# Advanced Algebraic Structures 

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(partly a translation of parts of older lecture notes in Dutch by L.N.M. van Geemen, H.W. Lenstra, F. Oort, and J. Top. )
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These lecture notes contain a translation into English of a number of chapters from the Dutch lecture notes Algebra II (Algebraic Structures) as they were used in the mathematics curriculum of Groningen University during the period 1993-2013. The original Dutch text may be found at http://www.math.rug.nl/~top/dic. pdf.

Both the present text and the original are loosely based on another Dutch text on Rings and Fields, called Algebra II, written in the late 1970's at the university of Amsterdam by Prof.dr. F. Oort and Prof.dr. H.W. Lenstra. In the 80's L.N.M. van Geemen at Utrecht university added some chapters to the text, and in the 90 's in Groningen I included a number of changes.

The present translation consists of two parts. The first one (Algebraic Structures) deals with basic concepts and properties of rings and fields as presented in the chapters $1-5,7-9$, and 12 of the Dutch notes. The second one (Advanced Algebraic Structures) discusses the chapters 8 (automorphisms and splitting fields, in slightly more detail than originally in order to facilitate treating Galois theory), a discussion of basic Galois theory based on previous lecture notes by Prof. M. van der Put and me, the chapters 11 (symmetric polynomials) and 10 (algebraically closed fields) as well as an extended version of chapter 14 (cyclotomic fields), and finally the chapters 6,13 on (projective) modules. A short introduction to tensor products is added.

# AUTOMORPHISMS OF FIELDS AND SPLITTING FIELDS 

## I. 1 Basic definitions

We repeat some basic definitions and properties of fields. In fact some of this was already used in the previous chapter.

A field $K$ has a smallest subfield, called the prime field of $K$. This prime field is either $\mathbb{Q}$, in which case the characteristic of $K$ is 0 , or $\mathbb{F}_{p}=\mathbb{Z} / p \mathbb{Z}$ for some prime number $p$, in which case the characteristic of $K$ is $p$.

Let $K \subset L$ be an extension of fields. Then $L$ can be seen as a vector space over $K$. The dimension of this vector space is denoted by $[L: K]$. In general $[L: K]$ is infinite. If $[L: K]$ is finite then $L$ is called a finite extension of $K$ of degree $[L: K]$.

Let $A$ be a subset of $L$, then $K(A)$ denotes the smallest subfield of $L$ containing both $A$ and $K$. Similarly, $K[A]$ denotes the smallest subring of $L$ containing $A$ and $K$. For $A=\left\{a_{1}, \ldots, a_{s}\right\}$ one writes $K(A)=K\left(a_{1}, \ldots, a_{s}\right)$ and $K[A]=K\left[a_{1}, \ldots, a_{s}\right]$. An element $a \in L$ is called algebraic over $K$ if there is a non-zero polynomial $f \in K[x]$ with $f(a)=0$. For algebraic $a$ there is a polynomial $F \neq 0$ of minimal degree such that $F(a)=0$. If $F$ is normalized to be monic, then $F$ is unique and is called the minimal polynomial of $a$ over $K$. Let $F$ have degree $n$, then $K(a)=K[a] \cong K[x] /(F)$ is a vector space over $K$ with dimension $n$ and with basis $1, a, \ldots, a^{n-1}$.

An element $a \in L$ which is not algebraic over $K$ is called transcendental over $K$. The obvious ring homomorphism $K[x] \rightarrow K[a]$ (i.e., $\sum c_{i} x^{i} \mapsto \sum c_{i} a^{i}$ ) is an isomorphism. Thus $K[a]$ is not a field. The field of fractions $K(a)$ of $K[a]$ is in this case isomorphic to $K(x)$, i.e., to the field of rational functions over $K$.

The field extension $K \subset L$ is called finitely generated if there are elements $a_{1}, \ldots, a_{s} \in L$ such that $L=K\left(a_{1}, \ldots, a_{s}\right)$. The elements $a_{1}, \ldots, a_{s}$ are called algebraically dependent over $K$ if there is a non-zero polynomial $f \in K\left[x_{1}, \ldots, x_{s}\right]$ with $f\left(a_{1}, \ldots, a_{s}\right)=0$. If such a polynomial $f$ does not exists, then $a_{1}, \ldots, a_{s}$ are called algebraically independent over $K$. In the latter case, the obvious homomorphism $K\left[x_{1}, \ldots, x_{s}\right] \rightarrow K\left[a_{1}, \ldots, a_{s}\right]$ is an isomorphism.Then $K \subset L$ is called a purely transcendental extension of $K$ of transcendence degree $s$ and $L$ is isomorphic to the field of fractions of the polynomial ring $K\left[x_{1}, \ldots, x_{s}\right]$. This field of fractions is denoted by $K\left(x_{1}, \ldots, x_{s}\right)$. It is an exercise (see Exercise 3 below) to show that a finitely generated extension $L$ of $K$ has an intermediate field $M$ (i.e., $K \subset M \subset L$ ) such that $K \subset M$ is purely transcendental and $M \subset L$ is a finite extension. The transcendence degree of $L$ over $K$ is defined as the transcendence degree of $M$ over $K$. It is not at all clear that this definition is a valid one. One has to show that it does not depend on the choice of the intermediate purely transcendental extension. This is, for instance, Theorem 25 in Chapter II of the book Commutative Algebra (Vol. 1) by O. Zariski and P. Samuel. A different proof which uses theory of "deriva-
tions" and which is valid only in characteristic zero, is to compute the dimension of the $M$-vector space consisting of all $K$-linear maps $D: M \rightarrow M$ which satisfy $D\left(m_{1} m_{2}\right)=m_{1} D\left(m_{2}\right)+m_{2} D\left(m_{1}\right)$.

## I. 2 Homomorphisms of fields

This section reviews some basic results on automorphisms of fields. If $K$ and $L$ are fields then $\phi: K \rightarrow L$ is called a homomorphism of fields (also called 'field homomorphism') if $\phi$ is a unitary ring homomorphism. In particular it satisfies $\phi(1)=1$.

A homomorphism of fields has the property

$$
\phi\left(\frac{1}{a}\right)=\phi(a)^{-1}, \quad \phi\left(\frac{a}{b}\right)=\frac{\phi(a)}{\phi(b)},
$$

since $\phi(a) \cdot \phi\left(a^{-1}\right)=\phi\left(a \cdot a^{-1}\right)=\phi(1)=1$ implies $\phi\left(a^{-1}\right)=(\phi(a))^{-1}$.
Using $\phi(1+1+\ldots+1)=1+1+\ldots+1$ it follows that the image of the prime field $K_{0}$ of $K$ is the prime field $L_{0}$ of $L$. Hence $\phi$ yields an isomorphism from $K_{0}$ to $L_{0}$. In case $K \subset L$ then $K_{0}=L_{0}$ and the restriction of $\phi$ to $K_{0}$ is the identity map.

The image $\phi(K)$ of a field homomorphism $\phi: K \rightarrow L$ is also a field. Every field homomorphism is injective, since the only ideals in $K$ are ( 0 ) and $K$ and $1 \notin \operatorname{Ker}(\phi)$. A field homomorphism need not be surjective. For example the inclusion $\mathbb{R} \hookrightarrow \mathbb{C}$ is not. Even if $K=L$ a field homomorphism needs not be surjective, as illustrated by Example I.2.2.

The composition of field homomorphisms

$$
K \xrightarrow{\phi} L \xrightarrow{\psi} M
$$

is also a field homomorphism, as one easily verifies.
An interesting and important field homomorphism exists in case $\operatorname{char}(K)=p$ :
I.2.1 Theorem. Let $K$ be a field such that $\operatorname{char}(K)=p>0$. Put

$$
F: K \longrightarrow K, \quad F: x \mapsto x^{p} .
$$

Then $F$ is a field homomorphism called the Frobenius homomorphism. In case $K$ is finite, $F$ is even a field automorphism.

