{G}_{i}$, the group $\mathrm{G}_{i}^{\prime} / \mathrm{G}_{i+1}^{\prime}$ can be identified with a subgroup of $\mathrm{G}_{i} / \mathrm{G}_{i+1}$. Hence it is cyclic of order $n_{i}^{\prime}$ dividing $n_{i}$ as well as $m$.

Set $\mathrm{L}_{i}^{\prime}=\operatorname{Fix}_{\mathrm{G}_{i}^{\prime}}\left(\mathrm{E}^{\prime}\right)$. Then

$$
\mathrm{K}^{\prime}=\mathrm{L}_{0}^{\prime} \subset \cdots \subset \mathrm{L}_{\mathrm{N}}^{\prime}=\mathrm{E}^{\prime}
$$

$\mathrm{E}^{\prime}$ is a Galois extension of $\mathrm{L}_{i}^{\prime}$ for all $i$, with $\operatorname{Aut}_{\mathrm{L}_{i}}\left(\mathrm{E}^{\prime}\right)=\mathrm{G}_{i}^{\prime}$, and for $i \leqslant \mathrm{~N}-1, \mathrm{~L}_{i+1}^{\prime}$ is a Galois extension of $\mathrm{L}_{i}^{\prime}$ with $\operatorname{Aut}_{\mathrm{L}_{i}^{\prime}}\left(\mathrm{L}_{i+1}^{\prime}\right)=\mathrm{G}_{i}^{\prime} / \mathrm{G}_{i+1}^{\prime} \approx \mathbb{Z} /\left(n_{i}^{\prime}\right)$.

By Proposition 5.8.3, the extension $\mathrm{L}_{i+1}^{\prime}$ of $\mathrm{L}_{i}^{\prime}$ is of the form $\mathrm{R}_{\mathrm{L}_{i}^{\prime}}\left(a_{i}, n_{i}^{\prime}\right)$. Hence $\mathrm{E}^{\prime}$ is a solvable extension of $\mathrm{K}^{\prime}$. As $\mathrm{K}^{\prime}$ is a solvable extension of K , so is the case of $\mathrm{E}^{\prime}$ and hence also of its sub-extension E .

### 5.8.11 Example of a Non Solvable Extension of $\mathbb{Q}$

Let E be the decomposition field of the polynomial $f=\mathrm{X}^{5}-3 \mathrm{X}-1$ over $\mathbb{Q}$ in $\mathbb{C}$. We show that E is a non solvable extension of $\mathbb{Q}$.
(1) $f$ is irreducible in $\mathbb{Q}[\mathrm{X}]$. The polynomial $f$ has a root $z_{1}$ such that $\left|z_{1}\right|<1$, the others have absolute value $>1$. Indeed, as $z$ travels once round the unit circle in $\mathbb{C}, f(z)=-3 z\left(1-\frac{z^{5}-1}{3 z}\right)$ travels round 0 once. Were $f$ reducible, there would be $f_{1}$ and $f_{2}$ of degree $\geqslant 1$ in $\mathbb{Z}[\mathrm{X}]$ such that $f=f_{1} . f_{2}$ (3.7.8, Lemma 1$)$. The leading coefficients and the constant terms of $f_{1}$ and $f_{2}$ should be 1 or -1 , contradicting the fact that one of them should have absolute value $>1$.
(2) $f$ has 3 real roots. Considering the sign of $f(x)$ for $x=-2,-1,1$ and 2 , shows that $f$ has at least 3 real roots. The derivative $f^{\prime}=5 \mathrm{X}^{4}-3$ has only 2 real roots, and so $f$ has at most 3 real roots.
(3) $\mathrm{G}=\operatorname{Aut}_{\mathbb{Q}}(\mathrm{E})$ is isomorphic to $\mathfrak{S}_{5}$. The group G acts on the set $\Delta$ of roots of $f$, which has 5 elements. Hence there is a homomorphism $\mathrm{G} \rightarrow \mathfrak{S}(\Delta)$, which is injective since $\Delta$ generates E as an extension of $\mathbb{Q}$. Complex conjugation induces a transposition $\tau$ permuting the two non real roots of $f$. The group G acts transitively on $\Delta$ since $f$ is irreducible (Proposition 5.6.5); if S is the stabilizer of a point of $\Delta$, then $\operatorname{Card}(\mathrm{G}) / \operatorname{Card}(\mathrm{S})=5$, and so $\operatorname{Card}(\mathrm{G})$ is a multiple of 5 . By Sylow's theorem, G contains a subgroup of order 5 , necessarily generated by an element of order 5
which can only be a 5 -cycle. But as in $\mathfrak{S}_{5}$, a subgroup in $\mathfrak{S}(\Delta)$ containing a 2 -cycle and a 5 -cycle is the whole group. Hence $G$ can be identified with $\mathfrak{S}(\Delta)$.
(4) The group $\mathfrak{S}_{5}$ is not solvable. This follows from Corollary 5.8.9. As a result, E is a non solvable extension of $\mathbb{Q}$ by Theorem 5.8.10.

## Exercises 5.8. (Solvability)

1.-Let $p$ be a prime. A finite group G is a $p$-group if the order of G is a power of $p$.
(a) Let G be a $p$-group and X a finite G -set; set $\mathrm{Y}=$ Fix $_{\mathrm{G}} \mathrm{X}$. Show that card $\mathrm{Y} \equiv$ card $\mathrm{X}(\bmod p)$.
(b) Making G act on itself by inner automorphisms, show that the centre of G is not trivial.
(c) Deduce that G is solvable.
2. (Cyclotomic polynomials)- Let $n$ be an integer, $\mathrm{G}_{n}$ the group of invertible elements of the ring $\mathbb{Z} /(n)$, and $\varphi_{n}$ its cardinality. Suppose $\mathrm{G}_{n}$ is cyclic, and let $\Omega$ be an algebraic closure of $\mathbb{Q}$. Set $\Phi_{n}(\mathrm{X})=\prod\left(\mathrm{X}-\zeta_{i}\right)$, where $\zeta_{i}$ are the primitive $n$-th roots of unity. The aim is to show that $\Phi_{n} \in \mathbb{Q}[X]$ is irreducible (this property holds even when $\mathrm{G}_{n}$ is not cyclic, but the proof is different).
(a) Show that $\mathrm{X}^{n}-1=\prod \Phi_{d}$, where $d$ is a divisor of $n$. Deduce that $\Phi_{n} \in \mathbb{Z}[\mathrm{X}]$ (use the results of 3.7.8). For any field K , let $\Phi_{n}^{\mathrm{K}}$ be the image of $\Phi_{n}$ in $\mathrm{K}[\mathrm{X}]$.
(b) Show that for any field K whose characteristic $p$ does not divide $n$, the roots of $\Phi_{n}^{\mathrm{K}}$ in an algebraic closure of K are the primitive $n$-th roots of unity.
(c) Prove that if there is a prime $p$ such that $\Phi_{n}^{\mathbb{F}_{p}}$ is irreducible, then the irreducibility of $\Phi_{n}$ in $\mathbb{Q}[X]$ follows (use 3.7 .8 to reduce from the case $\mathbb{Q}$ to $\mathbb{Z}$ ).
(d) Let $p$ be a prime whose class modulo $n$ generates $\mathrm{G}_{n}$. Show that the field $\mathbb{F}_{p^{r}}$ has no primitive $n$-th root of unity for $r<\varphi_{n}$. Deduce that $\Phi_{n}^{\mathbb{F}_{p}}$ is irreducible.
(e) Conclude by assuming the following result: an arithmetic progression whose first term and difference are coprime contains infinitely many primes.
(f) Show that $\mathrm{L}=\mathbb{Q}[\mathrm{X}] /\left(\Phi_{n}\right)$ is a Galois extension of $\mathbb{Q}$ and that its Galois group can be identified with $\mathrm{G}_{n}$. Each embedding $\sigma$ of L in $\Omega$ defines an isomorphism of L onto the sub-extension E of $\Omega$ generated by the $n$-th roots of unity. Hence there is an isomorphism $\sigma_{*}$ from $\mathrm{G}_{n}$ onto $\operatorname{Aut}_{\mathbb{Q}}(\mathrm{E})$. Does this isomorphism depend on $\sigma$ ?
3.-Let $K$ be of characteristic different from 2 and 3 and $P \in K[X]$ an irreducible polynomial of degree 3 . Let E be the extension generated in an algebraic closure of K by roots of P and the cubic roots $1, j, j^{2}$ of 1 .
(a) Show that E is a solvable extension of K . More precisely, show that $\widetilde{\mathrm{K}}=\mathrm{K}(j)$ has degree 2 or 1 over $K$ and that there is an extension $F$ of degree 2 or 1 of $\widetilde{K}$ such that $\mathrm{E}=\mathrm{F}(\sqrt[3]{a})$ for some $a \in \mathrm{~F}$ is an extension of degree 3 of F .

In the following, the aim is to reduce the search for the roots of P to extractions of square or cubic roots.
(b) Let $x \in \mathrm{E}$ be a root of P and $\sigma \in \mathrm{Aut}_{\mathrm{K}}(\mathrm{E})$ an element of order 3. The other two roots of P are $y=\sigma(x)$ and $z=\sigma^{2}(x)$. Let $\mathrm{F}, \mathrm{V}$ and W be the eigenspaces of $\sigma$ considered a $\widetilde{\mathrm{K}}$-linear map from E to E for the eigenvalues $1, j$ and $j^{2}$ respectively. Let $v$ and $w$ be the respective projections of $x$ onto V and W . Express $v$ and $w$ in terms of $x, y, z$ and $j$. Show that $\mathrm{F}=\mathrm{K}\left(v^{3}\right)=\mathrm{K}\left(w^{3}\right)$, and that $v w \in \mathrm{~K}$.
(c) Suppose $\mathrm{P}=\mathrm{X}^{3}+p \mathrm{X}+q$. Show that $x=v+w$ and $v w=-\frac{p}{3}$.
(d) Replacing $x$ by $v+w$ in $x^{3}+p x+q=0$, calculate $v^{3}+w^{3}$. Deduce that $v^{3}$ and $w^{3}$ are the two roots of a second degree equation over K and give this equation.
(e) Express $x$ by radicals (Cardano's formula).
4.-Let K be of characteristic different from 2 and $3, \mathrm{P}$ an irreducible polynomial of degree 4 over K , and $x, y, z, t$ the roots of P in an algebraic closure.
(a) Show that $u=x y+z t, v=x z+y t$ and $w=x t+y z$ are the roots of a degree 3 polynomial over K.
(b) Show that $x y$ and $z t$ are the roots of a second degree equation over $\mathrm{K}(u)$. Show that so are $x+y$ and $z+t$.
(c) Show that $x+y$ and $x y$ generated the same quadratic extension L of $\mathrm{K}(u)$. Show that $x$ and $y$ are the roots of a second degree equation over L.
(d) Let $\mathrm{E}=\mathrm{K}(x, y, z, t)$. Assume that the Galois group $\mathrm{Aut}_{\mathrm{K}}(\mathrm{E})$ is isomorphic to $\mathfrak{S}_{4}$. Give the list of subgroups of $\mathfrak{S}_{4}$ and for each of them, give the generators of the corresponding sub-extension of E .
(e) What can be said about the Galois group $\operatorname{Aut}_{K}(\mathrm{E})$ if $\mathrm{P}(\mathrm{X})=\mathrm{X}^{4}+p \mathrm{X}^{2}+q$ (square of a quadratic)? if $\mathrm{P}(\mathrm{X})=a \mathrm{X}^{4}+b \mathrm{X}^{3}+c \mathrm{X}^{2}+b \mathrm{X}+a$ (reciprocal polynomial)?
5.-(a) Let $K$ and $L$ be subfields of $\mathbb{C}$ such that $K \subset L$ and $[L: K]=2$. Show that every point of L can be obtained from points of K by a ruler and compass construction.
(b) Let E be a subfield of $\mathbb{C}$ generated by $e^{2 i \pi / 17}$. Show that E is a Galois extension of $\mathbb{Q}$ and that the Galois group $G=A u t_{\mathbb{Q}} E$ can be identified with a group of automorphisms of $\mathbb{Z} /(17)$. Using 5.2.10, deduce that G is isomorphic to a cyclic group of order $2^{k}, k \leqslant 4$ (in fact $\mathrm{G} \approx \mathbb{Z} /(16)$ ). Show that there is a sequence $\mathrm{E}_{0}, \ldots, \mathrm{E}_{k}$ of subfields of E with $\mathbb{Q}=\mathrm{E}_{0} \subset \mathrm{E}_{1} \subset \cdots \subset \mathrm{E}_{k}=\mathrm{E}$ and $\left[\mathrm{E}_{i}: \mathrm{E}_{i-1}\right]=2$. Are the subfields $\mathrm{E}_{i}$ Galois extensions of $\mathbb{Q}$ ?
(c) Show that a regular polygon with 17 sides can be constructed using ruler and compass (an explicit construction is not asked for). What properties of the number 17 have been used?

Give a ruler and compass construction of a regular pentagon.

### 5.9 Infinite Galois Theory

In this section, $\Omega$ will denote an algebraic closure of the field $\mathrm{K}, \Omega_{\text {sep }}$ the separable closure of K in $\Omega$, and G the group $\mathrm{Aut}_{\mathrm{K}} \Omega$.

### 5.9.1

Proposition $G=$ Aut $_{\mathrm{K}} \Omega_{\text {sep }}$.
In other words, every automorphism of $\Omega$ induces an automorphism of $\Omega_{\text {sep }}$, and every automorphism of $\Omega_{\text {sep }}$ extends uniquely to an automorphism of $\Omega$.

Proof The existence of an extension follows from 5.2.7, Corollary 2 . As $\Omega$ is a purely inseparable extension of $\Omega_{\text {sep }}$ by 5.5.4, uniqueness is a consequence of 5.5.3 applied to $\Omega_{\text {sep }}$ and $\Omega$.

### 5.9.2 Profinite Structure of G

Let $\mathcal{E}$ be the set of finite Galois extensions of K in $\Omega$. Ordered by inclusion, it is a directed set: the sub-extension of $\Omega$ generated by two finite Galois extensions is etale and invariant under $G$, hence Galois (5.6.3). For objects E and F of $\mathcal{E}$ such that $\mathrm{E} \subset \mathrm{F}$, every K -automorphism of F induces a K-automorphism of E . Hence there is a restriction homomorphism $\rho_{\mathrm{E}}^{\mathrm{F}}: \mathrm{Aut}_{\mathrm{K}} \mathrm{F} \rightarrow \mathrm{Aut}_{\mathrm{K}} \mathrm{E}$. By 5.6.4, $\Omega_{\text {sep }}$ is the union of $\mathrm{E} \in \mathcal{E}$. As a result,

$$
\mathrm{G}=\operatorname{Aut}_{\mathrm{K}} \Omega_{\text {sep }}=\lim _{\underset{\mathrm{E} \in \mathcal{E}}{ }} \operatorname{Aut}_{\mathrm{K}} \mathrm{E} .
$$

Equip G with the topology obtained by transferring the topology of the projective limite via this identification. The group $G$ then becomes a profinite group.

### 5.9.3 The Functor S

Let A be an algebraic and separable algebra over K (see 5.4.4). Set $\mathrm{S}(\mathrm{A})=$ Hom $_{\mathrm{K} \text { - } \Delta \ell g}$ $\left(\mathrm{A}, \Omega_{\text {sep }}\right)=\operatorname{Hom}_{\mathrm{K}-\mathscr{A l g}_{g}}(\mathrm{~A}, \Omega)$. Indeed, let $x \in \mathrm{~A}, f \in \mathrm{~S}(\mathrm{~A})$, and P be the minimal polynomial of $x$; since $\mathrm{P}(f(x))=0$, the minimal polynomial of $f(x)$ divides P , and $f(x)$ is separable.

For any finite subalgebra $B$ of $A$, the set $\operatorname{Hom}_{K-A t g}(B, \Omega)$ is finite. Then $\mathrm{S}(\mathrm{A})=\lim \operatorname{Hom}_{\mathrm{K}-\Delta l_{g}}(\mathrm{~B}, \Omega)$, the inverse limit being taken over all finite subalgebras of $A$. This identification turns $S(A)$ into a profinite space.

The group G acts on $\mathrm{S}(\mathrm{A})$ by $(g, f) \mapsto g \circ f$. This action is continuous. Indeed, for any etale algebra $B$, there is a finite Galois extension $E$ diagonalizing $B$; then the action of $G$ on $S(B)$ factorizes through $\mathrm{Aut}_{\mathrm{K}}(\mathrm{E})$, hence is continuous. Therefore the inverse limit $S(A)$ of finite $G$-sets $S(B)$, for finite subalgebras B of A, is an object of G- $\mathscr{P r o f}$ (2.8.6).

This defines a contravariant functor from the category $\mathscr{A}$ of separable K-algebraic algebras to the category G-Prof of profinite spaces on which the action of G is continuous.

### 5.9.4

$$
\text { Theorem The functor } \mathrm{S}: \mathscr{A} \rightarrow \mathrm{G}-\mathscr{P r o f} \text { is an anti-equivalence of categories. }
$$

Proof Let $\mathscr{D}$ be the category of etale K-algebras, and for any finite Galois extension E of K in $\Omega$, let $\mathscr{D}_{\mathrm{E}}$ be the sub-category of $\mathscr{D}$ consisting of algebras diagonalized by E. Set $\mathrm{G}_{\mathrm{E}}=$ Aut $_{\mathrm{K}} \mathrm{E}$. By Theorem 5.7.3., the functor $\mathrm{S}_{\mathrm{E}}: \mathrm{A} \rightarrow \operatorname{Hom}_{\mathrm{K}-\Delta \operatorname{leg}_{g}}(\mathrm{~A}, \mathrm{E})=$


As $\mathscr{D}=\underset{\longrightarrow}{\lim } \mathscr{D}_{\mathrm{E}}$, and $\mathrm{G}=\underset{\leftrightarrows}{\lim } \mathrm{G}_{\mathrm{E}}$, ù $\mathrm{G}-\operatorname{Set} f=\underset{\longrightarrow}{\lim } \mathrm{G}_{\mathrm{E}}-\operatorname{Set} \boldsymbol{\operatorname { S e t }}$ by 2.9.4. The functors $\mathrm{S}_{\mathrm{E}}$ mutually induce each other, and so by passing to the direct limit, defines an anti-equivalence from $\mathscr{D}$ onto G- $\operatorname{Set} \mathcal{f}$. This anti-equivalence is readily seen to be given by the functor S . It is also immediate that $\mathscr{A}$ can be identified with the category $\xrightarrow[\rightarrow]{\mathscr{D}}$ defined by $(\underset{\rightarrow}{\mathscr{D}})^{\circ}=\left(\mathscr{L}^{\circ}\right)$. Moreover, G- $\mathscr{P r O f}=\mathrm{G}-\operatorname{Set} f$ by 2.8.6. As a result, S is an anti-equivalence of categories.

### 5.9.5 Dictionary

Let A be a separable algebra over K .
(a) If $A=A_{1} \times A_{2}$, then $S(A)=S\left(A_{1}\right) \sqcup S\left(A_{2}\right)$. If $A=A_{1} \otimes A_{2}$, then $S(A)=$ $S\left(A_{1}\right) \times S\left(A_{2}\right)$. Indeed, the functor $S$ transforms sums into products and conversely.
(b) Let A and B be separable algebras. A homomorphism $f: \mathrm{A} \rightarrow \mathrm{B}$ is injective if and only if $f^{*}: \mathrm{S}(\mathrm{B}) \rightarrow \mathrm{S}(\mathrm{A})$ is surjective. Indeed consider the fibre product $A^{\prime}=A \times_{B} A$. The G-set $S\left(A^{\prime}\right)$ is sum of two copies of $S(A)$ amalgamated by $S(B)$, i.e. the quotient of $S(A) \sqcup S(A)=\{1,2\} \times S(A)$ by the equivalence relation identifying $(1, x)$ and $(2, x)$ for $x \in f^{*}(\mathrm{~S}(\mathrm{~B}))$, (it is the fibre product in the opposite category). Then $f$ is injective if and only if the two projections of $\mathrm{A}^{\prime}$ onto A agree, or equivalently the two injections of $\mathrm{S}(\mathrm{A})$ into $\mathrm{S}\left(\mathrm{A}^{\prime}\right)$ agree, which is the case if and only if $f^{*}$ is surjective.
(c) Similarly, $f$ is surjective if and only if $f^{*}$ is injective. A direct consequence of the definition of S is that if $f$ is surjective, then $f^{*}$ is injective. If $f$ is not surjective,
then $f$ factorizes in $\mathrm{A} \xrightarrow{\tilde{f}} \mathrm{~B}^{\prime} \xrightarrow{\iota} \mathrm{B}$, where $\iota$ is injective but not an isomorphism. Then $f^{*}=\tilde{f}^{*} \circ \iota^{*}$, and $\iota^{*}$ is surjective but not an isomorphism; hence $\iota^{*}$ is not injective since $\mathrm{S}(\mathrm{B})$ is compact, and $f^{*}$ is not injective.
(d) $A$ is an extension of $K$ if and only if $G$ acts transitively on $S(A)$. Indeed, $S(\Omega)=G$ and A is an extension of K if and only if there is an injective homomorphism from A to $\Omega$.
More generally, the orbits of the action of G on $\mathrm{S}(\mathrm{A})$ are in bijective correspondence with the maximal ideals of A.
(e) Assigning to each sub-extension $L$ of $\Omega_{\text {sep }}$ the stabilizer of the canonical injection $\iota_{\mathrm{L}} \in \mathrm{S}(\mathrm{L})$, i.e. $\operatorname{Aut}_{\mathrm{L}}\left(\Omega_{\text {sep }}\right)$, defines a bijection between the sub-extensions of $\Omega_{\text {sep }}$ and the closed subgroups of G. The inverse is given by $\mathrm{H} \mapsto \operatorname{Fix}_{\mathrm{H}}\left(\Omega_{\text {sep }}\right)$.
(f) The sub-extensions L and $\mathrm{L}^{\prime}$ of $\Omega_{\text {sep }}$ are isomorphic if and only if the sub-groups $\operatorname{Aut}_{\mathrm{L}}\left(\Omega_{\mathrm{sep}}\right)$ and $\mathrm{Aut}_{\mathrm{L}^{\prime}}\left(\Omega_{\mathrm{sep}}\right)$ of G are conjugate.

### 5.9.6

Proposition and Definition Let L be a sub-extension of $\Omega$. The following conditions are equivalent:
(i) L is separable and invariant under G ;
(ii) L is the union of finite Galois sub-extensions of $\Omega_{\text {sep }}$;
(iii) L is contained in $\Omega_{\text {sep }}$ and the sub-group $\operatorname{Aut}_{\mathrm{L}}(\Omega)$ of G is normal.

If these conditions hold, then L is said to be Galois.
Proof (ii) $\Rightarrow$ (i). Immediate from Theorem 5.6.3.
(i) $\Rightarrow$ (ii). Let $x \in \mathrm{~L}$. The extension generated by the conjugates of $x$ is etale and invariant under G, hence is Galois (5.6.3, Theorem).
(i) $\Rightarrow$ (iii). For $g \in \mathrm{G}, \operatorname{Aut}_{g(\mathrm{~L})} \Omega=g \cdot \mathrm{Aut}_{\mathrm{L}} \Omega \cdot g^{-1}$. If L is preserved by G, then $g \cdot$ Aut $_{\mathrm{L}} \Omega \cdot g^{-1}=\mathrm{Aut}_{\mathrm{L}} \Omega$ for all $g \in \mathrm{G}$.
(iii) $\Rightarrow$ (i). The same formula shows that $\mathrm{Aut}_{\mathrm{L}} \Omega=\operatorname{Aut}_{g(\mathrm{~L})} \Omega$. Hence, by (5.6.4, Corollary 2 ), $\mathrm{L}=g(\mathrm{~L})$.

## Reference

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## Chapter 6 Riemann Surfaces

## Introduction

We begin this chapter with the definition of a Riemann surface. Let B be a compact connected Riemann surface, and $\mathscr{M}(\mathrm{B})$ the field of meromorphic functions on B. If X is a compact connected Riemann surface over B (i.e. equipped with a non constant morphism $\pi: \mathrm{X} \rightarrow \mathrm{B})$, then the field $\mathscr{M}(\mathrm{X})$ is a finite extension of $\mathscr{M}(\mathrm{B})$. Moreover, there is a finite subset $\Delta$ of $B$ such that $X^{\prime}=X-\pi^{-1}(\Delta)$ is a connected finite cover of $\mathrm{B}-\Delta$. The functors $\mathrm{X} \mapsto \mathscr{M}(\mathrm{X})$ and $\mathrm{X} \mapsto \mathrm{X}^{\prime}$ give an equivalence between the category $\mathscr{V}_{\mathrm{B}}^{1}$ of compact connected Riemann surfaces over B , and respectively, the opposite category of the category of finite extensions of $\mathscr{M}(\mathrm{B})$, and the direct limit category of categories of connected covers of $B-\Delta$, where $\Delta \subset B$ is finite.

The study of Riemann surfaces can thus be undertaken either in a purely algebraic framework or a purely topological one. In each, the category $\mathscr{V}_{\mathrm{B}}^{1}$ can be identified (Chap. 4 for the topological framework, Chap. 5 for the algebraic one) with the category of finite sets acted upon by a profinite group $G$ which may be called G- $\mathscr{A l} \mathscr{Q}$ or G- $\mathscr{T} \boldsymbol{\not p}$. These two groups may be identified, and thus G- $\mathscr{A l} \mathscr{q}$ can be described topologically. In Section 6.4, we give a description of the Galois group of the algebraic closure of $\mathbb{C}(Z)=\mathscr{M}(\Sigma)$ (where $\Sigma$ is the Riemann sphere), which cannot be done purely algebraically.

In Section 6.7, we study the automorphism group of a Riemann surface, and assuming it to be finite we give an upper bound for its order when its genus is $\geqslant 2$. For this homology theory is required. We have chosen the theory obtained from triangulation. These tools are developed in Sections 6.5 and 6.6. In Section 6.8, we show that the bound obtained for the order of the automorphism group is not reached in genus 2, but is in genus 3. In Section 6.10 we show the existence of a Riemann surface of genus 3 whose automorphism group has order 168 . We construct this latter example in the framework of Poincaré geometry, which is introduced in Section 6.9.

### 6.1 Riemann Surfaces, Ramified Coverings

### 6.1.1 Riemann Surfaces

A Hausdorff topological space $X$ is said to be a topological surface if every point of $X$ has an open neighbourhood homeomorphic to some open subset of $\mathbb{R}^{2}$.

Let X be a topological surface. A complex chart for X is a homeomorphism $\varphi$ from an open subset $U$ of $X$ onto an open subset of $\mathbb{C} ; U$ is said to be the domain of the chart $\varphi$, and $\varphi$ a chart centered at $x_{0}$ if $\varphi\left(x_{0}\right)=0$ and if the image of $\varphi$ is a disk centered at 0 . Let $\varphi$ and $\psi$ be two charts for X with respective domains U and V . The map $\gamma: x \mapsto \psi\left(\varphi^{-1}(x)\right)$ is a homeomorphism from the open subset $\varphi(\mathrm{U} \cap \mathrm{V})$ of $\mathbb{C}$ onto the open subset $\psi(\mathrm{U} \cap \mathrm{V})$ of $\mathbb{C}$, and is called the transition map from $\varphi$ to $\psi$. A $\mathbb{C}$-analytic atlas on X is a family of charts whose domains cover X such that all transition maps are holomorphic, i.e. $\mathbb{C}$-analytic.

Two $\mathbb{C}$-analytic atlases $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ and $\left(\psi_{j}\right)_{j \in \mathrm{~J}}$ are equivalent if for all $(i, j)$, the transition maps from $\varphi_{i}$ to $\psi_{j}$ and from $\psi_{j}$ to $\varphi_{i}$ are holomorphic, i.e. if the family obtained by their union is also an atlas. This is an equivalence relation. A $\mathbb{C}$-analytic structure on X is an equivalence class of analytic atlases. A topological surface with an analytic structure is called a Riemann surface, or a $\mathbf{1}$-dimensional $\mathbb{C}$-analytic manifold, or a complex curve.

Let X be a Riemann surface whose structure is defined by an atlas $\left(\varphi_{i}\right)$, and $\mathrm{U}_{i}$ the domain of $\varphi_{i}$. Let Y be an open subset of X . The restrictions of $\varphi_{i}$ to $\mathrm{U}_{i} \cap \mathrm{Y}$ form an atlas for Y. The structure of Y defined by this atlas only depends on that of X . The structure on Y is said to be induced by that of X .

### 6.1.2 Example. The Riemann Sphere

Consider the set $\Sigma=\mathbb{C} \cup\{\infty\}$ and equip $\Sigma$ with the topology inducing the ordinary topology on $\mathbb{C}$, and for which the open neighbourhoods of $\infty$ are complements of the compact subsets of $\mathbb{C}$.

The space $\Sigma$ is homeomorphic to the sphere

$$
\mathbf{S}^{2}=\left\{(x, y, z) \in \mathbb{R}^{3} \mid x^{2}+y^{2}+z^{2}=1\right\} .
$$

Indeed, as can be checked, the stereographic projection

$$
(x, y, z) \mapsto \frac{x+i y}{1-z}, \quad(0,0,1) \mapsto \infty
$$

is a homeomorphism from $\mathrm{S}^{2}$ onto $\Sigma$.
The charts $\varphi_{0}=1_{\mathbb{C}}: \Sigma-\{\infty\} \rightarrow \mathbb{C}$ and $\varphi_{1}: \Sigma-\{0\} \rightarrow \mathbb{C}$, defined by $\varphi_{1}(z)=$ $\frac{1}{z}$ and $\varphi_{1}(\infty)=0$, form an analytic atlas on $\Sigma$, the transition map being $z \mapsto \frac{1}{z}$ from
$\mathbb{C}-\{0\}$ to itself. This atlas defines a $\mathbb{C}$-analytic structure on $\Sigma$. With this structure, $\Sigma$ is a compact Riemann surface called the Riemann sphere.

### 6.1.3 Analytic Maps

Let X be a Riemann surface and $\varphi$ a chart for X with domain U . If $f: \mathrm{X} \rightarrow \mathbb{C}$ is a function, the function $f \circ\left(\varphi^{-1}\right)$ defined on $\varphi(\mathrm{U})$ is called an expression of $f$ in the chart $\varphi$, and holomorphic if its expression in every chart in an atlas defining the structure de X is holomorphic, this property being independent of the choice of atlas. If Y is a Riemann surface, $\psi$ a chart with domain V and $f: \mathrm{X} \rightarrow \mathrm{Y}$ is a continuous map, the map $\psi \circ f \circ \varphi^{-1}$ defined on $\varphi\left(\mathrm{U} \cap f^{-1}(\mathrm{~V})\right)$ is called an expression of $f$ in the above charts, and analytic if there exist atlases $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ and $\left(\psi_{j}\right)_{j \in \mathrm{~J}}$ of X and Y such that its expression in the charts $\varphi_{i}$ and $\psi_{j}$ is holomorphic for all $(i, j) \in \mathrm{I} \times \mathrm{J}$. This property is independent of the equivalence class of the given atlases.

If $\mathrm{D} \subset \mathrm{X}$ is a discrete closed set, then a holomorphic function $f: \mathrm{X}-\mathrm{D} \rightarrow \mathbb{C}$ is called meromorphic on X with poles in D if for all $a \in \mathrm{D}$, there is a neighbourhood of $a$ in V and two holomorphic functions $g, h: \mathrm{V} \rightarrow \mathbb{C}$ such that $h^{-1}(0) \subset\{a\}$ and $f=\frac{g}{h}$ on $\mathrm{V}-\{a\}$ (then $f:=\frac{g}{h}$ on V ). The ring of meromorphic functions on X with poles in D is written $\mathscr{M}(\mathrm{X}, \mathrm{D})$. let

$$
\mathscr{M}(\mathrm{X})=\underset{\longrightarrow}{\lim } \mathscr{M}(\mathrm{X}, \mathrm{D})
$$

be the direct limit over the set of discrete closed subsets of X ; the elements of $\mathscr{M}(\mathrm{X})$ are called meromorphic functions on X .

### 6.1.4

Proposition and Definition Let X and Y be Riemann surfaces, $f: \mathrm{X} \rightarrow \mathrm{Y}$ an analytic map, $x_{0} \in \mathrm{X}$, and $y_{0}=f\left(x_{0}\right)$. Assume that $f$ is not constant in the neighbourhood of $x_{0}$. Let $\psi$ be a chart for Y centered at $y_{0}$. Then there is a chart $\varphi$ for X centered at $x_{0}$ such that the expression of $f$ in the charts $\varphi$ and $\psi$ is $z \mapsto z^{d}$, for some number $d$ independent of the choice of charts. It is called the ramification index of $f$ at $x_{0}$.

Proof Let $\varphi_{0}$ be an analytic chart for X in the neighbourhood of $x_{0}$ such that $\varphi_{0}\left(x_{0}\right)=$ 0 , and $\dot{f}$ the expression of $f$ in the charts $\varphi_{0}$ and $\psi$. Then $f(0)=0$. We give a series expansion of $\dot{f}$ in the neighbourhood of 0 . Let $\dot{f}(z)=\sum a_{k} z^{k}, d$ the smallest integer for which $a_{d} \neq 0$, and $c \in \mathbb{C}$ such that $c^{d}=a_{d}$. Then $\dot{f}(z)=(c z)^{d}(1+u(z))$ for some holomorphic $u$ such that $u(0)=0$. Since there is a holomorphic function $w \mapsto \sqrt[d]{w}$ in the neighbourhood of 1 , the function $1+u$ may be written as $h^{d}$ for some
holomorphic $h$ in the neighbourhood of 0 with $h(0)=1$. Hence $\dot{f}(z)=(c z \cdot h(z))^{d}$. Set $\varphi(x)=c \varphi_{0}(x) \cdot h\left(\varphi_{0}(x)\right)$. The implicit function theorem implies that $\varphi$ is an analytic chart in the neighbourhood of $x_{0}$, and $\psi(f(x))=(\varphi(x))^{d}$ for $x$ near $x_{0}$.

If $z \mapsto z^{d}$ is the expression of $f$, then there is a neighbourhood $\mathrm{U}_{0}$ of $x_{0}$ such that for any neighbourhood U of $x_{0}$ in $\mathrm{U}_{0}$, there is a neighbourhood V of $y_{0}$ such that for all $y \neq y_{0} \in \mathrm{~V}, d$ is the cardinality of $f^{-1}(y) \cap \mathrm{U}$. This gives a characterization of $d$ independent of the charts.

### 6.1.5

## Proposition Let


be a commutative diagram, where $\mathrm{X}, \mathrm{X}^{\prime}, \mathrm{B}$ and $\mathrm{B}^{\prime}$ are Riemann surfaces, $\alpha, \beta$, $\varphi$ analytic maps, and $\psi$ a continuous map. Assume that $\beta$ is not constant on any connected component of $\mathrm{X}^{\prime}$. Then $\psi$ is analytic.

Proof Since this is a local result, $\mathrm{X}^{\prime}$ and $\mathrm{B}^{\prime}$ may be assumed to be open in $\mathbb{C}$, $\beta: z \mapsto z^{d}$, and X connected. The function $\psi^{d}=\beta \circ \psi=\varphi \circ \alpha$ is analytic; so the set $\Delta=\psi^{-1}(0)=\left\{x \mid(\psi(x))^{d}=0\right\}$ is the whole of X , in which case $\psi$ is either constant, or a closed discrete set. Then $\psi$ is continuous on X , and analytic on $\mathrm{X}-\Delta$, so also on X by a well-known result on holomorphic functions. ${ }^{1}$

### 6.1.6

Proposition Let B be a connected Riemann surface, X a Riemann surface equipped with a proper analytic map $\pi: \mathrm{X} \rightarrow \mathrm{B}$. The following conditions are equivalent:
(i) there is no connected component of X on which $\pi$ is constant;
(ii) the fibres de $\pi$ are finite;
(iii) for all connected components $\mathrm{X}_{i}$ of $\mathrm{X}, \pi\left(\mathrm{X}_{i}\right)=\mathrm{B}$;
(iv) $\pi$ is open;
(v) for all $x \in \mathrm{X}$, there is a chart $\varphi$ for X in the neighbourhood of $x$ with $\varphi(x)=0$ and a chart $\psi$ for B with $\psi(\pi(x))=0$ such that the expression of $\pi$ in these charts takes the form $z \mapsto z^{d}$.

[^19]Proof (i) $\Rightarrow$ (v) follows from Proposition 6.1.4 by noting that, according to the analytic continuation principle, if $\pi$ is constant in the neighbourhood of $x$, then it also is on the connected component of $x$.
(v) $\Rightarrow$ (iv) since the map $z \mapsto z^{d}$ is open.
(iv) $\Rightarrow$ (iii): let $X_{i}$ be a connected component of $X$. Then $X_{i}$ is clopen and nonempty. Hence $\pi\left(\mathrm{X}_{i}\right)$ is clopen since $\pi$ is proper and nonempty; so $\pi\left(\mathrm{X}_{i}\right)=\mathrm{B}$.
(iii) $\Rightarrow$ (i) is obvious.
(v) $\Rightarrow$ (ii): each point of $X$ is isolated in its fibre, hence the fibres are discrete. As they are compact, they are finite.
(ii) $\Rightarrow$ (i) is obvious.

### 6.1.7 Ramified Coverings

Definition Let $B$ be a topological surface. A finite ramified covering $B$ is a topological surface X equipped with a proper continuous map $\pi: \mathrm{X} \rightarrow \mathrm{B}$ satisfying the following condition:
(RR) For all $x \in \mathrm{X}$, there is a complex chart for X centered at $x$ and a complex chart for B centered at $\pi(x)$ such that the expression of $\pi$ in these charts takes the form $z \mapsto z^{d}$.

Examples (1) Every finite covering of $B$ is a finite ramified covering of $B$. However a finite ramified covering $B$ is not in general a covering of $B$.
(2) Let B be a Riemann surface and X a Riemann surface equipped with a proper analytic map $\pi: \mathrm{X} \rightarrow \mathrm{B}$. Assume that there is no connected component of X on which $\pi$ is constant. Then X is a finite ramified covering of B and is called a finite ramified analytic covering of $B$.

Remarks (1) If X is a finite ramified covering of B , for any point $x$ the integer $d$ in condition (RR) is independent of the choice of charts, and is called the local degree or the ramification index of X at $x$. The set of points for which $d>1$ is discrete and closed in X , its projection $\Delta$ in B is discrete and closed in B and is called the ramification set and $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ is a finite covering of $\mathrm{B}-\Delta$.
(2) If X is a finite ramified analytic covering of B , then the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is open with finite fibres.
(3) If B is connected then for all closed discrete $\Delta \subset B, B-\Delta$ is connected. Indeed, if $h$ is a locally constant function on $\mathrm{B}-\Delta$, it extends to a locally constant function $\bar{h}$ on B ; then $\bar{h}$ is constant and so is $h$. Hence the degree of a finite ramified covering over a connected component of B is well defined.

### 6.1.8

Proposition Let B be a topological surface, $(\mathrm{X}, \pi)$ a ramified covering of $\mathrm{B}, \Delta \subset \mathrm{B}$ the ramification set of $\pi, b \in \mathrm{~B}, \varphi$ a chart for B centered at $b$ with domain U , and V a connected component of $\pi^{-1}(\mathrm{U})$. Set $\mathrm{U}^{*}=\mathrm{U}-\{b\}$. Suppose that $\mathrm{U}^{*} \cap \Delta=\varnothing$. Then the fibre $\mathrm{V}(b)=\mathrm{V} \cap \pi^{-1}(b)$ of V at $b$ only contains $x$, and there is a chart $\psi$ of X centered at $x$ with domain V such that the expression of $\pi$ in the charts $\psi, \varphi$ is $z \mapsto z^{d}$, where $d$ is the local degree of $\pi$ at $x$.

Lemma Any two connected d-fold coverings of $\mathrm{U}^{*}$ are isomorphic.
Proof of the Lemma The chart $\varphi$ induces a homeomorphism from $\mathrm{U}^{*}$ to a pointed disk $\mathrm{D}^{*}=\mathrm{D}-\{0\}$. For $\beta \in \mathrm{U}^{*}$, the fundamental group $\pi_{1}\left(\mathrm{U}^{*}, \beta\right)$ is isomorphic to $\mathbb{Z}$. By (4.6.9) the isomorphism classes of connected $d$-fold coverings of $\mathrm{U}^{*}$ are in bijective correspondence with (conjugation classes of) subgroups of index $d$ in $\mathbb{Z}$, and it is unique.

Proof of the Proposition The set $\mathrm{V}(b)$ is finite, and so $\mathrm{V}^{*}=\mathrm{V}-\mathrm{V}(b)$ is connected; equipped with the restriction of $\pi$ it is a connected component of $\mathrm{U}^{*}$ of some degree $d$. The image of $\varphi$ is a disk $\mathrm{D}=\mathrm{D}_{r}$. Set $r^{\prime}=r^{1 / d}, \mathrm{D}^{\prime}=\mathrm{D}_{r^{\prime}}$. The map $p: z \mapsto z^{d}$ turns $\mathrm{D}^{* *}$ into a connected component of $\mathrm{D}^{*}$ of degree $d$. By the lemma, the coverings $\mathrm{V}^{*}$ and $\varphi^{*}\left(\mathrm{D}^{\prime *}, p\right)$ of $\mathrm{U}^{*}$ are isomorphic. Thus there is a homeomorphism $\psi: \mathrm{V}^{*} \rightarrow \mathrm{D}^{*}$ such that $p \circ \psi=\varphi \circ p: \mathrm{V}^{*} \rightarrow \mathrm{D}^{*}$.

We show that $\mathrm{V}(b)$ is a singleton. Let $x_{i}$ be the points of $\mathrm{V}(b)$ with $i \in \mathrm{I}$, where I is finite. The map $\pi: \mathrm{V} \rightarrow \mathrm{U}$ is proper, and so there is a neighbourhood $\mathrm{U}^{\prime}$ of $b$ in U and a partition $\left(\mathrm{V}_{i}^{\prime}\right)_{i \in \mathrm{I}}$ of $\mathrm{V}^{\prime}=\mathrm{V} \cap \pi^{-1}\left(\mathrm{U}^{\prime}\right)$ into nonempty open sets such that $x_{i} \in \mathrm{~V}_{i}^{\prime}$ for all $i$. And we may assume that there is some $\rho>0$ such that $\mathrm{U}^{\prime}=\varphi^{-1}\left(\mathrm{D}_{\rho}\right)$. Then $\psi\left(\mathrm{V}_{i}^{\prime} \cap \mathrm{V}^{*}\right)$ form a partition of $\mathrm{D}_{\rho^{\prime}}^{*}$ into nonempty open sets, where $\rho^{\prime}=\rho^{1 / d}$. Hence, Card $\mathrm{V}(b)=1$.

Let $\mathrm{V}(b)=\{x\}$ and extend $\psi$ to a bijection from V onto $\mathrm{D}^{\prime}$, also denoted $\psi$, by setting $\psi(x)=0$. We show that it is a homeomorphism. If $x^{\prime} \in \mathrm{V}^{*}$ tends to $x$, then $\pi\left(x^{\prime}\right)$ tends to $b$, and $\varphi\left(\pi\left(x^{\prime}\right)\right)=\left(\psi\left(x^{\prime}\right)\right)^{d}$ tends to 0 ; so $\psi\left(x^{\prime}\right)$ tends to 0 . Hence $\psi$ is continuous. If $\psi\left(x^{\prime}\right)$ tends to 0 , then $\pi\left(x^{\prime}\right)=\varphi^{-1}\left(\left(\psi\left(x^{\prime}\right)\right)^{d}\right)$ tends to $b$, and so $x^{\prime}$ tends to $x$ since $\pi ; \mathrm{V} \rightarrow \mathrm{U}$ is proper and $\mathrm{V}(b)=\{x\}$. Hence $\psi^{-1}$ is continuous.

Then $\psi: \mathrm{V} \rightarrow \mathrm{D}^{\prime}$ is a chart for X centered at $x$, the expression of $\pi$ in this chart is $z \mapsto z^{d}$, and $d$ is necessarily the local degree of $\pi$ at $x$.

### 6.1.9

Let B be a topological surface, $a, b \in \mathrm{~B}$ with $a \neq b$, and $\gamma \in \pi_{1}(\mathrm{~B}-\{a\}, b)$, identified with the Poincaré group (4.9.9). The element $\gamma$ will be said to go round $a$ in B if there is a chart $\varphi: \mathrm{U} \rightarrow \mathbb{D}$ for B centered at $a$, a point $b^{\prime} \in \mathrm{U}$ and a class $\beta$ of paths from $b$ to $b^{\prime}$ such that $\gamma=\beta^{-1} \cdot \alpha \cdot \beta$, where $\alpha$ is the class of the loop in $\left(\mathrm{U}, b^{\prime}\right)$ whose image under $\varphi$ is $t \mapsto \varphi\left(b^{\prime}\right) \cdot e^{2 \pi i t}$.

Proposition Let B be a topological surface and X a finite ramified covering of $\mathrm{B} ; \Delta$ the ramification set. Let $b \in \mathrm{~B}-\Delta, a \in \Delta$ and $\gamma \in \pi_{1}(\mathrm{~B}-\Delta, b)$ an element going round $a$ in $\mathrm{B}-(\Delta-\{a\})$. Set $\sigma_{\gamma}$ be the permutation of $\mathrm{X}(b)$ defined by $\gamma$. Then the order of $\sigma_{\gamma}$ in the permutation group of $\mathrm{X}(b)$ is 1.c.m. $\left(r_{1}, \ldots, r_{k}\right)$, where $r_{i}$ are the ramification indices of X at points of $\mathrm{X}(a)$.

Proof The class of paths $\beta$ enables the identification of $\mathrm{X}(b)$ with $\mathrm{X}\left(b^{\prime}\right)$ and of $\pi_{1}(\mathrm{X}-\Delta, b)$ with $\pi_{1}\left(\mathrm{X}-\Delta, b^{\prime}\right)$; moreover these identifications are compatible with all actions. We may assume that $b=b^{\prime}, \beta=e_{b}$ and $\gamma=\alpha$.

Let $\mathrm{X}(a)=\left\{x_{1}, \ldots, x_{k}\right\}$. By Proposition 6.1.8, there are charts $\psi_{1}, \ldots, \psi_{k}$ de X with domains $\mathrm{V}_{1}, \ldots, \mathrm{~V}_{k}$, centered respectively at $x_{1}, \ldots, x_{k}$, such that the expression of the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ in the charts $\psi_{i}$ and $\varphi$ is $z \mapsto z^{r_{i}}$. The loop with expression $t \mapsto \varphi(b) \cdot e^{2 \pi i t}$ can be lifted in $\mathrm{V}_{i}$ to paths with expressions $t \mapsto \lambda_{i} e^{2 \pi i \frac{t}{r_{i}}}$, where $\lambda_{i}^{r_{i}}=\varphi(b)$. Hence the permutation $\sigma_{\alpha}$ preserves each set $\mathrm{V}_{i}(b)$, and its expression in $\mathrm{V}_{i}(b)$, identified by $\psi_{i}$ with the set of $r_{i}$-th roots of $\varphi(b)$, is multiplication by $e^{2 \pi i \frac{1}{r_{i}}}$. Hence $\left(\sigma_{\alpha}\right)^{\nu}=1$ if and only if $\nu$ is a common multiple to $r_{i}$.

### 6.1.10

Proposition Let X be a finite ramified covering of B . For any $\mathbb{C}$-analytic structure on B , there is a unique $\mathbb{C}$-analytic structure on X such that $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is analytic.

Proof Let $\delta$ be the ramification set of $\pi$. Let $\left(\varphi_{i}\right)_{i \in \mathrm{I}}$ be an atlas for B consisting of centered charts such that, for all $i, \varphi_{i}$ is a chart with domain $\mathrm{U}_{i}$ centered at $b_{i}$ with $\left(\mathrm{U}_{i}-\left\{b_{i}\right\}\right) \cap \Delta=\varnothing$. For each $i$, let $\mathrm{V}_{i, j}, j \in \mathrm{~J}_{i}$ be the connected components of $\pi^{-1}\left(\mathrm{U}_{i}\right)$, and $\mathrm{J}=\left\{(i, j) \mid j \in \mathbf{J}_{i}\right\}$. For each $(i, j) \in \mathbf{J}$, let $\psi_{i, j}$ be a chart of $\mathbf{X}$ such that the expression of $\pi$ in the charts $\psi_{i, j}$ and $\varphi_{i}$ is of the form $z \mapsto z^{d_{i, j}}$ (6.1.8, Proposition). The maps $\psi_{i, j}$ form an atlas for $\mathbf{X}$, and as a result of (6.1.5) the transition maps are analytic. This atlas defines an analytic structure on X with the desired property. Uniqueness follows from (6.1.5).

Corollary Let B be a Riemann surface. The forgetful functor identifies the category $\mathscr{V}_{\mathrm{B}}$ of finite analytically ramified covering over B with the category $\mathscr{R} \mathscr{R}$ of finite topological ramified coverings of B and continuous maps over B .

This is a restatement of the previous proposition for objets and of Proposition 6.1.5 for morphisms.

## 6.1 .11

Proposition Let B be a topological surface and $\Delta$ a closed discrete set in B . Then the functor $\rho:\left.\mathrm{X} \mapsto \mathrm{X}\right|_{\mathrm{B}-\Delta}$ from the category $\mathscr{R}_{\Delta}$ of finite ramified coverings
of B whose ramification set is contained in $\Delta$ to the category $\mathscr{C o v f}_{\mathrm{B}-\Delta}$ of finite coverings of $\mathrm{B}-\Delta$ is an equivalence of categories.

A quasi-inverse functor is given by the "end compactification" (see 6.1, Exercise 2).

Proof (a) We show that $\rho$ is essentially surjective. Suppose first that B is a disk D and $\Delta=\{0\}$. As $\rho$ commutes with finite sums, it suffices to show that every connected finite covering of $\mathrm{B}-\Delta=\mathrm{D}-\{0\}$ is the restriction of a ramified covering of D . Now, any two finite sets with the same cardinality on which $\mathbb{Z}$ acts transitively are isomorphic as $\mathbb{Z}$-sets. Hence by 4.6 .8 , every connected $d$-fold covering of $\mathrm{D}-\{0\}$ is isomorphic to ( $\widetilde{\mathrm{D}}-\{0\}, z \mapsto z^{d}$ ), where $\widetilde{\mathrm{D}}$ is the disk with radius $\sqrt[d]{r}, r$ being the radius of $D$, and this covering is the restriction of the ramified covering $\widetilde{D}$ of $D$.


$$
D-\{0\}
$$



Finally, in this case, a covering $X$ of $D-\{0\}$ is the restriction of a ramified covering Y of D obtained by adding a point over 0 for each connected component of X .
(b) In the general case, let X be a finite covering of $\mathrm{B}-\Delta$. Each $b \in \Delta$ has a neighbourhood U such that $\left.\mathrm{X}\right|_{\mathrm{U}-b}$ is the restriction of a finite ramified covering $\mathrm{Y}_{b}$ of U . Then $\mathrm{Y}_{b}=\left.\mathrm{X}\right|_{\mathrm{U}-b} \cup \mathrm{~F}_{b}$ for $\mathrm{F}_{b}=\mathrm{Y}_{b}(b)$. By the gluing lemma for topologies (4.4.6), the set $\mathrm{Y}=\mathrm{X} \sqcup \bigsqcup_{b \in \Delta} \mathrm{~F}_{b}$ can be equipped with a topology inducing that of X and that of all $\mathrm{Y}_{b}$. Then Y become a ramified covering of B whose restriction to $B-\Delta$ is $X$.
(c) $\rho$ is fully faithful. Let X and Y be objects of $\mathscr{R}_{\Delta}$ and $f$ a covering morphism from $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ to $\left.\mathrm{Y}\right|_{\mathrm{B}-\Delta}$. Let $x_{0} \in \mathrm{X}$ be such that $b_{0}=\pi\left(x_{0}\right) \in \Delta$, and $\mathrm{Y}\left(b_{0}\right)=\left\{y_{1}, \ldots, y_{k}\right\}$. Let $\mathrm{W}_{1}, \ldots, \mathrm{~W}_{k}$ be mutually disjoint respective neighbourhoods of $y_{1}, \ldots, y_{k}$ in Y . As Y is proper over B , there is an open neighbourhood U of $b_{0}$ in B such that $\left.\mathrm{Y}\right|_{\mathrm{U}} \subseteq \mathrm{W}_{1} \cup \ldots \cup \mathrm{~W}_{k}$. Let V be a neighbourhood of $x_{0}$ in X homeomorphic to a disc, such that $\mathrm{V} \cap \pi^{-1}(\Delta)=\left\{x_{0}\right\}$ and $\pi(\mathrm{V}) \subset \mathrm{U}$. Then $\mathrm{V}-\left\{x_{0}\right\}$ is connected and in $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$. The set $f\left(\mathrm{~V}-\left\{x_{0}\right\}\right)$ is connected and in $\mathrm{W}_{1} \cup \ldots \cup \mathrm{~W}_{k}$, hence in some $\mathrm{W}_{i}, \mathrm{~W}_{i_{0}}$ say. For any neighbourhood $\mathrm{U}^{\prime}$ of $b_{0}$ in B the set $f^{-1}\left(\left.\mathrm{~W}_{i_{0}}\right|_{U^{\prime}}\right)=\left.\mathrm{V}\right|_{\mathrm{U}^{\prime}}-\left\{x_{0}\right\}$ is the trace over $\mathrm{X}-\left\{x_{0}\right\}$ of a neighbourhood of $x_{0}$. However, as $\mathrm{U}^{\prime}$ varies, $\left.\mathrm{W}_{i_{0}}\right|_{\mathrm{U}^{\prime}}$ form a fundamental system of neighbourhoods of $y_{i_{0}}$. Hence the map $f$, extended by $x_{0} \mapsto y_{i 0}$, is continuous at $x_{0}$. Extending likewise at each point of $\pi^{-1}(\Delta)$ gives a continuous continuation of $f$ to a morphism from X to Y . As $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ is dense in X , this continuation is unique.

### 6.1.12

Theorem Let B be a Riemann surface, I the ordered set of discrete closed subsets of B and $\mathscr{V}_{\mathrm{B}}$ the category of finite analytically ramified coverings of B . Define a functor $\omega: \mathscr{V}_{\mathrm{B}} \rightarrow \lim _{\Delta \in \mathrm{I}} \mathscr{E}$ ovf $\mathcal{B}_{\mathrm{B}-\Delta}$ by associating to each object $\mathrm{X} \in \mathscr{V}_{\mathrm{B}}$ the underlying space of $\mathrm{X}_{\mathrm{B}-\Delta} \rightarrow$, where $\Delta$ is the ramification set of X . This functor is an equivalence of categories.

Proof By Corollary 6.1.10, the category $\mathscr{V}_{\text {B }}$ can be identified with

$$
\mathscr{R} \mathscr{R}=\bigcup_{\Delta} \mathscr{R} \mathscr{R}_{\Delta}=\lim _{\longrightarrow} \mathscr{R}_{\Delta}
$$

The theorem then reduces to Proposition 6.1 .11 by passing to the direct limit.

## Exercises 6.1. (Riemann surfaces)

1.-Show that the following spaces are topological surfaces:
(a) The torus $\mathbb{T}^{2}=S^{1} \times S^{1}=\mathbb{R}^{2} / \mathbb{Z}^{2}$ (it can be embedded in $\mathbb{R}^{4}$ :

$$
(\alpha, \theta) \mapsto(\cos \alpha, \sin \alpha, \cos \theta, \sin \theta)
$$

or in $\mathbb{R}^{3}$ :

$$
(\alpha, \theta) \mapsto((\mathrm{R}+r \cos \theta) \cos \alpha,(\mathrm{R}+r \cos \theta) \sin \alpha, r \sin \theta)
$$

with $r<\mathrm{R}$ ).
(b) The open Möbius strip

$$
\left\{\left.\left(\left(1+u \cos \frac{t}{2}\right) \cos t,\left(1+u \cos \frac{t}{2}\right) \sin t, u \sin \frac{t}{2}\right) \right\rvert\, u \in\right]-1,1[, t \in \mathbb{R}\}
$$

(c) The projective plane $\mathrm{P}^{2} \mathbb{R}$, i.e. the quotient of $\mathrm{S}^{2}$ by the equivalence relation identifying two opposite points. If $a \in \mathrm{P}^{2} \mathbb{R}$, then the space $\mathrm{P}^{2} \mathbb{R}-\{a\}$ is homeomorphic to the open Möbius strip.
(d) The half cone $\left\{(t \cos \theta, t \sin \theta, t) \mid t \in \mathbb{R}_{+}, \theta \in \mathbb{R}\right\} \subset \mathbb{R}^{3}$.
(e) Plücker's conoid:

$$
\{(1+u) \cos \theta,(1-u) \sin \theta, u) \mid u \in[-1,+1], \theta \in \mathbb{R} .\} \subset \mathbb{R}^{3}
$$

(f) The cone in $\mathbb{R}^{4}$ over a knot in $\mathbb{R}^{3}$ :
$\left\{(s(\mathrm{R}+r \cos 3 \theta) \cos 2 \theta, s(\mathrm{R}+r \cos 3 \theta) \sin 2 \theta, s r \sin 3 \theta, s) \mid s \in \mathbb{R}_{+}, \theta \in \mathbb{R}\right\} \subset \mathbb{R}^{4}$,
for fixed $r$ and R with $r<\mathrm{R}$.
2. (End compactification)-For a space X , let $\pi_{0}(\mathrm{X})$ denote the set of its connected components. Let $\mathscr{C}$ be the category of locally connected, locally compact spaces with finitely many connected components, and proper maps.
(a) Let $X \in \mathscr{C}$. Show that the relatively compact open subsets $U$ such that $\pi_{0}(\mathrm{X}-\mathrm{U})$ is finite form a cofinal set in the set of relatively compact subsets of X.
(b) Let $\mathrm{B}(\mathrm{X})=\underset{\leftarrow}{\lim } \pi_{0}(\mathrm{X}-\mathrm{K})$ be the inverse limit over the relatively compact subsets of X (the elements of $\mathrm{B}(\mathrm{X})$ are the ends of X ). Define a topology on $\widehat{\mathrm{X}}=$ $X \sqcup B(X)$ which makes it compact. Show that $X \mapsto \widehat{X}$ is a functor of $\mathscr{C}$ in the category of compact spaces.
(c) Give an example of a connected topological surface X such that $\widehat{\mathrm{X}}$ is not a topological surface.
(d) With the notation of Proposition 6.1.11, show that if $B$ is compact, X can be identified with the end compactification of $\left.X\right|_{\mathrm{B}-\Delta}$.
3.-(a) (Maximum principle) Let X be a Riemann surface, $x \in \mathrm{X}$ and $f: \mathrm{X} \rightarrow \mathbb{C}$ a holomorphic function. Deduce from 6.1.4 that $|f|$ has a nonzero upper or lower bound only if $f$ is constant in the neighbourhood of $x$.
(b) (Yet another proof of d'Alembert's theorem.) Show that if a holomorphic function $f: \mathbb{C} \rightarrow \mathbb{C}$ is proper (i.e. $|f(z)|$ tends to $\infty$ as $|z|$ tends to $\infty$ ), then it vanishes at least at one point. In particular every non constant polynomial has at least one root in $\mathbb{C}$.
(c) Show that if a holomorphic function $f: \mathbb{C} \rightarrow \mathbb{C}$ is proper, then it is necessarily a polynomial.

### 6.2 Ramified Coverings and Etale Algebras

### 6.2.1 Separation Theorem

The next result will be assumed:
Theorem Let X be a compact Riemann surface, and $a, b \in \mathrm{X}$ with $a \neq b$. Then there is a meromorphic function $f$ on X defined at $a$ and $b$ such that $f(a) \neq f(b)$.

For a proof, see 6.2, Exercise 1, or R. Gunning, Lectures on Riemann Surfaces [3].
Remark The theorem still holds when X is not compact, but the proof is different. If X is not compact but is connected, then there is a better result: the holomorphic functions separate the points.

### 6.2.2 Analytic Criterion for Meromorphism

Proposition Let X be a Riemann surface, D a discrete closed subset of $\mathrm{X}, f$ : $\mathrm{X}-\mathrm{D} \rightarrow \mathbb{C}$ a holomorphic function, $a \in \mathrm{X}, \varphi$ a chart centered at a with domain U . Then $f$ is meromorphic at a if and only if there is an open neighbourhood $\mathrm{U}^{\prime}$ of $a$ in U and constants $c$ and $k$ such that,

$$
\forall x \in \mathrm{U}^{\prime}-\{a\},|f(x)| \leqslant \frac{c}{|\varphi(x)|^{k}}
$$

Proof Assume that the condition holds. Then the function $x \mapsto(\varphi(x))^{k} f(x)$ is holomorphic on $\mathrm{U}-\{a\}$, and bounded in the neighbourhood of $a$ and so extends to a holomorphic function on U .

If $f$ is meromorphic at $a$, then in the neighbourhood of $a, f(x)=\frac{g(x)}{h(x)}$, where $h(x)=\varphi(x)^{k} u(x)$, with $u$ is holomorphic and $u(a) \neq 0$. Then $u$ is invertible in the neighbourhood of $a$ and $\frac{g}{u}$ is bounded by a constant $c$ in a neighbourhood $\mathrm{U}^{\prime}$ of $a$. Then

$$
|f(x)| \leqslant \frac{c}{|\varphi(x)|^{k}} \text { for } u \in \mathrm{U}^{\prime}-\{a\}
$$

### 6.2.3 The Functor $\mathscr{M}$

If X is a Riemann surface, let $\mathscr{M}(\mathrm{X})$ be the $\mathbb{C}$-algebra of meromorphic functions on X . If X is connected, then $\mathscr{M}(\mathrm{X})$ is a field. Let X and Y be two Riemann surfaces and $f: \mathrm{X} \rightarrow \mathrm{Y}$ an analytic map which is not constant on any connected component
of X. Define $f^{*}: \mathscr{M}(\mathrm{Y}) \rightarrow \mathscr{M}(\mathrm{X})$ by $f^{*}(h)=h \circ f$. If D is the set of poles of $h$, the function $f^{*}(h)$ is defined on $\mathrm{X}-f^{-1}(\mathrm{D})$. Then $(g \circ f)^{*}=f^{*} \circ g^{*}$.

Theorem Let B be a compact connected Riemann surface and X an analytically ramified d-fold covering $d$ of B . Then $\mathscr{M}(\mathrm{X})$ is an etale algebra of degree $d$ over $M$ (B).

Proof (a) Let $f \in \mathscr{M}(\mathrm{X})$. We show that $f$ is algebraic of degree $\leqslant d$ over $\mathscr{M}(\mathrm{B})$. This will follow if there is a polynomial $\mathrm{P} \in \mathscr{M}(\mathrm{B})[\mathrm{Z}]$ such that $\mathrm{P}(f)=0$. Let $\Delta_{r}$ (resp. $\Delta_{p}$ ) be the ramification set of X (resp. the projection of the set of poles of $f$ ), and $\Delta=\Delta_{r} \cup \Delta_{p}$. Let $b \in \mathrm{~B}-\Delta, \mathrm{X}(b)=\left\{x_{1}, \ldots, x_{d}\right\}$ and $a_{i}(b)$ the values of the elementary symmetric functions $f\left(x_{1}\right), \ldots, f\left(x_{d}\right)$, i.e.,

$$
a_{1}(b)=\sum_{i} f\left(x_{i}\right), a_{2}(b)=\sum_{i<j} f\left(x_{i}\right) f\left(x_{j}\right), \ldots, a_{d}(b)=\prod_{i} f\left(x_{i}\right)
$$

We show that for all $i, a_{i} \in \mathscr{M}(\mathrm{~B})$. All $a_{i}$ are meromorphic on $\mathrm{B}-\Delta_{r}$. Indeed, let $b \in \mathrm{~B}-\Delta_{r}$ and U a neighbourhood of $b$ such that $\left.\mathrm{X}\right|_{\mathrm{U}}$ is a trivial covering of U , i.e. $\left.\mathrm{X}\right|_{\mathrm{U}}=\mathrm{X}_{1} \cup \ldots \cup \mathrm{X}_{d}, \pi$ inducing an isomorphism from each $\mathrm{X}_{i}$ onto U . Define the meromorphic functions $f_{i}$ on U by $\left.f_{i} \circ \pi\right|_{\mathrm{x}_{i}}=\left.f\right|_{\mathrm{X}_{i}}$. Then

$$
\left.a_{1}\right|_{\mathrm{U}}=f_{1}+\cdots+f_{d},\left.a_{2}\right|_{\mathrm{U}}=\sum_{i<j} f_{i} f_{j}, \ldots,\left.a_{d}\right|_{\mathrm{U}}=f_{1} \ldots f_{d}
$$

Hence all $a_{i}$ are meromorphic on U . Let $\varphi$ be a chart centered at $b \in \Delta_{r}$ with domain U . For $x_{i} \in \mathrm{X}(b)$, let $\psi_{i}$ a chart of X centered at $x_{i}$ with domain $\mathrm{V}_{i}$ and such that the expression of $\pi$ in the charts $\psi_{i}$ and $\varphi_{i}$ is of the form $z \mapsto z^{d_{i}}$; then by 6.2.2 there is an open neighbourhood $\mathrm{V}_{i}$ of $x_{i}, \mathrm{~V}_{i}^{\prime} \subset \mathrm{V}_{i}$, and constants $c_{i}$ and $k_{i}$ such that, $\forall x_{i}^{\prime} \in \mathrm{V}_{i}^{\prime}$,

$$
\left|f\left(x_{i}^{\prime}\right)\right| \leqslant \frac{c_{i}}{\left|\psi_{i}\left(x_{i}^{\prime}\right)\right|^{k_{i}}}=\frac{c_{i}}{\left|\varphi\left(b^{\prime}\right)\right|^{k_{i} / d_{i}}} \text { where } b^{\prime}=\pi\left(x_{i}^{\prime}\right)
$$

Set $c=\sup _{i} c_{i}$ and $k=\sup _{i} \frac{k_{i}}{d_{i}}$. Then $\left|a_{j}\left(b^{\prime}\right)\right| \leqslant\binom{ d}{j} \frac{c^{j}}{\left|\varphi\left(b^{\prime}\right)\right|^{j k}}$ and $a_{i} \in \mathscr{M}(\mathrm{~B})$.
Set $\mathrm{P}(\mathrm{Z})=\mathrm{Z}^{d}+\sum_{i}(-1)^{i} a_{i} \mathrm{Z}^{d-i} \in \mathscr{M}(\mathrm{~B})[\mathrm{Z}]$. We show that $\mathrm{P}(f) \in \mathscr{M}(\mathrm{X})$ is null. Let $b \in \mathrm{~B}-\Delta$ and $x \in \mathrm{X}(b)$. The polynomial $\mathrm{P}_{d}(\mathrm{Z})=\mathrm{Z}^{d}+\sum_{i=1}^{d}(-1)^{i} a_{i}$ (b) $\mathrm{Z}^{d-i}$ has roots $f\left(x_{1}\right), \ldots, f\left(x_{d}\right)$, where $\mathrm{X}(b)=\left\{x_{1}, \ldots, x_{d}\right\}$. In particular,

$$
f(x)^{d}+\sum_{i=1}^{d}(-1)^{i} a_{i}(b) f(x)^{d-i}=0 .
$$

In other words, $\mathrm{P}(f)(x)=0$. As this holds for all $\left.x \in \mathrm{X}\right|_{\mathrm{B}-\Delta}, \mathrm{P}(f)=0$.
(b) We show that $\mathscr{M}(\mathrm{X})$ is etale of degree $\leqslant d$ over $\mathscr{M}(\mathrm{B})$. For this, X may be assumed to be connected. Then $\mathscr{M}(\mathrm{X})$ is an extension of $\mathscr{M}(\mathrm{B})$. It is the directed union of finitely generated sub-extensions $E_{\lambda}$ of $\mathscr{M}(B)$. Each $E_{\lambda}$ is finitely generated.

It is generated by algebraic elements and so is finite and hence etale since $\mathscr{M}(\mathrm{B})$ is of characteristic 0 , and so cyclic by the primitive element theorem, thus of degree $\leqslant d$. Therefore the degree of $\mathscr{M}(\mathrm{X})$ is at most $d$.
(c) We show that $\mathscr{M}(\mathrm{X})$ is of degree $\geqslant d$.

The separation theorem has the following corollary:
Lemma Let X be a Riemann surface and $a_{1}, \ldots, a_{d}$ distinct points of X . Then there is a meromorphic function $f$ on X , defined at each point $a_{1}, \ldots, a_{d}$ with distinct values $f\left(a_{1}\right), \ldots, f\left(a_{d}\right)$.

Proof of the Lemma We show that for all $i, j$ with $i \neq j$, there is a function $f_{i j}$ defined at $a_{1}, \ldots, a_{d}$ and such that $f\left(a_{i}\right) \neq f\left(a_{j}\right)$. Let $i \neq j$. By the separation theorem, there is a function $g$ defined at $a_{i}$ and $a_{j}$ with $g\left(a_{i}\right) \neq g\left(a_{j}\right)$. If necessary adding a constant, we may assume that $g\left(a_{k}\right) \neq 0$ for all $k \in\{1, \ldots, d\}$. Then $f_{i j}=\frac{1}{g}$ has the required property. Let A be the vector subspace of $\mathscr{M}(\mathrm{X})$ consisting of the functions defined at $a_{1}, \ldots, a_{d}$. For each $(i, j)$ the set $\mathrm{V}_{i j}$ of $f \in \mathrm{~A}$ such that $f\left(a_{i}\right)=f\left(a_{j}\right)$ is a strict vector subspace of A . Hence $\bigcup \mathrm{V}_{i j} \neq \mathrm{A}$ and any function $f \in \mathrm{~A}-\left(\bigcup \mathrm{V}_{i j}\right)$ has the required property.

End of the Proof of the Theorem Let $b_{0} \in \mathrm{~B}-\Delta_{r}, \mathrm{X}\left(b_{0}\right)=\left\{x_{1}, \ldots, x_{d}\right\}$, and $f \in$ $\mathscr{M}(\mathrm{X})$ a function defined at $x_{1}, \ldots, x_{d}$ with distinct values. We show that $f$ has degree $\geqslant d$ over $\mathscr{M}(\mathrm{B})$, i.e. that the minimal polynomial of $f$ over $\mathscr{M}(\mathrm{B})$ has degree $\geqslant d$. Let $\mathrm{P}=\sum_{0}^{k} c_{i} \mathrm{Z}^{i}$ be a polynomial with coefficients $c_{i} \in \mathscr{M}(\mathrm{~B})$ and such that $\mathrm{P}(f)=0$, i.e., for all $x \in \mathrm{X}$ and $b=\pi(x)$ for which $f(x)$ and $c_{i}(b)$ are defined, $\sum_{0}^{k} c_{i}(b)(f(x))^{i}=0$. We show that $k \geqslant d$. Note that for $b$ sufficiently near $b_{0}$, $f$ is defined on $\mathrm{X}(b)$ with distinct values at its points. If $b$ is sufficiently near $b_{0}$ and is not a pole for any $c_{i}$, then $\mathrm{P}_{b}=\sum c_{i}(b) \mathrm{Z}^{i} \in \mathbb{C}[\mathrm{Z}]$ has at least $d$ distinct roots: the values of $f$ at the points of $\mathrm{X}(b)$. Hence $k \geqslant d$. This result in particular applies to the minimal polynomial of $f$. Consequently, $\operatorname{deg}_{\mu_{(\mathrm{B})}} f \geqslant d$ and $\operatorname{deg}_{\mu_{(\mathrm{B})}}$ $\mathscr{M}(\mathrm{X}) \geqslant d$.

Remark Let X be a finite analytically ramified covering of $\mathrm{B}, b \in \mathrm{~B}-\Delta$, and $f \in \mathscr{M}(\mathrm{X})$ a function defined on $\mathrm{X}(b)$ with distinct values at its points. The above proof shows that $f$ generates $\mathscr{M}(\mathrm{X})$ as an algebra over $\mathscr{M}(\mathrm{B})$.

### 6.2.4

The main aim of this section is to prove the next result:

Theorem Let $\mathscr{V}_{\mathrm{B}}$ be the category of finite analytically ramified coverings of B and $\mathscr{E}_{\mathrm{B}}$ the category of etale algebras over $\mathrm{K}=\mathscr{M}(\mathrm{B})$. Then, the functor $\mathscr{M}: \mathrm{X} \mapsto \mathscr{M}(\mathrm{X})$ from $\mathscr{V}_{\mathrm{B}}$ to $\mathscr{E}_{\mathrm{B}}$ is an anti-equivalence of categories.

The proof will be given in 6.2.9.

Corollary 6.1 Let X and Y be finite analytically ramified coverings of B , and $\Delta \subset \mathrm{B}$ a finite set such that X and Y are unramified over $\mathrm{B}-\Delta$. Then $\operatorname{Hom}_{\mathscr{E}_{\mathrm{B}}}(\mathscr{M}(\mathrm{X}), \mathscr{M}(\mathrm{Y}))$ can be identified with $\operatorname{Hom}_{\mathscr{C o v f}(\mathrm{B}-\Delta)}\left(\left.\mathrm{Y}\right|_{\mathrm{B}-\Delta},\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}\right)$.
Corollary 6.2 Let X be a connected finite analytically ramified covering of B , and $\Delta \subset \mathrm{B}$ a finite set containing the ramification set of X . Then $\mathscr{M}(\mathrm{X})$ is Galois extension of $\mathscr{M}(\mathrm{B})$ if and only if $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ is a Galois covering of $\mathrm{B}-\Delta$. In this case, the algebraic Galois group $\mathrm{Aut}_{\mu(\mathrm{B})} \mathscr{M}_{(\mathrm{X})}$ can be identified with the opposite group to the topological Galois group $\mathrm{Aut}_{\mathrm{B}-\Delta}\left(\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}\right)$. In particular these Galois groups are isomorphic.

### 6.2.5 Construction of a Covering $\mathcal{S}(\mathbf{E}, \zeta)$

Let E be an etale algebra over $\mathscr{M}(\mathrm{B})$ and $\zeta \in \mathrm{E}$ a primitive element, i.e. generating E as algebra (5.4.11). Let

$$
\mathrm{P}(\mathrm{Z})=\mathrm{Z}^{d}+a_{1} \mathrm{Z}^{d-1}+\cdots+a_{d}
$$

be the minimal polynomial of $\zeta$ in E . Denote by $\Delta_{\mathrm{E}, \zeta}$ the set of elements $b \in \mathrm{~B}$ where either the coefficients of P are not defined or $\mathrm{P}_{b}=\mathrm{Z}^{d}+a_{1}(b) \mathrm{Z}^{d-1}+\cdots+a_{d}(b) \in$ $\mathbb{C}[Z]$ has multiple roots.

Lemma The set $\Delta_{\mathrm{E}, \zeta}$ is finite.
Proof Let $\mathrm{P}=\mathrm{Z}^{d}+a_{1} \mathrm{Z}^{d-1}+\cdots+a_{d}$ be a monic polynomial over $\mathscr{M}(\mathrm{B})$. Consider the discriminant $\sigma$ of P (3.7.12). It is in $\mathscr{M}(\mathrm{B})$. Denote by $\Delta^{\prime}$ the set of points of B where some $a_{i}$ is not defined, and let $b \in \mathrm{~B}-\Delta^{\prime}$. All the roots of $\mathrm{P}_{b} \in \mathbb{C}[\mathrm{Z}]$ are distinct if and only if its discriminant $\sigma(b)$ is nonzero.

If $\mathscr{M}(\mathrm{B})[\mathrm{Z}] /(\mathrm{P})$ is etale, then by $5.4 .4, \sigma$ is a nonzero element of $\mathscr{M}(\mathrm{B})$. Then the set $\Delta^{\prime \prime}$ of the zeros of $\sigma$ is finite, and $\Delta_{\mathrm{E}, \zeta}=\Delta^{\prime} \cup \Delta^{\prime \prime}$ is finite

Denote by $\mathcal{\delta}(\mathrm{E}, \zeta)$ the subspace of $(\mathrm{B}-\Delta) \times \mathbb{C}$ consisting of pairs $(b, z)$ such that $\mathrm{P}(b, z)=0$. By 4.3.4, Example $6, \mathcal{S}(\mathrm{E}, \zeta)$ is a finite covering of $\mathrm{B}-\Delta_{\mathrm{E}, \zeta}$.

### 6.2.6 Functoriality

Let E and F be two etale algebras over $\mathscr{M}(\mathrm{B}), \zeta$ and $\theta$ primitive elements of E and F respectively, P and Q their minimal polynomials over $\mathscr{M}(\mathrm{B})$ and $\alpha: \mathrm{E} \rightarrow \mathrm{F}$ an algebra homomorphism. We do not assume that $\alpha(\zeta)=\theta$. We are going to define a morphism $\alpha^{*}: \mathcal{S}(\mathrm{F}, \theta) \rightarrow \mathcal{S}(\mathrm{E}, \zeta)$ in the category $\lim \mathscr{C}$ ovf $(\mathrm{B}-\Delta)$.

The element $\alpha(\zeta) \in \mathrm{F}$ takes the form $\mathrm{R}(\theta)$, where $\mathrm{R} \in \mathscr{M}(\mathrm{B})[\mathrm{Z}]$. Let $\Delta_{\mathrm{R}}^{\prime}$ be the set of points of $B$ where the coefficients of $R$ are not all defined. Define a map $\widetilde{R}$ from $\left(\mathrm{B}-\Delta_{\mathrm{R}}^{\prime}\right) \times \mathbb{C}$ to itself by $\widetilde{\mathrm{R}}(b, z)=(b, \mathrm{R}(b, z))$.

Lemma and Definition The map $\widetilde{\mathrm{R}}$ induces a map from $\left.\mathcal{S}(\mathrm{F}, \theta)\right|_{\mathrm{B}-\Delta_{\mathrm{R}}^{\prime}}$ to $\mathcal{S}(\mathrm{E}, \zeta)$, which does not depend on the choice of R. It is written $\alpha^{*}$.

Proof Since $\mathrm{P}(\mathrm{R}(\theta))=\mathrm{P}(\alpha(\zeta))=\alpha(\mathrm{P}(\zeta))=0 \in \mathrm{~F}, \mathrm{P} \circ \mathrm{R}=\mathrm{Q} \cdot \mathrm{V}$ for some $\mathrm{V} \in$ $\mathscr{M}(\mathrm{B})[\mathrm{Z}]$. In other words, $\mathrm{P}(b, \mathrm{R}(b, z))=\mathrm{Q}(b, z) \mathrm{V}(b, z)$ for $b \in \mathrm{~B}$ such that the coefficients of $\mathrm{P}, \mathrm{Q}, \mathrm{R}, \mathrm{V}$ are defined and $z \in \mathbb{C}$. If $(b, z) \in \mathcal{S}(\mathrm{F}, \theta)$, then $\mathrm{Q}(b, z)=0$; so $\mathrm{P}(b, \mathrm{R}(b, z))=0$, and $\widetilde{\mathrm{R}}(b, z)=(b, \mathrm{R}(b, z)) \in \mathcal{\delta}(\mathrm{E}, \zeta)$.