Proof. Note that $F(1)=1$ and $F(a b)=(a b)^{p}=a^{p} b^{p}=F(a) F(b)$ (since $K$ is commutative). It remains to show that $F(a+b)=(a+b)^{p}$ equals $F(a)+F(b)=a^{p}+b^{p}$.

Newton's binomium formula, which holds in any commutative ring, implies:

$$
(a+b)^{p}=\sum_{k=0}^{p}\binom{p}{k} a^{k} b^{p-k}, \quad \text { with } \quad\binom{p}{k}=\frac{p!}{k!(p-k)!} \in \mathbb{Z} .
$$

The numerator of $\binom{p}{k}$ is divisible by $p$, and because $0<k<p$ and $p$ is prime, the denominator is not divisible by $p$. Hence

$$
(a+b)^{p}=a^{p}+b^{p}+p \cdot c
$$

for some $c \in K$. As $p=1+1+\ldots+1(p \times)$ and $\operatorname{char}(K)=p$ it follows that $p=0 \in K$. We conclude $(a+b)^{p}=a^{p}+b^{p}$.

In case $K$ is finite every injective map from $K$ to itself (so for example $F$ ) is surjective as well. Hence $F$ is bijective. The fact that $F^{-1}$ is a field homomorphism as well, is easily verified. This shows Theorem I.2.1
1.2.2 Example. Let $K=\mathbb{F}_{p}(T)$ be the field of rational functions (= quotients of polynomials) with coefficients in $\mathbb{F}_{p}$. Put $f(T)=\frac{a_{0}+a_{1} T+\ldots+a_{n} T^{n}}{b_{0}+b_{1} T+\ldots+b_{m} T^{m}} \in \mathbb{F}_{p}(T)$, then

$$
\begin{aligned}
F(f(T)) & =F\left(\frac{a_{0}+a_{1} T+\ldots+a_{n} T^{n}}{b_{0}+b_{1} T+\ldots+b_{m} T^{m}}\right) \\
& =\frac{F\left(a_{0}+a_{1} T+\ldots+a_{n} T^{n}\right)}{F\left(b_{0}+b_{1} T+\ldots+b_{m} T^{m}\right)} \\
& =\frac{F\left(a_{0}\right)+F\left(a_{1}\right) F(T)+\ldots+F\left(a_{n}\right) F\left(T^{n}\right)}{F\left(b_{0}\right)+F\left(b_{1}\right) F(T)+\ldots+F\left(b_{m}\right) F\left(T^{m}\right)} \\
& =\frac{a_{0}+a_{1} T^{p}+\ldots+a_{n} T^{p n}}{b_{0}+b_{1} T^{p}+\ldots+b_{m} T^{p m}}=f\left(T^{p}\right) .
\end{aligned}
$$

Here we used $F(a)=a$ for all $a \in \mathbb{F}_{p}$, the prime field of $\mathbb{F}_{p}(T)$. The image of the Frobenius homomorphism $F$ therefore consists of all rational functions in the variable $T^{p}$ with coefficiënts in $\mathbb{F}_{p}$. In particular $F$ is not surjective on $\mathbb{F}_{p}(T)$, for example $T \notin \operatorname{image}(F)$ (verify for yourself!).
I.2.3 Definition. If $L$ and $L^{\prime}$ are extensions of a field $K$ then a $K$-homomorphism $L \rightarrow L^{\prime}$ is a field homomorphism

$$
\phi: L \rightarrow L^{\prime} \quad \text { such that } \quad \phi_{\mid K}=\mathrm{id}_{K} .
$$

A $K$-isomorphism is a bijective $K$-homomorphism.
The fields $L$ and $L^{\prime}$ are called $K$-isomorphic (notation: $L \cong_{K} L^{\prime}$ ) if a $K$-isomorphism $L \rightarrow L^{\prime}$ exists.
A $K$-automorphism is a $K$-isomorphism with $L=L^{\prime}$.

If the fields $K$ and $L$ have the same prime field $K_{0}$, then every field homomorphism $K \rightarrow L$ is a $K_{0}$-homomorphism.
I.2.4 Definition. Given an extension $L$ of a field $K$ we write $\operatorname{Aut}_{K}(L)$ for the group of all $K$-automorphisms of $L$. Here the group law is the composition of maps, and the unit element is $\mathrm{id}_{L}$.
1.2.5 Theorem. Let $L$ be a finite extension of the field $K$ and $\alpha \in L$ with minimal polynomial $f_{K}^{\alpha} \in K[X]$.
(i) For every $\phi \in \operatorname{Aut}_{K}(L)$ we have that $\phi(\alpha)$ is a zero of $f_{K}^{\alpha}$.
(ii) If $m$ is the number of zeros of $f_{K}^{\alpha}$ in the field $K[\alpha] \subset L$, then \#Aut ${ }_{K}(K[\alpha])=m$. Here \#S the number of elements of a set $S$.
(iii) The number of $K$-automorphisms of $K[\alpha]$ is $\leq \operatorname{deg}\left(f_{K}^{\alpha}\right)$.

Proof. (i): Write $f_{K}^{\alpha}=X^{n}+a_{n-1} X^{n-1} \ldots+a_{1} X+a_{0}$ with $a_{i} \in K$. Applying $\phi \in A u t_{K}(L)$ to the equality $0=f_{K}^{\alpha}(\alpha)$ yields:

$$
\begin{aligned}
0 & =\phi\left(\alpha^{n}+a_{n-1} \alpha^{n-1}+\ldots a_{1} \alpha+a_{0}\right) \\
& =\phi(\alpha)^{n}+\phi\left(a_{n-1}\right) \phi(\alpha)^{n-1} \ldots \phi\left(a_{1}\right) \phi(\alpha)+\phi\left(a_{0}\right) \\
& =\phi(\alpha)^{n}+a_{n-1} \phi(\alpha)^{n-1}+\ldots+a_{1} \phi(\alpha)+a_{0} \\
& =f_{K}^{\alpha}(\phi(\alpha)),
\end{aligned}
$$

where we used $\phi\left(a_{i}\right)=a_{i}$ (note $a_{i} \in K$ and $\phi_{\mid K}=\operatorname{id}_{K}$ ). So $\phi(\alpha)$ is a zero of $f_{K}^{\alpha}$. (ii): Let $\left\{\alpha_{1}, \ldots, \alpha_{m}\right\} \subset K[\alpha]$ be the zeros of $f_{K}^{\alpha}$ in $K[\alpha]$, with $\alpha_{1}=\alpha$. We claim that the map

$$
\Delta: \operatorname{Aut}_{K}(K[\alpha]) \longrightarrow\left\{\alpha_{1}, \alpha_{2}, \ldots, \alpha_{m}\right\}, \quad \Delta(\phi):=\phi(\alpha)
$$

is a bijection.

From (i) we know that $\Delta$ is well defined. We now show that $\Delta$ is injective. Every $x \in K[\alpha]$ can be given in a unique way as $x=x_{0}+x_{1} \alpha+\ldots+x_{n-1} \alpha^{n-1}$ with $x_{i} \in K$. Then

$$
\begin{aligned}
\phi(x) & =\phi\left(x_{0}\right)+\phi\left(x_{1}\right) \phi(\alpha)+\ldots+\phi\left(x_{n-1}\right) \phi(\alpha)^{n-1} \\
& =x_{0}+x_{1} \phi(\alpha)+\ldots+x_{n-1} \phi(\alpha)^{n-1},
\end{aligned}
$$

so $\phi(x)$ is completely determined by $\phi(\alpha)=\Delta(\phi)$. Therefore, if $\Delta(\phi)=\Delta(\psi)$, then $\phi(x)=\psi(x)$ for all $x \in K[\alpha]$, so $\phi=\psi$.