Let $\mathrm{R}_{1} \in \mathscr{M}(\mathrm{~B})[\mathrm{Z}]$ be another polynomial such that $\alpha(\zeta)=\mathrm{R}_{1}(\theta)$. Then $\mathrm{R}-\mathrm{R}_{1}$ vanishes at $\theta$, and so is of the form $\mathrm{Q} \cdot \mathrm{W}$ with $\mathrm{W} \in \mathscr{M}(\mathrm{B})[\mathrm{Z}]$. If $\mathrm{Q}(b, z)=0$, $\mathrm{R}(b, z)=\mathrm{R}_{1}(b, z)$; so $\widetilde{\mathrm{R}}$ and $\widetilde{\mathrm{R}}_{1}$ agree on $\mathcal{S}(\mathrm{F}, \theta)$.

Let $\mathrm{E}, \mathrm{F}$ and G be etale algebras over $\mathscr{M}(\mathrm{B}), \zeta, \theta$ and $\lambda$ primitive elements of E , F and G respectively, $\alpha: \mathrm{E} \rightarrow \mathrm{F}$ and $\beta: \mathrm{F} \rightarrow \mathrm{G}$ homomorphisms. Then $(\beta \circ \alpha)^{*}=$ $\alpha^{*} \circ \beta^{*}$. Indeed, let $\mathrm{R}_{1}$ and $\mathrm{R}_{2} \in \mathscr{M}(\mathrm{~B})[\mathrm{Z}]$ be such that $\alpha(\zeta)=\mathrm{R}_{1}(\theta)$ and $\beta(\theta)=$ $\mathbf{R}_{2}(\lambda)$. So $\beta(\alpha(\zeta))=\beta\left(\mathbf{R}_{1}(\theta)\right)=\mathbf{R}_{1}(\beta(\theta))=\mathbf{R}_{1}\left(\mathbf{R}_{2}(\lambda)\right)$. Hence for $\mathbf{R}_{3}$ such that $\underset{\sim}{\beta} \circ \alpha(\zeta)=\mathrm{R}_{3}(\lambda)$, we may choose $\mathrm{R}_{3}=\mathrm{R}_{1} \circ \mathrm{R}_{2}$. Then $(\beta \circ \alpha)^{*}$ is induced by $\widetilde{\mathrm{R}}_{3}=$ $\widetilde{\mathrm{R}}_{1} \circ \widetilde{\mathrm{R}}_{2}$, and so $(\beta \circ \alpha)^{*}=\alpha^{*} \circ \beta^{*}$ since $\alpha^{*}$ and $\beta^{*}$ are respectively induced by $\widetilde{\mathrm{R}}_{1}$ and $\widetilde{\mathrm{R}}_{2}$.

### 6.2.7

Let X be a finite analytically ramified covering of B and $\Delta_{\mathrm{X}}$ its ramification set. Let $f \in \mathscr{M}(\mathrm{X})$ and $\Delta_{f}^{\prime}$ be the projection of the set of points where $f$ is not defined. Assume that there exists $b \in \mathrm{~B}-\Delta_{\mathrm{X}}$ such that $f$ is defined on $\mathrm{X}(b)$ with distinct values at its points. Then we know (6.2.3, Remark) that $f$ is a primitive element of $\mathscr{M}(\mathrm{X})$ and the covering $\mathcal{S}(\mathscr{M}(\mathrm{X}), f) \rightarrow \mathrm{B}-\Delta_{f}$ may be considered. Set $\Delta=$ $\Delta_{\mathrm{X}} \cup \Delta_{f}^{\prime} \cup \Delta_{f}$.

Proposition The map $(\pi, f):\left.\mathrm{X}\right|_{\mathrm{B}-\Delta} \rightarrow(\mathrm{B}-\Delta) \times \mathbb{C}$ induces a covering isomorphism from $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ onto $\left.\mathcal{S}(\mathscr{M}(\mathrm{X}), f)\right|_{\mathrm{B}-\Delta}$.

Proof Let P be the minimal polynomial of $f$ in $\mathscr{M}(\mathrm{X})$. For $x \in \mathrm{X}$ such that $f$ is defined at $x$ and the coefficients of P are defined at $b=\pi(x), \mathrm{P}(b, f(x))=0$. Hence $(\pi, f)$ is indeed a continuous map from $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ to $\mathcal{\delta}(\mathscr{M}(\mathrm{X}), f)$; since it commutes with the projections, it is a covering morphism. When $b \in \Delta$ (which can be shown never to occur), take $b^{\prime} \in \mathrm{B}-\Delta$ sufficiently near $b$ so that the assumptions made still hold for $b$. Then the covering morphism $(\pi, f)$ is bijective on the fibres at $b^{\prime}$. As a result of (4.3.6. Corollary 4.2) and (6.1.7 remarque 3 ), it is an isomorphism.

A Riemann surface of $X$ is said to be algebraic if there is a polynomial $Q \in \mathbb{C}[T, Z]$ (where $T$ and $Z$ are indeterminates), and finite sets $\Theta \subset X$ and $\Lambda \subset V=Q^{-1}(0)$ such that $\mathrm{X}-\Theta$ and $\mathrm{V}-\Lambda$ are isomorphic Riemann surfaces.

Corollary Every compact Riemann surface is algebraic.

Proof Let X be a compact Riemann surface and $\pi \in \mathscr{M}^{* *}(\mathrm{X})$ (i.e. $\pi$ is a meromorphic function on $X$ which is not constant on any connected component). Then (X, $\pi$ ) is an analytically ramified covering of $\Sigma$, and $\mathscr{M}(\mathrm{X})$ equipped with $\pi^{*}$ is an etale algebra over $\mathscr{M}(\Sigma)=\mathbb{C}(\mathrm{T})$. Choose a primitive element $f \in \mathscr{M}(\mathrm{X})$, and let $\mathrm{P} \in$ $\mathbb{C}(T)[Z]$ be its minimal polynomial. Then, there are $\mathrm{Q} \in \mathbb{C}[T, Z]$ and $\mathrm{D} \in \mathbb{C}[T]$ such that $\mathrm{P}(\mathrm{T}, \mathrm{Z})=\mathrm{Q}(\mathrm{T}, \mathrm{Z}) / \mathrm{D}(\mathrm{T})$. Define $\Delta$ as above. Then $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ is isomorphic to $\left.\mathcal{S}(\mathscr{M}(\mathrm{X}), f)\right|_{\mathrm{B}-\Delta}=\{(x, z) \in(\mathbb{C}-\Delta) \times \mathbb{C} \mid \mathrm{Q}(x, z)=0\}$.

### 6.2.8

Let E be an etale algebra over $\mathscr{M}(\mathrm{B})$ and $\zeta \in \mathrm{E}$ a primitive element. Consider the covering $\mathcal{S}(\mathrm{E}, \zeta)$ of $\mathrm{B}-\Delta$. By (6.1.12), there is an analytically ramified covering $\widehat{\mathcal{S}}(\mathrm{E}, \zeta)$ of B such that $\left.\widehat{\mathcal{S}}(\mathrm{E}, \zeta)\right|_{\mathrm{B}-\Delta}=\mathcal{S}(\mathrm{E}, \zeta)$. Denote by Z the map from $\mathcal{S}(\mathrm{E}, \zeta)$ to $\mathbb{C}$ induced by the second projection $\mathrm{B} \times \mathbb{C} \rightarrow \mathbb{C}$.
Proposition (a) The function Z is meromorphic on $\widehat{\mathcal{S}}(\mathrm{E}, \zeta)$.
(b) There is a unique isomorphism $\varphi: \mathrm{E} \rightarrow \mathscr{M}(\widehat{\mathcal{S}}(\mathrm{E}, \zeta))$ such that $\varphi(\zeta)=\mathrm{Z}$.

To prove (a), we use the criterion for meromorphism (6.2.2). To do so, we will need an upper bound.
Lemma (Upper bound of the roots of a polynomial) Let

$$
\mathrm{P}=\mathrm{Z}^{d}+a_{1} \mathrm{Z}^{d-1}+\cdots+a_{d}
$$

be a polynomial over $\mathbb{C}$ and $z$ a root of P . Then

$$
|z| \leqslant \sup \left(1,\left|a_{1}\right|+\cdots+\left|a_{d}\right|\right)
$$

Proof If $|z| \geqslant 1$, then $z=-a_{1}-\frac{a_{2}}{z}-\cdots-\frac{a_{d}}{z^{d-1}}, \quad$ and so $|z| \leqslant\left|a_{1}\right|+\cdots$ $+\left|a_{d}\right|$.
Proof of the Proposition (a) The function Z is clearly holomorphic over $\mathrm{B}-\Delta$. Let $b \in \Delta$ and $x \in \mathrm{X}(b)$, where $\mathrm{X}=\widehat{\mathcal{S}}(\mathrm{E}, \zeta)$. Let $\varphi$ be a chart for B centered at $b$ and $\psi$ a chart of X centered at $x$ such that the expression of $\pi$ in the charts $\psi$ and $\varphi$ is $z \mapsto z^{r}$. The coefficients $a_{1}, \ldots, a_{d}$ of the minimal polynomial P of $\zeta$ are meromorphic functions on B. Let $\dot{a}_{1}, \ldots, \dot{a}_{d}$ be their expression in the chart $\varphi$. For $x^{\prime}$ near $x$, the number $\mathrm{Z}\left(x^{\prime}\right)$ is a root of $\mathrm{Z}^{d}+a_{1}\left(b^{\prime}\right) \mathrm{Z}^{d-1}+\cdots+a_{d}\left(b^{\prime}\right)$, where $b^{\prime}=\pi\left(x^{\prime}\right)$, and so

$$
\left|\mathrm{Z}\left(x^{\prime}\right)\right| \leqslant \sup \left(1,\left|a_{1}\left(b^{\prime}\right)\right|+\cdots+\left|a_{d}\left(b^{\prime}\right)\right|\right)
$$

However, $a_{i}\left(b^{\prime}\right)=\dot{a}_{i}\left(\varphi\left(b^{\prime}\right)\right)=\dot{a}_{i}\left(\psi\left(x^{\prime}\right)^{r}\right)$. Being meromorphic, the functions $a_{i}$ can be bounded above as described in (6.2.2). Hence the same holds for Z , and Z is meromorphic.
(b) For $(b, z) \in \mathcal{S}(\mathrm{E}, \zeta), \mathrm{P}(b, z)=0$. In other words, $\mathrm{P}(\mathrm{Z})=0$ in $\mathscr{M}(\mathrm{X})$. We show that P is the minimal polynomial of Z in $\mathscr{M}(\mathrm{X})$. Let Q be a polynomial over $\mathscr{M}(\mathrm{B})$ such that $\mathrm{Q}(\mathrm{Z})=0$. We show that Q is a multiple of P . If necessary dividing Q by P , the degree of Q can be assumed to be strictly less than the degree $d$ of P . For all $b \in \mathrm{~B}-\Delta$ where the coefficients of Q are defined, the points of $\mathrm{X}(b)$ are the pairs $(b, z), z$ being a root of $\mathrm{P}_{b}$, and so the roots of $\mathrm{P}_{b}$ are roots of $\mathrm{Q}_{b}$. As there are $d$ distinct ones, $\mathrm{Q}_{b}=0$. Hence $\mathrm{Q}=0$.

So there is a unique homomorphism of $\mathscr{M}(\mathrm{B})$-algebras $\varphi: \mathrm{E} \rightarrow \mathscr{M}(\mathrm{X})$ such that $\varphi(\zeta)=\mathrm{Z}$, and $\varphi$ is injective. By assumption, the algebra E has degree $d$, and X is a ramified $d$-fold covering, and so by (6.2.3, Theorem) $\mathscr{M}(\mathrm{X})$ is an algebra of degree $d$. Hence $\varphi$ is an isomorphism.

### 6.2.9 Proof of Theorem 6.2.4

(a) The functor $\mathscr{M}$ is essentially surjective. Let E be an etale algebra over $\mathscr{M}(\mathrm{B})$, and $\zeta$ a primitive element of E. By Proposition 6.2.8, E is isomorphic to $\mathscr{M}(\widehat{\mathcal{S}}(\mathrm{E}, \zeta))$.
(b) The functor $\mathscr{M}$ is fully faithful. Let X and Y be finite analytically ramified covering of B and $u: \mathscr{M}(\mathrm{X}) \rightarrow \mathscr{M}(\mathrm{Y})$ a $\mathscr{M}(\mathrm{B})$-algebra homomorphism. We show that there is a morphism $h: \mathrm{Y} \rightarrow \mathrm{X}$ such that $h^{*}=u$. Let $f$ and $g$ be primitive elements of $\mathscr{M}(\mathrm{X})$ and $\mathscr{M}(\mathrm{Y})$ respectively, satisfying assumptions 6.2.7. The isomorphism $(\pi, f):\left.\mathrm{X}\right|_{\mathrm{B}-\Delta} \rightarrow \mathcal{S}(\mathscr{M}(\mathrm{X}), f)$ extends to an isomorphism $\tilde{f}: \mathrm{X} \rightarrow \widehat{\mathcal{S}}(\mathscr{M}(\mathrm{X}), f)$. Define $\tilde{g}: \mathrm{Y} \rightarrow \widehat{\mathcal{S}}(\mathscr{M}(\mathrm{Y}), g)$ likewise. By 6.1 .12 the morphism $u^{*}: \mathcal{S}(\mathscr{M}(\mathrm{Y}), g) \rightarrow \mathcal{S}(\mathscr{M}(\mathrm{X}), f)$ extends to a morphism which is also written $u^{*}: \widehat{\delta}(\mathscr{M}(\mathrm{Y}), g) \rightarrow \widehat{\delta}(\mathscr{M}(\mathrm{X}), f)$. Define $h: \mathrm{Y} \rightarrow \mathrm{X}$ by the commutativity of the diagram


We show that $h^{*}=u$. As $f$ is primitive, it suffices to show that $u(f)=h^{*}(f)=$ $f \circ h: \mathrm{Y} \rightarrow \mathbb{C}$. Now, $f \circ h=\mathrm{Z} \circ \tilde{f} \circ h=\mathrm{Z} \circ u^{*} \circ \tilde{g}$. For $y \in \mathrm{Y}(b), u^{*}(\tilde{g}(y))=$ $u^{*}(b, g(y))=\widetilde{\mathrm{R}}(b, g(y))=(b, \mathrm{R}(b, g(y))$, where R is a polynomial over $\mathscr{M}(\mathrm{B})$ such that $u(f)=\mathrm{R}(g)$. Hence $\mathrm{Z} \circ u^{*} \circ \tilde{g}=\mathrm{R}(g)=u(f)$, and so $h^{*}=u$.

We next show uniqueness. Let $h_{1}$ be such that $h_{1}^{*}=u$. Then $h_{1}^{*}(f)=h^{*}(f)$, in other words $f \circ h_{1}=f \circ h$, where $\tilde{f} \circ h_{1}=\tilde{f} \circ h$, and so $h_{1}=h$ since $\tilde{f}$ is an isomorphism.

Remark Consider the category $\mathscr{V}_{\mathrm{B}}$ of finite analytically ramified coverings of B , the category $\mathscr{E}_{\mathrm{B}}$ of etale algebras over $\mathscr{M}(\mathrm{B})$ and the category $\mathscr{R}=\lim \mathscr{C} \oslash \mho \mathcal{f}_{\mathrm{B}-\Delta}$. For each $\mathrm{E} \in \mathscr{E}_{\mathrm{B}}$ choose a primitive element. This defines a functor $\overrightarrow{\mathcal{S}}: \mathscr{E}_{\mathrm{B}} \rightarrow \mathscr{R}$. Then up to functorial isomorphism, there is a commutative diagram:

where $\omega$ is the forgetful functor. In this diagram, the three arrows are either equivalences or anti-equivalences of categories.

### 6.2.10 Dictionary

Let X and Y be finite analytically ramified coverings of S .
(a) The degree of the algebra $\mathscr{M}(\mathrm{X})$ over $\mathscr{M}(\mathrm{B})$ is equal to the degree of X over B , i.e. to the degree of the covering $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ of $\mathrm{B}-\Delta$ (see 6.2.3).
(b) $\mathscr{M}(\mathrm{X} \sqcup \mathrm{Y})=\mathscr{M}(\mathrm{X}) \times \mathscr{M}(\mathrm{Y})$. The algebra $\mathscr{M}(\mathrm{X})$ is an extension of $\mathscr{M}(\mathrm{B})$ if and only if X is connected. More generally, the connected components of X are in bijective correspondence with the maximal ideals of $\mathscr{M}(\mathrm{X})$.
(c) In general, the fibre product $\mathrm{X} \times_{B} \mathrm{Y}$ is not a ramified covering of B , but a finite analytically ramified covering Z of B satisfying $\left.\mathrm{Z}\right|_{\mathrm{B}-\Delta}=\left.\left(\mathrm{X} \times_{\mathrm{B}} \mathrm{Y}\right)\right|_{\mathrm{B}-\Delta}$, where $\Delta=\Delta_{\mathrm{X}} \cup \Delta_{\mathrm{Y}}$, may be constructed. Then $\mathscr{M}(\mathrm{Z})=\mathscr{M}(\mathrm{X}) \otimes_{\mathscr{M}(\mathrm{B})} \mathscr{M}(\mathrm{Y})$.
(d) Suppose that X is connected. Set $\Delta=\Delta_{\mathrm{X}} \cup \Delta_{\mathrm{Y}}$. Then $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ trivializes the covering $\left.\mathrm{Y}\right|_{\mathrm{B}-\Delta}$ is and only if $\mathscr{M}(\mathrm{X})$ diagonalizes $\mathscr{M}(\mathrm{Y})$. In particular $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta_{\mathrm{X}}}$ is a Galois covering if and only if $\mathscr{M}(\mathrm{X})$ is a Galois extension of $\mathscr{M}(\mathrm{B})$. In this case, $\operatorname{Aut}_{\mathscr{M}(\mathrm{B})}(\mathscr{M}(\mathrm{X}))=\operatorname{Aut}_{\mathrm{B}}(\mathrm{X})=\operatorname{Aut}_{\mathrm{B}-\Delta}\left(\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}\right)$ since the functor $\mathscr{M}$ is fully faithful.
(e) Let $f: \mathrm{Y} \rightarrow \mathrm{X}$ be morphism. Then $f^{*}: \mathscr{M}(\mathrm{X}) \rightarrow \mathscr{M}(\mathrm{Y})$ is injective (resp. surjective) if and only $f$ is surjective (resp. injective).
(f) Let $f \in \mathscr{M}(\mathrm{X})$. Then $f$ is a primitive element of $\mathscr{M}(\mathrm{X})$ if and only if there exists $b \in \mathrm{~B}-\Delta_{\mathrm{X}}$ such that $f$ defines an injective map from $\mathrm{X}(b)$ to $\mathbb{C}$. This condition has been seen to be sufficient (6.2.3, Remark). We show that it is necessary: suppose that $f$ is primitive and consider the morphism $(\pi, f):\left.\mathrm{X}\right|_{\mathrm{B}-\Delta} \rightarrow \mathrm{S}(\mathscr{M}(\mathrm{X}), f)$. Its image is a clopen subset $\mathrm{S}_{1}$ of $\mathrm{S}(\mathscr{M}(\mathrm{X}), f)$, which can be embedded in an analytic covering $\widehat{\mathrm{S}_{1}}$. The morphism $\tilde{f}: \mathrm{X} \rightarrow \widehat{\mathrm{S}_{1}}$ extending $(\pi, f)$ defines a homomor$\operatorname{phism} \tilde{f}^{*}: \mathscr{M}\left(\widehat{\mathrm{S}_{1}}\right) \rightarrow \mathscr{M}(\mathrm{X})$ which is surjective since $f$ is in the image. Hence $\operatorname{deg}_{\mathrm{B}}\left(\widehat{\mathrm{S}}_{1}\right) \geqslant \operatorname{deg}_{\mathrm{B}}(\mathrm{X})=\operatorname{deg}_{\mathrm{B}-\Delta}(\mathrm{S}(\mathscr{M}(\mathrm{X}), f))$, and so $\mathrm{S}_{1}=\mathrm{S}(\mathscr{M}(\mathrm{X}), f)$ and $f$ is an isomorphism. Thus, for $b \in \mathrm{~B}-\Delta$, the map $f$ induces an injection from $\mathrm{X}(b)$ to $\mathbb{C}$.

### 6.2.11

Let E be an etale algebra over $\mathscr{M}(\mathrm{B})$. Then, up to isomorphism, there is a unique analytically ramified covering X of B such that E is isomorphic to $\mathscr{M}(\mathrm{X})$. Let $\Delta_{\mathrm{E}}$ be the ramification set of $X$. It is said to be the ramification set of $E$.

If F is a subalgebra of E , then $\Delta_{\mathrm{F}} \subset \Delta_{\mathrm{E}}$. Indeed, if $\mathrm{F} \approx \mathscr{M}(\mathrm{Y})$ is a surjective morphism $f: \mathrm{X} \rightarrow \mathrm{Y}$, and if $b \notin \Delta_{\mathrm{E}}$, then there is a local section of X passing through all points of $\mathrm{X}(b)$. Hence composing with $f$ gives a local section of Y passing through all points of $\mathrm{Y}(b)$ and $b \notin \Delta_{\mathrm{F}}$.

Similarly, if $F$ is a quotient of $E$, then $\Delta_{\mathrm{F}} \subset \Delta_{\mathrm{E}}$.
If E and F are etale algebras over $\mathscr{M}(\mathrm{B})$, then $\Delta_{\mathrm{E} \times \mathrm{F}}=\Delta_{\mathrm{E}} \cup \Delta_{\mathrm{F}}$ and $\Delta_{\mathrm{E} \otimes \mathrm{F}} \subset$ $\Delta_{\mathrm{E}} \cup \Delta_{\mathrm{F}}$ by 6.2.10, (b) and (c). If E and F are nonzero, then $\Delta_{\mathrm{E} \otimes \mathrm{F}}=\Delta_{\mathrm{E}} \cup \Delta_{\mathrm{F}}$ since $E$ and $F$ can then be identified with subalgebras of $E \otimes F$.

If F and G are subalgebras of an etale algebra E over $\mathscr{M}(\mathrm{B})$, then $\Delta_{\mathrm{F} \cdot \mathrm{G}}=\Delta_{\mathrm{F}} \cup$ $\Delta_{G}$. Indeed $\Delta_{F \cdot G} \subset \Delta_{F} \cup \Delta_{G}$ since $F \cdot G$ is a quotient de $F \otimes G$, and the opposite inclusion holds since F and G are subalgebras of $\mathrm{F} \cdot \mathrm{G}$.

### 6.2.12

Proposition With the notation of (6.2.5) and (6.2.11), the set $\Delta_{\mathrm{E}}$ is the intersection of all $\Delta_{\mathrm{E}, \zeta}$, where $\zeta$ is a primitive element of the $\mathscr{M}(\mathrm{B})$-algebra E .

Proof For a primitive element $\zeta, \mathrm{S}(\mathrm{E}, \zeta)$ is an (unramified) covering of $\mathrm{B}-\Delta_{\mathrm{E}, \zeta}$, and so $\Delta_{\mathrm{E}} \subset \Delta_{\mathrm{E}, \zeta}$. Hence $\Delta_{\mathrm{E}} \subset \bigcap \Delta_{\mathrm{E}, \zeta}$.

Let $b \in \mathrm{~B}-\Delta_{\mathrm{E}, \zeta}$, and $\mathrm{X}(b)=\left\{x_{1}, \ldots, x_{d}\right\}$. By the separation theorem (6.2.1), there is a function $f \in \mathscr{M}(\mathrm{X})$ defined at $x_{1}, \ldots, x_{d}$, with distinct values at these points. Let $\mathrm{P} \in \mathscr{M}(\mathrm{B})[\mathrm{Z}]$ be the minimal polynomial of $f$ on $\mathscr{M}(\mathrm{B})$. For all $b^{\prime} \in \mathrm{B}-\Delta_{\mathrm{P}}$, $\mathrm{P}_{b^{\prime}} \in \mathbb{C}[\mathrm{Z}]$ vanishes at $f\left(x^{\prime}\right)$ for $x^{\prime} \in \mathrm{X}\left(b^{\prime}\right)$, and so the coefficients of $\mathrm{P}_{b^{\prime}}$ are the elementary symmetric functions of these numbers (with sign $(-1)^{d-i}$ ). Hence these coefficients are defined at $b$ and the roots of $\mathrm{P}_{b}$ are $f\left(x_{1}\right), \ldots, f\left(x_{d}\right)$. As these are distinct, the discriminant of $\mathrm{P}_{b}$ is nonzero and $b \notin \Delta_{\mathrm{E}, f}$. Since the degree of P is $d$, the function $f$ is a primitive element of E and $b \notin \bigcap \Delta_{\mathrm{E}, \zeta}$. Thus $\bigcap \Delta_{\mathrm{E}, \zeta} \subset \Delta_{\mathrm{E}}$.

## Exercises 6.2. (Ramified coverings and etale algebras)

1. (Proof of the separation theorem)—Let X be a Riemann surface. A differential form ${ }^{2} \alpha$ of degree 1 on X assigns to each point $x \in \mathrm{X}$ a $\mathbb{R}$-linear map $\alpha_{x}$ from the tangent space $\mathrm{T}_{x} \mathrm{X}$, a 1-dimensional $\mathbb{C}$-vector space, to $\mathbb{C}$. For a $\mathbb{C}^{1}$ function $f$ on X , the differential of $f$ is the 1 -form $d f$, where $d f_{x}(u)=u . f$ denotes the derivative of $f$ along the vector $u \in \mathrm{~T}_{x} \mathrm{X}$. If X is an open subset of $\mathbb{C}$, every 1 -form can be

[^20]written $\alpha=f_{1}(z) d x_{1}+f_{2}(z) d x_{2}$ in a chart for X with coordinate $z=x_{1}+i x_{2}$. In particular
$$
d f=\frac{\partial f}{\partial x_{1}} d x_{1}+\frac{\partial f}{\partial x_{2}} d x_{2}
$$

If for all $x$, the map $\alpha_{x}$ is $\mathbb{C}$-linear (resp. $\mathbb{C}$-antilinear), $\alpha$ is said to be of type $(1,0)$ (resp. $(0,1)$ ). The forms of type $(1,0)$ (resp. $(0,1))$ with respect to a chart are $g(z) d z$ (resp. $h(z) d \bar{z})$.

A differential 2-form $\eta$ assigns to each point $x \in \mathrm{X}$ an alternating $\mathbb{R}$-bilinear map $\eta_{x}: \mathrm{T}_{x} \mathrm{X} \times \mathrm{T}_{x} \mathrm{X} \rightarrow \mathbb{C}$. It is said to be real (resp. positive, resp. strictly positive) if

$$
(\forall x \in \mathrm{X})\left(\forall u \in \mathrm{~T}_{x} \mathrm{X}-\{0\}\right) \quad \eta_{x}(u, i u) \in \mathbb{R}(\text { resp. } \geqslant 0, \text { resp. }>0)
$$

If $\eta$ is a 2-form, then define $\int_{\mathrm{X}} \eta$ as in Bourbaki [6], $\S 10$.
If $\alpha$ and $\beta$ are 1 -forms, define the 2 -form $\alpha \wedge \beta$ by

$$
(\alpha \wedge \beta)_{x}(u, v)=\alpha_{x}(u) \beta_{x}(v)-\alpha_{x}(v) \beta_{x}(u)
$$

Finally the exterior derivative of a $\mathrm{C}^{1} 1$-form $\alpha$ is a 2-form written $d \alpha$. whose definition can be found in [4]. For practical calculations, in particular for question $b, \gamma$ ), it is sufficient to know the formulas $d(f d g)=d f \wedge d g$ (for $\mathrm{C}^{1}$ functions $f$, $g)$ and $\int_{\mathrm{X}} d \alpha=0$ for a $\mathrm{C}^{1} 1$-form $\alpha$ with compact support.

## A. Hilbert norm

[For the exercise, the results in this part may be assumed.]
Let $\mathbb{D} \subset \mathbb{C}$ be the open unit disc. If $f: \mathbb{D} \rightarrow \mathbb{C}$ is a continuous function then for $\mathrm{A} \subset \mathbb{D}$ set

$$
\|f\|_{\mathrm{A}}=\left(\int_{\mathrm{A}}|f(x+i y)|^{2} d x d y\right)^{\frac{1}{2}}
$$

If $\alpha=u d x+v d y$ is a differential form then set

$$
\|\alpha\|_{\mathrm{A}}=\left(\|u\|_{\mathrm{A}}^{2}+\|v\|_{\mathrm{A}}^{2}\right)^{\frac{1}{2}}
$$

(a) Show that if A and B are compact subsets in $\mathbb{D}$ such that $\mathrm{A} \subset \stackrel{\circ}{\mathrm{B}}$, then there is a constant $q$ such that for any $\mathrm{C}^{1}$ function $f$ on $\mathbb{D}$ (i.e. continuously $\mathbb{R}$-differentiable), there is a constant M satisfying

$$
\|f-\mathrm{M}\|_{\mathrm{A}} \leqslant q\|d f\|_{\mathrm{B}}
$$

Does this inequality continue to hold with $\mathrm{B}=\mathrm{A}$ when A is a disc?
(b) For a $\mathrm{C}^{1}$ function $f$ on $\mathbb{D}$, define

$$
\frac{\partial f}{\partial z} \text { and } \frac{\partial f}{\partial \bar{z}} \text { by } d f=\frac{\partial f}{\partial z} d z+\frac{\partial f}{\partial \bar{z}} d \bar{z}
$$

Let $h$ be a square-integrable function on $\mathbb{D}$. Show that the following conditions are equivalent:
(i) $h$ agrees nearly everywhere with a holomorphic function;
(ii) for any $\mathrm{C}^{1}$ function $f$ with compact support in $\mathbb{D}$,

$$
\int_{\mathbb{D}} h \frac{\partial f}{\partial \bar{z}} d x d y=0
$$

## B. Forms of type $(1,0)$ or $(0,1)$

Let X be a connected compact Riemann surface.
(a) Show that every 1 -form on X is the sum of a unique form of type $(1,0)$ and a unique form of type $(0,1)$.
(b) For every 1-form $\alpha$, define a form $* \alpha$ by

$$
(* \alpha)_{x}=-\tau \circ \alpha_{x} \circ \mathbf{J}
$$

where $\tau: \mathbb{C} \rightarrow \mathbb{C}$ is conjugation $z \mapsto \bar{z}$, and $\mathbf{J}$ multiplication by $i$. Show that, if $\alpha$ is of type $(1,0)$ (resp. $(0,1))$, then $* \alpha$ is of type $(0,1)($ resp. $(1,0))$. Show that, for every form $\alpha, \alpha \wedge * \alpha$ is real positive, and strictly positive at all points where $\alpha$ does not vanish.

Show that if $\alpha$ are $\beta$ two 1 -forms, then

$$
\alpha \wedge * \beta=\overline{\beta \wedge * \alpha}
$$

Set $(\alpha \mid \beta)=\int_{\mathrm{X}} \alpha \wedge * \beta$. With this scalar product, the space of continuous 1-forms is a pre-Hilbert space. Let $\mathrm{E}^{1}$ be its completion. Show that $\mathrm{E}^{1}=\mathrm{E}^{1,0} \oplus \mathrm{E}^{0,1}$, where $\mathrm{E}^{1,0}$ is the closure of the space of continuous forms of type $(1,0)$ and $\mathrm{E}^{0,1}$ is defined likewise, and that these subspaces are orthogonal.
(c) For a $\mathrm{C}^{1}$ function $f$, let $d^{\prime} f$ and $d^{\prime \prime} f$ be the parts of type $(1,0)$ and $(0,1)$ of $d f$. Show that $* d^{\prime} f=-i d^{\prime \prime} \bar{f}$ and $* d^{\prime \prime} f=i d^{\prime} \bar{f}$. Assuming that $f$ is $\mathrm{C}^{2}$, calculate $d(f d \bar{f})$, and show that $\left\|d^{\prime} f\right\|=\left\|d^{\prime \prime} f\right\|$. Does this inequality hold when $f$ is $\mathrm{C}^{1}$ ?

## C. Hilbert Complex

Let $\omega$ be a continuous 2-form strictly positive on X . Let $\mathrm{E}^{0}$ be the completion of the space of $\mathrm{C}^{1}$ functions with respect to the norm

$$
f \mapsto\left(\int|f|^{2} \omega+\int d f \wedge * d f\right)^{\frac{1}{2}}
$$

The maps $f \mapsto d f, f \mapsto d^{\prime} f$ and $f \mapsto d^{\prime \prime} f$ extend to maps $d: \mathrm{E}^{0} \rightarrow \mathrm{E}^{1}, d^{\prime}: \mathrm{E}^{0} \rightarrow$ $\mathrm{E}^{1,0}$, and $d^{\prime \prime}: \mathrm{E}^{0} \rightarrow \mathrm{E}^{0,1}$. Using A , (a) show that $d: \mathrm{E}^{0} \rightarrow \mathrm{E}^{1}$ is a strict morphism (i.e. that the algebraic isomorphism $\mathrm{E}^{0} / \operatorname{Ker} d \rightarrow \operatorname{Im} d$ is a homeomorphism). Generalize the equality proved in B , (c) to $f \in \mathrm{E}^{0}$. Deduce that $d^{\prime \prime}: \mathrm{E}^{0} \rightarrow \mathrm{E}^{0,1}$ is a strict morphism. What is its kernel?

## D. Finiteness theorem

(a) For 1 -forms $\alpha$ and $\beta$, set

$$
\langle\alpha, \beta\rangle=\int \alpha \wedge \beta
$$

Show that $(\alpha, \beta) \mapsto\langle\alpha, \beta\rangle$ extends to a bilinear form on $\mathrm{E}^{1} \times \mathrm{E}^{1} \rightarrow \mathbb{C}$, which defines an isomorphism between $\mathrm{E}^{1,0}$ and the dual of $\mathrm{E}^{0,1}$.
(b) Define $\mathrm{E}^{2}$ to be the dual of $\mathrm{E}^{0}$. Show that $\alpha \mapsto d \alpha$ extends to a continuous linear map $d: \mathrm{E}^{1} \rightarrow \mathrm{E}^{2}$ characterized by

$$
\langle f, d \alpha\rangle=-\langle d f, \alpha\rangle .
$$

(c) Let $d^{\prime \prime}: \mathrm{E}^{1,0} \rightarrow \mathrm{E}^{2}$ be the map induced by $d$. Show that its kernel is the space $\Omega$ of holomorphic differential forms on X, i.e. whose local expressions are $f(z) d z$ for $f$ holomorphic (use A, (b)).
d) Let H denote the cokernel of $d^{\prime \prime}: \mathrm{E}^{0} \rightarrow \mathrm{E}^{0,1}$. Show that the scalar product $(\alpha, \beta) \mapsto\langle\alpha, \beta\rangle$ defines an isomorphism between H and the dual of $\Omega$ (a particular case of the Serre duality theorem).
(e) Using (4.6, Exercise 13) and (4.9, Exercise 2), show that $\Omega$ is finite dimensional. Deduce that so is H (finiteness theorem).
(e') Using Montel's theorem for normal families, show that $\Omega$ is locally compact. Deduce another proof of the finiteness of the dimension of $\Omega$, hence of H .

## E. Construction of meromorphic functions

Let $x \in \mathrm{X}$ and $(\mathrm{U}, \mathrm{Z})$ an analytic chart for X centered at $x$.
Let $\eta$ be a $\mathrm{C}^{1}$ function with compact support in U , equal to 1 in the neighbourhood of $x$. Set

$$
\alpha_{n}=d^{\prime \prime}\left(\eta \cdot \frac{1}{\mathrm{Z}^{n}}\right) .
$$

(a) Show that the form $\alpha_{n}$ is continuous of type $(0,1)$ with support in U . Extend it by 0 to a form on X .
(b) Show that there exist nonzero numbers $c_{k}$ such that the form $\sum_{1}^{g+1} c_{k} \alpha_{k}$, where $g=\operatorname{dim} \mathrm{H}=\operatorname{dim} \Omega$ (analytic genus), can be written as $d^{\prime \prime} h$ with $h \in \mathrm{E}^{0}$.
(c) Show that $h$ is holomorphic on $\mathrm{X}-\mathrm{U}$ and in the neighbourhood of $x$. Set $f=h-\eta$. Show that $f$ is holomorphic on $\mathrm{U}-\{x\}$, meromorphic on X and has a pole of order $\geqslant 1$ and $\leqslant g+1$ at $x$.
(d) Let $y \neq x \in \mathrm{X}$. Show that there is a meromorphic function $f$ on X , such that $f(y)=1$ with a pole at $x$. Conclude by considering $\frac{1}{f}$.

### 6.3 Extensions of $\mathbb{C}$ with Transcendence Degree 1

### 6.3.1

Consider the Riemann sphere $\Sigma$ (6.1.2).
Proposition $\underset{\sim}{\text { Let }} \mathrm{X}$ be a Riemann surface. Extend every meromorphic function $f$ on X to a map $\tilde{f}: \mathrm{X} \rightarrow \Sigma$ by assigning the value $\infty$ to the poles of $f$. This gives a bijection $f \mapsto \tilde{f}$ from $\mathscr{M}(\mathrm{X})$ onto the set of analytic maps from X to $\Sigma$, non constant and distinct from $\infty$ on every connected component of X .

Proof Let $f \in \mathscr{M}(\mathrm{X})$ and $x_{0} \in \mathrm{X}$. In the neighbourhood of $x_{0}, f=\frac{g}{h}$ for some holomorphic $g$ and $h$. If necessary simplifying the fraction, $g$ and $h$ can be assumed not to both vanish at $x_{0}$. If $h\left(x_{0}\right) \neq 0$, then the function $\varphi_{0} \circ \tilde{f}=f$ is holomorphic in the neighbourhood of $x_{0}$. If $g\left(x_{0}\right) \neq 0, \varphi_{1} \circ \tilde{f}=\frac{h}{g}$ is holomorphic in the neighbourhood of $x_{0}$. Hence $f$ is analytic.

Let $h: \mathrm{X} \rightarrow \Sigma$ be an analytic map and set $\mathrm{D}=h^{-1}(\infty)$. Assume that there is no connected component of X in D . Then $f=\varphi_{0} \circ h: \mathrm{X}-\mathrm{D} \rightarrow \mathbb{C}$ is meromorphic on X ; indeed, $f$ is holomorphic on $\mathrm{X}-\mathrm{D}$. Let $x_{0} \in \mathrm{D}$; the function $g=\varphi_{1} \circ h$ is holomorphic and is not identically zero in the neighbourhood of $x_{0}$, and $f$ agrees with the meromorphic function $\frac{1}{g}$.

Remark Let $f=\frac{g}{h}$ be a meromorphic function and $x \in \mathrm{X}$. If $g$ has a zero of order $k>0$ at $x$ and $h(x) \neq 0$, then $f(x)=0$; so $\tilde{f}(x)=0$. The ramification index of X over $\Sigma$ at $x$ for $\tilde{f}$ is $k$. Indeed $f=\varphi^{k} u$ for some chart $\varphi$ centered at $x$ and $u(x) \neq 0$. Then $\frac{u}{g}=v^{k}$ in the neighbourhood of $x$ for some holomorphic $v$ such that $v(x) \neq 0$. Then $f=\psi^{k}$, where $\psi=\varphi v$, and $\psi$ is another chart centered at $x$. Similarly, if $g(x) \neq 0$ and $h$ has a zero of order $k>0$, then $\tilde{f}(x)=\infty$ and the ramification index of X at $x$ is $k$.

### 6.3.2 Meromorphic Functions on the Riemann Sphere

Let $\mathrm{P} \in \mathbb{C}[Z]$. The polynomial P defines a function $\mathbb{C} \rightarrow \mathbb{C}$, meromorphic on $\Sigma$. Indeed, it is holomorphic on $\mathbb{C}$, and in the neighbourhood of $\infty, \mathrm{P} \circ \varphi_{1}=$ $\left(z \mapsto \mathrm{P}\left(\frac{1}{z}\right)\right)$ is a rational fraction. More generally, any rational fraction defines a meromorphic function on $\Sigma$.

Theorem Every meromorphic function on the Riemann sphere $\Sigma$ is a rational fraction.

Proof Let U be an open subset of $\Sigma$ and $f$ a meromorphic function on U . Let $a \in \mathrm{U}$ be a pole of $f$. Define the polar part $\mathrm{P}_{a} f$ of $f$ at $a$ as follows: if $a \neq \infty$, then $f$
can be uniquely written as $f(z)=\sum_{k=1}^{N} \frac{c_{k}}{(z-a)^{k}}+h(z)$ with $h$ holomorphic in the neighbourhood of $a$; then the polar part is $\sum_{k=1}^{N} \frac{c_{k}}{(z-a)^{k}}$.

If $a=\infty, f$ can be uniquely written as

$$
f(z)=\sum_{k=1}^{\mathrm{N}} c_{k} z^{k}+h(z)
$$

where $h$ extends to a holomorphic function in the neighbourhood of $\infty$ in $\Sigma$; then the polar part is $\sum_{k=1}^{\mathrm{N}} c_{k} z^{k}$. In both cases, $a$ is the unique pole of the rational fraction $\mathrm{P}_{a} f$, but is not a pole of $f-\mathrm{P}_{a} f$.

Let $f$ be a meromorphic function on $\Sigma$ and $a_{1}, \ldots, a_{n}$ the poles of $f$. Then $h=f-\sum_{1}^{\mathrm{N}} \mathrm{P}_{a_{i}} f$ is holomorphic on $\Sigma$, and necessarily constant. Indeed, $|h|$ reaches its maximum at some point $x$. By the maximum principle, $h$ is constant in its neighbourhood, and hence, by analytic continuation, constant on $\Sigma$. (Liouville's theorem could also be used: every bounded holomorphic function on $\mathbb{C}$ is constant). Thus, $f$ is a rational fraction.

### 6.3.3 Example: Homographies

A rational fraction $f=x \mapsto \frac{a x+b}{c x+d}$, where $a d-b c \neq 0$, defines an analytic map from $\Sigma$ to $\Sigma$ called a homography. All homographies are automorphisms of $\Sigma$. The inverse of $f$ is given by $y \mapsto \frac{-d y+b}{c y-a}$.
Proposition All automorphisms of $\Sigma$ are homographies.
Proof Let $f$ be an automorphism of $\Sigma$. Write $f$ as $\frac{\mathrm{P}}{\mathrm{Q}}$, where P and Q are coprime polynomials in $\mathbb{C}[\mathrm{X}]$. The zeros of P being the zeros of $f, \mathrm{P}$ has at most one zero, which must be simple by Remark 6.3.1. Hence the degree of $P$ is at most 1 . Similarly $f$ can only have one pole, and it is simple. Hence the degree of Q is at most 1 . Thus $f=x \mapsto \frac{a x+b}{c x+d}$ for some $a, b, c, d$ such that $a d-b c \neq 0$, for otherwise $f$ would be constant.

Corollary The automorphism group of $\Sigma$ is generated by the translations $x \mapsto x+b$, the homotheties $x \mapsto$ ax and the map $x \mapsto \frac{1}{x}$.

### 6.3.4 Transcendence degree of $\mathscr{M}(\mathbf{X})$

Proposition Let X be a nonempty connected compact Riemann surface. Then $\mathscr{M}(\mathrm{X})$ is a finitely generated extension of $\mathbb{C}$ with transcendence degree 1 .

Proof Let $f$ be a non-constant meromorphic function on X (such functions exist by 6.2.1). Consider X as a finite analytically ramified covering of $\Sigma$ by $f$ (6.3.1). Then $\mathscr{M}(\mathrm{X})$ is a finitely generated extension of $\mathscr{M}(\Sigma)=\mathbb{C}(\mathrm{Z})$ by (6.2.3).

### 6.3.5

Let $\boldsymbol{U}$ denote the category of connected compact Riemann surfaces and non-constant analytic (hence surjective) maps, and $\mathscr{K}$ the category of finitely generated extensions of $\mathbb{C}$ with transcendence degree 1 .

Theorem The functor $\mathscr{M}: \mathscr{U} \rightarrow \mathscr{K}$ is an anti-equivalence of categories.
Proof (a) The functor $\mathscr{M}$ is essentially surjective. Let $\mathrm{E} \in \mathscr{K}$, then E is isomorphic to a finite extension of $\mathbb{C}(\mathrm{Z})=\mathscr{M}(\Sigma)$ (6.3.4). It is necessarily etale since $\mathbb{C}(\mathrm{Z})$ is of characteristic 0 . So E is isomorphic to an algebra $\mathscr{M}(\mathrm{X})$, where X is a finite analytically ramified covering of $\Sigma$ (6.2.4), which is necessarily connected since E is a field.
(b) The map $f \mapsto f^{*}$ from $\operatorname{Hom}_{\mathscr{U}}(\mathrm{X}, \Sigma)$ to $\operatorname{Hom}_{\mathscr{K}}(\mathscr{M}(\Sigma), \mathscr{M}(\mathrm{X}))$ is bijective. The set $\operatorname{Hom}_{\mathscr{U}}(\mathrm{X}, \Sigma)$ can be identified with the set of non constant meromorphic functions on $\mathbf{X}$, as well as with the set of transcendent functions $f \in$ $\mathscr{M}(\mathrm{X})$ over $\mathbb{C}$ since $\mathbb{C}$ is algebraically closed. The set $\operatorname{Hom}_{\mathscr{K}}(\mathscr{M}(\Sigma), \mathscr{M}(\mathrm{X}))=$ $\operatorname{Hom}_{\mathbb{C}-\mathscr{A} \mathscr{g}_{g}}(\mathbb{C}(\mathrm{Z}), \mathscr{M}(\mathrm{X}))$ can also be identified with the set of transcendental elements of $\mathscr{M}(\mathrm{X})$ over $\mathbb{C}$. The identification of $\operatorname{Hom}_{\mathscr{U}}(\mathrm{X}, \Sigma)$ with $\operatorname{Hom}_{\mathscr{K}}(\mathscr{M}(\Sigma)$, $\mathscr{M}(\mathrm{X}))$ thus obtained can be checked to be given by the map $f \mapsto f^{*}$.
(c) The functor $\mathscr{M}$ is fully faithful. Let X and Y be objects of $\mathscr{U}$ and $h$ : $\mathscr{M}(\mathrm{X}) \rightarrow \mathscr{M}(\mathrm{Y})$ a $\mathbb{C}$-algebra homomorphism. Let $u: \mathrm{X} \rightarrow \Sigma$ be a non constant analytic map. The homomorphism $h \circ u^{*}: \mathscr{M}(\Sigma) \rightarrow \mathscr{M}(\mathrm{Y})$ is of the form $v^{*}$, where $v \in \operatorname{Hom}_{\mathscr{U}}(\mathrm{Y}, \Sigma)$ by (b). If $f: \mathrm{Y} \rightarrow \mathrm{X}$ is a morphism such that $h=f^{*}$, then necessarily $u \circ f=v$. Indeed, $(u \circ f)^{*}=f^{*} \circ u^{*}=h \circ u^{*}=v^{*}$, and so $u \circ f=v$ by (b).

Consider $\mathscr{M}(\mathrm{X})$ and $\mathscr{M}(\mathrm{Y})$ as algebras over $\mathscr{M}(\Sigma)$ using $u^{*}$ and $v^{*}$. Then $h$ is an algebra homomorphism. By Theorem 6.2.4, there is a unique $\Sigma$-morphism $f: \mathrm{Y} \rightarrow \mathrm{X}$ such that $f^{*}=h$.

Corollary The category of compact Riemann surfaces and analytic maps non constant on the connected components is anti-equivalent by $\mathscr{M}$ to the category of $\mathbb{C}$ algebras that are finite products of finitely generated extensions of $\mathbb{C}$ with transcendence degree 1 .

## Exercises 6.3. (Extensions of $\mathbb{C}$ with transcendence degree 1)

1.-Embed the space $X$ of 4.5 , Exercise 3 in $\mathbb{C}$ in such a way that $A$ and $B$ go round the points 1 and -1 respectively.

Give the inverse image of X under the map $x \mapsto \frac{-x^{3}+3 x}{2}$.

Interpret the analogy between 4.5, Exercise 3 and 5.6 and 5.7, Exercise 2.
2.-Changing 5.6 and 5.7, Exercise 3, accordingly, construct an example of a Galois covering Z of a Galois covering Y of a space X with Galois covering Z .
3. (Algebraic definition of ramification)— Let $E$ be an extension of $\mathbb{C}$ with transcendence degree 1 . Any map $\varphi: \mathrm{E} \rightarrow \Sigma=\mathbb{C} \cup\{\infty\}$ such that

$$
\varphi(x+y)=\varphi(x)+\varphi(y) \text { if }(\varphi(x), \varphi(y)) \neq(\infty, \infty)
$$

and

$$
\varphi(x y)=\varphi(x) \cdot \varphi(y) \text { if }(\varphi(x), \varphi(y)) \neq(0, \infty) \text { and } \neq(\infty, 0)
$$

is called a place of E .
(a) Show that if $\mathrm{E}=\mathscr{M}(\mathrm{X})$ then, for all $x \in \mathrm{X}$ the map $\delta_{x}: f \mapsto f(x)$ is a place of E , and that this gives a bijection from X onto the set of places of E . (Consider first the case $\mathrm{E}=\mathbb{C}(\mathbb{Z})=\mathscr{M}(\Sigma)$.)
(b) Show that if E is finitely generated, then the set X of places of E can be equipped with the structure of a compact Riemann surface identifying E with $\mathscr{M}(\mathrm{X})$.
(c) Let $\varphi$ be a place of E , and $\mathfrak{D}_{\varphi}$ the subring of E consisting of $f \in \mathrm{E}$ such that $\varphi(f) \neq \infty$. Set $\mathfrak{m}_{\varphi}=\varphi^{-1}(0)$. Show that $\mathfrak{m}_{\varphi}$ is the unique maximal ideal of $\mathfrak{D}_{\varphi}$. Show that, if E is finitely generated, then $\mathfrak{D}_{\varphi}$ is a principal ring. Describe the ideals of $\mathfrak{D}_{\varphi}$.
(d) Let X and Y be connected compact Riemann surfaces, $h: \mathrm{Y} \rightarrow \mathrm{X}$ a surjective analytic map, $y \in \mathrm{Y}$ and $x=h(y)$. Set $\mathrm{E}=\mathscr{M}(\mathrm{X}), \mathrm{F}=\mathscr{M}(\mathrm{Y})$, so that E can be identified with a subextension of $\mathrm{F}, \mathfrak{D}_{x}=\mathfrak{D}_{\delta_{x}} \subset \mathrm{E}$; define $\mathfrak{D}_{y}, \mathfrak{m}_{x}$ and $\mathfrak{m}_{y}$ similarly. Show that $\mathfrak{m}_{x} \mathfrak{D}_{y}=\mathfrak{m}_{y}{ }^{k}$, where $k$ is the ramification index of Y over X at $y$.

### 6.4 Determination of Some Galois Groups

### 6.4.1 Free Profinite Groups

For a finite set J , let $\mathrm{L}(\mathrm{J})$ be the free group on $\mathrm{J}(4.7 .4)$ and $\left(\mathrm{X}_{i}\right)_{i \in \mathrm{~J}}$ the canonical basis for $L(J)$. Let $\widehat{L}(J)$ be the profinite completion of $L(J)$ (2.9.5). For $J \subset J^{\prime}$, define $\rho_{J}^{J^{\prime}}$ : $\mathrm{L}\left(\mathbf{J}^{\prime}\right) \rightarrow \mathrm{L}(\mathrm{J})$ by $\rho_{\mathrm{J}}^{\mathrm{J}^{\prime}}\left(\mathrm{X}_{i}\right)=\mathrm{X}_{i}$ if $i \in \mathrm{~J}$ and $\rho_{\mathrm{J}}^{\mathrm{J}^{\prime}}\left(\mathrm{X}_{i}\right)=e$ if $i \notin \mathrm{~J}$. Passing to profinite completions gives a continuous homomorphism $\widehat{\rho}_{\mathrm{J}}^{\prime}: \widehat{\mathrm{L}}\left(\mathrm{J}^{\prime}\right) \rightarrow \widehat{\mathrm{L}}(\mathrm{J})$.

For an arbitrary set $I$, the inverse limit $\widehat{L}(I)$ of $\widehat{L}(J)$ for finite $J$ in $I$ is called the free profinite group on I.

For finite $\mathrm{J} \subset \mathrm{I}$, define the elements $\left(\mathrm{X}_{i}^{\mathrm{J}}\right) \in \widehat{\mathrm{L}}(\mathrm{J}), i \in \mathrm{I}$, by taking $\mathrm{X}_{i}^{\mathrm{J}}$ to be the image of $\mathrm{X}_{i} \in \mathrm{~L}(\mathrm{~J})$ if $i \in \mathrm{~J}$ and $e$ if $i \notin \mathrm{~J}$. Passing to the inverse limit gives elements $\left(\widehat{\mathrm{X}}_{i}\right) \in \widehat{\mathrm{L}}(\mathrm{I}), i \in \mathrm{I}$.

Let G be a profinite group. A family $\left(x_{i}\right)_{i \in \mathrm{I}}$ of elements of G tends to $e$ if for all neighbourhoods U of $e$ in G , the set $\left\{i \mid x_{i} \notin \mathrm{U}\right\}$ is finite. Equivalently, if $\mathrm{G}=\underset{\leftrightarrows}{\lim } \mathrm{G}_{\lambda}$,
where all $\mathrm{G}_{\lambda}$ are finite groups, then the family $\left(\pi_{\lambda}\left(x_{i}\right)\right)_{i \in \mathrm{I}}$ has finite support for all $\lambda, \pi_{\lambda}: \mathrm{G} \rightarrow \mathrm{G}_{\lambda}$ being the canonical map.

## Proposition Let I be a set.

(a) The covariant factor $\Phi: \mathscr{E}_{\nu \mathcal{A}} \rightarrow \mathscr{E} n s$ which assigns to every profinite group G the set of families $\left(x_{i}\right)_{i \in \mathrm{I}}$ of elements of G tending to $e$ is represented by $\widehat{\mathrm{L}}(\mathrm{I})$ equipped with ( $\widehat{\mathrm{X}}_{i}$ ).
(b) Let $\mathrm{L}^{\prime}$ be a profinite group and $\left(\mathrm{X}_{i}^{\prime}\right)$ a family of elements of $\mathrm{L}^{\prime}$ tending to $e$. Assume that for all finite groups G and families $\left(x_{i}\right)_{i \in \mathrm{I}}$ of elements of G with finite support, there is a unique continuous homomorphism $f: \mathrm{L}^{\prime} \rightarrow \mathrm{G}$ such that $f\left(\mathrm{X}_{i}^{\prime}\right)=$ $x_{i}$ for all $i$. Then there is a unique isomorphism from $\widehat{\mathrm{L}}(\mathrm{I})$ onto $\mathrm{L}^{\prime}$ transforming $\widehat{\mathrm{X}}_{i}$ into $\mathrm{X}_{i}^{\prime}$.

Proof (a) Under the assumptions of (b), the profinite group $\mathrm{L}^{\prime}$ equipped of $\left(\mathrm{X}_{i}^{\prime}\right)$ represents $\Phi$. Indeed, if $\mathrm{G}=\operatorname{\operatorname {lim}} \mathrm{G}_{\lambda}$, where all $\mathrm{G}_{\lambda}$ are finite groups and the family $\left(x_{i}\right)_{i \in \mathrm{I}}$ tends to $e$ in G , then for all $\lambda$ the family $\left(\pi_{\lambda}\left(x_{i}\right)\right)_{i \in \mathrm{I}}$ has finite support in $\mathrm{G}_{\lambda}$; hence there is a unique $f_{\lambda}: \mathrm{L}^{\prime} \rightarrow \mathrm{G}_{\lambda}$ such that $f_{\lambda}\left(\mathrm{X}_{i}^{\prime}\right)=\pi_{\lambda}\left(x_{i}\right)$. Passing to the inverse limit gives a unique $f: \mathrm{L}^{\prime} \rightarrow \mathrm{G}$ such that $f\left(\mathrm{X}_{i}^{\prime}\right)=x_{i}$.

The profinite group $\widehat{\mathrm{L}}(\mathrm{I})$ equipped with $\left(\widehat{\mathrm{X}}_{i}\right)_{i \in \mathrm{I}}$ satisfies the assumptions of $(\mathrm{b})$. First, the family $\left(\widehat{\mathrm{X}}_{i}\right)$ is readily seen to tend to $e$ in $\widehat{\mathrm{L}}(\mathrm{I})$. Let G be a finite group and $\left(x_{i}\right)_{i \in \mathrm{I}}$ a family of elements of G with finite support J . The homomorphism $u: \mathrm{L}(\mathrm{J}) \rightarrow \mathrm{G}$ defined by $u\left(\mathrm{X}_{i}\right)=x_{i}$ for $i \in \mathrm{~J}$ extends to profinite completions in $\hat{u}: \widehat{\mathrm{L}}(\mathrm{J}) \rightarrow \widehat{\mathrm{G}}=\mathrm{G}$, and $f=\hat{u} \circ \widehat{\rho}_{\mathrm{J}}: \widehat{\mathrm{L}}(\mathrm{I}) \rightarrow \mathrm{G}$ satisfies $f\left(\widehat{\mathrm{X}}_{i}\right)=x_{i}$.
(b) If $f^{\prime}: \widehat{\mathrm{L}}(\mathrm{I}) \rightarrow \mathrm{G}$ satisfies $f^{\prime}\left(\widehat{\mathrm{X}}_{i}\right)=x_{i}$ for $i \in \mathrm{I}$, then $f^{\prime}$ factorizes through some $\widehat{\mathrm{L}}\left(\mathrm{J}^{\prime}\right)$. Indeed $\operatorname{Ker} f^{\prime}$ is a neighbourhood of $e$ in $\widehat{\mathrm{L}}(\mathrm{I})$, and so contains a set $\widehat{\rho}_{\mathrm{J}}{ }^{-1}(\mathrm{U})$, where U is a neighbourhood of $e$ in $\widehat{\mathrm{L}}\left(\mathrm{J}^{\prime}\right)$, and hence necessarily $\operatorname{Ker} \widehat{\rho}_{J^{\prime}}$. For $i \notin \mathbf{J}^{\prime}, \mathrm{X}_{i}^{J^{\prime}}=e$. So $x_{i}=e$ and $i \notin \mathbf{J}$; thus $\mathbf{J} \subset \mathbf{J}^{\prime}$. It follows that $f^{\prime}=f$.

The essential point of this section is to prove that if $\Omega$ is an algebraic closure of $\mathbb{C}(\mathrm{Z})$, then the Galois group $\operatorname{Aut}_{\mathbb{C}(\mathrm{Z})}(\Omega)$ is isomorphic to the free profinite group on the set $\mathbb{C}$.

### 6.4.2 The Homomorphism $\boldsymbol{\theta}_{x}$

Let B be a connected compact Riemann surface, $\Delta$ a finite subset of B and $b \in \mathrm{~B}-\Delta$. Let E be a finite Galois extension of $\mathscr{M}(\mathrm{B})$ such that $\Delta_{\mathrm{E}} \subset \Delta$ and X is an analytically ramified covering of B equipped with an isomorphism $\iota: \mathscr{M}(\mathrm{X}) \rightarrow \mathrm{E}$ identifying these two algebras. Let $x \in \mathrm{X}(b)$.

Define a homomorphism $\theta_{x}: \pi_{1}(\mathrm{~B}-\Delta, b) \rightarrow \mathrm{Aut}_{\mu_{(\mathrm{B})} \mathrm{E} \text { as follows: the fun- }}$ damental group $\pi_{1}(\mathrm{~B}-\Delta, b)$ acts on $\mathrm{X}(b)$ since $\left.\mathrm{X}\right|_{\mathrm{B}-\Delta}$ is a covering. Let $\gamma \in$ $\pi_{1}(\mathrm{~B}-\Delta, b) ; \operatorname{as~}^{\operatorname{Aut}_{\mathrm{B}}}(\mathrm{X})$ acts simply transitively on $\mathrm{X}(b)$, there is a unique automorphism $g$ of X such that $g(x)=\gamma \cdot x$. Let $\theta_{x}(\gamma)$ be the automorphism $g^{*}$ of $\mathrm{E}=\mathscr{M}(\mathrm{X})$. For $\gamma^{\prime} \in \pi_{1}(\mathrm{~B}-\Delta, b)$, let $g^{\prime}$ be the automorphism of X corresponding to $\gamma^{\prime}$ and $g^{\prime \prime}$
the one corresponding to $\gamma \cdot \gamma^{\prime}$. Then $g^{\prime \prime}(x)=\gamma \cdot \gamma^{\prime} \cdot x=\gamma \cdot g^{\prime}(x)=g^{\prime}(\gamma \cdot x)=$ $g^{\prime}(g(x))$. So $g^{\prime \prime}=g^{\prime} \circ g$ and $\theta_{x}\left(\gamma \cdot \gamma^{\prime}\right)=g^{\prime * *}=g^{*} \circ g^{\prime *}=\theta_{x}(\gamma) \circ \theta_{x}\left(\gamma^{\prime}\right)$.

Remarks (1) The homomorphism $\theta_{x}$ is surjective and its kernel is the stabilizer of $x$ in $\pi_{1}(\mathrm{~B}-\Delta, b)$; it is normal and of finite index.
(2) The homomorphism $\theta_{x}$ depends on the choice of $x$.

## Functoriality

Let E and F be finite Galois extensions of $\mathscr{M}(\mathrm{B})$ such that $\Delta_{\mathrm{E}} \subset \Delta$ and $\Delta_{\mathrm{F}} \subset \Delta$. Let X and Y be analytically ramified coverings of B equipped with isomorphisms identifying $\mathscr{M}(\mathrm{X})$ with E and $\mathscr{M}(\mathrm{Y})$ with F . Let $f: \mathrm{E} \rightarrow \mathrm{F}$ be a homomorphism. There is a unique B -morphism $u: \mathrm{Y} \rightarrow \mathrm{X}$ such that $u^{*}=f$.

Proposition Choose $y \in \mathrm{Y}(b)$ and set $x=u(y)$. Then $f$ is compatible with the actions defined by $\theta_{x}$ and $\theta_{y}$; in other words, for $t \in \mathrm{E}$ and $\gamma \in \pi_{1}(\mathrm{~B}-\Delta, b)$, $f\left(\theta_{x}(\gamma, t)\right)=\theta_{y}(\gamma, f(t))$.

Proof Let $\gamma \in \pi_{1}(\mathrm{~B}-\Delta, b)$, and take $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{X}), g_{1} \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{Y})$ such that $g(x)=$ $\gamma \cdot x$ and $g_{1}(y)=\gamma \cdot y$. Then

$$
g(u(y))=g(x)=\gamma \cdot x=\gamma \cdot u(y)=u(\gamma \cdot y)=u\left(g_{1}(y)\right) .
$$

As F in an extension, $\left.\mathrm{Y}\right|_{\mathrm{B}-\Delta}$ is a connected covering, and by (4.3.5) the diagram

is commutative. Applying the functor $\mathscr{M}$ gives a commutative diagram


### 6.4.3 Galois Group of $\Omega_{\Delta}$

Let B be a connected Riemann surface, $\Delta$ a finite subset of $\mathrm{B}, b \in \mathrm{~B}-\Delta$ and $\Omega$ an algebraic closure of $\mathscr{M}(\mathrm{B})$. Denote by $\Omega_{\Delta}$ the union of finite sub-extensions E of $\Omega$ such that $\Delta_{\mathrm{E}} \subset \Delta$. By 6.2 .11 , it is a directed union, and so $\Omega_{\Delta}$ is a sub-extension of $\Omega$. It is obviously preserved by automorphisms of $\Omega$, and separable since the
characteristic is 0 . Hence $\Omega_{\Delta}$ is a Galois extension, usually infinite, of $\mathscr{M}(\mathrm{B})$. By 6.2.11, for any finite extension F of $\mathscr{M}(\mathrm{B})$ in $\Omega_{\Delta}, \Delta_{\mathrm{F}} \subset \Delta$.

Let $\left(\mathrm{E}_{i}\right)$ be a directed family of finite Galois extensions of $\mathscr{M}(\mathrm{B})$ such that $\bigcup \mathrm{E}_{i}=$ $\Omega_{\Delta}$. For each $i$, identify $\mathrm{E}_{i}$ with an algebra $\mathscr{M}\left(\mathrm{X}_{i}\right)$, where $\mathrm{X}_{i}$ is a finite analytically ramified covering of B . For $i \leqslant j, \mathrm{E}_{i} \subset \mathrm{E}_{j}$; so there is a B -morphism $p_{i}^{j}: \mathrm{X}_{j} \rightarrow \mathrm{X}_{i}$ such that $\left(p_{i}^{j}\right)^{*}$ is a canonical injection of $\mathrm{E}_{i}$ in $\mathrm{E}_{j}$. These $\mathrm{X}_{i}$ and $p_{i}^{j}$ form a projective system of ramified coverings of $B$.

By Tychonoff's theorem, $\lim X_{i}(b) \neq$. Let $x \in \lim X_{i}(b)$, and for all $i$ let $p_{i}(x)$ be the image of $x$ in $X_{i}(b)$. Passing to the inverse limit starting from $\theta_{p_{i}(x)}$ gives a homomorphism $\theta_{x}: \pi_{1}(\mathrm{~B}-\Delta, b) \rightarrow \operatorname{Aut}_{\mu(\mathrm{B})} \Omega_{\Delta}=\lim _{\mathrm{Aut}_{\mathscr{M}(\mathrm{B})} \mathrm{E}_{i} \text {. As the Galois group }}$ Aut ${ }_{\mu(\mathrm{B})} \Omega_{\Delta}$ is profinite, the homomorphism $\theta_{x}$ induces a continuous homomorphism $\widehat{\theta}_{x}$ of the profinite completion of $\pi_{1}(\mathrm{~B}-\Delta, b)$ in $\mathrm{Aut}_{\mu_{(\mathrm{B})}} \Omega_{\Delta}$.

Proposition The homomorphism $\widehat{\theta_{x}}: \widehat{\pi_{1}}(\mathrm{~B}-\Delta, b) \rightarrow$ Aut $_{M_{(\mathrm{B})} \Omega_{\Delta}}$ is an isomorphism.

Proof The family ( $\mathrm{E}_{i}$ ) may be assumed to consist of all the finite Galois extensions in $\Omega_{\Delta}$. For all $i, \theta_{p_{i}(x)}$ gives an isomorphism $h_{i}: \mathrm{G} / \mathrm{N}_{i} \rightarrow$ Aut ${ }_{\mu(\mathrm{B})} \mathrm{E}_{i}$, where $\mathrm{G}=$ $\pi_{1}(\mathrm{~B}-\Delta, b)$ and $\mathrm{N}_{i}$ is the stabilizer of $p_{i}(x)$ in G. Passing to the projective limit, we get an isomorphism $h: \underset{\leftarrow}{\lim } \mathrm{G} / \mathrm{N}_{i} \rightarrow$ Aut $\Omega_{\Delta}$.

Let N be a normal subgroup of finite index in G . Then there is a finite Galois covering Y of $\mathrm{B}-\Delta$ and $y \in \mathrm{Y}(b)$ with stabilizer N (4.6.11 and 4.6.10). Y can be completed to an analytically ramified covering $\widehat{\mathrm{Y}}$ of B . Then $\mathscr{M}(\widehat{\mathrm{Y}})$ is a Galois extension of $\mathscr{M}(\mathrm{B})$, which can be embedded in $\Omega$. It then is in $\Omega_{\Delta}$, and thus can be identified with some $\mathrm{E}_{i}$, and $\mathrm{N}=\mathrm{N}_{i}$.

Therefore $\lim \mathrm{G} / \mathrm{N}_{i}$ is a profinite completion of G , and $h=\widehat{\theta_{x}}$; so $h$ is an isomorphism.

### 6.4.4 Fundamental Groups with Respect to Germs of Paths

We now give a description of

$$
\operatorname{Aut}_{\mathscr{M}(\mathrm{B})} \Omega=\lim _{\leftarrow} \mathrm{Aut}_{\mathscr{M}(\mathrm{B})} \Omega_{\Delta}
$$

the inverse limit being taken over all finite $\Delta$ in B . The problem lies in the choice of the basepoint. Indeed, fixing $b \in \mathrm{~B}$ prevents us from considering the sets $\Delta$ containing $b$. To overcome this problem we replace the basepoint by a germ of path B .

Let $\beta:[0,1] \rightarrow \mathrm{B}$ be an injective continuous map. For all finite $\Delta \subset \mathrm{B}$, let $c_{\Delta}$ be the largest $c$ such that $\beta(] 0, c[) \subset \mathrm{B}-\Delta$. Denote by $\pi_{1}(\mathrm{~B}-\Delta, \beta)$ the quotient of the set of pairs $(t, \gamma)$ where $t \in] 0, c_{\Delta}\left[\right.$ and $\gamma \in \pi_{1}(\mathrm{~B}-\Delta, \beta(t))$ by the equivalence relation identifying $\left(t_{1}, \gamma_{1}\right)$ to $\left(t_{2}, \gamma_{2}\right)$ if $\gamma_{2}$ is the image of $\gamma_{1}$ by the isomorphism from $\pi_{1}\left(\mathrm{~B}-\Delta, \beta\left(t_{1}\right)\right)$ onto $\pi_{1}\left(\mathrm{~B}-\Delta, \beta\left(t_{2}\right)\right)$ defined by $\left.\beta\right|_{\left[t_{1}, t_{2}\right]}$ (4.6.6). It has a natural group structure, and for all $t \in] 0, c_{\Delta}$ [ there is a canonical isomorphism $\chi_{t}$
from $\pi_{1}(\mathrm{~B}-\Delta, \beta(t))$ onto $\pi_{1}(\mathrm{~B}-\Delta, \beta)$. Besides the group $\pi_{1}(\mathrm{~B}-\Delta, \beta)$ is the (trivial) direct limit of $\pi_{1}(\mathrm{~B}-\Delta, \beta(t))$.

Let $X$ be a finite analytically ramified covering of $B$ and for $\Delta$ take its ramification set. Denote by $\Gamma(\beta, X)$ the set of continuous coverings of $\left.\beta\right|_{] 0, c_{\Delta}[ }$ in X . If X is of degree $d$, then it is a set with $d$ elements since $\left.\beta^{*} \mathrm{X}\right|_{] 0, c_{\Delta}[ }$ is a trivial $d$-fold covering.
Proposition Assume X is Galois. Let $\xi \in \Gamma(\beta, \mathrm{X})$. There is a unique homomorphism $\theta_{\xi}: \pi_{1}(\mathrm{~B}-\Delta, \beta) \rightarrow \mathrm{Aut}_{\mathscr{M}(\mathrm{B})} \mathscr{M}^{(\mathrm{X})}$ such that for all $\left.t \in\right] 0, c_{\Delta}\left[, \theta_{\xi} \circ \chi_{t}=\theta_{\xi(t)}\right.$.
Proof It suffices to show that, for $t$ and $t^{\prime}$ in $] 0, c_{\Delta}\left[\right.$ such that $t<t^{\prime}$, there is a commutative diagram

where $\beta_{*}(\gamma)=\left.\left.\beta\right|_{\left[t, t^{\prime}\right]} \cdot \gamma \cdot \beta\right|_{\left[t^{\prime}, t\right]}$, denoting by $\left.\beta\right|_{\left[t^{\prime}, t\right]}$ the inverse path of $\left.\beta\right|_{\left[t, t^{\prime}\right]}$. Let $\gamma \in \pi_{1}(\mathrm{~B}-\Delta, \beta(t))$ and $\gamma^{\prime}=\beta_{*}(\gamma)$. Let $\eta$ be a path with initial point $\xi(t)$ lifting a loop from the class $\gamma$. The endpoint of $\eta$ is $\gamma \cdot \xi(t)$. Let $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{X})$ be such that $g(\xi(t))=\gamma \cdot \xi(t)$. Then the path $\eta^{\prime}=\left.g\left(\left.\xi\right|_{\left[t, t^{\prime}\right]}\right) \cdot \eta \cdot \xi\right|_{\left[t^{\prime}, t\right]}$ is a path lifting a path from the class $\gamma^{\prime}$, with initial point $\xi\left(t^{\prime}\right)$ and endpoint $g\left(\xi\left(t^{\prime}\right)\right)$; so $\gamma^{\prime} \cdot \xi\left(t^{\prime}\right)=$ $g\left(\xi\left(t^{\prime}\right)\right)$. Hence the same $g \in \operatorname{Aut}_{\mathrm{B}}(\mathrm{X})$ maps $\xi(t)$ onto $\gamma \cdot \xi(t)$ and $\xi\left(t^{\prime}\right)$ onto $\gamma^{\prime} \cdot \xi\left(t^{\prime}\right)$. Therefore $\theta_{\xi(t)}(\gamma)=\theta_{\xi\left(t^{\prime}\right)}\left(\gamma^{\prime}\right)$.

## Functoriality

Let X and $\mathrm{X}^{\prime}$ be finite analytically ramified Galois coverings, $\Delta$ and $\Delta^{\prime}$ their ramification sets, and $f: \mathrm{X}^{\prime} \rightarrow \mathrm{X}$ a B-morphism. The morphism $f$ is necessarily surjective, $\Delta \subset \Delta^{\prime}$ and $c_{\Delta^{\prime}} \leqslant c_{\Delta}$. For all $\xi^{\prime} \in \Gamma\left(\beta, \mathrm{X}^{\prime}\right)$ there is a unique $\xi \in \Gamma(\beta, \mathrm{X})$ such that $\left.\xi\right|_{] 0, c_{\Delta^{\prime}}[ }=f \circ \xi^{\prime}$. This element $\xi$ is written $f_{*}\left(\xi^{\prime}\right)$.

If $\xi=f_{*}\left(\xi^{\prime}\right)$, then the homomorphism $f^{*}: \mathscr{M}(\mathrm{X}) \rightarrow \mathscr{M}\left(\mathrm{X}^{\prime}\right)$ is compatible with the actions of $\pi_{1}(\mathrm{~B}-\Delta, \beta)$ defined by $\theta_{\xi}$ and $\theta_{\xi^{\prime}}$. This is the functoriality Proposition 6.4.2 applied to $\xi(t)$ and $\xi^{\prime}(t)$ for $\left.t \in\right] 0, c_{\Delta}[$.

### 6.4.5 Galois Group of the Algebraic Closure of $\mathscr{M}(\mathbf{B})$

If $\Delta$ and $\Delta^{\prime}$ are finite subsets of B such that $\Delta \subset \Delta^{\prime}$, then there is a unique homomorphism $r_{\Delta}^{\Delta^{\prime}}: \pi_{1}\left(\mathrm{~B}-\Delta^{\prime}, \beta\right) \rightarrow \pi_{1}(\mathrm{~B}-\Delta, \beta)$, making the diagram

commutative for all $t \in] 0, c_{\Delta}[$. As $\Delta$ runs through the directed set of finite subsets of $\mathrm{B}, \pi_{1}(\mathrm{~B}-\Delta, \beta)$ form an inverse system. This is also the case of their profinite completions $\widehat{\pi}_{1}(\mathrm{~B}-\Delta, \beta)$.

Let $\Omega$ be an algebraic closure of $\mathscr{M}(\mathrm{B})$ and $\left(\mathrm{E}_{i}\right)$ an increasing family including all finite Galois sub-extensions of $\Omega$. Set $\Delta_{i}=\Delta_{\mathrm{E}_{i}}$. Identify each $\mathrm{E}_{i}$ with an extension $\mathscr{M}\left(\mathrm{X}_{i}\right)$, where $\mathrm{X}_{i}$ is a finite analytically ramified covering of B. By Tychonoff's theorem (1.7.6 and 2.5.14), $\Gamma=\lim _{\leftarrow} \Gamma\left(\beta, \mathbf{X}_{i}\right) \neq \varnothing$. Let $\xi=\left(\xi_{i}\right) \in \Gamma$. For each $i$, the homomorphism

$$
\theta_{\xi_{i}}: \pi_{1}\left(\mathrm{~B}-\Delta_{i}, \beta\right) \rightarrow \operatorname{Aut}_{\mu(\mathrm{B})} \mathrm{E}_{i}
$$

induces

$$
\widehat{\theta_{\xi_{i}}}: \widehat{\pi_{1}}\left(\mathrm{~B}-\Delta_{i}, \beta\right) \rightarrow \operatorname{Aut}_{M(\mathrm{~B})} \mathrm{E}_{i} .
$$

For $i \leqslant j$, the diagram

$$
\begin{gathered}
\widehat{\pi_{1}}\left(\mathrm{~B}-\Delta_{j}, \beta\right) \xrightarrow{\widehat{\theta_{\xi_{j}}}} \operatorname{Aut}_{\mu(\mathrm{B})} \mathrm{E}_{j} \\
\begin{array}{r}
r_{\Delta_{i}}^{\Delta_{j}} \downarrow
\end{array} \\
\widehat{\pi_{1}}\left(\mathrm{~B}-\Delta_{i}, \beta\right) \xrightarrow{\widehat{\theta_{\xi_{i}}}} \underset{\downarrow}{ } \mathrm{Aut}_{\mu(\mathrm{B})} \mathrm{E}_{i}
\end{gathered}
$$

where $\rho$ is the restriction, is commutative by the functoriality seen in 6.4.4,.
Theorem The homomorphism $\widehat{\theta_{\xi}}: \lim _{\leftrightarrows} \widehat{\pi}_{1}(\mathrm{~B}-\Delta, \beta) \rightarrow \mathrm{Aut}_{M_{(\mathrm{B})} \Omega \text { induced by } \theta_{\xi_{i}}}$ by passing to the inverse limit is an isomorphism.

Proof For each $\Delta$, let $\xi_{\Delta}$ be the image of $\xi$ in the inverse limit of $\Gamma\left(\beta, \mathrm{X}_{i}\right)$ for all $i$ such that $\Delta_{i} \subset \Delta$. There is an isomorphism $\widehat{\theta_{\xi_{\Delta}}}: \widehat{\pi_{1}}(\mathrm{~B}-\Delta, \beta) \rightarrow$ Aut $\mu_{(\mathrm{B})} \Omega_{\Delta}$. This can be seen by identifying $\pi_{1}(\mathrm{~B}-\Delta, \beta)$ with $\pi_{1}(\mathrm{~B}-\Delta, \beta(t))$ for some $\left.t \in\right] 0, c_{\Delta}[$, and by applying Proposition 6.4.3. The theorem follows by passing to the inverse limit.

### 6.4.6 Galois Group of the Algebraic Closure of $\mathbb{C}(\mathbf{Z})$

We now apply the previous theorem to the case when B is the Riemann sphere $\Sigma$. For all $\Delta \subset \mathbb{C}$, set $\bar{\Delta}=\Delta \cup\{\infty\} \subset \Sigma$. These $\bar{\Delta}$ form a cofinal system of the finite subsets of $\Sigma$. Hence

$$
\operatorname{Aut}_{\mathbb{C}(\mathrm{Z})} \Omega=\operatorname{Aut}_{\mu(\Sigma)} \Omega \approx \lim _{\leftarrow} \widehat{\pi}_{1}(\Sigma-\bar{\Delta}, \beta)=\lim _{\leftarrow} \widehat{\pi}_{1}(\mathbb{C}-\Delta, \beta)
$$

where $\beta$ is a continuous injective map from $[0,1]$ to $\mathbb{C}$.
We may assume that the image of $\beta$ only meets each real line of $\mathbb{C}$ at finitely many points. This is for example the case when $\beta$ describes an arc of a circle.

Proposition The profinite group $\lim \widehat{\pi}_{1}(\mathbb{C}-\Delta, \beta)$ is isomorphic to the free profinite group $\widehat{\mathrm{L}}(\mathbb{C})$ on the set $\mathbb{C}$.


Proof For finite $\Delta \subset \mathbb{C}$ and $b \in \mathbb{C}-\Delta$, the group $\pi_{1}(\mathbb{C}-\Delta, b)$ is isomorphic to the free group $\mathrm{L}(\Delta)$ on $\Delta$. Likewise for $\pi_{1}(\mathbb{C}-\Delta, \beta)$. We describe such an isomorphism, i.e. we choose a basis for $\pi_{1}(\mathbb{C}-\Delta, \beta)$. Let $b \in \mathbb{C}$ be such that $b, a_{1}, a_{2}$ are never colinear for $a_{1} \in \Delta, a_{2} \in \Delta, a_{1} \neq a_{2}$. Then define $\gamma_{b, a} \in \pi_{1}(\mathbb{C}-\Delta, b)$ for all $a \in \Delta$ as follows: $\gamma_{b, a}$ is the class of a loop consisting of going along a straight line from $b$ to $a$, stopping at a distance $\varepsilon$ from $a$, then describing an anticlockwise circle of radius $\varepsilon$ around $a$ and returning to $b$ along a straight line, $\varepsilon$ being chosen to be strictly smaller than $\left|a^{\prime}-a\right|$ for all $a^{\prime} \in \Delta-\{a\}$ (Fig. 1). Then $\left(\gamma_{b, a}\right)_{a \in \Delta}$ form a basis for $\pi_{1}(\mathbb{C}-\Delta, b)$.

Next, or all $\Delta$, define a basis $\left(\gamma_{\beta, a}\right)_{a \in \Delta}$ for $\pi_{1}(\mathbb{C}-\Delta, \beta)$. Given the assumption on $\beta$, there exists nonzero $c^{\prime} \leqslant c_{\Delta}$ such that $\beta(] 0, c^{\prime}[)$ never meets any line connecting two points of $\Delta$. Let $c_{\Delta}^{\prime}$ to be the greatest among all $c^{\prime}$. As long as $t$ is in $] 0, c_{\Delta}^{\prime}[$, $\chi_{t}\left(\gamma_{\beta(t), a}\right) \in \pi_{1}(\mathbb{C}-\Delta, \beta)$ is independent of the choice of $t$; we write it $\gamma_{\beta, a}$ (Fig.2).

For $\Delta^{\prime} \supset \Delta$, the homomorphism $r_{\Delta}^{\Delta^{\prime}}: \pi_{1}\left(\mathbb{C}-\Delta^{\prime}, \beta\right) \rightarrow \pi_{1}(\mathbb{C}-\Delta, \beta)$ maps $\gamma_{\beta, a} \in \pi_{1}\left(\mathbb{C}-\Delta^{\prime}, \beta\right)$ to $\gamma_{\beta, a} \in \pi_{1}(\mathbb{C}-\Delta, \beta)$ if $a \in \Delta$ and to $e$ if $a \notin \Delta$. In other words, there is a commutative diagram

where $\rho_{\Delta}^{\Delta^{\prime}}$ is the homomorphism defined in 6.4.1 and $\alpha_{\Delta}$ the isomorphism defined by the basis $\left(\gamma_{\beta, a}\right)_{a \in \Delta}$. Passing to profinite limits and then to the projective limit over $\Delta$ gives an isomorphism $\hat{\alpha}: \widehat{\mathrm{L}}(\mathbb{C}) \rightarrow \lim _{\longleftarrow} \widehat{\pi}_{1}(\mathbb{C}-\Delta, \beta)$.

Corollary 6.3 The Galois group $\mathrm{Aut}_{\mathbb{C}(\mathrm{Z})} \Omega$ is isomorphic to $\widehat{\mathrm{L}}(\mathbb{C})$.
Corollary 6.4 Every finite groups is isomorphic to the automorphism group of a finite extension of $\mathbb{C}(Z)$.

## Exercises 6.4. (Some Galois groups)

1.-(a) Embed $k_{0}=\mathbb{C}(Z)$ in $k=\mathbb{C}(X)$ (which is isomorphic to it) by $Z=X^{3}-3 X$.

Show that $k$ is not a Galois extension. Consider the set $\Gamma$ of $(z, y) \in \mathbb{C}^{2}$ such that the equation $\mathrm{X}^{3}-3 \mathrm{X}=z$ has roots $x^{\prime}$ and $x^{\prime \prime}$ with $y=2 x^{\prime}-x^{\prime \prime}$, equipped with the projection $\pi:(z, y) \mapsto z$. Show that $\Gamma$ can be completed to a ramified covering $\widehat{\Gamma}$ of $\Sigma$ and that $\mathscr{M}(\widehat{\Gamma})$ is the Galois closure of $k$ over $k_{0}$.
(b) Show that there is an involutive automorphism $\tau$ of $\widehat{\Gamma}$ whose fixed points are $a_{1}=(2,5)$ and $a_{2}=(-2,-5)$.
(c) Let $e \in \Gamma$ be a point over $0, \alpha_{1}$ and $\alpha_{2}$ paths from $e$ to $a_{1}$ and $a_{2}$ respectively, and avoiding the other ramification points. Let $\gamma_{1}$ be a loop in $\mathbb{C}-\{2,-2\}$ consisting of a path from 0 to a point near 2 given by $\pi \circ \alpha_{1}$, then a circular path around 2 and finally the former path back to 0 . Similarly, define $\gamma_{2}$ as a loop around the point -2 . How do $\gamma_{1}$ and $\gamma_{2}$ act on $\Gamma(0)$ ?
(d) Show that, although the loops $\gamma_{1}$ and $\gamma_{2}$ go around the points 2 and -2 respectively, they do not generate the group $\pi_{1}(\mathbb{C}-\{2,-2\}, 0)$. Represent these loops in $\mathbb{C}-\{2,-2\}$.
2.-Let $K$ be an algebraically closed field of characteristic 0 (not necessarily equal to $\mathbb{C}$ ), and $\Omega$ an algebraically closure of $K(Z)$. The aim is to show that the Galois group $\operatorname{Aut}_{\mathrm{K}(\mathrm{Z})} \Omega$ is (not naturally) isomorphic to the free profinite group on the set K.
(a) Show that if a profinite group $G$ has a basis with $n$ elements, then every family of $n$ elements generating G is a basis.
(b) Let X be a set, $\left(\left(\mathrm{G}_{\mathrm{E}}\right),\left(\rho_{\mathrm{F}}^{\mathrm{E}}\right)_{\mathrm{F} \subset \mathrm{E}}\right)$ an inverse system of profinite groups, indexed by the finite subsets $E$ of $X$. Set $G_{X}=\lim G_{E}$, and define the maps $\rho_{\mathrm{E}}^{\mathrm{X}}: \mathrm{G}_{\mathrm{X}} \rightarrow \mathrm{G}_{\mathrm{E}}$ by passing to the limit. Let Y be either a finite subset of X or $\mathrm{Y}=\mathrm{X}$. A family $\left(g_{x}\right)_{x \in \mathrm{Y}}$ of elements in $\mathrm{G}_{\mathrm{Y}}$ will be said to be adapted if $\rho_{\mathrm{E}}^{\mathrm{Y}}\left(g_{x}\right)=e$ for every finite subset E of Y such that $x \notin \mathrm{E}$.

Show that if the maps $\rho_{\mathrm{F}}^{\mathrm{E}}$ are surjective and each $\mathrm{G}_{\mathrm{E}}$ has an adapted basis, then so does the profinite group $G_{X}$ (use Tychonoff's theorem and a).

Show that every adapted basis for $\mathrm{G}_{\mathrm{E}}$ stems from an adapted basis for $\mathrm{G}_{\mathrm{X}}$.
(c) If $k$ is a finite extension of $\mathrm{K}(\mathrm{Z})$, let $\mathrm{R}(k)$ be its ramification set (6.3, Exercise 3). It is a finite set in $\mathrm{K} \cup\{\infty\}$.

Let $\mathrm{E} \subset \mathrm{K}$ be a finite set, $\Omega_{\mathrm{E}}$ the union of the finite extensions of $\mathrm{K}(\mathrm{Z})$ in $\Omega$ such that $\mathrm{R}(k) \subset \mathrm{E} \cup\{\infty\}$. Set

$$
\mathrm{G}_{\mathrm{E}}=\mathrm{Aut}_{\mathrm{K}(\mathrm{Z})} \Omega_{\mathrm{E}} .
$$

Show that, if $\mathrm{K}=\mathbb{C}$, the profinite group $\mathrm{G}_{\mathrm{E}}$ has an adapted basis.
(d) Let K and $\mathrm{K}^{\prime}$ be algebraically closed of characteristic 0 and $\iota: \mathrm{K}^{\prime} \rightarrow \mathrm{K}$ an embedding enabling the identification of $K^{\prime}$ with a subfield of $K$. Set

$$
k_{0}=\mathrm{K}(\mathrm{Z}) \text { and } k_{0}^{\prime}=\mathrm{K}^{\prime}(\mathrm{Z}) ;
$$

the field $k_{0}^{\prime}$ is a subfield of $k_{0}$, and any algebraic closure $\Omega$ of $k_{0}$ contains an algebraic closure $\Omega^{\prime}$ of $k_{0}^{\prime}$. Let E be a subset of $\mathrm{K}^{\prime}$. Define $\Omega_{\mathrm{E}}$ and $\Omega_{\mathrm{E}}^{\prime}$ as in (c). Show that $\Omega_{\mathrm{E}}^{\prime}=\Omega^{\prime} \cap \Omega_{\mathrm{E}}$, and that the restriction morphism $\iota^{*}: \mathrm{G}_{\mathrm{E}} \rightarrow \mathrm{G}_{\mathrm{E}}^{\prime}$ is an isomorphism.

Deduce that (c) continues to hold for subfields K of $\mathbb{C}$, and for arbitrary algebraically closed fields $K$ of characteristic 0 .
(e) Conclude.
3.-Which profinite groups are isomorphic to the Galois group of an extension of $\mathbb{C}(\mathrm{Z})$ ?

### 6.5 Triangulation of Riemann Surfaces

### 6.5.1 Definition of a Triangulation

The reference segment is $[0,1]$. The reference triangle T is the triangle in $\mathbb{C}$ with vertices $\alpha_{0}=0, \alpha_{1}=1$ and $\alpha_{2}=j+1=e^{i \pi / 3}$ i.e. the convex envelope of $\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$. We parametrize T by the maps $\iota_{0}, \iota_{1}, \iota_{2}$ defined by $\iota_{0}(t)=t$, $\iota_{1}(t)=1+j t, \iota_{2}(t)=j+1+j^{2} t$. Let $\stackrel{\circ}{\mathrm{T}}$ denote the interior of T .


Let X be a compact topological surface. A triangulation $\tau$ of X is given by

- three finite sets $\mathrm{I}_{0}, \mathrm{I}_{1}, \mathrm{I}_{2}$;
- a family $\left(s_{i}\right)_{i \in \mathrm{I}_{0}}$ of points of X, called the vertices of $\tau$;
- a family $\left(a_{j}\right)_{j \in \mathrm{I}_{1}}$ of injective continuous maps from $[0,1]$ to X , called the edges of $\tau$;
- a family $\left(f_{k}\right)_{k \in \mathrm{I}_{2}}$ of injective continuous maps from T to X , called the faces of $\tau$.

For an edge $a$, the point $a(0)$ (resp. $a(1)$ ) is called the origin (resp. endpoint) of $a$. The sides of a face $f$ are $f \circ \iota_{\lambda}$. The subsets $\mathrm{K}_{0}(\tau)=\left\{s_{i}\right\}, \mathrm{K}_{1}(\tau)=\bigcup a_{j}([0,1])$ and $\mathrm{K}_{2}(\tau)=\bigcup f_{k}(\mathrm{~T})$ are called the skeletons of $\tau$.

The above are subject to the following conditions:
( $\mathrm{T}_{1}$ ) The origin and endpoint of an edge $a_{j}$ are vertices.
$\left(\mathrm{T}_{2}\right) a_{j}(] 0,1[)$ are disjoint and do not meet $\mathrm{K}_{0}$.
$\left(\mathrm{T}_{3}\right)$ A side of a face $f_{k}$ is either an edge $a_{j}$ or a flipped edge $t \mapsto a_{j}(1-t)$.
( $\mathrm{T}_{4}$ ) $f_{k}(\stackrel{\circ}{\mathrm{~T}})$ are disjoint and do not meet $\mathrm{K}_{1}$.
( $\mathrm{T}_{5}$ ) $\mathrm{K}_{2}=\mathrm{X}$.
For $j \in \mathrm{I}_{1}$ and $k \in \mathrm{I}_{2}$, define the incidence number $\varepsilon(j, k)$ by

$$
\varepsilon(j, k)= \begin{cases}1 & \text { if } a_{j} \text { is a side of } f_{k} \\ -1 & \text { if the flipped edge } a_{j} \text { is a side of } f_{k} \\ 0 & \text { otherwise }\end{cases}
$$

For $i \in \mathrm{I}_{0}$ and $j \in \mathrm{I}_{1}$, define the incidence number $\varepsilon(i, j)$ by

$$
\varepsilon(i, j)= \begin{cases}1 & \text { if } s_{i} \text { is the endpoint of } a_{j} \\ -1 & \text { if } s_{i} \text { is the origin of } a_{j} \\ 0 & \text { otherwise }\end{cases}
$$



## Example Tetrahedral triangulation of the sphere $\mathbf{S}^{\mathbf{2}}$

$$
\begin{array}{lr}
a_{1}(0)=s_{0} a_{1}(1)=s_{1} & f_{0} \circ \iota_{0}=a_{4} f_{0} \circ \iota_{1}=a_{5} f_{0} \circ \iota_{2}=a_{6} \\
a_{2}(0)=s_{0} a_{2}(1)=s_{2} & f_{1} \circ \iota_{0}=a_{2} f_{1} \circ \iota_{1}=a_{4} f_{1} \circ \iota_{2}=a_{3} \circ \tau \\
a_{3}(0)=s_{0} a_{3}(1)=s_{3} \quad f_{2} \circ \iota_{0}=a_{3} f_{2} \circ \iota_{1}=a_{5} f_{2} \circ \iota_{2}=a_{1} \circ \tau \\
a_{4}(0)=s_{2} a_{4}(1)=s_{3} & f_{3} \circ \iota_{0}=a_{1} f_{3} \circ \iota_{1}=a_{6} f_{3} \circ \iota_{2}=a_{2} \circ \tau \\
a_{5}(0)=s_{3} a_{5}(1)=s_{1} & \text { where } \tau(t)=1-t .
\end{array}
$$

Incidence numbers:
On $\mathrm{I}_{0} \times \mathrm{I}_{1}$,

| $j i$ | 0 | 1 | 2 | 3 |
| :---: | ---: | ---: | ---: | ---: |
| 1 | -1 | 1 | 0 | 0 |
| 2 | -1 | 0 | 1 | 0 |
| 3 | -1 | 0 | 0 | 1 |
| 4 | 0 | 0 | -1 | 1 |
| 5 | 1 | 0 | 0 | -1 |
| 6 | -1 | 1 | 0 | 0 |


| on $\mathrm{I}_{1} \times \mathrm{I}_{2}$, |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $k^{j}$ | 1 | 2 | 3 | 4 | 5 | 6 |
| 0 | 0 | 0 | 0 | 1 |  | 1 |
| 1 | 0 | 1 | -1 | 1 | 0 | 0 |
| 2 | -1 | 0 | 1 | 0 | 1 | 0 |
| 3 |  | -1 | 0 | 0 | 0 | 1 |

### 6.5.2 Direct $\mathbf{C}^{1}$ Triangulations

Let X be a Riemann surface. A triangulation $\tau$ a of X is called ${ }^{3}$ a $\operatorname{direct} \mathrm{C}^{1}$ triangulation if the following conditions hold:
$\left(\mathrm{TD}_{1}\right)$ For $j \in \mathrm{I}_{1}$, the map $a_{j}$ is $\mathrm{C}^{1}$ with nonzero derivative on $] 0,1[$, i.e. for $t \in$ ]0, $1\left[\right.$ there is an analytic chart $\varphi$ for X centered at $a_{j}(t)$ such that $\varphi \circ a_{j}$ : $] 0,1[\rightarrow \mathbb{C}$ defined in the neighbourhood of $t$ is continuously differentiable in the neighbourhood of $t$ with nonzero derivative at $t$ (this condition then holds for all charts).
$\left(\mathrm{TD}_{2}\right)$ For $k \in \mathrm{I}_{2}$, the map $f_{k}$ is $\mathrm{C}^{1}$ on $\mathrm{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ in the real sense. Moreover, let $t \in \mathrm{~T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ and $\varphi$ an analytic chart for X centered at $f_{k}(t)$, then the real functions $u$ and $v$ defined in the neighbourhood of $t$ by $f(x+i y)=$ $u(x, y)+i v(x, y)$ satisfy

$$
\operatorname{det}\left(\begin{array}{ll}
\frac{\partial u}{\partial x} & \frac{\partial v}{\partial x} \\
\frac{\partial u}{\partial y} & \frac{\partial v}{\partial y}
\end{array}\right)>0
$$

(this condition then holds for all charts, for if $h: \mathbb{C} \rightarrow \mathbb{C}$ is a $\mathbb{C}$-linear map, then $h$ can be written $z \cdot 1_{\mathbb{C}}$, and, as an $\mathbb{R}$-linear map from $\mathbb{R}^{2}$ to itself, the determinant of $h$ is $z \bar{z}>0$ ).

[^21]We next show that every compact Riemann surface has a direct $\mathrm{C}^{1}$ triangulation.
Remark Identifying the Riemann sphere $\Sigma=\mathbb{C} \cup\{\infty\}$ with the sphere $S^{2}$, the triangulation in Example 6.5.1 is not direct (incidence table and Exercise 1).

### 6.5.3 Existence of Triangulations of the Riemann Sphere

Proposition Let $\Delta$ be a finite subset of the Riemann sphere $\Sigma$. Then there is a direct $\mathrm{C}^{1}$ triangulation $\tau$ of $\Sigma$ such that $\mathrm{K}_{0}(\tau)$ contains $\Delta$.

This proposition is an immediate consequence of the next two lemmas.
Lemma 6.1 There is a triangulation of $\Sigma$.
Proof The map $f_{0}$ which assigns to the point $(\rho, \theta)$ the point with polar coordinates $\left(\rho \frac{\cos (\theta-\pi / 6)}{\cos \pi / 6}, \frac{3}{2} \theta\right)$ is a homeomorphism from T onto the quarter circle $\{z=x+i y \mid$ $|z| \leqslant 1, x \geqslant 0, y \geqslant 0\}$. Define $f_{1}, \ldots, f_{7}$ by

$$
f_{1}(t)=i f_{0}(t), f_{2}(t)=-f_{0}(t), f_{3}(t)=-i f_{0}(t), f_{k+4}(t)=\frac{1}{f_{k}(t)}
$$



Check that $\left(f_{k}\right)_{k=0, \ldots, 7}$ are the faces of a direct $\mathbf{C}^{1}$ triangulation of $\Sigma$ with 8 faces, 12 edges and the 6 vertices $0,1, i,-1,-i, \infty$.

Lemma 6.2 Let X be a compact Riemann surface, $\tau$ a direct $\mathrm{C}^{1}$ triangulation of X and $s$ a point of X . Then there is a direct $\mathrm{C}^{1}$ triangulation $\tau^{\prime}$ of X such that $\mathrm{K}_{0}\left(\tau^{\prime}\right)=\mathrm{K}_{0}(\tau) \cup\{s\}$.

Proof If $s \in \mathrm{~K}_{0}$, take $\tau^{\prime}=\tau$.
If $s \in \mathrm{~K}_{1}-\mathrm{K}_{0}$, then there exists $\left.u \in\right] 0,1\left[\right.$ such that $s=a_{j}(u)$. For all $k \in \mathrm{I}_{2}$ such that $\varepsilon(j, k) \neq 0, v=f_{k}^{-1}(s)$ belongs to a side of the reference triangle T . Let $\alpha$ and $\beta$ be the endpoints of this side, $\gamma$ the opposite vertex, $\sigma_{k}^{\prime}$ and $\sigma_{k}^{\prime \prime}$ direct affine homeomorphisms from T onto the triangles with vertices $\alpha, \gamma, v$ and $\beta, \gamma$, $v$ respectively, and $\lambda_{k}:[0,1] \rightarrow \mathrm{T}$ an affine map such that $\lambda_{k}(0)=\gamma, \lambda_{k}(1)=v$. Replace the edge $a_{j}$ by $a_{j}^{\prime}: t \mapsto a_{j}(u t)$ and $a_{j}^{\prime \prime}: t \mapsto a_{j}(1-t+u t)$; for each $k$ such that $\varepsilon(j, k) \neq 0$, replace the face $f_{k}$ by $f_{k}^{\prime}=f_{k} \circ \sigma_{k}^{\prime}$ and $f_{k}^{\prime \prime}=f_{k} \circ \sigma_{k}^{\prime \prime}$ and add an edge $f_{k} \circ \lambda_{k}$.

If $s \in \mathrm{~K}_{2}-\mathrm{K}_{1}$, then there is a unique pair $(k, t) \in \mathrm{I}_{2} \times \stackrel{\circ}{\mathrm{T}}$ such that $s=f_{k}(t)$. Let $\sigma$ (resp. $\sigma^{\prime}$, resp. $\sigma^{\prime \prime}$ ) be an affine map from T to T sending $\alpha_{0}, \alpha_{1}, \alpha_{2}$ to $t, \alpha_{0}, \alpha_{1}$ (resp. $t, \alpha_{1}, \alpha_{2}$, resp. $t, \alpha_{2}, \alpha_{0}$ ) respectively, and $\lambda, \lambda^{\prime}, \lambda^{\prime \prime}$ affine maps from [0, 1] to T sending 0 to $t$ and 1 to $\alpha_{0}, \alpha_{1}, \alpha_{2}$ respectively. Replace the face $f_{k}$ by faces $f_{k} \circ \sigma$, $f_{k} \circ \sigma^{\prime}, f_{k} \circ \sigma^{\prime \prime}$ and add three edges $f_{k} \circ \lambda, f_{k} \circ \lambda^{\prime}, f_{k} \circ \lambda^{\prime \prime}$.

In all cases, we obtain a triangulation $\tau^{\prime}$ answering the question.

### 6.5.4 Lifting of a Triangulation

Let X and Y be compact topological surfaces equipped with the respective triangulations $\tau$ and $\tau^{\prime}$. A continuous map $h: \mathrm{X} \rightarrow \mathrm{Y}$ is compatible with triangulations $\tau$ and $\tau^{\prime}$ if for all vertices $s_{i}$ (resp. all edges $a_{j}$, resp. all faces $f_{k}$ ) of $\tau$, the image $h\left(s_{i}\right)$ (resp. $h \circ a_{j}$, resp. $h \circ f_{k}$ ) is a vertex (resp. an edge, resp. a face) of $\tau^{\prime}$.

Proposition Let B be a compact topological surface, X a finite ramified covering of $\mathrm{B}, \Delta$ the ramification set of X , and let $\tau$ be a triangulation of B such that $\Delta \subset \mathrm{K}_{0}(\tau)$. Then there is a triangulation $\tau^{\prime}$ of X , unique up to permutation of indices, such that the projection $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is compatible with $\tau^{\prime}$ and $\tau$.

Proof Let $f$ be a face (resp. an edge) of $\tau$. The space $f^{*}(\mathrm{X})$ restricted to $\mathrm{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ (resp. ]0,1[) is a contractible covering of $\mathrm{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ (resp. $] 0,1[)$. Hence this covering has $d$ sections where $d$ is the degree of X over B . Corresponding to them there are $d$ liftings $g_{\nu}: \mathrm{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\} \rightarrow \mathrm{X}$ (resp. $] 0,1[\rightarrow \mathrm{X}$ ) of $f$ with $\nu \in\{1, \ldots, d\}$.

The next lemma shows that this liftings extend by continuity to T (resp. [0, 1]).
Lemma Let Z be a topological space, $z_{0} \in \mathrm{Z}$. Suppose that $z_{0}$ has a fundamental system of neighbourhoods $\mathrm{W}_{i}$ such that for all $i, \mathrm{~W}_{i}-\left\{z_{0}\right\}$ is connected and nonempty. Let $\varphi: \mathrm{Z} \rightarrow \mathrm{B}$ and $\psi: \mathrm{Z}-\left\{z_{0}\right\} \rightarrow \mathrm{X}$ be continuous maps such that $\pi \circ \psi=\left.\varphi\right|_{\mathrm{Z}-\left\{\mathrm{Z}_{0}\right\}}$. Then $\psi$ can be uniquely extended to a continuous map $\bar{\psi}: \mathrm{Z} \rightarrow \mathrm{X}$ such that $\pi \circ \psi=\varphi$

Proof of the Lemma Let $b_{0}=\varphi\left(z_{0}\right), \mathrm{X}\left(b_{0}\right)=\left\{x_{1}, \ldots, x_{r}\right\}$, and $\mathrm{U}_{1}, \ldots, \mathrm{U}_{r}$ disjoint open neighbourhoods of respectively $x_{1}, \ldots, x_{r}$ in X . As $\pi$ is proper, there
is a neighbourhood V of $b_{0}$ in B such that $\pi^{-1}(\mathrm{~V}) \subset \mathrm{U}_{1} \cup \ldots \cup \mathrm{U}_{r}$. As $\varphi$ is continuous, there exists $\mathrm{W}_{i}$ such that $\varphi\left(\mathrm{W}_{i}\right) \subset \mathrm{V}$. Then for $p=1, \ldots, r$, the open sets $\psi^{-1}\left(\mathrm{U}_{p}\right) \cap \mathrm{W}_{i}$ form a partition of $\mathrm{W}_{i}-\left\{z_{0}\right\}$. Hence only one of them is nonempty. Without loss of generality, assume it is $\psi^{-1}\left(\mathrm{U}_{1}\right) \cap \mathrm{W}_{i}$. We then show that $\psi(z)$ tends to $x_{1}$ as $z$ tends to $z_{0}$ while remaining distinct from $z_{0}$. As $\mathrm{V}^{\prime}$ runs through a fundamental system of neighbourhoods of $b_{0}$ in B , the set $\pi^{-1}\left(\mathrm{~V}^{\prime}\right) \cap \mathrm{U}_{1}$ runs through a fundamental system of neighbourhoods of $x_{1}$ in X . Indeed, if $\mathrm{U}_{1}^{\prime}$ is a neighbourhood of $x_{1}$ in X , then $\mathrm{U}_{1}^{\prime} \cup \mathrm{U}_{2} \cup \ldots \cup \mathrm{U}_{r}$ contains some set $\pi^{-1}\left(\mathrm{~V}^{\prime}\right)$ and $\pi^{-1}\left(\mathrm{~V}^{\prime}\right) \cap \mathrm{U}_{1}=\mathrm{U}_{1}^{\prime}$. Then $\psi^{-1}\left(\pi^{-1}\left(\mathrm{~V}^{\prime}\right) \cap \mathrm{U}_{1}\right)=\varphi^{-1}\left(\mathrm{~V}^{\prime}\right) \cap \psi^{-1}\left(\mathrm{U}_{1}\right) \supset \varphi^{-1}\left(\mathrm{~V}^{\prime}\right) \cap\left(\mathrm{W}_{i}-\left\{z_{0}\right\}\right)$. Hence the desired extension of $\psi$ is $\psi$ defined by $\psi\left(z_{0}\right)=x_{1}$.

End of the Proof of the Proposition Applying the lemma to $\mathrm{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ (resp. ]0, 1[) three times (resp. twice) lifts each face $f_{k}$ (resp. each edge $a_{j}$ ) of $\tau$ to $d$ faces (resp. $d$ edges) in X. Take as vertices the points of X over the vertices of $\tau$. These vertices, edges and faces form a triangulation $\tau^{\prime}$ of X unique up to permutation of indices with the desired property.

Remarks (1) Suppose that B is a Riemann surface and X an analytic covering of B. Then if $\tau$ is $\mathrm{C}^{1}$ and direct, so is $\tau^{\prime}$.

Indeed, this property concerns the behaviour of $\tau$ away from vertices, and there $\pi: \mathrm{X} \rightarrow \mathrm{B}$ is a local isomorphism.
(2) Let $k_{0}(\tau), k_{1}(\tau), k_{2}(\tau)$ be the respective number of elements of $\mathrm{I}_{0}, \mathrm{I}_{1}, \mathrm{I}_{2}$. Then

$$
k_{2}\left(\tau^{\prime}\right)=d \cdot k_{2}(\tau), k_{1}\left(\tau^{\prime}\right)=d \cdot k_{1}(\tau), k_{0}\left(\tau^{\prime}\right)=d \cdot k_{0}(\tau)-\sum_{x \in \mathbb{R}}(e(x)-1)
$$

where R is the set of ramification points of $\mathrm{X}, e(x)$ denoting the ramification index of $x$.

### 6.5.5

Theorem Let X be a compact Riemann surface and $\Delta \subset \mathrm{X}$ a finite set. There is a direct $\mathrm{C}^{1}$ triangulation $\tau$ of X such that $\mathrm{K}_{0}(\tau) \supset \Delta$.

Proof Let $f \in \mathscr{M}(\mathrm{X})$ be a meromorphic function which is not constant on any of the connected components (such functions exist by the separation Theorem 6.2.1). There is an analytic map $\pi: \mathrm{X} \rightarrow \Sigma$ corresponding to $f$. By Proposition 6.5.3, there is a direct $\mathrm{C}^{1}$ triangulation $\tau_{0}$ of $\Sigma$ such that $\mathrm{K}_{0}\left(\tau_{0}\right)$ contains the projection of $\Delta$ and the ramification set of X . By Proposition 6.5.4, there is a triangulation $\tau$ of X such that $\tau$ is compatible with $\tau$ and $\tau_{0}$, which implies that $\mathrm{K}_{0}(\tau) \supset \pi^{-1}(\pi(\Delta)) \supset \Delta$. By 6.5.4, remark 1 , the triangulation $\tau$ is direct and $\mathrm{C}^{1}$.

### 6.5.6

Proposition Let X be a Riemann surface, $\tau$ a direct $\mathrm{C}^{1}$ triangulation. For each edge $a_{j}$ there are exactly two faces $f_{k}$ such that $\varepsilon(j, k) \neq 0$. For one of them (which will be said to be on the left of $a_{j}$ ), $\varepsilon(j, k)=+1$, for the other (which will be said to be on the right), $\varepsilon(j, k)=-1$.

Lemma 6.3 Let J be a neighbourhood of 0 in $\mathbb{R}$ and $a: \mathrm{J} \rightarrow \mathbb{R}^{2} a \mathrm{C}^{1}$ map such that $a(0)=0$ and $a^{\prime}(0) \neq 0$. Then there is a neighbourhood U of 0 in $\mathbb{R}^{2}$, a neighbourhood $\mathbf{J}^{\prime}$ of 0 in $\mathbf{J}$, and a diffeomorphism $\psi$ from U onto an open subset of $\mathbb{R}^{2}$ with Jacobian $>0$ such that $\psi\left(a\left(\mathbf{J}^{\prime}\right)\right)=\psi(\mathrm{U}) \cap(\mathbb{R} \times\{0\})$ and for all $t \in \mathbf{J}^{\prime}$, $\psi(a(t))=(t, 0)$.

Proof Let $\lambda: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ be a linear map with determinant $>0$ such that $\lambda\left(a^{\prime}(0)\right)=$ $(\xi, \eta)$ for some $\xi>0$. Then $\lambda(a(t))=(u(t), v(t))$ for some $u^{\prime}(0)=\xi>0$. So $u$ induces a diffeomorphism from an open interval $\mathrm{J}^{\prime}$ of 0 in J onto a neighbourhood $\mathbf{J}^{\prime \prime}$ of 0 in $\mathbb{R}$. Set $\theta(x, y)=\left(u^{-1}(x), y-v\left(u^{-1}(x)\right)\right)$ for $x \in \mathbf{J}^{\prime \prime}, y \in \mathbb{R}$. Then $\mathrm{U}=$ $\lambda^{-1}\left(\mathbf{J}^{\prime \prime} \times \mathbb{R}\right)$ and $\psi=\theta \circ \lambda$ are as desired, the Jacobian of $\theta$ being $\operatorname{det}\left(\begin{array}{c}u^{\prime-1} * \\ 0\end{array} 1\right)=$ $u^{\prime-1}>0$.

Lemma 6.4 Let V be a neighbourhood of 0 in $\mathbb{R} \times \mathbb{R}_{+}$, such that $\mathrm{V}=\mathrm{V} \cap(\mathbb{R} \times$ $\mathbb{R}_{+}^{*}$ ) is connected. Let $f: \mathrm{V} \rightarrow \mathbb{R}^{2}$ be a $\mathrm{C}^{1}$ injective map with Jacobian $>0$ on V , such that $f(0)=0$. Set $a(t)=f(t, 0)$. Let $\mathrm{U}, \psi$ and $\mathbf{J}^{\prime}$ be as in Lemma 6.3. Assume that $\mathrm{J}^{\prime} \times\{0\} \subset \mathrm{V}$ and $f(\mathrm{~V}) \subset \mathrm{U}$. Then $\psi(f(\mathrm{~V}))$ is a neighbourhood of 0 in $\mathbb{R} \times \mathbb{R}_{+}$.

Proof Let $\mathrm{D} \subseteq \mathrm{V}$ be a half-disk centered at 0 closed in $\mathbb{R} \times \mathbb{R}_{+}$and let $\stackrel{\circ}{\mathrm{D}}$ denote the interior of D in $\mathbb{R}^{2}$. The boundary $\partial \mathrm{D}=\mathrm{D}-\stackrel{\circ}{\mathrm{D}}$ is of the form $\mathrm{J} \cup \mathrm{L}$, where J is the segment and L the closed half-circle. Set $g=\psi \circ f$. The set $g(\mathrm{~L})$ is compact in $\mathbb{R}^{2}$ and $0 \notin g(\mathrm{~L})$. Let $\mathrm{D}^{\prime}$ be a closed half-circle centered at 0 in $\mathbb{R} \times \mathbb{R}_{+}$such that $\stackrel{\circ}{\mathrm{D}^{\prime}} \cap g(\mathrm{~L})=\varnothing$. As $g(\mathrm{~J}) \subset \mathbb{R} \times\{0\}, g(\mathrm{D}) \cap \stackrel{\circ}{\mathrm{D}^{\prime}}=g(\stackrel{\circ}{\mathrm{D}}) \cap \stackrel{\circ}{\mathrm{D}^{\prime}}$. Now, $g(\mathrm{D})$ is compact, and so closed, and $g(\stackrel{\circ}{\mathrm{D}})$ is open in $\mathbb{R}^{2}$ by the local inversion theorem (implicit function
 $g(\mathrm{D}) \supset \stackrel{\circ}{\mathrm{D}^{\prime}}$ or $g(\mathrm{D}) \cap \stackrel{\circ}{\mathrm{D}^{\prime}}=\varnothing$. Define $u, v \in[0, \varepsilon[\subset \mathbb{R}$ by $g(0, t)=(u(t), v(t))$. As $g(t, 0)=(t, 0)$, the Jacobian of $g$ at 0 is $\operatorname{det}\left(\begin{array}{ll}1 & u^{\prime}(0) \\ 0 & v^{\prime}(0)\end{array}\right)=v^{\prime}(0)$. Thus $v^{\prime}(0)>0$. For sufficiently small $t>0, u(t)$ is small and so is $v(t)>0$; thus $g(0, t) \in \stackrel{\circ}{\mathrm{D}^{\prime}}$ and


Indeed $g(\stackrel{\circ}{\mathrm{~V}})$ is connected and does not meet $\mathbb{R} \times\{0\}$ since $g$ is injective and $(\mathbb{R} \times\{0\}) \cap \mathrm{U}=g\left(\mathrm{~J}^{\prime} \times\{0\}\right) . \quad$ As $\quad$ a $\quad$ result, $\quad g(\stackrel{\circ}{\mathrm{~V}}) \subseteq \mathbb{R} \times \mathbb{R}_{+}^{*}, \quad$ and $\quad g(\mathrm{~V}) \subset$ $\mathbb{R} \times \mathbb{R}_{+}$.

Proof of the Proposition Let $j \in \mathrm{I}_{1}$, and set $x=a_{j}\left(\frac{1}{2}\right)$. Let $\varphi$ be an analytic chart for X centered at $x$. Define $\tilde{a}:]-\frac{1}{2}, \frac{1}{2}\left[\rightarrow \mathbb{R}^{2}\right.$ in the neighbourhood of 0 by

$$
\tilde{a}(t)=\varphi\left(a_{j}\left(t+\frac{1}{2}\right)\right)
$$

Let $\mathbf{J}^{\prime}, \mathrm{U}$ and $\psi$ satisfy the conditions of Lemma 6.3 for $\tilde{a}$. Let $f_{k}$ be a face of $\tau$ such that $\varepsilon(j, k) \neq 0$. Then $x_{k}=f_{k}^{-1}(x)$ is the midpoint of a side of T, i.e. $x_{k}=$ $\iota_{\nu}\left(\frac{1}{2}\right)$ for some $\nu \in\{0,1,2\}$. Let $h_{k}: \mathbb{C} \rightarrow \mathbb{C}$ be a $\mathbb{C}$-affine map such that $h_{k}(t)=$ $\iota_{\nu}\left(t+\frac{1}{2}\right)$. Then $h_{k}$ defines a homeomorphism from a neighbourhood $\mathrm{V}_{k}$ of 0 in $\mathbb{R} \times \mathbb{R}_{+}$onto a neighbourhood W of $x_{k}$ in T . We may assume that $\mathrm{W} \subset f_{k}^{-1}\left(\varphi^{-1}(\mathrm{U})\right)$, $\mathrm{W} \cap \iota_{\lambda}([0,1])=\varnothing$ for $\lambda \neq \nu, \mathrm{W} \cap \stackrel{\circ}{\mathrm{T}}$ connected and $\mathrm{V}_{k} \supset \mathrm{~J}^{\prime}$ if $\varepsilon(j, k)=1$, while $\mathrm{V}_{k} \supset_{\tilde{J}^{\prime}}-\mathrm{J}^{\prime}$ if $\varepsilon(j, k)=-1$. Set $\widetilde{f}_{k}=\varphi \circ f_{k} \circ h_{k}: \mathrm{V}_{k} \rightarrow \mathrm{U} \subset \mathbb{R}$. If $\varepsilon(j, k)=1$, then $\psi \circ \widetilde{f}_{k}(t)=(t, 0)$. Then, by Lemma $6.4, \psi \circ \stackrel{\tilde{f}}{k}\left(\mathrm{~V}_{k}\right)$ is a neighbourhood of 0 in $\mathbb{R} \times \mathbb{R}_{+}$. If $\varepsilon(j, k)=-1$, then $\psi \circ \widetilde{f}_{k}(t)=(-t, 0)$, and by Lemma 6.4 applied to $-\psi, \psi \circ \widetilde{f}_{k}\left(\mathrm{~V}_{k}\right)$ is a neighbourhood of 0 in $\mathbb{R} \times \mathbb{R}_{-}$.

Let $\mathrm{I}_{2}^{\prime}$ denote the set of $k$ such that $\varepsilon(j, k) \neq 0$. Then

$$
\Omega=\bigcup_{k \in \mathrm{I}_{2}^{\prime}} f_{k} \circ h_{k}\left(\mathrm{~V}_{k}\right)
$$

is a neighbourhood of $x$ in X . Indeed, taking $\mathrm{V}_{k}=\varnothing$ for $k \notin \mathrm{I}_{2}^{\prime}$, $C \Omega \subset \bigcup_{k \in \mathrm{I}_{2}} f_{k}\left(\mathrm{~T}-h_{k}\left(\mathrm{~V}_{k}\right)\right)$, and this union is a compact set not containing $x$. Hence $\bigcup \psi \circ \widetilde{f}_{k}\left(\mathrm{~V}_{k}\right)=\psi(\varphi(\Omega))$ is a neighbourhood of 0 in $\mathbb{R}^{2}$.

As all $(u, v)$ with $v>0(\operatorname{resp} v<0)$ belong to a neighbourhood of 0 in $\mathbb{R}^{2}$, there exists $k$ such that $\varepsilon(j, k)=+1$ (resp. -1 ).

If $\varepsilon(j, k)=\varepsilon\left(j, k^{\prime}\right)=1$ with $k \neq k^{\prime}$, then the set $\psi\left(\tilde{f}_{k}\left(\mathrm{~V}_{k}\right)\right) \cap \psi\left(\widetilde{f_{k^{\prime}}}\left(\mathrm{V}_{k^{\prime}}\right)\right)$ is a neighbourhood of 0 in $\mathbb{R} \times \mathbb{R}_{+}$, and so contains points $(u, v)$ with $v>0$, but such a point is of the form $\varphi\left(x^{\prime}\right)$ with $x^{\prime} \in f_{k}(\stackrel{\circ}{\mathrm{~T}}) \cap f_{k^{\prime}}(\stackrel{\circ}{\mathrm{T}})$, which is impossible by condition $\left(\mathrm{T}_{4}\right)$ of 6.5.1. Hence there is a unique $k$ such that $\varepsilon(j, k)=1$ and similarly a unique $k$ such that $\varepsilon(j, k)=-1$.

Remark This proposition can be generalized to oriented topological surfaces, but the proof is clearly harder.

## Exercises 6.5. (Triangulation of a Riemann surface)

1.-Let X be a topological surface and $\tau=\left(\mathrm{I}_{0}, \mathrm{I}_{1}, \mathrm{I}_{2},\left(s_{i}\right),\left(a_{j}\right),\left(f_{k}\right)\right)$ a triangulation of X .
(a) Show that each edge is the side of two faces, i.e. that, for all $j \in \mathrm{I}_{2}$, there are exactly two values of $k$ such that $\varepsilon(j, k) \neq 0$.
(b) Suppose that X is a $\mathrm{C}^{1}$ surface and that $\tau$ is $\mathrm{C}^{1}$. Show that the following conditions are equivalent:
$\left(\forall j \in \mathrm{I}_{1}\right)\left(\exists k, k^{\prime} \in \mathrm{I}_{2}\right) \varepsilon(j, k)=1$ and $\varepsilon\left(j, k^{\prime}\right)=-1$;
(ii) there is a unique orientation of X such that every face of $\tau$ is a map preserving the orientation of $\stackrel{\circ}{\mathrm{T}}$;
(c) Generalize the definition of triangulations to surfaces with boundaries and give a triangulation of a closed Möbius strip. Using this triangulation, show that the Möbius strip is not orientable.
2.-Let X be a compact topological surface and $\tau$ a triangulation of X . The combinatorial data of $\tau$ are

- the finite sets $\mathrm{I}_{0}, \mathrm{I}_{1}, \mathrm{I}_{2}$;
- the incidence functions $\varepsilon^{0,1}: \mathrm{I}_{0} \times \mathrm{I}_{1} \rightarrow\{0,1,-1\}$ and $\varepsilon^{1,2}: \mathrm{I}_{1} \times \mathrm{I}_{2} \rightarrow\{0,1,-1\}$;
- For $(j, k) \in \mathrm{I}_{1} \times \mathrm{I}_{2}$ with $\varepsilon(j, k) \neq 0, \nu \in\{0,1,2\}$ such that the side $f_{k} \circ \iota_{\nu}$ of $f_{k}$ is either the edge $a_{j}$ or the flip of the edge $a_{j}$.
(a) Let T be the reference triangle (6.5.1) and $\mathrm{F}: \mathrm{I}_{2} \times \mathrm{T} \rightarrow \mathrm{X}$ the map defined by $\mathrm{F}(k, t)=f_{k}(t)$. Show that F is surjective, and determine from the combinatorial data the equivalence relation on $\mathrm{I}_{2} \times \mathrm{T}$ induced by F .
(b) Show that any two compact topological surfaces admitting triangulations and with the same combinatorial data are homeomorphic.


### 6.6 Simplicial Homology

In this section X denotes a compact Riemann surface and $\tau=\left(\left(s_{i}\right)_{i \in \mathrm{I}_{0}},\left(a_{j}\right)_{j \in \mathrm{I}_{1}}\right.$, $\left.\left(f_{k}\right)_{k \in \mathbf{I}_{2}}\right)$ a direct $\mathbf{C}^{1}$ triangulation of $\mathbf{X}$. The results may be generalized to an arbitrary oriented topological surface. The proofs use Proposition 6.5.6.

### 6.6.1 Chain Complex Associated to a Triangulation

For $\nu=0,1,2$, let $\mathrm{C}_{\nu}$ be the free $\mathbb{Z}$-module $\mathbb{Z}^{\mathrm{I}_{\nu}}$ on $\mathrm{I}_{\nu}$. For $\nu \neq 0,1,2$, set $\mathrm{C}_{\nu}=0$. Denote by $\left[s_{i}\right]$ (resp. $\left[a_{j}\right]$, resp. $\left[f_{k}\right]$ ) the basis element corresponding to $i$ (resp. $j$, resp. $k$ ). Define $\partial_{\nu}: \mathrm{C}_{\nu} \rightarrow \mathrm{C}_{\nu-1}$ by

$$
\partial_{2}\left[f_{k}\right]=\sum_{j} \varepsilon(j, k)\left[a_{j}\right] \text { and } \partial_{1}\left[a_{j}\right]=\sum_{i} \varepsilon(i, j)\left[s_{i}\right]=\left[a_{j}(1)\right]-\left[a_{j}(0)\right] .
$$

The $\partial_{\nu}$ are called boundary operators. The boundary of a face may be described as consisting of three edges with their incidence number, and the boundary of an edge as the "endpoint minus the origin".

Proposition and Definition $\partial_{1} \circ \partial_{2}=0$.
The chain complex associated to $\tau$ is $0 \rightarrow \mathrm{C}_{2} \xrightarrow{\sigma_{2}} \mathrm{C}_{1} \xrightarrow{\sigma_{1}} \mathrm{C}_{0} \rightarrow 0$, and $\mathrm{H}_{\nu}(\mathrm{X}, \tau)=$ $\operatorname{Ker} \partial_{\nu} / \operatorname{Im} \partial_{\nu+1}$ is said to be the homology of X defined by $\tau$.

Proof Let $k \in \mathrm{I}_{2}$. Then

$$
\partial_{1} \circ \partial_{2}\left[f_{k}\right]=\sum_{j} \varepsilon(j, k)\left(\left[a_{j}(1)\right]-\left[a_{j}(0)\right]\right)
$$

Let $\nu \in\{0,1,2\}=\mathbb{Z} /(3), \alpha_{\nu}=\left[f_{k} \circ \iota_{\nu}(0)\right]$, and define $j_{\nu}$ by $f_{k} \circ \iota_{\nu}=a_{j_{\nu}}$ or the flip of $a_{j_{\nu}}$. In both cases $\varepsilon\left(j_{\nu}, k\right) \partial_{1}\left(a_{j_{\nu}}\right)=\alpha_{\nu+1}-\alpha_{\nu}$. If $j$ is not of the form $j_{\nu}$, then $\varepsilon(j, k)=0$. Hence $\partial_{1} \circ \partial_{2}\left[f_{k}\right]=\sum\left(\alpha_{\nu+1}-\alpha_{\nu}\right)=0$.

Remark As may be shown, $\mathrm{H}_{\nu}(\mathrm{X}, \tau)$ can be identified with the singular homology of X, and so does not depend on $\tau$ (cf. Zisman [7], 6.3.2.7). A particular case of this result will be proved in 6.6.4; see also 6.6.2 and 6.6, Exercise 1 .

### 6.6.2 Connected Case

Proposition If $\mathrm{X}(\neq \varnothing)$ is connected, then
(a) $\mathrm{H}_{0}(\mathrm{X}, \tau)=\mathbb{Z}$;
(b) $\mathrm{H}_{2}(\mathrm{X}, \tau)=\mathbb{Z}$.

The class of $\left[s_{i}\right]$ is the same for all $i \in \mathrm{I}_{0}$, and this element is a generator of $\mathrm{H}_{0}(\mathrm{X}, \tau)$. Then $\sum\left[f_{k}\right] \in \operatorname{Ker} \partial_{2}$, and the class of this element is a generator of $\mathrm{H}_{2}(\mathrm{X}, \tau)$.

Proof (a) $\mathrm{H}_{0}(\mathrm{X}, \tau)=$ Coker $\partial_{1}$. Define $\varepsilon: \mathrm{C}_{0} \rightarrow \mathbb{Z}$ by $\varepsilon\left(\left[s_{i}\right]\right)=1$ for all i. As $\mathrm{I}_{0} \neq \varnothing$, the map $\varepsilon$ is surjective, and clearly $\varepsilon \circ \partial_{1}=0$; hence $\operatorname{Im} \partial_{1} \subset \operatorname{Ker} \varepsilon$. We show the converse inclusion. The submodule $\operatorname{Ker} \varepsilon$ is generated by the $\left[s_{i}\right]-\left[s_{i^{\prime}}\right]$. Indeed, let $i_{0} \in \mathrm{I}$ and $c=\sum \lambda_{i}\left[s_{i}\right] \in \operatorname{Ker} \varepsilon$. Then $\sum \lambda_{i}=0$; so

$$
\lambda_{i_{0}}=-\sum_{i \neq i_{0}} \lambda_{i}, \text { and } c=\sum_{i \neq i_{0}} \lambda_{i}\left(\left[s_{i}\right]-\left[s_{i_{0}}\right]\right) .
$$

A vertex $s_{i}$ will be said to be connected to a point $x \in \mathrm{X}$ if there is a sequence $i_{0}=i, i_{1}, \ldots, i_{\mathrm{N}}$ such that, for $n=1, \ldots, \mathrm{~N}, s_{i_{n-1}}$ and $s_{i_{n}}$ are not endpoints of the same edge and $s_{i_{\mathrm{N}}}$ is the vertex of a face or of an edge containing $x$. If two vertices $s_{i}$ and $s_{i^{\prime}}$ are connected, then $\left[s_{i}\right]-\left[s_{i^{\prime}}\right] \in \operatorname{Im} \partial_{1}$.

For $i \in \mathrm{I}_{0}$, le $\mathrm{X}_{i}$ be the set of points of X connected to $s_{i}$. These $\mathrm{X}_{i}$ form a finite partition $X$ into nonempty closed sets, and so there is only one and any two vertices are connected. Consequently, $\operatorname{Im} \partial_{1}=\operatorname{Ker} \varepsilon$, and so $\mathrm{H}_{0}(\mathrm{X}, \tau)=\operatorname{Coker} \partial_{1}=\operatorname{Im} \varepsilon=\mathbb{Z}$.
(b) $\mathrm{H}_{2}=\operatorname{Ker} \partial_{2}$. Let $c=\sum \lambda_{k}\left[f_{k}\right] \in \mathrm{C}_{2}$. If all $\lambda_{k}$ are equal, it follows from Proposition 6.5.6 that $\partial_{2} c=0$. Conversely, suppose that $\partial_{2} c=0$. If the faces $f_{k}$ and $f_{k^{\prime}}$ have a common edge, then $\lambda_{k}=\lambda_{k^{\prime}}$. For all $n$, the union $X_{n}$ of $f_{k}(\mathrm{~T})$ for all $k$ such that $\lambda_{k}=n$ is a closed set. For $n \neq n^{\prime}, \mathrm{X}_{n} \cap \mathrm{X}_{n^{\prime}} \subset \mathrm{K}_{0}$, and so $\mathrm{X}_{n}^{\prime}=\mathrm{X}_{n}-\mathrm{K}_{0}$ form a finite partition of $X-K_{0}$ into closed sets. By 6.1.7, Remark 3, the space $X-K_{0}$ is connected; hence there is only one $X_{n}^{\prime} \neq \varnothing$, all $\lambda_{k}$ are equal, and $\left.\mathrm{H}_{2}(\mathrm{X}, \tau)\right)$ can be identified with $\mathbb{Z}$.

Corollary For X not necessarily connected,

$$
\mathrm{H}_{0}(\mathrm{X}, \tau)=\mathrm{H}_{2}(\mathrm{X}, \tau)=\mathbb{Z}^{\pi_{0}(\mathrm{X})}
$$

where $\pi_{0}(\mathrm{X})$ denotes the set of connected components of X .

### 6.6.3 Barycentric Subdivision

From the triangulation $\tau$, we define a triangulation $\tau^{\prime}$ of X . Let $\mathrm{I}_{0,1}$ (resp. $\mathrm{I}_{1,2}$ ) be the set of $(i, j) \in \mathrm{I}_{0} \times \mathrm{I}_{1}$ (resp. $\mathrm{I}_{1} \times \mathrm{I}_{2}$ ) such that $\varepsilon(i, j) \neq 0$. Let $\mathrm{I}_{0,2}$ be the set of $(i, k) \in \mathrm{I}_{0} \times \mathrm{I}_{2}$ such that $s_{i}$ is a vertex of $f_{k}$, i.e. $s_{i}=f_{k}\left(\alpha_{\nu}\right)$ for some $\nu \in\{0,1,2\}$.

Set $\mathrm{I}_{0}^{\prime}=\mathrm{I}_{0} \sqcup \mathrm{I}_{1} \sqcup \mathrm{I}_{2}$, and

$$
\begin{cases}s_{i}^{\prime}=s_{i} & \text { if } i \in \mathrm{I}_{0} \\ s_{j}^{\prime}=a_{j}\left(\frac{1}{2}\right) & \text { if } j \in \mathrm{I}_{1} \\ s_{k}^{\prime}=f_{k}(\mathrm{G}) & \text { if } k \in \mathrm{I}_{2}\end{cases}
$$

where G is the centre of mass of T .
Set $\mathrm{I}_{1}^{\prime}=\mathrm{I}_{0,1} \sqcup \mathrm{I}_{1,2} \sqcup \mathrm{I}_{0,2}$. For $(i, j) \in \mathrm{I}_{0,1}$, set

$$
a_{i, j}^{\prime}(t)= \begin{cases}a_{j}\left(\frac{t}{2}\right) & \text { if } s_{i}=a_{j}(0) \\ a_{j}\left(1-\frac{t}{2}\right) & \text { if } s_{i}=a_{j}(1)\end{cases}
$$

For $(j, k) \in \mathrm{I}_{1,2}$ or $\mathrm{I}_{0,2}$, set

$$
a_{j, k}^{\prime}=f_{k} \circ h,
$$

where $h:[0,1] \rightarrow \mathrm{T}$ is the affine map defined by $h(0)=f_{k}^{-1}\left(s_{j}^{\prime}\right)$ and $h(1)=\mathrm{G}$.
Let $\mathrm{I}_{2}^{\prime}$ be the set of $(i, j, k) \in \mathrm{I}_{0} \times \mathrm{I}_{1} \times \mathrm{I}_{2}$ such that $(i, j) \in \mathrm{I}_{0,1}$ and $(j, k) \in \mathrm{I}_{1,2}$. For $(i, j, k) \in \mathrm{I}_{2}^{\prime}$, set

$$
f_{i, j, k}^{\prime}=f_{k} \circ h
$$

where $h: \mathrm{T} \rightarrow \mathrm{T}$ is a direct affine homeomorphism from T onto the triangle with vertices $f_{k}^{-1}\left(s_{i}\right), f_{k}^{-1}\left(s_{j}^{\prime}\right)$, G.

This gives a direct $\mathrm{C}^{1}$ triangulation $\tau^{\prime}=\left(\left(s_{i}^{\prime}\right)_{i \in \mathrm{I}_{0}^{\prime}},\left(a_{j}^{\prime}\right)_{j \in \mathrm{I}_{1}^{\prime}},\left(f_{k}^{\prime}\right)_{k \in \mathrm{I}_{2}^{\prime}}\right)$ of X , called the barycentric subdivision of $\tau$. In this triangulation, each face of $\tau$ has been replaced by 6 faces configured as:


Let $\mathrm{C}^{\prime}$. be the chain complex associated to the triangulation $\tau^{\prime}$. Define a homomorphism $\sigma: \mathrm{C} . \rightarrow \mathrm{C}^{\prime}$. by setting $\sigma\left(\left[f_{k}\right]\right)=\sum\left[f_{i, j, k}^{\prime}\right]$, the sum being over all $(i, j)$ such that $(i, j, k) \in \mathrm{I}_{2}^{\prime}$, and $\sigma\left(\left[a_{j}\right]=\left[a_{i, j}^{\prime}\right]-\left[a_{i^{\prime}, j}^{\prime}\right]\right.$, where $i$ and $i^{\prime}$ are defined by $s_{i}=a_{j}(0)$ and $s_{i^{\prime}}=a_{j}(1), \sigma\left(\left[s_{i}\right]\right)=\left[s_{i}^{\prime}\right]$. It can be checked that $\sigma$ is a morphism of chain complexes.

### 6.6.4

Proposition For all $\nu$, the homomorphism $\sigma$ induces an isomorphism

$$
\sigma_{*}: \mathrm{H}_{\nu}(\mathrm{X}, \tau) \xrightarrow{\approx} \mathrm{H}_{\nu}\left(\mathrm{X}, \tau^{\prime}\right)
$$

Remark This proposition is obvious except for $\nu=1$.
Proof For C. and $\mathrm{C}^{\prime}$, define respective chain subcomplexes $\mathrm{F}_{\lambda} \mathrm{C}$. and $\mathrm{F}_{\lambda} \mathrm{C}^{\prime}$, where $\lambda=0,1,2$, as follows.

$$
\mathrm{F}_{\lambda} \mathrm{C}_{\nu}=\mathrm{C}_{\nu} \text { for } \nu \leqslant \lambda, \quad \mathrm{F}_{\lambda} \mathrm{C}_{\nu}=0 \text { for } \nu>\lambda ;
$$

$\mathrm{F}_{\lambda} \mathrm{C}^{\prime}$. is the submodule of $\mathrm{C}^{\prime}$. generated by the vertices, edges and faces of $\tau^{\prime}$ contained in the $\lambda$-skeleton of $\tau$. More explicitly

$$
\begin{aligned}
& \mathrm{F}_{0} \mathrm{C}_{0}^{\prime} \text { is generated by }\left[s_{i}^{\prime}\right] \text { for } i \in \mathrm{I}_{0}, \\
& \mathrm{~F}_{0} \mathrm{C}_{1}^{\prime}=\mathrm{F}_{0} \mathrm{C}_{2}^{\prime}=0, \\
& \mathrm{~F}_{1} \mathrm{C}_{0}^{\prime} \text { is generated by }\left[s_{i}^{\prime}\right] \text { for } i \in \mathrm{I}_{0} \sqcup \mathrm{I}_{1}, \\
& \mathrm{~F}_{1} \mathrm{C}_{1}^{\prime} \text { is generated by }\left[a_{i, j}^{\prime}\right] \text { for }(i, j) \in \mathrm{I}_{0,1}, \\
& \mathrm{~F}_{1} \mathrm{C}_{2}^{\prime}=0, \\
& \mathrm{~F}_{2} \mathrm{C}_{\nu}^{\prime}=\mathrm{C}_{\nu}^{\prime} .
\end{aligned}
$$

It can be checked that $F_{\lambda} C$. and $F_{\lambda} C^{\prime}$. are subcomplexes, i.e. that

$$
\partial\left(\mathrm{F}_{\lambda} \mathrm{C}_{\nu}\right) \subset \mathrm{F}_{\lambda} \mathrm{C}_{\nu-1} ;
$$

likewise for $\mathrm{C}^{\prime}$. It can also be checked that $\sigma\left(\mathrm{F}_{\lambda} \mathrm{C}.\right) \subset \mathrm{F}_{\lambda} \mathrm{C}^{\prime}$. Set $\mathrm{G}_{\lambda} \mathrm{C} .=\mathrm{F}_{\lambda} \mathrm{C}_{.} / \mathrm{F}_{\lambda-1} \mathrm{C}$. and define $\mathrm{G}_{\lambda} \mathrm{C}^{\prime}$. similarly. The morphism $\sigma$ gives a morphism of chain complexes from $\mathrm{G}_{\lambda} \mathrm{C}$. to $\mathrm{G}_{\lambda} \mathrm{C}^{\prime}$, which in turn gives a homomorphism $\sigma_{*}: \mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C} . \rightarrow \mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C}^{\prime}$. by passing to the homology.

Lemma The homomorphism $\sigma_{*}: \mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C} . \rightarrow \mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C}^{\prime}$. is an isomorphism.
Proof of the Lemma As $\mathrm{G}_{\lambda} \mathrm{C} .=\cdots \rightarrow 0 \rightarrow \mathrm{C}_{\lambda} \rightarrow 0 \rightarrow \cdots, \mathrm{H}_{\lambda} \mathrm{G}_{\lambda} \mathrm{C} .=\mathrm{C}_{\lambda}$ and $\mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C}$. $=0$ for $\nu \neq \lambda$. The chain complex $\mathrm{G}_{\lambda} \mathrm{C}^{\prime}$. is of the form $\bigoplus_{i \in \mathrm{I}_{\lambda}} \mathrm{K}_{i}$, where the chain complex $\mathrm{K}_{i}$ is as follows

$$
\begin{array}{r}
0 \rightarrow \mathbb{Z}^{6} \rightarrow \mathbb{Z}^{6} \rightarrow \mathbb{Z} \rightarrow 0 \text { if } \lambda=2 \\
0 \rightarrow 0 \rightarrow \mathbb{Z}^{2} \rightarrow \mathbb{Z} \rightarrow 0 \text { if } \lambda=1 \\
0 \rightarrow 0 \rightarrow \mathbb{Z} \rightarrow 0 \text { if } \lambda=0
\end{array}
$$

It may be checked that $\mathrm{H}_{\nu} \mathrm{K}_{i}=\mathbb{Z}$ for $\nu=\lambda$ and $\mathrm{H}_{\nu} \mathrm{K}_{i}=0$ for $\nu \neq \lambda$. Hence,

$$
\mathrm{H}_{\lambda} \mathrm{G}_{\lambda} \mathrm{C}_{.}^{\prime}=\mathbb{Z}^{\mathrm{I}_{\lambda}}=\mathrm{C}_{\lambda}=\mathrm{H}_{\lambda} \mathrm{G}_{\lambda} \mathrm{C} . \text { and } \mathrm{H}_{\nu} \mathrm{G}_{\lambda} \mathrm{C}_{.}^{\prime}=0 \text { pour } \nu \neq \lambda .
$$

The above identification can be checked to be indeed given by $\sigma_{*}$, which therefore is an isomorphism.

End of the Proof of the Proposition We show by induction on $\lambda$ that $\sigma_{*}: \mathrm{H}_{\nu} \mathrm{F}_{\lambda} \mathrm{C}$. $\rightarrow \mathrm{H}_{\nu} \mathrm{F}_{\lambda} \mathrm{C}^{\prime}$. is an isomorphism. For $\lambda=-1$, everything vanishes. Consider the commutative diagram of complexes


Passing to the associated long exact sequences gives the following commutative diagram:


By the lemma, $u_{1}$ and $u_{4}$ are isomorphisms, by the induction hypothesis so are $u_{2}$ and $u_{5}$. Applying the five lemma ${ }^{4}$ then shows that $u_{3}$ is an isomorphism, proving our claim. The proposition follows from taking $\lambda=2$.

[^22]
### 6.6.5 Dual of a Chain Complex

We construct a complex $\mathrm{C}^{\prime \prime}$, and a homomorphism $\sigma^{\prime}: \mathrm{C}^{\prime \prime}$. $\rightarrow \mathrm{C}^{\prime}$. which is shown to give an isomorphism onto the homology. The complex $\mathrm{C}^{\prime \prime}$. is not associated to a triangulation, but to a "cell decomposition", a notion that has not been formally introduced. It is a triangulation where triangles are replaced by polygons. The cell decomposition considered, called the dual of the triangulation $\tau$, is obtained by assembling in a cell all the faces of $\tau^{\prime}$ containing a vertex of $\tau$.


We continue with the definition of $\mathrm{C}^{\prime \prime}$ and $\sigma^{\prime}$. Set

$$
\mathrm{I}_{0}^{\prime \prime}=\mathrm{I}_{2}, \quad \mathrm{I}_{1}^{\prime \prime}=\mathrm{I}_{1}, \quad \mathrm{I}_{2}^{\prime \prime}=\mathrm{I}_{0}, \quad \text { and } \quad \mathrm{C}_{\nu}^{\prime \prime}=\mathrm{Z}^{\mathrm{I}_{\nu}^{\prime \prime}}
$$

Let $\left[s_{k}^{\prime \prime}\right]$ be the elements of the canonical basis $\mathrm{C}_{0}^{\prime \prime},\left[a_{j}^{\prime \prime}\right]$ those of $\mathrm{C}_{1}^{\prime \prime}$ and $\left[f_{i}^{\prime \prime}\right]$ those of $\mathrm{C}_{2}^{\prime \prime}$.

Define $\sigma^{\prime}: \mathrm{C}^{\prime \prime} \rightarrow \mathrm{C}^{\prime}$. as follows: for $i \in \mathrm{I}_{0}$,

$$
\sigma^{\prime}\left(\left[f_{i}^{\prime \prime}\right]\right)=\sum\left[f_{i, j, k}^{\prime}\right]
$$

the sum being over all $(j, k)$ such that $(i, j, k) \in \mathrm{I}_{2}^{\prime}$; for $j \in \mathrm{I}_{1}$,

$$
\sigma^{\prime}\left(\left[a_{j}^{\prime \prime}\right]\right)=\left[a_{j, k}^{\prime}\right]-\left[a_{j, k^{\prime}}^{\prime}\right]
$$

where $k$ is the index of the face of $\tau$ on the left of $a_{j}$ and $k^{\prime}$ that of the face on the right; for $k \in \mathrm{I}_{2}, \sigma^{\prime}\left(\left[s_{k}^{\prime \prime}\right]\right)=\left[s_{k}^{\prime}\right]$.

Define $\partial_{\nu}: \mathrm{C}_{\nu}^{\prime \prime} \rightarrow \mathrm{C}_{\nu-1}^{\prime \prime}$ so that $\sigma^{\prime}$ is a morphism of chain complexes. This forces

$$
\partial_{2}\left[f_{i}^{\prime \prime}\right]=\sum-\varepsilon(i, j)\left[a_{j}^{\prime \prime}\right] \text { and } \partial_{1}\left(\left[a_{j}^{\prime \prime}\right]\right)=\left[s_{k}^{\prime \prime}\right]-\left[s_{k^{\prime}}^{\prime \prime}\right],
$$

where $k$ (resp. $k^{\prime}$ ) is the index of the face on the left (resp. on the right) of $a_{j}$.

### 6.6.6

Proposition For all $\nu$, the homomorphism $\sigma^{\prime}$ induces an isomorphism

$$
\sigma_{*}^{\prime}: \mathrm{H}_{\nu} \mathrm{C}_{.}^{\prime \prime} \rightarrow \mathrm{H}_{\nu} \mathrm{C}_{.}^{\prime}=\mathrm{H}_{\nu}\left(\mathrm{X}, \tau^{\prime}\right)
$$

Proof For $\lambda=0,1,2$, define the respective chain subcomplexes $\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime}$. and $\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime}$. of $\mathrm{C}^{\prime}$. and $\mathrm{C}^{\prime \prime}$. as follows:

$$
\mathrm{F}_{\lambda}^{\prime} \mathrm{C}_{\nu}^{\prime \prime}=\mathrm{C}_{\nu}^{\prime \prime} \text { for } \nu \leqslant \lambda ; \mathrm{F}_{\lambda}^{\prime} \mathrm{C}_{\nu}^{\prime \prime}=0 \text { for } \nu>\lambda ;
$$

$\mathrm{F}_{\nu}^{\prime} \mathrm{C}^{\prime}$. is the submodule of $\mathrm{C}^{\prime}$. generated by the vertices, edges and faces of $\tau^{\prime}$ contained in the $\lambda$-skeleton of the dual cell decomposition of $\tau$. More precisely:

$$
\begin{aligned}
& \mathrm{F}_{0}^{\prime} \mathrm{C}_{0}^{\prime} \text { is generated by }\left[s_{k}^{\prime}\right] \text { for } k \in \mathrm{I}_{2}, \\
& \mathrm{~F}_{0}^{\prime} \mathrm{C}_{1}^{\prime}=\mathrm{F}_{0}^{\prime} \mathrm{C}_{2}^{\prime}=0, \\
& \mathrm{~F}_{1}^{\prime} \mathrm{C}_{0}^{\prime} \text { is generated by }\left[s_{j}^{\prime}\right] \text { for } j \in \mathrm{I}_{1} \sqcup \mathrm{I}_{2}, \\
& \mathrm{~F}_{1}^{\prime} \mathrm{C}_{1}^{\prime} \text { is generated by }\left[a_{j, k}^{\prime}\right] \text { for }(j, k) \in \mathrm{I}_{1,2}, \\
& \mathrm{~F}_{1}^{\prime} \mathrm{C}_{2}^{\prime}=0, \\
& \mathrm{~F}_{2}^{\prime} \mathrm{C}_{\nu}^{\prime}=\mathrm{C}_{\nu}^{\prime} .
\end{aligned}
$$

It can be checked that $\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime}$, and $\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime}$ are chain subcomplexes and $\sigma^{\prime}\left(\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime}\right) \subset \mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime}$. Set $\mathrm{G}_{\lambda}^{\prime} \mathrm{C}^{\prime}=\mathrm{F}_{\lambda}^{\prime} \mathrm{C}^{\prime} . / \mathrm{F}_{\lambda-1}^{\prime} \mathrm{C}^{\prime}$. and likewise define $\mathrm{G}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime}$. Then $\mathrm{G}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime}=\cdots \rightarrow 0 \rightarrow$ $\mathrm{C}_{\lambda}^{\prime \prime} \rightarrow 0 \rightarrow \cdots$. As in Proposition 6.6.4, it suffices to show that $\sigma_{*}^{\prime}: \mathrm{H}_{\nu} \mathrm{G}_{\lambda}^{\prime} \mathrm{C}^{\prime \prime} \rightarrow$ $\mathrm{H}_{\nu} \mathrm{G}_{\lambda}^{\prime} \mathrm{C}^{\prime}$ is an isomorphism.

This is immediate for $\lambda=0$, and for $\lambda=1$ the proof is similar to that of Lemma 6.6.4. The chain complex $\mathrm{G}_{2}^{\prime} \mathrm{C}^{\prime}$. is of the form $\sum_{i \in \mathrm{I}_{0}} \mathrm{~K}_{i}^{\prime}$ where $\mathrm{K}_{i}^{\prime}$ is the chain complex $0 \rightarrow \mathbb{Z}^{\mathrm{J}_{2}(i)} \rightarrow \mathbb{Z}^{\mathrm{J}_{1}(i)} \rightarrow \mathbb{Z} \rightarrow 0$, where $\mathrm{J}_{2}(i)$ is the set of $(i, j, k) \in \mathrm{I}_{2}^{\prime}$ with $(j, k) \in$ $\mathrm{I}_{1} \times \mathrm{I}_{2}$, and $\mathrm{J}_{1}(i)$ is the set of $(i, j) \in \mathrm{I}_{0,1} \sqcup \mathrm{I}_{0,2}$ with $j \in \mathrm{I}_{1} \sqcup \mathrm{I}_{2}$. The proposition then follows from the following lemma:
Lemma $\mathrm{H}_{0}\left(\mathrm{~K}_{i}^{\prime}\right)=\mathrm{H}_{1}\left(\mathrm{~K}_{i}^{\prime}\right)=0$, and $\mathrm{H}_{2}\left(\mathrm{~K}_{i}^{\prime}\right)$, generated by

$$
\sum_{(i, j, k) \in \mathrm{J}_{2}(i)}\left[f_{i, j, k}^{\prime}\right]
$$

is isomorphic to $\mathbb{Z}$.
Proof of the Lemma For $(i, j) \in \mathrm{J}_{1}(i)$, the element $\partial_{1}\left[a_{i, j}^{\prime}\right]=\left[s_{j}^{\prime}\right]-\left[s_{i}\right]$ gives -1 in $\mathbb{Z}=\left(\mathrm{K}_{i}^{\prime}\right)_{0}$. Since $s_{i}$ is not isolated, $\mathrm{J}_{1}(i) \neq \varnothing$, and so

$$
\partial_{1}:\left(\mathrm{K}_{i}^{\prime}\right)_{1}=\mathbb{Z}^{\mathrm{J}_{1}(i)} \rightarrow \mathbb{Z}=\left(\mathrm{K}_{i}^{\prime}\right)_{0}
$$

is surjective and $\mathrm{H}_{0}\left(\mathrm{~K}_{i}^{\prime}\right)=0$.

We show that $\mathrm{H}_{1}\left(\mathrm{~K}_{i}^{\prime}\right)=0$. The submodule $\operatorname{Ker} \partial_{1}: \mathbb{Z}^{\mathrm{J}_{1}(i)} \rightarrow \mathbb{Z}$ is generated by all $\left[a_{j}^{\prime}\right]-\left[a_{j^{\prime}}^{\prime}\right]$, where $j, j^{\prime} \in \mathrm{J}_{1}(i)$ (see proof of 6.6.2). Two edges $a_{j}^{\prime}$ and $a_{j^{\prime}}^{\prime}$ with $j, j^{\prime} \in \mathrm{J}_{1}(i)$ are said to be connected if there is a finite sequence $j_{0}=j, j_{1}, \ldots, j_{\mathrm{N}}=$ $j^{\prime}$ in $\mathrm{J}_{1}(i)$ such that $a_{j_{n-1}}^{\prime}$ and $a_{j_{n}}^{\prime}$ are two sides of a same face $f_{k_{n}}^{\prime}$ with $k_{n} \in \mathrm{~J}_{2}(i)$ for $n=1, \ldots, \mathrm{~N}$. Then, $\partial_{2}\left[f_{k_{n}}^{\prime}\right]$ having the same image in $\mathrm{G}_{2}^{\prime} \mathrm{C}^{\prime}$. as $\pm\left(\left[a_{j_{n}}^{\prime}\right]-\left[a_{j_{n-1}}^{\prime}\right]\right)$, $\left[a_{j}^{\prime}\right]-\left[a_{j^{\prime}}^{\prime}\right] \in \operatorname{Im} \partial_{2}:\left(\mathrm{K}_{i}^{\prime}\right)_{2} \rightarrow\left(\mathrm{~K}_{i}^{\prime}\right)_{1}$ if $a_{j}^{\prime}$ and $a_{j^{\prime}}^{\prime}$ are connected. The set $\mathrm{V}=$ $\bigcup_{k \in \mathrm{~J}_{2}(i)} f_{k}^{\prime}(\mathrm{T})$ is a neighbourhood of $s_{i}$ in X . The space $\mathrm{V}-\left\{s_{i}\right\}=\bigcup\left(f_{k}^{\prime}(\mathrm{T})-\left\{s_{i}\right\}\right)$ is connected. Indeed, as X is a topological surface, there is a neighbourhood U of $s_{i}$ such that $\mathrm{U}-\left\{s_{i}\right\}$ is connected, and each $\left(f_{k}^{\prime}(\mathrm{T})-\left\{s_{i}\right\}\right)$ is connected and meets $\mathrm{U}-\left\{s_{i}\right\}$.

A point $x \in \mathrm{~V}-\left\{s_{i}\right\}$ is said to be connected to an edge $a_{j}^{\prime}$ with $j \in \mathrm{~J}_{1}(i)$ if $a_{j}^{\prime}$ is connected to a side $a_{j^{\prime}}^{\prime}$ of a face containing $x$. For all $j \in \mathrm{~J}_{1}(i)$, let $\mathrm{W}_{j}$ be the set of all $x \in \mathrm{~V}-\left\{s_{i}\right\}$ connected to $a_{j}^{\prime}$. The set of $\mathrm{W}_{j}$ partitions $\mathrm{V}-\left\{s_{i}\right\}$ into finite sets, and hence there is only one, and any two arbitrary edges $a_{j}^{\prime}$ and $a_{j^{\prime}}^{\prime}$ with $j, j^{\prime} \in \mathrm{J}_{1}(i)$ are connected. Therefore $\mathrm{H}_{1}\left(\mathrm{~K}_{i}^{\prime}\right)=0$.

As each face $f_{k}^{\prime}$ with $k \in \mathrm{~J}_{2}(i)$ contains 2 edges $a_{j}^{\prime}$ with $j \in \mathrm{~J}_{1}(i)$ (the third one defining an element of $\mathrm{F}_{1} \mathrm{C}_{1}^{\prime}$ ) and as each edge $a_{j}^{\prime}$ with $j \in \mathrm{~J}_{1}(i)$ is contained in 2 faces, the sets $\mathrm{J}_{1}(i)$ and $\mathrm{J}_{2}(i)$ have the same number $n$ of elements. The $\mathbb{Z}$-module $\operatorname{Im} \partial_{2}=\operatorname{Ker} \partial_{1}$, being a submodule of a free $\mathbb{Z}$-module of rank $n$, is free, and $\operatorname{Im} \partial_{1}$ is of rank 1, while $\operatorname{Ker} \partial_{1}$ is of rank $n-1$. Hence $\operatorname{Ker} \partial_{2}$ is free of rank 1. The element $\sum_{k \in \mathrm{~J}_{2}(i)} f_{k}^{\prime}$ is in Ker $\partial_{1}$ and generates it since it is not divisible by an integer $>1$.

This completes the proof of the proposition.

### 6.6.7 Remark

(1) The isomorphisms $\sigma_{*}$ and $\sigma_{*}^{\prime}$ identify $\mathrm{H}_{\nu}\left(\mathrm{X}, \tau^{\prime}\right)=\mathrm{H}_{\nu}\left(\mathrm{C}^{\prime}\right)$ and $\mathrm{H}_{\nu}\left(\mathrm{C}^{\prime \prime}\right)$ with $\mathrm{H}_{\nu}(\mathrm{X}, \tau)=\mathrm{H}_{\nu}(\mathrm{C}$.$) .$
(2) Identifying $\mathrm{C}_{\nu}^{\prime \prime}$ with the dual of $\mathrm{C}_{2-\nu}$, i.e. with $\operatorname{Hom}_{\mathbb{Z}}\left(\mathrm{C}_{2-\nu}, \mathbb{Z}\right), \partial_{2}^{\prime \prime}: \mathrm{C}_{2}^{\prime \prime} \rightarrow$ $\mathrm{C}_{1}^{\prime \prime}$ becomes the opposite of the transpose of $\partial_{1}: \mathrm{C}_{1} \rightarrow \mathrm{C}_{0}$, and $\partial_{1}^{\prime \prime}: \mathrm{C}_{1}^{\prime \prime} \rightarrow \mathrm{C}_{0}^{\prime \prime}$ the transpose of $\partial_{2}: \mathrm{C}_{2} \rightarrow \mathrm{C}_{1}$. Indeed, the entries of the matrix representing $\partial_{\nu}: \mathrm{C}_{\nu} \rightarrow$ $\mathrm{C}_{\nu-1}$ are the incidence numbers; those of the the matrix representing $\partial_{2}^{\prime \prime}$ are $-\varepsilon(i, j)$, while those of the matrix representing $\partial_{1}^{\prime \prime}$ are $\varepsilon(j, k)$.

### 6.6.8

Proposition The $\mathbb{Z}$-module $\mathrm{H}_{1}(\mathrm{X}, \tau)$ is finitely generated and free.
Lemma 6.5 Let A be a principal ring, E and F finitely generated free A -modules and $u: \mathrm{E} \rightarrow \mathrm{F}$ a homomorphism. Consider the transpose of $u^{\top}: \mathrm{F}^{\top} \rightarrow \mathrm{E}^{\top}$ of $u$. The modules Coker $u$ and Coker $u^{\top}$ have isomorphic torsion subsets.