In order to prove surjectivity, given any zero $\beta \in K[\alpha]$ of $f_{K}^{\alpha}$ we have to construct a $K$-automorphism $\phi$ with $\phi(\alpha)=\beta$. This is done as follows. Take

$$
\operatorname{ev}_{\beta}: K[X] \longrightarrow K[\alpha]
$$

the evaluation homomorphism in $\beta$. Since $\beta$ is a zero of $f_{K}^{\alpha} \in K[X]$ it follows that $f_{K}^{\alpha} \in \operatorname{Ker}\left(\Phi_{\beta}\right)$. However $f_{K}^{\alpha}$ is irreducible in $K[X]$, hence $\operatorname{Ker}\left(\operatorname{ev}_{\beta}\right)=\left(f_{K}^{\alpha}\right)$. So $\operatorname{ev}_{\beta}$ induces an injective $K$-homomorphism $\overline{\operatorname{ev}_{\beta}}: K[X] /\left(f_{K}^{\alpha}\right) \rightarrow K[\alpha]$ with image $K[\beta] \subset K[\alpha]$. The first isomorphism theorem for rings now implies $K[\beta] \cong K[X] /\left(f_{K}^{\alpha}\right) ;$ the given isomorphism is $K$-linear, hence $K[\beta]$ is a linear subspace of $K[\alpha]$ with dimension $\operatorname{dim}_{K}\left(K[X] /\left(f_{K}^{\alpha}\right)\right)=\operatorname{dim}_{K}(K[\alpha])$. This dimension is finite and therefore $K[\beta]=K[\alpha]$; in other words $\overline{\mathrm{ev}_{\beta}}$ is an isomorphism. Analogous to the above, write $\overline{\mathrm{ev}_{\alpha}}$ for the isomorphism $K[X] /\left(f_{K}^{\alpha}\right) \xlongequal{\cong} K[\alpha]$ induced by the evaluation homomorphism $\mathrm{ev}_{\alpha}$. We have isomorphisms:

and hence a $K$-automorphism $\phi:=\overline{\operatorname{ev}_{\beta}} \circ{\overline{\operatorname{ev}_{\alpha}}}^{-1}: \quad K[\alpha] \stackrel{\cong}{\Longrightarrow} K[\alpha]$. The definition of evaluation homomorphisms shows

$$
\overline{\mathrm{ev}}_{\alpha}-1(\alpha)=X+\left(f_{K}^{\alpha}\right) \quad \text { and } \quad \overline{\operatorname{ev}_{\beta}}\left(X+\left(f_{K}^{\alpha}\right)\right)=\beta,
$$

hence $\phi(\alpha)=\beta$.
(iii) is immediate from (ii). This finishes the proof.
I.2.6 Example. Let $f \in K[X]$ be a monic irreducible polynomial of degree 2 and put $L=K[X] /(f)$. We write $\alpha:=X+(f) \in L$ which is a zero of $f$ in $L=K[\alpha]$. Writing $f=X^{2}+a X+b$ one calculates (using long division if necessary) that in $L[X]$ :

$$
X^{2}+a X+b=(X-\alpha)(X-(-a-\alpha))
$$

We now distinguish two cases:

$$
\text { (i) } \alpha=-a-\alpha, \quad \text { (ii) } \alpha \neq-a-\alpha \text {. }
$$

In the first case $2 \alpha=-a$ follows. If $2 \neq 0$ in $K$ then $\alpha=-a / 2 \in K$ contradicting the irreducibility of $f$. Hence, the first case is only possible for $\operatorname{char}(K)=2$ and $a=0$. An example of such an irreducible $f$ is $X^{2}+T \in \mathbb{F}_{2}(T)[X]$, a polynomial over the field of rational functions in the variable $T$ with coefficients in $\mathbb{F}_{2}$. (Were $f$ reducible then $(g / h)^{2}=T$ so $g^{2}=T h^{2}$ for certain $g$, $h \in \mathbb{F}_{2}[T]$; comparing degrees shows this is impossible.)

In the remaining case (ii) $f$ has two distinct zeros in $L$, hence \#Aut ${ }_{K}(L)=2$ by TheoremI.2.5 All groups consisting of two elements are isomorphic, so

$$
\operatorname{Aut}_{K}(L)=\left\{i d_{L}, \phi\right\} \cong \mathbb{Z} / 2 \mathbb{Z}
$$

Here $\phi(\alpha)=-\alpha-\alpha$ and $\phi(x+y \alpha)=\phi(x)+\phi(y) \phi(\alpha)=x+y(-a-\alpha)$ for $x, y \in K)$.
A well known example is $K=\mathbb{R}$ and $f=X^{2}+1$, here $\phi$ is complex conjugation.
Another example is given by the finite fields $L=\mathbb{F}_{p}[X] /\left(X^{2}-d\right)$ with $d$ a nonsquare in $\mathbb{F}_{p}$, in particular we must have $p>2$ here since every $d \in \mathbb{F}_{2}$ is (its own) square. Then

$$
\phi: L \rightarrow L, \quad \phi: x+y \alpha \mapsto x-y \alpha, \quad x, y \in \mathbb{F}_{p}
$$

is the nontrivial field automorphism.
On the other hand Theorem I.2.1 yields an automorphism as well, namely the Frobenius automorphism $F$. We have

$$
F(x+y \alpha)=F(x)+F(y) F(\alpha)=x+y \alpha^{p}
$$

since $x, y \in \mathbb{F}_{p}$, the prime field of $L$. Claim:

$$
F=\phi
$$

As \#Aut ${\underset{F}{p}}(L)=2$ it suffices to show $F \neq i d_{L}$. In case $F=\mathrm{id}_{L}$ the polynomial $X^{p}-X$ would have \# $L=p^{2}$ zeros in $L$, a contradiction. Conclusion: $F \neq \mathrm{id}_{L}$ and therefore $F=\phi$.

From $F=\phi$ one deduces $\alpha^{p}=F(\alpha)=\phi(\alpha)=-\alpha$. As $p>2$ is prime and therefore odd, we write $p=2\left(\frac{p-1}{2}\right)+1$ with $\frac{p-1}{2} \in \mathbb{Z}$. Then

$$
-\alpha=\alpha^{p}=\left(\alpha^{2}\right)^{\frac{p-1}{2}} \alpha=d^{\frac{p-1}{2}} \alpha .
$$

This shows $d^{\frac{p-1}{2}}=-1$ in $\mathbb{F}_{p} \subset L$.

## I. 3 Solving polynomial equations; splitting fields

Let $f \in K[x]$ be a non-constant polynomial. In general, $f$ does not split as a product of linear factors in $K[x]$ because $K$ does not always contain all the solutions of $f(a)=0$. We would like to "find" or to "construct" a larger field which contains all the roots of $f$. This is formalized in the next definition.
I.3.1 Definition. A splitting field $L$ of $f$ over $K$ is a field extension such that:
(a) $f$ splits in $L[x]$ as a product of linear factors.
(b) Let $a_{1}, \ldots, a_{s}$ denote the zeros of $f$ in $L$, then $L=K\left(a_{1}, \ldots, a_{s}\right)$.

## I.3.2 Proposition.

(1) A splitting field exists.
(2) Let $L_{1}, L_{2}$ be two splitting fields for $f$ over $K$. Then there exists a $K$-linear isomorphism of the fields $L_{1}, L_{2}$.