Proof Let $\left(f_{1}, \ldots, f_{q}\right)$ be a basis for F adapted to $\operatorname{Im} u$ (3.5.2 and 3.5.7), and $a_{1}, \ldots, a_{r} \in \mathrm{~A}$ such that $\left(a_{1} f_{1}, \ldots, a_{r} f_{r}\right)$ is a basis for $\operatorname{Im} u$. Let $e_{1}, \ldots, e_{r}$ be representatives of $a_{1} f_{1}, \ldots, a_{r} f_{r}$ in E . Completing to a basis for E by taking a basis $\left(e_{r+1}, \ldots, e_{p}\right)$ for $\operatorname{Ker} u$, the matrix representing $u$ is

$$
{ }_{q-r}\left\{\left(\begin{array}{ccccc}
a_{1} & & 0 & 0 & \\
\vdots & \ddots & \vdots & \vdots & \\
\overbrace{-} & \begin{array}{c}
0 \\
0
\end{array} & a_{r} & 0 & \\
0 & 0 & 0 & 0 \\
0 & & \vdots & \vdots & \ddots
\end{array}\right)\right.
$$

and Coker $u$ is isomorphic to $\bigoplus_{i=1}^{r} \mathrm{~A} /\left(a_{i}\right) \oplus \mathrm{A}^{q-r}$. The matrix representing $u^{\top}$ is

$$
r\left\{\begin{aligned}
p-r
\end{aligned}\left\{\begin{array}{ccccc}
a_{1} & & 0 & 0 & \\
\vdots & \ddots & \vdots & \vdots & \\
\overbrace{2} & \vdots \\
0 & & a_{r} & 0 & \\
0 & & 0 & 0 & \\
\vdots & & \vdots & \vdots & \ddots
\end{array}\right)\right.
$$

Coker $u^{\top}$ is isomorphic to $\bigoplus_{i=1}^{r} \mathrm{~A} /\left(a_{i}\right) \oplus \mathrm{A}^{p-r}$. In both cases, the torsion subset is $\bigoplus_{i=1}^{r} \mathrm{~A} /\left(a_{i}\right)$.

Lemma 6.6 Let A be a principal ring, E. a chain complex of finitely generated free A-modules and $\mathrm{E}^{\top}$ its dual. Then $\mathrm{H}_{n} \mathrm{E}$. and $\mathrm{H}^{n+1} \mathrm{E}^{\top}$ have isomorphic torsion subsets.

Proof There are exact sequences

$$
\begin{gathered}
0 \rightarrow \mathrm{H}_{n} \mathrm{E} . \rightarrow \text { Coker } d_{n+1} \rightarrow \operatorname{Im} d_{n} \rightarrow 0 \\
0 \rightarrow \mathrm{H}^{n+1} \mathrm{E}^{\top} \rightarrow \operatorname{Coker} d^{n} \rightarrow \operatorname{Im} d^{n+1} \rightarrow 0
\end{gathered}
$$

As $\operatorname{Im} d_{n} \subset \mathrm{E}_{n-1}$ is torsion free, $\mathrm{H}_{n} \mathrm{E}$. and Coker $d_{n+1}$ have the same torsion subset. Likewise, $\mathrm{H}^{n+1} \mathrm{E}^{\top}$ and Coker $d^{n}$ have the same torsion subset. As $d^{n}: \mathrm{E}^{n} \rightarrow$ $\mathrm{E}^{n+1}$ is the transpose of $d_{n+1}: \mathrm{E}_{n+1} \rightarrow \mathrm{E}_{n}$, Lemma 6.6 follows from Lemma 6.5, $\square$

Remark The isomorphisms whose existence was proved in Lemmas 6.5 and 6.6 depend on the choice of an adapted basis. There are no natural ones. (see 2.3, Exercise 4, c).

Proof of the Proposition The module $\mathrm{H}_{1}(\mathrm{X}, \tau)=\mathrm{H}_{1} \mathrm{C}$. is finitely generated since it is a quotient of a submodule of $\mathrm{C}_{1}$. The torsion subset of $\mathrm{H}_{1}(\mathrm{X}, \tau)=\mathrm{H}_{1} \mathrm{C}$. is
isomorphic to that of $\mathrm{H}^{2} \mathrm{C}^{\top}$. Changing the sign of the differential, which does not impact the homology, enables the identification of $\mathrm{C}^{\top}$ with $\mathrm{C}^{\prime \prime}$, and

$$
\mathrm{H}^{2} \mathrm{C}^{\top}=\mathrm{H}_{0} \mathrm{C}_{.}^{\prime \prime}=\mathrm{H}_{0} \mathrm{C} .=\mathbb{Z}^{\pi_{0}(\mathrm{X})}
$$

So its torsion subset is 0 . Hence $\mathrm{H}_{1}(\mathrm{X}, \tau)$ is finitely generated and torsion free over $\mathbb{Z}$, and thus is free (3.5.8, Corollary 3.2).

### 6.6.9 Intersection Product

Define a bilinear map $\left(c, c^{\prime \prime}\right) \mapsto c \cdot c^{\prime \prime}$ from $\mathrm{C}^{\nu} \times \mathrm{C}_{2-\nu}^{\prime \prime}$ to $\mathbb{Z}$ by

$$
\begin{aligned}
{\left[s_{i}\right] \cdot\left[f_{i^{\prime \prime}}^{\prime \prime}\right] } & =\delta_{i}^{i^{\prime}}, \\
{\left[a_{j}\right] \cdot\left[a_{j^{\prime \prime}}^{\prime \prime}\right] } & =\delta_{j}^{j^{\prime}}, \\
{\left[f_{k}\right] \cdot\left[s_{k^{\prime \prime}}^{\prime \prime}\right] } & =\delta_{k}^{k^{\prime}},
\end{aligned}
$$

where $\delta_{i}^{i^{\prime}}=1$ if $i=i^{\prime}$ and 0 otherwise.
For $c \in \mathrm{C}_{\nu}$ and $c^{\prime \prime} \in \mathrm{C}_{3-\nu}, \partial c \cdot c^{\prime \prime}=(-1)^{\nu} c \cdot \partial c^{\prime \prime}$. Indeed,

$$
\partial\left[a_{j}\right] \cdot\left[f_{i}^{\prime \prime}\right]=-\varepsilon(i, j)=-\left[a_{j}\right] \cdot \partial\left[f_{i}^{\prime \prime}\right]
$$

and

$$
\partial f_{k} \cdot\left[a_{j}^{\prime \prime}\right]=\varepsilon(j, k)=\left[f_{k}\right] \cdot \partial\left[a_{j}^{\prime \prime}\right] .
$$

Hence restriction and passage to the quotient gives a map from $\mathrm{H}_{\nu} \mathrm{C} . \times \mathrm{H}_{2-\nu} \mathrm{C}^{\prime \prime}$, to $\mathbb{Z}$ or $\mathrm{H}_{\nu}(\mathrm{X}, \tau) \times \mathrm{H}_{2-\nu}(\mathrm{X}, \tau) \rightarrow \mathbb{Z}$, called the intersection product. We also write it $\left(\gamma, \gamma^{\prime}\right) \mapsto \gamma \cdot \gamma^{\prime}$, and $\gamma \cdot \gamma^{\prime}$ is called the intersection number of $\gamma$ and $\gamma^{\prime}$.

Remarks 1. For $\nu=0$ or $2, \mathrm{H}_{0}(\mathrm{X})$ and $\mathrm{H}_{2}(\mathrm{X})$ can be identified with $\mathbb{Z}$ (6.6.2, proposition) if X is connected, the intersection becoming the multiplication $\mathbb{Z} \times \mathbb{Z} \rightarrow \mathbb{Z}$.
2. We have already mentioned in (6.6.1) that $\mathrm{H}_{\nu}(\mathrm{X}, \tau)$ is independent of $\tau$. Hence $\mathrm{H}_{\nu}(\mathrm{X})$ is well defined. For $\nu=1$, the intersection product on $\mathrm{H}_{1}(\mathrm{X})$ can be interpreted as follows. Let $\tau=\left(\mathrm{I}_{0}, \mathrm{I}_{1}, \mathrm{I}_{2},\left(s_{i}\right),\left(a_{j}\right),\left(f_{k}\right)\right)$ and $\tilde{\tau}=\left(\widetilde{\mathrm{I}}_{0}, \widetilde{\mathrm{I}}_{1}, \widetilde{\mathrm{I}}_{2},\left(s_{\tilde{i}}\right),\left(a_{\tilde{j}}\right),\left(f_{\tilde{k}}\right)\right)$ be two mutually transverse $\mathrm{C}^{1}$ triangulations of X . This means that $\left\{s_{i}\right\}_{i \in \mathrm{I}_{0}} \cap\left\{s_{i}\right\}_{\tilde{i} \widetilde{\mathrm{I}}_{0}}=$ $\varnothing$ and that for $j \in \mathrm{I}_{1}$ and $\tilde{j} \in \widetilde{\mathrm{I}}_{1}$, the curves $\mathrm{A}_{j}=a_{j}(] 0,1[)$ and $\mathrm{A}_{\tilde{J}}=a_{\tilde{j}}(] 0,1[)$ intersect transversally at finitely many points. For $z=a_{j}(t)=a_{\tilde{j}}(\tilde{t}) \in \mathrm{A}_{j} \cap \mathrm{~A}_{\tilde{j}}$, set $\theta(z)=1$ if the basis $\left(a_{j}^{\prime}(t), a_{\tilde{j}}^{\prime}(\tilde{t})\right)$ for the tangent space $\mathrm{T}_{z} \mathrm{X}$ consisting of the derived vectors is direct; set $\theta(z)=-1$ if it is inverse, and $a_{j} \cdot a_{\tilde{J}}=\sum_{z \in \mathrm{~A}_{j} \cap \mathrm{~A}_{j}} \theta(z)$. For $\xi, \eta \in \mathrm{H}_{1}(\mathrm{X})$, consider $\xi$ as an element in $\mathrm{H}_{1}(\mathrm{X}, \tau)$ and $\eta$ as an element in $\mathrm{H}_{1}(\mathrm{X}, \tilde{\tau})$. Choose respective representatives $\sum \lambda_{j}\left[a_{j}\right] \in \mathrm{C}_{1}(\mathrm{X}, \tau)$ and $\sum \mu_{\tilde{j}}\left[a_{\tilde{j}}\right] \in \mathrm{C}_{1}(\mathrm{X}, \tilde{\tau})$ of $\xi$ and $\eta$. The intersection product can be shown to be given by

$$
\xi \cdot \eta=\sum \lambda_{j} \mu_{\tilde{j}}\left(\left[a_{j}\right] \cdot\left[a_{\tilde{j}}\right]\right)
$$

Theorem (Poincaré duality) The intersection product defines an isomorphism from $\mathrm{H}_{2-\nu}(\mathrm{X}, \tau)$ onto the dual of $\mathrm{H}_{\nu}(\mathrm{X}, \tau)$.

In other words, for any linear form $u: \mathrm{H}_{\nu}(\mathrm{X}, \tau) \rightarrow \mathbb{Z}$, there is a unique $\gamma^{\prime} \in$ $\mathrm{H}_{2-\nu}(\mathrm{X}, \tau)$ such that $(\forall \gamma \in \mathrm{H}(\mathrm{X}, \tau)) u(\gamma)=\gamma \cdot \gamma^{\prime}$. This theorem is in fact trivial except for $\nu=1$.

Lemma Let A be a principal ring, E. a chain complex of free A-modules such that $\mathrm{H}_{n} \mathrm{E}$. are free for all $n$, and $\mathrm{E}^{\top}$ its dual. Then $\mathrm{H}^{n} \mathrm{E}^{\top}$ can be identified with the dual of $\mathrm{H}_{n} \mathrm{E}$. for all $n$.

Proof There are exact sequences

$$
0 \rightarrow \mathrm{H}_{n} \mathrm{E} . \rightarrow \text { Coker } d_{n+1} \rightarrow \operatorname{Im} d_{n} \rightarrow 0
$$

and

$$
0 \rightarrow \mathrm{H}_{n-1} \text { E. } \rightarrow \text { Coker } d_{n} \rightarrow \operatorname{Im} d_{n-1} \rightarrow 0
$$

These are split since the submodules $\operatorname{Im} d_{n}$ of $\mathrm{E}_{n-1}$ and and $\operatorname{Im} d_{n-1}$ of $\mathrm{E}_{n-2}$ are free. Hence taking the transpose of the former gives an exact sequence $0 \rightarrow\left(\operatorname{Im} d_{n}\right)^{\top} \rightarrow$ $\operatorname{Ker}\left(d_{n+1}^{\top}\right) \rightarrow\left(\mathrm{H}_{n} \mathrm{E} \text {. }\right)^{\top} \rightarrow 0$. As for the second sequence it shows that Coker $d_{n}$ is free. So the exact sequence $0 \rightarrow \operatorname{Im} d_{n} \rightarrow \mathrm{E}_{n-1} \rightarrow$ Coker $d_{n} \rightarrow 0$ is split, and the restriction $\mathrm{E}_{n-1}^{\top} \rightarrow\left(\operatorname{Im} d_{n}\right)^{\top}$ is surjective. Thus $\left(\operatorname{Im} d_{n}\right)^{\top}$ has the same image in $\operatorname{Ker}\left(d_{n+1}^{\top}\right)$ as $\mathrm{E}_{n-1}^{\top}$ under $d_{n}^{\top}$, and $\left(\mathrm{H}_{n} \mathrm{E} \text {. }\right)^{\top}$ can be identified with the $n$-th homology of $E^{\top}$.

Proof of the Theorem Renumbering and changing the sign of the differentials, which leaves the homology invariant, the complex $\mathrm{C}^{\prime \prime}$. can be identified with $\mathrm{C}^{\top}$. By the lemma, $\mathrm{H}_{2-\nu}(\mathrm{X}, \tau)=\mathrm{H}_{2-\nu} \mathrm{C}^{\prime \prime}$. can therefore be identified with the dual of $\mathrm{H}_{\nu}(\mathrm{X}, \tau)$, and this identification can be checked to be given by the intersection product.

### 6.6.10

Proposition The intersection product $\mathrm{H}_{1}(\mathrm{X}, \tau) \times \mathrm{H}_{1}(\mathrm{X}, \tau) \rightarrow \mathbb{Z}$ is an alternating bilinear form.

Proof It suffices to show that if $c \in \mathrm{C}_{1}$ and $c^{\prime \prime} \in \mathrm{C}_{1}^{\prime \prime}$ are such that $\sigma c-\sigma^{\prime} c^{\prime \prime}=\partial b$ with $b \in \mathrm{C}_{2}^{\prime}$, then $c \cdot c^{\prime \prime}=0$. Let $c_{j}$ be the coefficient of $\left[a_{j}\right]$ in $c$, etc. Let $j \in \mathrm{I}_{1}$, and define $i^{\prime}, i^{\prime \prime}, k^{\prime}, k^{\prime \prime}$ by

$$
\varepsilon\left(i^{\prime}, j\right)=-1, \varepsilon\left(i^{\prime \prime}, j\right)=1, \varepsilon\left(j, k^{\prime}\right)=1, \varepsilon\left(j, k^{\prime \prime}\right)=-1
$$

Then

$$
\begin{aligned}
& c_{j}=(\sigma c)_{i^{\prime} j}=-(\sigma c)_{i^{\prime \prime} j}=(\partial b)_{i^{\prime} j}=-(\partial b)_{i^{\prime \prime} j} \\
&=b_{i^{\prime} j k^{\prime}}-b_{i^{\prime} j k^{\prime \prime}}=b_{i^{\prime \prime} j k^{\prime}}-b_{i^{\prime \prime} j k^{\prime \prime}} \\
& c_{j}^{\prime \prime}=\left(\sigma^{\prime} c^{\prime \prime}\right)_{j k^{\prime}}=-\left(\sigma^{\prime} c^{\prime \prime}\right)_{j k^{\prime \prime}}=-(\partial b)_{j k^{\prime}}=(\partial b)_{j k^{\prime \prime}} \\
&=b_{i^{\prime} j k^{\prime}}-b_{i^{\prime \prime} j k^{\prime}}=b_{i^{\prime} j k^{\prime \prime}}-b_{i^{\prime \prime} j k^{\prime \prime}}
\end{aligned}
$$

Hence

$$
\begin{aligned}
2 c_{j} c_{j}^{\prime \prime} & =\left(b_{i^{\prime} j k^{\prime}}-b_{i^{\prime} j k^{\prime \prime}}\right)\left(b_{i^{\prime} j k^{\prime}}+b_{i^{\prime} j k^{\prime \prime}}-\left(b_{i^{\prime \prime} j k^{\prime}}+b_{i^{\prime \prime} j k^{\prime \prime}}\right)\right) \\
& =b_{i^{\prime} j k^{\prime}}^{2}-b_{i^{\prime} j k^{\prime \prime}}^{2}-\left(b_{i^{\prime \prime} j k^{\prime}}-b_{i^{\prime \prime} j k^{\prime \prime}}\right)\left(b_{i^{\prime \prime} j k^{\prime}}+b_{i^{\prime \prime} j k^{\prime \prime}}\right) \\
& =b_{i^{\prime} j k^{\prime}}^{2}-b_{i^{\prime} j k^{\prime \prime}}^{2}-b_{i^{\prime \prime} j k^{\prime}}^{2}+b_{i^{\prime \prime} j k^{\prime \prime}}^{2}=\sum_{(i, k) \in \mathrm{I}_{0} \times \mathrm{I}_{2}}-\varepsilon(i, j) \varepsilon(j, k) b_{i j k}^{2}
\end{aligned}
$$

Therefore, summing over $j$,

$$
\begin{equation*}
2 c \cdot c^{\prime \prime}=\sum_{(i, j, k)}-\varepsilon(i, j) \varepsilon(j, k) b_{i j k}^{2} \tag{*}
\end{equation*}
$$

Note that, for $(i, k) \in \mathrm{I}_{0,2}$, there are two values $j^{\prime}$ and $j^{\prime \prime}$ de $j$ such that $\varepsilon(i, j) \varepsilon(j, k) \neq$ 0 , one with +1 (corresponding to face $f_{i, j, k}^{\prime}$ on the left of edge $a_{i k}^{\prime}$ ) and one with -1 (corresponding to the righthand face). As $(\partial b)_{i k}=0, b_{i j^{\prime} k}=b_{i j^{\prime \prime} k}$. Hence, the terms in the sum $(*)$ mutually cancel each other. So $2 c \cdot c^{\prime \prime}=0 \in \mathbb{Z}$, and $c \cdot c^{\prime \prime}=0$.

### 6.6.11

Theorem and Definition The $\mathbb{Z}$-module $\mathrm{H}_{1}(\mathrm{X}, \tau)$ is free of even rank. The number $g$ defined by $\operatorname{dim} \mathrm{H}_{1}(\mathrm{X}, \tau)=2 g$ is called the genus of X .

Thanks to Propositions 6.6.9 and 6.6.10, it suffices to prove the following lemma:
 alternating bilinear form on E defining an isomorphism from E onto its dual, then $n$ is even.

Proof Let $\left(e_{i}\right)$ be a basis for E and let $\left(e_{i}^{\prime}\right)$ be the dual basis. Let $u$ be an alternating bilinear form defining an isomorphism $\tilde{u}: \mathrm{E} \rightarrow \mathrm{E}^{\top}$. The homomorphism $\tilde{u}$ has matrix $\mathrm{M}=\left(u_{i, j}\right)$ where $u_{i, j}=u\left(e_{i}, e_{j}\right)$. Then $\mathrm{M}^{\top}=-\mathrm{M}$, where $\operatorname{det} \mathbf{M}=\operatorname{det} \mathbf{M}^{\top}=\operatorname{det}(-\mathbf{M})=(-1)^{n} \operatorname{det} \mathbf{M}$. As $\operatorname{det} \mathbf{M}$ is invertible, $n$ is even.

This completes the proof of the theorem.
Remark The genus $g$ of $X$ can be shown not to depend on the choice of the triangulation $\tau$ (6.6, Exercise 1). In what follows, this result will be assumed.

### 6.6.12 The Euler-Poincaré Characteristic

Let A be an integral Noetherian ring and E. a complex of A-modules. Assume that all $\mathrm{H}_{i} \mathrm{E}$. are finitely generated and only finitely many nonzero. Then the Euler-Poincaré characteristic of E . is the number

$$
\chi\left(\mathrm{E}_{.}\right)=\sum(-1)^{i} \mathrm{r} k\left(\mathrm{H}_{i} \mathrm{E} .\right)
$$

Proposition Assume that all $\mathrm{E}_{i}$ are finitely generated and only finitely many nonzero. Then

$$
\chi(\mathrm{E} .)=\sum(-1)^{i} \mathrm{r} k\left(\mathrm{E}_{i}\right)
$$

Proof As rk $\left(\mathrm{H}_{i} \mathrm{E}.\right)=\mathrm{r} k\left(\operatorname{Ker} \partial_{i}\right)-\mathrm{rk} \partial_{i+1}=\mathrm{r} k \mathrm{E}_{i}-\mathrm{r} k \partial_{i}-\mathrm{r} k \partial_{i+1}$,

$$
\begin{aligned}
\chi(\mathrm{E} .) & =\sum(-1)^{i} \mathrm{r} k\left(\mathrm{E}_{i}\right)-\sum(-1)^{i} \mathrm{r} k \partial_{i}-\sum(-1)^{i-1} \mathrm{r} k \partial_{i} \\
& =\sum(-1)^{i} \mathrm{r} k\left(\mathrm{E}_{i}\right)
\end{aligned}
$$

The Euler-Poincaré characteristic of X is the number

$$
\chi(\mathrm{X})=\sum(-1)^{i} \operatorname{dim} \mathrm{H}_{i}(\mathrm{X}, \tau)
$$

By the above proposition, $\chi(\mathrm{X})=k_{0}-k_{1}+k_{2}$, where $k_{0}, k_{1}$ and $k_{2}$ are respectively the number of vertices, edges and faces of $\tau$. If X is connected, then $\chi(\mathrm{X})=2-2 g$ by 6.6 .2 and 6.6.11. If X is not connected, then $\chi(\mathrm{X})=\sum \chi\left(\mathrm{X}_{i}\right)$, where the $\mathrm{X}_{i}$ are the connected components of X.

As seen, the Euler-Poincaré characteristic of a compact Riemann surface is always even.

### 6.6.13 The Riemann-Hurwitz Formula

Theorem Let B be a compact Riemann surface and X a finite analytically ramified $d$-fold covering of B. Then

$$
\chi(\mathrm{X})=d \cdot \chi(\mathrm{~B})-\sum_{x \in \mathrm{R}}\left(e_{x}-1\right)
$$

where R is the set of ramification points of X , and $e_{x}$ the ramification index of $x$.
This theorem is an immediate consequence of 6.5.4, Remark 2.

Corollary 6.5 (a) The number $\sum_{x \in \mathrm{R}}\left(e_{x}-1\right)$ is even.
(b) If X and B are connected, then

$$
g(\mathrm{X})-1=d \cdot(g(\mathrm{~B})-1)+\frac{1}{2} \sum_{x \in \mathrm{R}}\left(e_{x}-1\right)
$$

where $g(\mathrm{X})$ (resp. $g(\mathrm{~B})$ ) is the genus of X (resp. of B ).
Corollary 6.6 Let X and Y be connected compact Riemann surfaces. If $g(\mathrm{Y})>$ $g(\mathrm{X})$, every analytic map from X to Y is constant.

Proof Let $f: \mathrm{X} \rightarrow \mathrm{Y}$ be a non constant map of degree $d$. Then either $g(\mathrm{X})-1 \geqslant$ $d \cdot(g(\mathrm{Y})-1) \quad$ or $\quad g(\mathrm{Y})-1 \geqslant 0 \quad$ since $g(\mathrm{Y})>g(\mathrm{X}) \geqslant 0 ; \quad$ so $\quad g(\mathrm{X})-1 \geqslant$ $d \cdot(g(\mathrm{Y})-1) \geqslant g(\mathrm{Y})-1$, contrary to assumption.

### 6.6.14 Genus, Uniformization

(1) The Riemann sphere is of genus 0 . Every connected compact Riemann surface of genus 0 is isomorphic to the Riemann sphere (6.6, Exercise 7).
(2) Let $\Gamma \subset \mathbb{C}$ be a closed subgroup isomorphic to $\mathbb{Z}^{2}$; the quotient group $\mathbb{C} / \Gamma$ is a connected compact Riemann surface of genus 1 . Such a surface is called a complex torus. Every connected compact Riemann surface of genus 1 can be shown to be isomorphic to a complex torus.
(3) Every connected compact Riemann surface can be shown (6.6, Exercise 6) to be homeomorphic to a surface of the following type (a torus with $g$ holes, where $g$ is the genus of the surface).


Tore à 0 trou $=$ sphère


Tore à 1 trou $=$ tore


Tore à 2 trous


Tore à 3 trous
(4) If $X$ is a connected compact Riemann surface, then it is isomorphic to $\widetilde{X} / G$, where $\widetilde{X}$ is a universal covering of $X$ (hence a simply connected Riemann surface), and $G$ a subgroup of Aut $\widetilde{X}$ acting freely on $\widetilde{X}$.

Every simply connected Riemann surface is isomorphic to the Riemann sphere $\Sigma$, to the complex plane $\mathbb{C}$ or to the unit disk $\mathbb{D}$ (uniformization theorem, 6.6, Exercise 8). If $X$ has genus 0 , then $X \approx \widetilde{X} \approx \Sigma$ and $G=$ Id. If $X$ has genus 1 , then $\widetilde{X} \approx \mathbb{C}$, the group $G$ is isomorphic to $\mathbb{Z}^{2}$ and X is a complex torus. If X has genus $g \geqslant 2$, then $\widetilde{\mathrm{X}} \approx \mathbb{D}$; the automorphisms of $\mathbb{D}$ will be studied in $\S 6.9$.

## Exercises 6.6. (Homology)

1. (Homology and fundamental group)-(a) Let B be a topological space and G a group. A covering of B equipped with a continuous action of G inducing a simply transitive action of G on each fibre is called a principal G -covering of B .
$(\alpha)$ Show that if $\left(\mathrm{B}, b_{0}\right)$ is a connected pointed space admitting a universal covering, then the isomorphism classes of principal pointed G-coverings of $\left(\mathrm{B}, b_{0}\right)$ are in bijective correspondence with the homomorphisms from $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ to G . Study the functoriality of this correspondence in relation to group homomorphisms.
$(\beta)$ Show that the isomorphism classes of principal G-coverings correspond to the conjugation classes of homomorphisms from $\pi_{1}\left(\mathrm{~B}, b_{0}\right)$ to G . Study the case where G is commutative.
(b) $\alpha$ ) Let X be a surface, $\tau$ a triangulation de X , and C . the chain complex associated to $\tau$. Let G be a commutative group. Defined the chain complex $\mathrm{C}^{*}(\mathrm{G})$ by $\mathrm{C}^{\nu}(\mathrm{G})=\operatorname{Hom}_{\mathbb{Z}}\left(\mathrm{C}_{\nu}, \mathrm{G}\right)$.

Show that the isomorphism classes of principal G-coverings of X are in bijective correspondence with the elements of $\mathrm{H}^{1}\left(\mathbf{C}^{\cdot}(\mathrm{G})\right)$ (without necessarily assuming that X is connected)
$(\beta)$ Show that $\mathrm{H}^{\nu}\left(\mathrm{C}^{*}(\mathrm{G})\right)=\operatorname{Hom}_{\mathbb{Z}}\left(\mathrm{H}_{\nu}(\mathrm{X}, \tau) ; \mathrm{G}\right)$.
$(\gamma)$ Suppose that X is connected and let $x_{0} \in \mathrm{X}$. Show that $\mathrm{H}_{1}(\mathrm{X}, \tau)$ and the quotient of $\pi_{1}\left(\mathrm{X}, x_{0}\right)$ by its commutator group are two representatives of the same covariant functor from $\mathbb{Z}-\mathscr{M} \boldsymbol{O} \boldsymbol{d}$ to $\mathscr{E}$ ns. Deduce that these groups are isomorphic. Give an isomorphism by describing its action.
( $\delta$ ) Show that the genus of a Riemann surface is independent of the choice of triangulation.
2. (Hyperelliptic curves)-Let $\mathrm{P} \in \mathbb{C}[\mathrm{X}]$ be a degree $d$ polynomial with distinct roots, $\Gamma=\left\{(x, y) \in \mathbb{C}^{2} \mid y^{2}=\mathrm{P}(x)\right\}$ and $\pi$ the projection $(x, y) \mapsto x$ from $\Gamma$ to $\mathbb{C}$.
(a) Show that $\Gamma$ is a ramified covering of $\mathbb{C}$ extending to a ramified covering $\widehat{\Gamma}$ of $\Sigma$. Show that $\widehat{\Gamma}$ is (resp. is not) ramified over $\infty$ if $d$ is odd (resp. even).
(b) What is the genus of $\widehat{\Gamma}$ ?
3.-Let $d$ be an integer $\geqslant 1$ and $\Gamma=\left\{(x, y) \in \mathbb{C}^{2} \mid x^{d}+y^{d}=1\right\}$.
(a) Show that $\Gamma$ together with $\pi:(x, y) \mapsto x$ is a Galois ramified covering of $\mathbb{C}$. What is its degree? How many ramification points does it have? What is its ramification index?
(b) Show that $\Gamma$ extends to a ramified covering $\widehat{\Gamma}$ of $\Sigma$ unramified over $\infty$. What is the genus of $\widehat{\Gamma}$ ?
4.-Let $\mathrm{P} \in \mathbb{C}[\mathrm{X}, \mathrm{Y}]$ be a degree $d$ polynomial, and $\Gamma$ the set of $(x, y) \in \mathbb{C}^{2}$ such that $\mathrm{P}(x, y)=0$ (algebraic curve with equation P ).
(a) Suppose that the coefficient of $\mathrm{Y}^{d}$ in P is nonzero, and that the polynomial $y \mapsto \mathrm{P}_{d}(1, y)$ has distinct roots, where $\mathrm{P}_{d}$ denotes the homogeneous part of P of degree $d$. Show that the curve $\Gamma$ has $d$ distinct non vertical asymptotes, and that the ramified covering $(\Gamma, \pi)$ of $\mathbb{C}$, where $\pi$ is the projection $(x, y) \mapsto x$, extends to a ramified covering $\widehat{\Gamma}$ of $\Sigma$, unramified over $\infty$.
(b) Let $\sigma \in \mathbb{C}[\mathrm{X}]$ be the discriminant of P considered an element of $(\mathbb{C}[\mathrm{X}])[\mathrm{Y}]$ (3.7.12). What is the degree of the polynomial $\sigma$ ?
(c) Suppose that $\sigma$ only has simple roots. What is the ramification of $\Gamma$ over a root of $\sigma$ ? What is the genus of $\widehat{\Gamma}$ ? Is the ramified covering $\widehat{\Gamma}$ Galois?
5. (Fundamental group of a torus with $p$ holes.) -Let $\mathrm{D}_{0}$ be a disk with centre 0 and radius R in $\mathbb{R}^{2}, \mathrm{D}_{1}, \ldots, \mathrm{D}_{p}$ the disks of centre $\mathrm{C}_{i}=\left(0, c_{i}\right)$ with $-\mathrm{R}<c_{1}<\cdots<$ $c_{p}<\mathrm{R}$, and sufficiently small radius $r$ so that the disks $\mathrm{D}_{i}$ are disjoint and in the interior of $\mathrm{D}_{0}$. Set $a_{0}=(-\mathrm{R}, 0), a_{i}=\left(-r, c_{i}\right)$ and $\mathrm{A}=\mathrm{D}_{0}-\left(\bigcup \stackrel{\circ}{\mathrm{D}}_{i}\right)$.


Let $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ be a continuous function, $>0$ on $\AA$ and $<0$ on $\mathbb{R}^{2}-\mathrm{A}, \mathrm{X}$ the set of $(x, y, z) \in \mathbb{R}^{3}$ such that $z^{2}=f(x, y)$ and $\pi$ the projection $(x, y, z) \mapsto(x, y)$ of X onto A . The surface X is called a torus with $p$ holes.
(a) Let $u_{i}$ be the rectilinear path from $a_{0}$ to $a_{i}, v_{i}$ the loop

$$
t \mapsto\left(-r \cos t, c_{i}-r \sin t\right)
$$

in A. Denote by $\Gamma$ the union of images of all $u_{i}$ and $v_{i}$, and set $\mathrm{U}=\mathrm{A}-\Gamma$. Show that U is contractible, and that there is a contraction

$$
\varphi:[0,1] \times \mathrm{U} \rightarrow \mathrm{U}
$$

onto a point of $\partial \mathrm{D}_{0}$ inducing a contraction of $\partial \mathrm{D}_{0}-\left\{a_{0}\right\}$. Set $\tilde{\mathrm{U}}=\pi^{-1}(\mathrm{U})$. Show that $\widetilde{U}$ is contractible.
(b) Set $\mathrm{X}_{+}=\mathrm{X} \cap\left(\mathbb{R}^{2} \times \mathbb{R}_{+}\right)$and $\mathrm{X}_{-}=\mathrm{X} \cap\left(\mathbb{R}^{2} \times \mathbb{R}_{-}\right)$. Let $u_{i+}$ and $u_{i-}$ be the paths lifting $u_{i}$ in $X_{+}$and $X_{-}$respectively. Set $\widetilde{\Gamma}=\pi^{-1}(\Gamma)$, and in $\pi_{1}\left(\widetilde{\Gamma}, a_{0}\right)$, let $\alpha_{i}$ be the class of the loop $u_{i-}^{-1} \cdot u_{i+}$ and $\beta_{i}$ the class of $u_{i+}^{-1} \cdot v_{i} \cdot u_{i+}$. Show that $\pi_{1}\left(\widetilde{\Gamma}, a_{0}\right)$ is a free group having $\left(\alpha_{1}, \beta_{1}, \ldots, \alpha_{p}, \beta_{p}\right)$ as basis.
(c) Let $\varepsilon>0$, and V the set of points of A at a distance $<\varepsilon$ from $\Gamma$. Show that, if $\varepsilon$ is chosen sufficiently small, which will be henceforth assumed, then there is a deformation retraction of V onto $\Gamma$ (4.3.16).

Set $\tilde{\mathrm{V}}=\pi^{-1}(\mathrm{~V})$. Show that there is deformation retraction of $\tilde{\mathrm{V}}$ onto $\widetilde{\Gamma}$.
(d) Show that $\mathrm{U} \cap \mathrm{V}=\mathrm{V}-\Gamma$ is homeomorphic to the product of a closed interval $\underset{\sim}{\mathrm{U}} \cap \widetilde{\mathrm{V}}$ and open interval, and that $\tilde{\mathrm{U}} \cap \widetilde{\mathrm{V}}$ is homeomorphic to an annulus. Let $a_{0}^{\prime} \in$ $\widetilde{\mathrm{U}} \cap \widetilde{\mathrm{V}} \cap \partial \mathrm{D}_{0}$ be a point over $a_{0}$, and $\gamma \in \pi_{1}\left(\widetilde{\mathrm{U}} \cap \tilde{\mathrm{V}}, a_{0}^{\prime}\right)$ an element generating the group. What is the image of $\gamma$ in $\pi_{1}\left(\widetilde{\mathrm{~V}}, a_{0}^{\prime}\right)=\pi_{1}\left(\widetilde{\Gamma}, a_{0}\right)$ ?
(e) Using Van Kampen's theorem, show that $\pi_{1}\left(\mathrm{X}, a_{0}\right)$ can be identified with the quotient of a free group on $\left(\alpha_{1}, \beta_{1}, \ldots, \alpha_{p}, \beta_{p}\right)$ by the normal subgroup generated by

$$
\alpha_{p} \beta_{p}^{-1} \alpha_{p}^{-1} \beta_{p} \ldots \alpha_{1} \beta_{1}^{-1} \alpha_{1}^{-1} \beta_{1} .
$$

(f) Find a basis $\alpha_{1}^{\prime}, \beta_{1}^{\prime}, \ldots, \alpha_{p}^{\prime}, \beta_{p}^{\prime}$ for $\pi_{1}\left(\widetilde{\Gamma}, a_{0}\right)$ such that $\pi_{1}\left(\mathrm{X}, a_{0}\right)$ can be identified with the quotient of a free group with this basis by the normal subgroup generated by

$$
\alpha_{1}^{\prime} \beta_{1}^{\prime} \alpha_{1}^{\prime-1} \beta_{1}^{\prime-1} \ldots \alpha_{p}^{\prime} \beta_{p}^{\prime} \alpha_{p}^{\prime-1} \beta_{p}^{\prime-1}
$$

6. (Classification of surfaces.) - We keep the notation of the previous exercise, but denote the set A by $\mathrm{A}_{p}$ and the torus X with $p$ holes $\mathrm{X}_{p}$.

A set of the form $\mathrm{S}=\mathrm{H} \cap \mathrm{D}_{0}$ will be called a segment in $\mathrm{A}_{p}$ when H is the open half-plane such that $\mathrm{S} \neq \varnothing, \overline{\mathrm{S}} \neq \mathrm{D}_{0}, \overline{\mathrm{~S}} \cap \mathrm{D}_{i}=\varnothing$ for $1 \leqslant i \leqslant p$.
(a) Show that if $S$ is a segment of $A_{p}$, then the inverse image $\widetilde{\mathrm{S}}$ of S in $\mathrm{X}_{p}$ is homeomorphic to an open disc, and that $\widetilde{\widetilde{S}}=\pi^{-1}(\overline{\mathrm{~S}})$ is homeomorphic to a closed disc.
(b) Let $\mathrm{S}_{1}, \ldots, \mathrm{~S}_{q}$ be segments of $\mathrm{A}_{p}$ such that $\overline{\mathrm{S}}_{1}, \ldots, \overline{\mathrm{~S}}_{q}$ are disjoint, $\mathrm{A}_{p, q}=$ $\mathrm{A}_{p}-\left(\mathrm{S}_{1} \cup \cdots \cup \mathrm{~S}_{q}\right)$, and $\mathrm{X}_{p, q}$ the inverse image of $\mathrm{A}_{p, q}$ in $\mathrm{X}_{p}$. Then $\mathrm{X}_{p, q}$ will be said to be a torus with $p$ holes without $q$ disks. The boundary $\partial \mathbf{X}_{p, q}$ is the boundary of $\mathrm{X}_{p, q}$ in $\mathrm{X}_{p}$.

Show that $\mathrm{X}_{p, q}$ is a surface with topological boundary, i.e. that any point has a neighbourhood homeomorphic to an open subset of $\mathbb{R}^{2}$ or of $\mathbb{R}_{+} \times \mathbb{R}$.

Show that the homeomorphism class of $\mathrm{X}_{p, q}$ does not depend on the integers $p$ and $q$, and that if $\sigma$ is a permutation of $\{1, \ldots, q\}$, there is a homeomorphism $h$ of $\mathrm{X}_{p}$ inducing a homeomorphism of $\mathrm{X}_{p, q}$, and such that $h\left(\widetilde{\mathrm{~S}}_{j}\right)=\widetilde{\mathrm{S}}_{\sigma(j)}$ (reduce to the case where $\sigma$ is a transposition, and all $\mathrm{S}_{j}$ are in $\left.[r, \mathrm{R}] \times \mathbb{R}\right)$.
(c) Let $p, q, p^{\prime}, q^{\prime}$ be four integers, $\mathrm{I}_{1}$ and $\mathrm{I}_{2}$ arcs (i.e. subspaces homeomorphic to $[0,1]$ ) contained in the boundary of $\mathrm{X}_{p, q}$ and $\mathrm{X}_{p^{\prime}, q^{\prime}}$ respectively, D a disc, $\mathrm{J}_{1}$ and $\mathrm{J}_{2}$ disjoint arcs in the boundary of D , and Y the quotient of the disjoint union $\mathrm{X}_{p, q} \sqcup$
$\mathrm{X}_{p^{\prime}, q^{\prime}} \sqcup \mathrm{D}$ by an equivalence relation identifying $\mathrm{I}_{1}$ with $\mathrm{J}_{1}$, and $\mathrm{I}_{2}$ with $\mathrm{J}_{2}$, along homeomorphisms.

Show that the homeomorphism class of Y only depends on the integers $p, q, p^{\prime}$, $q^{\prime}$, and that Y is homeomorphic to $\mathrm{X}_{p+p^{\prime}, q+q^{\prime}-1}$.
(d) Let $w_{j}$ be the rectilinear path in $\mathrm{A}_{p, q}$ along the boundary of $\mathrm{S}_{j}$ keeping $\mathrm{S}_{j}$ to the right; $w_{j+}$ and $w_{j-}$ paths in $\mathrm{X}_{p, q}$ lifting $w_{j}$ to $\mathrm{X}_{p+}$ and $\mathrm{X}_{p-}$ respectively. The connected component of the boundary of $X_{p, q}$ consisting of the boundary $\widetilde{S}_{j}$ is oriented by the loop $w_{j+} \circ w_{j-}^{-1}$. Let D be a disk with oriented boundary. Let $\mathrm{I}_{1}$ and $\mathrm{I}_{2}$ be two disjoint arcs in $\partial \mathrm{X}_{p, q}, \mathrm{~J}_{1}$ and $\mathrm{J}_{2}$ two disjoint arcs in $\partial \mathrm{D}, h_{1}: \mathrm{I}_{1} \rightarrow \mathrm{~J}_{1}$ and $h_{2}: \mathrm{I}_{2} \rightarrow \mathrm{~J}_{2}$ homeomorphisms reversing the orientation. Let Y be the quotient of $\mathrm{X}_{p, q} \sqcup \mathrm{D}$ by the equivalence relation identifying $\mathrm{I}_{1}$ with $\mathrm{J}_{1}$ by $h_{1}$ and $\mathrm{I}_{2}$ with $\mathrm{J}_{2}$ by $h_{2}$. Show that Y is homeomorphic to $\mathrm{X}_{p, q+1}$ if $\mathrm{I}_{1}$ and $\mathrm{I}_{2}$ are in the same connected component of $\partial \mathbf{X}_{p, q}$, and to $\mathrm{X}_{p+1, q-1}$ otherwise.
(e) Let X be a compact surface with an oriented triangulation $\tau$ (i.e. satisfying the conclusions of Proposition 6.5.6, which is the case if X is a Riemann surface and $\tau$ $\mathrm{C}^{1}$ and direct). Let $f_{1}, \ldots, f_{k}$ be the faces of the dual cell decomposition of $\tau$ (see 6.6.5), and $\mathrm{Y}_{\nu}$ the union of the faces $f_{1}, \ldots, f_{\nu}$. Show by induction on $\nu$ that $\mathrm{Y}_{\nu}$ is the disjoint union of surfaces with boundaries of the form $\mathrm{X}_{p, q}$. Deduce that if X is nonempty and connected, then it is homeomorphic to a torus with $p$ holes for some value of $p$.
7.-Let $X$ be a connected compact Riemann surface and $x \in X$. Suppose $X$ has genus 0 .
(a) Show that all finite unramified coverings of X are trivial (the Riemann-Hurwicz theorem may be used).
(b) Show that all homomorphisms from $\pi_{1}(\mathrm{X}, x)$ to the additive group $\mathbb{C}$ are trivial (6.6, exercise 1 , may be used).
(c) With the notation of ( 6.2 , Exercise $1, \mathrm{D}, \mathrm{c}$ ), show that $\Omega=0$. Applying the result of E , (c) of this exercise, deduce that there is a meromorphic function $f$ on X , holomorphic on $\mathrm{X}-\{x\}$ with a simple pole at $x$. Show that $f$ is necessarily an isomorphism from X onto the Riemann sphere.
8. (Proof of the uniformization theorem) - Let X be a non compact, connected Riemann surface. The aim is to show that if X is simply connected, then X is isomorphic to $\mathbb{C}$ or to the disk $\mathbb{D}$.

The language of singular homology is freely used: for any space $\mathrm{Y}, \mathrm{H}_{k}(\mathrm{Y})$ is the $k$-dimensional homology group of Y with coefficients in $\mathbb{Z}$. If Y is connected and locally simply path connected and if $y_{0} \in \mathrm{Y}$, then the group $\mathrm{H}_{1}(\mathrm{Y})$ can be identified with the quotient of $\pi_{1}\left(\mathrm{Y}, y_{0}\right)$ by its commutator group (see 2.7, Exercise 1 ). If Y is a non compact connected surface with boundary, then $\mathrm{H}_{2}(\mathrm{Y})=\mathrm{H}_{2}(\mathrm{Y}, \partial \mathrm{Y})=0$.

The aim is to show that $X$ is isomorphic to $\mathbb{C}$ or $\mathbb{D}$ if $H_{1}(X)=0$, a weaker assumption than $\pi_{1}(\mathrm{Y})=0$ (hence a stronger result).

The surface X need not be assumed to be the countable union of compact spaces. This indeed is not assumed but only for the challenge. Most applications (for example if X is a universal covering of a compact surface) this is known beforehand. This mainly entails proceeding with care.

## A. Pieces

A 2-dimensional connected compact $\mathbb{R}$-submanifold P with $\mathrm{C}^{\infty}$ boundary in X is called a piece. The boundary (resp. interior) of P is written $\partial \mathrm{P}$ (resp. $\stackrel{\circ}{\mathrm{P}}$ ).

The piece P is said to be full if $\mathrm{X}-\mathrm{P}$ has no relatively compact connected component. Given an arbitrary piece $P$, the union of $P$ and of the relatively compact connected components is said to be "P filled" and is written $\hat{\mathrm{P}}$. Show that it is a full piece.
(a) Show that every compact subset of $X$ is contained in a full piece.
(b) Suppose that $H_{1}(X)=0$ and that P is a full piece in X . Show that $\mathrm{H}_{2}(\mathrm{X}, \mathrm{P})=$ $H_{2}(X-\stackrel{\circ}{P}, \partial P)=0$. Using the exact sequence $H_{2}(X, P) \rightarrow H_{1}(P) \rightarrow H_{1}(X)$, show that $\mathrm{H}_{1}(\mathrm{P})=0$.

## B. Harmonic functions

If U is an open subset of X and $h: \mathrm{U} \rightarrow \mathbb{R}$ a $\mathrm{C}^{2}$ function, then the Laplacian of $h$ is the differential 2-form $\Delta h$ on U whose expression in a $\mathbb{C}$-analytic chart is $\left(\frac{\partial^{2} \dot{h}}{\partial x^{2}}+\frac{\partial^{2} \dot{h}}{\partial y^{2}}\right) d x \wedge d y$, where $\dot{h}$ is the expression of $h$ (check that this form is independent of the choice of the chart). When $\Delta h=0, h$ is called harmonic.

Show that the real part of a holomorphic function is harmonic, and that if $U$ is simply connected and open, then every harmonic function $h$ on $U$ is the real part of a holomorphic function $f$ on U , unique up to addition of an imaginary constant. Then $g=\operatorname{Im} f$ is said to be a harmonic conjugate to $h$. Show that all harmonic functions are $\mathbb{R}$-analytic and satisfy the maximum principle.

## C. Green's functions

Given a piece P and a point $a \in \stackrel{\circ}{\mathrm{P}}$, a function $\mathrm{G}: \mathrm{P}-\{a\} \rightarrow \mathbb{R}$ is a Green's function on P with respect to $a$ if
(G1) G is continuous on $\mathrm{P}-\{a\}$ and vanishes on $\partial \mathrm{P}$;
(G2) G is harmonic on $\stackrel{\circ}{\mathrm{P}}-\{a\}$;
(G3) If $\varphi$ is a $\mathbb{C}$-analytic chart centered at $a$, the function $G$ takes the form $-\log |\varphi|+h$ in the neighbourhood of $a$, where $h$ is harmonic in the neighbourhood of $a$ (check that this condition is independent of the choice of the chart $\varphi$ ).

Show uniqueness: there is at most one Green's function on P with respect to $a$.
We next aim to show existence (note that X has been assumed to be non compact, and so $\partial \mathrm{P} \neq \varnothing$. There are two ways of proceeding. They will be described in ( $\mathrm{C}^{\prime}$ ) and ( $\mathrm{C}^{\prime \prime}$ ).

## C'. Poisson's equation

(a) Show that, for any $\mathrm{C}^{\infty}$ differential 2-form $\omega$ with compact support in $\stackrel{\circ}{\mathrm{P}}$, there is a $\mathrm{C}^{\infty}$ function $h: \mathrm{P} \rightarrow \mathbb{R}$ such that $\Delta h=\omega$ and $h=0$ on $\partial \mathrm{P}$ (this requires quite powerful tools from analysis, see [9]).
(b) Let $\varphi$ be a $\mathbb{C}$-analytic chart for X centered at $a$, defined in a neighbourhood V of $a$ contained in $\stackrel{\circ}{\mathrm{P}}$, and $\eta: \mathrm{V} \rightarrow \mathbb{R}$ a $\mathrm{C}^{\infty}$ function with compact support in V , equal to 1 in a neighbourhood $\mathrm{V}^{\prime}$ of $a$. Define the function $g$ on $\mathrm{P}-\{a\}$ by $g=-\eta \log |\varphi|$ on $\mathrm{V}-\{a\}$ and $g=0$ on $\mathrm{P}-\mathrm{V}$, and the differential form $\omega$ on P by $\omega=\Delta g$ on $\mathrm{P}-\{a\}$ and $\omega=0$ on $\mathrm{V}^{\prime}$. Show that there is a function $h: \mathrm{P}-\{a\} \rightarrow \mathbb{R}$ such that $\Delta h=\omega$ and $h=0$ on $\partial \mathbf{P}$. Show that $\mathrm{G}=g-h$ is then a Green's function on P with respect to $a$.

## $C^{\prime \prime}$. Perron families

Show that any continuous function $u: S^{1} \rightarrow \mathbb{R}$ has a unique continuous extension $h: \overline{\mathbb{D}} \rightarrow \mathbb{R}$ harmonic on $\mathbb{D}$. Show that $h(0)=\int_{0}^{1} u\left(e^{2 i \pi t}\right) d t$ and that the gradient of $h$ satisfies $\left|\nabla_{0} h\right| \leq \sup _{z \in \mathrm{~S}^{1}}|u(z)|$.

Let V be an open subset of $\mathbb{C}$ and $h: \mathrm{V} \rightarrow \mathbb{R}$ a continuous function. Show that $h$ is harmonic if and only if for any closed disk $\overline{\mathrm{D}}_{c, r}$ in $\mathrm{V}, h(c)=\int_{0}^{1} h(c+r$. $\left.e^{2 i \pi t}\right) d t$. If for every closed disk $\overline{\mathrm{D}}_{c, r}$ in $\mathrm{V}, h(c) \leq \int_{0}^{1} h\left(c+r \cdot e^{2 i \pi t}\right) d t(\operatorname{resp} . h(c) \geq$ $\int_{0}^{1} h\left(c+r \cdot e^{2 i \pi t}\right) d t$ ), $h$ is called subharmonic (resp. superharmonic). Show that a locally subharmonic function is subharmonic.

Let U be an open subset of X . A continuous function $h: \mathrm{U} \rightarrow \mathbb{R}$ is said to be subharmonic (resp. superharmonic) if its expression in every $\mathbb{C}$-analytic chart is subharmonic (resp. superharmonic).
(a) A set $\mathscr{F}$ of continuous functions $\mathrm{U} \rightarrow \mathbb{R}$ is a Perron family on U if
(P1) $h \in \mathscr{F} \Rightarrow h$ is subharmonic;
(P2) $\left(h_{1} \in \mathscr{F}\right.$ and $\left.h_{2} \in \mathscr{F}\right) \Rightarrow \sup \left(h_{1}, h_{2}\right) \in \mathscr{F}$;
(P3) For $h \in \mathscr{F}$ and $\Delta \subset \mathrm{U}$ in the domain of the chart $\varphi$, such that $\varphi(\Delta)$ is a closed disc, the continuous function $\tilde{h}$ which agrees with $h$ on $\mathrm{U}-\Delta$ and is harmonic on $\stackrel{\circ}{\Delta}$ is in $\mathscr{F}$.

Show that if U is connected and $\mathscr{F}$ is a Perron family on U , then $\sup _{h \in \mathscr{F}} h$ is either harmonic or identically $+\infty$ (Perron's theorem).
(b) Let $u: \partial \mathrm{P} \rightarrow \mathbb{R}$ be a continuous function. Let $\mathscr{F}_{u}$ be the set of continuous functions $h: \mathrm{P} \rightarrow \mathbb{R}$, subharmonic on $\stackrel{\circ}{\mathrm{P}}$ and such that $h \leq u$ on $\partial \mathrm{P}$. Show that the restrictions of these $h \in \mathscr{F}_{u}$ to $\stackrel{\circ}{\mathrm{P}}$ form a Perron family on $\stackrel{\circ}{\mathrm{P}}$, and that $h_{u}=\sup _{h \in \mathscr{F}_{u}} h$ is a continuous function on P , equal to $u$ on $\partial \mathrm{P}$ and harmonic on $\stackrel{\circ}{\mathrm{P}}$ (solution to the Dirichlet problem).

In particular, if Q is a piece with non connected boundary and B a connected component of $\partial \mathrm{Q}$, then there is a function $h_{\mathrm{B}}$ continuous on Q , harmonic on $\stackrel{\circ}{\mathrm{Q}}, 1$ on B and 0 on all other connected components of $\partial \mathrm{Q}$.
(c) Let $\varphi: \mathrm{V} \xrightarrow{\approx} \mathrm{D}_{r}$ be a chart for X centered at $a$ with $\mathrm{V} \subset \stackrel{\circ}{\mathrm{P}}, r^{\prime}$ and $r^{\prime \prime}$ such that $0<r^{\prime}<r^{\prime \prime}<r$. Set $\mathrm{V}^{\prime}=\varphi^{-1}\left(\mathrm{D}_{r^{\prime}}\right)$ and $\mathrm{V}^{\prime \prime}=\varphi^{-1}\left(\mathrm{D}_{r^{\prime \prime}}\right), \Gamma^{\prime}=\partial \mathrm{V}^{\prime}$ and $\Gamma^{\prime \prime}=\partial \mathrm{V}^{\prime \prime}$. Let $h_{\Gamma^{\prime}}$ be the continuous function on $\mathrm{P}-\mathrm{V}^{\prime}$, harmonic on $\stackrel{\circ}{\mathrm{P}}-\overline{\mathrm{V}^{\prime}}, 1$ on $\partial \mathrm{V}^{\prime}$ and 0 on $\partial \mathrm{P}$. For $m>0$, define $\tilde{\lambda}_{m}: \mathrm{V}-\{a\} \rightarrow \mathbb{R}$ by $\tilde{\lambda}_{m}=-\log |\varphi|+\log r^{\prime}+m$, and
$\lambda_{m}: \mathrm{P}-\{a\} \rightarrow \mathbb{R}$ by $\lambda_{m}=\tilde{\lambda}_{m}$ on $\overline{\mathrm{V}^{\prime}}-\{a\}$ and $\lambda_{m}=m \cdot h_{\Gamma^{\prime}}$ on $\mathrm{P}-\mathrm{V}^{\prime}$, so that $\lambda_{m}=m$ on $\partial \mathrm{V}^{\prime}$. Show that for sufficiently large $m, m . h_{\Gamma^{\prime}} \leq \tilde{\lambda}_{m}$ on $\Gamma^{\prime \prime}$, and that $\lambda_{m}$ is therefore superharmonic. Choose $m$ satisfying these properties.

Let $\mathscr{G}$ be the set of continuous functions on $\mathrm{P}-\{a\}$, subharmonic on $\stackrel{\circ}{\mathrm{P}}-\{a\}$, vanishing on $\partial \mathrm{P}$ and bounded above by $-\log |\varphi|+\mathrm{O}(1)$ in the neighbourhood of $a$. Show that $g \in \mathscr{G} \Rightarrow g \leq \lambda_{m}$ and that $\mathrm{G}=\sup _{g \in \mathscr{G}} g$ is a Green's function on P with respect to $a$.

## D. Uniformization of a piece

Here P is assumed to be a piece in X such that $\mathrm{H}_{1}(\mathrm{P})=0$ and $a \in \stackrel{\circ}{\mathrm{P}}$. The aim is to show the existence of a homeomorphism $\phi_{\mathrm{P}}$ from P onto the closed disk $\overline{\mathbb{D}}$ inducing a $\mathbb{C}$-analytic isomorphism from $\stackrel{\circ}{\mathrm{P}}$ onto $\mathbb{D}$, with $\phi_{\mathrm{P}}(a)=0$.
(a) Let $\Delta$ be a neighbourhood of $a$ transformed into a closed disk by a chart centered at $a$. By considering the exact sequence

$$
\mathrm{H}_{1}(\mathrm{P}) \rightarrow \mathrm{H}_{1}(\mathrm{P}, \Delta) \rightarrow \mathrm{H}_{0}(\Delta) \rightarrow \mathrm{H}_{0}(\mathrm{P}),
$$

show that $\mathrm{H}_{1}(\mathrm{P}-\{a\}, \Delta-\{a\})=\mathrm{H}_{1}(\mathrm{P}, \Delta)=0$. Show that the natural map $\mathrm{H}_{1}(\partial \Delta) \rightarrow \mathrm{H}_{1}(\mathrm{P}-\{a\})$ is surjective. Show that if $\alpha$ is a closed differential 1-form on $\mathrm{P}-\{a\}$ then, for any loop $\gamma$ in $\mathrm{P}-\{a\}$, the value of the integral $\int_{\gamma} \alpha$ is an integer multiple of $\int_{\partial \Delta} \alpha$.
(b) Show that Green's function $\mathrm{G}=\mathrm{G}_{\mathrm{P}, a}$ has a $\mathbb{R} / 2 \pi \mathbb{Z}$-valued harmonic conjugate, i.e. that there is a continuous function $\Theta: \mathrm{P}-\{a\} \rightarrow \mathbb{R} / 2 \pi \mathbb{Z}$ such that any lifting $\tilde{\Theta}: U \rightarrow \mathbb{R}$ of $\Theta$ to a simply connected open set is a harmonic conjugate of G.

Show that $\phi=e^{\mathrm{G}+i \Theta}$, extended by $\phi(a)=0$, has the desired property.

## E. Erdös-Kocbe

(a) Let $f: \mathbb{D} \rightarrow \mathbb{D}$ be a holomorphic map such that $f(0)=0$ and $\left|f^{\prime}(0)\right|=\rho$ is near 1. Set $g(z)=\frac{f(z)}{z}$ and $g(0)=\rho$. Show that $g$ is a holomorphic map from $\mathbb{D}$ to $\mathbb{D}$. Show that, for $|z| \leq r,|g(z)-\rho| \leq r \frac{1-\rho^{2}}{1-\rho r}$ and $|f(z)-\rho z| \leq r^{2} \frac{1-\rho^{2}}{1-\rho r}$.

Let R and $\mathrm{R}^{\prime}$ be such that $0<\mathrm{R}<\mathrm{R}^{\prime}$, and let $\psi: \mathrm{D}_{\mathrm{R}} \rightarrow \mathrm{D}_{\mathrm{R}^{\prime}}$ be a holomorphic map such that $\psi(0)=0$ and $\psi^{\prime}(0)=1$. Show that

$$
|\psi(z)-z| \leq \frac{\mathrm{R}^{\prime 2}-\mathrm{R}^{2}}{\mathrm{R}^{\prime}(\mathrm{R}-r)}
$$

(b) Let $f: \mathrm{D}_{\mathrm{R}} \rightarrow \mathbb{C}$ be an injective holomorphic map such that $f(0)=0$ and $f^{\prime}(0)=1$. Show that $f\left(\mathrm{D}_{\mathrm{R}}\right) \supset \mathrm{D}_{\mathrm{R} / 4}$ (Kæbe quarter theorem). Show that $f\left(\mathrm{D}_{\mathrm{R} / 16}\right) \subset$ $\mathrm{D}_{\mathrm{R} / 4}$, and that the function $g$ defined by $g(z)=\frac{f(z)-z}{z^{2}}$ is bounded above by $80 / \mathrm{R}$ on $\mathrm{D}_{\mathrm{R} / 16}$, and so $|f(z)-z| \leq 80 \frac{r^{2}}{\mathrm{R}}$ for $|z| \leq r \leq \frac{\mathrm{R}}{16}$.

## F. Uniformization of X

Suppose X is a non compact, connected Riemann surface such that $\mathrm{H}_{1}(\mathrm{X})=0$ and $a \in \mathrm{X}$. Choose $v(\neq 0)$ in the tangent space $\mathrm{T}_{a} \mathrm{X}$. For each full piece P such that $a \in \stackrel{\circ}{\mathrm{P}}$, there is a radius $\mathrm{R}(\mathrm{P})$ and an isomorphism $\Phi_{\mathrm{P}}: \stackrel{\circ}{\mathrm{P}} \rightarrow \mathrm{D}_{\mathrm{R}(\mathrm{P})}$ such that the image of $v$ under the tangent linear map $\mathrm{T}_{a} \Phi_{\mathrm{P}}$ is 1 . For each P , the pair $\left(\mathrm{R}_{\mathrm{P}}, \Phi_{\mathrm{P}}\right)$ is unique.

Let $\mathrm{R}(\mathrm{X})$ be the upper bound of $\mathrm{R}(\mathrm{P})$ for all full pieces P such that $a \in \stackrel{\circ}{\mathrm{P}}$. The aim is to show that $X$ is isomorphic to $\mathbb{D}$ if $R(X)<\infty$ and to $\mathbb{C}$ if $R(X)=\infty$.
(a) Let $\left(\mathrm{P}_{n}\right)$ be a sequence of full pieces such that $\mathrm{P}_{n} \subset \stackrel{\circ}{\mathrm{P}}_{n+1}$ for all $n$ and the radii $\mathrm{R}_{n}=\mathrm{R}\left(\mathrm{P}_{n}\right)$ tend to $\mathrm{R}=\mathrm{R}(\mathrm{X})$. Set $\mathrm{X}^{\prime}$ to be the union of all $\stackrel{\circ}{\mathrm{P}}_{n}$.

Using the inequalities proved in ( E ), show that, for any compact subset K of $\mathrm{X}^{\prime}$, for sufficiently large $n, \Phi_{n}$ form a Cauchy sequence with respect to uniform convergence on K . Show that passing to the limit gives an isomorphism from $\mathrm{X}^{\prime}$ onto $D_{R}$ if $R<\infty$, and onto $\mathbb{C}$ if $R=\infty$.
(b) The aim is to show by contradiction that $\mathrm{X}^{\prime}=\mathrm{X}$. Let $b \in \mathrm{X}-\mathrm{X}^{\prime}$. Construct a sequence of full pieces $\left(\tilde{\mathrm{P}}_{n}\right)$ such that, for all $n, \tilde{\mathrm{P}}_{n}$ is in the interior of $\tilde{\mathrm{P}}_{n+1}$ and $\mathrm{P}_{n} \cup\{b\}$ in the interior of $\tilde{\mathrm{P}}_{n}$. This gives an open subset $\tilde{X}^{\prime}$ strictly containing $\mathrm{X}^{\prime}$, and an isomorphism $\Psi$ from $D_{R}$ or from $\mathbb{C}$ onto one of its open strict subsets, tangent to the identity at 0 . Deduce a contradiction.
(c) Using the transfinite line (1.6, Exercise 1) or (4.9, Exercise 3), give examples of surfaces having an $\mathbb{R}$-analytic structure, but not any $\mathbb{C}$-analytic structure.

### 6.7 Finite Automorphism Groups of Riemann Surfaces

### 6.7.1 Quotient Riemann Surface

Proposition Let X be a connected Riemann surface and G a finite automorphism group of X . Then the quotient $\mathrm{Y}=\mathrm{X} / \mathrm{G}$ is a topological surface and Y has a unique $\mathbb{C}$-analytic structure $\sigma$ such that the canonical map $\chi: \mathrm{X} \rightarrow \mathrm{Y}$ is analytic. Let $\mathrm{Y}^{\prime}$ be an open subset of Y and Z a Riemann surface; a continuous map $f: \mathrm{Y}^{\prime} \rightarrow \mathrm{Z}$ is analytic with respect to $\sigma$ if an only if $f \circ \chi: \chi^{-1}\left(\mathrm{Y}^{\prime}\right) \rightarrow \mathrm{Z}$ is analytic.

Lemma Let X be a connected Riemann surface, $x_{0}$ a point of X and H a finite automorphism group of X fixing $x_{0}$. Then H is cyclic and there is chart $\varphi$ for X centered at $x_{0}$, with domain U preserved by H , such that the expressions of the elements de H are the maps $z \mapsto a \cdot z$, where $a$ is a d-th root of unity, $d$ denoting the order of H .

Proof Let $\varphi$ be a chart for X centered at $x_{0}$, with domain $\mathrm{U}_{0}$. Set $\mathrm{U}_{1}=\bigcap_{g \in \mathrm{H}} g \mathrm{U}_{0}$, so that $\mathrm{U}_{1}$ is preserved by H . For $x \in \mathrm{U}_{1}$, set $p(x)=\prod_{g \in \mathrm{H}} \varphi(g \cdot x)$. The function $p$ is analytic on $\mathrm{U}_{1}$ and has $x_{0}$ as a zero of order $d$. Hence in a neighbourhood $\mathrm{U}_{2}$ of $x_{0}$ which may be assumed to be preserved by H , it can be written as $p=\varphi_{1}^{d}$, where $\varphi_{1}$ has $x_{0}$ as zero of order 1 . Then the function $\varphi_{1}$ induces a homeomorphism
from a neighbourhood $\mathrm{U}_{3}$ of $x_{0}$, which may be assumed to be preserved by H , onto a neighbourhood of 0 in $\mathbb{C}$, containing the open disk with radius $r$. As $p(g \cdot x)=p(x)$ for $g \in \mathrm{H}$, the open set U of $x$ such that $\left|\varphi_{1}(x)\right|<r$, i.e. $|p(x)|<r^{d}$, is preserved by H , and $\varphi_{1}$ induces a chart for X centered at $x_{0}$, with domain U .

For $x \in \mathrm{U}-\left\{x_{0}\right\}$ and $g \in \mathrm{H}, \varphi_{1}(g \cdot x)^{d}=p(g \cdot x)=p(x)=\varphi_{1}(x)^{d}$, and so $\varphi_{1}(g \cdot x)=a_{g, x} \varphi_{1}(x)$, where $a_{g, x}$ is a $d$-th root of unity. The map $x \mapsto a_{g, x}$ is continuous, hence locally constant, hence constant since $\mathrm{U}-\left\{x_{0}\right\}$ is connected, and $\varphi_{1}(g \cdot x)=a_{g} \varphi_{1}(x)$ for $x \in \mathrm{U}$. The map $g \mapsto a_{g}$ from H to the group $\mu_{d}$ of $d$-th roots of unity is a homomorphism. The expression of $g$ in the chart $\varphi_{1}$ is multiplication by $a_{g}$. If $a_{g}=1$, then the map $g$ agrees with $1_{\mathrm{X}}$ on U , hence, by analytic continuation, also on X. The homomorphism $g \mapsto a_{g}$ is therefore injective and as H and $\mu_{d}$ are of order $d$, it is bijective.

Proof of the Proposition Let $y \in \mathrm{Y}$, and $\chi^{-1}(y)=\left\{x_{1}, \ldots, x_{r}\right\}$. Denote by $S$ the stabilizer of $x_{1}$; the set $\chi^{-1}(y)$ can be identified with $\mathrm{G} / \mathrm{S}$. Let $\varphi_{1}$ be a chart centered at $x_{1}$, with domain $\mathrm{U}_{1}$, satisfying the properties of the lemma for S . As $\mathrm{U}_{1}$ is preserved by S , the open set $g \cdot \mathrm{U}_{1}$ only depends on the class of $g$ in $\mathrm{G} / \mathrm{S}$. Denote it by $\mathrm{U}_{k}$ if $g\left(x_{1}\right)=x_{k}$. If necessary by shrinking $\mathrm{U}_{1}$, all $\mathrm{U}_{k}$ may be assumed to be disjoint. The open set $\mathrm{W}=\bigcup \mathrm{U}_{k}$ is saturated and so $\mathrm{V}=\chi(\mathrm{W})$ is open in Y . The map $\chi$ induces a homeomorphism from $\mathrm{U}_{1} / \mathrm{S}$ onto V . There is a unique map $\psi: \mathrm{V} \rightarrow \mathbb{C}$ making the diagram

commutative. This map $\psi$ is a homeomorphism from V onto an open disk of $\mathbb{C}$, i.e. a chart for Y centered at $y$. Besides the following property holds: let $\mathrm{V}^{\prime}$ be an open subset of V and Z a Riemann surface; if $f: \mathrm{V}^{\prime} \rightarrow \mathrm{Z}$ is a continuous map then, $f \circ \psi^{-1}: \psi\left(\mathrm{V}^{\prime}\right) \rightarrow \mathrm{Z}$ is analytic if and only if so is $f \circ \chi: \chi^{-1}\left(\mathrm{~V}^{\prime}\right) \rightarrow \mathrm{Z}$.

The charts thereby obtained form an analytic atlas. The analytic structure $\sigma$ thus defined on Y has the required properties.

If $\sigma^{\prime}$ is another analytic structure on Y such that $\chi: \mathrm{X} \rightarrow\left(\mathrm{Y}, \sigma^{\prime}\right)$ is analytic, then the identity map from $(\mathrm{Y}, \sigma)$ to $\left(\mathrm{Y}, \sigma^{\prime}\right)$ is analytic and bijective, and so is an isomorphism, and $\sigma^{\prime}=\sigma$.

Remark The above proof shows that, for all $x \in \mathrm{X}$, the ramification index of $x$ for $\chi: \mathrm{X} \rightarrow \mathrm{Y}$ is the order $d$ of the stabilizer of $x$.

### 6.7.2 Genus $g \geqslant 2$ Case

We show in 6.9.17 that if X is a connected compact Riemann surface of genus $\geqslant 2$, then the automorphism group of X is finite. Here, without using this result, we prove
the following theorem (which, given finiteness, provides an upper bound for the number of automorphisms).

Theorem Let X be a connected compact Riemann surface with genus $g \geqslant 2$ and G a finite automorphism group of X of order $n$. Then

$$
n \leqslant 84 \cdot(g-1)
$$

Proof $\operatorname{Set} \mathrm{Y}=\mathrm{X} / \mathrm{G}$ and let $g^{\prime}$ be the genus of Y. For each $y \in \mathrm{Y}$, the points of $\chi^{-1}(y)$ have conjugate stabilizers, and so the same ramification index $s_{y}$, and $n=r_{y} s_{y}$, where $r_{y}$ is the number of distinct points over $y$. Let $\Delta$ be the ramification set of X ; if $\Delta=\left\{y_{1}, \ldots, y_{k}\right\}$, write $s_{i}$ for $s_{y_{i}}$. As

$$
g-1=n \cdot\left(g^{\prime}-1\right)+\frac{1}{2} \sum_{y \in \Delta} \frac{n}{s_{y}}\left(s_{y}-1\right)=n\left(\left(g^{\prime}-1\right)+\frac{1}{2} \sum_{y \in \Delta}\left(1-\frac{1}{s_{y}}\right)\right),
$$

$n=\frac{g-1}{\mathrm{~A}}$ with $\mathrm{A}=g^{\prime}-1+\frac{1}{2} \sum_{1}^{k}\left(1-\frac{1}{s_{i}}\right)$, where $k$ is the cardinality of $\Delta$. We find a lower bound for A . Now, $\mathrm{A}>0$ holds. If $g^{\prime}>1$, then $\mathrm{A} \geqslant 1$, if $g^{\prime}=1$, then $k>0$, so that $\mathrm{A} \geqslant \frac{1}{4}$. Suppose that $g^{\prime}=0$. Then

$$
2 \mathrm{~A}=k-2-\sum \frac{1}{s_{i}},
$$

and so $k \geqslant 3$ since $\mathrm{A}>0$.
For $k=3$,

$$
2 \mathrm{~A}=1-\sum \frac{1}{s_{i}}
$$

The greatest value $<1$ of $\frac{1}{s_{1}}+\frac{1}{s_{2}}+\frac{1}{s_{3}}$, with integers $s_{i} \geqslant 2$, is $1-\frac{1}{42}$. It is reached for $(2,3,7)$, and so $A \geqslant \frac{1}{84}$.

For $k=4$,

$$
2 \mathrm{~A}=2-\sum \frac{1}{s_{i}}
$$

the greatest value $<2$ of $\sum_{1}^{4} \frac{1}{s_{i}}$, with integers $s_{i} \geqslant 2$, is $2-\frac{1}{6}$. It is reached for $(2,2,2,3)$, and so $\mathrm{A} \geqslant \frac{1}{12}$.

For $k \geqslant 5$,

$$
\mathrm{A}=-1+\frac{1}{2} \sum_{1}^{k}\left(1-\frac{1}{s_{i}}\right) \geqslant-1+\frac{k}{4} \geqslant \frac{1}{4} .
$$

In all cases, $A \geqslant \frac{1}{84}$, and so $n \leqslant 84(g-1)$.

Remark As will be seen in section 8, this bound is not reached for $g=2$. There is a ramified covering X of the Riemann sphere (6.8, Exercise 1 ) for which it is for $g=3$; in section 10 we construct X as a quotient of the disk.

### 6.7.3 Genus 1 Case

Let X be a connected compact Riemann surface of genus 1 and G a finite automorphism group of X of order $n$. Set $\mathrm{Y}=\mathrm{X} / \mathrm{G}$ and let $g^{\prime}$ be the genus of Y . As

$$
\begin{gathered}
n\left(g^{\prime}-1+\frac{1}{2} \sum_{k}\left(1-\frac{1}{s_{i}}\right)\right)=0 \\
g^{\prime}=1-\frac{1}{2} \sum_{1}^{k}\left(1-\frac{1}{s_{i}}\right)
\end{gathered}
$$

since $n \neq 0$ since it is the order of a group. Hence, either $g^{\prime}=1$ or 0 .
If $g^{\prime}=1, k=0$, i.e. G acts freely on X . In fact, as already mentioned, X is a complex torus, i.e. isomorphic to $\mathbb{C} / \Gamma$ with $\Gamma \approx \mathbb{Z}^{2}$. Then $G$ can be identified with a subgroup of $\mathbb{C} / \Gamma$ acting by translations. The order $n$ of $G$ is arbitrary.

If $g^{\prime}=0$, then

$$
\sum_{1}^{k}\left(1-\frac{1}{s_{i}}\right)=2
$$

and so $k=3$ or 4 since $\frac{1}{2} \leqslant 1-\frac{1}{s_{i}}<1$.
for $k=4, s_{i}=2$ for all $i$.
For $k=3, \frac{1}{s_{1}}+\frac{1}{s_{2}}+\frac{1}{s_{3}}=1$. The possible solutions are $(2,3,6),(2,4,4)$ and $(3,3,3)$. These solutions can be shown to be effectively realizable (6.7, Exercise 1).

### 6.7.4 Genus 0 Case

Let X be a connected compact Riemann surface of genus 0 and G a finite automorphism group of X of order $n$. Set $\mathrm{Y}=\mathrm{X} / \mathrm{G}$; the surface Y is necessarily of genus 0 . The Riemann-Hurwitz formula gives

$$
\sum_{1}^{k}\left(1-\frac{1}{s_{i}}\right)=2-\frac{2}{n}
$$

As $\frac{1}{2} \leqslant 1-\frac{1}{s_{i}}<1$, then $1<k<4$, and so $k=2$ or 3 if $n>1$.

If $k=2$, then $\frac{1}{s_{1}}+\frac{1}{s_{2}}=\frac{2}{n}$. But as $s_{i}$ is the order of a subgroup of G, $s_{i} \leqslant n$, and so $s_{1}=s_{2}=n$. There are two points $x_{1}, x_{2}$ in X fixed by G , the group G is cyclic of order $n$ (6.7.1, Lemma), and G acts freely on $\mathrm{X}-\left\{x_{1}, x_{2}\right\}$.

Example. X is the Riemann sphere, G the set of maps $z \mapsto \omega z$, where $\omega$ is an $n$-th root of unity. In fact, this example is typical of the general situation (6.7, Exercise 3).

If $k=3$, then $\frac{1}{s_{1}}+\frac{1}{s_{2}}+\frac{1}{s_{3}}=1+\frac{2}{n}$. The solutions of this equation with $s_{i} \geqslant 2$ are:
(2, 2, $d$ ) with $n=2 d$,
$(2,3,3)$ with $n=12$,
$(2,3,4)$ with $n=24$,
$(2,3,5)$ with $n=60$.
Let us study the case $(2,2, d), n=2 d$ (dihedral case): Let $y_{3} \in \mathrm{Y}$ correspond to $s_{3}=d$. There are two points $x^{\prime}$ and $x^{\prime \prime}$ in the class $y_{3}$. The stabilizer H of $x^{\prime}$ is cyclic of order $d$ and index 2 in G, and hence normal; the stabilizer of $x^{\prime \prime}$, being a conjugate of H , is H . Let $\sigma \in \mathrm{G}-\mathrm{H}$, then $\sigma\left(x^{\prime}\right)=x^{\prime \prime}$. The automorphism $\sigma$ has at least one fixed $x_{1}$ since there is no nonzero finite group $\mathrm{G}^{\prime}$ acting freely on X (otherwise $k\left(\mathrm{G}^{\prime}\right)=0$ ). Since $\chi^{-1}\left(y_{3}\right)=\left\{x^{\prime}, x^{\prime \prime}\right\}, \chi\left(x_{1}\right) \neq y_{3}$; hence it may be assumed that $\chi\left(x_{1}\right)=y_{1}$.

As $s_{1}=2$, the stabilizer of $x_{1}$ has order 2 , and so $\sigma$ is of order 2 .
Hence there is a cyclic subgroup H of the group G of order $d$ such that all elements of $\mathrm{G}-\mathrm{H}$ are of order 2; this is only possible for the dihedral group, i.e. the crossed product of $\mathbb{Z} /(2)$ and $\mathbb{Z} /(d), \mathbb{Z} /(2)$ acting on $\mathbb{Z} /(d)$ by $a \mapsto-a$ (see 6.10.1).

Cases $(2,3,3),(2,3,4)$ and $(2,3,5)$ are realizable as finite automorphism groups of the Riemann sphere, and uniquely so up to conjugation. The groups obtained are respectively isomorphic to $\mathfrak{A}_{4}$ (tetrahedral group), $\mathfrak{S}_{4}$ (cubic or octahedral group) $\mathfrak{A}_{5}$ (dodecahedral or icosahedral group) (Exercises 3, 4 and 5).