Proof. (1) One uses induction with respect to the degree of $f$. Take a (monic) irreducible factor $g$ (in $K[x]$ ) of $f$. Then $K_{1}:=K[y] /(g(y))$ contains a zero, say $\alpha$ of $g$. Thus $f$ factors in $K_{1}[x]$ as $(x-\alpha) h$. By induction, a splitting field $L$ for $h$ over $K_{1}$ exists. It is easily seen that $L$ is also a splitting field for $f$ over $K$. We note, in passing, that $[L: K]<\infty$.
(2) We will make a proof of a somewhat more general statement:

Two fields $K_{1}$ and $K_{2}$ and an isomorphism $\phi_{0}: K_{1} \rightarrow K_{2}$ are given. Extend $\phi_{0}$ to a ring isomorphism $K_{1}[x] \rightarrow K_{2}[x]$ by $\phi_{0}\left(\sum a_{i} x^{i}\right)=\sum \phi_{0}\left(a_{i}\right) x^{i}$. Suppose $f_{1} \in K_{1}[x]$ and $f_{2} \in K_{2}[x]$ satisfy $\phi_{0}\left(f_{1}\right)=f_{2}$. Let $L_{1}, L_{2}$ denote two splitting fields for $f_{1}$ and $f_{2}$ over
$K_{1}$ and $K_{2}$, respectively. Then $\phi_{0}$ extends to an isomorphism between the field $L_{1}$ and $L_{2}$.

We will construct the desired isomorphism $\phi: L_{1} \rightarrow L_{2}$ step by step. The lowest level, $\phi_{0}: K_{1} \rightarrow K_{2}$ is given. Consider an irreducible factor $g$ of $f_{1}$ and a zero $\alpha \in L_{1}$ of $g$. Then $\phi_{0}(g)$ is an irreducible factor of $\phi_{0}\left(f_{1}\right)=f_{2}$. Hence $\phi_{0}(g)$ splits completely over $L_{2}$. Choose a $\beta$ in $L_{2}$ with $\phi_{0}(g)(\beta)=0$. Define $\phi_{1}: K_{1}(\alpha) \rightarrow K_{2}(\beta)$ by the formula $\sum_{i=0}^{n-1} a_{i} \alpha^{i} \mapsto \sum_{i=0}^{n-1} \phi_{0}\left(a_{i}\right) \beta^{i}$, where all $a_{i} \in K_{1}$ and where $n$ is the degree of $g$. It is easily verified that $\phi_{1}$ is indeed an isomorphism of fields. Now replace $K_{1}, K_{2}, f_{1}, f_{2}$ by $K_{1}(\alpha), K_{2}(\beta), f_{1} /(X-\alpha), f_{2} /(X-\beta)$. The fields $L_{1}, L_{2}$ are splitting fields over $K_{1}(\alpha)$ and $K_{2}(\beta)$ for the polynomials $f_{1} /(X-\alpha)$ and $f_{2} /(X-\beta)$. Induction finishes the proof.
I.3.3 Notation. For a field $K$ and a nonzero $f \in K[x]$ the splitting field of $f$ over $K$ (which by Proposition I.3.2 is unique up to $K$-isomorphisms) is denoted $\Omega_{K}^{f}$.

A polynomial $f \in K[x] \backslash K$ is called separable if the roots of $f$ in any field extension $L$ of $K$ are distinct. It follows from Proposition I.3.2 that it suffices to verify this for a splitting field $L$ of $f$ over $K$.
I.3.4 Corollary. Let $L \supset K$ be the splitting field of a separable nonconstant polynomial $f \in K[x]$. The number of $K$-linear automorphisms of $L$ equals $[L: K]$.

Proof. In the situation given in the proof of part (2) of Proposition I.3.2, we compute the dimensions and count the number of choices for extensions. We work again in the more general situation. Suppose $\phi_{0}: K_{1} \rightarrow K_{2}$ is given. The polynomials $f_{1}$ and $f_{2}$ are by assumption separable. Thus $g$ and $\phi_{0}(g)$ are separable. It is obvious that $\phi_{1}: K_{1}(\alpha) \rightarrow L_{2}$ should map $\alpha$ to a root $\beta$ of $\phi_{0}(g)$. The number of possibilities for $\beta$ is the degree of $\phi_{0}(g)$ (which equals the degree of $g$ ). Thus for $\phi_{1}$ there are $\operatorname{deg}(g)$ possibilities. One has $\left[L_{1}: K_{1}\right]=\left[L_{1}: K_{1}(\alpha)\right]\left[K_{1}(\alpha): K_{1}\right]$. By induction the number of extensions $L_{1} \rightarrow L_{2}$ of a given $\phi_{1}$ is equal to [ $L_{1}: K_{1}(\alpha)$ ]. This completes the proof.

1. Let $K$ be a field. Show that the map

$$
\phi: K(X) \longrightarrow K(X), \quad f \mapsto f(X+1)
$$

is a field automorphism. Find the order of $\phi$ in the $\operatorname{group} \operatorname{Aut}(K(X))$. (The answer depends on the characteristic $\operatorname{char}(K) \ldots$ )
2. Describe a splitting field of $X^{2}-101$ over $\mathbb{Q}$.
3. Let $L$ be a splitting field of $f$ over $K$ and $f=\prod_{i=1}^{n}\left(X-\alpha_{i}\right)$ in $L[X]$. Prove that $L=K\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n-1}\right)$ (so $\alpha_{n}$ is omitted!).
4. Suppose $f \in K[X]$ is monic of degree $n$. Prove: $\left[\Omega_{K}^{f}: K\right]$ divides $n$ ! (Hint: use the construction in the proof of I.3.2.)
5. Prove that $L=\mathbb{Q}(\sqrt[4]{2}, i)$ is a splitting field of $X^{4}-2$ over $\mathbb{Q}$. Determine $[L: \mathbb{Q}]$ and $\sharp \operatorname{Aut}(L)$.
Prove that $\operatorname{Aut}_{\mathbb{Q}(i)}(L) \cong \mathbb{Z} / 4 \mathbb{Z}$.
6. Let $\zeta \in \mathbb{C}$ be a zero of $f=X^{4}+X^{3}+X^{2}+X+1$. Show that $\zeta^{5}=1$ and that $\zeta^{2}, \zeta^{3}, \zeta^{4}$ are zeros of $f$ as well. Prove that $\mathbb{Q}(\zeta)$ is a splitting field of $f$ over $\mathbb{Q}$. Determine $\operatorname{Aut}(\mathbb{Q}(\zeta))$.
7. Prove that $\Omega_{\mathbb{Q}}^{X^{2}-2} \neq \Omega_{\mathbb{Q}}^{X^{2}-3}$ and that $\Omega_{K}^{X^{2}-\overline{2}} \simeq \Omega_{K}^{X^{2}-\overline{3}}$ for $K=\mathbb{F}_{5}$.
8. Prove that $\mathbb{Q}(i, \sqrt{2})$ is a splitting field of $f_{\mathbb{Q}}^{i+\sqrt{2}}$ over $\mathbb{Q}$.

Show that $\operatorname{Aut}(\mathbb{Q}(i, \sqrt{2})) \simeq \mathbb{Z} / 2 \mathbb{Z} \times \mathbb{Z} / 2 \mathbb{Z}$.
9. Let $L$ be a splitting field of $f$ over $K$ and put $n=\operatorname{deg}(f)$.
(a) Prove: any $K$-automorphism of $L$ permutes the zeros of $f$ in $L$,
(b) Prove: the group $\operatorname{Aut}_{K}(L)$ is isomorphic to a subgroup of $S_{n}$;
(c) Show that \#Aut ${ }_{K}(L)$ is a divisor of $n$ !.
10. Prove: $\operatorname{Aut}(\mathbb{Q}(\sqrt[3]{2}, i \sqrt{3})) \simeq S_{3}$.
11. Take $f=X^{3}+X^{2}-2 X-1 \in \mathbb{Q}[X]$ and let $\alpha \in \Omega_{\mathbb{Q}}^{f}$ be a zero of $f$.

Compute $f_{\mathbb{Q}}^{\alpha^{2}-2}$ and show that $\mathbb{Q}(\alpha)=\Omega_{\mathbb{Q}}^{f}$. Deduce that $\operatorname{Aut}\left(\Omega_{\mathbb{Q}}^{f}\right) \cong \mathbb{Z} / 3 \mathbb{Z}$.