## Exercises 6.7. (Finite Automorphism Groups)

1. (Automorphisms of complex tori)—Let X be a complex torus, i.e. a surface of the form $\mathbb{C} / \Gamma$, where $\Gamma$ is a subgroup of $\mathbb{C}$ isomorphic to $\mathbb{Z}^{2}$.
(a) Let $f: \mathrm{X} \rightarrow \mathrm{X}$ be an automorphism with respect to the Riemann surface structure. Show that there is a holomorphic function $\tilde{f}: \mathbb{C} \rightarrow \mathbb{C}$ such that the diagram

commutes. Show that $\tilde{f}^{\prime}(x+\gamma)=\tilde{f}^{\prime}(x)$ for $\gamma \in \Gamma$ (where $\tilde{f}^{\prime}$ denotes the derivative of $\tilde{f}$ ). Using the maximum principle (6.1, Exercise 3), deduce that $\tilde{f}$ is of the form $x \mapsto a x+b$, with $a \Gamma=\Gamma$.
(b) Show that the number $\nu$ of nonzero elements of $\Gamma$ with minimum absolute value can be 2,4 or 6 . For how many elements $a \in \mathbb{C}$, does $a \Gamma=\Gamma$ ?
(c) Let $\mathrm{G}_{2}$ be the group consisting of 1 and $x \mapsto-x$. How many points in X are fixed by $\mathrm{G}_{2}$ ? What is the genus of $\mathrm{X} / \mathrm{G}_{2}$ ? How many elements does the ramification set of $\mathrm{X} \rightarrow \mathrm{X} / \mathrm{G}_{2}$ have, and what is their ramification index?
(d) Same questions for the group $\mathrm{G}_{4}$ generated by $x \mapsto i x$ in the case $\nu=4$, for the groups $\mathrm{G}_{3}$ generated by $x \mapsto e^{\frac{2 i \pi}{3}} x$ and $\mathrm{G}_{6}$ generated by $x \mapsto e^{\frac{i \pi}{3}} x$ in the case $\nu=6$.
(e) Let $G$ be a finite automorphism group of $X$, and $\mathrm{G}^{\prime}$ the subgroup consisting of the translations in $G$. Set $X^{\prime}=X / G^{\prime}$. Show that $X^{\prime}$ is a complex torus, that $X$ is an (unramified) covering of $X^{\prime}$, and that if $\mathrm{G}^{\prime} \neq \mathrm{G}$, then the Riemann surface of $\mathrm{X}^{\prime}$, equipped with the equivalence relation induced by $G$, is isomorphic to one of the situations studied in (c) and (d).
2.-Let Y be a complex torus, $a$ and $b$ distinct points of Y.
(a) Show that $\pi_{1}(\mathrm{Y}-\{a\}, b)$ is a free group with a 2 element basis.
(b) Construct a group G of order 8 , with two elements $u$ and $v$ such that $u v u^{-1} v^{-1}$ is of order 2.
(c) Construct a ramified Galois covering X of Y of degree 8 whose ramification set is $\{a\}$, with ramification index 2 over $a$.
(d) What is the genus of Y? What is the order of the group $\operatorname{Aut}_{\mathrm{X}}(\mathrm{Y})$ ?
2. (Groups $\mathrm{SO}_{3}$ and Aut $\Sigma$ )—Let $\mathrm{SL}_{2} \mathbb{C}$ be the group of $2 \times 2$ matrices with complex entries having determinant 1.
(a) Show that the automorphism group Aut $\Sigma$ of the Riemann sphere can be identified with $\mathrm{SL}_{2} \mathbb{C} /\{+1,-1\}$ (see 6.3.3).
(b) Let G be a finite subgroup of Aut $\Sigma$. Show that there is a Hermitian form $h$ on $\mathbb{C}^{2}$ such that all elements $g \in \mathrm{G}$ are induced by an element $\tilde{g} \in \mathrm{SL}_{2} \mathbb{C}$ preserving $h$.
(c) The stereographic projection (see 6.1.2) identifies $\Sigma$ with the sphere $S^{2} \subset \mathbb{R}^{3}$.

Show that the group $\mathrm{SO}_{3}$ of direct isometries of $\mathrm{S}^{2}$ can be identified with a subgroup of Aut $\Sigma$.
(d) Show that for any finite subgroup $G$ of $\operatorname{Aut} \Sigma$, there is a subgroup $\mathrm{G}^{\prime}$ conjugate to G contained in the group $\mathrm{SO}_{3}$ of isometries of $\mathrm{S}^{2}$ preserving the orientation.

Show that if G is cyclic (resp. dihedral), $\mathrm{G}^{\prime}$ can regarded as consisting of rotations about the axis $\mathrm{O} z$ (resp. contained in the group generated by the rotations about the $z$-axis $\mathrm{O} z$ and the symmetry about the $x y$-plane).
4. (Subgroups of $\mathrm{SO}_{3}$ and regular polyhedra)—Let S be the Euclidean unit sphere of $\mathbb{R}^{3}$ and G a nontrivial finite subgroup of $\mathrm{SO}_{3}$.
(a) Let $g \in \mathrm{G}$ be a non-identity element, D the axis of rotation $g$. Show that, if D is invariant under G , then G is either cyclic or dihedral (see 6.10.1). Show that, if all non-identity element of G have order 2 , then G is either cyclic of order 2 or dihedral
of order 4 , isomorphic to $\mathbb{Z} /(2) \times \mathbb{Z} /(2)$, consisting of 1 and of the symmetries about three orthogonal lines.

In the following, it is assumed that G is neither cyclic nor dihedral. Let $h \in \mathrm{G}$ be an element of order $k>2$, D the axis of rotation $h, a$ one of the two points of $\mathrm{D} \cap \mathrm{S}$, H the subgroup generated by $h$. Let $\Omega$ be the set $\mathrm{G} \cdot a$ of the images of $a$, and $\delta$ the minimum distance between two distinct points of $\Omega$; two points of $\Omega$ are said to be connected if their distance is $\delta$.
(b) Show that there is a point connected to $a$. Noting that the set of points of S whose distance from $a$ is $\delta$ is a circle with radius $<\delta$, show that there are at most 5 points connected to $a$. Deduce that the order $k$ of H can only be 3,4 or 5 , and that H acts transitively on the set of points connected to $a$.
(c) Show that G acts transitively on the set of pairs of connected points. If $b$ and $c$ are connected, show that the symmetry about the line generated by $b+c$ is contained in G.
(d) For connected $b$ and $c$, let the $\operatorname{arc} \gamma_{b, c}$ be the shortest path from $b$ to $c$ in S , and $\Gamma$ the union of $\gamma_{b, c}$. Show that $\Gamma$ is invariant under $G$ and that $G$ acts transitively on the connected components of $S-\Gamma$.
(e) Let C be a connected component of $\mathrm{S}-\Gamma$. Show that the stabilizer $\mathrm{H}^{\prime}$ of C in G acts transitively on $\Phi=\overline{\mathrm{C}} \cap \Omega$. Show that $\Phi$ is contained in a plane, and consists of the vertices of a regular polygon F . Show that the order $k^{\prime}$ of $\mathrm{H}^{\prime}$ must be either 3, 4 or 5 .
(f) Show that the angles of F are $<\frac{2 \pi}{k}$. Deduce that

$$
(k-2)\left(k^{\prime}-2\right)<4
$$

What are the possibilities for the pair $\left(k, k^{\prime}\right)$ ?
(g) Show that the convex envelope P of $\Omega$ is a polyhedron whose boundary is $\bigcup_{g \in \mathrm{G}} g \mathrm{~F}$, and that G is the stabilizer of P in $\mathrm{SO}_{3}$. Show that the numbers $k$ and $k^{\prime}$ define P up to isometry and G up to conjugation. How many conjugation classes of finite subgroups that are neither cyclic nor dihedral are there in $\mathrm{SO}_{3}$ ?
(h) Show that G is isomorphic to the quotient $\Lambda$ of the free group $\mathrm{L}(\alpha, \beta, \gamma)$ by the normal subgroup generated by $\alpha^{k}, \beta^{k^{\prime}}, \gamma^{2}$ and $\alpha \beta \gamma$ (use the proof given in 6.10).
5. (Dodecahedral and icosahedral group) - Let $P$ be a regular dodecahedral, $A$ the set of edges of P and G the group of direct isometries of P . Let C be a set with 5 elements, called "colours".
(a) Colour the edges of P using all the colours of C (i.e. define a map $\gamma$ from A to C ) in such a way that for each face $f$ of P and each side $a$ of $f$, the edge $a^{\prime}$ not contained in $f$ terminating at the vertex of $f$ opposite to $a$ is of the same colour as $a$. Show that this colouring is unique up to permutation of C .
(b) Show that for all $g \in \mathrm{G}$, there is a permutation $\tilde{g}$ of C making the diagram

commute.
(c) Show that $g \mapsto \tilde{g}$ is an isomorphism from $G$ onto the group $\mathfrak{A}_{5}$ of even permutations of C .
(d) Show that the set of faces of $P$ can be identified with the quotient of $\mathfrak{A}_{5}$ by the group of circular permutations.
6. (Construction of a Riemann surface of genus 2 with automorphism group isomorphic to $\mathrm{SL}_{2} \mathrm{~F}_{3}$ )-(a) Let T be a regular tetrahedron. Show that the group H of isometries of T preserving the orientation is isomorphic to the group $\mathfrak{A}_{4}$ of even permutations of the vertices of $T$. Set $B=T / H$. Show that $T$ is a ramified covering of $B$ and that $B$ is homeomorphic to the Riemann sphere.

Identify B with the Riemann sphere by a homeomorphism sending the vertices of $T$ onto $\infty$.
(b) Show that there is a ramified 2 -fold covering X of T whose ramification set is the set of the midpoints of the edges of T , and that it is unique up to isomorphism. Show that for any isometry $f$ of T , there is an $f$-morphism $\tilde{f}$ of X . Deduce that X is a ramified Galois covering of $B$. Set $G=\operatorname{Aut}_{B} X$. Show that $H$ is isomorphic to a quotient of G. Find the order of G.
(c) Let $L$ be the inverse image in $X$ of $\infty \in B$, and $V=\{O\} \sqcup L$, so that $V$ has 9 elements.

Let $F_{3}$ be the field $\mathbb{Z} /(3)$. Show that $V$ has a unique $F_{3}$-vector space structure such that O is the origin and the affine lines are the sets obtained

- by taking O and the 2 points of X over a vertex of T ,
or
- by taking the images of 3 vertices of T under a continuous section of X over the face containing these vertices.
(d) Let $\mathrm{GL}(\mathrm{V})$ be the automorphism group of the vector space V and $\mathrm{SL}(\mathrm{V})$ the subgroup of $\mathrm{GL}(\mathrm{V})$ consisting of the elements with determinant 1 . An element $g \in \mathrm{GL}(\mathrm{V})$ defines a permutation of the set $\mathrm{P}(\mathrm{V})$ of lines through O . Show that $g \in \mathrm{SL}(\mathrm{V})$ if and only if this permutation is even.
(e) Show that each element of G defines an automorphism of V, giving a homomorphism $\varphi: \mathrm{G} \rightarrow \mathrm{GL}(\mathrm{V})$.

Show that $\varphi(\mathrm{G}) \subset \mathrm{SL}(\mathrm{V})$ and that $\varphi$ is injective, so that $\varphi$ is an isomorphism from $G$ onto $\operatorname{SL}(\mathrm{V})$.

### 6.8 Automorphism Groups: The Genus 2 Case

### 6.8.1 Homogeneity of Topological Surfaces

Proposition Let X be a topological surface, A a closed subset of X, $x$ and $y$ two points in the same connected component of $\mathrm{X}-\mathrm{A}$. Then there is a homeomorphism $h: \mathrm{X} \rightarrow \mathrm{X}$ restricting to the identity on A and such that $h(x)=y$.

Proof (a) A particular case. Suppose that $\mathrm{U}=\mathrm{X}-\mathrm{A}$ is homeomorphic to the disk $\mathbb{D}$, and that there is a homeomorphism $\varphi: \mathrm{U} \rightarrow \mathbb{D}$ extending to a homeomorphism from $\overline{\mathrm{U}}$ onto $\overline{\mathbb{D}}$. It then suffices to show that if $a, b \in \mathbb{D}$, then there is a homeomorphism $h: \overline{\mathbb{D}} \rightarrow \overline{\mathbb{D}}$ inducing the identity on $\mathrm{S}^{1}=\partial \overline{\mathbb{D}}$ and such that $h(a)=b$. Define $f_{a}:[0,1] \times \mathrm{S}^{1} \rightarrow \overline{\mathbb{D}}$ by $f_{a}(t, u)=(1-t) a+t u$. The map $f_{a}$ is surjective and, as $[0,1] \times S^{1}$ is compact, it identifies $\overline{\mathbb{D}}$ with the quotient space of $[0,1] \times S^{1}$ by the equivalence relation identifying the points of $\{0\} \times S^{1}$. Define $f_{b}$ likewise. It has the same property. There is a homeomorphism $h$ making the diagram

commute and with the desired properties.
(b) General case. Write $x \sim y$ if there is a homeomorphism $h: \mathrm{X} \rightarrow \mathrm{X}$ inducing the identity on A and such that $h(x)=y$. It is an equivalence relation. The particular case implies that all $x \in \mathrm{X}-\mathrm{A}$ have a neighbourhood U which does not meet A and such that all $y \in \mathrm{U}$ is equivalent to $x$. The classes of points of $\mathrm{X}-\mathrm{A}$ form a partition of X - A into open subsets. Any two points in a connected component are in the same class.

Remark This proof applies to all dimensions.

### 6.8.2

Corollary 6.7 Let X be a connected topological surface, $\Delta \subset \mathrm{X}$ a finite set, $x$ and $y$ points of $\mathrm{X}-\Delta$. Then there is a homeomorphism $h: \mathrm{X} \rightarrow \mathrm{X}$ inducing the identity on $\Delta$ and such that $h(x)=y$.

Indeed, $\mathrm{X}-\Delta$ is connected (6.1.7, Remark 3).
Corollary 6.8 Let X be a topological surface, U a connected open subset of X , $x_{1}, \ldots, x_{k}, y_{1}, \ldots, y_{k}$ distinct points of U . Then there is a homeomorphism $h$ of X restricting to the identity on $\mathrm{X}-\mathrm{U}$ and such that

$$
h\left(x_{1}\right)=y_{1}, \ldots, h\left(x_{k}\right)=y_{k}
$$

Proof For each $i$, let $h_{i}$ be a homeomorphism of X restricting to the identity on $\mathrm{X}-\mathrm{U} \cup\left\{y_{1}, \ldots, y_{i-1}, x_{i+1}, \ldots, x_{k}\right\}$ and such that $h\left(x_{i}\right)=y_{i}$. Then $h=h_{k} \circ \cdots \circ$ $h_{1}$ has the desired property.

Corollary 6.9 Let X be a nonempty connected topological surface, and $x_{1}, \ldots, x_{k}$ points of X . Then there is an open subset U of X homeomorphic to a disk and containing $x_{1}, \ldots, x_{k}$.

Proof The points $x_{i}$ may be assumed to be distinct. Let V be an open subset of X homeomorphic to a disc, not containing any $x_{i}$, and $y_{1}, \ldots, y_{k}$ be distinct points of V. If $h: \mathrm{X} \rightarrow \mathrm{X}$ is a homeomorphism such that $h_{i}\left(x_{i}\right)=y_{i}$, then the open subset $\mathrm{U}=h^{-1}(\mathrm{~V})$ has the desired property.

### 6.8.3 Coalescence of Ramifications

Proposition and Definition Let B be a topological surface, X a finite ramified covering of B and $\Delta_{\mathrm{X}}$ its ramification set. Let U be an open subset of B homeomorphic to the disk $\mathbb{D}$ and containing $\Delta_{\mathrm{x}}$, and a a point of U . Then there is a finite ramified covering $\widetilde{\mathrm{X}}$ of B , unique up to isomorphism, such that $\Delta_{\widetilde{\mathrm{X}}} \subset\{a\}$ and $\left.\widetilde{\mathrm{X}}\right|_{\mathrm{W}}=\left.\mathrm{X}\right|_{\mathrm{W}}$ for a neighbourhood W of $\mathrm{X}-\mathrm{U}$.
$\widetilde{\mathrm{X}}$ is said to be obtained from X by a coalescence of ramifications at a.
Proof Let $\varphi: \mathrm{U} \rightarrow \mathbb{D}$ be a homeomorphism, and $\mathrm{Q} \subset \mathbb{D}$ an annulus such that $\mathbb{D}-\mathrm{Q}$ is a compact set containing $\varphi\left(\Delta_{\mathrm{X}}\right)$ and $\varphi(a)$. Set $\mathrm{V}=\varphi^{-1}(\mathrm{Q})$ and $\mathrm{W}=$ $(\mathrm{B}-\mathrm{U}) \cup \mathrm{V}$. For $b \in \mathrm{Q}, \pi_{1}(\mathrm{Q}, b)=\pi_{1}(\mathbb{D}-\{\varphi(a)\}, b)=\mathbb{Z}$, and $\iota_{*}: \pi_{1}(\mathrm{Q}, b) \rightarrow$ $\pi_{1}(\mathbb{D}-\{\varphi(a)\}, b)$ is an isomorphism. Hence

$$
\iota^{*}: \mathscr{C o v}(\mathbb{D}-\{\varphi(a)\}) \rightarrow \mathscr{C o v}(\mathrm{Q})
$$

is an equivalence of categories. In particular every covering of $Q$ extends to a covering of $\mathbb{D}-\varphi(a)$. Hence, every covering of V extends to a covering of $\mathrm{U}-\{a\}$. The covering $\left.\mathrm{X}\right|_{\mathrm{V}}$ extends to a covering Y of $\mathrm{U}-\{a\}$. Gluing $\left.\mathrm{X}\right|_{\mathrm{W}}$ and Y over V gives a finite covering of $\mathrm{B}-\{a\}$, extending to a ramified covering $\widetilde{\mathrm{X}}$ of B (6.1.11), with the desired property.

If both $\widetilde{\mathrm{X}}_{1}$ and $\widetilde{\mathrm{X}}_{2}$ have this property, then there is a unique isomorphism from $\left.\widetilde{\mathrm{X}}_{1}\right|_{\mathrm{U}}$ onto $\left.\widetilde{\mathrm{X}}_{2}\right|_{\mathrm{U}}$ inducing the identity over V , hence a unique isomorphism from $\widetilde{\mathrm{X}}_{1}$ onto $\widetilde{\mathrm{X}}_{2}$ inducing the identity over W .

### 6.8.4

With the notation of 6.8.3, let $a_{1}, \ldots, a_{k}$ be the elements of $\Delta_{\mathrm{x}}$ and $b \in \mathrm{~W}$.

Proposition There are elements $\gamma_{1}, \ldots, \gamma_{k}$ in $\pi_{1}\left(\mathrm{~B}-\Delta_{\mathrm{X}}, b\right)$ such that $\gamma_{i}$ goes round $a_{i}$ (see 6.1.9, Proposition) in $\mathrm{B}-\left(\Delta_{\mathrm{X}}-\left\{a_{i}\right\}\right)$, and there exists $\gamma \in \pi_{1}(\mathrm{~B}-\{a\}, b)$ going round $a$ in B , such that the action of $\gamma$ on $\mathrm{X}(b)=\widetilde{\mathrm{X}}(b)$ defined by the covering $\left.\widetilde{\mathrm{X}}\right|_{\mathrm{B}-\{a\}}$ agrees with the action of the product $\gamma_{k} \ldots \gamma_{1}$ on $\mathrm{X}(b)$ defined by the covering $\left.X\right|_{B-\Delta_{X}}$.

Proof Keeping the notation of the former proof, we may assume that $b \in \mathrm{~V}$, and that $a_{1}, \ldots, a_{k}$ are not in $\overline{\mathrm{V}}$. If necessary multiplying $\varphi$ by a constant with absolute value 1 , we may assume that $b^{\prime}=\varphi(b) \in \mathbb{R}^{*}$, and y (6.8.2, Corollary 6.8), that $a_{\nu}^{\prime}=\varphi\left(a_{\nu}\right)$ are purely imaginary and $a_{\nu}^{\prime}=\lambda_{\nu} i$ with $\lambda_{\nu}<\lambda_{\nu+1}$.

Let $\rho>0$ be such that $\rho \leqslant \frac{1}{2}\left|a_{\nu}^{\prime}-a_{\nu+1}^{\prime}\right|$ for all $\nu$. For each $\nu$, let $\gamma_{\nu}^{\prime}$ be a loop in $\mathbb{D}-\left\{a_{1}^{\prime}, \ldots, a_{k}^{\prime}\right\}$ of the form $\beta_{\nu}^{-1} \cdot \alpha_{\nu} \cdot \beta_{\nu}$, where $\beta_{\nu}$ is a path from $b^{\prime}$ to $a^{\prime}-\rho$ such that the real part of $\beta_{\nu}(t)$ is $<0$ for all $t$, and $\alpha_{\nu}(t)=a^{\prime}-\rho \cdot e^{2 \pi i t}$. Let $\gamma^{\prime}$ be the loop $t \mapsto b^{\prime} \cdot e^{2 \pi i t}$ in Q .

Lemma The loops $\gamma_{k}^{\prime} \ldots \gamma_{1}^{\prime}$ and $\gamma^{\prime}$ are homotopic in $\mathbb{D}-\left\{a_{1}^{\prime}, \ldots, a_{k}^{\prime}\right\}$.
Proof The loop $\alpha_{\nu}$ is of the form

$$
\alpha_{\nu}^{(3)} \cdot \alpha_{\nu}^{(2)} \cdot \alpha_{\nu}^{(1)},
$$

where $\alpha_{\nu}^{(1)}$ and $\alpha_{\nu}^{(3)}$ are paths in $\mathbb{R}_{-}+i \mathbb{R}$ and $\alpha_{\nu}^{(2)}$ in $\mathbb{R}_{+}+i \mathbb{R}$.
Similarly, $\gamma^{\prime}$ is of the form $\gamma^{(3)} \cdot \gamma^{(2)} \cdot \gamma^{(1)}$. Then

$$
\gamma_{k}^{\prime} \ldots \gamma_{1}^{\prime}=\beta_{k}^{-1} \ldots \alpha_{\nu+1}^{(2)} \cdot \alpha_{\nu+1}^{(1)} \cdot \beta_{\nu+1} \cdot \beta_{\nu}^{-1} \cdot \alpha_{\nu}^{(3)} \cdot \alpha_{\nu}^{(2)} \ldots \beta_{1} .
$$

Now, $\alpha_{\nu+1}^{(1)} \cdot \beta_{\nu+1} \cdot \beta_{\nu}^{-1} \cdot \alpha_{\nu}^{(3)}$ is homotopic by barycentric subdivision to the rectilinear path $\varepsilon_{\nu}$ de $a_{\nu}^{\prime}+\rho i a_{\nu+1}^{\prime}-\rho i$, and so

$$
\gamma_{k}^{\prime} \ldots \gamma_{1}^{\prime} \simeq \beta_{k}^{-1} \cdot \alpha_{k}^{(3)} \cdot \alpha_{k}^{(2)} \cdot \varepsilon_{k-1} \cdot \alpha_{k-1}^{(2)} \ldots \varepsilon_{1} \cdot \alpha_{1}^{(2)} \cdot \alpha_{1}^{(1)} \cdot \beta_{1}
$$

In a similar way, $\beta_{k}^{-1} \cdot \alpha_{k}^{(3)}, \alpha_{k}^{(2)} \cdot \varepsilon_{k-1} \ldots \varepsilon_{1} \cdot \alpha_{1}^{(2)}$ and $\alpha_{1}^{(1)} \cdot \beta_{1}$ are respectively homotopic to $\gamma^{(3)} \cdot \varepsilon_{k}, \varepsilon_{k}^{-1} \cdot \gamma^{(2)} \cdot \varepsilon_{0}^{-1}$ and $\varepsilon_{0} \cdot \gamma^{(1)}$ where $\varepsilon_{0}$ and $\varepsilon_{k}$ are rectilinear paths. Hence $\gamma_{k}^{\prime} \ldots \gamma_{1}^{\prime} \simeq \gamma^{\prime}$.


End of the Proof of the Proposition The loops $\gamma_{\nu}=\varphi^{-1} \circ \gamma_{\nu}^{\prime}$ and $\gamma=\varphi^{-1} \circ \gamma^{\prime}$ satisfy the required assumptions. The loops $\gamma_{k} \ldots \gamma_{1}$ and $\gamma$ are homotopic in $\mathrm{B}-\Delta_{\mathrm{X}}$, and so act similarly on $\mathrm{X}(b)$. The actions of $\gamma$ on $\mathrm{X}(b)$ defined by the coverings X and $\widetilde{\mathrm{X}}$ agree since $\gamma$ is a loop in W and $\left.\widetilde{\mathrm{X}}\right|_{\mathrm{W}}=\left.\mathrm{X}\right|_{\mathrm{W}}$.

### 6.8.5 Coverings Ramified at a Single Point

Proposition Let B be a connected compact Riemann surface of genus $0, a \in \mathrm{~B}$, and X be a finite ramified covering of B such that $\Delta_{\mathrm{X}} \subset\{a\}$. Then X is a trivial covering of B . In particular $\Delta_{\mathrm{X}}=\varnothing$.
$1^{s t}$ proof (assuming known that B is homeomorphic to $\mathrm{S}^{2}$ ) The space $\mathrm{B}-\{a\}$ is homeomorphic to $\mathbb{R}^{2}$, and so is simply connected. The covering $\left.X\right|_{\mathrm{B}-\{a\}}$ is trivial, and so extends to a trivial covering of B . As the extension to a ramified covering is unique up to isomorphism (6.1.11), X is a trivial covering.
$2^{e}$ proof (without the above assumption) We may assume that $\mathrm{X} \neq \varnothing$ and that it is connected. Let $\mathrm{X}(a)=\left\{x_{1}, \ldots, x_{k}\right\}, r_{i}$ the ramification index of $x_{i}$ and $d$ the degree of X. Then $d=r_{1}+\cdots+r_{k}$. By the Riemann-Hurwitz formula (6.6.13),

$$
2-2 g(\mathrm{X})=2 d-\sum_{1}^{k}\left(r_{i}-1\right)=2 d-\sum r_{i}+k=d+k
$$

As $2-2 g(\mathrm{X}) \leqslant 2, d \geqslant 1$ and $k \geqslant 1$. Therefore $d=1$ and $k=1$.

### 6.8.6

Theorem No connected compact Riemann surface of genus 2 has an automorphism group of order 84.

Proof (By contradiction) Let X be a Riemann surface of genus 2 and G an automorphism group of $X$ of order 84 . Set $Y=X / G$. Then $X$ is a ramified Galois covering of Y , which implies (6.7.2) that Y has genus 0 and that X is ramified over 3 points $a_{1}, a_{2}$ and $a_{3}$ of Y, with ramification index 2 over $a_{1}, 3$ over $a_{2}$ and 7 over $a_{3}$.

Let $\widetilde{\mathrm{X}}$ be a ramified covering of Y obtained from X by coalescing the ramifications at a point $a \in \mathrm{Y}$. Let $b \in \mathrm{Y}-\left\{a, a_{1}, a_{2}, a_{3}\right\}, \gamma_{1}, \gamma_{2}, \gamma_{3} \in \pi_{1}\left(\mathrm{Y}-\left\{a_{1}, a_{2}, a_{3}\right\}, b\right)$ and $\gamma \in \pi_{1}(\mathrm{Y}-\{a\}, b)$ satisfying the conditions of Proposition 6.8.4. Let $\Gamma$ be the image of $\pi_{1}\left(\mathrm{Y}-\left\{a_{1}, a_{2}, a_{3}\right\}, b\right)$ in the permutation group $\mathrm{X}(b)$, and $\sigma_{i}$ the image of $\gamma_{i}$. By Proposition 6.8.5, X is a trivial covering, and so $\gamma$ acts trivially on $\mathbf{X}(b)$ and $\sigma_{3} \circ \sigma_{2} \circ \sigma_{1}=1$. By Proposition 6.1.9, $\sigma_{1}, \sigma_{2}, \sigma_{3}$ have respective orders 2, 3 and 7 . Indeed, by (4.6.11), the group $\Gamma$ is isomorphic to $G$, and so has 84 elements. The theorem then follows from the next lemma:

### 6.8.7

Lemma There is no group $\Gamma$ of order 84 having elements $u$, $v$, $w$ of respective orders 2,3 and 7 with $w=v \cdot u$.

Proof (By contradiction) Set $\mathrm{F}=\Gamma /(w)$. The set F has 12 elements and $\Gamma$ acts on F. As $v$ has order 3, the orbits of $(v)$ can have 1 or 3 elements. However, for all $x \in \mathrm{~F}$, the stabilizer of $x$ is a conjugate of $(w)$, and so is of order 7 and cannot contain $v$. Hence there are 4 orbits of $(v)$, each with 3 elements. For the same reason, there are 6 orbits of ( $u$ ), each with 2 elements.

We next give the fixed points of $w=v u$. If $x$ is one such point, then $v u \cdot x=x$, and so $u \cdot x=v^{2} x$. Hence $x=u v^{2} \cdot x$ and $v u v^{2} \cdot x=v \cdot x \neq v^{2} \cdot x$, since the orbit of $x$ has 3 distinct elements. Hence $v^{2} \cdot x$ is not a fixed point. Neither is $y=v \cdot x$, for otherwise $x=v^{2} \cdot y$ would not be fixed. Hence the orbits of $(v)$ contain at most one fixed point of $w$, which therefore fixes at most 4 fixed elements.

However the orbits of $(w)$ have 1 or 7 elements, and there can at most be one with 7 elements since $F$ only has 12 elements. So $w$ fixes at least 5 points, and thus $5 \leqslant 4$.

## Exercises 6.8. (Automorphism group, genus 2)

1.-Let $e, a_{1}, \ldots, a_{k}$ be distinct points on the Riemann sphere $\Sigma$.
(a) Show that $\pi_{1}\left(\Sigma-\left\{a_{1}, \ldots, a_{k}\right\}, e\right)$ is a free group with $k-1$ generators. More precisely, there are $\gamma_{1}, \ldots, \gamma_{k} \in \pi_{1}\left(\Sigma-\left\{a_{1}, \ldots, a_{k}\right\}, e\right)$ such that $\left(\gamma_{1}, \ldots, \gamma_{k-1}\right)$ is a basis for this group, and $\gamma_{i}$ go round $a_{i}$ in $\left(\Sigma-\left\{a_{1}, \ldots, a_{k}\right\}\right) \cup\left\{a_{i}\right\}$ and $\gamma_{k}$. $\gamma_{k-1} \cdots \gamma_{1}=e$.
(b) Let G be a finite group generated by $\alpha_{1}, \ldots, \alpha_{k-1}$ and set $\alpha_{k}=\alpha_{1}^{-1} \ldots \alpha_{k-1}^{-1}$. Show that there is a ramified Galois covering $X$ of $\Sigma$ with group G, whose ramification set is in $\left\{a_{1}, \ldots, a_{k}\right\}$ with ramification over $a_{i}$ having index the order $s_{i}$ of $\alpha_{i}$ in G.
(c) Give the genus of X in terms of the order $n$ of G and of the indices $s_{i}$.

### 6.9 Poincaré Geometry

The universal covering of a Riemann surface $X$ of genus $\geqslant 2$ can be shown to be isomorphic to a disk, so that X becomes a quotient of the disk by an automorphism group. In this section, we study the automorphisms of the disk. There is a metric on $\mathbb{D}$ with respect to which every automorphism of $\mathbb{D}$ is an isometry (Poincaré metric).

### 6.9.1 Homographic Transformations of Circles

A circle on the Riemann sphere $\Sigma=\mathbb{C} \cup\{\infty\}$ is a set which is a circle in $\mathbb{C}$, and so of the form $L \cup\{\infty\}$, where $L$ is a line in $\mathbb{C}$. The complement of a circle on $\Sigma$ has
two connected components, called the disks limited by this circle. Hence a disk on $\Sigma$ is either an open disk in $\mathbb{C}$, or an open half-plane in $\mathbb{C}$, or the complement in $\Sigma$ of a closed disk in $\mathbb{C}$.

Proposition The homographic image of a circle (resp. a disk) on $\Sigma$ is a circle (resp. a disk).

Proof The homography group is generated by the maps $x \mapsto a x+b$ and the map $h: x \mapsto \frac{1}{x}$. The proposition obviously holds for the former ones. We show that $h$ sends a circle C on $\Sigma$ to a circle. First suppose that C is a circle in $\mathbb{C}$ which not passing through 0 , and that $a$ and $b$ are the endpoints of the diameter of C through 0 (an arbitrary diameter if C is centered at 0 ). Then $x \in \mathrm{C}$ if and only if $\frac{x-a}{x-b} \in i \mathbb{R} \cup\{\infty\}$. As $\frac{1 / x-1 / a}{1 / x-1 / b}=\frac{b}{a} \frac{x-a}{x-b}$ and $\frac{b}{a} \in \mathbb{R}^{*}, x \in \mathrm{C}$ if and only if $\frac{1}{x}$ is in the circle of diameter $\left[\frac{1}{a}, \frac{1}{b}\right]$.

If $\mathbb{C}$ is a circle in $\mathbb{C}$ passing through 0 , then let $a$ be such that $(0, a)$ is a diameter. So $x \in \mathrm{C}$ if and only if $\frac{x-a}{x} \in i \mathbb{R} \cup\{\infty\}$, i.e. if $\frac{1}{x}$ is in the inverse image of $i \mathbb{R} \cup\{\infty\}$ under $y \mapsto 1-a y$. This method gives all the circles $L \cup\{\infty\}$, where $L$ is a line in $\mathbb{C}$ not passing through 0 . As $h$ is involutive, the image of such a circle must be a circle of $\mathbb{C}$ passing through 0 .

If $\mathrm{C}=\mathrm{L} \cup\{\infty\}$, where $\mathrm{L}=\lambda \mathbb{R}$ is a line passing through 0 , then the image of C is $\frac{1}{\lambda} \mathbb{R} \cup\{\infty\}$.

Finally, as $h$ is a homeomorphism of $\Sigma$, the images of the connected components of $\Sigma-\mathrm{C}$ are the connected components of $\Sigma-h(\mathrm{C})$, and so are disks.

### 6.9.2

Let $\mathbb{D}$ be the unit disk of $\mathbb{C}$ and $S^{1}$ the unit circle of $\mathbb{C}$, i.e. $\mathbb{D}=\{z| | z \mid<1\}$ and $S^{1}=\{z| | z \mid=1\}$.

Proposition Let $h: \Sigma \rightarrow \Sigma$ be a homography. Then, $h\left(\mathrm{~S}^{1}\right)=\mathrm{S}^{1}$ if and only if $h$ is of the form

$$
x \mapsto \lambda \frac{x-a}{1-\bar{a} x} \text { with } a \in \mathbb{C}-\mathrm{S}^{1} \text { and } \lambda \in \mathrm{S}^{1}, \text { or } x \mapsto \frac{\lambda}{x} \text { with } \lambda \in \mathrm{S}^{1}
$$

Proof Suppose that $h$ is a homography. Let $b=h^{-1}(\infty)$. Set $h_{0}(x)=\frac{x-a}{1-\bar{a} x}$ with $a=\frac{1}{b}$ if $b \neq 0$ and $h_{0}(x)=\frac{1}{x}$ if $b=0$. Now, $h_{0}(b)=\infty$ and $h_{0}\left(S^{1}\right)=S^{1}$. Set $u=h \circ h_{0}^{-1}$.

The map $u$ is a homography such that $u(\infty)=\infty$, and hence of the form $z \mapsto \lambda z+\beta$, and $u\left(\mathrm{~S}^{1}\right)=\mathrm{S}^{1}$, which implies that $|\lambda|=1$ and $\beta=0$; so $h=\lambda h_{0}$.

Conversely, for $x \in \mathrm{~S}^{1}, 1-\bar{a} x=x(\bar{x}-\bar{a})$. Thus $|1-\bar{a} x|=|x-a|$ and $|h(x)|=1$. The case $h=\left(x \mapsto \frac{\lambda}{x}\right)$ is immediate.

Corollary 6.10 Let $h$ be a homography. Then $h(\mathbb{D})=\mathbb{D}$ if and only if $h$ is of the form

$$
h_{a, \lambda}: x \mapsto \lambda \frac{x-a}{1-\bar{a} x} \text { with } a \in \mathbb{D} \text { and } \lambda \in \mathrm{S}^{1}
$$

Proof If $h\left(\mathrm{~S}^{1}\right)=\mathrm{S}^{1}$, the image of $\mathbb{D}$ is either $\mathbb{D}$ or $\Sigma-\overline{\mathbb{D}}$. If $h=h_{a, \lambda}$, then $h(0)=$ $-\lambda a$ and $h(\mathbb{D})=\mathbb{D}$ if $a \in \mathbb{D}$.

Conversely, if $h(\mathbb{D})=\mathbb{D}$, then $h\left(\mathrm{~S}^{1}\right)=\mathrm{S}^{1}$ since $\mathrm{S}^{1}$ is the boundary of $\mathbb{D}$; $h$ cannot be of the form $\frac{\lambda}{x}$, and so is of the form $h_{a, \lambda}$, and $a=h^{-1}(0) \in \mathbb{D}$.

Corollary 6.11 The homographies $h$ such that $h(\mathbb{D})=\mathbb{D}$ form a group acting transitively on $\mathbb{D}$.

Indeed the orbit of 0 is $\mathbb{D}$ since $h_{a, \lambda}^{-1}(0)=a$ for $a \in \mathbb{D}$.

### 6.9.3 Automorphisms of $\mathbb{D}$

An automorphism of $\mathbb{D}$ is an automorphism of D as a Riemann surface, i.e. a biholomorphic bijection (it follows from 6.1.6 that any holomorphic bijection is biholomorphic).

Theorem Every automorphism of $\mathbb{D}$ is induced by a homography of $\Sigma$ preserving $\mathbb{D}$.

Proof Let $f$ be an automorphism of $\mathbb{D}$. We first suppose that $f(0)=0$. Then there is a holomorphic $u$ with $|u(z)|$ tending to 1 as $|z|$ tends to 1 such that $f(z)=z \cdot u(z)$. By the maximum principle, $|u(0)| \leqslant 1$, and so $\left|f^{\prime}(0)\right| \leqslant 1$ since $f^{\prime}(0)=u(0)$. Set $g=f^{-1}$. Similarly, $\left|g^{\prime}(0)\right| \leqslant 1$. As $g^{\prime}(0)=\frac{1}{f^{\prime}(0)},\left|f^{\prime}(0)\right|=1$ and $|u(0)|=1$. Hence $u$ is a constant $\lambda$ with absolute value 1 , and $f(z)=\lambda z$.

In the general case, set $a=f^{-1}(0)$ and let $f_{0}$ be an automorphism of $\mathbb{D}$ induced by a homography and such that $f(a)=0$ (see 6.9.2, Corollary 6.11). Set $g=f \circ f_{0}^{-1}$. Then $g$ is an automorphism of $\mathbb{D}$ such that $g(0)=0$, and so is $g: z \mapsto \lambda z$ with $\lambda \in \mathrm{S}^{1}$, and $f=\lambda f_{0}$.

### 6.9.4 Riemannian Metrics

Let U be an open subset of $\mathbb{C}$. A Riemannian metric ${ }^{5}$ on U is given by a continuous function $\mu: \mathrm{U} \rightarrow \mathbb{R}_{+}^{*}$. If $\gamma:[a, b] \rightarrow \mathrm{U}$ is a continuous differentiable path, then the length of $\gamma$ with respect to the Riemannian metric defined by $\mu$ is $\mathrm{L}_{\mu}(\gamma)=\int_{a}^{b}\left|\gamma^{\prime}(t)\right| \mu(\gamma(t)) d t$. For $x, y \in \mathrm{U}$, the distance from $x$ to $y$ with respect to

[^23]the Riemannian metric defined by $\mu$ is the lower bound $d_{\mu}(x, y)$ of the lengths of paths continuously differentiable from $x$ to $y$ in U .

Let $\mathrm{U}_{1}$ and $\mathrm{U}_{2}$ be open subsets of $\mathbb{C}, \mu_{1}$ and $\mu_{2}$ functions defining Riemannian metrics on $\mathrm{U}_{1}$ and $\mathrm{U}_{2}$ respectively. Let $f: \mathrm{U}_{1} \rightarrow \mathrm{U}_{2}$ be a biholomorphic bijection. Consider the following conditions:
(i) $\left(\forall x \in \mathrm{U}_{1}\right) \quad \mu_{1}(x)=\left|f^{\prime}(x)\right| \mu_{2}(f(x))$;
(ii) For any continuously differentiable path $\gamma$ in $\mathrm{U}_{1}$,

$$
\mathrm{L}_{\mu_{1}}(\gamma)=\mathrm{L}_{\mu_{2}}(f \circ \gamma) ;
$$

(iii) for $x, y \in \mathrm{U}_{1}, d_{\mu_{1}}(x, y)=d_{\mu_{2}}(f(x), f(y))$.
(i) $\Rightarrow$ (ii) $\Rightarrow$ (iii) are immediate. ${ }^{6}$ These conditions can be shown to be equivalent (6.9, Exercise 3). Given the function $\mu_{2}$, there is a unique function $\mu_{1}$ satisfying (i); the Riemannian metric defined by $\mu_{1}$ is the inverse image under $f$ of the Riemannian metric defined by $\mu_{2}$ and we write $\mu_{1}=f^{*} \mu_{2}$.

### 6.9.5 Poincaré Metric

Theorem and Definition There is a unique Riemannian metric on $\mathbb{D}$ defined by a function $\mu$ such that $\mu(0)=1$, and such that, for any automorphism $f$ of $\mathbb{D}, f^{*} \mu=\mu$.

This Riemannian metric is called the Poincaré metric on $\mathbb{D}$. It is given by $\mu(x)=$ $\frac{1}{1-x \bar{x}}$.

Proof (1) Uniqueness. Let $\mu$ be another Riemannian metric with the required property and $x \in \mathbb{D}$. As $h_{x, \lambda}$ is an automorphism $\mathbb{D}$ sending $x$ to 0 ,

$$
\mu(x)=\left|h_{x, \lambda}^{\prime}(x)\right| \mu(0)=\left|h_{x, \lambda}^{\prime}(x)\right|=\frac{1}{|1-x \bar{x}|}=\frac{1}{1-x \bar{x}} .
$$

(2) Existence. The function $\mu: x \mapsto \frac{1}{1-x \bar{x}}$ is a strictly positive continuous function on $\mathbb{D}$. Consider the Riemannian metric defined by $\mu$. Let $f$ be an automorphism of $\mathbb{D}$ and $x \in \mathbb{D}$. There exists $\lambda \in \mathbb{C}$ with absolute value 1 such that $h_{y, 1} \circ f=h_{x, \lambda}$, where $y=f(x)$. Hence,

$$
\left|h_{y, 1}^{\prime}(y)\right| \cdot\left|f^{\prime}(x)\right|=\left|h_{x, \lambda}^{\prime}(x)\right|
$$

and so $\left|f^{\prime}(x)\right|=\frac{\left|h_{x, \lambda}^{\prime}(x)\right|}{\left|h_{y, 1}^{\prime}(y)\right|}=\frac{\mu(x)}{\mu(y)}$ and $f^{*} \mu=\mu$.

[^24]
### 6.9.6

Proposition Let $x$ and $y$ be two points of $\mathbb{R} \cap \mathbb{D}=]-1,1\left[\right.$ with $x \leqslant y$, and $\gamma_{0}$ the path $t \mapsto t$ from $[x, y]$ to $\mathbb{D}$. For the Poincaré metric $\mu$ on $\mathbb{D}, d_{\mu}(x, y)=\mathrm{L}_{\mu}\left(\gamma_{0}\right)=$ Argth $y$ - Argth $x$.

Proof Now, $\mathrm{L}_{\mu}\left(\gamma_{0}\right)=\int_{x}^{y} \frac{1}{1-t^{2}} d t=$ Argth $y-\operatorname{Argth} x$.
Let $\gamma$ be a continuously differentiable path from $x$ to $y$. Write

$$
\gamma(t)=\gamma_{1}(t)+i \gamma_{2}(t) \text { with } \gamma_{1}(t) \text { and } \gamma_{2}(t) \text { reals. }
$$

As $\gamma_{1}^{\prime}(t)=\operatorname{Re} \gamma^{\prime}(t),\left|\gamma_{1}^{\prime}(t)\right| \leqslant\left|\gamma^{\prime}(t)\right|$, and $\left|\gamma_{1}(t)\right| \leqslant|\gamma(t)|$; so

$$
\frac{1}{1-\left(\gamma_{1}(t)\right)^{2}} \leqslant \frac{1}{1-|\gamma(t)|^{2}}
$$

and

$$
\begin{aligned}
\mathrm{L}_{\mu}(\gamma)=\int_{x}^{y} \frac{\left|\gamma^{\prime}(t)\right|}{1-|\gamma(t)|^{2}} & d t \geqslant \mathrm{~L}_{\mu}\left(\gamma_{1}\right)=\int_{x}^{y} \frac{\left|\gamma_{1}^{\prime}(t)\right|}{1-\left(\gamma_{1}(t)\right)^{2}} d t \\
& \geqslant \int_{x}^{y} \frac{\gamma_{1}^{\prime}(t)}{1-\left(\gamma_{1}(t)\right)^{2}} d t=\int_{x}^{y} \frac{d u}{1-u^{2}}=\mathrm{L}_{\mu}\left(\gamma_{0}\right)
\end{aligned}
$$

The distance $d_{\mu}$ is called the Poincaré distance.
Remark The above proof shows that if $\gamma$ is a continuously differentiable path from $x$ to $y$ whose image is not contained in $\mathbb{R}$, then $\mathrm{L}_{\mu}(\gamma)>d_{\mu}(x, y)$.

Corollary 6.12 Let $x \in \mathbb{D}$. For the Poincaré metric $\mu$,

$$
d_{\mu}(0, x)=\operatorname{Argth}|x|
$$

Indeed, $x=\lambda|x|$ with $\lambda \in \mathrm{S}^{1}$. As $z \mapsto \lambda z$ is an isometry with respect to the Poincaré distance, $d_{\mu}(0, x)=d_{\mu}(0,|x|)=$ Argth $|x|$.

Corollary 6.13 The topology defined on $\mathbb{D}$ by the Poincaré distance is the Euclidean topology.

Indeed, the neighbourhoods of 0 are the same. For all $x \in \mathbb{D}$, the map $h_{x, 1}$ is a homeomorphism with respect to the Euclidean topology and an isometry with respect to the Poincaré distance, transforming $x$ into 0 . Hence the neighbourhoods of $x$ are the same.

Corollary 6.14 Every closed subset of $\mathbb{D}$ bounded with respect to the Poincaré distance is compact.

### 6.9.7 Poincaré Circles

The Poincaré circle $\mathrm{C}_{a, r}$ with centre $a \in \mathbb{D}$ and radius $r$ is the set of $x \in \mathbb{D}$ such that $d_{\mu}(a, x)=r$.

Proposition Every Poincaré circle is an Euclidean circle contained in $\mathbb{D}$.
Proof By 6.9.6, Corollary 6.12, the Poincaré circle $\mathrm{C}_{0, r}$ is the Euclidean circle with centre 0 and radius th $r$. With the notation of 6.9.2, Corollary 6.10, $\mathrm{C}_{a, r}=h_{a, 1}^{-1}\left(\mathrm{C}_{0, r}\right)$, and so, by Proposition 6.9.1, $\mathrm{C}_{a, r}$ is a circle in $\mathbb{D}$.

Remarks (1) If $a \neq 0$, then it is not the Euclidean centre of the circle $\mathrm{C}_{a, r}$.
(2) The Poincaré circle $\mathrm{C}_{a, r}$ is in the Euclidean annulus consisting of $u$ such that

$$
\left(1-|a|^{2}\right) \frac{\rho}{1+\rho}<|u-a|<\left(1-|a|^{2}\right) \frac{\rho}{1-\rho}
$$

where $\rho=\operatorname{th} r$. Indeed, if $u \in \mathrm{C}_{a, r}$, then $u=h_{a, 1}^{-1}(z)=\frac{a+z}{1+\bar{a} z}$, with $z \in \mathrm{C}_{0, r}$. Hence $|z|=\rho$, and $u-a=\frac{(1-a \bar{a}) z}{1+\bar{a} z}$; so $|u-a|=\frac{1-|a|^{2}}{|1+\bar{a} z|}|z|$, with $1-\rho<|1+\bar{a} z|<1+\rho$. Therefore, the Euclidean radius $r_{e}$ of $\mathrm{C}_{a, r}$ satisfies

$$
\left(1-|a|^{2}\right) \frac{\rho}{1+\rho}<r_{e}<\left(1-|a|^{2}\right) \frac{\rho}{1-\rho} .
$$

In particular the Euclidean diameter of $\mathrm{C}_{a, r}$ tends to 0 as $|a|$ tends to 1 , the radius $r$ staying bounded (see M.C. Escher's engraving "Circle Limit III", opposite 6.10.18, where all the fish have the same Poincaré diameter).

### 6.9.8 Geodesic Paths

Let U be an open subset of $\mathbb{C}$ and $\mu$ a Riemannian metric on U . A continuously differentiable map $\gamma:[a, b] \rightarrow \mathrm{U}$ is a geodesic path if for all intervals $\left[a^{\prime}, b^{\prime}\right] \subset$ $[a, b]$,

$$
d_{\mu}\left(\gamma\left(a^{\prime}\right), \gamma\left(b^{\prime}\right)\right)=\mathrm{L}_{\mu}\left(\left.\gamma\right|_{\left[a^{\prime}, b^{\prime}\right]}\right)=\left|b^{\prime}-a^{\prime}\right| .
$$

Lemma For any $x, y \in \mathbb{D}$, there is a unique geodesic path $\gamma:\left[0, d_{\mu}(x, y)\right] \rightarrow \mathbb{D}$ from $x$ to $y$ with respect to the Poincaré metric. Its image is an arc of circle of $\Sigma$ orthogonal to $S^{1}$ passing through $x$ and $y$.

Proof Suppose first that $x=0$ and $y \in \mathbb{R}_{+}$. Then the map $s \mapsto$ th $s$ from [ 0 , Argth $t$ ] to $\mathbb{D}$ is a geodesic path from 0 to $y$ and by (6.9.6, Remarque) it is unique. Its image is an interval of $\mathbb{R}$.

In the general case, there is an automorphism $h$ of $\mathbb{D}$ such that $h(x)=0$ and $h(y) \in \mathbb{R}_{+}$. If $\gamma_{0}$ is the geodesic path from 0 to $h(y)$, then $\gamma=h^{-1} \circ \gamma_{0}$ is a geodesic
path from $x$ to $y$ and it is the only one. The automorphism $h$ extends to a homography, which we also write $h: \Sigma \rightarrow \Sigma$, and the image of $\gamma$ is an arc of circle $h^{-1}(\mathbb{R} \cup\{\infty\})$. This circle is orthogonal to $S^{1}$ since $\mathbb{R}$ is orthogonal to $S^{1}$ and $h$ preserves angles, its tangent linear map being $\mathbb{C}$-linear, and so is a similitude.

### 6.9.9 Geodesics

Definition A geodesic in $\mathbb{D}$ is any set of the form $\gamma(\mathbb{R})$, where $\gamma: \mathbb{R} \rightarrow \mathbb{D}$ is a map such that, for all $[a, b] \subset \mathbb{R},\left.\gamma\right|_{[a, b]}$ is a geodesic path with respect to the Poincaré metric.

Proposition The geodesics of $\mathbb{D}$ are the sets of the form $\Gamma \cap \mathbb{D}$, where $\Gamma$ is a circle orthogonal to $\mathrm{S}^{1}$.

Proof (a) Any geodesic is of the form $\Gamma \cap \mathbb{D}$ : let $\gamma: \mathbb{R} \rightarrow \mathbb{D}$ be as in the above definition. For each $[a, b], \gamma([a, b])$ is in a circle $\Gamma_{a, b}$. For $\left[a^{\prime}, b^{\prime}\right] \subset[a, b], \Gamma_{a, b} \cap$ $\Gamma_{a^{\prime}, b^{\prime}} \supset \gamma\left[a^{\prime}, b^{\prime}\right]$ containing more than 3 points, and so $\Gamma_{a, b}=\Gamma_{a^{\prime}, b^{\prime}}$. Hence $\Gamma_{a, b}$ does not depend on the choice of $[a, b]$. We write it $\Gamma$. Then $\gamma(\mathbb{R}) \subset \Gamma$ and $\gamma$ is an isometry from Euclidean $\mathbb{R}$ onto a subset of $\Gamma \cap \mathbb{D}$ equipped with the Poincaré distance. This subset is homeomorphic to $\mathbb{R}$, and so is an open arc. It is complete and hence closed in $\Gamma \cap \mathbb{D}$. As $\Gamma \cap \mathbb{D}$ is connected and $\gamma(\mathbb{R}) \neq \varnothing, \gamma(\mathbb{R})=\Gamma \cap \mathbb{D}$.
(b) Any set $\Gamma \cap \mathbb{D}$ is a geodesic: indeed, $\mathbb{R} \cap \mathbb{D}$ is a geodesic since it is the image of $t \mapsto$ th $t$. Let $\Gamma$ be a circle orthogonal to $\mathrm{S}^{1}$; choose $x \in \Gamma \cap \mathbb{D}$. Let $h$ be a homography preserving $\mathbb{D}$ and such that $h(x)=0$. Then $h(\Gamma)$ is a circle on $\Sigma$ passing through 0 and orthogonal to $S^{1}$. Hence $h(\Gamma)$ must be of the form $L \cup\{\infty\}$, where $L$ is a line passing through 0 , and we may assume that $\mathrm{L}=\mathbb{R}$. Then $\Gamma \cap \mathbb{D}=h^{-1}(\mathbb{R} \cap \mathbb{D})$ is a geodesic.

## 6.9 .10

Corollary 6.15 The geodesics passing through 0 are the sets of the form $\mathrm{L} \cap \mathbb{D}$, where L is a line passing through 0 .

Corollary 6.16 There is a a unique geodesic through two distinct points of $\mathbb{D}$.
Corollary 6.17 There is a unique geodesic passing through a tangent point to a given differentiable curve. In particular, two distinct geodesics are never tangent.

Corollary 6.18 The automorphism group of $\mathbb{D}$ acts transitively on the set of geodesics of $\mathbb{D}$.

### 6.9.11 Isometries of $\mathbb{D}$

Let $\tau$ be the bijection $x \mapsto \bar{x}$ from $\mathbb{D}$ onto $\mathbb{D}$. For any continuously differentiable path $\gamma$ in $\mathbb{D}, \mathrm{L}_{\mu}(\tau \circ \gamma)=\mathrm{L}_{\mu}(\gamma)$. Hence $\tau$ is an isometry with respect to the Poincaré distance.

Proposition Any isometry $f: \mathbb{D} \rightarrow \mathbb{D}$ with respect to the Poincaré distance is either an automorphism of $\mathbb{D}$ or of the form $\tau \circ h$, where $h$ is an automorphism of $\mathbb{D}$.

Lemma 6.7 Let $a, b \in \mathbb{D}, a \neq b$, and $r, r^{\prime} \in \mathbb{Z}, r, r^{\prime}>0$. Then $\mathrm{C}_{a, r} \neq \mathrm{C}_{b, r^{\prime}}$.
Proof Let $\Delta=\gamma(\mathbb{R})$ be the geodesic passing through $a$ and $b$, with $\gamma$ satisfying the conditions of definition 6.9.9; define $\alpha$ and $\beta$ by $\gamma(\alpha)=a$ and $\gamma(\beta)=b$. The circle $\mathrm{C}_{a, r}$ meets $\Delta$ at two points $\gamma\left(\alpha_{1}\right)$ and $\gamma\left(\alpha_{2}\right)$ with

$$
\alpha_{1}=\alpha-r \quad \text { and } \quad \alpha_{2}=\alpha+r ;
$$

so $\alpha=\frac{1}{2}\left(\alpha_{1}+\alpha_{2}\right)$. Similarly,

$$
\Delta \cap \mathrm{C}_{b, r^{\prime}}=\left\{\gamma\left(\beta_{1}\right), \gamma\left(\beta_{2}\right)\right\} \text { with } \beta=\frac{1}{2}\left(\beta_{1}+\beta_{2}\right)
$$

If $\mathrm{C}_{a, r}=\mathrm{C}_{b, r^{\prime}}$, then $\alpha=\beta$, and thus $a=b$.
Lemma 6.8 Let $f: \mathbb{D} \rightarrow \mathbb{D}$ be an isometry fixing two distinct points $a$ and $b$ of $\mathbb{R} \cap \mathbb{D}$. Then $f=1_{\mathbb{D}}$ or $f=\tau$.

Let $x \in \mathbb{D}-\mathbb{R}$ and $r=d_{\mu}(a, x), r^{\prime}=d_{\mu}(b, x)$. The point $f(x)$ is in $\mathrm{C}_{a, r} \cap \mathrm{C}_{b, r^{\prime}}$. Now, by Lemma 6.7, $\mathrm{C}_{a, r}$ and $\mathrm{C}_{b, r^{\prime}}$ are distinct Euclidean circles with intersection consisting of $x$ and $\tau(x)$. Hence $f(x)=x$ or $f(x)=\bar{x}$. We may write $f(x)=$ $\operatorname{Re}(x)+\varepsilon(x) \operatorname{Im}(x) \cdot i$, where $\varepsilon(x)= \pm 1$. As $f$ is continuous, the map $\varepsilon: \mathbb{D}-\mathbb{R} \rightarrow$ $\{+1,-1\}$ is continuous, hence constant on each half-disc. This constant is the same on both, for otherwise $f(x)=f(\bar{x})$ and $f$ would not be injective. Therefore, $f$ agrees with either $1_{\mathbb{D}}$ or $\tau$ on $\mathbb{D}-\mathbb{R}$, hence on $\mathbb{D}$ by continuity.
Proof of the Proposition Choose a point $c \in \mathbb{D} \cap \mathbb{R}_{+}^{*}$, set $a=f^{-1}(0)$ and $b=$ $f^{-1}(c)$. There is an automorphism $h$ de $\mathbb{D}$ such that $h(a)=0$ and $h(b) \in \mathbb{R}_{+}$. As $d_{\mu}(0, h(b))=d_{\mu}(a, b)=d_{\mu}(0, c), h(b)=c$. Set $g=h \circ f^{-1}$. Then $g(0)=0$, $g(c)=c$ and $g$ is an isometry. By Lemma $6.8, g=1_{\mathbb{D}}$ or $g=\tau$, and so $f=h$ or $f=\tau \circ h$.

### 6.9.12

Proposition and Definition Let $\Gamma$ be a geodesic of $\mathbb{D}$. There is a unique isometry $\sigma_{\Gamma}$ such that $\Gamma$ is the set of points fixed by $\sigma_{\Gamma}$. The latter, called the geodesic symmetry, is antiholomorphic.

Proof The case $\Gamma=\mathbb{R} \cap \mathbb{D}$ follows from 6.9.11, Lemma 6.8: then $\sigma_{\Gamma}=\tau$. The general case follows from 6.9.10, Corollary 6.16.


Remarks (1) The symmetry $\sigma_{\Gamma}$ is induced by inversion with respect to the circle containing $\Gamma$.
(2) The map $\sigma_{\Gamma}$ can be described as follows: let $x \in \mathbb{D}$; there is a unique geodesic H orthogonal to $\Gamma$ and passing through $x$. It crosses $\Gamma$ at a point $h$. if $x \neq h, y=\sigma_{\Gamma}(x)$ is the point of H distinct from $x$ such that $d_{\mu}(h, y)=d_{\mu}(h, x)$.

## 6.9 .13

Proposition and Definition Let $a \in \mathbb{D}, \theta \in \mathbb{R} / 2 \pi \mathbb{Z}$; there is a unique isometry $r_{a, \theta}$ of $\mathbb{D}$ fixing $a$, whose tangent map at a is the rotation $t \mapsto e^{i \theta} t$. It is called the Poincaré rotation around the centre $a$ with angle $\theta$ and is holomorphic. Any automorphism of $\mathbb{D}$ fixing $a$ is of the form $r_{a, \theta}$.

Proof The proposition is immediate if $a=0$; the isometry $r_{a, \theta}$ is then the Euclidean rotation around O of angle $\theta$.

The general case follows by transference using an automorphism $h$ of $\mathbb{D}$ such that $h(a)=0$, noting that

$$
\left(h^{\prime}(a)\right)^{-1} \cdot 1 \circ e^{i \theta} \cdot 1 \circ h^{\prime}(a) \cdot 1=e^{i \theta} \cdot 1,
$$

Corollary 6.19 Let $\Gamma$ and $\Gamma^{\prime}$ be geodesics of $\mathbb{D}$ intersecting at a point a with angle $\theta$. Then $\sigma_{\Gamma^{\prime}} \circ \sigma_{\Gamma}=r_{a, 2 \theta}$.

Corollary 6.20 Let $a \in \mathbb{D}, \theta \in \mathbb{R} / 2 \pi \mathbb{Z}$ and $\Gamma$ a geodesic through a. Then $r_{a, \theta} \circ \sigma_{\Gamma}$, (resp. $\sigma_{\Gamma} \circ r_{a, \theta}$ ) is the geodesic symmetry

$$
r_{a, \frac{\theta}{2}}(\Gamma) \quad\left(\operatorname{resp} \cdot r_{a,-\frac{\theta}{2}}(\Gamma)\right) .
$$

### 6.9.14 Geodesically Convex Sets

Let $a$ and $b$ be points of $\mathbb{D}$ and $\gamma$ the geodesic path from $a$ to $b$.
The image de $\gamma$ is called the geodesic segment with endpoints $a$ and $b$, and is written $[a, b]$. Let $\alpha$ and $\beta$ be such that $\gamma(\alpha)=a$ and $\gamma(\beta)=b$. For $u, v \in \mathbb{R}_{+}$such that $u+v=1$, we call the point $\gamma(u \alpha+v \beta)$ the geodesic barycentre of $a$ and $b$ with coefficients $u$ and $v$, and write it $\operatorname{Bar}(a, u ; b, v)$.

The map $(a, b, u) \mapsto \operatorname{Bar}(a, u ; b, 1-u)$ from $\mathbb{D} \times \mathbb{D} \times[0,1]$ to $\mathbb{D}$ is continuous. Contrary to what happens in Euclidean geometry, if $u, v$ and $w$ are three numbers $>0$ such that $u+v+w=1$, then the points

$$
\operatorname{Bar}(a, u ;(b, v ; c, w))=\operatorname{Bar}\left(a, u ; \operatorname{Bar}\left(b, \frac{v}{v+w} ; c, \frac{w}{v+w}\right), v+w\right)
$$

and

$$
\operatorname{Bar}((a, u ; b, v) ; c, w)=\operatorname{Bar}\left(\operatorname{Bar}\left(a, \frac{u}{u+v} ; b, \frac{v}{u+v}\right), u+v ; c, w\right)
$$

are usually distinct.
A subset A of $\mathbb{D}$ is said to be geodesically convex if for all $a, b \in \mathrm{~A}$, the geodesic segment $[a, b]$ is in A. All geodesically convex sets are connected.

Let $\Gamma$ be a geodesic of $\mathbb{D}$. The complement of $\Gamma$ has two connected components $\mathrm{D}_{1}$ and $\mathrm{D}_{2}$. The sets $\mathrm{H}_{1}=\mathrm{D}_{1} \cup \Gamma$ and $\mathrm{H}_{2}=\mathrm{D}_{2} \cup \Gamma$ are called the half-disks limited by $\Gamma$. They are closed and their boundary is $\Gamma$. They are geodesically convex: indeed, let $a, b \in \mathrm{H}_{1}$ and $\gamma:[u, v] \rightarrow \mathbb{D}$ be the geodesic path from $a$ to $b$. Suppose that for some $w \in[u, v], \gamma(w) \in \mathrm{D}_{2}$. Then there exist $u_{1} \in\left[u, w\left[\right.\right.$ and $\left.\left.v_{1} \in\right] w, v\right]$ such that $\gamma\left(u_{1}\right) \in \Gamma$, and $\gamma\left(v_{1}\right) \in \Gamma$, which is impossible since the geodesic through $a$ and $b$ can only meet $\Gamma$ at one point if it is distinct from $\Gamma$.

### 6.9.15 Geodesic Triangles

Let $a, b, c \in \mathbb{D}$ be points not on the same geodesic. Let $\Gamma_{b, c}$ denote the geodesic through $b$ and $c$ and $\mathrm{H}_{a}$ the half-disk limited by $\Gamma_{b, c}$ containing $a$. Define $\mathrm{H}_{b}$ and $\mathrm{H}_{c}$ likewise. The geodesic triangle with vertices $a, b, c$ is $\mathrm{T}=\mathrm{H}_{a} \cap \mathrm{H}_{b} \cap \mathrm{H}_{c}$. The set T is geodesically convex, and hence is connected. As it is closed, the formula $\partial(A \cap B)=(\partial A \cap B) \cup(A \cap \partial B)$ which holds for any two arbitrary closed sets implies that its boundary is

$$
\partial \mathrm{T}=[a, b] \cup[b, c] \cup] c, a] .
$$

The set T is bounded with respect to the Poincaré distance. Indeed, for all $x \in \mathrm{~T}$, the Poincaré circle through $x$ and with centre $a$ contains points that are not in T since it contains points of each half-disk limited by $\Gamma_{a, b}$; hence it meets $\partial \mathrm{T}$, and $d_{\mu}(a, x)$ is bounded above by

$$
\sup _{y \in \partial \mathrm{~T}} d_{\mu}(a, y)
$$

a number independent of $x$, which is finite since $\partial \mathrm{T}$ is compact. So the set T is compact (6.9.6, Corollary 6.14).

### 6.9.16 Sum of the Angles of a Geodesic Triangle

Let $a, b, c \in \mathbb{D}, u$ (resp. $v$ ) the derivative at the origin of the geodesic path from $a$ to $b$ (resp. from $a$ to $c$ ). There is a unique $\alpha \in]-\pi,+\pi\left[\right.$ such that $v=e^{i \alpha} \cdot u$. We say that $|\alpha|$ is the (unorieunorientednted) angle at $a$ of the geodesic triangle T with vertices $a, b, c$.

Theorem (a) The sum of the three angles of a geodesic triangle is strictly less than $\pi$.
(b) Given $\alpha$, $\beta$, $\gamma$ in $] 0, \pi[$ such that $\alpha+\beta+\gamma<\pi$, there is a geodesic triangle in $\mathbb{D}$ with angles $\alpha, \beta$ and $\gamma$.

Proof (a) If $f$ is an automorphism of $\mathbb{D}$, the angle of T at $a$ equals the angle of $f(\mathrm{~T})$ at $f(a)$. Hence we may assume that $a=0$. Then the geodesics $\Gamma_{a, b}$ and $\Gamma_{a, c}$ are Euclidean line segments, and $a$ is in the region outside the circle containing $\Gamma_{b, c}$. Let $\omega$ be the Euclidean centre of this circle, $\alpha^{\prime}$ and $\beta^{\prime}$ the angles at $b$ and $c$ of the Euclidean triangle with vertices $a, b, c$.

Then $\alpha=\alpha^{\prime}-\frac{1}{2} \delta$ and $\beta=\beta^{\prime}-\frac{1}{2} \delta$, where $\delta$ in the Euclidean angle $b \omega c$. Hence $\alpha+\beta+\gamma<\alpha+\beta^{\prime}+\gamma^{\prime}=\pi$.
(b) Set

$$
\delta=\pi-(\alpha+\beta+\gamma), \beta^{\prime}=\beta+\frac{1}{2} \delta \text { and } \gamma^{\prime}=\gamma+\frac{1}{2} \delta .
$$

We first construct an Euclidean triangle with vertices $a, b, c$ having angles $\alpha, \beta^{\prime}, \gamma^{\prime}$, then an isosceles Euclidean triangle with vertices $\omega, b, c$, the angle at the vertex $\omega$ being $\delta$, and such that the points $a$ and $\omega$ are on opposite sides of the line $b c$. Let $\Omega$ be the circle through $b$ and $c$ centered at $\omega$. As $\alpha=\pi-\delta-\beta-\gamma<\pi-\frac{\delta}{2}$, the point $a$ is outside the region encircled by $\Omega$. Let $\Theta_{1}$ be the circle centered at $a$ orthogonal to $\Omega$. As $\frac{\delta}{2}<\beta^{\prime}, b$ is in the region inside $\Theta_{1}$, and so is $c$. If necessary by applying a similitude, we may assume that $a=0$ and $\Theta_{1}=\mathrm{S}^{1}$. The geodesic triangle with vertices $a, b, c$ has the required property.

### 6.9.17 Finiteness of the Automorphism Group ( $g \geqslant 2$ )

Theorem Let X be a compact Riemann surface whose universal covering $\tilde{\mathrm{X}}$ is isomorphic to $\mathbb{D}$. Then the automorphism group $\mathrm{G}=\operatorname{Aut}(\mathrm{X})$ is finite.

Proof Let $\pi: \mathbb{D} \rightarrow \mathrm{X}$ be an analytic map transforming $\mathbb{D}$ into a universal covering of X. Then X can be identified with $\mathbb{D} / \Gamma$, where $\Gamma=\operatorname{Aut}_{X}(\mathbb{D})$ (4.5.1, Remark). Equip $\mathbb{D}$ with the Poincare metric $d_{\mathbb{D}}$. The group $\Gamma$ acts on $\mathbb{D}$ by isometries and equip X with the quotient metric

$$
d_{\mathrm{X}}(x, y)=\inf _{\tilde{y} \in \pi^{-1}(y)} d_{\mathbb{D}}(\tilde{x}, \tilde{y})
$$

where $\tilde{x}$ is an arbitrary point of $\pi^{-1}(x)$.
Equip G with the metric of uniform convergence

$$
d_{\mathrm{G}}\left(g, g^{\prime}\right)=\sup _{x \in \mathrm{X}} d_{\mathrm{X}}\left(g(x), g^{\prime}(x)\right)
$$

The group G acts on X by isometries and X is compact; so G is compact by the Ascoli theorem. We show that it is discrete.

Set $\Gamma^{*}=\Gamma-\left\{1_{\mathbb{D}}\right\}$, and for $u \in \mathbb{D}$, let $\tilde{\rho}(u)=\inf _{\gamma \in \Gamma^{*}} d_{\mathbb{D}}(u, \gamma(u))$. The map $\tilde{\rho}: \mathbb{D} \rightarrow \mathbb{R}_{+}^{*}$ is continuous, and $\tilde{\rho}(u)$ only depends on $x=\pi(u)$. Hence there is a continuous map $\rho: \mathrm{X} \rightarrow \mathbb{R}_{+}^{*}$ such that $\tilde{\rho}=\rho \circ \pi$. Set $\delta=\inf _{x \in \mathrm{X}} \rho(x)$. As X is compact, $\delta>0$.

Lemma Let $g \in$ G. If $d_{\mathrm{G}}\left(1_{\mathrm{X}}, g\right)<\frac{\delta}{2}$, then $g=1_{\mathbb{D}}$.
Proof of the Lemma Assume that $d_{\mathrm{G}}\left(1_{\mathrm{X}}, g\right)=\alpha<\frac{\delta}{2}$. For $u \in \mathbb{D}$, there is a unique $v \in \mathbb{D}$ such that $\pi(v)=g(\pi(u))$ and $d_{\mathbb{D}}(u, v) \leq \alpha$, and the map $\tilde{g}: u \mapsto v$ thus defined is continuous. Since $\pi \circ \tilde{g}=g \circ \pi, \tilde{g}: \mathbb{D} \rightarrow \mathbb{D}$ is holomorphic. With respect to the Poincaré metric, $d_{\mathbb{D}}(u, \tilde{g}(u)) \leq \frac{\delta}{2}$. With respect to the Euclidean metric, $|\tilde{g}(u)-u|$ tends to 0 as $|u|$ tends to 1 (6.9.8, Remark 2). But the function $u \mapsto \tilde{g}(u)-u$ is holomorphic, and hence by the maximum principle is null. So $\tilde{g}=1_{\mathbb{D}}$ and $g=1_{\mathrm{X}}$.

End of the the proof of the theorem The lemma implies that the identity element $1_{\mathrm{X}}$ in the group G is isolated. Hence G is discrete. As it is compact, it is finite.

## Exercises 6.9. (Poincaré geometry)

1. (The upper half-plane.) - Let $\Pi$ be the set of $z \in \mathbb{C}$ such that $\operatorname{Im}(z)>0$.
(a) Show that $z \mapsto \frac{z-i}{z+i}$ is an isomorphism from $\Pi$ onto $\mathbb{D}$.
(b) Show that the homographies preserving $\Pi$ are those of the form

$$
z \mapsto \frac{a z+b}{c z+d} \text { with } a, b, c, d \in \mathbb{R} \text { and } a d-b c>0 .
$$

What are the automorphisms of $\Pi$ (as a Riemann surface)?
(c) Show that there is a unique Riemannian metric on $\Pi$, invariant under the automorphisms, and define it.
2.-Prove the triangular inequality with respect to the distance defined by a Riemannian metric.
3. (Isometries with respect to the Riemannian metric) - Let U be an open subset of $\mathbb{C}$ and $\mu: \mathrm{U} \rightarrow \mathbb{R}_{+}^{*}$ a continuous function.
(a) Let $a \in \mathrm{U}$ and $r>0$ be such that the Euclidean disk $\mathrm{D}_{a, r}$ is contained in U . Let $m$ and $\mathrm{M} \in \mathbb{R}_{+}$be such that

$$
m \leqslant \mu(x) \leqslant \mathrm{M} \text { for all } x \in \mathrm{D}_{a, r}
$$

Show that for all $x \in \mathrm{D}_{a, r}$ and all differentiable path $\gamma$ from $a$ to $x$ in U ,

$$
\mathrm{L}_{\mu}(\gamma) \geqslant m|x-a|
$$

Deduce that

$$
\left(\forall x \in \mathrm{D}_{a, r}\right) \quad m|x-a| \leqslant d_{\mu}(a, x) \leqslant \mathrm{M}|x-a| .
$$

(b) Show that, if $\gamma$ is a differentiable path in U with $\gamma(0)=a$, then the function $t \mapsto d_{\mu}(a, \gamma(t))$ is differentiable at 0 with derivative $\mu(a) \cdot\left|\gamma^{\prime}(0)\right|$.
(c) Deduce that (iii) $\Rightarrow$ (i) in 6.9.4.
(d) With the notation of 6.9 .4 , let $g: \mathrm{U}_{1} \rightarrow \mathrm{U}_{2}$ be a not necessarily holomorphic $\mathrm{C}^{1}$ real map. Assume that $g$ satisfies relation (iii).

Show that $\frac{\mu_{2}(g(x))}{\mu_{1}(x)} \mathrm{T}_{x} g$ is an $\mathbb{R}$-linear isometry of $\mathbb{C}$. Deduce that $g$ is either holomorphic or antiholomorphic.
(e) Does this result continue to hold if $g$ is no longer assumed to be a $\mathrm{C}^{1}$ real map?
4.-Let $a, b, c \in \mathbb{D}$ be points not on the same geodesic. Show that the map $\varphi$ : $(u, v, w) \mapsto \operatorname{Bar}(a, u ;(b, v ; c, w))$ from $\mathrm{T}_{0}=\left\{(u, v, w) \in \mathbb{R}_{+}^{3} \mid u+v+w=1\right\}$ to $\mathbb{D}$ is a homeomorphism from the Euclidean triangle $\mathrm{T}_{0}$ onto the geodesic triangle T with vertices $a, b, c$.
5.-Let $a$ and $b$ be distinct points of $\mathbb{D}$. Show that the set of points of $\mathbb{D}$ equidistant from $a$ and $b$ with respect to the Poincare distance is a geodesic, which will be called the bisector of $(a, b)$.

Let $a, b, c \in \mathbb{D}, \mathrm{M}_{1}, \mathrm{M}_{2}$ and $\mathrm{M}_{3}$ the bisectors of $(a, b),(b, c)$ and $(c, a)$. Show that, if $\mathbf{M}_{1}$ meets $\mathbf{M}_{2}$ at a point $\omega$, then the geodesic $\mathbf{M}_{3}$ passes through $\omega$. Give an example where $\mathrm{M}_{1} \cap \mathrm{M}_{2}=\varnothing$. Show that, in all cases, the circles containing $\mathrm{M}_{1}, \mathrm{M}_{2}$ and $\mathrm{M}_{3}$ are in the same linear sheaf (set of circles orthogonal to two given circles).
6.-Let $\Gamma=\gamma(\mathbb{R})$ be a geodesic with $\gamma$ as in Definition 6.9.9, and $\lambda \in \mathbb{R}$.
(a) Show that there is a unique automorphism $f$ of $\mathbb{D}$ such that $f(\gamma(t))=\gamma(t+\lambda)$ for all $t \in \mathbb{R}$. Then $f$ is said to be the translation of axis $\Gamma$ oriented by $\gamma$ and unit $\lambda$.
(b) Let $x \in \mathbb{D}$. Show that there is a unique geodesic H through $x$ and orthogonal to $\Gamma$. Give an explicit description of $f(x)$.
(c) Let $\mathrm{D}_{1}$ be one of the half-disks limited by $\Gamma$ and $\rho>0$. Show that the set of points of $\mathrm{D}_{1}$ at a (Poincaré) distance $\rho$ from $\Gamma$ is the intersection of $\mathbb{D}$ with a circle of $\Sigma$ which is not orthogonal to the limit circle $S^{1}$. Show that, if $\lambda \neq 0$, then $\Gamma$ is the only geodesic preserved by $f$.
(d) Show that, for all $\mathrm{M} \notin \Gamma, d_{\mu}(\mathrm{M}, f(\mathrm{M}))>|\lambda|$.
7.-Let $f$ be an antiholomorphic isometry of $\mathbb{D}$. Show that $f$ extends to a homeomorphism (also written $f$ ) of $\overline{\mathbb{D}}=\mathbb{D} \cup \mathrm{S}^{1}$. Show that $\left.f\right|_{\mathbf{S}^{1}}$ necessarily fixes two distinct points $u$ and $v$.

Let $\Gamma$ be the geodesic of $\mathbb{D}$ with endpoints $u, v$.
Show that $f$ is the composition of a translation of axis $\Gamma$ and of a symmetry with respect to $\Gamma$.
8. (Automorphisms of $\mathbb{D}$ )-(a) Let $\varphi$ be an isomorphism of the half-plane $\Pi$ (6.9, Exercise 1) on $\mathbb{D}$ and $\lambda \in \mathbb{R}, \mathrm{T}_{\lambda}$ the translation $x \mapsto x+\lambda$ from $\Pi$ to $\Pi$, and $a=$ $\varphi(\infty) \in \mathrm{S}^{1}$.

Show that $\mathrm{P}=\varphi \circ \mathrm{T}_{\lambda} \circ \varphi^{-1}$ is an automorphism of $\mathbb{D}$ whose extension to $\overline{\mathbb{D}}$ has a unique fixed point $a$. It is said to be a parabolic transformation. Show that all tangent circles to $\mathrm{S}^{1}$ at $a$ and in $\overline{\mathbb{D}}$ are preserved by P , and that if $\Gamma$ is a geodesic with endpoint $a$, then so is $\mathrm{P}(\Gamma)$.
(b) Show that an automorphism $f$ of $\mathbb{D}$ is either a rotation, or a translation of axis a geodesic of $\mathbb{D}$, or a parabolic transformation (consider the fixed points of the homography extending $f$ ).
(c) Define a topology on the group $\operatorname{Aut}(\mathbb{D})$ and give a description of the 1-parameter subgroups, i.e. of the continuous homomorphisms from $\mathbb{R}$ to $\operatorname{Aut}(\mathbb{D})$.
9.-Let $a, b, c \in \mathbb{D}$ be points not on the same geodesic. Denote by $\mathrm{M}_{a}$ and $\mathrm{M}_{a}^{\prime}$ the interior and exterior bisector geodesics of ( $[a, b],[a, c]$ ); similarly for $\mathbf{M}_{b}$, etc.

Show that $\mathrm{M}_{a}, \mathrm{M}_{b}$ and $\mathrm{M}_{c}$ are concurrent.
Show that if $\mathbf{M}_{a}$ and $\mathbf{M}_{b}^{\prime}$ intersect at a point $\alpha$, then the geodesic $\mathbf{M}_{c}^{\prime}$ passes through $\alpha$. Show that in general the circles containing $\mathrm{M}_{a}, \mathrm{M}_{b}^{\prime}$ and $\mathrm{M}_{c}^{\prime}$ are in a same linear sheaf.
10.-(a) Show that in Poincaré geometry, the three heights of a triangle lie on three circles in the same linear sheaf.
(b) Does the same hold for the three bisectors?
11.-Let $A, A^{\prime}, B, B^{\prime}$ be four points of $\mathbb{D}, I$ the midpoint of $\left[A, A^{\prime}\right]$ (i.e. $\operatorname{Bar}\left(A, \frac{1}{2}\right.$; $\left.A^{\prime}, \frac{1}{2}\right)$ and $J$ the midpoint of $\left.\left[B, B^{\prime}\right]\right)$. Show that

$$
d_{\mu}(\mathrm{I}, \mathrm{~J}) \leqslant \frac{1}{2}\left(d_{\mu}(\mathrm{A}, \mathrm{~B})+d_{\mu}\left(\mathrm{A}^{\prime}, \mathrm{B}^{\prime}\right)\right) .
$$

(Show that $\mathrm{C}=r_{\mathrm{J}, \pi}\left(\mathrm{A}^{\prime}\right)$ can be obtained from A by applying a translation of axis the geodesic $\Delta$ passing through I and J and of unit $2 d_{\mu}(\mathrm{I}, \mathrm{J})$, then use 6.9, Exercise 6, (d).)