## GALOIS THEORY

## II. 1 Galois extensions and examples

For our purposes a pleasant definition of a Galois extension $K \subset L$ is:
II.1.1 Definition. $L \supset K$ is called a Galois extension of $K$ if $L$ is a splitting field of a separable polynomial $f$ over $K$.
II.1.2 Definition. The Galois group $\operatorname{Gal}(L / K)$ of the extension $K \subset L$ is the group of all $K$-linear automorphisms of $L$.

Observe that indeed $\operatorname{Gal}(L / K)$ is a group, with composition of automorphisms as group law and the identity automorphism as the unit element.
II.1.3 Lemma. Let $K \subset L$ be a Galois extension and suppose that $a \in L$ is invariant under the action of $\operatorname{Gal}(L / K)$. Then $a \in K$.

Proof. From Corollary I.3.4 we know that $[L: K]=\# \operatorname{Gal}(L / K)$. Let $f$ be a polynomial in $K[x]$ such that $L$ is its splitting field over $K$. Observe that $L$ is also the splitting field of $f$ over the field $K(a)$. The assumption that $a \in L$ is invariant implies $\operatorname{Gal}(L / K)=\operatorname{Gal}(L / K(a))$. Thus $[L: K]=[L: K(a)]$ and hence $[K(a): K]=1$, which means that $a \in K$.
II.1.4 Corollary. Let $K \subset L$ be a Galois extension and $a \in L$. The minimal polynomial $F$ of a over $K$ is separable and all its roots are in $L$.

Proof. Let the orbit $\operatorname{Gal}(L / K) a=\{g a \mid$ for all $g \in \operatorname{Gal}(L / K)\}$ be $\left\{a_{1}, \ldots, a_{s}\right\}$. Consider the polynomial $G:=\left(x-a_{1}\right) \cdots\left(x-a_{s}\right)=x^{s}+b_{s-1} x^{s-1}+\cdots+b_{1} x+b_{0}$ in $L[x]$. This polynomial is invariant under the action of $\operatorname{Gal}(L / K)$ and hence Lemma II.1.3 implies that $G \in K[x]$. Clearly $F$ divides $G$ and therefore it has the required properties.
II.1.5 Proposition. Let $K \subset L$ be a finite extension. The following are equivalent:
(1) $K \subset L$ is a Galois extension.
(2) For every element $a \in L$, the minimal polynomial $F \in K[x]$ of a has the property that all its roots lie in $L$ and are simple.

Proof. (1) $\Rightarrow(2)$ is the statement of Corollary II.1.4.
(2) $\Rightarrow$ (1). Take elements $a_{1}, \ldots, a_{s} \in L$ such that $L=K\left(a_{1}, \ldots, a_{s}\right)$. Let $F_{i}$ be the minimal polynomial of $a_{i}$ over $K$. We may suppose that $F_{1}, \ldots, F_{t}$ are the distinct
elements in $\left\{F_{1}, \ldots, F_{s}\right\}$. Put $F=F_{1} \cdots F_{t}$. Then each $F_{i}$ is separable and has all its roots in $L$. Then also $F$ is separable and all its roots are in $L$. Moreover the set of the roots of $F$ contains $\left\{a_{1}, \ldots, a_{s}\right\}$. Thus $L$ is the splitting field of $F$ over $K$.

In several textbooks on Galois theory, part (2) of the above proposition is used as the definition of a Galois extension. To be more precise: a finite extension $L / K$ is called normal if for every $a \in L$ the minimal polynomial $F$ of $a$ over $K$ splits in $L[x]$ as a product of linear factors. The extension is called separable if for every $a \in L$, the minimal polynomial $F$ of $a$ over $K$ is separable. The extension is called Galois if it is both normal and separable.
II.1.6 Remark. Let $f \in K[x]$ be a separable polynomial and let $L \supset K$ be a splitting field. The roots of $f$ in $L$ are say $\left\{a_{1}, \ldots, a_{n}\right\}$; note that $n=\operatorname{deg}(f)$. Every $\sigma \in \operatorname{Gal}(L / k)$ permutes this set. Thus we find a homomorphism $\operatorname{Gal}(L / K) \rightarrow S_{n}$. This homomorphism is injective. Indeed, if $\sigma\left(a_{i}\right)=a_{i}$ for all $i$, then $\sigma$ is the identity since $L=K\left(a_{1}, \ldots, a_{n}\right)$.
II.1.7 Example. The Galois group of the splitting field of $x^{4}-2$ over $\mathbb{Q}$ is in this way isomorphic with the subgroup

$$
\{(1),(24),(1234),(12)(34),(13)(24),(13),(1432),(14)(23)\}
$$

of $S_{4}$. Here we identify $k \in\{1,2,3,4\}$ with the zero $a_{k}=i^{k-1} \sqrt[4]{2}$ of $x^{4}-2$. A permutation sending $k$ to $\ell$ corresponds to an automorphism sending $a_{k}$ to $a_{\ell}$.

We now give some more examples and elementary properties.
II.1.8 Example. $\mathbb{Q}(\sqrt[3]{2}) / \mathbb{Q}$ is not normal.

This follows using Proposition II.1.5. Indeed, the minimal polynomial of $\sqrt[3]{2}$ over $\mathbb{Q}$ is $x^{3}-2$. One can consider $\mathbb{Q}(\sqrt[3]{2})$ as a subfield of $\mathbb{R}$. The other two zeros of $x^{3}-2$ in $\mathbb{C}$ are $\omega \sqrt[3]{2}$ and $\omega^{2} \sqrt[3]{2}$, where $\omega=e^{2 \pi i / 3}$. They do not lie in $\mathbb{R}$ and hence also not in $\mathbb{Q}(\sqrt[3]{2})$.
II.1.9 Example. If $[L: K]=2$ and the characteristic of $K$ is $\neq 2$, then $L / K$ is Galois. Indeed, choose $\alpha \in L \backslash K$. Its minimal polynomial $F$ has the form $F=x^{2}+a x+b$. The polynomial $F$ splits in $L[x]$ as $(x-\alpha)(x-\beta)$ with $\beta+\alpha=-\alpha$. If $\alpha$ would equal $\beta$, then one finds the contradiction $\alpha=-\frac{a}{2} \in K$. It follows that $L$ is the splitting field of the separable polynomial $F$ over $K$.
II.1.10 Example. If the characteristic of $K$ is 0 , then every finite $L / K$ is separable. To see this, choose $a \in L$. Its minimal polynomial $F \in K[x]$ is irreducible. The derivative $F^{\prime}=\frac{d F}{d x}$ of $F$ is not 0 and thus the g.c.d. of $F$ and $F^{\prime}$ is 1 . By Exercise 4 of Section I.3, one has that $F$ is separable.
II.1.11 Example. A field $K$ of characteristic $p>0$ is called perfect if every element of $K$ is a pth power. Every finite extension $L / K$ of a perfect field is separable.
Proof: choose $a \in L$ with minimal polynomial $F \in K[x]$. If $F$ is not separable then the g.c.d. of $F$ and $F^{\prime}$ is not 1 . Since $F$ is irreducible this implies that $F^{\prime}=0$. Thus $F$ contains only $p$ the powers of $x$, i.e., $F=\sum a_{n} x^{p n}$ with $a_{n} \in K$. Each $a_{n}$ is written as $b_{n}^{p}$ with $b_{n} \in K$. Then $F=\left(\sum b_{n} x^{n}\right)^{p}$, which contradicts the fact that $F$ is irreducible.
II.1.12 Example. Every finite field is perfect. The field $\mathbb{F}_{p}(t)$ is not perfect. Indeed, the finite fields of characteristic $p$ are the $\mathbb{F}_{q}$ with $q$ a power of $p$. The "Frobenius map" Fr:z $\rightarrow z^{p}$ on any field of characteristic $p>0$ is a homomorphism
of the additive group. The kernel is 0 since $z^{p}=0$ implies $z=0$. By counting, one sees that $F r: \mathbb{F}_{q} \rightarrow \mathbb{F}_{q}$ is bijective. And $F r: \mathbb{F}_{p}(t) \rightarrow \mathbb{F}_{p}(t)$ is not surjective since $t$ is not in the image.
II.1.13 Example. Put $K=\mathbb{F}_{p}(t)$ and let $F=x^{p}-t$. The splitting field of $F$ is not separable over $K$.
Indeed, the derivative of $F$ is 0 . We note that the splitting field in question can be identified with $\mathbb{F}_{p}(x)$.