### 6.10 Tiling of the Disk

### 6.10.1 The Dihedral Group

Let X be a topological space. A tiling of X is a locally finite family $\left(\mathrm{A}_{i}\right)_{i \in \mathrm{I}}$ of closed subsets of X with disjoint interiors whose union is X .

Let $n$ be an integer. Let $\Delta_{1}$ and $\Delta_{2}$ be two lines in $\mathbb{C}, \frac{\pi}{n}$ the angle between them, $\Pi_{1}$ and $\Pi_{2}$ two closed half-planes limited by $\Delta_{1}$ and $\Delta_{2}$ respectively, $\Pi_{2}$ being obtained from $\Pi_{1}$ by a rotation of angle $\frac{\pi}{n}+\pi$. Set $\mathrm{S}=\Pi_{1} \cap \Pi_{2}, \Delta_{i}^{+}=\Delta_{i} \cap \mathrm{~S}$. Let $u$ and $v$ be the symmetries with respect to $\Delta_{1}$ and $\Delta_{2}$ respectively, $G$ the isometry group of $\mathbb{C}$ (with respect to the Euclidean metric) generated by $u$ and $v$.

Proposition and Definition (a) The group G has $2 n$ elements and is said to be the dihedral group of order $2 n$.
(b) The family $(g(\mathrm{~S}))_{g \in \mathrm{G}}$ is a tiling of $\mathbb{C}$.
(c) Let $g, g^{\prime} \in \mathrm{G}$ and $x, x^{\prime} \in \mathrm{S}$. If $g(x)=g^{\prime}\left(x^{\prime}\right)$, then $x=x^{\prime}$; and $g(x)=g^{\prime}(x)$ if and only if

```
\(-g^{\prime}=g\) if \(x \in \stackrel{\circ}{\mathrm{~S}}\);
\(-g^{\prime}=g\) or \(g^{\prime}=g w\left(\right.\) resp. gv) if \(x \in \Delta_{1}^{+}-\{0\}\) (resp. \(\left.\Delta_{2}^{+}-\{0\}\right)\).
```

Lemma Let G be a group generated by the two elements $u$ and $v$ such that

$$
u^{2}=v^{2}=(v u)^{n}=e
$$

(identity element). Then any element of G can be written as $(v u)^{k}$ or $(v u)^{k} v$, with $0 \leqslant k \leqslant n-1$. In particular, G has at most $2 n$ elements.

Proof of the Lemma $(v u)^{k}$ can be defined for $k \in \mathbb{Z} /(n)$. Then,

$$
\begin{gathered}
u=u^{-1}=(v u)^{n-1} v, \quad v=v^{-1}=(v u)^{0} v, \quad(v u)^{k}(v u)^{k^{\prime}}=(v u)^{k+k^{\prime}}, \\
(v u)^{k} v \cdot(v u)^{k^{\prime}}=(v u)^{k}(u v)^{k^{\prime}-1} u=(v u)^{n+k-k^{\prime}} v, \\
(v u)^{k} v \cdot(v u)^{k^{\prime}} v=(v u)^{n+k-k^{\prime}} .
\end{gathered}
$$

The set H of elements of the above form is stable, and contains $u, v$ and their inverses; so $\mathrm{H}=\mathrm{G}$.

Proof of the Proposition We may assume that $\Delta_{1}=\mathbb{R}, \Delta_{2}=e^{\frac{\pi}{n} i} \mathbb{R}$, and that S is the set of $\rho . e^{i \theta}$ with $\rho \in \mathbb{R}_{+}, \theta \in\left[0, \frac{\pi}{n}\right]$. Then $v \circ u$ is the rotation of angle $\frac{2 \pi}{n}$ and $(v \circ u)^{n}=1$. If $x=\rho . e^{i \theta}$, then $(v \circ u)^{k}(x)=\rho . e^{i \theta^{\prime}}$ with $\theta^{\prime}=\frac{2 k \pi}{n}+\theta$ and $(v \circ u)^{k} \circ$ $v(x)=\rho . e^{i \theta^{\prime}}$ with $\theta^{\prime}=\frac{(2 k+1) \pi}{n}+\frac{\pi}{n}-\theta$.

For all $y \in \mathbb{C}^{*}$, there exist $g \in \mathrm{G}$ and $x \in \mathrm{~S}$ such that $y=g(x)$. Indeed, if $y=$ $\rho . e^{i \theta^{\prime}}$, then $\theta^{\prime}=\frac{m}{n} \pi+\theta$ for some $m \in \mathbb{Z} /(2 n)$ and $\theta \in\left[0, \frac{\pi}{n}\right]$. If $\theta^{\prime}$ is not a multiple of
$\frac{\pi}{n}$, then the pair $(m, \theta)$ is unique. Now, $y=(v \circ u)^{k}(x)$ with $x=\rho . e^{i \theta}$ if $m=2 k$, and $y=(v \circ u)^{k} v(x)$ with $x=\rho . e^{i\left(\frac{\pi}{n}-\theta\right)}$ if $m=2 k+1$, and $y$ can be uniquely written as $g(x)$. This already shows that all $(v \circ u)^{k}(\stackrel{\circ}{\mathrm{~S}})$ and $(v \circ u)^{k} \circ v(\stackrel{\circ}{\mathrm{~S}})$ are disjoint; i.e. G has $2 n$ elements and all $g(\stackrel{\circ}{\mathrm{~S}})$ are disjoint.

If $\theta^{\prime}=\frac{2 k \pi}{n}\left(\operatorname{resp} \frac{(2 k+1) \pi}{n}\right)$, then $y$ can be written as $g(x)$ in two different ways, namely $y=(v \circ u)^{k}(x)$ and $(v \circ u)^{k-1} \circ v(x)$ with $x=\rho$, (resp. $y=(v \circ u)^{k}(x)$ and $(v \circ u)^{k} \circ v(x)$ with $x=\rho . e^{i \frac{\pi}{n}}$ ). This proves (c) and completes the proof of (b).

### 6.10.2 A Tiling of $\mathbb{D}$

Let $\operatorname{Ismet}(\mathbb{D})$ denote the isometry group of the disk $\mathbb{D}$ with respect to the Poincaré metric, and $\sigma_{a, b}$ the Poincaré symmetry with respect to the geodesic through $a$ and $b$. Let $n, p$ and $q$ be integers such that

$$
\frac{1}{n}+\frac{1}{p}+\frac{1}{q}<1
$$

and let T be a geodesic triangle in $\mathbb{D}$ with angles $\frac{\pi}{n}, \frac{\pi}{p} \frac{\pi}{q}$, and vertices $a, b, c$ (Theorem 6.9.16). We give a presentation by generators and relations of the subgroup of Ismet $(\mathbb{D})$ generated by the symmetries $\sigma_{b, c}, \sigma_{c, a}$ and $\sigma_{a, b}$. Let $\mathrm{L}=\mathrm{L}(\mathrm{U}, \mathrm{V}, \mathrm{W})$ be the free group on three indeterminates. Define

$$
\psi: \mathrm{L} \rightarrow \operatorname{Ismet}(\mathbb{D}) \text { by } \psi(\mathrm{U})=\sigma_{b, c}, \psi(\mathrm{~V})=\sigma_{c, a} \text { and } \psi(\mathrm{W})=\sigma_{a, b}
$$

Then $\psi\left(\mathrm{U}^{2}\right)=\psi\left(\mathrm{V}^{2}\right)=\psi\left(\mathrm{W}^{2}\right)=\psi\left((\mathrm{UV})^{q}\right)=\psi\left((\mathrm{VW})^{n}\right)=\psi\left((\mathrm{WU})^{p}\right)=1$. Let N be the normal subgroup of L generated by $\mathrm{U}^{2}, \mathrm{~V}^{2}, \mathrm{~W}^{2},(\mathrm{UV})^{q},(\mathrm{VW})^{n}$ and (WU) ${ }^{p}$. Set $\mathrm{G}=\mathrm{L} / \mathrm{N}$ and let $\chi: \mathrm{L} \rightarrow \mathrm{G}$ be the canonical map. Define $\varphi: \mathrm{G} \rightarrow \operatorname{Ismet}(\mathbb{D})$ by the commutative diagram:


The group $G$ acts on $\mathbb{D}$ by isometries via $\varphi$.
Let $u, v, w$ be the images of $\mathrm{U}, \mathrm{V}, \mathrm{W}$ in G .
Theorem The family $(g \cdot \mathrm{~T})_{g \in \mathrm{G}}$ of the transforms of T is a tiling of $\mathbb{D}$. In particular, $\varphi$ is injective.

The proof of this theorem is quite long. It will be completed in 6.10.9. We begin by constructing a space E containing T and acted on by G , for which it will be easy to
see that for $g \in \mathrm{G}, g \cdot \mathrm{~T}$ form a tiling of E . Define a map $\pi: \mathrm{E} \rightarrow \mathbb{D}$ compatible with the actions of G. The difficulty lies in showing that $\pi$ is a homeomorphism.

### 6.10.3 Construction of E

Let $\mathscr{R}$ be the equivalence relation on $\mathrm{G} \times \mathrm{T}$ identifying $(g, x)$ with $(g u, x)$ if $x \in$ [ $b, c$ ], with $(g v, x)$ if $x \in[a, c]$, with $(g w, x)$ if $x \in[a, b]$. (In other words, $\mathscr{R}$ is the equivalence relation generated by the above pairs). Denote by E the quotient space $(\mathrm{G} \times \mathrm{T}) / \mathscr{R}$. The map $(g, x) \mapsto g \cdot x$ from $\mathrm{G} \times \mathrm{T}$ to $\mathbb{D}$ induces a continuous map $\pi: \mathrm{E} \rightarrow \mathbb{D}$.

By passing to the quotient, the action $(h,(g, x)) \mapsto(h g, x)$ of G on $\mathrm{G} \times \mathrm{T}$ gives an action of G on E and $\pi$ is compatible with the actions of G on E and $\mathbb{D}$.

Let $\mathrm{G}_{a}$ (resp. $\mathrm{G}_{b}$, resp. $\mathrm{G}_{c}$ ) be the subgroup of G generated by $v$ and $w$ (resp. $w$ and $u$, resp. $u$ and $v$ ). For $x \in] a, b[$, (resp. $] b, c[$, resp. $] c, a\left[\right.$ ), set $\mathrm{G}_{x}=\{e, w\}$ (resp. $\{e, u\}$, resp. $\{e, v\})$; for $x \in \stackrel{\circ}{\mathrm{~T}}$, set $\mathrm{G}_{x}=\{e\}$. If $(g, x) \sim\left(g^{\prime}, x^{\prime}\right)$ with respect to $\mathscr{R}$, then $x=x^{\prime}$. Indeed, the projection $\mathrm{pr}_{2}: \mathrm{G} \times \mathrm{T} \rightarrow \mathrm{T}$ is compatible with $\mathscr{R}$. Now, $(g, x) \sim\left(g^{\prime}, x\right)$ if and only if $g^{\prime} \in g \cdot \mathrm{G}_{x}$. Indeed, the relation thus described is an equivalence relation which implies $\mathscr{R}$ and is satisfied by all pairs generating $\mathscr{R}$.

Let $\chi$ be the canonical map from $\mathrm{G} \times \mathrm{T}$ to E . Then the map $\iota: x \mapsto \chi(e, x)$ from T to E is injective, and identifies T with its image in E .

### 6.10.4 The Star of a in E

For $x \in \mathrm{~T}$, let $\mathrm{E}_{x}$ be the image of $\mathrm{G}_{x} \times \mathrm{T}$ in E ; this set is called the star of $x$. An open star of $x$ is the image $\mathrm{E}_{x}^{\prime}$ of $\mathrm{G}_{x} \times \mathrm{T}_{x}^{\prime}$ in E , where $\mathrm{T}_{x}^{\prime}$ is T without the sides that do not contain $x$ (for example $\mathrm{T}_{a}^{\prime}=\mathrm{T}-[b, c], \mathrm{T}_{x}^{\prime}=\mathrm{T}-([b, c] \cup[c, a])$ if $\left.x \in\right] a, b[$, $\mathrm{T}_{x}^{\prime}=\stackrel{\circ}{\mathrm{T}}$ if $x \in \stackrel{\circ}{\mathrm{~T}}$ ).

For all $x$, the open start of $x$ is open in E since it is the image of a saturated open subset of $\mathrm{G} \times \mathrm{T}$; the closed star of $x$ is a compact neighbourhood of $x$.

Lemma 6.9 The map $\pi$ induces a homeomorphism from $\mathrm{E}_{a}$ onto a neighbourhood of a in $\mathbb{D}$.

Proof (a) Reduction to the case $a=0$. Let $h$ be an automorphism of $\mathbb{D}$ such that $h(a)=0$. Set $\widetilde{\mathrm{T}}=h(\mathrm{~T})$, and let $\widetilde{\mathrm{E}}$ be the quotient of $\mathrm{G} \times \widetilde{\mathrm{T}}$ by the equivalence relation obtained by transferring $\mathscr{R}$. There is a commutative diagram

where the horizontal arrows are homeomorphisms. The space $\widetilde{\mathrm{E}}$ is obtained from $\widetilde{\mathrm{T}}$ in the same manner in which E is obtained from T. Hence we may assume that $a=0$.
(b) Proof of the lemma in this case. The geodesics $\Gamma_{a, c}$ and $\Gamma_{a, b}$ are then part of Euclidean lines $\Delta_{1}$ and $\Delta_{2}$, the Poincaré symmetries $\sigma_{a, c}$ and $\sigma_{a, b}$ are induced by Euclidean symmetries with respect to $\Delta_{1}$ and $\Delta_{2}$, and T is a neighbourhood of 0 in a sector $S$ of angle $\frac{\pi}{n}$.

The group $\mathrm{G}_{a}$ can be identified with the dihedral group of order $2 n$. The equivalence relation induced on $\mathrm{G}_{a} \times \mathrm{T}$ by the map $(g, x) \mapsto g \cdot x$ is the one described in Proposition 6.10.1, (c); it implies the one induced by $\mathscr{R}$. Hence, $\left.\pi\right|_{\mathrm{E}_{a}}$ is injective. As T is a neighbourhood of $a$ in S , the set $\pi\left(\mathrm{E}_{a}\right)=\bigcup_{g \in \mathrm{G}_{a}} g \cdot \mathrm{~T}$ is a neighbourhood of $a$ in $\bigcup_{g \in \mathrm{G}_{a}} g \cdot \mathrm{~S}=\mathbb{C}$.

As $\mathrm{G}_{a}$ is finite, $\mathrm{G}_{a} \times \mathrm{T}$ is compact, and so $\mathrm{E}_{a}$ is quasi-compact; and as $\mathbb{D}$ is Hausdorff, $\pi$ induces a homeomorphism from $\mathrm{E}_{a}$ onto its image.

### 6.10 .5

Lemma 6.10 The map $\pi: \mathrm{E} \rightarrow \mathbb{D}$ is etale.
Proof The set $\stackrel{\circ}{\mathrm{T}}^{\circ}$ in open in E, and $\pi$ induces the identity on $\stackrel{\circ}{\mathrm{T}}$. Hence $\pi$ is etale on $\stackrel{\circ}{\mathrm{T}}$.

For $x \in] a, b\left[\right.$, the map $\left.\pi\right|_{\mathrm{E}_{x}}$ is injective since $\mathrm{T}-[a, b]$ and its image under $\sigma_{a, b}$ are in disjoint open half-disks, and so are disjoint.

As T is a neighbourhood of $x$ in the closed half-disk limited by $\Gamma_{a, b}$, the set $\pi\left(\mathrm{E}_{x}\right)=\mathrm{T} \cup \sigma_{a, b}(\mathrm{~T})$ is a neighbourhood of $x$ in $\mathbb{D}$. As $\mathrm{E}_{x}$ in quasi-compact, $\pi$ induces a homeomorphism from $\mathrm{E}_{x}$ onto its image.

Hence, $\pi$ is etale at $x$. Similarly, $\pi$ is etale at all points either of $] b, c$ [ or of $] c, a[$.
By Lemma 1 (6.10.4), $\pi$ is etale at $a$. Similarly, $\pi$ is etale at $b$ and at $c$, hence on T.

Let $m=\chi(g, x) \in \mathrm{E}$. There is a commutative diagram

where the horizontal arrows are homeomorphisms, and $\pi$ is etale at $x$. Hence $\pi$ is etale at $m=g \cdot x$.

### 6.10.6

Lemma 6.11 There exists $r>0$ such that, with respect to the Poincare metric, for all $t \in \mathrm{E}$, there is a continuous section s of E defined on the ball of radius $r$ centered at $\pi(t)$, with $s(\pi(t))=t$.

Proof For $t \in \mathrm{~T}$, let $\mathrm{R}(t)$ be the upper bound in $\mathbb{R}_{+}$of all $r$ for which there is a continuous section $s$ of E defined on the ball $\mathrm{D}_{t, r}$ of radius $r$ centered at $t$ relative to the Poincaré metric, with $s\left(t^{\prime}\right)=t^{\prime}$ for all $t^{\prime} \in \mathrm{D}_{t, r} \cap \mathrm{~T}$. As E is etale, $\mathrm{R}(t)>0$ for all $t$. For $t^{\prime} \in \mathrm{D}_{t, r}$,

$$
\mathrm{D}_{t^{\prime}, r-d\left(t, t^{\prime}\right)} \subset \mathrm{D}_{t, r} ;
$$

Hence $\mathrm{R}\left(t^{\prime}\right) \geqslant \mathrm{R}(t)-d\left(t, t^{\prime}\right)$, and the function R is either constant, equal to $+\infty$, or Lipschitz with constant 1. In particular, it is continuous; as T is compact, it has a minimum $\mathrm{R}_{0}>0$ and any $r<\mathrm{R}_{0}$ has the required property for all $t \in \mathrm{~T}$.

Let $t=g \cdot t_{0} \in \mathrm{E}$ with $t_{0} \in \mathrm{~T}$, and $s_{0}: \mathrm{D}_{t_{0}, r} \rightarrow \mathrm{E}$ be a continuous section with $s_{0}\left(t_{0}\right)=t_{0}$. Then $s: x \mapsto g \cdot s_{0}\left(g^{-1}(x)\right)$ is a continuous section of $\mathrm{D}_{\pi(t), r}$ in E with $s(\pi(t))=t$. Hence there is some $r<\mathrm{R}_{0}$ has the required property for all $t \in \mathrm{E}$.

### 6.10.7

Lemma 6.12 Let B be a metric space and E an etale Hausdorff space over B. Suppose there exists $r>0$ such that, for all $t \in \mathrm{E}$, there is a continuous section $s$ of E defined on the ball of radius $r$ centered at $\pi(t)$, with $s(\pi(t))=t$. Suppose also that the intersection of two balls of B is connected. Then E is a covering of B .

Proof Let $x \in \mathrm{~B}$ and U be the ball of radius $r$ centered at $x$. Set $\mathrm{F}=\mathrm{E}(x)$. For all $t \in \mathrm{~F}$, there is a continuous section $s_{t}: \mathrm{U} \rightarrow \mathrm{E}$ such that $s_{t}(x)=t$. For $t \neq t^{\prime}$, $s_{t}(y) \neq s_{t^{\prime}}(y)$ for all $y \in \mathrm{U}$. Indeed, U is connected and the set of $y$ such that $s_{t}(y)=$ $s_{t^{\prime}}(y)$ is open since E is etale, closed since E is Hausdorff, and does not contain $x$; so is empty. The map $\varphi:(y, t) \mapsto s_{t}(y)$ from $\mathrm{U} \times \mathrm{F}$ to $\left.\mathrm{E}\right|_{\mathrm{U}}$ is injective, a U-morphism, and is open. We show that it is surjective:

Let $\left.t^{\prime} \in \mathrm{E}\right|_{\mathrm{U}}, x^{\prime}=\pi\left(t^{\prime}\right)$ and $\mathrm{U}^{\prime}$ the ball of radius $r$ centered at $x^{\prime}$. There is a continuous section $s^{\prime}: \mathrm{U}^{\prime} \rightarrow \mathrm{E}$ with $s^{\prime}\left(x^{\prime}\right)=t^{\prime}$. Since $x^{\prime} \in \mathrm{U}, x \in \mathrm{U}^{\prime}$. Set $t=s^{\prime}(x)$. The continuous sections $s^{\prime}$ and $s_{t}$ agree at $x$; as E is etale and Hausdorff and $\mathrm{U} \cap \mathrm{U}^{\prime}$ is connected, they agree on $\mathrm{U} \cap \mathrm{U}^{\prime}$. In particular, $\varphi\left(x^{\prime}, t\right)=s_{t}\left(x^{\prime}\right)=t^{\prime}$. Hence, $\varphi$ is a homeomorphism and $\left.\mathrm{E}\right|_{\mathrm{U}}$ is a trivial covering.

### 6.10 .8

Lemma 6.13 The space E is Hausdorff and connected.
Proof (a) E is Hausdorff: The projection $\mathrm{pr}_{2}: \mathrm{G} \times \mathrm{T} \rightarrow \mathrm{T}$ is compatible with the equivalence relation $\mathscr{R}$, hence defines a continuous map $q: \mathrm{E} \rightarrow \mathrm{T}$. Therefore, $g \cdot t$ and $g^{\prime} \cdot t^{\prime}$ can be separated for $t \neq t^{\prime}$. If $g \cdot t \neq g^{\prime} \cdot t$, then the sets $g \cdot \mathrm{E}_{t}^{\prime}$ and $g^{\prime} \cdot \mathrm{E}_{t}^{\prime}$ are open and disjoint.
(b) E connectedness: Let A be a clopen subset of E . For all $g$, the image of $\{g\} \times \mathrm{T}$ is in A or in $\mathrm{E}-\mathrm{A}$. Let M be the set of $g \in \mathrm{G}$ for which the image of $\{g\} \times \mathrm{T}$ is in A. If $g \in \mathrm{M}$, then $g u \in \mathrm{M}$ since $(g u, x)$ and $(g, x)$ have the same image for $x \in[b, c]$. As $u^{-1}=u$ in G, $g u^{-1} \in$ M. Similarly, $g v, g v^{-1}, g w, g w^{-1}$ are in M. Since G is generated by $u, v, w, \mathrm{M}=\mathrm{G}$ or $\mathrm{M}=\varnothing$, and so $\mathrm{A}=\mathrm{E}$ or $\mathrm{A}=\varnothing$.

### 6.10.9 Proof of Theorem 6.10.2

In the space $\mathbb{D}$ equipped with the Poincare distance, the balls are the Euclidean disks. Hence the intersection of two balls is connected. Thanks to lemmas 3 and 5, the assumptions of lemma 4 hold. Therefore $E$ is a covering of $\mathbb{D}$. As $\mathbb{D}$ is simply connected, this covering is trivial. As E connected, it is of degree 1. In other words, $\pi$ is a homeomorphism.

### 6.10.10 Crossed Products

Let A and B be two groups and $(b, a) \mapsto{ }^{b} a$ an action of B on A by automorphisms (hence ${ }^{b b^{\prime}} a={ }^{b}\left({ }^{b^{\prime}} a\right)$ and ${ }^{b}\left(a a^{\prime}\right)={ }^{b} a \cdot{ }^{b} a^{\prime}$ for $a, a^{\prime} \in \mathrm{A}, b, b^{\prime} \in \mathrm{B}$ ). For the sake of convenience, we assume that $\mathrm{A} \cap \mathrm{B}=\{e\}$, where $e$ is the identity element common to A and B .

Proposition and Definition There is a group C such that A is a normal subgroup of $\mathrm{C}, \mathrm{B}$ a subgroup of $\mathrm{C}, b \cdot a \cdot b^{-1}={ }^{b}$ a for $a \in \mathrm{~A}$ and $b \in \mathrm{~B}$, and such that canonical map $\chi \circ \iota: \mathrm{B} \rightarrow \mathrm{C} / \mathrm{A}$ is an isomorphism. If C and $\mathrm{C}^{\prime}$ satisfy these conditions, then there is a unique isomorphism from C onto $\mathrm{C}^{\prime}$ inducing the identity on A and B .

We say that C is thecrossed product (or semi-direct product) of A by B. It is written $\mathrm{C}=\mathrm{A} \rtimes \mathrm{B}$.

Proof (a) Uniqueness: Let C be a group with the desired property. The map $(a, b) \mapsto$ $a \cdot b$ from $\mathrm{A} \times \mathrm{B}$ to C is bijective. For $a, a^{\prime} \in \mathrm{A}$ and $b, b^{\prime} \in \mathrm{B}, a \cdot b \cdot a^{\prime} \cdot b^{\prime}=a$. ${ }^{b} a^{\prime} \cdot b \cdot b^{\prime}$. Hence C can be identified with the product set $\mathrm{A} \times \mathrm{B}$ equipped with the composition law

$$
\left((a, b),\left(a^{\prime}, b^{\prime}\right)\right) \mapsto\left(a^{b} a^{\prime}, b b^{\prime}\right)
$$

(b) Existence: The composition law defined by

$$
(a, b) \cdot\left(a^{\prime}, b^{\prime}\right)=\left(a^{b} a^{\prime}, b b^{\prime}\right)
$$

on $\mathrm{A} \times \mathrm{B}$ is associative with

$$
(a, b) \cdot\left(a^{\prime}, b^{\prime}\right) \cdot\left(a^{\prime \prime}, b^{\prime \prime}\right)=\left(a^{b} a^{\prime} b b^{\prime} a^{\prime \prime}, b b^{\prime} b^{\prime \prime}\right)
$$

and each element has an inverse

$$
(a, b)^{-1}=\left(b^{b^{-1}} a^{-1}, b^{-1}\right)
$$

The set $\mathrm{A} \times \mathrm{B}$ equipped with this law is a group C with the desired property.

### 6.10 .11

Consider the group $\mathrm{G}=\mathrm{L}(\mathrm{U}, \mathrm{V}, \mathrm{W}) / \mathrm{N}$ defined in (6.10.2), and the homomorphism $\varepsilon: \mathrm{G} \rightarrow\{+1,-1\}$ defined by $\varepsilon(u)=\varepsilon(v)=\varepsilon(w)=-1$.

Set $\mathrm{G}_{+}=\operatorname{Ker} \varepsilon=\varphi^{-1}($ Aut $\mathbb{D})$. The group $\mathrm{G}_{+}$is a normal subgroup of index 2 in G , and G can be identified with the crossed product $\mathrm{G}_{+} \rtimes(u)$, where $(u)=\{e, u\}$.

Consider the elements $\alpha=v w, \beta=w u$ and $\gamma=u v$ in $\mathrm{G}_{+}$. As $\alpha^{n}=\beta^{p}=\gamma^{q}=$ $\alpha \beta \gamma=e$, the canonical map

$$
\psi_{1}: \mathrm{L}(\alpha, \beta, \gamma) \rightarrow \mathrm{G}_{+}
$$

induces $\varphi_{1}: \Lambda \rightarrow \mathrm{G}_{+}$, where $\Lambda$ is the quotient of $\mathrm{L}(\alpha, \beta, \gamma)$ by the normal subgroup generated by $\alpha^{n} \beta^{p}$, $\gamma^{q}$ and $\alpha \beta \gamma$.

Proposition The homomorphism $\varphi_{1}: \Lambda \rightarrow \mathrm{G}_{+}$is an isomorphism.
Proof Define an action of $(u)$ on $\Lambda$ such that $\varphi_{1}$ is a $(u)$-homomorphism. The action of $(u)$ on $\mathrm{G}_{+}$by inner automorphism in G satisfies
(*)

$$
\left\{\begin{array}{l}
{ }^{u} \alpha=\beta^{-1} \alpha^{-1} \beta, \\
{ }^{u} \beta=\beta^{-1}, \\
{ }^{u} \gamma=\gamma^{-1} .
\end{array}\right.
$$

Define an endomorphism $x \mapsto{ }^{u} x$ of $\mathrm{L}(\alpha, \beta, \gamma)$ by formulas $(*)$.
It can be checked to be involutive, i.e. ${ }^{u}\left({ }^{u} x\right)=x$. Hence there is an action of $(u)$ on $\mathrm{L}(\alpha, \beta, \gamma)$. Passing to the quotient, it can be checked to define an action of $(u)$ on $\Lambda$ by automorphisms. By construction, $\varphi_{1}$ is a (u)-homomorphism.

The homomorphism $\varphi_{1}$ defines a homomorphism

$$
\overline{\varphi_{1}}: \Lambda \rtimes(u) \rightarrow \mathrm{G}_{+} \rtimes(u)=\mathrm{G} .
$$

We show that $\overline{\varphi_{1}}$ is an isomorphism by constructing an inverse. The homomorphism $\mathrm{U} \mapsto u, \mathrm{~V} \mapsto u \cdot \gamma, \mathrm{~W} \mapsto \beta \cdot u$ from $\mathrm{L}(\mathrm{U}, \mathrm{V}, \mathrm{W})$ to $\Lambda \rtimes(u)$ can be checked to pass to the quotient and to define a homomorphism $\sigma: \mathrm{G} \rightarrow \Lambda \rtimes(u)$. As can be verified, this is an inverse of $\overline{\varphi_{1}}$.

In the commutative diagram

$\overline{\varphi_{1}}$ is an isomorphism, and so $\varphi_{1}$ is an isomorphism.

### 6.10.12

Proposition Let $g \in \mathrm{G}_{+}$. If $\varphi(g) \in$ Aut $\mathbb{D}$ has a fixed point in $\mathbb{D}$, then $g$ is the conjugate in $\mathrm{G}_{+}$of a power of either $\alpha, \beta$ or $\gamma$.

Proof We may assume that $g \neq e$. Let $x \in \mathbb{D}$ be such that $g \cdot x=x$. By Theorem 6.10.2, $x=h \cdot t$ for some $h \in \mathrm{G}, t \in \mathrm{~T}$. Set $g^{\prime}=h^{-1} g h$; then $g^{\prime} \cdot t=t . t \notin \stackrel{\circ}{\mathrm{~T}}$ since the images of $\stackrel{\circ}{\mathrm{T}}$ are disjoint. If $t \in] a, b\left[\right.$, then $g^{\prime}=w$, which is impossible since $g \in \mathrm{G}_{+}$, which is normal. Similarly $\left.t \notin\right] b, c[$ and $t \notin] c, a[$. Hence $t$ is a vertex. It follows from the construction of $\mathrm{E}(6.10 .3)$ that the stabilizer of $a$ in $\mathrm{G}^{\text {is }} \mathrm{G}_{a}$. Hence, if $t=a$, then $g^{\prime} \in \mathrm{G}_{a} \cap \mathrm{G}_{+}=\left\{1, \alpha, \ldots, \alpha^{n-1}\right\}$. Similarly, if $t=b$ (resp. $c$ ), then $g^{\prime}$ is a power of $\beta$ (resp. $\gamma$ ). This gives the result when $h \in \mathrm{G}_{+}$. Otherwise $h=h^{\prime} u$, and $g$ is conjugate in $\mathrm{G}_{+}$to

$$
u \alpha^{k} u^{-1}=\beta^{-1} \alpha^{n-k} \beta, \text { or to } u \beta^{k} u^{-1}=\beta^{p-k}, \text { or to } u \gamma^{k} u^{-1}=\gamma^{q-k}
$$

### 6.10 .13

Proposition Let H be a normal subgroup of $\mathrm{G}_{+}$. Suppose that the images of $\alpha, \beta$ and $\gamma$ in $\mathrm{G}_{+} / \mathrm{H}$ are elements of orders $n, p$ and $q$ respectively.
(a) The group H acts freely on $\mathbb{D}$ and $\mathbb{D}$ is a covering of the quotient $\mathrm{X}=\mathbb{D} / \mathrm{H}$. The space X is Hausdorff.
(b) There is a unique Riemann surface structure on X such that the canonical map $\chi: \mathbb{D} \rightarrow \mathrm{X}$ is analytic.
(c) The space X is compact if and only if H is of finite index in G .

Proof (a) Let $h \in \mathrm{H}$ be such that $\varphi(h)$ has a fixed point in $\mathbb{D}$. By Proposition 6.10.12, $h$ is conjugate to some $\alpha^{k}, \beta^{k}$ or $\gamma^{k}$. Its image $\chi(h) \in \mathrm{G} / \mathrm{H}$ is then conjugate to $\chi(\alpha)^{k}$ (resp. $\chi(\beta)^{k}$, resp. $\left.\chi(\gamma)^{k}\right)$; As it is zero, $k$ is a multiple of $n$ (resp. $p$, resp. $q$ ), and $h=1$. Hence H acts freely on $\mathbb{D}$.

For $t \in \mathrm{~T}, g \in \mathrm{G}$ and $h \neq e$ in $\mathrm{H}, h \cdot g \cdot \mathrm{E}_{t}^{\prime} \cap g \cdot \mathrm{E}_{t}^{\prime}=\varnothing$ since $g^{-1} h g$ is a nonzero element of H , hence is not in $\mathrm{G}_{t}$. Thus the assumptions of Theorem 4.3.13 hold, and $\mathbb{D}$ is a covering of X . As $\pi$ identifies E with $\mathbb{D}$, the continuous map from E to T induced by $\mathrm{pr}_{2}: \mathrm{G} \times \mathrm{T} \rightarrow \mathrm{T}$ defines, by passing to the quotient, a continuous map $\mathrm{X} \rightarrow \mathrm{T}$. Thus the images of $g \cdot t$ and $g^{\prime} \cdot t^{\prime}$ can be separated for $t \neq t^{\prime} \in \mathrm{T}$. If $g \cdot t$ and $g^{\prime} \cdot t$ have distinct images in X , then

$$
g^{\prime} \notin g \cdot \mathrm{H} \cdot \mathrm{G}_{t}=\mathrm{H} \cdot g \cdot \mathrm{G}_{t},
$$

and the open saturated subsets $\mathrm{H} \cdot g \cdot \mathrm{E}_{t}^{\prime}$ and $\mathrm{H} \cdot g^{\prime} \cdot \mathrm{E}_{t}^{\prime}$ of $\mathbb{D}$ are disjoint. Hence X is Hausdorff.
(b) Every continuous section $\sigma: \mathrm{U} \rightarrow \mathbb{D}$, where U is an open subset of X is a chart with domain U , and these charts form an atlas, the change of charts being induced by automorphisms of $\mathbb{D}$. The map $\chi$ is analytic with respect to the Riemann surface structure defined by this atlas, and it follows from 6.1.10 that this structure is the only one with the required property.
(c) Let $\left(g_{i}\right)_{i \in \mathrm{I}}$ be a family of representatives for the elements of $\mathrm{G} / \mathrm{H}$. The image of $\bigcup_{i} g_{i} \cdot \mathrm{~T}$ is X . If $\mathrm{G} / \mathrm{H}$ is finite, then so is I , and both $\bigcup_{i} g_{i} \cdot \mathrm{~T}$ and X are compact.

Let $t \in \stackrel{\circ}{\mathrm{~T}}$. The set $\mathrm{G} \cdot t$ is discrete and closed in $\mathbb{D}$; hence so is its image in X . If X is compact, this image is finite; hence, as it is in bijective correspondence with $\mathrm{G} / \mathrm{H}, \mathrm{G} / \mathrm{H}$ is finite if X is compact.

### 6.10 .14

We keep the assumptions of 6.10.13 and further assume that H is of finite index ${ }^{7}$ in G . The group $\Gamma=\mathrm{G}_{+} / \mathrm{H}$ can be identified with an automorphism group of the Riemann surface X , and the space $\mathrm{Y}=\mathrm{X} / \Gamma$ can be identified with $\mathbb{D} / \mathrm{G}_{+}$. Moreover, Y has a natural Riemann surface structure (6.7.1), and $X$ is a ramified Galois covering of Y. Let $\pi$ be the canonical map from X onto Y , and $\hat{a}, \hat{b}$ and $\hat{c}$ the images of $a, b$ and $c$ under $\pi \circ \chi$.

[^25]Proposition The ramification set of X over Y is $\{\hat{a}, \hat{b}, \hat{c}\}$. For $\xi \in \mathrm{X}$ over $\hat{a}$ (resp. $\hat{b}$, resp. $\hat{c}$ ), the ramification index of X over Y at $\xi$ is $n$ (resp. p, resp. q).

Proof For $\xi=\chi(x) \in \mathrm{X}$, the stabilizer $\Gamma_{\xi}$ of $\xi$ in $\Gamma$ can be identified with the stabilizer $\left(\mathrm{G}_{+}\right)_{x}$ of $x$ in $\mathrm{G}_{+}$. Indeed, if $\chi(g) \cdot \xi=\xi$, then $g \cdot x=h \cdot x$ for some $h \in \mathrm{H}$, and $h^{-1} g \in \mathrm{G}_{x}$. Hence the canonical map $\mathrm{G}_{x} \rightarrow \Gamma_{\xi}$ is surjective; its kernel is $\mathbb{G}_{x} \cap \mathrm{H}$ and it is trivial since H acts freely.

Write $x$ as $g \cdot t$ with $g \in \mathrm{G}, t \in \mathrm{~T}$. The group $\left(\mathrm{G}_{+}\right)_{x}$ is conjugate to $\left(\mathrm{G}_{+}\right)_{t}$, hence nontrivial if and only if $t$ is one of the 3 vertices of T. If $t$ is a vertex, then $x=g^{\prime} \cdot t$ with $g^{\prime} \in \mathrm{G}_{+}$. So $\xi=\chi\left(g^{\prime}\right) \cdot \chi(t)$, and $\pi(\xi)=\hat{a}, \hat{b}$ or $\hat{c}$. The ramification index of X at $\xi$ is the cardinal of $\left(\mathrm{G}_{+}\right)_{t}$, i.e. $n$ if $\pi(\xi)=\hat{a}, p$ if $\pi(\xi)=\hat{b}$ and $q$ if $\pi(\xi)=\hat{c}$.

### 6.10 .15

Proposition Under the assumptions of 6.10.14, the Riemann surface Y has genus 0.
Proof We construct a triangulation of Y. We keep the notation of this section as well as that of 6.5.1. Let $\theta: \mathbb{T} \rightarrow \mathrm{T}$ be a $\mathrm{C}^{1}$ homeomorphism on $\mathbb{T}-\left\{\alpha_{0}, \alpha_{1}, \alpha_{2}\right\}$ with Jacobian $>0$, and such that $\theta \circ \iota_{\nu}$ is a geodesic path for $\nu=1,2$ or 3 . Set

$$
\theta_{1}(t)=u(\theta((1+j) \bar{t}))
$$

Set $f_{0}=\pi \circ \chi \circ \theta$ and $f_{1}=\pi \circ \chi \circ \theta_{1}$. Let $a_{1}$ (resp. $a_{2}$, resp. $a_{3}$ ) be the image under $\pi \circ \chi$ of the geodesic path from $a$ to $b$ (resp. from $b$ to $c$, resp. from $c$ to $a$ ), and $s_{1}=\hat{a}, s_{2}=\hat{b}, s_{3}=\hat{c}$. This gives a triangulation $\tau$ of Y with 3 vertices, 3 edges and 2 faces. The Euler-Poincaré characteristic of Y (equipped with $\tau$ ) is 2, and so its genus is 0 .

Remark Using this triangulation, Y can be checked to be homeomorphic to the Riemann sphere.

### 6.10 .16

Let $d$ be the cardinal of $\Gamma$, i.e. the index of H in $\mathrm{G}_{+}$.
Corollary The genus $g$ of X is $1+\frac{d}{2}\left(1-\frac{1}{n}-\frac{1}{p}-\frac{1}{q}\right)$.
Proof There are $\frac{d}{n}$ ramification points of index $n, \frac{d}{p}$ of index $p$ and $\frac{d}{q}$ of index $q$. The Riemann-Hurwitz formula 6.6 .13 gives

$$
g-1=-d+\frac{1}{2}\left(\frac{d}{n}(n-1)+\frac{d}{p}(p-1)+\frac{d}{q}(q-1)\right) .
$$

### 6.10 .17

Theorem Let $g, n, p$ and $q$ be integers such that $g \geqslant 2, \frac{1}{n}+\frac{1}{p}+\frac{1}{q}<1$, and $d=$ $\frac{2 g-2}{1-\frac{1}{n}-\frac{1}{p}-\frac{1}{q}}$ is integral. Assume that there is a group $\Gamma$ of order $d$ with three elements $\hat{\alpha}, \hat{\beta}, \hat{\gamma}$ of respective orders n, $p$ and $q$ satisfying $\hat{\alpha} \hat{\beta} \hat{\gamma}=e$. Then there is a Riemann surface X of genus $g$ such that $\Gamma$ is identified with an automorphism group of X .

Proof By Proposition 6.10.11, $\Gamma$ can be identified with a quotient of $\mathrm{G}_{+}$in such a way that $\hat{\alpha}, \hat{\beta}$ and $\hat{\gamma}$ are the respective images of $\alpha, \beta$ and $\gamma$. Then the assumptions 6.10.14 hold and the Riemann surface $X=\mathbb{D} / H$ has the required property.

### 6.10 .18

Corollary There is a Riemann surface of genus 3 with an automorphism group of order 168.

This corollary follows from Theorem 6.10.17 applied to $g=3, n=2, p=3$, $q=7$, which implies $d=168$, and from the following lemma:

Lemma Let $k=\mathbb{F}_{2}$ be the field $\mathbb{Z} /(2)$. The group $\Gamma=\operatorname{GL}\left(k^{3}\right)$ of invertible $3 \times 3$ matrices with entries in $k$ has 168 elements. There are three elements $\mathrm{A}, \mathrm{B}, \mathrm{C}$ in $\Gamma$ of respective order 2,3 and 7 such that $\mathrm{A} \cdot \mathrm{B} \cdot \mathrm{C}=\mathrm{I}$.

Proof There is a bijection between the group $\Gamma$ and the set of bases $\left(\mathrm{V}_{1}, \mathrm{~V}_{2}, \mathrm{~V}_{3}\right)$ for $k^{3}$. For $\mathrm{V}_{1}$, there are 7 possible choices; $\mathrm{V}_{1}$ being chosen, there are 6 ways of choosing $V_{2}$ and, $V_{1}$ and $V_{2}$ being chosen, 4 of choosing $V_{3}$. The order of $G$ is therefore $7 \times 6 \times 4=168$.

Take

$$
A=\left(\begin{array}{lll}
1 & 0 & 1 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right) \quad B=\left(\begin{array}{lll}
0 & 0 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right)
$$

in $\Gamma$ of respective orders 2 and 3. The "sagittal diagram" (Oh modern mathematics, what have you got into?)

shows that $\mathrm{A} \circ \mathrm{B}$ has order 7 . Then $\mathrm{C}=(\mathrm{A} \circ \mathrm{B})^{-1}$ has the required property.

## Exercises 6.10. (Disk tilings)

1.-Observe M. C. Escher's engraving. Ignoring the colours, it represents a tiling of the disk of the form $(g \mathrm{P})_{g \in \mathrm{G}_{+}}$, where $\mathrm{G}_{+}$is defined as in 6.10 .11 (the fundamental domain has the shape of a fish).
(a) What are the values of $n, p, q$ ?
(b) Let P be a fish. The elements $g \in \mathrm{G}_{+}$for which $g \mathrm{P}$ has the same colour as P form a subgroup of H . What is the index of H ?

Is the equivalence relation " $g \mathrm{P}$ and $g^{\prime} \mathrm{P}$ are of the same colour" between elements $g, g^{\prime} \in \mathrm{G}, g^{\prime} \in g \mathrm{H}$ or $g^{\prime} \in \mathrm{H} g$ ?

Is the union of four fish of different colours necessarily a fundamental domain for the action of H ?
(c) Give the permutations of $\mathrm{G} / \mathrm{H}=\{g \mathrm{H}\}_{g \in \mathrm{G}}$ defined by the action of G. How many are there?
(d) Consider the subgroup $\Gamma$ of $\mathrm{G}_{+}$consisting of the elements preserving the colouring. Show that $\Gamma$ is the largest normal subgroup of $\mathrm{G}_{+}$in H . Give its index in $\mathrm{G}_{+}$.
(e) Show that $\Gamma$ acts freely on $\mathbb{D}$ and that the Riemann surface $X=\mathbb{D} / \Gamma$ is compact.
(f) Give the equivalence relation identifying the space $B=\mathbb{D} / G_{+}$with the quotient of $P$. Show that $B$ is homeomorphic to the sphere $S^{2}$, and so is a Riemann surface of genus 0 . What are the ramification indices of $\mathbb{D}$ over $B$ ? What is the genus of $\mathrm{B} / \mathrm{H}$ ? of $\mathrm{B} / \Gamma$ ? (The Riemann-Hurwitz formula may be applied.)
(g) Are the spines of the fish geodesics?


The tiling of the Theorem 6.10 .2 for $n=2, p=3, q=7$

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## Chapter 7 <br> Dessins d'Enfants

## Introduction

The Galois group $\mathbb{G}=$ Aut $_{\mathbb{Q}}(\overline{\mathbb{Q}})$, where $\overline{\mathbb{Q}}$ is an algebraic closure of $\mathbb{Q}$, let us say in $\mathbb{C}$, fascinates arithmeticians. This profinite group is hard to grasp. It certainly embeds in the product of $\mathfrak{S}\left(\mathrm{P}^{-1}(0)\right)$, as P runs through the set of irreducible polynomials of $\mathbb{Z}[\mathrm{X}]$, but this set is itself not easy to understand if only because there is no natural way of numbering the roots of a polynomial.

The aim of Grothendieck's theory of dessins d'enfants is to make $\mathbb{G}$ act on finite sets defined combinatorially in order to gather information about it.

The "arithmetic" Riemann surfaces are those that can be defined over a "number field", i.e. a finite extension of $\mathbb{Q}$. In section 1, we explain what a Riemann surface defined over a subfield of $\mathbb{C}$ is. The theory begins with Belyi's theorem, according to which any arithmetic Riemann surface is a covering of the Riemann sphere ramified only over 0,1 and $\infty$. The classification of these coverings becomes a combinatorial problem.

### 7.1 Definability

### 7.1.1 Defining Polynomials

The algebra of meromorphic functions for the Riemann sphere $\Sigma \mathscr{M}(\Sigma)$ can be identified with the field $\mathbb{C}(\mathrm{X})$ of rational fractions in one indeterminate X .

Let ( $\mathrm{S}, \pi$ ) be an analytic ramified covering of $\Sigma$, and consider $\mathscr{M}(\mathrm{S})$ as a $\mathbb{C}(\mathrm{X})$-algebra using $\pi^{*}: \mathbb{C}(\mathrm{X}) \rightarrow \mathcal{M}(\mathrm{S})$. This is the situation in subsection 6.2.5, whose notation we keep (with Y instead of Z). Then $\mathscr{M}(\mathrm{S})$ is an etale algebra over $\mathbb{C}(\mathrm{X})$ (Theorem 6.2.3). If $\eta \in \mathscr{M}(\mathrm{S})$ is a primitive element with minimal polynomial

$$
\mathrm{P}(\mathrm{Y})=\mathrm{Y}^{d}+a_{d-1}(\mathrm{X}) \mathrm{Y}^{d-1}+\cdots+a_{0}(\mathrm{X}) \in \mathbb{C}(\mathrm{X})[\mathrm{Y}]
$$

then the algebra $\mathscr{M}(\mathrm{S})$ can be identified with $\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$, and S is a compactification of

$$
\mathrm{S}^{\prime}=\mathrm{V}_{\mathrm{P}}^{\prime}=\left\{(x, y) \in \mathbb{C}^{2} \mid x \in \mathbb{C}-\Delta_{\mathrm{P}}, \mathrm{P}(x, y)=0\right\}
$$

where $\Delta_{\mathrm{P}}=\Delta_{\mathscr{M}(\mathrm{S}), \eta}$ is the set of $x \in \mathbb{C}$ for which $\mathrm{P}(x, \mathrm{Y}) \in \mathbb{C}[\mathrm{Y}]$ either has a multiple root or is not defined.

Each $a_{i} \in \mathbb{C}(\mathrm{X})$ can be uniquely written as $\frac{p_{i}}{q_{i}}$ with $p_{i}$ and $q_{i}$ coprime in $\mathbb{C}[\mathrm{X}]$ and $q_{i}$ monic. Let $b_{d}$ be the (monic) l.c.m. of the polynomials $q_{i}, b_{i}=b_{d} a_{i} \in \mathbb{C}[\mathrm{X}]$ and $\mathrm{Q}=b_{d} \mathrm{Y}^{d}+b_{d-1} \mathrm{Y}^{d-1}+\cdots+b_{0}$. The polynomial $\mathrm{Q} \in \mathbb{C}[\mathrm{X}, \mathrm{Y}]$ is called the polynomial obtained from $\mathrm{P} \in \mathbb{C}(\mathrm{X})[\mathrm{Y}]$ by eliminating denominators. It will be said to be a defining polynomial of $(\mathrm{S}, \pi)$.

The polynomial Q has the following property: the decomposition of Q into prime factors has no repetition, nor any factor such as $\mathrm{X}-a$. A polynomial $\mathrm{Q} \in \mathbb{C}[\mathrm{X}, \mathrm{Y}]$ with this property is called regular in Y . The set $\Delta_{\mathrm{Q}}=\Delta_{\mathrm{P}}$ is the set of $x \in \mathbb{C}$ for which $\mathrm{Q}(x, \mathrm{Y}) \in \mathbb{C}[\mathrm{Y}]$ has either a multiple root or Y -degree strictly less than that of Q . The space $\mathrm{S}^{\prime}$ can be identified with

$$
\mathrm{V}_{\mathrm{Q}}^{\prime}=\mathrm{V}_{\mathrm{P}}^{\prime}=\left\{(x, y) \in \mathbb{C}^{2} \mid x \notin \Delta_{\mathrm{Q}}, \mathrm{Q}(x, y)=0\right\}
$$

Conversely, if $\mathrm{Q} \in \mathbb{C}[\mathrm{X}, \mathrm{Y}]$ is a regular polynomial in Y , then the set

$$
\mathrm{V}_{\mathrm{Q}}^{\prime}=\left\{(x, y) \in\left(\mathbb{C}-\Delta_{\mathrm{Q}}\right) \times \mathbb{C} \mid \mathrm{Q}(x, y)=0\right\}
$$

together with $\pi:(x, y) \mapsto x$, is a finite covering of $\mathbb{C}-\Delta_{Q}$. By 6.1.10, this covering is induced by a finite analytic ramified covering $\widehat{\mathrm{V}_{\mathrm{Q}}}$ of $\Sigma$. $\widehat{\mathrm{V}_{\mathrm{Q}}}$ is said to be the analytic ramified covering defined by Q .

Proposition The closure of $\mathrm{S}^{\prime}$ in $\mathbb{C}^{2}$ is the algebraic set

$$
\mathrm{V}_{\mathrm{Q}}=\left\{(x, y) \in \mathbb{C}^{2} \mid \mathrm{Q}(x, y)=0\right\}
$$

Proof The set $\mathrm{V}_{\mathrm{Q}}$ is closed and contains $\mathrm{V}_{\mathrm{Q}}^{\prime}$, hence so does its closure. We show the converse inclusion, which concerns the continuity of the roots of the polynomial (see Chap.3, Exercise 8). Let $\left(x_{0}, y_{0}\right)$ be a point such that $\mathrm{Q}\left(x_{0}, y_{0}\right)=0$, and $\left(x_{n}\right)$ a sequence in $\mathbb{C}-\Delta_{\mathrm{Q}}$ tending to $x_{0}$. Let $\nu$ be the multiplicity of $y_{0}$ as a root of $\mathrm{Q}\left(x_{0}, \mathrm{Y}\right)$. For sufficiently small $r>0$, the point $\mathrm{Q}\left(x_{0}, y\right)$ wraps around $0 \nu$ times as $y$ runs through the boundary of the disc $\mathrm{D}_{y_{0}, r}$. Hence the same holds for $\mathrm{Q}\left(x_{n}, y\right)$ for sufficiently large $n$. So for sufficiently large $n, \mathrm{Q}\left(x_{n}, \mathrm{Y}\right)$ has a root in $\mathrm{D}_{y_{0}, r}$. Hence there is a sequence $\left(y_{n}\right)$ tending to $y_{0}$ such that $\left(x_{n}, y_{n}\right) \in \mathrm{V}_{\mathrm{Q}}^{\prime}$ for all sufficiently large $n$.

### 7.1.2 Definability of a Ramified Covering

Let $K$ be a subfield of $\mathbb{C}$ and $A$ an etale algebra over $K(X)$. Then the algebra $A_{\mathbb{C}}=$ $\mathbb{C}(X) \otimes_{K(X)} A$ is an etale algebra over $\mathbb{C}(X)$. By (6.2.4), as a $\mathbb{C}(X)$-algebra, it is isomorphic to $\mathscr{M}(\mathrm{S})$, where ( $\mathrm{S}, \pi$ ) is an analytic ramified covering of $\Sigma$. A primitive
element $\eta$ of the $\mathrm{K}(\mathrm{X})$-algebra A remains primitive in the $\mathbb{C}(\mathrm{X})$-algebra- $\mathrm{A}_{\mathbb{C}}$, and S can be taken to be the Riemann surface $\widehat{\mathcal{S}}\left(\mathrm{A}_{\mathbb{C}}, \eta\right)$.

Proposition and Definition Let $(\mathrm{S}, \pi)$ be an analytic ramified covering of $\Sigma$ and K a subfield of $\mathbb{C}$. The following conditions are equivalent:
(i) there is a defining polynomial $\mathrm{Q} \in \mathrm{K}[\mathrm{X}, \mathrm{Y}]$ of $(\mathrm{S}, \pi)$;
(ii) there is a $\mathrm{K}(\mathrm{X})$-subalgebra A in $\mathscr{M}(\mathrm{S})$ such that the natural map $\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})}$ $\mathrm{A} \rightarrow \mathscr{M}(\mathrm{S})$ is an isomorphism.

If these conditions hold, then $(\mathrm{S}, \pi)$ is said to be definable over K . An algebra A satisfying condition (ii) is called a K-structure on (S, $\pi$ ).

Proof (a) (i) $\Rightarrow$ (ii): Set $\mathrm{A}=\mathrm{K}(\mathrm{X})[\mathrm{Y}] /(\mathrm{Q})$. Then $\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{Q})$ can be identified with $\mathscr{M}(\mathrm{S})$.
(b) (ii) $\Rightarrow$ (i): Let $y$ be a primitive element of A over $\mathrm{K}(\mathrm{X}), \mathrm{P} \in \mathrm{K}(\mathrm{X})[\mathrm{Y}]$ its minimal polynomial and $\mathrm{Q} \in \mathrm{K}[\mathrm{X}, \mathrm{Y}]$ obtained by eliminating the denominators. Then Q is a defining polynomial.

Remark $\mathbb{C} \otimes_{\mathrm{K}} \mathrm{K}[\mathrm{X}]=\mathbb{C}[X]$, but in general $\mathbb{C} \otimes_{\mathrm{K}} \mathrm{K}(X) \neq \mathbb{C}(X)$. The algebra $\mathbb{C} \otimes_{\mathrm{K}} \mathrm{K}(\mathrm{X})$ is the ring of fractions with numerators in $\mathbb{C}[X]$ and denominators in $K[X]$; its an integral ring and $\mathbb{C}(X)$ can be identified with its field of fractions.

### 7.1.3 Semi-definability

Proposition and Definition Let $(\mathrm{S}, \pi)$ be an analytic ramified covering of $\Sigma$ and K a subfield of $\mathbb{C}$. The following conditions are equivalent:
(i) ( $\mathrm{S}, \pi$ ) is definable over a finite extension L of K in $\mathbb{C}$;
( $i^{\prime}$ ) ( $\mathrm{S}, \pi$ ) is definable over the algebraic closure $\overline{\mathrm{K}}$ of K in $\mathbb{C}$;
(ii) $(\mathrm{S}, \pi)$ is isomorphic to a ramified covering (union of connected components) of an analytic ramified covering definable over K .

If these conditions hold, $(\mathrm{S}, \pi)$ will be said to be semi-definable over K .
Proof The equivalence (i) $\Leftrightarrow$ (i') is immediate.
(a) (i) $\Rightarrow$ (ii): Let A be an etale algebra over $\mathrm{L}(\mathrm{X})$, where L is a finite extension of K in $\mathbb{C}$, such that $\mathscr{M}(\mathrm{S})$ can be identified with $\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{L}(\mathrm{X})} \mathrm{A}$. Let $\eta$ be a primitive element of A and P its minimal polynomial, so that A can be identified with $\mathrm{L}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$. Let E be a Galois extension of K in $\mathbb{C}$ containing L , and G the Galois group of $\mathrm{Aut}_{\mathrm{K}}(\mathrm{E})$. As $G$ acts on the coefficients of the polynomial of $\mathrm{E}(\mathrm{X}, \mathrm{Y})$, it can be made to act on $\mathrm{E}(\mathrm{X}, \mathrm{Y})$ by leaving fixed the indeterminates. Let $g_{*} \mathrm{Q}$ denote the image of Q under $g$; note that if Q is monic, then so is $g_{*} \mathrm{Q}$.

Let $\widetilde{\mathrm{P}}$ be the l.c.m. of $g_{*} \mathrm{P}$ : it is the product of polynomials $g_{*} \mathrm{Q}$ with $g \in \mathrm{G}$ and Q an irreducible factor of P , taken without repetition.

Hence $\widetilde{P}=P \cdot P_{1}$ with $P_{1}$ and $P$ coprime. Then $\widetilde{A}_{E}=E(X)[Y] /(\widetilde{P})$ is isomorphic to $A_{E} \times \mathrm{A}_{\mathrm{E}}^{\prime}$, where $\mathrm{A}_{\mathrm{E}}=\mathrm{E}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})=\mathrm{E}(\mathrm{X}) \otimes_{\mathrm{L}(\mathrm{X})} \mathrm{A}$ and $\mathrm{A}_{\mathrm{E}}^{\prime}=\mathrm{E}(\mathrm{X})[\mathrm{Y}] /\left(\mathrm{P}_{1}\right)$. Set $\quad \widetilde{\mathrm{A}}_{\mathbb{C}}=\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \widetilde{\mathrm{A}}_{\mathrm{E}}, \quad \mathrm{A}_{\mathbb{C}}^{\prime}=\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}_{\mathrm{E}}^{\prime}, \widetilde{\mathrm{S}}=\widehat{\mathcal{S}}\left(\tilde{\mathrm{A}}_{\mathbb{C}}, \eta\right)$ 6.2.9, $\quad \mathrm{S}^{\prime}=$ $\widehat{\mathcal{S}}\left(\mathrm{A}_{\mathbb{C}}^{\prime}, \eta\right)$, with $\eta$ denoting the image of Y . Then $\widetilde{\mathrm{A}}_{\mathbb{C}}$ is isomorphic to $\mathrm{A}_{\mathbb{C}} \times \mathrm{A}_{\mathbb{C}}^{\prime}$ and $\tilde{\mathrm{S}}$ is a sum, i.e. a disjoint union of S and $\mathrm{S}^{\prime}$.

Moreover $\widetilde{\mathrm{P}} \in \mathrm{E}(\mathrm{X})[\mathrm{Y}]$ is left invariant by G ; hence $\widetilde{\mathrm{P}} \in \mathrm{K}(\mathrm{X})[\mathrm{Y}]$ and $\widetilde{\mathrm{S}}$ together with its projection onto $\Sigma$ is definable over K.
(b) (ii) $\Rightarrow$ (i): Assume that $S$ is a clopen subset of $\widetilde{S}$ definable over K. We may assume that $\widetilde{\mathrm{S}}=\widehat{\delta}(\mathrm{A}, \tilde{\eta})$, where $\mathrm{A}=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\widetilde{\mathrm{P}})$ and $\tilde{\eta}$ is the image of Y . Hence $\widetilde{\mathrm{S}}=\mathrm{S} \sqcup \mathrm{S}_{1}$.

As in (6.2.5), let $\Delta_{\widetilde{\mathrm{P}}}$ be the set of $x \in \mathbb{C}$ for which either one of the coefficients of $\widetilde{\mathrm{P}}(x, \mathrm{Y})$ is not defined or this polynomial is a multiple root; it is a finite set. For all $x \in \mathbb{C}-\Delta_{\widetilde{\mathrm{P}}}$,

$$
\widetilde{\mathrm{P}}(x, \mathrm{Y})=\prod_{y \in \widetilde{\mathrm{~S}}(x)}(\mathrm{Y}-y)=\left(\prod_{y \in \mathrm{~S}(x)}(\mathrm{Y}-y)\right) \cdot\left(\prod_{y \in \mathrm{~S}_{1}(x)}(\mathrm{Y}-y)\right)=\mathrm{P}(x) \cdot \mathrm{P}_{1}(x)
$$

where $\mathrm{S}(x)=\tilde{\eta}\left(\pi_{\mathrm{S}}^{-1}(x)\right)$, etc.
The coefficients of $\mathrm{P}(x)$ are the symmetric functions of $z_{i}$ in $\mathrm{S}(x)$; their dependence on $x$ is holomorphic on $\mathbb{C}-\Delta_{\widetilde{\mathrm{P}}}$, with growth of the order of $\mathrm{O}\left(\frac{1}{|x-a|^{k}}\right)$ in the neighbourhood of $a \in \delta_{\widetilde{\mathrm{P}}}$, and polynomial at infinity. Thanks to (6.2.2, Proposition), $\mathrm{P}, \mathrm{P}_{1} \in \mathbb{C}(\mathrm{X})[\mathrm{Y}]$. Then $\mathrm{S}=\widehat{\mathcal{S}}(\mathrm{A}, \eta)$, where $\mathrm{A}=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$ and $\eta$ is the image of Y.

We show that the coefficients of P are algebraic over K. Let $\Omega$ be an algebraic closure of $\mathbb{C}(X)$. Since $\mathrm{P}=\prod_{i \in \mathrm{I}}\left(\mathrm{Y}-\eta_{i}\right)$ with algebraic $\eta_{i}$ over $\mathrm{K}(\mathrm{X})$ and $\mathrm{I} \subset \widetilde{\mathrm{I}}$, $\widetilde{\mathrm{P}}=\prod_{i \in \widetilde{\mathrm{I}}}\left(\mathrm{Y}-\eta_{i}\right)$ in $\Omega[\mathrm{Y}]$ Hence the coefficients of P as a polynomial in Y are algebraic elements of $\mathbb{C}(X)$ over $K(X)$. The proposition then follows from the next lemma:

Lemma For all algebraic polynomials $f \in \mathbb{C}(\mathrm{X})$ over $\mathrm{K}(\mathrm{X})$, there is a finite extension L of K such that $f \in \mathrm{~L}(\mathrm{X})$.
Proof of the Lemma $f$ can be uniquely written as $f=\frac{\mathrm{Q}(\mathrm{X})}{\mathrm{Q}_{0}(\mathrm{X})}$ with Q and $\mathrm{Q}_{0}$ coprime in $\mathbb{C}[X]$ and $Q_{0}$ monic. Set $G=\operatorname{Aut}_{K}(\mathbb{C})$ and make $G$ act on $\mathbb{C}(X)$ by leaving $X$ fixed. The G-orbit of $f$ is finite, and so the same holds for those of $\mathrm{Q}, \mathrm{Q}_{0}$ and of each of their coefficients. Hence, each coefficients of $Q$ and $Q_{0}$ is algebraic over $K$, and L can be taken to be the extension of K in $\mathbb{C}$ generated by these coefficients.

### 7.1.4 Definability for Riemann Surfaces

Definition Let $S$ be a compact Riemann surface and $K$ a subfield of $\mathbb{C}$. $S$ is said to be definable (resp. semi-definable) over K if there is a function $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that the ramified covering $(\mathrm{S}, f)$ is definable (resp. semi-definable) over K .

We remind the reader that $f \in \mathscr{M}^{* *}(\mathrm{~S})$ means that there is no connected component of S on which $f$ is constant. A subalgebra A of $\mathscr{M}(\mathrm{S})$ is called a K -structure on S if there is a function $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that A is a K -structure on ( $\mathrm{S}, f$ ).

Proposition Let A be a K -subalgebra of $\mathscr{M}(\mathrm{S})$. Then, A is a K -structure on S if and only if
(i) A is a finite product of extension fields of K with transcendence degree 1 ;
(ii) $\lambda f \mapsto \lambda f$ extends to an isomorphism from the total ring offractions of $\mathbb{C} \otimes_{\mathrm{K}} \mathrm{A}$ onto $\mathscr{M}(\mathrm{S})$.

Proof Suppose that A is a K-structure for $(\mathrm{S}, \pi)$, where $\pi=f$. Then A is an etale algebra over $K(X)$, hence a finite product of extensions of $K$ with transcendence degree 1 . The algebra $\mathscr{M}(\mathrm{S})=\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}$ is the algebra of fractions with numerators in A and denominators in $\mathbb{C}(\mathrm{X})-\{0\}$. Hence it is contained in the total ring of fractions of A, itself contained in $\mathscr{M}$ (S). Hence A satisfies (i) and (ii).

Conversely, suppose that A satisfies (i); let $f \in \mathrm{~A}$ be transcendental over K such that A is a finite algebra over $\mathrm{K}(f)=\mathrm{K}(\mathrm{X})$. Let $g$ be a primitive element of $K(X)$-algebra $A$. Then $A$ can be identified with $K(X)[Y] /(P)$, where $P$ is the minimal polynomial of $g$. The algebra $\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$ is the ring of fractions with numerators in $\mathbb{C}[\mathrm{X}, \mathrm{Y}] /(\mathrm{Q})$ (where Q is obtained from P by eliminating denominators) and denominators in $\mathrm{K}[\mathrm{X}]-\{0\}$. Hence it is contained in the total ring of fractions of $\mathbb{C} \otimes_{K} A$. As it is an etale algebra over $K(X)$, it is a product of fields; it contains $\mathbb{C} \otimes_{K} A$, hence also its total ring of fractions. Thus the total ring of fractions of $\mathbb{C} \otimes_{\mathrm{K}} \mathrm{A}$ is must be $\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})=\mathscr{M}(\mathrm{S})$, and A is a K-structure on ( $\mathrm{S}, f$ ).

Remarks (1) For any homographic function $f \in \mathscr{M}(\Sigma)$, i.e. of the form $\frac{a \mathrm{X}+b}{c \mathrm{X}+d}$ with $a, b, c, d \in \mathbb{C}$ and $a d-b c \neq 0$, the algebra $\mathrm{K}(f)$ is a K -structure on $\Sigma$. Then, $\mathrm{K}(f)=\mathrm{K}(g)$ if and only if there exist $\alpha, \beta, \gamma, \delta \in \mathrm{K}$ such that $\alpha \delta-\beta \gamma \neq 0$ and $g=\frac{\alpha f+\beta}{\gamma f+\delta}$. If K is countable (for example $\overline{\mathbb{Q}}$ ), this gives uncountably many distinct K -structures on $\Sigma$. These algebras are however isomorphic.
(2) If $S$ is connected with genus $\geqslant 2$ (or more generally if each connected component of S is of genus $\geqslant 2$ ), and if K is algebraically closed, then there is at most one K-structure on $S$ (Exercise 7.1.3).

## Exercices 7.1 (Definability)

1.-Let $S$ be a compact Riemann surface.
(a) Show that endowing $S$ with an $\mathbb{R}$-structure is equivalent to defining an antiholomorphic homeomorphism $\tau: \mathrm{S} \rightarrow \mathrm{S}$ such that $\tau \circ \tau=\mathrm{I}_{\mathrm{S}}$.
(b) The real points of $S$ are then the fixed points of $\tau$. Give examples related to the same Riemann surface where the set $S(\mathbb{R})$ of real points in empty, and of others where it is not.
(c) Show that the number of $\mathbb{R}$-structures on S is bounded above by the cardinality of the automorphism group of S. Give examples with arbitrarily large numbers of $\mathbb{R}$-structures.
2.-Let S be a Riemann surface with a defining polynomial $\mathrm{Y}^{2}-f(\mathrm{X})$, where $f \in \mathbb{Q}[\mathrm{X}]$ is non constant.
(a) Show that for $a \in \mathbb{Q}^{*}$, the $\mathbb{Q}$-subalgebra $\mathrm{A}_{a}$ generated by X and $\sqrt{a}$. Y defines a $\mathbb{Q}$-structure on S .
(b) Show that this gives infinitely many distinct $\mathbb{Q}$-structures on S .
(c) Give examples where these distinct $\mathbb{Q}$-structures are isomorphic, and example where they are not.
3.-Let $S$ be a compact Riemann surface all of whose connected components have genus $\geqslant 2$ and K an algebraically closed subfield of $\mathbb{C}$ (for example the algebraic closure $\overline{\mathbb{Q}}$ of $\mathbb{Q}$ in $\mathbb{C}$ ). The aim is to show that there is at most one K -structure on S .

Let $A_{1}$ and $A_{2}$ be $K$-structures on $S$, and $Q_{1}, Q_{2} \in K[X, Y]$ defining polynomials of respective Y-degrees $d_{1}$ and $d_{2}$. The identity of $\mathscr{M}(\mathrm{S})$ gives an isomorphism $\phi$ from the total ring of fractions of $\mathbb{C} \otimes_{K} A_{1}$ onto that of $\mathbb{C} \otimes_{K} A_{2}$. The isomorphism $\phi$ is given by rational fractions $r_{i, \varepsilon}$ with $\varepsilon=0,1$, such that

$$
\phi(\mathrm{X})=\sum_{i=0}^{d_{2}-1} r_{i, 0}(\mathrm{X}) \cdot \mathrm{Y}^{i}, \quad \phi(\mathrm{Y})=\sum_{i=0}^{d_{2}-1} r_{i, 1}(\mathrm{X}) \cdot \mathrm{Y}^{i}
$$

Consider the algebraic Galois group $G=\operatorname{Aut}_{\mathrm{K}}(\mathbb{C})$. The action of G on the coefficients defines rational fractions ${ }^{g} r_{i, \varepsilon} \in \mathbb{C}(\mathrm{X})$ for $g \in \mathrm{G}$. Show that these ${ }^{g} r_{i, \varepsilon}$ define an automorphism of S. Using (6.9.17), show that the G-orbit of the coefficients of $r_{i, \varepsilon}$ is finite. Deduce that these are elements of K and that $\mathrm{A}_{1}=\mathrm{A}_{2}$.

### 7.2 Belyi's Theorem

### 7.2.1 Statement of the Theorem

For $f \in \mathscr{M}^{* *}(\mathrm{~S})$ (see 7.1.4), let $\Delta(f)$ be the set of critical values of $f$ considered an analytic map from S to the Riemann sphere $\Sigma$. A polynomial $f \in \mathbb{C}[\mathrm{X}]$ will be considered a map $\Sigma \rightarrow \Sigma$; in particular $\infty \in \Delta(f)$ if the degree of $f$ is $>1$.

The algebraic closure of a subfield K of $\mathbb{C}$ in $\mathbb{C}$ will be denoted $\overline{\mathrm{K}}$.
Definition A compact Riemann surface is said to be arithmetic if it is semi-definable over $\mathbb{Q}$.

The aim of this section is to prove the following theorem:
Theorem Let S be a compact Riemann surface. The following conditions are equivalent:
(i) S is arithmetic;
(ii) there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset \overline{\mathbb{Q}} \cup\{\infty\}$;
(iii) there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset\{0,1, \infty\}$.

The implication (iii) $\Rightarrow$ (ii) is trivial. The equivalence (i) $\Leftrightarrow$ (ii) is standard, but its detailed proof is long. The originality of the result due to Belyi lies in the implication (ii) $\Rightarrow$ (iii).

The equivalence (i) $\Leftrightarrow$ (ii) is a particular case of the more general statement:
Proposition Let S be a compact Riemann surface and K a subfield of $\mathbb{C}$. The following conditions are equivalent:
(i) S is semi-definable over K ;
(ii) there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset \overline{\mathrm{K}} \cup\{\infty\}$.

In particular, if K is algebraically closed, for example if it is $\overline{\mathbb{Q}}$, then S is definable over K if and only if there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset \mathrm{K} \cup\{\infty\}$.

Implication (i) $\Rightarrow$ (ii) will be proved in 7.2.2. Implication (ii) $\Rightarrow$ (i) will be proved in 7.2.5 using a density result stated in 7.2.3 as well as classic results of algebraic geometry. We will also give a method to obtain the number of distinct roots of a polynomial (7.2.4).

### 7.2.2 Implication $(i) \Rightarrow$ (ii)

Suppose that $\mathrm{S}=\widehat{\mathrm{V}_{\mathrm{Q}}}$ (see 7.1.1), $\mathrm{Q} \in \mathrm{L}[\mathrm{X}, \mathrm{Y}]$, is regular of degree $d$ in Y , where L is a finite extension of K . Then $\mathscr{M}(\mathrm{S})=\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$, where P is the polynomial Q made monic in $\mathbb{C}(\mathrm{X})[\mathrm{Y}]$. Set $\mathrm{A}=\mathrm{L}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$ and let $f$ be an arbitrary element of $\mathrm{A}^{* *}=\mathrm{A} \cap \mathscr{M}^{* *}(\mathrm{~S})$. We show that $\Delta(f) \subset \overline{\mathrm{L}} \cup\{\infty\}=\overline{\mathrm{K}} \cup\{\infty\}$.

Let $c \in \mathrm{~S}$ be a critical point of $f$ and set $v=f(c)$. The function $f$ is representable by a unique polynomial $\mathrm{F} \in \mathrm{L}(\mathrm{X})[\mathrm{Y}]$ of degree $<d$ :

$$
\mathrm{F}=f_{0}+\cdots+f_{d-1} \mathrm{Y}^{d-1}
$$

1st case: $c$ is a point $\left(x_{c}, y_{c}\right)$ of $\mathrm{V}_{\mathrm{Q}}^{\prime}$ and $x_{c}$ is not a pole of any $f_{i}$.
The derivative $g=\mathrm{D}_{\mathrm{S}} f$ de $f$ with respect to X along S is given by:

$$
g=\frac{\partial \mathrm{F}}{\partial \mathrm{X}}-\frac{\partial \mathrm{Q} / \partial \mathrm{X}}{\partial \mathrm{Q} / \partial \mathrm{Y}} \cdot \frac{\partial \mathrm{~F}}{\partial \mathrm{Y}}
$$

Hence it is in A.
The resultant $h=\operatorname{Res}(g, \mathrm{P})$ in is $\mathrm{L}(\mathrm{X})$. As $h(x)=\prod_{y \in \mathrm{~S}(x)} g(x, y)$ for $x \in \mathbb{C}-\Delta_{\mathrm{Q}}$, and $h\left(x_{c}\right)=0, x_{c}$ is algebraic over L . Since $\mathrm{Q}\left(x_{c}, y_{c}\right)=0$, the ordinate $y_{c}$ is also algebraic over L , and the same holds for $v=\widehat{\mathrm{F}}\left(x_{c}\right)$.

2nd case: The image of $c$ under the map $(\mathrm{X}, \mathrm{Y}): \widehat{\mathrm{V}_{\mathrm{Q}}} \rightarrow \Sigma^{2}$ is a point $\left(x_{c}, y_{c}\right)$ with $x_{c} \in \Delta_{\mathrm{Q}}$, where $x_{c}$ is a pole of some $f_{i}$.

In this case, $x_{c}$ is obviously algebraic over L , and S can be analytically parametrized in the neighbourhood of $c$ by T , where $\mathrm{T}^{k}=\mathrm{X}-x_{c}$. Then Y can be expressed
as a series in T , with coefficients following from an algebraic calculation. Hence $\mathrm{Y} \in \overline{\mathrm{L}}[[\mathrm{T}]]$. The function $f=\mathrm{F}(\mathrm{X}, \mathrm{Y})$ is represented by an element of $\overline{\mathrm{L}}((\mathrm{T}))$, i.e. by a fraction with numerator in $\overline{\mathrm{L}}[[\mathrm{T}]]$ and denominator a power of T . And $v=f(c)$, obtained by setting $\mathrm{T}=0$, is in $\overline{\mathrm{L}} \cup\{\infty\}$.

This completes the proof of implication (i) $\Rightarrow$ (ii) in (7.2.1, proposition).

### 7.2.3 Density

Proposition Let K be a subfield of $\mathbb{C}$, and $\overline{\mathrm{K}}$ its algebraic closure in $\mathbb{C}$. Let $\mathbf{f}=$ $\left(f_{1}, \ldots, f_{p}\right) \in \mathrm{K}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]^{p}$, and M the set of points $z=\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n}$ such that $f_{1}(z)=\cdots=f_{p}(z)=0$. Then $\mathrm{M} \cap \overline{\mathrm{K}}^{n}$ is dense in M .

Proof The set M is algebraic in $\mathbb{C}^{n}$. The set $\mathrm{M}^{*}$ of regular points of M is open and dense in $\mathrm{M}[1,2]$. Let $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathrm{M}^{*}$, and $d$ be the dimension of $\mathrm{M}^{*}$ at $x$. There is a partition of $\{1, \ldots, n\}$ in $\mathrm{J}^{\prime}$ and $\mathrm{J}^{\prime \prime}$ with $\operatorname{Card}\left(\mathrm{J}^{\prime}\right)=d$ such that, identifying $\mathbb{C}^{n}$ with $\mathbb{C}^{\mathrm{J}^{\prime}} \times \mathbb{C}^{\mathrm{J}^{\prime \prime}}$, the vector space $\mathrm{T}_{x} \mathrm{M}^{*}$ is the graph of a linear map $\mathbb{C}^{\mathrm{J}^{\prime}} \rightarrow \mathbb{C}^{\mathrm{J}^{\prime \prime}}$.

We then write $x=\left(x^{\prime}, x^{\prime \prime}\right)$. By the implicit function theorem, there are neighbourhoods $\mathrm{U}^{\prime}$ of $x^{\prime}$ and $\mathrm{U}^{\prime \prime}$ of $x^{\prime \prime}$ such that $\mathrm{M}^{*} \cap\left(\mathrm{U}^{\prime} \times \mathrm{U}^{\prime \prime}\right)$ is the graph of a $\mathbb{C}$-analytic map $h: \mathrm{U}^{\prime} \rightarrow \mathrm{U}^{\prime \prime}$. This map is algebraic over $\mathrm{K}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{~J}^{\prime}}\right]$, in the following sense: for all $j \in \mathbf{J}^{\prime \prime}$, the map $h_{j}: \mathrm{U}^{\prime} \rightarrow \mathbb{C}$ satisfies a relation

$$
a_{j, m_{j}} h_{j}^{m_{j}}+\cdots+a_{j, 0}=0
$$

where all $a_{j, k}$ are in $\mathrm{K}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{~J}^{\prime}}\right]$ and $a_{j, m_{j}} \neq 0$. Hence $y^{\prime} \in \mathrm{U}^{\prime} \cup \overline{\mathrm{K}}^{\mathbf{J}^{\prime}}$ such that $a_{j, m_{j}}\left(y^{\prime}\right) \neq 0$ for all $j$, the coordinates of the point $y^{\prime \prime}=h\left(y^{\prime}\right)$ are in $\overline{\mathrm{K}}$, and $y=\left(y^{\prime}, y^{\prime \prime}\right) \in \mathbf{M}^{*} \cap\left(\mathrm{U}^{\prime} \times \mathrm{U}^{\prime \prime}\right) \cap \overline{\mathrm{K}}^{n}$. Therefore, $\mathbf{M}^{*} \cap\left(\mathrm{U}^{\prime} \times \mathrm{U}^{\prime \prime}\right) \cap \overline{\mathrm{K}}^{n}$ is dense in $\mathbf{M}^{*} \cap\left(\mathrm{U}^{\prime} \times \mathrm{U}^{\prime \prime}\right)$. It follows that $\mathrm{M}^{*} \cap \overline{\mathrm{~K}}^{n}$ is dense in $\mathbf{M}^{*}$, hence so in $\mathbf{M}$.