## II. 2 Galois correspondence and primitive elements

Let a Galois extension $L / K$ with Galois group $G$ be given. One considers two sets, $\mathscr{M}$, the set of the intermediate fields, i.e., the fields $M$ with $K \subset M \subset L$ and the set $\mathscr{G}$ of the subgroups of $G$.

There are two maps between those sets, $\alpha: \mathscr{M} \rightarrow \mathscr{G}$, defined by $\alpha(M)=G a l(L / M)$ and $\beta: \mathscr{G} \rightarrow \mathscr{M}$, defined by $\beta(H)=L^{H}$, i.e., the subfield of $L$ consisting of the elements which are invariant under the action of $H$. It is obvious that the two maps reverse inclusions. What is often called the "main theorem of Galois theory" is
II.2.1 Theorem (The Galois correspondence).
(1) $\alpha$ and $\beta$ are each others inverse.
(2) The subgroup $H \in \mathscr{G}$ is normal if and only if $\beta(H)$ is a normal (or Galois) extension of $K$.
(3) Suppose that $H$ is a normal subgroup of $G$. Then $M / K$ is a Galois extension with Galois group G/H.

We will first prove a lemma.
II.2.2 Lemma. Suppose that the finite extension $L / K$ has the property that there are only finitely many intermediate fields. Then there is an $\alpha \in L$ with $L=K(\alpha)$.

Proof. For a finite field $K$ the proof is quite easy. $L$ is also a finite field and it is known that the multiplicative group $L^{*}$ of $L$ is cyclic, i.e., there is an element $\xi$ with $L^{*}=\left\{\xi^{n} \mid n \in \mathbb{Z}\right\}$. Clearly $L=K(\xi)$. Suppose now that $K$ is infinite and that $L \neq K$. Let $n \geq 1$ be minimal such that $L=K\left(a_{1}, \ldots, a_{n}\right)$ for certain elements $a_{1}, \ldots, a_{n}$. We have to show that $n=1$. Suppose $n \geq 2$ and consider for every $\lambda \in K$ the element $b_{\lambda}:=a_{1}+\lambda a_{2}$. The fields $K\left(b_{\lambda}\right)$ are intermediate fields for $L / K$ and thus there are $\lambda_{1} \neq \lambda_{2}$ with $M:=K\left(b_{\lambda_{1}}\right)=K\left(b_{\lambda_{2}}\right)$. The field $M$ contains $a_{1}+\lambda_{1} a_{2}$ and $a_{1}+\lambda_{2} a_{2}$. It follows that $M=K\left(a_{1}, a_{2}\right)$ and that $L=K\left(b_{\lambda_{1}}, a_{3}, \ldots, a_{n}\right)$. This contradicts the minimality of $n$.

Proof. of Theorem II.2.1
(1) $\beta \alpha(M)$ is the field $L^{\text {Gal(L/M) }}$. Applying Lemma II.1.3 to the Galois extension $L / M$ implies that $L^{G a l(L / M)}=M$. This implies that $\alpha$ is injective. Since $\mathscr{G}$ is finite, also $\mathscr{M}$ is finite. Take $H \in \mathscr{G}$ and put $M=\beta(H)=L^{H}$. Clearly $H \subset \alpha \beta(H)$. We have to prove equality. According to Lemma II.2.2 one can write $L=M(a)$ for some $a \in L$. Consider the polynomial $G=\prod_{\sigma \in H}(x-\sigma(a)) \in L[x]$. This polynomial is invariant under the action of $H$. Therefore its coefficients are also invariant under $H$ and belong to $M$. Then $G \in M[x]$ and the minimal polynomial of $a$ over $M$ divides $G$. Thus $[L: M] \leq \# H$ and $\# G a l(L / M) \leq \# H$. Since $H$ is a subgroup of $G a l(L / M)$ one finds the required equality $H=G a l(L / M)$.
(2) The subgroup $H$ corresponds to $M=\beta(H)=L^{H}$. For any $\sigma \in G$ one has

$$
\beta\left(\sigma H \sigma^{-1}\right)=\sigma(M)
$$

Thus $H$ is normal if and only if $\sigma(M)=M$ for all $\sigma \in G$. Using that $L / K$ is normal and separable, one finds that the latter property of $M$ is equivalent with $M / K$ normal (or also Galois).
(3) $H$ is supposed to be normal and thus $M:=\beta(H)=L^{H}$ satisfies $\sigma(M)=M$ for all $\sigma \in G$. Thus one can define a restriction map $\operatorname{Gal}(L / K) \rightarrow G a l(M / K)$. The kernel of this homomorphism is $H$. The map is surjective since $\#(G / H)$ is equal to $[M: K]$.
II.2.3 Remark. The proof of II.2.1 2 ) shows a small fact that deserves to be noted. Namely, if $L / K$ is a finite Galois extension and $\sigma \in G:=\operatorname{Gal}(L / K)$, then for every intermediate field $M$, so $K \subset M \subset L$, also $\sigma(M)$ is an intermediate field. Moreover since $\sigma$ is a $K$-linear automorphism, $\operatorname{dim}_{K}(M)=\operatorname{dim}_{K}(\sigma(M))$.

If this field $M$ equals $L^{H}$ for some subgroup $H \subset G$, then $L / M$ is Galois and $H=\operatorname{Gal}(L / M)$, and in particular $\# H=[L: M]=[L: K] / \operatorname{dim}_{K}(M)$.

The equality $\beta\left(\sigma H \sigma^{-1}\right)=\sigma(M)$ states that the subgroup $\sigma H \sigma^{-1}$ via the Galois correspondence is associated to the subfield $\sigma(M)$. In other words, the action of $\sigma$ on the set $\mathscr{M}$ of intermediate fields yields via the Galois correspondence the action 'conjugate by $\sigma$ ' on the set $\mathscr{G}$ of subgroups.
II.2.4 Corollary (The theorem of the primitive element). For every finite separable extension $L / K$ there is a cyclic element, i.e., an element $a \in L$ with $L=K(a)$.

Proof. According to Lemma II.2.2, it suffices to show that there are finitely many intermediate fields for $L / K$. If we can show that $L$ lies in the splitting field $\tilde{L}$ of a separable polynomial, then there are only finitely many intermediate fields for $\tilde{L} / K$ and then also for $L / K$. Let $L / K$ be generated by elements $a_{1}, \ldots, a_{s}$. The minimal polynomial of $a_{i}$ over $K$ is denoted by $F_{i}$. We may suppose that $F_{1}, \ldots, F_{t}$ are the distinct minimal polynomials, then $F:=F_{1} \cdots F_{t}$ is separable and every $a_{i}$ is a zero of $F$. The splitting field $\tilde{L}$ of $F$ contains $L$ and we are done.

## II. 3 Exercises

1. Let $K \subset L$ be fields and $a \in L$. Show that $K[a]$ is a field if and only if $a$ is algebraic over $K$.
2. Let $K \subset L$ be fields and $S \subset L$ a non-empty finite subset. Show that $K[S]$ is a field if and only if every element of $S$ is algebraic over $K$.
Find a counterexample with $S=L$ to show that the condition that $S$ is finite cannot be missed.
3. Let $K \subset L$ be a finitely generated field extension. Prove the existence of an intermediate field $M$ such that $K \subset M$ is purely transcendental (or $K=M$ ) and $[L: M]<\infty$.
4. Prove that $f \in K[x] \backslash K$ is separable if and only if $f$ and its derivative $f^{\prime}=\frac{d}{d x} f$ are relatively prime.
5. Let $L$ be a splitting field of $f$ over $K$. Let $n$ be the degree of $f$. Prove that $[L: K] \leq n!$. Try to make examples with $K=\mathbb{Q}$ and $f$ of degree 3 with $[L: \mathbb{Q}]=6$.
6. Let $L$ be the splitting field over $\mathbb{Q}$ of the polynomial $x^{3}-3$. Produce an explicit splitting field $L \subset \mathbb{C}$ and find the ( $\mathbb{Q}$-linear) field automorphisms of $L$.
7. The same question for the polynomial $x^{8}-1$ over $\mathbb{Q}$.
8. Show that a finite extension of a finite field is Galois.
9. Show that $\mathbb{F}_{p^{n}}$ is the splitting field of the polynomial $x^{p^{n}}-x$ over $\mathbb{F}_{p}$. Prove that $\mathbb{F}_{p^{n}} / \mathbb{F}_{p}$ is Galois and that its Galois group is $\left\{1, F r, \ldots, F r^{n-1}\right\}$, where $F r$ is the Frobenius map on $\mathbb{F}_{p^{n}}$ defined by $\operatorname{Fr}(z)=z^{p}$.
10. Show that $\mathbb{Q}(\sqrt[4]{3}, i)$ is a Galois extension of $\mathbb{Q}(i)$ and compute the corresponding Galois group.
11. Determine the Galois group of $x^{12}-2$ over $\mathbb{Q}$.
12. Take $q=p^{n}$ for some prime $p$, and let $K \supset \mathbb{F}_{q}$ and $a \in K$, with the property that the polynomial $X^{q}-X+a$ does not split completely in $K[X]$. Let $\alpha$ be a zero of this polynomial in some splitting field.
(a) Show that $K(\alpha) \neq K$.
(b) Show that the extension $K(\alpha) \supset K$ is Galois.
(c) Show that every $\sigma \in \operatorname{Gal}(K(\alpha) / K)$ satisfies $\sigma(\alpha)=\alpha+t$ for some $t \in \mathbb{F}_{q}$.
(d) Show that $\sigma \mapsto \sigma(\alpha)-\alpha$ defines an injective homomorphism of groups: $\operatorname{Gal}(K(\alpha) / K) \rightarrow\left(\mathbb{F}_{q},+, 0\right)$.
(e) In the special case that $K$ is itself a finite field, observe that the Galois group of a finite extension of finite fields is cyclic, and deduce that $X^{q}-X+a$ factors as a product of $q / p$ distinct irreducible polynomials of degree $p$ in $K[X]$.
(f) Take $K=\mathbb{F}_{q}(t)$ and $X^{q}-X+t$. Explain why this polynomial is irreducible in $K[X]$, and determine the Galois group of its splitting field over $K$.
13. Prove that $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ is a Galois extension of $\mathbb{Q}$. Determine the Galois group, all intermediate subfields and a primitive element.
14. The same questions for $\mathbb{Q}\left(e^{2 \pi i / 3}, \sqrt[3]{3}\right) / \mathbb{Q}$.
15. Find all subfields of $\mathbb{Q}(\sqrt[4]{2}, i)$ and a primitive element for each of them. Which of these fields are normal over $\mathbb{Q}$ ?
16. One of the theorems in the Ph.D. thesis of Amol Sasane (Groningen, 2001, advisor Prof. dr. R. Curtain) states that $\tan (\pi / 2001) \notin \mathbb{Q}$. Prove this theorem.
17. (a) Prove that for an odd prime number $p$, the field $\mathbb{Q}\left(\zeta_{p}\right)$ has a unique subfield $K$ with $[K: \mathbb{Q}]=2$.
(b) Find a condition on $p$ such that $K \subset \mathbb{R}$.

Hint: complex conjugation is an element of $\operatorname{Gal}\left(\mathbb{Q}\left(\zeta_{p}\right) / \mathbb{Q}\right)$. The question whether or not $K$ is a real field is the same as the question whether complex conjugation is an element of the subgroup $\operatorname{Gal}\left(\mathbb{Q}\left(\zeta_{p}\right) / K\right)$.
(c) Write $\varepsilon(n):=1$ if $n \bmod p$ is a square in $\mathbb{F}_{p}^{*}$ and $\varepsilon(n):=-1$ otherwise. Show that $K / \mathbb{Q}$ is generated by the element $\sum_{n=1}^{p-1} \varepsilon(n) \zeta_{p}^{n}$.
18. Prove that the regular 7 -gon cannot be constructed by ruler and compass. Explanation: The admissible operations with ruler and compass are: drawing a line through two points, drawing a circle with given center and radius, intersecting two lines, interseting a line with a circle and intersecting two circles. Hint: Identify the plane with $\mathbb{C}$. The points $\mathbb{Q} \subset \mathbb{C}$ are given. Prove
that every "constructable" point $a \in \mathbb{C}$ gives rise to a "tower of field extensions" $K_{n}:=\mathbb{Q}(a) \supset K_{n-1} \supset \cdots \supset K_{0}=\mathbb{Q}$ with $\left[K_{i+1}: K_{i}\right]=2$ for every $i$.
19. Cyclic extensions.

Let $K$ be a field and $n>1$ an integer. Suppose that the characteristic of $K$ is 0 or $p$ with $p \not X n$ and that $K$ contains all the $n^{\text {th }}$ roots of unity. In this exercise we want to prove the following statement:
$E \supset K$ is a Galois extension with a cyclic Galois group of order $n$ if and only if $E=K(\alpha)$ where $\alpha$ is a root of an irreducible polynomial $x^{n}-a \in K[x]$.
(a) Let $f(x)=x^{n}-a \in K[x]$ be an irreducible polynomial. Show that $f$ is separable. Show that the splitting field $E$ of $f(x)$ over $K$ is of the form $K(\alpha)$ with $\alpha$ a root of $f(x)=0$. Furthermore, show that the Galois group of $E$ over $K$ is generated by the map defined by $\alpha \mapsto \zeta \alpha$ where $\zeta$ is a primitive $n^{\text {th }}$ root of unity. (b) Let $E$ be a Galois extension of $K$ with a cyclic Galois group of order $n$. Let $\sigma$ generate the Galois group.
(i) One considers $\sigma$ as a $K$-linear map on $E$ as vector space over $K$. Prove that every eigenvalue $\lambda$ of $\sigma$ satisfies $\lambda^{n}=1$ and thus belongs to $K$.
(ii) Prove that there is a basis of eigenvectors for $\sigma$. (Hint: Jordan normal form).
(iii) Prove that every eigenvalue has multiplicity 1. (Hint: if $\sigma e_{i}=\lambda e_{i}$ for $i=1,2$ and $e_{1} \neq 0 \neq e_{2}$, then $\left.\sigma\left(\frac{e_{1}}{e_{2}}\right)=\frac{e_{1}}{e_{2}}\right)$.
(iv) Prove that there is an $\alpha \in E$ with $\alpha \neq 0, \sigma(\alpha)=\zeta \alpha$ and $\zeta$ a primitive $n^{\text {th }}$ root of unity.
(v) Show that the $\sigma^{i}(\alpha), i=0, \ldots, n-1$, are all distinct and that therefore the minimal polynomial of $\alpha$ over $K$ is $x^{n}-a$ with $a=\alpha^{n} \in K$.