### 7.2.4 Polynomials with $k$ Distinct Roots

Proposition Let $\mathrm{Q} \in \mathbb{C}[\mathrm{X}]$ be a polynomial of degree d, with $k$ distinct roots. Then the function $h_{\mathrm{Q}}: \mathbb{C} \rightarrow \mathbb{C}$ assigning to the discriminant $\operatorname{discr}(\mathrm{Q}-t)$ has a zero of order $d-k$ at 0 .

Proof Let $\alpha_{1}, \ldots, \alpha_{k}$ be the distinct roots of Q with multiplicities $\mu_{1}, \ldots, \mu_{k}$. For $t$ near 0 , the polynomial $\mathrm{Q}-t$ has $\mu_{i}$ roots $\xi_{i, 1}, \ldots, \xi_{i, \mu_{i}}$ near $\alpha_{i}$, with $\xi_{i, j}-\alpha_{i}$ infinitely small of the order of $|t|^{1 / \mu_{i}}$ as $t$ tends to 0 , as well as $\xi_{i, j}-\xi_{i, j^{\prime}}$ for $j \neq j^{\prime}$. Indeed, by proposition 6.1.4, there is a chart $\varphi$ de $\mathbb{C}$ in the neighbourhood of $\alpha_{i}$ such that $\mathrm{Q}(z)=\varphi(z)^{\mu_{i}}$ for $z$ in the domain of $\varphi$; then let $\xi_{i, j}$ be $\varphi^{-1}\left(s_{j}\right)$, where the $s_{j}$ are the $\mu_{i}$ th roots of $t$. The product $\prod_{j \neq j^{\prime}}\left(\xi_{i, j}-\xi_{i, j^{\prime}}\right)$ is of the order of $\left(|t|^{1 / \mu_{i}}\right)^{\mu_{i}^{2}-\mu_{i}}=|t|^{\mu_{i}-1}$;
and $\xi_{i, j}-\xi_{i^{\prime}, j^{\prime}}$ is of the order of 1 for $i \neq i^{\prime}$. Hence

$$
\operatorname{discr}(\mathrm{Q}-t)=\prod_{(i, j) \neq\left(i^{\prime}, j^{\prime}\right)}\left(\xi_{i, j}-\xi_{i^{\prime}, j^{\prime}}\right)
$$

is of the order of $|t|^{\nu}$, where $\nu=\sum \mu_{i}-1=d-k$.
Corollary In the space of polynomials $\mathrm{Q} \in \mathbb{C}[\mathrm{X}]$ of degree $\leqslant d$, identifiable with $\mathbb{C}^{d+1}$, the set of polynomials of degree $<d$ or with at most $k$ distinct roots is a $\mathbb{Q}$-algebraic set.

Proof This is the set on which polynomial functions $\mathrm{Q} \mapsto h_{\mathrm{Q}}^{(r)}(0)$ with $0 \leqslant r<$ $d-k$ over $\mathbb{Z}$ vanish.

### 7.2.5 Implication (ii) $\Rightarrow$ ( $i$ )

Proposition Let S be a Riemann surface, K an algebraically closed subfield of $\mathbb{C}$ and $f \in \mathscr{M}^{* *}(\mathrm{~S})$. Suppose that $\Delta(f) \subset \mathrm{K} \cup\{\infty\}$. Then S is definable over K .

Proof Let $\mathrm{P}(\mathrm{X}, \mathrm{Y})=a_{d}(\mathrm{X}) \mathrm{Y}^{d}+\cdots+a_{0}(\mathrm{X})$ be a defining polynomial of $(\mathrm{S}, f)$. With the notation of 7.1.1,

$$
\widehat{\Delta}_{\mathrm{P}}=\Delta_{\mathrm{P}} \cup\{\infty\} \supset \Delta(f)
$$

Let $\xi_{1}, \ldots, \xi_{s}$ be the points of $\widehat{\Delta}_{\mathrm{P}}$, with $\Delta(f)=\left\{\xi_{1}, \ldots, \xi_{r}\right\}$. For $1 \leqslant i \leqslant s$, let $\mathrm{A}_{i}$ be disjoint closed discs centered at $\xi_{i}$ if $\xi_{i} \in \mathbb{C}$, of the form $\Sigma-\mathbb{D}_{\mathrm{R}}$ if $\xi_{i}=\infty$. Set $\mathrm{A}_{i}^{*}=\mathrm{A}_{i}-\left\{\xi_{i}\right\}$ and let $\Gamma_{i}$ be the boundary of $\mathrm{A}_{i}$.

For $1 \leqslant i \leqslant r$, the ramified covering ( $\mathrm{S}, f$ ) induces a non trivial covering of $\Gamma_{i}$; with the notation of 6.2 .5 , for $r<i \leqslant s$, the following holds:

- if $\xi_{i} \in \Delta^{\prime}$, then one of $\frac{a_{i}}{a_{d}}$ is not defined; so $a_{d}\left(\xi_{i}\right)=0$, and the point $a_{d}(x) \in \mathbb{C}^{*}$ goes round $\nu_{i} \neq 0$ times as $x$ runs through $\Gamma_{i}$;
- if $\xi_{i}=\infty$, one of $a_{j}(x)$ goes round a non-trivial number of times;
- if $\xi_{i} \in \Delta^{\prime \prime}-\Delta^{\prime}$, then the set $\pi^{-1}\left(\mathrm{~A}_{i}\right) \cap \mathrm{V}_{\mathrm{P}}$ is the union of graphs of $d$ holomorphic functions $\eta_{1}, \ldots, \eta_{d}: \mathrm{A}_{i} \rightarrow \mathbb{C}$, and there are two indices $j \neq j^{\prime}$ such that $\eta_{j^{\prime}}(x)-\eta_{j}(x)$ goes round $\nu_{i} \neq 0$ times as $x$ runs through $\Gamma_{i}$.

For $0 \leqslant i \leqslant d$, let $e_{i}$ be the degree of $a_{i} \in \mathbb{C}[\mathrm{X}], e=e_{d}$, and E the set of polynomials $\mathrm{Q}=b_{d}(\mathrm{X}) \mathrm{Y}^{d}+\cdots+b_{0}(\mathrm{X}) \in \mathbb{C}[\mathrm{X}, \mathrm{Y}]$ such that $b_{i}$ is of degree $\leqslant e_{i}$ for all $i$, and $b_{d}$ is monic of degree $e$. This set can be identified with $\mathbb{C}^{\mathrm{N}}$, where $\mathrm{N}=\sum\left(e_{i}+1\right)-1$. Let $\mathrm{E}_{\mathrm{K}}$ be the set of $\mathrm{Q} \in \mathrm{E}$ with coefficients in K .

Let $M$ be the set of $\mathrm{Q} \in \mathrm{E}$ such that $\widehat{\Delta}_{\mathrm{Q}} \supset \Delta(f)$ and $\operatorname{Card}\left(\Delta_{\mathrm{Q}}\right) \leqslant \operatorname{Card}\left(\Delta_{\mathrm{P}}\right)$.
Lemma 7.1 For $\mathrm{Q} \in \mathrm{M}$ sufficiently near P , as a ramified covering of $\Sigma$, the Riemann surface $\widehat{\mathrm{V}_{\mathrm{Q}}}$ is isomorphic to $\widehat{\mathrm{V}_{\mathrm{P}}}$, hence to S equipped with $f$.

Proof The set $\mathrm{B}=\Sigma-\bigcup \stackrel{\circ}{\mathrm{A}}_{i}$ is compact, and $\mathrm{X} \mapsto \operatorname{discrP}(\mathrm{X})$ as well as $a_{d}$ de not vanish on $B$. Let $\Omega$ be a connected relatively compact neighbourhood of P in E such that for $\mathrm{Q} \in \Omega$, the functions $\mathrm{X} \mapsto \operatorname{discr} \mathrm{Q}(\mathrm{X})$ and $b_{d}$ do not vanish on B . Then, by (4.3.12, Proposition),

$$
\mathrm{Z}=\{(\mathrm{Q}, x, y) \in \Omega \times \mathrm{B} \times \mathbb{C} \mid \mathrm{Q}(x)(y)=0\}
$$

together with the projection $\pi:(\mathrm{Q}, x, y) \mapsto(\mathrm{Q}, x)$ is a covering of $\Omega \times \mathrm{B}$. Indeed, $\pi$ is etale by the implicit function theorem, and proper since the functions $b_{i}$ are bounded and $b_{d}$ is bounded below by a strictly positive constant. The injections $x \mapsto(\mathrm{P}, x)$ and $x \mapsto(\mathrm{Q}, x)$ are homotopic; so $\mathrm{V}_{\mathrm{P}}$ and $\mathrm{V}_{\mathrm{Q}}$ induce coverings isomorphic to B (4.3.17, Corollary 1).

In particular, for all $i$, the coverings of $\Gamma_{i}$ induced by $\mathrm{V}_{\mathrm{P}}$ and $\mathrm{V}_{\mathrm{Q}}$ are isomorphic. Each set $\mathrm{A}_{i}$ has at least one point of $\widehat{\Delta}_{\mathrm{Q}}$. Indeed, for $1 \leqslant i \leqslant r, \mathrm{~V}_{\mathrm{Q}}$ induces a non trivial covering of $\Gamma_{i}$; for $i$ such that $\xi_{i} \in \Delta^{\prime} \cup\{\infty\}, b_{i}(x)$ stays near $a_{i}(x)$ as $x$ runs through $\Gamma_{i}$, and goes round $0 \nu_{i}$ times; for $i$ such that $\xi_{i} \in \Delta^{\prime \prime}-\Delta^{\prime}$, there are functions $\eta_{1}^{\mathrm{Q}}, \ldots, \eta_{d}^{\mathrm{Q}}$ near $\eta_{1}, \ldots, \eta_{d}$, where $\eta_{j^{\prime}}^{\mathrm{Q}}(x)-\eta_{j}(x)$ goes round $\nu_{i}$ times as $x$ runs through $\Gamma_{i}$.

As $\operatorname{Card}\left(\widehat{\Delta}_{\mathrm{Q}}\right) \leqslant \operatorname{Card}\left(\widehat{\Delta}_{\mathrm{P}}\right)$, each set $\mathrm{A}_{i}$ has exactly one point $\xi_{i}^{\mathrm{Q}}$ of $\widehat{\Delta}_{\mathrm{Q}}$, and the covering of $\overline{\mathrm{A}}_{i}-\left\{\xi_{i}^{\mathrm{Q}}\right\}$ induced by $\mathrm{V}_{\mathrm{Q}}$ is determined by its restriction to $\Gamma_{i}$.

For $r<i \leqslant s$, this covering is trivial; Hence $\mathrm{V}_{\mathrm{Q}}$ induces an (unramified) covering of $\Sigma-\bigcup_{i=1}^{r} \mathrm{~A}_{i}$, and an isomorphism between the coverings induced by $\widehat{\mathrm{V}_{\mathrm{P}}}$ and $\widehat{\mathrm{V}_{\mathrm{Q}}}$ over $\Sigma-\bigcup \mathrm{A}_{i}$ extends to coverings induced over $\Sigma-\bigcup_{i=1}^{r} \mathrm{~A}_{i}$.

For $1 \leqslant i \leqslant r, \xi_{i}^{\mathrm{Q}}=\xi_{i}$ and an isomorphism between the coverings induced by $\widehat{\mathrm{V}_{\mathrm{P}}}$ and $\widehat{\mathrm{V}_{\mathrm{Q}}}$ over $\Sigma-\bigcup_{i=1}^{r} \mathrm{~A}_{i}$ extends to an isomorphism of ramified coverings of $\Sigma$ between $\widehat{V_{P}}$ and $\widehat{V_{Q}}$,

Lemma 7.2 The set $\mathrm{M} \cap \mathrm{E}_{\mathrm{K}}$ is dense in M .
Proof The set M is K-algebraic. Indeed, the coefficients of the discriminant $\delta_{\mathrm{Q}}=$ $\operatorname{discr}_{\mathrm{Y}} \mathrm{Q} \in \mathbb{C}[\mathrm{X}]$ are polynomials whose coefficients are integral with respect to the coefficients of Q . For each $\xi \in \Delta(f)-\{\infty\} \subset \mathrm{K}$, the set of $\mathrm{Q}=b_{d} \mathrm{Y}^{d}+\cdots+b_{0} \in$ E such that $\left(b_{d} \delta_{\mathrm{Q}}\right)(\xi)=0$ is K-algebraic. As a result of (7.2.4, Corollary) the set of $\mathrm{Q} \in \mathrm{E}$ such that $\operatorname{Card}\left(\Delta_{\mathrm{Q}}\right) \leqslant \operatorname{Card}\left(\Delta_{\mathrm{P}}\right)$ is $\mathbb{Q}$-algebraic. The intersection M of all these sets is therefore a K -algebraic subset of E .

The lemma then follows from (7.2.3, Proposition).
End of the Proof of the Proposition With the above notation, S is isomorphic to $\widehat{\mathrm{V}_{\mathrm{P}}}$, hence also to $\widehat{\mathrm{V}_{\mathrm{Q}}}$ for $\mathrm{Q} \in \mathrm{M}$ sufficiently near P (Lemma 7.1). By Lemma 7.2, $\mathrm{Q} \in \mathrm{M} \cap \mathrm{E}_{\mathrm{K}}$ may be chosen sufficiently near P for this to hold. Then $\widehat{\mathrm{V}_{\mathrm{Q}}}$ is defined on K .

This completes the proof of (7.2.1, Proposition).

### 7.2.6 The Action $f$ :

Let S and $\mathrm{S}^{\prime}$ be compact Riemann surfaces and $f: \mathrm{S} \rightarrow \mathrm{S}^{\prime}$ an analytic map. For any finite subset A of S, set

$$
f_{!}(\mathrm{A})=f(\mathrm{~A}) \cup \Delta(f)
$$

where $\Delta(f)$ is the set of critical values of $f$. Then

$$
\begin{gathered}
f_{!}(\varnothing)=\Delta(f), \\
f_{!}(\mathrm{A} \cup \mathrm{~B})=f_{!}(\mathrm{A}) \cup f_{!}(\mathrm{B})=f_{!}(\mathrm{A}) \cup f(\mathrm{~B}), \\
(g \circ f)!(\mathrm{A})=g_{!}\left(f_{!}(\mathrm{A})\right) .
\end{gathered}
$$

In particular,

$$
\Delta(g \circ f)=g_{!}(\Delta(f))
$$

We will apply this to the case where $\mathrm{S}=\mathrm{S}^{\prime}=\Sigma$ and $f$ is a polynomial of degree $d \geqslant 2$. Then $\Delta(f)=\Delta^{*}(f) \cup\{\infty\}$, where $\Delta^{*}(f)$ is the set of critical values of $f: \mathbb{C} \rightarrow \mathbb{C}$.

### 7.2.7 Rational Critical Values

Proposition For a compact Riemann surface of S, (ii) $\Rightarrow$ (ii'), where:
(ii) there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset \overline{\mathbb{Q}} \cup\{\infty\}$,
(ii') there exists $f \in \mathscr{M}^{* *}(\mathrm{~S})$ such that $\Delta(f) \subset \mathbb{Q} \cup\{\infty\}$.
Lemma Let $\mathrm{A} \subset \mathbb{Q} \cup\{\infty\} \cup \mathrm{P}^{-1}(0)$ be a finite set, with $\mathrm{P} \in \mathbb{Q}[\mathrm{X}]$ of degree $d$. Then there is a polynomial $h$ such that $h_{!}(\mathrm{A}) \subset \mathbb{Q} \cup\{\infty\} \cup \mathrm{P}_{1}^{-1}(0)$ with $\mathrm{P}_{1} \in \mathbb{Q}[\mathrm{X}]$ of degree $\leqslant d-1$.

Proof In fact, $h=\mathrm{P}$ is suitable. Set $\mathrm{A}=\mathrm{A}^{\prime} \cup \mathrm{A}^{\prime \prime}$, where $\mathrm{A}^{\prime} \subset \mathrm{P}^{-1}(0)$ and $\mathrm{A}^{\prime \prime} \subset$ $\mathbb{Q} \cup\{\infty\}$. Then, $\mathrm{P}\left(\mathrm{A}^{\prime}\right) \subset\{0\}$ and $\mathrm{P}\left(\mathrm{A}^{\prime \prime}\right) \subset \mathbb{Q} \cup\{\infty\}$. The set $\Delta^{*}(\mathrm{P})$ of the finite critical values of P has at most $d-1$ elements, and is invariant under $\mathrm{Aut}_{\mathbb{Q}}(\mathbb{C})$, and equals $\mathrm{P}_{1}^{-1}(0)$, for some $\mathrm{P}_{1} \in \mathbb{Q}[\mathrm{X}]$ of degree $\leqslant d-1$. And
$\mathrm{P}_{!}(\mathrm{A})=\mathrm{P}(\mathrm{A}) \cup \Delta(\mathrm{P}) \subset \mathbb{Q} \cup\{\infty\} \cup \mathrm{P}_{1}^{-1}(0)$.
A repeated application of this lemma gives:
Corollary Let $\mathrm{A} \subset \overline{\mathbb{Q}} \cup\{\infty\}$ be a finite set. Then there is a polynomial $h \in \mathbb{Q}[\mathrm{X}]$ such that $h_{!}(\mathrm{A}) \subset \mathbb{Q} \cup\{\infty\}$.

Proof of the Proposition Let $f$ be such that $\Delta(f) \subset \overline{\mathbb{Q}} \cup\{\infty\}$. By the above corollary, there exists $h \in \mathbb{Q}[\mathrm{X}]$ such that $\Delta(h \circ f)=h_{!}(\Delta(f)) \subset \mathbb{Q} \cup\{\infty\}$.

### 7.2.8 Reduction of the Number of Critical Values

Proposition Let $\mathrm{A} \subset \mathbb{Q} \cup\{\infty\}$ be a finite set such that $\infty \in \mathrm{A}$ and $\operatorname{Card}(\mathrm{A})>3$. There is a polynomial $h \in \mathbb{Q}[\mathrm{X}]$ for which Card $\left(h_{!}(\mathrm{A})\right)<\operatorname{Card}(\mathrm{A})$.

Lemma For $f=f_{p, q}=\mathrm{X}^{p}(1-\mathrm{X})^{q}$ and $\mathrm{A}=\left\{0, \frac{p}{p+q}, 1, \infty\right\}$, $\operatorname{Card}\left(f_{!}(\mathrm{A})\right)=3$.
Proof The critical points of $f$ in $\mathbb{C}$ are 0,1 and $\frac{p}{p+q}$; hence $\Delta^{*}(f)=\left\{0, f\left(\frac{p}{p+q}\right)\right\}$, and $f_{!}(\mathrm{A})=f(\mathrm{~A})=\left\{0, f\left(\frac{p}{p+q}\right), \infty\right\}$.

Proof of the Proposition Let $\mathrm{A}=\left\{a_{1}, a_{2}, \ldots, a_{n-1}, \infty\right\}$ with $a_{1}<a_{2}<a_{3}$. Let $\alpha$ be the affine map $\mathbb{C} \rightarrow \mathbb{C}$ such that $\alpha\left(a_{1}\right)=0$ and $\alpha\left(a_{3}\right)=1$. Then $\left.\alpha\left(a_{2}\right) \in \mathbb{Q} \cap\right] 0,1[$ is of the form $\frac{p}{p+q}$. Take $f=\mathrm{X}^{p}(1-\mathrm{X})^{q}$ and $h=f \circ \alpha$. Then $h_{!}\left(\left\{a_{1}, a_{2}, a_{3}, \infty\right\}\right)=$ $f_{!}\left(\left\{0, \frac{p}{p+q}, \infty\right\}\right)$, and $\left.h_{!}(\mathrm{A})=h_{!}\left(\left\{a_{1}, a_{2}, a_{3}, \infty\right\}\right) \cup h\left(a_{4}, \ldots, a_{n-1}\right\}\right)$ has at most $3+n-4=n-1$ elements.

Corollary 7.1 For any finite set $\mathrm{A} \subset \mathbb{Q} \cup\{\infty\}$, there is a polynomial $h \in \mathbb{Q}[\mathrm{X}]$ such that $h_{!}(\mathrm{A}) \subset\{0,1, \infty\}$.

Proof Repeated application of the proposition gives a polynomial $H \in \mathbb{Q}[X]$ such that $\operatorname{Card}\left(\mathrm{H}_{!}(\mathrm{A})\right) \leqslant 3$. This implies $\infty \in \mathrm{H}_{!}(\mathrm{A})$, and so $\mathrm{H}_{!}(\mathrm{A}) \subset\left\{b_{1}, b_{2}, \infty\right\}$ with $b_{1}<b_{2}$. Let $\beta$ be affine and such that $\beta\left(b_{1}\right)=0$ and $\beta\left(b_{2}\right)=1$. For $h=\beta \circ \mathbf{H}$, $h_{!}(\mathrm{A}) \subset\{0,1, \infty\}$.

Corollary 7.2 For any finite set $\mathrm{A} \subset \overline{\mathbb{Q}} \cup\{\infty\}$, there is a polynomial $h \in \mathbb{Q}[\mathrm{X}]$ such that $h_{!}(\mathrm{A}) \subset\{0,1, \infty\}$.

Proof The lemma follows from Corollary 7.1 and (7.2.7 Corollary).
End of the Proof of Belyi's Theorem We keep the notation of (7.2.1) and (7.2.7). It remains to prove that (ii') $\Rightarrow$ (iii). If $g \in \mathscr{M}^{* *}(\mathbf{S})$ such that $\Delta(g) \subset \mathbb{Q} \cup\{\infty\}$, then by Corollary 7.1, there exists $h \in \mathbb{Q}[\mathrm{X}]$ for which $h_{!}(\Delta(g)) \subset\{0,1, \infty\}$. Then $\Delta(h \circ g)=h_{!}(\Delta(g)) \subset\{0,1, \infty\}$.

### 7.3 Equivalence Between Various Categories

In this section, we define equivalences between several categories seemingly from different domains: algebra, topology, combinatorics, complex analysis. This should establish correspondences between these different domains, thereby contributing to their enrichment.

### 7.3.1 Ramification Locus of an Algebra

Let $K$ be a subfield of $\mathbb{C}$ and $A$ an etale $K(X)$-algebra. Denote by $A_{\mathbb{C}}$ the $\mathbb{C}(X)$-algebra $\mathbb{C}(X) \otimes_{K(X)} A$. With the notation of (6.2.11), set $\Delta_{A}=\Delta_{A_{C}}$ : hence it is the ramification set of an analytic ramified covering ( $\mathrm{S}, \pi$ ) of $\Sigma$ such that as a $\mathbb{C}(\mathrm{X})$-algebra, $\mathscr{M}(\mathrm{S})$ is isomorphic to $\mathrm{A}_{\mathbb{C}}$.

### 7.3.2 Transfer of $\Delta$

Let K be a field of characteristic 0 and A an etale $\mathrm{K}(\mathrm{X})$-algebra; set $\alpha$ to be the ring homomorphism $\mathrm{K}(\mathrm{X}) \rightarrow$ A defining the algebra structure of A . Let $\iota_{1}$ and $\iota_{2}$ be two embeddings of K into $\mathbb{C}$. For $i=1,2$, set $\mathrm{K}_{i}=\iota_{i}(\mathrm{~K})$. Extend $\iota_{i}$ to an embedding $\mathrm{K}(\mathrm{X}) \rightarrow \mathbb{C}(\mathrm{X})$ by $\iota_{i}(\mathrm{X})=\mathrm{X}$. Let $\mathrm{A}_{i}$ be the $\mathrm{K}_{i}(\mathrm{X})$-algebra obtained by equipping the ring A of $\alpha_{i}=\alpha \circ \iota_{i}^{-1}: \mathrm{K}_{i} \rightarrow \mathrm{~A}$, and $\Delta_{i}$ the subset $\Delta_{\mathrm{A}_{i}}$ of $\mathbb{C}$.
Proposition Set $\sigma=\iota_{2} \circ \iota_{1}^{-1}: \mathrm{K}_{1} \rightarrow \mathrm{~K}_{2}$, and let $\tau$ be an algebraic automorphism of $\mathbb{C}$ inducing $\sigma$. Then $\tau\left(\Delta_{1}\right)=\Delta_{2}$.

Proof Extend $\sigma$ to $\mathrm{K}_{1}(\mathrm{X})$ and $\tau$ to $\mathbb{C}(\mathrm{X})[\mathrm{Y}]$ by $\sigma(\mathrm{X})=\mathrm{X}, \tau(\mathrm{X})=\mathrm{X}$ and $\tau(\mathrm{Y})=\mathrm{Y}$. Then $\tau$ is a $\sigma$-morphism. Hence there is a morphism

$$
\tau_{*}=\tau \otimes 1_{\mathrm{A}}:\left(\mathrm{A}_{1}\right)_{\mathbb{C}}=\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}_{1}(\mathrm{X})} \mathrm{A}_{1} \rightarrow \mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}_{2}(\mathrm{X})} \mathrm{A}_{2}=\left(\mathrm{A}_{2}\right)_{\mathbb{C}}
$$

If $\eta_{1}$ is a primitive element in $\left(\mathrm{A}_{1}\right)_{\mathbb{C}}$, then $\eta_{2}=\tau_{*}\left(\eta_{1}\right)$ is primitive in $\left(\mathrm{A}_{2}\right)_{\mathbb{C}}$, and the automorphism $\tau_{*}$ of $\mathbb{C}(\mathrm{X})[\mathrm{Y}]$ inducing $\tau$ on $\mathbb{C}$ and fixing X and Y maps the minimal polynomial $\mathrm{P}_{1}$ of $\eta_{1}$ onto the minimal polynomial $\mathrm{P}_{2}$ of $\eta_{2}$. The definition of $\Delta_{\mathrm{E}, \eta}$ in (6.2.5) being of an algebraic nature,

$$
\Delta_{\left(\mathrm{A}_{2}\right) \mathrm{c}, \eta_{2}}=\tau\left(\Delta_{\left(\mathrm{A}_{1}\right) \mathrm{c}, \eta_{1}}\right)
$$

Thanks to (6.2.12, Proposition), taking the intersection over all primitive elements gives

$$
\Delta_{\left(\mathrm{A}_{2}\right) \mathrm{C}}=\tau\left(\Delta_{\left(\mathrm{A}_{1}\right) \mathrm{C}}\right)
$$

> cqfd

### 7.3.3

Corollary Let K be a subfield of $\mathbb{C}$ and A an etale $\mathrm{K}(\mathrm{X})$-algebra. Then the set $\Delta_{\mathrm{A}}$ is contained in $\overline{\mathrm{K}}$ and generates a Galois extension of K .
Proof For any algebraic automorphism $\tau$ of $\mathbb{C}$ inducing the identity on $\mathrm{K}, \tau\left(\Delta_{\mathrm{A}}\right)=$ $\Delta_{\mathrm{A}}$. As $\Delta_{\mathrm{A}}$ is finite, the elements of $\Delta_{\mathrm{A}}$ only have finitely many conjugates. Hence they are algebraic over K . The extension generated by $\Delta_{\mathrm{A}}$ is invariant under Aut $_{\mathrm{K}}(\overline{\mathrm{K}})$, and so is Galois.

### 7.3.4 The Category $\mathscr{A}$

Let $\mathscr{A}$ denote the category of etale algebras A over $\overline{\mathbb{Q}}(\mathrm{X})$ such that $\Delta_{\mathrm{A}} \subset\{0,1, \infty\}$. The morphisms are the homomorphisms of $\overline{\mathbb{Q}}(\mathrm{X})$-algebras.

The Galois group $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ acts on the right on $\mathscr{A}$ in the following sense: any element $g \in \mathbb{G}$ defines a (covariant) functor $g^{*}: \mathscr{A} \rightarrow \mathscr{A},(g . h)^{*}=h^{*} \circ g^{*}$ and $\left(1_{\overline{\mathbb{Q}}}\right)^{*}=1_{\mathscr{A}}$. If A is an etale $\overline{\mathbb{Q}}(\mathrm{X})$-algebra, denoting by $\alpha: \overline{\mathbb{Q}}(\mathrm{X}) \rightarrow \mathrm{A}$ the ring homomorphism defining its structure, the algebra $g^{*} \mathrm{~A}$ is the ring A equipped with $\alpha \circ g$. By (7.3.2, Proposition), $\Delta_{g^{*} \mathrm{~A}}=\Delta_{\mathrm{A}}$. Hence $g^{*} \mathrm{~A}$ is indeed an object of $\mathscr{A}$. If $\varphi: \mathrm{A} \rightarrow \mathrm{B}$ is a morphism of $\frac{9}{\mathbb{Q}}(\mathrm{X})$-algebras, then the same map is a morphism $g^{*} \mathrm{~A} \rightarrow g^{*} \mathrm{~B}$. Hence, we may set $g^{*} \varphi=\varphi$.

Remark For each $g$ the functor $g^{*}$ is a covariant functor from $\mathscr{A}$ to itself, but $g^{*}$ depends "contravariantly" on $g$.

### 7.3.5 Characterization of the Elements of an Etale Algebra

Proposition Let K be an algebraically closed subfield of $\mathbb{C}$ and A an etale $\mathrm{K}(\mathrm{X})$-algebra. Set $\mathrm{E}=\mathbb{C}(\mathrm{X}) \otimes_{\mathrm{K}(\mathrm{X})} \mathrm{A}$; identify $\mathrm{K}(\mathrm{X})$ and A with subsets of E by $f \mapsto f 1$ and $a \mapsto 1 a$. Then A is the algebraic closure $\mathrm{K}(\mathrm{X})$ in E .

Proof The algebra A is finite over $\mathrm{K}(\mathrm{X})$, and hence is in the algebraic closure of $\mathrm{K}(\mathrm{X})$ in E . We show that no $f \in \mathrm{E}-\mathrm{A}$ is algebraic over $\mathrm{K}(\mathrm{X})$.

Identify A with $\mathrm{K}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$, where P is a polynomial of degree $d=[\mathrm{A}: \mathrm{K}(\mathrm{X})]$ (minimal polynomial of a primitive element). Then $\left(1, y, \ldots, y^{d-1}\right)$ is a basis for A as a $\mathrm{K}(\mathrm{X})$-vector space, and of E as a $\mathbb{C}(\mathrm{X})$-vector space. Any $f \in \mathrm{E}$ can be uniquely written as $f_{0}+f_{1} y+\cdots+f_{d-1} y^{d-1}$. If $f \in \mathrm{E}-\mathrm{A}$, then for some $i, f_{i}$ is in $\mathbb{C}(\mathrm{X})-\mathrm{K}(\mathrm{X})$. This $f_{i}$ can be uniquely written as $\frac{p}{q}$, where $p, q \in \mathbb{C}[\mathrm{X}]$ such that g.c.d. $(p, q)=1$ and $q$ monic: at least one of the coefficients of $p$ or $q$ is in $\mathbb{C}-K$, hence transcendental over $K$. The orbit of this coefficient under the action of $\mathrm{G}_{\mathrm{K}}=\mathrm{Aut}_{\mathrm{K}}(\mathbb{C})$ is infinite. Hence so is the orbit of $f$ under the action of $\mathrm{G}_{\mathrm{K}}$ acting on $\mathbb{C}(\mathrm{X})[\mathrm{Y}] /(\mathrm{P})$ by keeping fixed X and Y . As E is a finite product of extensions of $\mathrm{K}(\mathrm{X})$, this shows that $f$ is transcendental over $\mathrm{K}(\mathrm{X})$.

### 7.3.6 The Categories $\mathscr{E}$ and $\mathscr{V}$

Let $\mathscr{E}$ be the category of etale $\mathbb{C}(\mathrm{X})$-algebras E such that $\Delta_{\mathrm{E}} \subset\{0,1, \infty\}$, the morphisms being the homomorphisms of $\mathbb{C}(\mathrm{X})$-algebras.

Let $\mathscr{V}$ be the category $\mathscr{V}_{\Sigma,\{0,1, \infty\}}$ of analytic ramified coverings X of the Riemann sphere $\Sigma$ such that $\Delta_{\mathrm{X}} \subset\{0,1, \infty\}$. The functor $\mathscr{M}$ which assigns to each

Riemann surface S the algebra $\mathscr{M}(\mathrm{S})$ of meromorphic functions on S defines an anti-equivalence of categories $\mathscr{M}: \mathscr{V} \rightarrow \mathscr{E}$ (6.2.4). (6.2.11).

As a result of (7.2.5), every algebra object E of $\mathscr{E}$ contains a $\overline{\mathbb{Q}}(\mathrm{X})$-subalgebra A such that $E$ is identified with $A_{\mathbb{C}}=\mathbb{C}(X) \otimes_{\overline{\mathbb{Q}}(X)} A$, and by (7.3.5) this subalgebra is the algebraic closure of $\overline{\mathbb{Q}}(X)$ in $E$. Hence, the functor of scalar extension $A \mapsto A_{\mathbb{C}}$ from $\mathscr{A}$ to $\mathscr{E}$ and the functor $\alpha: \mathscr{E} \rightarrow \mathscr{A}$ which assigns to each $\mathbb{C}(\mathrm{X})$-algebra E the algebraic closure of $\overline{\mathbb{Q}}(\mathrm{X})$ in E are quasi-inverse equivalences of categories.

### 7.3.7 A Diagram of Categories

By (6.1.10, Corollary), the forgetful functor $\mathscr{V} \rightarrow \mathscr{R} \mathscr{R} .=\mathscr{R}_{\Sigma,\{0,1, \infty\}}$ (category of finite topological ramified coverings $S$ of $\Sigma$ such that $\Delta_{\mathrm{S}} \subset\{0,1, \infty\}$ ) is an equivalence of categories. By (6.1.11, Proposition) the restriction functor $\left.S \mapsto S\right|_{\mathbb{C}-\{0,1, \infty\}}$ is an equivalence of categories from $\mathscr{R} \mathscr{R}$. onto the category $\mathscr{C} \boldsymbol{\bullet} \boldsymbol{\mathcal { R }} .=\mathscr{C} \boldsymbol{\bullet} \boldsymbol{\mathcal { F }}_{\mathbb{C}-\{0,1\}}$ of finite coverings of $\mathbb{C}-\{0,1\}$.

Choose a basepoint $z_{0}$ in $\mathbb{C}-\{0,1\}$, for example $z_{0}=\frac{1}{2}$, and let $\Gamma$ be the fundamental group $\pi_{1}\left(\mathbb{C}-\{0,1\}, z_{0}\right)$. This is a (non commutative) free group on 2 generators $\gamma_{0}$ and $\gamma_{1}$, represented by loops around 0 and 1 respectively (6.4). Its profinite completion $\widehat{\Gamma}=\widehat{\Gamma}_{T}=\widehat{\pi}_{1}\left(\mathbb{C}-\{0,1\}, z_{0}\right)$ is a profinite free group on 2 generators, which we also write $\gamma_{0}$ and $\gamma_{1}$. By (4.5.3, Theorem), the functor $\mathbf{S} \mapsto \mathbf{S}\left(z_{0}\right)$ defines an equivalence from $\mathscr{C o v f}$. onto the categories $\Gamma$ - $\operatorname{Set} \mathcal{f}$ (resp. $\widehat{\Gamma}$ - $\operatorname{Set} \mathcal{f}$ ) of finite sets on which $\Gamma$ acts (resp. $\widehat{\Gamma}$ acts continuously), in other words of finite sets equipped with two permutations.

Let $\Omega$ be an algebraic closure of $\overline{\mathbb{Q}}(\mathrm{X})$, and $\Omega_{\{0,1, \infty\}}$ the (by (6.2.11) directed) union sub-extension in $\Omega$ of sub-extensions that are objets of $\mathscr{A}$. Let $\widehat{\Gamma}_{a \notin g}$ be the Galois group of Aut $_{\overline{\mathbb{Q}}(\mathrm{X})} \Omega_{\{0,1, \infty\}}$. By (5.9.4 Theorem), the functor $\mathrm{A} \mapsto \operatorname{Hom}_{\overline{\mathbb{Q}}(\mathrm{X})-A \mathscr{A}}$ (A, $\Omega_{\{0,1, \infty\}}$ ) defines an anti-equivalence of categories from $\mathscr{A}$ onto $\widehat{\Gamma}_{a \ell g}$ - $\mathcal{S e t} \mathcal{f}$.

This gives a diagram of functors all of which are equivalences or anti-equivalences of categories (the arrows representing contravariant functors are circled):

$$
\widehat{\Gamma}_{a t g}-\operatorname{Set} f \leftarrow-\mathscr{A} \rightleftharpoons \mathscr{E} \stackrel{M}{\leftarrow}-\mathscr{V} \rightarrow \mathscr{R} . \rightarrow \mathscr{C o v f} . \rightarrow \Gamma-\operatorname{Set} f \leftarrow \widehat{\Gamma}-\operatorname{Set} \mathcal{R}
$$

In the following, we expand this diagram.
Remarks (1) The categories $\widehat{\Gamma}_{\boldsymbol{a f} \ell}-\operatorname{Set} \mathcal{f}$ and $\widehat{\Gamma}$ - $\operatorname{Set} \boldsymbol{f}$ being equivalent, the profinite groups $\widehat{\Gamma}$ and $\widehat{\Gamma}_{\boldsymbol{a} g}$ are isomorphic (2.9.8, Corollary); in particular $\widehat{\Gamma}_{\boldsymbol{a} \ell}$ is a profinite free group on 2 generators.
(2) The group $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ acts on $\mathscr{A}$. Hence it also acts on all categories in the above diagram. For example, to each $g \in \mathbb{G}$ can be assigned a functor $g_{*}$ from the category $\mathscr{C}$ ovf. to itself in such a way that the functors $(g \cdot h)_{*}$ and $g_{*} \circ h_{*}$ are isomorphic for $g$ and $h$ in $\mathbb{G}$, and $1_{*} \approx 1_{\mathscr{C} \text { ouf }}$. These are merely isomorphic functors following from the transference of the action of $\mathbb{G}$ on $\mathscr{A}$ using equivalences of categories without quasi-inverses. There are also anti-equivalences, but the functors
$g_{*}$ are always covariant, and $(g . h)_{*} \approx g_{*} \circ h_{*}$ is always the case. In other words the group $\mathbb{G}$ always acts on the left.

Proposition Let E and $\mathrm{E}^{\prime}$ be finite $\widehat{\Gamma}$-sets with $\mathrm{E}^{\prime}=g_{*}(\mathrm{E})$ for $g \in \mathbb{G}$. Then there is an automorphism $\phi$ of $\widehat{\Gamma}$ such that $\mathrm{E}^{\prime} \approx \phi^{*} \mathrm{E}$. In particular there is a bijection $\beta: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ mapping $\widehat{\Gamma}$ (or $\Gamma$ ) onto the the group $\Gamma_{\mathrm{E}^{\prime}}$ in $\mathfrak{S}\left(\mathrm{E}^{\prime}\right)$, which is the image of $\Gamma_{\mathrm{E}}$ (similarly defined) under $f \mapsto \beta \circ f \circ \beta^{-1}$.

Proof The element $g \in \mathbb{G}$ defines an equivalence $g_{*}$ from the category $\widehat{\Gamma}$ - $\operatorname{set} \boldsymbol{f}$ to itself. By (2.9.8, Corollary) there is an automorphism $\phi$ of $\widehat{\Gamma}$ such that the functors $g_{*}$ and $\phi^{*}$ are isomorphic. An isomorphism from $\phi^{*}$ onto $g_{*}$ gives a bijection $\beta: \mathrm{E} \rightarrow \mathrm{E}^{\prime}$ such that $\beta(\phi(g) \cdot x)=g \cdot \beta(x)$. Letting $\sigma_{g}$ (resp. $\sigma_{g}^{\prime}$ ) be the action of $g \in \Gamma$ on E (resp. on $\mathrm{E}^{\prime}$ ), $\beta \circ \sigma_{\phi(g)}=\sigma_{g}^{\prime} \circ \beta$, i.e. $\sigma_{g}^{\prime}=\beta \circ \sigma_{\phi(g)} \circ \beta^{-1}$.

### 7.3.8 Cyclic Order

Let $O$ be a finite set of cardinality $k$. A cyclic order on $O$ is an equivalence class consisting of bijections $\varphi: \mathrm{O} \rightarrow \mathbb{Z} /(k)$ for the equivalence relation identifying $\varphi$ with $\varphi+a, a$ being a constant. The bijections in this class are called the cyclic charts for the cyclically ordered set O .

The map $\sigma=\varphi \circ \mathrm{T} \circ \varphi^{-1}: \mathrm{O} \rightarrow \mathrm{O}$, where T is the translation $t \mapsto t+1$ in $\mathbb{Z} /(k)$, is independent of the choice of the cyclic chart $\varphi$. Each $x$ has a next element $\sigma(x)$. If E is a finite set and $f$ a permutation of E , i.e. a bijection $\mathrm{E} \rightarrow \mathrm{E}$, then there is a unique cyclic order on each orbit O of $f$ such that the restriction of $f$ to O is the map taking elements to the next one. Defining a cyclic order on O is thus equivalent to defining a permutation of O acting transitively (assuming that $\mathrm{O} \neq \varnothing$ : there is no cyclic order on $\varnothing!$ )

Let O and $\mathrm{O}^{\prime}$ be two cyclically ordered sets of cardinalities $k$ and $k^{\prime}$. A map $h: \mathrm{O}^{\prime} \rightarrow \mathrm{O}$ will be said to be ramified if $h \circ \sigma^{\prime}=\sigma \circ h$, where $\sigma$ and $\sigma^{\prime}$ are the maps assigning the next element in O and $\mathrm{O}^{\prime}$. As a result, $k^{\prime}$ is a multiple of $k$, and the expression $\varphi^{-1} \circ \varphi^{\prime}$ in the cyclic charts is of the form $i \mapsto \chi(i)+a$, for some canonical map $\chi: \mathbb{Z} /\left(k^{\prime}\right) \rightarrow \mathbb{Z} /(k)$ and constant $a \in \mathbb{Z} /(k)$.

### 7.3.9 Graphs

Here we give a different definition of graphs from the one used in § 4.8.
A finite graph consists of:
(1) a finite set N whose elements are called vertices;
(2) a finite set A whose elements are called directed edges;
(3) two maps $\varepsilon_{0}, \varepsilon_{1}: \mathrm{A} \rightarrow \mathrm{N}$ called source and range;
(4) a fixed point free reversing involution $\tau: \mathrm{A} \rightarrow \mathrm{A}$ such that $\varepsilon_{0} \circ \tau=\varepsilon_{1}$ (whence $\varepsilon_{1} \circ \tau=\varepsilon_{0}$ ).

In this chapter, for simplicity's sake a finite graph will be called a graph. The main differences with the definition given in (4.8.1) are:
(1) $\left(\varepsilon_{0}, \varepsilon_{1}\right): \mathrm{A} \rightarrow \mathrm{N}^{2}$ is not assumed to be injective (two vertices may be connected by several edges), so A cannot be identified to a subset of $\mathrm{N} \times \mathrm{N}$,
(2) the image of $\left(\varepsilon_{0}, \varepsilon_{1}\right)$ is not assumed to avoid the diagonal: loops are allowed,
(3) N and A are assumed to be finite.

The geometric realization $|\mathrm{H}|$ of a graph $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}\right)$ is the quotient of $\mathrm{N} \sqcup(\mathrm{A} \times[0,1])$ by the equivalence relation generated by $(a, 0) \sim \varepsilon_{0}(a),(a, 1) \sim$ $\varepsilon_{1}(a)$ and $(a, t) \sim(\tau(a), 1-t)$ for $a \in \mathrm{~A}$ and $t \in[0,1]$.

The elements of $\mathrm{A} / \tau$ are called (undirected) edges.

### 7.3.10 The Category $\mathscr{H}$

A bicoloured graph is a graph $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}, \tau\right)$ equipped with a map $\pi$ : $\mathrm{N} \rightarrow\{0,1\}$ such that $\pi\left(\varepsilon_{1}(a)\right) \neq \pi\left(\varepsilon_{0}(a)\right)$ for $a \in \mathrm{~A}$. A vertex $s$ is called white if $\pi(s)=0$, black if $\pi(s)=1$. The set A is the disjoint union of the set $\mathrm{A}^{+}$of outgoing directed edges (white source, black range) and of the set $\mathrm{A}^{-}$of incoming directed edges (black source, white range).

The category of bicoloured graphs cyclically ordered at each vertex. will be denoted $\mathscr{H}$. Hence an object of $\mathscr{H}$ is a bicoloured graph $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}, \tau, \pi\right)$ equipped, for each vertex $s \in \mathrm{~N}$, with a cyclic order $\omega_{s}$ on the set $\mathrm{A}_{s}$ of directed edges with source $s$.

Defining cyclic orders $\omega_{s}$ is equivalent to defining the map $\sigma: \mathrm{A} \rightarrow \mathrm{A}$ assigning to each directed edge $a$, the next one with respect to the cyclic order $\omega_{\varepsilon_{0}(a)}$, and subject to the condition $\varepsilon_{0} \circ \sigma=\varepsilon_{0}$ and to acting transitively on the fibres of $\varepsilon_{0}$.

A morphism $f:\left(\mathrm{H}^{\prime},\left(\omega_{s^{\prime}}^{\prime}\right)\right) \rightarrow\left(\mathrm{H},\left(\omega_{s}\right)\right)$ is given by two maps $f_{\mathrm{N}}: \mathrm{N}^{\prime} \rightarrow \mathrm{N}$ and $f_{\mathrm{A}}: \mathrm{A}^{\prime} \rightarrow$ A such that $\varepsilon_{0} \circ f_{\mathrm{A}}=f_{\mathrm{N}} \circ \varepsilon_{0}^{\prime}, \varepsilon_{1} \circ f_{\mathrm{A}}=f_{\mathrm{N}} \circ \varepsilon_{1}^{\prime}, f_{\mathrm{A}} \circ \tau^{\prime}=\tau \circ f_{\mathrm{A}}^{\prime}, \pi \circ$ $f_{\mathrm{N}}=\pi^{\prime}$, and such that $f_{\mathrm{A}}$ induces a branching map $f_{s}^{\prime}: \mathrm{A}_{s^{\prime}}^{\prime} \rightarrow \mathrm{A}_{f\left(s^{\prime}\right)}$ for all vertices $s^{\prime}$ of $\mathrm{H}^{\prime}$ (in other words, $f_{\mathrm{A}} \circ \sigma^{\prime}=\sigma \circ f_{\mathrm{A}}$ ).

Remark A graph with a cyclic order at the vertices has no isolated vertex: the map $\varepsilon_{0}: \mathrm{A} \rightarrow \mathrm{N}$ is surjective since the emptyset has no cyclic order, and so is $\varepsilon_{1}$ since $\varepsilon_{1}=\varepsilon_{0} \circ \tau$. The map $f_{\mathrm{N}}$ in a morphism of $\mathscr{H}$ is thus determined by $f_{\mathrm{A}}$.

Let $\sigma_{0}\left(\right.$ resp. $\left.\sigma_{1}\right)$ be the permutation of $\mathrm{A}^{+}$induced by $\sigma($ resp. by $\tau \circ \sigma \circ \tau)$.
Proposition The functor $\left(\left(\mathrm{H},\left(\omega_{s}\right)\right) \mapsto\left(\mathrm{A}^{+}, \sigma_{0}, \sigma_{1}\right)\right.$ is an equivalence of categories from $\mathscr{H}$ onto the category $\mathscr{H}$ of finite sets equipped with two permutations $\sigma_{0}, \sigma_{1}$.

Proof Define $(\mathrm{A}, \sigma)$ from $\left(\mathrm{B}, \sigma_{0}, \sigma_{1}\right)$, where $\sigma_{0}$ and $\sigma_{1}$ are permutations of B , by $\mathrm{A}=\mathrm{B} \times\{0,1\}$ and $\sigma(b, i)=\left(\sigma_{i}(b), i\right)$. Then define the object $\left(\mathrm{H},\left(\omega_{s}\right)\right)$ of $\mathscr{H}$ by taking $\mathrm{N}_{0}$ to be the set $\mathrm{B} / \sigma_{0}$ of orbits of $\sigma_{0}, \mathrm{~N}_{1}=\mathrm{B} / \sigma_{1}$, and $\mathrm{N}=\mathrm{N}_{0} \sqcup \mathrm{~N}_{1}$. Set $\tau(b, i)=(b, 1-i)$. Take $\varepsilon_{i}(b, j)$ to be the class of $b$ in $\mathrm{N}_{i+j(\bmod 2)}, \pi=0$ on $\mathrm{N}_{0}$
and 1 on $\mathrm{N}_{1}$, and $\omega_{s}$ to be the cyclic order with respect to which the next element is given by the restriction of $\sigma$.

This defines a functor $\breve{\mathscr{H}} \rightarrow \mathscr{H}$. It is a quasi-inverse of the functor $\mathscr{H} \rightarrow \breve{\mathscr{H}}$.

### 7.3.11 The Functor $\theta: \mathscr{R} \mathscr{R} . \rightarrow \mathscr{H}$

If $\left(\mathrm{S}, \pi_{\mathrm{S}}\right)$ is an object of $\mathscr{R} \mathscr{R} .=\mathscr{R}_{\Sigma,\{0,1, \infty\}}$, then the space $\pi_{\mathrm{S}}^{-1}([0,1])$ is identified with the geometric realization of the bicoloured graph $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}, \tau, \pi\right)$ defined as follows:

- $\mathrm{N}=\pi_{\mathrm{S}}^{-1}(\{0,1\})$; the map $\pi: \mathrm{N} \rightarrow\{0,1\}$ is the restriction of $\pi_{\mathrm{S}}$;
- $\mathrm{A}^{+}$(resp. $\mathrm{A}^{-}$) is the set of continuous maps $a:[0,1] \rightarrow \mathrm{S}$ such that $\pi(a(t))=t$ (resp. $=1-t$ ),
- $\tau(a)$ is the map $t \mapsto a(1-t)$,
- $\varepsilon_{0}(a)=a(0), \varepsilon_{1}(a)=a(1)$.

Indeed, the geometric realization $|\mathrm{H}|$ of the graph H is a quotient of $\mathrm{A} \times[0,1]$, the map $(a, t) \mapsto a(t)$ of $\mathrm{A} \times[0,1]$ passes to the quotient and defines a continuous $\operatorname{map}|\mathrm{H}| \rightarrow \pi_{\mathrm{S}}^{-1}([0,1])$. It can be checked to be bijective; as $|\mathrm{H}|$ is compact, it is a homeomorphism.

Define as follows a cyclic order at the vertices of H. If $s \in \mathrm{~S}(0)$ (resp. $\mathrm{S}(1)$ ), there is (6.1.8) a chart $\varphi$ for S in the neighbourhood of $s$ such that $\pi(x)=(\varphi(x))^{d}$ (resp. $\left.\pi(x)=(1-\varphi(x))^{d}\right)$, where $d$ is the local degree of the projection $\pi: \mathrm{S} \rightarrow \Sigma$ at $s$. The directed edges $a_{\ell}$ with source $s$ are given by $\varphi\left(a_{\ell}(t)\right)=e^{2 i \pi \ell / d} t$, for $\ell \in \mathbb{Z} /(d)$. The bijection $\ell \mapsto a_{\ell}$ defines a cyclic order $\omega_{s}$ on $\mathrm{A}_{s}$, independent of the choice of $\varphi$. The bicoloured graph H , equipped with $\left(\omega_{s}\right)_{s \in \mathrm{~N}}$, is an object of $\mathscr{H}$ which will be written $\theta(\mathbf{S})$.

If $f: \mathrm{S}^{\prime} \rightarrow \mathrm{S}$ is a morphism of $\mathscr{R} \mathscr{R}$., then the map $f$ induces $f_{\mathrm{N}}: \mathrm{N}^{\prime} \rightarrow \mathrm{N}$ and $f_{\mathrm{A}}: \mathrm{A}^{\prime} \rightarrow \mathrm{A}$. These define a morphism $f_{*}$ of bicoloured graphs. The map $f$ turns $\mathrm{S}^{\prime}$ into a ramified covering of $S$ (this can be either directly seen or by using (6.1.10) and (6.1.5). Consequently, for $s^{\prime} \in \mathrm{S}^{\prime}, f_{\mathrm{A}}$ induces a branching map $\mathrm{A}_{s^{\prime}}^{\prime} \rightarrow \mathrm{A}_{f\left(s^{\prime}\right)}$. Hence $f_{*}$ is a morphism of $\mathscr{H}$, and $\theta$ becomes a functor $\mathscr{R} \mathscr{R} . \rightarrow \mathscr{H}$.

We construct a quasi inverse functor $\eta$ of a functor $\dot{\theta}$ related to $\theta$, and then deduce that $\theta$ is an equivalence of categories.

### 7.3.12 The Functor $\eta$

Let $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}, \tau, \pi,\left(\omega_{s}\right)\right)$ be an object of $\mathscr{H}$, i.e. a bicoloured graph cyclically ordered at each vertex. We construct a ramified covering $\dot{\mathrm{S}}$ of $\mathbb{C}$ from $\mathscr{H}$, with $\Delta_{\dot{\mathrm{s}}} \subset\{0,1\}$.

Let $\overline{\mathbb{H}}$ be the closed upper half-plane: $\overline{\mathbb{H}}=\{x+i y \mid x \in \mathbb{R}, y \geqslant 0\}$, and consider $\Xi=\mathrm{A} \times \overline{\mathbb{H}}$. Define $\pi_{\Xi}: \Xi \rightarrow \mathbb{C}$ by $\pi_{\Xi}(a, z)=z$ if $a$ is an outgoing edge, $\pi_{\Xi}(a, z)=1-z$ if $a$ is an incoming edge. Equip $\Xi$ with the equivalence relation generated by:

- $(a, t) \sim(\tau(a), 1-t)$ for $t \in[0,1]$,
- $(a, t) \sim(\tau(\sigma(a)), 1-t)$ for $t \leqslant 0$, where $\sigma(a)$ is the directed edge along $a$ with respect to the cyclic order at the source vertex of $a$.

These conditions imply $(a, 0) \sim\left(a^{\prime}, 0\right)$ if $a$ and $a^{\prime}$ have the same source, and $(a, 1) \sim\left(a^{\prime}, 1\right)$ if $a$ and $a^{\prime}$ have the same range. Indeed, $(a, 0) \sim(\tau(\sigma(a)), 1) \sim$ $(\sigma(a), 0),(a, 1) \sim(\tau(a), 0) \sim\left(\sigma_{1}(a), 1\right)$ where $\sigma_{1}=\tau \circ \sigma \circ \tau$, and $\sigma$ (resp. $\left.\sigma_{1}\right)$ acts transitively on the set of directed edges having a given source (resp. range).

Let $S$ be the quotient space of $\Xi$ by the equivalence relation thus defined. Two equivalent points have the same image under the map $\pi_{\Xi}$. Hence there is a continuous map $\pi_{\mathrm{S}}: \mathrm{S} \rightarrow \mathbb{C}$. Let $\eta(\mathrm{H})$ be the space S equipped with $\pi_{\mathrm{S}}$.

Proposition (a) The space S equipped with $\pi_{\mathrm{S}}$ is a finite ramified covering of $\mathbb{C}$, with $\Delta_{\mathrm{S}} \subset\{0,1\}$.
(b) If $f=\left(f_{\mathrm{N}}, f_{\mathrm{A}}\right): \mathrm{H}^{\prime} \rightarrow \mathrm{H}$ is a morphism in $\mathscr{H}$, by passing to the quotient derived from $f_{\mathrm{A}} \times 1_{\overline{\mathbb{H}}}: \mathrm{A}^{\prime} \times \overline{\mathbb{H}} \rightarrow \mathrm{A} \times \overline{\mathbb{H}}$ gives a morphism $f_{*}: \mathrm{S}^{\prime} \rightarrow \mathrm{S}$ in $\mathscr{R}_{\mathbb{C},\{0,1\}}$. This defines a covariant functor $\eta: \mathscr{H} \rightarrow \mathscr{R}_{\mathbb{C},\{0,1\}}$.
(c) The functor $\theta$ is of the form $\dot{\theta} \circ \rho$, where $\rho: \mathscr{R} \mathscr{R} \rightarrow \mathscr{R}_{\mathbb{C},\{0,1\}}$ is the restriction functor, which is an equivalence of categories. The functor $\eta$ is a quasi-inverse of $\dot{\theta}$.

Proof (a) The restriction of $S$ to $(\mathbb{C}-\mathbb{R}) \cup] 0,1[$ is a quotient of $\mathrm{A} \times(\mathbb{H} \cup] 0,1[)$ by the equivalence relation identifying $(a, t)$ with $(\tau(a), 1-t)$ for $t \in] 0,1[$; it is a trivial covering of the open subset $(\mathbb{C}-\mathbb{R}) \cup] 0,1\left[\right.$, with fibre $\mathrm{A}^{+}$.

The restriction of $S$ to $(\mathbb{C}-\mathbb{R}) \cup \mathbb{R}_{-}^{*}$ is the quotient of

$$
\left(\mathrm{A}^{+} \times\left(\mathbb{H} \cup \mathbb{R}_{-}^{*}\right) \cup\left(\mathrm{A}^{-} \times(\mathbb{H} \cup] 1, \infty[)\right.\right.
$$

by the equivalence relation identifying $(a, t)$ with $(\tau(\sigma(a)), 1-t)$ for $a \in \mathrm{~A}^{+}, t<0$; it is a trivial covering of the open subset $(\mathbb{C}-\mathbb{R}) \cup \mathrm{R}_{-}^{*}$, with fibre $\mathrm{A}^{+}$.

The restriction of $S$ to $(\mathbb{C}-\mathbb{R}) \cup] 1, \infty[$ is the quotient of

$$
\left(\mathrm{A}^{-} \times\left(\mathbb{H} \cup \mathbb{R}_{-}^{*}\right) \cup\left(\mathrm{A}^{+} \times(\mathbb{H} \cup] 1, \infty[)\right.\right.
$$

by the equivalence relation identifying $(a, t)$ with $(\tau(\sigma(a)), 1-t)$ for $a \in \mathrm{~A}^{-}, t<0$; it is a trivial covering of the open subset $(\mathbb{C}-\mathbb{R}) \cup] 1, \infty\left[\right.$, with fibre $\mathrm{A}^{+}$.

Hence, $S$ induces a covering of $\mathbb{C}-\{0,1\}$.
Let $\chi$ be the quotient map $\Xi \rightarrow \mathrm{S}, a \in \mathrm{~A}$, and $d$ the number of directed edges with source $s=\varepsilon_{0}(a)$. Define $\varphi_{a}$ in the neighbourhood of $\dot{s}=\chi(a, 0)$ in S by $\varphi_{a}\left(\chi\left(\sigma^{\ell}(a)\right), z\right)=e^{2 i \pi \ell / d} \cdot z^{1 / d}$, and $\varphi_{a}\left(\chi\left(\tau \sigma^{\ell}(a)\right), z\right)=e^{2 i \pi \ell / d} \cdot(1-z)^{1 / d}$, computed with $\operatorname{Arg}(z) \in[0, \pi]$ and $\operatorname{Arg}(1-z) \in[-\pi, 0]$ for $z \in \overline{\mathbb{H}}$. Then $\varphi_{a}$ is a chart for S in the neighbourhood of $\dot{s}$, and letting $\psi$ denote the chart for $\mathbb{C}$ which is the identity if $a \in \mathrm{~A}^{+}$and $z \mapsto 1-z$ if $a \in \mathrm{~A}_{-}$, the expression of $\pi_{\mathrm{S}}$ in the charts $\varphi_{a}$ and $\psi$ is $z \mapsto z^{d}$.

The cyclic order defined by this chart on the set of directed edges originating at $s$ agrees with $\omega_{s}$.
(b) It suffices to check.
(c) It is a matter of defining for an object H of $\mathscr{H}$ an isomorphism $\mathrm{H} \rightarrow \mathrm{H}^{\prime}=$ $\dot{\theta}(\eta(\mathrm{H}))$, and for a ramified covering S of $\mathbb{C}$ with $\Delta_{\mathrm{S}} \subset\{0,1\}$ an isomorphism $\mathrm{S} \rightarrow \mathbf{S}^{\prime}=\eta(\dot{\theta}(\mathrm{S}))$.

Let $\mathrm{H}=(\mathrm{N}, \mathrm{A}, \ldots)$ be an object of $\mathscr{H}$; set $\left(\mathrm{S}, \pi_{\mathrm{S}}\right)=\eta(\mathrm{H})$ with $\mathrm{S}=\Xi / \sim$ as above, and $\mathrm{H}^{\prime}=\left(\mathrm{N}^{\prime}, \mathrm{A}^{\prime}, \ldots\right)=\theta\left(\mathrm{S}, \pi_{\mathrm{S}}\right)$. The set $\mathrm{N}_{0}^{\prime}=\pi_{\mathrm{S}}^{-1}(0)$ is the quotient of $\mathrm{A}^{+}$ by the equivalence relation identifying $a$ with $a^{\prime}$ if $\varepsilon_{0}(a)=\varepsilon_{0}\left(a^{\prime}\right)$, hence can be identified with $\mathrm{N}_{0}$. Likewise, $\mathrm{N}_{1}^{\prime}$ can be identified with $\mathrm{N}_{1}$. Assigning to $a \in \mathrm{~A}$ the map $t \mapsto \chi(a, t)$ defines a bijection $\mathrm{A} \rightarrow \mathrm{A}^{\prime}$ compatible with $\varepsilon_{0}, \varepsilon_{1}, \tau$ and $\pi$. Hence there is an isomorphism $\mathrm{H} \rightarrow \mathrm{H}^{\prime}$ to $\mathscr{H}$.

Let $\left(\mathrm{S}, \pi_{\mathrm{S}}\right)$ be a ramified covering of $\mathbb{C}$ with $\Delta_{\mathrm{S}} \subset\{0,1\}, \mathrm{H}=(\mathrm{N}, \mathrm{A}, \ldots)=\dot{\theta}(\mathrm{S})$ and $\left(\mathrm{S}^{\prime}, \pi_{\mathrm{S}^{\prime}}\right)=\eta(\mathrm{H})$. The space $\mathrm{S}^{\prime}$ is a quotient of $\Xi^{\prime}=\mathrm{A} \times \overline{\mathbb{H}}$. For $a \in \mathrm{~A}$, the map $a:[0,1] \rightarrow \mathrm{S}$ is a unique continuous extension $\widehat{a}$ to $\overline{\mathbb{H}}$ such that $\pi_{\mathrm{S}} \circ \widehat{a}(z)=z$ if $a \in \mathrm{~A}^{+}$, and $1-z$ if $a \in \mathrm{~A}^{-}$. Hence there is a continuous map $\Xi^{\prime} \rightarrow \mathrm{S}$ compatible with the equivalence relation, and so giving a map $S^{\prime} \rightarrow S$. It can be checked that the latter is a homeomorphism over $\mathbb{C}$, and that the entire construction depends functoriality on S .

Corollary The functor $\theta: \mathscr{R} \mathscr{R} . \rightarrow \mathscr{H}$ is an equivalence of categories.

### 7.3.13 Action of $\mathbb{G}$ on $\mathscr{H}$

So there is a diagram of equivalences of categories

$$
\underset{\downarrow}{\Gamma}-\operatorname{setf}^{1}
$$

$$
\widehat{\Gamma}_{a t g}-\operatorname{Set} f \longleftrightarrow 0 \rightarrow \mathscr{A} \longleftrightarrow \mathscr{E} \longleftrightarrow 0 \rightarrow \mathscr{V} \longleftrightarrow \mathscr{R} . \longleftrightarrow \mathscr{H}
$$

The group $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ acts on $\mathscr{A}$, hence also on all the categories of the above diagram, in particular on $\mathscr{H}$. Here too it is an action in a weak sense, as in (7.3.7). However $\mathbb{G}$ acts on the set $\mathscr{X}$ of isomorphism classes of objects of $\mathscr{H}$.

The valency of a vertex $s$ in a graph H which is an object of $\mathscr{H}$ is the number $v_{\mathrm{H}}(s)$ of directed edges having source $s$. Letting X, E and A denote the objects of $\mathscr{R} \mathscr{R}$., $\mathscr{E}$, $\mathscr{A}$ corresponding to H , the valency $v_{\mathrm{H}}(x)$ is the local degree at $s_{\mathrm{X}}$ corresponding to $s$ under the projection $X \rightarrow \Sigma$. It can also be defined algebraically: it is the dimension of the vector space $\mathrm{A}_{s} / \mathfrak{m}_{\pi(s)} \cdot \mathrm{A}_{s}$ over $\overline{\mathbb{Q}}$, where $\mathrm{A}_{s}$ is the subalgebra of A consisting of functions defined at $s$, and $\mathfrak{m}_{\pi(s)}$ the ideal consisting of functions of $\overline{\mathbb{Q}}(\mathrm{X})$ defined and vanishing at $\pi(s)$.

A valency function $\nu_{\mathrm{H}}$ is the map $\{0,1\} \times \mathbb{N} \rightarrow \mathbb{N}$ assigning to $(u, v)$ the number of vertices $s$ of H such that $\pi(s)=u$ and $\nu_{\mathrm{H}}(s)=v$. For $g \in \mathbb{G}$ and H an object of $\mathscr{H}$, the graphs H and $g^{*}(\mathrm{H})$ have the same valency function.

Several equivalence relations can be considered on the set $\mathscr{X}$ : Given $x, x^{\prime} \in \mathscr{X}$ represented by the objets H and $\mathrm{H}^{\prime}$ of $\mathscr{H}$, we write

- $x \sim_{g b} x^{\prime}$ or $\mathrm{H} \sim_{g b} \mathrm{H}^{\prime}$, if H and $\mathrm{H}^{\prime}$ are isomorphic as bicoloured graphs, irrespective of the cyclic order at the vertices;
- $x \sim_{\nu} x^{\prime}$ or $\mathrm{H} \sim_{\nu} \mathrm{H}^{\prime}$, if H and $\mathrm{H}^{\prime}$ have the same valency function, i.e. $\nu_{\mathrm{H}}=\nu_{\mathrm{H}^{\prime}}$.
- $x \sim_{\mathbb{G}} x^{\prime}$ or $\mathrm{H} \sim_{\mathbb{G}} \mathrm{H}^{\prime}$, if there exists $g \in \mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ mapping $x$ onto $x^{\prime}$ (i.e. H into an object of $\mathscr{H}$ isomorphic to $\mathrm{H}^{\prime}$ ).

The relation $\sim_{\nu}$ is weaker than the other two: $x \sim_{g b} x^{\prime} \Rightarrow x \sim_{\nu} x^{\prime}$ and $x \sim_{\mathbb{G}} x^{\prime} \Rightarrow$ $x \sim_{\nu} x^{\prime}$ (this latter implication follows from the algebraic definition of valency).

Clearly, $x \sim_{\nu} x^{\prime} \nRightarrow x \sim_{g b} x^{\prime}$. In section(7.5) we will see that $x \sim_{g b} x^{\prime} \nRightarrow x \sim_{\mathbb{G}}$ $x^{\prime}$, (and necessarily $x \sim_{\nu} x^{\prime} \nRightarrow x \sim_{\mathbb{G}} x^{\prime}$ ), and that $x \sim_{\mathbb{G}} x^{\prime} \nRightarrow x \sim_{g b} x^{\prime}$ (in particular the action of $\mathbb{G}$ is not trivial). In fact, already in section(7.4), we will see that the action of $\mathbb{G}$ on $\mathscr{X}$ is faithful.

All this continues to hold in the full subcategory $\mathscr{T}$ of $\mathscr{H}$ whose objets are the trees, and in the set $\mathscr{y}$ of isomorphism classes of objects of $\mathscr{T}$.

### 7.3.14 A Numerical Equivalence Class

The graph along $\mathrm{H}_{1}$ has

- 4 white vertices of valency 1 ;
- 2 white vertices of valency 3 ;
- 1 white vertices of valency 4;
- 7 black vertices of valency 2 ;


There are 31 isomorphism classes of objects of $\mathscr{H}$ in the same class of $\sim_{\nu}$ (same valency function).

A bicoloured graph all of whose black vertices have valency 2 can be uniquely obtained from an uncoloured graph-or, equivalently, monocoloured white-by subdividing each edge by a black vertex. Moreover, a cyclic order at the vertices of the
monocoloured graph extends uniquely to the bicoloured graph. This gives a bijection between the set $X_{2}$ of isomorphism classes of objects of $\mathscr{H}$ with bivalent black vertices, and the set $\mathscr{Z}$ of isomorphism classes of monocoloured white graphs, cyclically ordered at each vertex.

The 31 points of $\mathscr{Z}$ corresponding to those of the class of $\sim_{\nu}$ from $\mathrm{H}_{1}$ to $\mathscr{X}_{2}$ are represented by the graphs of the page 429 .

These objects are divided into 8 classes of $\sim_{g b}$ :
$\left\{\mathrm{G}_{1}, \mathrm{G}_{2}, \mathrm{G}_{3}\right\},\left\{\mathrm{G}_{4}, \mathrm{G}_{5}\right\},\left\{\mathrm{G}_{6}, \mathrm{G}_{11}\right\}$,
$\left\{\mathrm{G}_{7}, \mathrm{G}_{8}, \mathrm{G}_{9}, \mathrm{G}_{10}, \mathrm{G}_{12}, \mathrm{G}_{13}\right\},\left\{\mathrm{G}_{14}, \mathrm{G}_{15}\right\}$,
$\left\{G_{16}, \ldots, G_{21}\right\},\left\{G_{22}, \ldots, G_{28}\right\},\left\{G_{29}, \mathbb{G}_{30}, G_{31}\right\}$.
This does not tell us much on the action of the group $\mathbb{G}$. Proposition 7.3 .7 gives a necessary condition for two dessins d'enfants to be in the same orbit of the action of $\mathbb{G}$. Two new invariants will give other necessary conditions: the extended valency function and the automorphism group.

### 7.3.15 The Extended Valency Function

To each object H of $\mathscr{H}$ corresponds an object X of the ramified covering $\mathscr{R} \mathscr{R}$. of $\Sigma$, such that $\Delta_{\mathrm{X}} \subset\{0,1, \infty\}$. The integer $\nu_{\mathrm{H}}(0, k)\left(\right.$ resp. $\left.\nu_{\mathrm{H}}(1, k)\right)$ is the number of points of X of local degree $k$ over 0 (resp. 1). Similarly, let $\nu_{\mathrm{H}}(\infty, k)$ be the number of points of $X$ of local degree $k$ over $\infty$. This defines an extended valency function $\widehat{\nu_{\mathrm{H}}}:\{0,1, \infty\} \times \mathbb{N} \rightarrow \mathbb{N}$.

If $\mathrm{H}^{\prime}=g^{*} \mathrm{H}$ with $g \in \mathbb{G}$, then $\widehat{\nu_{\mathrm{H}^{\prime}}}=\widehat{\nu_{\mathrm{H}}}$. This follows from the algebraic definition of the local degree. Define the strengthened numerical equivalence $\sim_{\widehat{\nu}}$ by setting $\mathrm{H} \sim_{\widehat{\nu}} \mathrm{H}^{\prime}$ if $\widehat{\nu_{\mathrm{H}^{\prime}}}=\widehat{\nu_{\mathrm{H}}}$. This equivalence relation is stronger that $\sim_{\nu}$ but weaker than $\sim_{\mathbb{G}}$. As will be seen in the example of 7.3.14, it is neither stronger nor weaker than $\sim_{g b}$.

Let $\mathrm{H}=\left(\mathrm{N}, \mathrm{A}, \varepsilon_{0}, \varepsilon_{1}, \tau,\left(\omega_{s}\right)_{s \in \mathrm{~N}}\right)$ be an object of $\mathscr{H}$ and X the object corresponding to $\mathscr{R} \mathscr{R}$.. The number of edges is

$$
a=\sum_{k \in \mathbb{N}} k \cdot \nu_{\mathrm{H}}(0, k)=\sum_{k \in \mathbb{N}} k \cdot \nu_{\mathrm{H}}(1, k) .
$$

This is also the degree of X as a ramified covering of $\Sigma$. The number $f$ of points of X over $\infty$ equals the number of connected components of $X-|H|=X-\pi^{-1}([0,1])$; it is also the number of cycles of the permutation $\tau \circ \sigma$ of A , and $\nu_{\mathrm{H}}(\infty, k)$ is the number of $2 k$-cycles.

The space X is connected if and only if this is the case of H . If H is connected, then X is a surface of genus $g$ given by the Euler formula $2-2 g=n-a+f$, where $n=n_{0}+n_{1}$ is the number of vertices, $n_{0}$ the number of white vertices and $n_{1}$ the number of black ones.

In the examples of 7.3.14, $n_{0}=7, n_{1}=7, a=14, f=2, g=0$. The extended valency function is given by the local degrees of the points of X over $\infty$; it varies with the example but the sum total is always 14 :
춫

$\mathrm{G}_{6}$



$\mathrm{G}_{16}$


$\mathrm{G}_{17}$

$\mathrm{G}_{18}$

13

$\mathrm{G}_{8}$

$\mathrm{G}_{9}$

$\mathrm{G}_{10}$
$\bigcirc$
$\mathrm{G}_{14}$



$\underbrace{\text { C }}_{\mathrm{G}_{30}}$

$\mathrm{G}_{24}$
$\uparrow \mathcal{A}$
$\mathrm{G}_{20} \quad \mathrm{G}_{21}$

$\mathrm{G}_{19}$

$\mathrm{G}_{25}$




C

- in $G_{1}$ to $G_{10}$, there is a point of degree 1 and a point of degree 13 ;
- in $\mathrm{G}_{14}, \mathrm{G}_{16}, \mathrm{G}_{17}$ and $\mathrm{G}_{29}$, there is a point of degree 2 and a point of degree 12 ;
- in $G_{12}, G_{13}$ and $G_{22}$, there is a point of degree 3 and a point of degree 11 ;
- in $G_{15}, G_{18}, G_{19}, G_{20}$ and $G_{30}$, there is a point of degree 4 and a point of degree 10;
- in $\mathrm{G}_{23}, \mathrm{G}_{24}$ and $\mathrm{G}_{26}$, there is a point of degree 5 and a point of degree 9 ;
- in $G_{21}$ and $G_{31}$, there is a point of degree 6 and a point of degree 8 ;
- in $\mathrm{G}_{11}, \mathrm{G}_{25}, \mathrm{G}_{27}$ and $\mathrm{G}_{28}$, there are two points of degree 7 .

The class of $\sim_{\nu}$ considered decomposes into 7 classes of $\sim_{\widehat{\nu}}$. There are at least 7 classes of $\sim_{\mathbb{G}}$.

The only graphs having a non trivial automorphism (as object of $\mathscr{H}$ ) are $\mathrm{G}_{11}, \mathrm{G}_{27}$ and $\mathrm{G}_{28}$ : for these three objets there is an automorphism of order 2 exchanging the two points over $\infty$. Hence $G_{25}$ is not in the same class of $\sim_{\mathbb{G}}$ as these three, and the lower bound becomes 8 .

All the graphs $G_{1}$ to $G_{10}$ can be shown to be in the same class of $\sim_{\mathbb{G}}$. Complex conjugation induces an element of $\mathbb{G}$ exchanging $\mathrm{G}_{12}$ and $\mathrm{G}_{13}, \mathrm{G}_{16}$ and $\mathrm{G}_{17}, \mathrm{G}_{18}$ and $\mathrm{G}_{19}, \mathrm{G}_{23}$ and $\mathrm{G}_{24}, \mathrm{G}_{27}$ and $\mathrm{G}_{28}$.

Hence, the number of classes of $\sim_{\mathbb{G}}$ in the given class of $\sim_{\nu}$ is bounded above by 17.

## Exercise 7.3. (Equivalence between various categories)

1.-Let G be a group and $\mathscr{C}$ a category. We use the following language:

An exact action $\rho$ of G on $\mathscr{C}$ consists of a functor $\rho(g): \mathscr{C} \rightarrow \mathscr{C}$ for all $g \in \mathrm{G}$ such that for $g_{1}, g_{2} \in \mathrm{G}, \rho\left(g_{1} . g_{2}\right)=\rho\left(g_{1}\right) \circ \rho\left(g_{2}\right)$.

A weak action $\rho$ of G on $\mathscr{C}$ consists of of a functor $\rho(g): \mathscr{C} \rightarrow \mathscr{C}$ for all $g \in \mathrm{G}$ such that for $g_{1}, g_{2} \in \mathrm{G}$, the functors $\rho\left(g_{1} . g_{2}\right)$ and $\rho\left(g_{1}\right) \circ \rho\left(g_{2}\right)$ are isomorphic.

Two weak actions $\rho$ and $\rho^{\prime}$ are equivalent if for all $g \in \mathrm{G}$, the functors $\rho(g)$ and $\rho^{\prime}(g)$ are isomorphic. A strong action is a weak action equivalent to an exact action.
(A) Suppose that $\mathscr{C}$ is a category with a single object $\omega$ such that $\operatorname{Hom}(\omega ; \omega)$ is a group $F$. Set $H=\operatorname{Aut}(F)$, and let $J$ be the subgroup of $H$ consisting of the inner automorphisms of F .
(a) Show that J is a normal subgroup of H .
(b) Show that the exact actions of G on $\mathscr{C}$ correspond to homomorphisms $\mathrm{G} \rightarrow \mathrm{H}$. Show that the equivalence classes of weak actions correspond to homomorphisms $\mathrm{G} \rightarrow \mathrm{H} / \mathrm{J}$.
(c) Give an example of a weak action which is not strong.
(B) Show that the actions of $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ on the categories described in Remark 2 of 7.3.7 are strong actions.
2.—Show that the composite functor $\breve{\theta}: \mathscr{R} \mathscr{R} . \rightarrow \breve{\mathscr{H}}$ of $\theta: \mathscr{R} \mathscr{R} . \rightarrow \mathscr{H}$ and of the functor $\mathscr{H} \rightarrow \breve{\mathscr{H}}$ defined in 7.3.10 assign to $(\mathrm{S}, f)$ the set $\mathrm{S}^{-1}\left(\frac{1}{2}\right)$ equipped with the permutations induced by the loops around 0 and 1 respectively.

### 7.4 Belyi Polynomials

### 7.4.1 Ramified Polynomial Coverings

Proposition and Definition Let $(\mathrm{S}, \pi)$ be an analytic ramified covering of the Riemann sphere $\Sigma$. The following conditions are equivalent:
(i) S is homeomorphic to the sphere $\mathrm{S}^{2}$ and the set $\pi^{-1}(\infty)$ is reduced to a point;
(ii) there is an isomorphism $\phi: \mathrm{S} \rightarrow \Sigma$ and a polynomial $f \in \mathbb{C}[\mathrm{X}]$ such that $\pi=f \circ \phi$.