## ROOTS OF UNITY

## III. 1 Cyclotomic fields over the rationals

Let $n \geq 1$ be an integer and write $\zeta=\zeta_{n}=e^{2 \pi i / n} \in \mathbb{C}$. The subfield $\mathbb{Q}(\zeta) \subset \mathbb{C}$ is the splitting field of $x^{n}-1$ over $\mathbb{Q}$, since all the roots of $x^{n}-1$ are $\zeta^{k}, k=0,1, \ldots, n-1$. For $n \geq 3$, one calls $\mathbb{Q}(\zeta)$ the nth cyclotomic field. Since $x^{n}-1 \in \mathbb{Q}[x]$ is separable, $\mathbb{Q}(\zeta) \supset \mathbb{Q}$ is a Galois extension. We will determine the degree and the Galois group of this extension.

The minimal polynomial $\Phi_{n}$ of $\zeta$ over $\mathbb{Q}$ is called the $n$th cyclotomic polynomial. An explicit general formula for $\Phi_{n}$ is not available. Still we will need some information on $\Phi_{n}$ in order to calculate the Galois group of $\mathbb{Q}(\zeta) / \mathbb{Q}$. Two tools from standard courses in algebra that we will use are:
(1) Lemma of Gauss: Let $f \in \mathbb{Z}[x]$ be monic and let $g, h \in \mathbb{Q}[x]$ be monic, too. Then $f=g h$ implies that $g, h \in \mathbb{Z}[x]$.
(2) Eisenstein's criterion: Let $f=a_{n} x^{n}+\cdots+a_{0} \in \mathbb{Z}[x]$ be such that all its coeffcients, with the exception of $a_{n}$, are divisible by a prime number $p$ and $p^{2}$ does not divide $a_{0}$. Then $f$ is irreducible in $\mathbb{Z}[x]$ and in $\mathbb{Q}[x]$.
The proof of both statements uses "reduction modulo $p$ ", i.e., the ring homomorphism $\mathbb{Z}[x] \rightarrow \mathbb{F}_{p}[x]$, given by $f=\sum a_{n} x^{n} \mapsto \bar{f}:=\sum \bar{a}_{n} x^{n}$, where for $a \in \mathbb{Z}$ one has written $\bar{a}$ for its image in $\mathbb{F}_{p}=\mathbb{Z} / p \mathbb{Z}$. Reduction modulo $p$ will also be used in the proof of the next proposition, as well as the identity $\bar{f}\left(x^{p}\right)=\bar{f}^{p}$ for $\bar{f} \in \mathbb{F}_{p}[x]$.

## III.1.1 Proposition.

(1) $\Phi_{n} \in \mathbb{Z}[x]$.
(2) For a prime number $p$ one has

$$
\Phi_{p}=\frac{x^{p}-1}{x-1}=x^{p-1}+x^{p-2}+\cdots+x+1 .
$$

(3) For any integer $n \geq 1$ one has

$$
\Phi_{n}=\prod_{1 \leq j<n, \operatorname{gcd}(j, n)=1}\left(x-\zeta_{n}^{j}\right) .
$$

This means that the zeros of $\Phi_{n}$ are precisely all elements of order $n$ in $\mathbb{C}^{\times}$. Moreover, the degree of $\Phi_{n}$ is $\phi(n)$, where $\phi$ is Euler's phi-function defined by $\phi(n)=\#(\mathbb{Z} / n \mathbb{Z})^{\times}$.

Proof. (1) follows from the Lemma of Gauss applied to $f=X^{n}-1$ and $g=\Phi_{n}$.
(2) follows from (3), however, we now give a direct proof. Write $f(x)=x^{p-1}+\ldots+x+1$ and make the substitution $x=1+t$. Then

$$
f(1+t)=\frac{(1+t)^{p}-1}{(1+t)-1}=\sum_{i=1}^{p}\binom{p}{i} t^{i-1}
$$

satisfies Eisenstein's criterion and is therefore irreducible. It follows that $f$ itself is irreducible, and since $f=\left(x^{p}-1\right) /(x-1)$, the number $\zeta_{p}$ is a zero of $f$. Hence $f=\Phi_{p}$ which is what we wanted to prove.
(3) $x^{n}-1=\Pi_{1 \leq j \leq n}\left(x-\zeta^{j}\right)$ and $\Phi_{n}$ can only contain the factors $\left(x-\zeta^{j}\right)$ with $\operatorname{gcd}(j, n)=1$. Indeed, $\operatorname{gcd}(j, n)=d>1$ implies that $\zeta^{j}$ is a root of $x^{n / d}-1$. From $\Phi_{n}\left(\zeta^{j}\right)=0$ it would follow that $\Phi_{n}$ divides $x^{n / d}-1$ which leads to the contradiction $\zeta^{n / d}=1$.

In order to see that every $x-\zeta^{j}$ with $\operatorname{gcd}(j, n)=1$ is a factor of $\Phi_{n}$, we use a trick. Decompose $j$ as a product $p_{1} \cdots p_{t}$ of (not necessarily distinct) prime factors. The $p_{i}$ do not divide $n$ since $\operatorname{gcd}(j, n)=1$. Write $\zeta^{j}=\left(\cdots\left(\left(\zeta^{p_{1}}\right)^{p_{2}}\right) \cdots\right)^{p_{t}}$. We claim that the following statement holds:

$$
{ }^{(*)} \text { if } p \text { does not divide } n \text { and if } \Phi(\alpha)=0 \text {, then } \Phi\left(\alpha^{p}\right)=0 \text {. }
$$

Using this assertion one finds

$$
\Phi_{n}(\zeta)=0 \Rightarrow \Phi_{n}\left(\zeta^{p_{1}}\right)=0 \Rightarrow \Phi_{n}\left(\zeta^{p_{1} p_{2}}\right)=0 \Rightarrow \cdots \Rightarrow \Phi_{n}\left(\zeta^{j}\right)=0 .
$$

We will prove (*) by deriving a contradiction from the assumptions: $p \nmid n$ and $\Phi_{n}(\alpha)=0$ and $\Phi_{n}\left(\alpha^{p}\right) \neq 0$.

The equality $x^{n}-1=\Phi_{n} \cdot f$ with $f\left(\alpha^{p}\right)=0$ is clear. Now $f\left(\alpha^{p}\right)=0$ means that $\alpha$ is a zero of $f\left(x^{p}\right)$. Therefore $\Phi_{n}$ divides $f\left(x^{p}\right)$. According to the Lemma of Gauss, this division takes place in the ring $\mathbb{Z}[x]$. Hence $\Phi_{n}\left(x^{p}\right) \Phi_{n}(x)$ divides $x^{p n}-1$ in the ring $\mathbb{Z}[x]$. After reduction modulo $p$, one finds that $\bar{\Phi}_{n}^{p+1}$ divides $x^{p n}-1=\left(x^{n}-1\right)^{p}$ in $\mathbb{F}_{p}[x]$. However $p \nmid n$ and $x^{n}-1$ has only simple roots. The multiplicity of any root of $\left(x^{n}-1\right)^{p}$ is $p$ and this is the contradiction that we wanted to find.

The proof of the following corollary is left to the reader.
III.1.2 Corollary. The Galois group $\operatorname{Gal}\left(\mathbb{Q}\left(\zeta_{n}\right) / \mathbb{Q}\right)$ is isomorphic to the group $(\mathbb{Z} / n \mathbb{Z})^{\times}$. This isomorphism identifies the unit $a \bmod n$ with the field automorphism sending $\zeta_{n}$ to $\zeta_{n}^{a}$.
III.1.3 Proposition (Formulas for $\Phi_{n}$ ).

1. $x^{n}-1=\prod_{d \mid n} \Phi_{d}$ (i.e., the product over all divisors $d$ of $n$ ).
2. $\Phi_{n}=\prod_{d \mid n}\left(x^{d}-1\right)^{\mu(n / d)}$, where $\mu$ is the Möbius function given by $\mu(1)=1, \mu(n)=0$ if $n$