If these conditions hold, then $(\mathrm{S}, \pi)$ is said to be a ramified polynomial covering.
Proof The implication $($ ii $) \Rightarrow$ (i) is immediate. We show that $($ i $) \Rightarrow$ (ii). If the Riemann surface of $S$ is homeomorphic to $S^{2}$, it is then isomorphic to $\Sigma$ by the Riemann uniformization theorem. If moreover $\pi^{-1}(\infty)$ is reduced to a point $e$, then there is an isomorphism $\phi: \mathrm{S} \rightarrow \Sigma$ such that $\phi(e)=\infty$. Then $f=\pi \circ \phi^{-1}$ is an analytic map $\Sigma \rightarrow \Sigma$ such that $f^{-1}(\infty)=\{\infty\}$, hence a polynomial.

Remark The isomorphism $\phi$ is uniquely defined up to left composition by an affine isomorphism of $\Sigma$, and the polynomial $f$ up to right composition: if $\pi=f_{1} \circ \phi_{1}=$ $f_{2} \circ \phi_{2}$ as above, then there is an affine isomorphism A: $\boldsymbol{\mathrm { A }} \rightarrow \Sigma$ such that $f_{2}=$ $f_{1} \circ \mathrm{~A}$ and $\phi_{2}=\mathrm{A}^{-1} \circ \phi_{1}$.

### 7.4.2

Let $(\mathrm{S}, \pi)$ be an analytic ramified covering of $\Sigma$ such that $\Delta_{\pi} \subset\{0,1, \infty\}, \mathrm{H}$ the graph $\pi^{-1}([0,1])$, and set $S^{\prime}=S-\pi^{-1}(\infty)$.

Proposition ( $\mathrm{S}, \pi$ ) is a ramified polynomial covering if and only if the graph H is a tree (i.e. is simply connected).

Lemma The graph H is a deformation retract of $\mathrm{S}^{\prime}$.
Proof Let $\overline{\mathbb{H}}($ resp. $-\overline{\mathbb{H}})$ be the closed upper (resp. lower) half-plane in $\mathbb{C}$. Let $\rho$ : $\mathbb{C} \rightarrow[0,1]$ be a continuous retraction such that $\rho\left(\mathbb{R}_{-}\right)=\{0\}$ and $\rho([1, \infty[)=\{1\}$. Let $(t, z) \mapsto \rho^{t}(z)$ be a homotopy from $\mathrm{I}_{\mathbb{C}}$ to $\rho$ such that each $\rho^{t}$ preserves $\mathbb{R}_{-}$and $\left[0, \infty\left[\right.\right.$, and fixes each point of $[0,1]$. Set $S^{+}=\pi^{-1}(\overline{\mathbb{H}})$ and $S^{-}=\pi^{-1}(-\overline{\mathbb{H}})$. These may be written

$$
\mathrm{S}^{+}=\bigcup_{i \in \mathrm{I}^{+}} \sigma_{i}(\overline{\mathbb{H}}), \quad \mathrm{S}^{-}=\bigcup_{i \in \mathrm{I}^{-}} \sigma_{i}(-\overline{\mathbb{H}}),
$$

where $\sigma_{i}$ are the continuous sections. Set $\rho_{i}^{t}=\sigma_{i} \circ \rho^{t} \circ \pi$ onto the image of $\sigma_{i}$. Then, for each $t$, the maps $\rho_{i}^{t}$ can be glued together to give a map $\tilde{\rho}^{t}: \mathrm{S}^{\prime} \rightarrow \mathrm{S}^{\prime}$, and this defines a homotopy from $\mathrm{I}_{\mathrm{S}^{\prime}}$ to a retraction $\tilde{\rho}$ from $\mathrm{S}^{\prime}$ onto H .

Proof of the Proposition If $(S, \pi)$ is polynomial, then $S^{\prime}$ is isomorphic to $\mathbb{C}$, hence is simply connected, and H is a tree. If H is a tree, then $\mathrm{S}^{\prime}$ is simply connected, hence isomorphic to $\mathbb{D}, \mathbb{C}$ or $\Sigma$ by the Riemann uniformization theorem. The space $S^{\prime}$ is a nonempty ramified covering of $\mathbb{C}$; so $\pi: S^{\prime} \rightarrow \mathbb{C}$ is surjective, $S^{\prime}$ is not compact, and $\Sigma$ is excluded. Every bounded holomorphic function on $S^{\prime}$ is constant. Indeed, if $\zeta$ is such a function, then the coefficients of the minimal polynomial of $\zeta$ over $\mathbb{C}[\mathrm{X}]$ are bounded holomorphic functions on $\mathbb{C}$, hence constant, and $\zeta$ can only take finitely many values; as $S^{\prime}$ is connected, $\zeta$ is constant. Hence, $\mathbb{D}$ is excluded, and $S^{\prime}$ is isomorphic to $\mathbb{C}$.

The set $S-S^{\prime}=\pi^{-1}(\infty)$ has exactly one point. Indeed, there is an open neighbourhood U of $\infty$ in $\Sigma$ such that each connected component of $\pi^{-1}(\mathrm{U})$ contains $a$ point of $\pi^{-1}(\infty)$ and a connected component of $\pi^{-1}(\mathrm{U}) \cap S^{\prime}$. Hence there is a bijection between $\pi^{-1}(\infty)$ and the set of relatively compact connected components of $\mathrm{S}^{\prime} \cap \pi^{-1}(\mathrm{U})=\mathrm{S}^{\prime}-\pi^{-1}(\mathbb{C}-\mathrm{U}$ ) (see 6.1, Exercise 2, end compactification). Now, $\mathbb{C}-\mathrm{U}$ is compact; so is $\pi^{-1}(\mathbb{C}-\mathrm{U})$ since $\pi$ is proper, and the complement of a compact subset of $\mathbb{C}$ has exactly one relatively compact connected component; so does $S^{\prime}$. Hence Card $\pi^{-1}(\infty)=1$.

Therefore, $S$ is isomorphic to $\Sigma$.

### 7.4.3 Belyi Polynomials

If $\Delta(f) \subset\{0,1, \infty\}, f \in \mathbb{C}[\mathrm{X}]$ is called a Belyi polynomial. For any subfield K of $\mathbb{C}$, let $\mathscr{B} \boldsymbol{e l} \underset{\sim}{ } \boldsymbol{\ell}(\mathrm{K})$ denote the set of Belyi polynomials with coefficients in K , and $\mathscr{A} \mathcal{P}(\mathrm{K})$ the group of affine transformations $z \mapsto a z+b$ with $(a, b) \in \mathrm{K}^{*} \times \mathrm{K}$. It acts on the right on $\mathscr{B} \boldsymbol{e} \boldsymbol{\not p} \boldsymbol{\ell}(\mathrm{K})$ by $(\mathrm{A}, f) \mapsto f \circ \mathrm{~A}$.

Proposition If $\mathrm{K} \subset \mathbb{C}$ is an algebraically closed subfield, then the set of isomorphism classes of ramified polynomial coverings $(\mathrm{S}, \pi)$ of $\Sigma$ such that $\Delta_{\pi} \subset\{0,1, \infty\}$ can be identified with the quotient $\mathscr{B}$ elpol $(\mathrm{K}) / \mathscr{A P f}(\mathrm{K})$.

Proof Let $(\mathrm{S}, \pi)$ be a ramified polynomial $d$-fold covering of $\Sigma$ such that $\Delta_{\pi} \subset$ $\{0,1, \infty\}$. It is definable over $\overline{\mathbb{Q}}(7.2 .5)$, and let $\mathrm{Q} \in \overline{\mathbb{Q}}[\mathrm{X}, \mathrm{Y}]$ be a defining polynomial. As (S, $\pi$ ) is polynomial, there is an isomorphism $\phi: \mathrm{S} \rightarrow \Sigma$ mapping $\pi^{-1}(\infty)$ onto $\infty$. The inverse $\Psi$ de $\phi$ induces a map $\psi: \Sigma^{\prime} \rightarrow \mathrm{V}_{\mathrm{Q}}$, where $\Sigma^{\prime}$ is the complement of a finite set. The map $\psi$ is of the form $(f, g)$, where $f \in \mathbb{C}[T]$ is a Belyi polynomial and $g \in \mathbb{C}(\mathrm{~T})$.

The isomorphism $\phi$ is not uniquely defined, but its choice determines $\psi, f$ and $g$.
Lemma $\phi$ may be chosen so that $f \in \overline{\mathbb{Q}}(\mathrm{~T})$.
Proof Let $\left(x_{0}, y_{0}\right)$ and $\left(x_{1}, y_{1}\right)$ be distinct points of $\mathrm{V}_{\mathrm{Q}}$ such that $x_{0}$ and $x_{1}$ are in $\overline{\mathbb{Q}}-\Delta_{\mathrm{Q}}$. Then $y_{0}, y_{1} \in \overline{\mathbb{Q}}$ since $\mathrm{Q}\left(x_{i}, y_{i}\right)=0$. Choose $\phi$ so that $\psi(0)=\left(x_{0}, y_{0}\right)$ and $\psi(1)=\left(x_{1}, y_{1}\right)$. Then $f \in \mathscr{B} \boldsymbol{\ell} p \neq \ell(\mathbb{C})$ and $g \in \mathbb{C}(\mathrm{~T})$ are characterized by $\mathrm{Q}(f(\mathrm{~T}), g(\mathrm{~T}))=0$ and $\operatorname{deg}(f)=d$. These are algebraic conditions defined on $\overline{\mathbb{Q}}$.

Hence the coefficients of $f$ are left invariant by all algebraic automorphisms from $\mathbb{C}$ onto $\overline{\mathbb{Q}}$; hence they are in $\overline{\mathbb{Q}}$.

End of the Proof of the Proposition We have shown that there exist $f_{0} \in$ $\mathscr{B}$ elpol $(\overline{\mathbb{Q}}) \subset \mathscr{B}$ elpol $(\mathrm{K})$ and an isomorphism $\phi_{0}: \mathrm{S} \rightarrow \Sigma$ such that $f_{0} \circ \phi_{0}=$ $\pi$. If $f \in \mathscr{B} \boldsymbol{e l} \boldsymbol{2} \boldsymbol{\ell}(\mathrm{~K})$ and $\phi: \mathrm{S} \xrightarrow{\approx} \Sigma$ satisfy $f \circ \phi=\pi$, then there exists A of the form $t \mapsto a t+b$ such that $\phi_{0}=\mathrm{A} \circ \phi$. So $\pi=f_{0} \circ \mathrm{~A} \circ \phi=f \circ \phi$, and hence $f=f_{0} \circ \mathrm{~A}$ since $\phi$ is surjective.

Comparing the leading coefficients gives $a^{d} \in \mathrm{~K}$; thus $a \in \mathrm{~K}$ since K is algebraically closed. Comparing the coefficients of the terms of degree $d-1$ shows that $b \in \mathrm{~K}$.

### 7.4.4 Statement of the Faithfulness Theorem

By 7.4.3, proposition and 7.3.13, the set of isomorphism classes of ramified polynomial coverings (S, $\pi$ ) of $\Sigma$ such that $\Delta_{\pi} \subset\{0,1, \infty\}$ can be naturally identified with the following sets:

- the quotient $\mathscr{B e l p o l}(\mathbb{C}) / \operatorname{APf}(\mathbb{C})$;
- the quotient $\mathscr{B}$ elpol $(\overline{\mathbb{Q}}) / \operatorname{APf}(\overline{\mathbb{Q}})$;
- the set $\mathscr{T}$ of isomorphism classes of bicoloured trees cyclically ordered at each vertex.

The Galois group Aut $(\overline{\mathbb{Q}})$ acts on $\overline{\mathbb{Q}}[\mathrm{X}]$ by $(g, f) \mapsto g_{*}(f)$ (action on the coefficients with $\left.g_{*}(\mathrm{X})=\mathrm{X}\right)$. The set $\mathscr{B} \boldsymbol{e} \boldsymbol{\ell}$ р $\boldsymbol{\ell}(\overline{\mathbb{Q}})$ is invariant; so is $\mathscr{A P P}(\overline{\mathbb{Q}})$; Hence Aut $(\overline{\mathbb{Q}})$ acts on the quotient $\mathscr{B} \boldsymbol{e l} \boldsymbol{p o l}(\overline{\mathbb{Q}}) / \mathscr{A} \mathcal{P f}(\overline{\mathbb{Q}})$ and on $\mathscr{T}$.

The following theorem states that this action is faithful.
Theorem (Lenstra-Schneps) For all nontrivial $g \in$ Aut $_{\mathbb{Q}}(\overline{\mathbb{Q}})$, there is a Belyi polynomial $f \in \mathscr{B} \boldsymbol{e l}_{\boldsymbol{p}} \boldsymbol{\ell}(\overline{\mathbb{Q}})$ such that $g_{*}(f)$ cannot be written as $f \circ \mathrm{~A}$, where $\mathrm{A} \in \operatorname{APf}(\overline{\mathbb{Q}})$.

The proof will be given in (7.4.6). It is taken from [3] and rests on a uniqueness property up to action of $\mathscr{A} \notin f$ of the factorization of polynomials by composition. It will be given in (7.4.5).

### 7.4.5 Comparing Factorizations

Proposition Let K be a subfield of $\mathbb{C}, f_{1}, f_{2}, g_{1}$ and $g_{2}$ polynomials of degrees $>0$ such that $g_{1} \circ f_{1}=g_{2} \circ f_{2}$. Suppose that $f_{1}$ and $f_{2}$ have the same degree $d$ with coefficients in K , and that $g_{1}$ and $g_{2}$ have the same degree $d^{\prime}$. Then there is an automorphism $\mathrm{A} \in \mathscr{A} \mathcal{P}(\mathrm{K})$ such that $f_{2}=\mathrm{A} \circ f_{1}$ and $g_{2}=g_{1} \circ \mathrm{~A}^{-1}$.

Remark Set $h=g_{1} \circ f_{1}=g_{2} \circ f_{2}$. As $\operatorname{deg}(h)=\operatorname{deg}\left(f_{1}\right) \cdot \operatorname{deg}\left(g_{1}\right)=\operatorname{deg}\left(f_{2}\right)$. $\operatorname{deg}\left(g_{2}\right)$, each of the assumptions $\operatorname{deg}\left(f_{1}\right)=\operatorname{deg}\left(f_{2}\right)$ and $\operatorname{deg}\left(g_{1}\right)=\operatorname{deg}\left(g_{2}\right)$ implies the other.

For any polynomial $f \in \mathbb{C}[\mathrm{X}]$, define the equivalence relation $\sim_{f}$ on $\mathbb{C}$ by $x \sim_{f}$ $y \Leftrightarrow f(x)=f(y)$. We will say it is defined by $f$.

Lemma Under the assumptions of the proposition, $f_{1}$ and $f_{2}$ define the same relation on $\mathbb{C}$.

Proof Set $h=g_{1} \circ f_{1}=g_{2} \circ f_{2}$; let $\mathrm{L}_{1}, \mathrm{~L}_{2}$ and M be the graphs of the equivalence relations defined respectively by $f_{1}, f_{2}$ and $h: \mathrm{M}=\left\{(x, y) \in \mathbb{C}^{2} \mid h(x)=h(y)\right\}$ and $\mathrm{L}_{j}=\left\{(x, y) \in \mathbb{C}^{2} \mid f_{j}(x)=f_{j}(y)\right\}$. Then $\mathrm{L}_{j} \subset \mathrm{M}$.

Let $\pi$ be the projection $(x, y) \mapsto x$ and $\Delta_{\mathrm{M}}$ the set of $x \in \mathbb{C}$ such that the polyno$\operatorname{mial} h(\mathrm{Y})-h(x) \in \mathbb{C}[\mathrm{Y}]$ has a multiple root. As $\Delta_{\mathrm{M}}=h^{-1}\left(\Delta_{h}\right)$, where $\Delta_{h}$ is the set of critical values of $h$, the set $\Delta_{\mathrm{M}}$ is finite. The space M equipped with $\pi$ induces a (non ramified) covering $\mathrm{M}^{\prime}$ of $\mathbb{C}-\Delta_{\mathrm{M}}$ of degree $d \cdot d^{\prime}$. Defining $\Delta_{\mathrm{L}_{j}}$ likewise, $\Delta_{\mathrm{L}_{j}} \subset \Delta_{\mathrm{M}}$. Set $\mathrm{L}_{j}^{\prime}=\mathrm{L}_{j}-\pi^{-1}\left(\Delta_{\mathrm{M}}\right)$. The spaces $\mathrm{L}_{1}^{\prime}$ and $\mathrm{L}_{2}^{\prime}$ are subcoverings of $\mathrm{M}^{\prime}$ over $\mathbb{C}-\Delta_{\mathrm{M}}$.

In the neighbourhood of infinity, the covering $\mathrm{M}^{\prime}$ is trivial: for sufficiently large R , there are $d \cdot d^{\prime}$ analytic sections $\tau_{k}: \mathbb{C}-\mathbb{D}_{\mathrm{R}} \rightarrow \mathrm{M}^{\prime}$, for $0 \leqslant k \leqslant d \cdot d^{\prime}-1$, such that the restriction of $\mathrm{M}^{\prime}$ to $\mathbb{C}-\mathbb{D}_{\mathrm{R}}$ is the union of their images, $\frac{\tau_{k}(x)}{x}$ tending to $e^{2 i \pi \frac{k}{d . d^{\prime}}}$ as $x$ tends to $\infty$. This follows from the fact that the part of highest degree in $h(\mathrm{Y})-h(\mathrm{X})$ is of the form $c \cdot\left(\mathrm{Y}^{d d^{\prime}}-\mathrm{X}^{d d^{\prime}}\right)$.

Similarly the restriction of $\mathrm{L}_{j}^{\prime}$ to $\mathbb{C}-\mathbb{D}_{\mathrm{R}}$ is the union of the images of $d$ sections $\sigma_{k}^{j}$ for $1 \leqslant k \leqslant d-1$, with $\frac{\sigma_{k}^{j}(x)}{x}$ tending to $e^{2 i \pi \frac{k}{d}}$ as $x$ tends to $\infty$. However $\mathrm{L}_{j}^{\prime}$ is a subcovering of $\mathrm{M}^{\prime}$. This forces $\sigma_{k}^{1}=\sigma_{k}^{2}=\tau_{k . d^{\prime}}$. So $\mathrm{L}_{1}^{\prime}$ and $\mathrm{L}_{2}^{\prime}$ agree over $\mathbb{C}-\mathbb{D}_{\mathrm{R}}$. As these are subcoverings of $\mathrm{M}^{\prime}$ over $\mathbb{C}-\Delta_{\mathrm{M}}$, the latter being connected, $\mathrm{L}_{1}^{\prime}=\mathrm{L}_{2}^{\prime}$. By (7.1.1, Proposition), $\mathrm{L}_{j}$ is the closure of $\mathrm{L}_{j}^{\prime}$. Thus $\mathrm{L}_{1}=\mathrm{L}_{2}$.

Proof of the Proposition A non constant polynomial $f: \mathbb{C} \rightarrow \mathbb{C}$ is an open map, and so defines a homeomorphism $f_{*}: \mathbb{C} / \sim_{f} \rightarrow \mathbb{C}$. The maps $\left(f_{1}\right)_{*}$ and $\left(f_{2}\right)_{*}$ are homeomorphisms of $\mathbb{C} / \sim_{f_{1}}=\mathbb{C} / \sim_{f_{2}}$ over $\mathbb{C}$, and $\mathrm{A}=\left(f_{2}\right)_{*} \circ\left(f_{1}\right)_{*}^{-1}$ is a homeomorphism $\mathbb{C} \rightarrow \mathbb{C}$ satisfying $f_{2}=\mathrm{A} \circ f_{1}$.

In the neighbourhood of a point $z$ which is not a critical value of $f_{1}, \mathrm{~A}=f_{2} \circ \psi$, where $\psi$ is a branch of $f_{1}^{-1}$. Hence A is holomorphic. The map A , being continuous and holomorphic on the complement of a finite set, is holomorphic. Similarly for $\mathrm{A}^{-1}$, and so A is a $\mathbb{C}$-analytic automorphism of $\mathbb{C}$; therefore $\mathrm{A} \in \mathscr{A} \mathcal{P f}(\mathbb{C})$.

Setting $\mathrm{A}(z)=a . z+b, f_{2}=a \cdot f_{1}+b$. Comparing the leading terms, $a \in \mathrm{~K}$, and $b=f_{2}-a \cdot f_{1} \in \mathrm{~K}$.

### 7.4.6 Proof of the Faithfulness Theorem

Let $\gamma \in \operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ be a nontrivial element. Then there exists $\alpha_{1} \in \overline{\mathbb{Q}}$ such that $\alpha_{2}=$ $\gamma\left(\alpha_{1}\right) \neq \alpha_{1}$.

Let $f_{1} \in \overline{\mathbb{Q}}[\mathrm{X}]$ be a polynomial whose derivative $f_{1}^{\prime}$ has 0,1 and $\alpha_{1}$ as roots with distinct multiplicities $m_{0}, m_{1}, m_{\alpha}$. The polynomial $f_{2}=\gamma_{*}\left(f_{1}\right)$ has a derivative $f_{2}^{\prime}$ whose roots are 0,1 and $\alpha_{2}$ with the same multiplicities $m_{0}, m_{1}, m_{\alpha}$.

Now, $\Delta\left(f_{1}\right) \subset \overline{\mathbb{Q}} \cup\{\infty\}$. By (7.2.8, Corollary), there exists $g \in \mathbb{Q}[\mathrm{X}]$ such that $\Delta\left(g \circ f_{1}\right)=g_{!}\left(\Delta\left(f_{1}\right)\right) \subset\{0,1, \infty\}$. Then $h_{1}=g \circ f_{1} \in \mathscr{B} \boldsymbol{e \ell} \boldsymbol{p} \circ \boldsymbol{\ell}(\overline{\mathbb{Q}})$. So $\Delta\left(f_{2}\right)=$ $\gamma\left(\Delta\left(f_{1}\right)\right)$, and $h_{2}=g \circ f_{2} \in \mathscr{B}$ el р。人 $(\overline{\mathbb{Q}})$.

The faithfulness theorem then follows from the next lemma:
Lemma The polynomial $h_{2}$ cannot be written as $h_{1} \circ \mathrm{~A}$ with $\mathrm{A} \in \operatorname{APf}(\mathbb{C})$.
Proof Suppose that $h_{2}=h_{1} \circ \mathrm{~A}$ with $\mathrm{A} \in \operatorname{APf}(\mathbb{C})$, i.e. $g \circ f_{2}=g \circ f_{1} \circ \mathrm{~A}$. By (7.4.5, Proposition), there exists $\mathrm{B} \in \mathscr{A} \mathcal{P}(\mathbb{C})$ such that $f_{2}=\mathrm{B} \circ f_{1} \circ \mathrm{~A}$ and $g=$ $g \circ \mathrm{~B}^{-1}$.

Write $\mathrm{A}(z)=a \cdot z+a_{0}$ and $\mathrm{B}(z)=b \cdot z+b_{0}$. Then $f_{2}=\mathrm{B} \circ f_{1} \circ \mathrm{~A}$ has derivative $f_{2}^{\prime}=a b\left(f_{1}^{\prime} \circ \mathrm{A}\right)$. The roots of $f_{1}^{\prime}$ are the images under A of the roots of $f_{2}^{\prime}$ with their multiplicities preserved. Hence, $\mathrm{A}(0)=0, \mathrm{~A}(1)=1, \mathrm{~A}\left(\alpha_{2}\right)=\alpha_{1}$. The first two equalities imply that A is the identity, and then $\alpha_{1}=\alpha_{2}$ by the last one, giving a contradiction.

This completes the Proof of Theorem 7.4.4.

### 7.5 Two Examples

In this section we give two examples:

- an example of two Belyi polynomials giving non-isomorphic trees, regardless of the colouring and of the cyclic order at vertices, but which are conjugate under the action of the group $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$;
- an example, taken from Zapponi [4], of two Belyi polynomials giving isomorphic bicoloured trees with distinct cyclic order at each vertex, and which are not conjugate under the action of $\mathbb{G}$.


### 7.5.1 First Example: Setting of the Problem

Consider the following bicoloured trees


They have the same valency function: they both have two white vertices of valency 2 , a white vertex of valency 3 , three black vertices of valency 1 and two black vertices of valency 2 . Up to isomorphism, these are the only bicoloured trees with this valency function, and this remains true if the cyclic order at each vertex is taken account of.

Are they in the same orbit under the action of $\mathbb{G}=$ Aut $_{\mathbb{Q}}(\overline{\mathbb{Q}})$ ? In other words, is there some $g \in \mathbb{G}$ such that $g^{*}\left(\mathrm{~T}_{1}\right)=\mathrm{T}_{2}$ ? Proposition 7.5 .4 will provide the answer.

### 7.5.2 The Polynomial $\mathbf{Q}_{c}$

The bicoloured trees $T_{1}$ and $T_{2}$ correspond to the Belyi polynomials $P_{1}$ and $P_{2}$ of degree 7 , with one triple root and two double ones, unique up to pre-composition with an affine map. Thus the origin can be taken to be the triple root, and then such a polynomial can be written as $\mathrm{P}=a z^{3}\left(z^{2}-2 b z+c\right)^{2}$.

For $\mathrm{P}_{1}$ and $\mathrm{P}_{2}, b \neq 0$. Otherwise P would be an odd polynomial and 1 would be a critical value only if it so was -1 . Hence we may assume that $b=1$ since the action of a homothety on the variable can be expressed as $b \mapsto \lambda b, c \mapsto \lambda^{2} c$.

Let $\mathrm{Q}_{c}$ be the polynomial $z^{3}\left(z^{2}-2 z+c\right)^{2}$. Let us find the values of $c$ for which $\mathrm{Q}_{c}$ has two nonzero equal critical values, given by two distinct critical points. For such a value $v$ of $c, \mathrm{P}=\frac{1}{v} \mathrm{Q}_{c}$ with equal to one of $\mathrm{P}_{1}, \mathrm{P}_{2}$; and $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$ can be obtained in this manner.

### 7.5.3 Study of $\mathbf{Q}=\mathbf{Q}_{c}$

Derivative. $\mathrm{Q}^{\prime}(z)=z^{2}\left(z^{2}-2 z+c\right) \mathrm{G}(z)$, where $\mathrm{G}(z)=7 z^{2}-10 z+3 c$.
Critical points (other than the roots). They are given by the equation $\mathrm{G}(\alpha)=0$. The reduced discriminant is $\Delta=25-21 c$. Set $\delta=\sqrt{\Delta}$. Then $c=\frac{1}{21}\left(25-\delta^{2}\right)$; the critical points other than the roots are $\alpha=\frac{1}{7}(5+\delta)$ and $\alpha^{\prime}=\frac{1}{7}(5-\delta)$.

We express the critical value N corresponding to $\alpha$ in terms of $\delta$, that corresponding to $\alpha^{\prime}$ will follow by replacing $\delta$ by $-\delta$ :
$\alpha^{2}=\frac{1}{49}\left(25+10 \delta+\delta^{2}\right) ;$
$\alpha^{3}=\frac{1}{343}\left(125+75 \delta+15 \delta^{2}+\delta^{3}\right)$;
$\alpha^{2}-2 \alpha+c=\frac{4}{147}\left(10-3 \delta-\delta^{2}\right)$;
$\left(\alpha^{2}-2 \alpha+c\right)^{2}=\frac{16}{9.7^{4}}\left(100-60 \delta-11 \delta^{2}+6 \delta^{3}+\delta^{4}\right)$.
Non-zero critical values.

$$
\begin{aligned}
& \mathrm{Q}(\alpha)=\alpha^{3}\left(\alpha^{2}-2 \alpha+c\right)^{2} \\
& \quad=\frac{16}{9.7^{7}}\left(12500-4375 \delta^{2}-875 \delta^{3}+350 \delta^{4}-154 \delta^{5}+21 \delta^{6}+\delta^{7}\right)
\end{aligned}
$$

Difference of these critical values.
$\mathrm{Q}(\alpha)-\mathrm{Q}\left(\alpha^{\prime}\right)=\frac{32}{9.7^{7}} \delta^{3}\left(-875-154 \Delta+\Delta^{2}\right)$.
The values of $c$ for which $\mathrm{Q}_{c}$ has two nonzero equal critical values correspond to the values of $\delta$ for which this expression vanishes.

### 7.5.4 Interpretation of These Results

The value $\delta=0$ gives $c=\frac{25}{21}$. For this value, $\mathrm{Q}_{c}$ has a double critical point (hence of local degree 3) with nonzero critical value. The Belyi polynomial $\mathrm{P}=\frac{1}{12500} \mathrm{Q}_{c}$ gives the tree


It is the only bicoloured tree with this valency function.
The polynomial $R(\Delta)=\Delta^{2}-154 \Delta-875$ has two real roots $\Delta^{+}>0$ and $\Delta^{-}<$ 0 . In both cases, $c=\frac{1}{21}(25-\Delta) \in \mathbb{R}$, the coefficients of $\mathrm{Q}_{c}$ and P are real, and the tree realized in $\mathbb{C}$ is symmetric with respect to $\mathbb{R}$. The root $\Delta^{+}$gives two nonzero real values for $\delta$, hence two distinct real values for the critical point $\alpha$, i.e. two white points on the real axis; it corresponds to the tree $\mathrm{T}_{1}$. The root $\Delta^{-}$gives two distinct complex conjugate values for $\delta$, hence also for $\alpha$; it corresponds to the tree $\mathrm{T}_{2}$.

The reduced discriminant of R is $77^{2}+875=6804=18^{2} \cdot 21$. It follows that
Proposition An element $g \in \mathbb{G}=$ Aut $_{\mathbb{Q}}(\overline{\mathbb{Q}})$ exchanges the trees $\mathrm{T}_{1}$ and $\mathrm{T}_{2}$ if and only if it exchanges the two square roots of 21 .

### 7.5.5 Second Example: "Leila's bouquet"

Consider the bicoloured tree L

with 1 black pentavalent vertex $\mathrm{O}, 15$ monovalent black vertices, and 5 white vertices $\mathrm{S}_{1}, \ldots, \mathrm{~S}_{5}, \mathrm{~S}_{i}$ being of valency $d_{i}=i+1$. Up to isomorphism, L is the only bicoloured
tree with this valency function. Let $a_{i}$ denote the edge from O to $\mathrm{S}_{i}$. For any cyclic order $\omega$ on $\{1, \ldots, 5\}$, let $\mathrm{L}_{\omega}$ be an objet of $\mathscr{H}$ obtained by equipping L with the order $\omega$ on the edges $a_{i}$ (the cyclic order at white vertices is unimportant). This gives, up to isomorphism, all the objects of $\mathscr{H}$ with this valency function. A cyclic order will be defined by a bijection $\varphi:\{1, \ldots, 5\} \rightarrow\{1, \ldots, 5\}$.

Are all the objects $L_{\omega}$ in the same orbit under the action of $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$ ? The corollary of the above Theorem 7.5.8 shows that this is not the case. The Proof of Theorem 7.5.8 will be completed in 7.5.16; to state this theorem, we need to define two rationals $\delta(\mathrm{P})=\delta_{\boldsymbol{a} \boldsymbol{Z}}(\mathrm{P})$ and $\delta_{\text {top }}(\mathrm{P})$, and these definitions rest on the calculation given in 7.5.6.

The figure below shows two examples of objects $L_{\omega}$.

$\varphi(i)=i$

$\varphi(1,2,3,4,5)=(2,3,1,5,4)$

### 7.5.6 Computing Discriminants

Let P be a degree 20 Belyi polynomial such that $\mathrm{P}^{-1}([0,1])$ is a bicoloured tree isomorphic to L , in other words such that $\theta(\mathrm{P})$ is an object of $\mathscr{H}$ isomorphic to one of $\mathrm{L}_{\omega}$. For each $\omega$, there is a unique such polynomial P up to pre-composition with an affine transformation. Hence $P$ may be chosen to be monic with the black pentavalent vertex O at 0 .

Letting $x_{i}$ be the root of P corresponding to the white vertex $\mathrm{S}_{i}$,

$$
\mathrm{P}(\mathrm{Z})=\prod_{i=1}^{5}\left(\mathrm{Z}-x_{i}\right)^{d_{i}} \quad \text { where } \quad d_{i}=i+1
$$

Here we use the "algebraic" numbering of $x_{i}$ given by their multiplicity as roots. This numbering provides a cyclic order which may be different from the "topologi-
cal" cyclic order induced by the extension of the tree in

$\mathbb{C}$.
Let Q be the monic polynomial having $x_{i}$ as simple roots:

$$
\mathrm{Q}(\mathrm{Z})=\prod_{i=1}^{5}\left(\mathrm{Z}-x_{i}\right)
$$

Let $\Delta$ be the discriminant of Q :

$$
\Delta=\prod_{i, j \mid i \neq j}\left(x_{i}-x_{j}\right)=\prod_{i, j \mid i<j}\left(x_{i}-x_{j}\right)^{2}
$$

(note that there are 10 pairs $(i, j)$ such that $i<j$ and that 10 is even. This is why the latter equality has $a+$ sign).

## Proposition

$$
\Delta=\left(\frac{200}{3}\right)^{2} \cdot \prod_{i} x_{i}^{4}
$$

## Lemma

$$
\frac{\mathrm{P}^{\prime}(\mathrm{Z})}{\mathrm{P}(\mathrm{Z})}=\sum_{i} \frac{d_{i}}{\mathrm{Z}-x_{i}}=20 \frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}
$$

Proof (a) As $\mathrm{P}(\mathrm{Z})=\prod\left(\mathrm{Z}-x_{i}\right)^{d_{i}}, \frac{\mathrm{P}^{\prime}}{\mathrm{P}}=\log ^{\prime}(\mathrm{P})=\sum \frac{d_{i}}{\mathrm{Z}-x_{i}}$.
(b) The only zero of the rational fraction $\frac{\mathrm{P}^{\prime}}{\mathrm{P}}$ in $\mathbb{C}$ is one of multiplicity 4 at the point 0 , and the only pole is a simple at each $x_{i}$; hence $\frac{\mathrm{P}^{\prime}}{\mathrm{P}}$ is of the form $k \cdot \frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}$. As Z tends to infinity, $\frac{\mathrm{P}^{\prime}}{\mathrm{P}}$ is equivalent to $\frac{20}{\mathrm{Z}}$, and soù $k=20$.

## Proof of the Proposition

$$
20 \mathrm{Z}^{4}=\mathrm{Q} \cdot \frac{\mathrm{P}^{\prime}}{\mathrm{P}}=\sum_{i} d_{i} \cdot \frac{\mathrm{Q}}{\mathrm{Z}-x_{i}}=\sum_{i} d_{i} \prod_{j \mid j \neq i}\left(\mathrm{Z}-x_{j}\right)
$$

All terms of the sum vanish at $\mathrm{Z}=x_{i_{0}}$ except the one corresponding to $i=i_{0}$; so

$$
20 x_{i}^{4}=d_{i} \prod_{j \mid j \neq i}\left(x_{i}-x_{j}\right)
$$

Taking the product over the different values of $i$ gives:

$$
20^{5} \cdot \prod_{i} x_{i}^{4}=\left(\prod_{i} d_{i}\right) \cdot \Delta
$$

Hence, as $\prod d_{i}=6!=720=2^{4} \cdot 3^{2} \cdot 5$,

$$
\frac{\Delta}{\prod x_{i}^{4}}=\frac{2^{10} 5^{5}}{2^{4} 3^{2} 5}=\frac{2^{6} 5^{4}}{3^{2}}=\left(\frac{2^{3} 5^{2}}{3}\right)^{2}=\left(\frac{200}{3}\right)^{2}
$$

### 7.5.7 The Invariant $\delta_{\text {alg }}$

If P is a Belyi polynomial satisfying the conditions of 7.5.6, define $\delta(\mathrm{P})=\delta_{\text {alg }}(\mathrm{P})$ by

$$
\prod_{i, j \mid i<j}\left(x_{i}-x_{j}\right)=\delta(\mathrm{P}) \cdot \prod_{i} x_{i}^{2}
$$

It follows from Proposition 7.5.6. that $\delta_{\boldsymbol{a} \ell g}(\mathrm{P})= \pm \frac{200}{3}$. In particular $\delta_{\boldsymbol{a} \ell g}(\mathrm{P})$ is rational.

Proposition The rational $\delta_{\text {alg }}(\mathrm{P})$ only depends on the orbit of the isomorphism class of the object $\theta(\mathrm{P})$ of $\mathscr{H}$ under the action of $\mathbb{G}=\operatorname{Aut}_{\mathbb{Q}}(\overline{\mathbb{Q}})$.

Proof Let $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$ be two polynomials satisfying the conditions of (7.5.6). Let $x_{i}^{(1)}$ and $x_{i}^{(2)}$ be the respective roots of $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$ of multiplicity $i+1$.
(a) If the objects $\theta\left(\mathrm{P}_{1}\right)$ and $\theta\left(\mathrm{P}_{2}\right)$ of $\mathscr{H}$ are isomorphic, then there exists $\lambda$ such that $\mathrm{P}_{2}(\mathrm{Z})=\mathrm{P}_{1}(\lambda \mathrm{Z})$. So $x_{i}^{(1)}=\lambda x_{i}^{(2)}$, and hence

$$
\begin{gathered}
\prod_{i<j}\left(x_{i}^{(1)}-x_{j}^{(2)}\right)=\lambda^{10} \prod_{i<j}\left(x_{i}^{(2)}-x_{j}^{(2)}\right) \\
\prod\left(x_{i}^{(1)}\right)^{2}=\lambda^{10} \prod\left(x_{i}^{(2)}\right)^{2} .
\end{gathered}
$$

Finally, $\delta_{a t g}\left(\mathrm{P}_{1}\right)=\delta_{\boldsymbol{a} \boldsymbol{g}}\left(\mathrm{P}_{2}\right)$.
(b) If there is an automorphism $g \in \mathbb{G}$ such that $\mathrm{P}_{2}=g_{*}\left(\mathrm{P}_{1}\right)$, then $x_{i}^{(2)}=g\left(x_{i}^{(1)}\right)$. So $\delta\left(\mathrm{P}_{2}\right)=g\left(\delta\left(\mathrm{P}_{1}\right)\right)$. As $\delta\left(\mathrm{P}_{1}\right) \in \mathbb{Q}$, this implies $\delta_{\boldsymbol{a} \ell g}\left(\mathrm{P}_{2}\right)=\delta_{\boldsymbol{a} \ell g}\left(\mathrm{P}_{1}\right)$.

### 7.5.8 Definition of $\delta_{\text {top }}(\mathbf{P})$

Under the assumptions of (7.5.6), for a permutation $\varphi$ of $\{1, \ldots, 5\}$, define $\delta_{\varphi}(\mathrm{P})$ by

$$
\prod_{i, j \mid \varphi(i)<\varphi(j)}\left(x_{i}-x_{j}\right)=\delta_{\varphi}(\mathrm{P}) \cdot \prod_{i} x_{i}^{2}
$$

Then $\delta_{\varphi}(\mathrm{P})=\varepsilon(\varphi) \cdot \delta(\mathrm{P})$, where $\varepsilon(\varphi)$ is the signature of $\varphi$. In particular, $\delta_{\varphi}(\mathrm{P})$ only depends on the cyclic order defined by $\varphi$ on the set of $x_{i}$ (note that 5 is odd!). We write $\delta_{\xi}(\mathrm{P})=\delta_{\varphi}(\mathrm{P})$ if $\xi$ is the cyclic order defined by $\varphi$.

Set $\delta_{\text {top }}(\mathrm{P})=\delta_{\omega}(\mathrm{P})$, where $\omega$ is the cyclic order on all $a_{i}$ such that $\theta(\mathrm{P})$ is isomorphic to $\mathrm{L}_{\omega}$, i.e. the cyclic order of the graph $\mathrm{H}=\mathrm{P}^{-1}([0,1])$ at 0 induced by its inclusion in $\mathbb{C}$.

Theorem (Zapponi) For any Belyi polynomial P satisfying the conditions of (7.5.6),

$$
\delta_{\text {top }}(\mathrm{P})=-\frac{200}{3}
$$

Corollary Let $\omega_{1}$ and $\omega_{2}$ be two cyclic orders on $\{1, \ldots, 5\}$, and $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$ the Belyi polynomials satisfying conditions of 7.5 .6 and realizing $\mathrm{L}_{\omega_{1}}$ and $\mathrm{L}_{\omega_{2}}$ respectively. Let $\varphi_{1}$ and $\varphi_{2}$ be bijections from $\{1, \ldots, 5\}$ onto $\mathbb{Z} /(5)$ defining $\omega_{1}$ and $\omega_{2}$. Set $\sigma=$ $\varphi_{2} \circ \varphi_{1}^{-1}$ (so that $\sigma$ transforms $\omega_{1}$ into $\omega_{2}$ ).

If $\varepsilon(\sigma)=-1$, then $\delta_{\text {alg }}\left(\mathrm{P}_{2}\right)=-\delta_{\boldsymbol{a t g}}\left(\mathrm{P}_{1}\right)$, and there is no $g \in \mathbb{G}$ transforming $\mathrm{P}_{1}$ into $\mathrm{P}_{2}$.

The rest of this section is dedicated to the proof of the theorem. It will be completed in (7.5.16). If A is an affine transformation, then $\delta_{\text {top }}(\mathrm{P})=\delta_{\text {top }}(\mathrm{P} \circ \mathrm{A})$. The focus will be on Q (equipped with a numbering of its roots) and not on P . We construct a connected space $\mathscr{B}^{+}$such that, for all P , there is an affine A such that the polynomial Q corresponding to $\mathrm{P} \circ \mathrm{A}$ is in $\mathscr{B}^{+}$; we also construct $\theta_{\text {top }}: \mathscr{B}^{+} \rightarrow \mathbb{R}^{*}$ such that $\delta_{\text {topp }}(\mathrm{P})=\theta_{\text {top }}(\mathrm{Q})$. The function $\theta_{\text {topp }}$, being continuous, its sign remains invariant; this is why the sign of $\delta_{\text {top }}(\mathrm{P})$ does not depend on P .

If P is a Belyi polynomial satisfying the conditions of (7.5.6), then

$$
\frac{1}{20} \frac{\mathrm{P}^{\prime}(\mathrm{Z})}{\mathrm{P}(\mathrm{Z})}=\frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}=\sum \frac{r_{i}}{\mathrm{Z}-x_{i}}
$$

with $r_{i}=\frac{1}{20} \cdot d_{i}$. We investigate the most general situation where there is a monic polynomial Q of degree 5 , such that $f_{\mathrm{Q}}(\mathrm{Z})=\frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}$ can be written as $\sum \frac{r_{i}}{\mathrm{Z}-x_{i}}$, with $r_{i} \in \mathbb{R}_{+}^{*}$.

### 7.5.9 Notation

Let $\mathscr{A}$ be the set of monic polynomials Q of degree 5 with 5 distinct roots, and $\widetilde{\mathscr{A}}$ the set of $\widetilde{\mathbb{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right)$ such that $\mathrm{Q} \in \mathscr{A}$ and $\mathrm{Q}^{-1}(0)=\left\{x_{1}, \ldots, x_{5}\right\}$. The sets $\mathscr{A}$ and $\tilde{\mathscr{A}}$ are $\mathbb{C}$-analytic submanifolds of dimension 5 in $\mathbb{C}^{6}$ and $\mathbb{C}^{11}$ respectively, and the projection $\varpi:\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \mapsto \mathrm{Q}$ makes $\tilde{\mathscr{A}}$ a covering of degree $5!=120$ of $\mathscr{A}$ (there are 5 ! ways to number the roots).

For $\mathrm{Q} \in \mathscr{A}$, set $f_{\mathrm{Q}}(\mathrm{Z})=\frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z}}$. The rational fraction $f_{\mathrm{Q}}$ has simple poles at the zeros $x_{i}$ of Q , and it is of the form $\frac{1}{\mathrm{Z}}+\mathrm{O}\left(\frac{1}{\mathrm{Z}^{2}}\right)$ as $\mathrm{Z} \rightarrow \infty$. Hence

$$
f_{\mathrm{Q}}(\mathrm{Z})=\sum \frac{r_{i}}{\mathrm{Z}-x_{i}},
$$

where $r_{i}$ are given by

$$
r_{i}=r_{i}(\widetilde{\mathrm{Q}})=\frac{x_{i}^{4}}{\mathrm{Q}^{\prime}\left(x_{i}\right)}=\frac{x_{i}^{4}}{\prod_{j \mid j \neq i}\left(x_{i}-x_{j}\right)}
$$

They are nonzero and $\sum r_{i}=1$. This defines a map

$$
\rho: \widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \mapsto\left(r_{1}, \ldots, r_{5}\right)
$$

from $\tilde{\mathscr{A}}$ to the subset $\mathscr{W}$ of $\mathbb{C}^{5}$ consisting of $r=\left(r_{i}\right)$ such that $(\forall i) r_{i} \neq 0$ and $\sum r_{i}=$ 1. The set $\mathscr{W}$ is a submanifold of dimension 4 in $\mathbb{C}^{5}$, and the map $\rho$ is analytic.

Let $\mathscr{I}$ be the set of $r=\left(r_{i}\right) \in \mathscr{W}$ such that there is a nonempty subset J of $\{1, \ldots, 5\}$ with $\sum_{i \in \mathrm{~J}} r_{i}=0$, and $\mathscr{W}^{*}=\mathscr{W}-\mathscr{Z}$. Set $\widetilde{\mathscr{A}}^{*}=\rho^{-1}\left(\widetilde{\mathscr{W}}^{*}\right) \subset \widetilde{\mathscr{A}}, \mathscr{A}^{*}=\varpi\left(\widetilde{\mathscr{A}}^{*}\right)$, and $\mathscr{W}^{+}=\left\{r=\left(r_{i}\right) \in \mathscr{W} \mid(\forall i) r_{i} \in \mathbb{R}_{+}^{*}\right\}$. Let $\mathscr{W}^{+} \subset \mathscr{W}^{*}$. Set $\tilde{\mathscr{A}}^{+}=\rho^{-1}\left(\mathscr{W}^{+}\right)$and $\mathscr{A}^{+}=\varpi\left(\tilde{\mathscr{A}}^{+}\right)$.

The polynomial Q is determined by the $x_{i}$. Multiplying the $x_{i}$ by a same factor $\lambda \neq 0$, changes the polynomial Q to $\lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right)$, the fraction $f_{\mathrm{Q}}$ to $\frac{1}{\lambda} f_{\mathrm{Q}}\left(\frac{\mathrm{Z}}{\lambda}\right)$, and leaves the $r_{i}$ invariant.

Let $\mathscr{B}$ be the set of $\mathrm{Q} \in \mathscr{A}$ such that $\mathrm{Q}(0)=-1$, i.e. $\Pi x_{i}=1$. Set $\widetilde{\mathscr{B}}=\varpi^{-1}(\mathscr{B})$, $\widetilde{\mathscr{B}}^{*}=\widetilde{\mathscr{A}}^{*} \cap \widetilde{\mathscr{B}}, \widetilde{\mathscr{B}}^{+}=\widetilde{\mathscr{A}}^{+} \cap \widetilde{\mathscr{B}}, \mathscr{B}^{*}=\mathscr{A}^{*} \cap \mathscr{B}, \mathscr{B}^{+}=\mathscr{A}^{+} \cap \mathscr{B}$.

The map $\left(\lambda, \mathrm{Q},\left(x_{i}\right)\right) \mapsto\left(\lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right),\left(\lambda x_{i}\right)\right)$ makes $\mathbb{C}^{*} \times \widetilde{\mathscr{B}} \mathrm{a} 5$-fold covering of $\tilde{\mathscr{A}}$.

### 7.5.10 A Covering Property

Proposition The map $\rho$ makes $\widetilde{\mathscr{B}}^{*}$ into a finite covering of $\mathscr{W}^{*}$.
Corollary The map $\rho$ makes $\widetilde{\mathscr{B}}^{+}$a trivial finite covering of $\mathscr{W}^{+}$.
Note that $\mathscr{W}^{+}=\{r \in] 0,1\left[^{5} \mid \sum r_{i}=1\right\}$ is contractible.
The proposition then follows by (4.3.12, Proposition) and by Lemmas (7.5.11) and (7.5.12) given below.

### 7.5.11 The Tangent Linear Map $\rho$

Lemma For all $\widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \in \widetilde{\mathscr{B}}$, the tangent linear map $\mathrm{T}_{\widetilde{\mathrm{Q}}} \rho$ induces an isomorphism $\mathrm{T}_{\widetilde{\mathrm{Q}}} \widetilde{\mathscr{B}} \rightarrow \mathrm{T}_{r} \mathscr{W}$.

Proof (a) Differentiation of $\mathrm{Q} \mapsto f_{\mathrm{Q}}$ on $\mathscr{A}$.
Let $\mathrm{Q} \in \mathscr{A}$ and consider a family $\left(\mathrm{Q}_{\varepsilon}\right)_{\varepsilon \in \mathrm{V}}$, where $\mathrm{Q}_{\varepsilon}=\mathrm{Q}+\varepsilon \mathrm{H}$, with H a polynomial of degree $\leqslant 4$, and V a neighbourhood of 0 in $\mathbb{C}$, which will be assumed to be connected and which we will allow ourselves to contract.

For $z \in \mathbb{C}$ such that $\mathrm{Q}(z) \neq 0$, there is a finite expansion

$$
f_{\mathrm{Q}_{\varepsilon}}(z)=f_{\mathrm{Q}}(z)-\frac{z^{4} \mathrm{H}(z)}{\mathrm{Q}(z)^{2}} \varepsilon+\mathrm{O}\left(\varepsilon^{2}\right) .
$$

It holds uniformly as $z$ runs through a compact subset of $\mathbb{C}-\mathrm{Q}^{-1}(0)$.
(b) Differentiation of $\rho$ on $\tilde{\mathscr{A}}$.

Let $\widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \in \tilde{\mathscr{A}}$. By the implicit functions theorem, there are holomorphic functions $\xi_{i}$ on V such that $\xi_{i}(0)=x_{i}$ and $\mathrm{Q}_{\varepsilon}\left(\xi_{i}(\varepsilon)\right)=0$. Set $\widetilde{\mathrm{Q}}_{\varepsilon}=$ $\left(\mathrm{Q}_{\varepsilon}, \xi_{1}(\varepsilon), \ldots, \xi_{5}(\varepsilon)\right)$. Since

$$
\rho_{i}\left(\widetilde{\mathrm{Q}}_{\varepsilon}\right)=\frac{1}{2 i \pi} \int_{\gamma_{i}} \mathrm{Q}_{\varepsilon}(z) d z
$$

where $\gamma_{i}$ is a loop only going round the roots $x_{i}$, and doing so only once. It follows that

$$
\rho_{i}\left(\widetilde{\mathrm{Q}}_{\varepsilon}\right)=\rho_{i}\left(\mathrm{Q}_{\varepsilon}\right)-\operatorname{Residue}_{x_{i}}\left(\frac{\mathrm{Z}^{4} \mathrm{H}}{\mathrm{Q}^{2}}\right) \cdot \varepsilon+\mathrm{O}\left(\varepsilon^{2}\right)
$$

Therefore the tangent linear map to $\rho$ is

$$
\mathrm{T}_{\widetilde{\mathrm{Q}}} \rho: \mathrm{H} \mapsto\left(-\operatorname{Residue}_{x_{i}}\left(\frac{\mathrm{Z}^{4} \mathrm{H}}{\mathrm{Q}^{2}}\right)\right)_{i \in\{1, \ldots, 5\}}
$$

(c) Injectivity of $\mathrm{T}_{\widetilde{\mathrm{Q}}} \rho$ on $\mathrm{T}_{\widetilde{\mathrm{Q}}} \widetilde{\mathscr{B}}$.

The space $\mathrm{T}_{\widetilde{\mathbb{Q}}} \widetilde{\mathscr{B}}=\mathrm{T}_{\mathrm{Q}} \mathscr{B}$ consists of all $\mathrm{H} \in \mathbb{C}[Z]$ of degree $\leqslant 4$ such that $\mathrm{H}(0)=$ 0 . Let $\mathrm{H} \in \mathrm{T}_{\widetilde{\mathrm{Q}}} \widetilde{\mathscr{B}}$ and suppose that

$$
\text { ( } \forall i) \quad \operatorname{Residue}_{x_{i}}\left(\frac{\mathrm{Z}^{4} \mathrm{H}}{\mathrm{Q}^{2}}\right)=0
$$

Then $\frac{\mathrm{Z}^{4} \mathrm{H}}{\mathrm{Q}^{2}}$ is of the form $\sum \frac{a_{i}}{\left(\mathrm{Z}-x_{i}\right)^{2}}=\mathrm{R}^{\prime}$, where $\mathrm{R}=\sum \frac{-a_{i}}{\mathrm{Z}-x_{i}}$.
The rational fraction $\mathrm{R}_{0} \stackrel{\left(2-x_{i}\right.}{=}-\mathrm{R}(0)$ is bounded at infinity and has simple poles only at $x_{i}$; hence it is of the form $\frac{\mathrm{G}}{\mathrm{Q}}$, where G is a polynomial of degree $\leqslant 5$. As $R(0)=0$ and $R^{\prime}=\frac{Z^{4} H}{Q^{2}}$ has a zero of order $\geqslant 4$ at $0, R_{0}$ has a zero of order $\geqslant 5$ at 0 . Therefore $\mathrm{R}_{0}$ is of the form $\lambda \frac{\mathrm{Z}^{5}}{\mathrm{Q}(\mathrm{Z})}$, and

$$
\frac{\mathrm{Z}^{4} \mathrm{H}}{\mathrm{Q}^{2}}=\mathrm{R}_{0}^{\prime}=\lambda \frac{5 \mathrm{Z}^{4} \mathrm{Q}-\mathrm{Z}^{5} \mathrm{Q}^{\prime}}{\mathrm{Q}^{2}}
$$

so $\mathrm{H}=\lambda\left(5 \mathrm{Q}-\mathrm{ZQ}^{\prime}\right)$ and $\mathrm{H}(0)=-5 \lambda$. Hence as $\mathrm{H}(0)=0, \lambda=0$ and $\mathrm{H}=0$.
Therefore the map $\mathrm{T}_{\widetilde{\mathrm{Q}}} \widetilde{\mathscr{B}} \rightarrow \mathrm{T}_{r} \mathscr{W}$ induced by $\mathrm{T}_{\widetilde{\mathrm{Q}}} \rho$ is injective. The initial and end spaces having dimension 4 , it is an isomorphism.

### 7.5.12 A Properness Property

Lemma The map $\rho: \widetilde{\mathscr{B}}^{*} \rightarrow \mathscr{W}^{*}$ is proper.

Proof For $\widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \in \widetilde{\mathscr{A}}$ set $\mathrm{R}(\mathrm{Q})=\sup \left(\left|x_{i}\right|\right)$ and $\mathrm{N}(\widetilde{\mathrm{Q}})=(\mathrm{N}(\mathrm{Q})$, $\left.\mathrm{Y}_{1}, \ldots, \mathrm{Y}_{5}\right)$, where $\mathrm{Y}_{i}=\frac{x_{i}}{\mathrm{R}(\mathrm{Q})}$ and $\mathrm{N}(\mathrm{Q})=\prod\left(\mathrm{Z}-\mathrm{Y}_{i}\right)$. Then $r_{i}=$ Residue $_{x_{i}} f_{\mathrm{Q}}=$ Residue $_{\mathrm{Y}_{i}} f_{\mathrm{N}(\mathrm{Q})}$, where $f_{\mathrm{N}(\mathrm{Q})}=\frac{\mathrm{Z}^{4}}{\mathrm{~N}(\mathrm{Q})}$.

Let $\left(\widetilde{\mathrm{Q}}_{n}\right)=\left(\left(\mathrm{Q}_{n}, x_{1}^{(n)}, \ldots, x_{5}^{(n)}\right)\right)$ be a sequence in $\widetilde{\mathscr{B}}$ such that $\rho\left(\widetilde{\mathrm{Q}}_{n}\right)$ tends to a point $r=\left(r_{i}\right) \in \mathscr{W}^{*}$. If necessary extracting a subsequence, $x_{i}^{(n)}$ may be supposed to converge to a limit $x_{i} \in \Sigma=\mathbb{C} \cup\{\infty\}$. We need to show that all $x_{i}$ are in $\mathbb{C}^{*}$ and are distinct.

If all $x_{i}$ are finite, then $\prod x_{i}=1$, and so $x_{i} \in \mathbb{C}^{*}$ for all $i$, and

$$
(\forall i) \quad \prod_{j \mid j \neq i}\left(x_{i}-x_{j}\right)=\frac{x_{i}^{4}}{r_{i}} \neq 0 \text {; }
$$

hence the $x_{i}$ are distinct.
Suppose that some $x_{i}$ equals $\infty$. Then $\mathrm{R}_{n}=\mathrm{R}\left(\mathrm{Q}_{n}\right) \rightarrow \infty$, and $\mathrm{Y}_{i}^{(n)}=\frac{x_{i}^{(n)}}{\mathrm{R}_{n}}$ satisfies $\left|\mathrm{Y}_{i}^{(n)}\right| \leqslant 1$. If necessary extracting a subsequence, $\mathrm{Y}_{i}^{(n)}$ may be supposed to converge to a limit $\mathrm{Y}_{i} \in \mathbb{C} . \mathrm{As} \prod \mathrm{Y}_{i}=\lim \mathrm{R}_{n}^{-5}=0$, the set J of $i$ such that $\mathrm{Y}_{i}=0$ is nonempty, and its cardinal $k$ satisfies $0<k<5$.

The polynomial $\mathrm{N}\left(\mathrm{Q}_{n}\right)$ tends to $\prod\left(\mathrm{Z}-\mathrm{Y}_{i}\right)=\mathrm{Z}^{k} \cdot \prod_{i \notin \mathrm{~J}}\left(\mathrm{Z}-\mathrm{Y}_{i}\right)$, and the rational fraction $f_{\mathrm{N}(\mathrm{Q})}$ tends to

$$
g=\frac{\mathrm{Z}^{4-k}}{\prod_{i \notin \mathrm{~J}}\left(\mathrm{Z}-\mathrm{Y}_{i}\right)}
$$

which has no pole at 0 . Hence $\sum_{i \in \mathrm{~J}} r_{i}^{(n)} \rightarrow 0$, and $\sum_{i \in \mathrm{~J}} r_{i}=0$, giving a contradiction.

This completes the Proof of Proposition 7.5.10.

### 7.5.13 A Uniqueness Property

Proposition Let $\widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \in \widetilde{\mathscr{B}}$ and $r=\left(r_{i}\right)=\rho(\widetilde{\mathrm{Q}})$. If $r_{i}=\frac{1}{5}$ for all $i$, then $\mathrm{Q}=\mathrm{Z}^{5}-1$.
Proof The function $5 f_{\mathrm{Q}}=\frac{5 \mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}=\sum \frac{1}{\mathrm{Z}-x_{i}}$ is the logarithmic derivative of $\mathrm{Q}=$ $\prod\left(\mathrm{Z}-x_{i}\right)$. So $\mathrm{Q}^{\prime}=5 \mathrm{Z}^{4}$, and Q is of the form $\mathrm{Z}^{5}-c$. As $c=\prod x_{i}=1$, where $\mathrm{Q}(\mathrm{Z})=\mathrm{Z}^{5}-1$,
Corollary The covering map $\rho: \widetilde{\mathscr{B}}^{*} \rightarrow \mathscr{W}^{*}$ is of degree $5!=120$.

### 7.5.14 Topological Numbering

For $\mathrm{Q} \in \mathscr{B}^{+}$, there is a natural topological way to number the roots of Q .

First let $\mathrm{Q} \in \mathscr{B}$ be arbitrary. On a sufficiently small connected neighbourhood V of 0 , there is a unique holomorphic function F such that $\mathrm{F}^{\prime}=f_{\mathrm{Q}}=\frac{\mathrm{Z}^{4}}{\mathrm{Q}}$ and $\mathrm{F}(0)=0$. As $f_{\mathrm{Q}}(z)=-z^{4}+\mathrm{O}\left(z^{5}\right), \mathrm{F}(z)=-\frac{1}{5} z^{5}+\mathrm{O}\left(z^{6}\right)$. If necessary contracting V , we find a holomorphic map $\phi: \mathrm{V} \rightarrow \mathbb{C}$ such that $\mathrm{F}(z)=-\frac{1}{5}(\phi(z))^{5}, \phi(0)=0$ and $\phi^{\prime}(0)=1$, and $\phi$ may be assumed to be an isomorphism from V onto a disc $\mathbb{D}_{\ell}$.

Now suppose that $\mathrm{Q} \in \mathscr{B}^{+}$. Then there are $r_{i}>0$ such that $f_{\mathrm{Q}}(\mathrm{Z})=\sum \frac{r_{i}}{\mathrm{Z}-x_{i}}$; set $h(z)=\sum r_{i} \log \left|z-x_{i}\right|$ and $\mathrm{U}=\{z \mid h(z)<0\} . \mathrm{On} \mathrm{V}, h=\mathfrak{R}(\mathrm{F})$.

The set U is an open relatively compact subset in $\mathbb{C}$ since as $z$ tends to infinity, $h(z) \sim \log |z|$ tends to infinity. The open $\mathrm{V} \cap \mathrm{U}=\left\{z \mid \Re\left((\phi(z))^{5}>0\right\}\right.$ is the union of 5 sectors $S_{1}, \ldots, S_{5}$, where

$$
\mathrm{S}_{k}=\left\{\left.z \in \mathrm{~V}-\{0\}| | \arg (\phi(z))-\frac{2 \pi k}{5} \right\rvert\,<\frac{\pi}{10}\right\}
$$

(here, for $\theta \in \mathbb{T}=\mathbb{R} / \mathbb{Z}$, set $|\theta|=\inf |t|$, where $t$ is a representative of $\theta$ ).
Proposition (a) For $\mathrm{Q} \in \mathscr{B}^{+}$, the open subset U has five connected components $\left(\mathrm{U}_{i}\right)_{i \in \mathbb{Z} /(5)}$ with $\mathrm{U}_{i} \cap \mathrm{~V}=\mathrm{S}_{i}$;
(b) Each connected component $\mathrm{U}_{i}$ contains a unique root $x_{i}^{\text {top }}$ of Q .
(c) As Q varies in $\mathscr{B}^{+}$, for each $i \in \mathbb{Z} /(5)$ the root $x_{i}^{\text {topp }}(\mathrm{Q})$ depends continuously on Q .

Proof (a) and (b) Let $\mathrm{U}_{i}$ be the connected component of U containing $\mathrm{S}_{i}$. The $\mathrm{U}_{i}$ are distinct, for otherwise there would be a Jordan curve in $\overline{\mathrm{U}}$ around some connected component of $\mathrm{V}-\overline{\mathrm{U}}$, contradicting the maximum principle.

Thanks to this principle (or rather to the minimum one), each connected component of U contains a root of Q ; so U at most 5 connected components.

Hence, U has exactly 5 connected components, each meets V along a unique $\mathrm{S}_{i}$ and contains a unique root of Q .
(c) We make Q vary in $\mathscr{B}^{+}$; We will write $\mathrm{S}_{i}(\mathrm{Q}), x_{i}^{\text {top }}(\mathrm{Q})$, etc. Let $\mathrm{Q}_{0} \in \mathscr{A}^{+}$and $i \in \mathbb{Z} /(5)$. We choose a point $a \in \mathrm{~S}_{i}\left(\mathrm{Q}_{0}\right)$ and a path $\gamma$ from $a$ to $x_{i}^{\text {top }}\left(\mathrm{Q}_{0}\right)$ in $\mathrm{U}_{\mathrm{Q}_{0}}$. The implicit function theorem enables us to find a continuous function $\mathrm{Q} \mapsto \xi(\mathrm{Q})$ defined in the neighbourhood of $\mathrm{Q}_{0}$, such that $\xi(\mathrm{Q})$ is a root of Q and $\xi\left(\mathrm{Q}_{0}\right)=x_{i}^{\text {top }}\left(\mathrm{Q}_{0}\right)$, and a path $\gamma_{\mathrm{Q}}$ from $a$ to $\xi(\mathrm{Q})$ in $\mathbb{C}$ continuously dependent on Q . For Q sufficiently near $\mathrm{Q}_{0}, a \in \mathrm{~S}_{i}(\mathrm{Q})$, and the path $\gamma_{\mathrm{Q}}$ is contained in $\mathrm{U}_{\mathrm{Q}}$. Then $\xi(\mathrm{Q})$ is in the same connected component as $a \in \mathrm{~S}_{i}(\mathrm{Q})$; thus $\xi(\mathrm{Q})=x_{i}^{\text {top }}(\mathrm{Q})$, which implies that $x_{i}^{\text {top }}(\mathrm{Q})$ is continuously dependent on Q .

Corollary 7.3 The space $\mathscr{B}^{+}$is homeomorphic to $\mathscr{W}^{+}$. In particular, it is connected.
Proof Let E be a connected component of $\mathscr{B}^{+}$. The map $\mathrm{Q} \mapsto\left(\mathrm{Q}, x_{0}^{\text {top }}(\mathrm{Q}), \ldots\right.$, $\left.x_{4}^{\text {top }}(\mathrm{Q})\right)$ is a continuous section $\mathscr{B}^{+} \rightarrow \widetilde{\mathscr{B}}^{+}$; it induces a homeomorphism from E onto a connected component of $\widetilde{\mathscr{B}}^{+}$, and $\rho$ induces a homeomorphism from this component onto $\mathscr{W}^{+}$since the latter is connected. This gives a homeomorphism $\psi: \mathrm{E} \rightarrow \mathscr{W}^{+}$.

In particular there exists $\mathbf{Q}_{1} \in \mathscr{B}^{+}$such that $\psi\left(\mathbf{Q}_{1}\right)=\rho\left(\widetilde{\mathbf{Q}}_{1}\right)=\left(\frac{1}{5}, \ldots, \frac{1}{5}\right)$ for some $\widetilde{\mathrm{Q}}_{1}$ over $\mathrm{Q}_{1}$. However this implies that $\mathrm{Q}_{1}=\mathrm{Z}^{5}-1$. Hence $\mathscr{B}^{+}$only has one connected component, namely $\mathrm{Z}^{5}-1$.

Corollary 7.4 The space $\mathscr{A}^{+}$is connected.
Proof The map $(\lambda, \mathrm{Q}) \mapsto \lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right)$ from $\mathbb{C}^{*} \times \mathscr{B}^{+}$to $\mathscr{A}^{+}$is continuous and surjective, since there is a lifting $\left(\lambda, \mathrm{Q},\left(x_{i}\right)\right) \mapsto\left(\lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right),\left(\lambda x_{i}\right)\right)$, which is a covering. Hence $\mathscr{A}^{+}$is connected.

### 7.5.15 The Invariant $\theta_{\text {top }}(\mathrm{Q})$

Let $\widetilde{\mathrm{Q}}=\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) \in \tilde{\mathscr{A}}$ and $r=\left(r_{i}\right)=\rho(\widetilde{\mathrm{Q}})$. As $f_{\mathrm{Q}}(\mathrm{Z})=\frac{\mathrm{Z}^{4}}{\mathrm{Q}(\mathrm{Z})}=\sum \frac{r_{i}}{\mathrm{Z}-x_{i}}$ and $r_{i}^{-1}=x_{i}^{-4} \prod_{j \mid j \neq i}\left(x_{i}-x_{j}\right), \prod_{i}^{-1}=\prod_{i, j \mid j \neq i}\left(x_{i}-x_{j}\right)$.

Set

$$
\theta(\widetilde{\mathrm{Q}})=\prod x_{i}^{-2} \cdot \prod_{i, j \mid i<j}\left(x_{i}-x_{j}\right) .
$$

The number of factors in this product is 10 , which is even. Thus $(\theta(\widetilde{\mathbb{Q}}))^{2}=\prod r_{i}^{-1}$. If $\sigma$ is a permutation of $\{1, \ldots, 5\}$, then

$$
\theta\left(\mathrm{Q}, x_{\sigma(1)}, \ldots, x_{\sigma(5)}\right)=\varepsilon(\sigma) \cdot \theta\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right) .
$$

As 5 is odd, any cyclic permutation is even, and $\theta(\widetilde{\mathrm{Q}})$ only depends on the chosen cyclic order on the roots of Q .

So $\theta\left(\lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right), \lambda x_{1}, \ldots, \lambda x_{5}\right)=\theta\left(\mathrm{Q}, x_{1}, \ldots, x_{5}\right)$.
If $\widetilde{\mathrm{Q}} \in \widetilde{\mathscr{A}}^{+}$, then $(\theta(\widetilde{\mathrm{Q}}))^{2}=\prod r_{i}^{-1} \in \mathbb{R}_{+}^{*}$, and so $\theta(\widetilde{\mathrm{Q}}) \in \mathbb{R}^{*}$.
For $\mathrm{Q} \in \mathscr{B}^{+}$, set $\theta_{\operatorname{tap}}(\mathrm{Q})=\theta\left(\mathrm{Q},\left(x_{i}^{\text {totp }}(\mathrm{Q})\right)_{i \in \mathbb{Z} /(5)}\right)$. For $\mathrm{Q} \in \mathscr{A}^{+}$, there exists $\lambda$ such that $\mathrm{Q}_{1}=\lambda^{5} \mathrm{Q}\left(\frac{\mathrm{Z}}{\lambda}\right) \in \mathscr{B}^{+}$. Then set $\theta_{\text {top }}(\mathrm{Q})=\theta_{\text {topp }}\left(\mathrm{Q}_{1}\right)$; there are 5 possible values for $\lambda$, but they all give the same cyclic order on the set of roots of Q , so that $\theta_{\text {top }}(\mathrm{Q})$ is well defined.
Proposition For all $\mathrm{Q} \in \mathscr{A}^{+}, \theta_{\text {top }}(\mathrm{Q})<0$.
Proof The map $\theta_{\text {top }}: \mathscr{A}^{+} \rightarrow \mathbb{R}^{*}$ is continuous by (7.5.14, Proposition), (c); hence its sign remains constant. It suffices to check that $\theta_{\text {top }}\left(\mathrm{Z}^{5}-1\right)<0$.

Setting $x_{k}=e^{2 i \pi k / 5}$, for $k=-2, \ldots, 2$,

$$
\begin{aligned}
\theta_{\text {top }}\left(\mathrm{Z}^{5}-1\right)= & \left(\left(x_{-2}-x_{-1}\right)\left(x_{1}-x_{2}\right)\right) \cdot\left(\left(x_{-2}-x_{0}\right)\left(x_{0}-x_{2}\right)\right) \cdot\left(\left(x_{-1}-x_{0}\right)\right. \\
& \left.\left(x_{0}-x_{1}\right)\right) \cdot\left(\left(x_{-1}-x_{2}\right)\left(x_{2}-x_{1}\right)\right) \cdot\left(\left(x_{-2}-x_{2}\right)\left(x_{-1}-x_{1}\right)\right) .
\end{aligned}
$$

In this product of 5 factors, the first 4 are $<0$ as they are the opposite of the product of two conjugates, and the last one is a product of two purely imaginary numbers of the same sign. Hence $\theta_{\text {top }}\left(Z^{5}-1\right)<0$.


### 7.5.16 Relation Between $\delta_{\text {top }}(\mathrm{P})$ and $\theta_{\text {topp }}(\mathrm{Q})$

Let P be a Belyi polynomial satisfying the conditions of (7.5.6). Then $\mathrm{Q}(\mathrm{Z})=$ $\prod\left(\mathrm{Z}-x_{i}\right) \in \mathscr{A}$, with $r_{i}=\frac{d_{i}}{20}$. The function $h(\mathrm{Z})=\sum r_{i} \log \left|\mathrm{Z}-x_{i}\right|$ equals $\frac{1}{20} \sum d_{i}$ $\log \left|\mathrm{Z}-x_{i}\right|=\frac{1}{20} \log |\mathrm{P}(\mathrm{Z})|$. Hence the edges connected to 0 in the graph $\mathrm{P}^{-1}([0,1])$ are each contained in a connected component of the open subset $\mathrm{U}=\{\mathrm{Z} \mid h(\mathrm{Z})<0\}$, and the cyclic order defined on the set of $x_{i}$ by the topological numbering of (7.5.14) agrees with that defined by the branches of the graph $\mathrm{P}^{-1}(0)$.

Therefore

$$
\theta_{\text {top }}(\mathrm{Q})=\delta_{\text {top }}(\mathrm{P}) .
$$

By Proposition 7.5.15, $\delta_{\text {top }}(\mathrm{P})=\theta_{\text {top }}(\mathrm{Q})<0$, proving Theorem 7.5.8.

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## Index of Notation

| Classical notation |  |
| :---: | :---: |
| Card X: | cardinal of X . |
| A - B : | $A \cap C B$. |
| $\mathfrak{P}(\mathrm{X})$ : | set of all subsets of X. |
| $x \sim_{\mu} y:$ | $x$ and $y$ are equivalent with respect to relation $\mu$. |
| ヨl: | there exists a unique |
| inf, sup: | infimum, supremum. |
| 1.c.m.: | least common multiple. |
| g.c.d.: | greatest common divisor. |
| $\mathrm{L}(\mathrm{E}, \mathrm{F})$ : | space of linear maps from E to F. |
| $\ell^{2}$ : | space of square integrable sequences. |
| $\langle x, y\rangle$ : | scalar product. |
| $\widehat{\mathrm{E}}$ : | completion of E. |
| $\mathrm{A}\left[\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right]\right]:$ | ring of formal series with coefficients in A . |
| $\mathbb{C}\left\{\mathrm{X}_{1}, \ldots, \mathrm{X}_{n}\right\}$ : | convergent power series in the neighbourhood of O . |
| $\aleph_{0}$ : | countably infinite cardinal. |
| $\mathrm{P}^{n} \mathbb{R}, \mathrm{P}^{n} \mathbb{C}$ : | $n$-dimensional real projective space, resp. complex. |
| $\mathrm{S}^{n}$ : | unit sphere in $\mathbb{R}^{n+1}$. |
| X/G: | quotient of $X$ by the equivalence relation whose classes are the G-orbits. |
| $\mathrm{O}_{n}\left(\right.$ resp. $\left.\mathrm{U}_{n}\right)$ : | group of linear isometries of $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ ). |
| $\mathfrak{S}(\mathrm{E})$ : | permutation group of E . |
| $\mathfrak{A}(\mathrm{E})$ : | alternating group. |
| $\begin{aligned} & \mathrm{SO}_{n}\left(\text { resp. } \mathrm{SU}_{n}\right) \text { : } \\ & \mathfrak{P}(\mathrm{E}): \end{aligned}$ | subgroup of $\mathrm{O}_{n}$ (resp. $\mathrm{U}_{n}$ ) of matrices with determinant +1 . set of all subsets of E . |

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[^0]:    ${ }^{1}$ See N. Bourbaki, Algebra, ch. 1 to 3 (1970), ch. 3, §7, n ${ }^{\circ} 4$, (A III.80).

[^1]:    ${ }^{2}$ More precisely, assume there is an universe $\mathrm{U}(2.1 .3)$ such that $\mathscr{C} \subset \mathrm{U}$. For the definition of $\mathscr{C}$, we only consider inverse systems $\left(X_{i}\right)_{i \in \mathrm{I}}$ such that $\mathrm{I} \in \mathrm{U}$. Then $\mathscr{C} \subset \mathrm{U}$.

[^2]:    ${ }^{1}$ Queysanne [1], §186, Lelong-Ferrand and Arnaudiès [2], cor. of th. 4.7.2, p. 151.
    ${ }^{2}$ Lelong-Ferrand and Arnaudiès [3], th. 4.11.2, p. 387.

[^3]:    ${ }^{3}$ Queysanne [1], §192, Exercise 5, p. 413, Lelong-Ferrand and Arnaudiès [2], chap. 4, Exercise 17, p. 476.

[^4]:    ${ }^{4}$ It is usually written $x \mid y$. The disadvantage of the symbol | is that crossing it is problematic and it cannot be reversed.

[^5]:    ${ }^{5}$ Abbreviations of "greatest common divisor" and "least common multiple".

[^6]:    ${ }^{6}$ The main difference between modules and vector spaces is that there are torsion modules and submodules that are not direct factors.

[^7]:    ${ }^{7}$ Here the assumption of the commutativity of A is essential.
    ${ }^{8}$ We use the usual notation $u^{*}, v_{*}$ which omits any mention of E and F (see 2.2.1).

[^8]:    ${ }^{9}$ Queysanne [1], §183, Lelong-Ferrand and Arnaudiès [2], 4.1, p. 123.
    ${ }^{10}$ Queysanne [1], §186, [2], 4.7, p. 151.
    ${ }^{11}$ Queysanne [1], §135, [2], th. 8.3.1, p. 253.

[^9]:    ${ }^{12}$ Bourbaki [4], chap. 2, §7, th. 3 (A II.96).

[^10]:    ${ }^{13} n>0$ for otherwise det $\varphi=1$. If $h a_{i}^{j}=0$, for all $(i, j), \varphi\left(h e_{i}\right)=0$ and $\varphi$ is not injective.

[^11]:    ${ }^{14}$ Guichardet [5], 1.2.

[^12]:    ${ }^{15}$ Cf. Queysanne [1], Example 488, p. 509 or [2], def. 11.3.2, p. 343
    ${ }^{16}$ Queysanne [1], Exerc. 486 and 495, [2], § 11.4, p. 348.

[^13]:    ${ }^{17}$ Frisch [6].

[^14]:    ${ }^{18} \mathrm{P}_{1}$ not necessarily homogeneous.
    ${ }^{19}$ This theorem is known as the finite basis theorem (the term basis being at that time used for generating family).

[^15]:    ${ }^{20}$ If $\mathrm{B}=\mathrm{A}\left[\left(\mathrm{X}_{i}\right)_{i \in \mathrm{I}}\right]$ and $\forall i, x_{i}=\mathrm{X}_{i}$, then $f=1_{\mathrm{A}\left[\left(\mathrm{X}_{i}\right]\right.}$; so $\mathrm{P}\left(\left(\mathrm{X}_{i}\right)\right)=\mathrm{P}$.
    ${ }^{21}$ Queysanne [1], § 192, Corollary 1, p. 413, Lelong-Ferrand and Arnaudiès [2], theorem 4.5.2, p. 139.

[^16]:    ${ }^{1}$ Both open and closed, Translator.
    ${ }^{2}$ We speak of an $n$-fold covering when $\operatorname{deg}_{b}(\mathrm{X})=n$, Translator.

[^17]:    ${ }^{3}$ The reader should be aware of possible conflicts with the usual meaning of the term "simply", for example: I do not assume X simply connected, but simply to be connected.

[^18]:    ${ }^{1}$ This convention is not usual.

[^19]:    ${ }^{1}$ Cartan [1], 2.7, Corollary of th. 2.7.4, p. 74; see also [2], 8.9, making sure that $a_{n}=0$ for $n<0$.

[^20]:    ${ }^{2}$ See Cartan [4] for the notion of basis for differential forms on manifolds, and Weil [5] for notions specific to complex analytic manifolds.

[^21]:    ${ }^{3}$ The definitions given in this section are specific to the situation. Their generalization requires changes.

[^22]:    ${ }^{4}$ See Zisman [7], 4.1.3 or Bourbaki [8], Chap. 10, § 1, cor. 3. (A X-7).

[^23]:    ${ }^{5}$ The Riemannian metrics considered are those that are compatible with the complex structure.

[^24]:    ${ }^{6}$ A not necessarily holomorphic map $f$ is an isometry with respect to the distances defined by $\mu_{1}$ and $\mu_{2}$ if it satisfies condition (iii). It may be shown that, as a result, $f$ is either holomorphic, or antiholomorphic. We will prove this in the particular case of the Poincaré distance.

[^25]:    ${ }^{7}$ In fact, this assumption can be shown to be superfluous for the results of this subsection and of 6.10.15.

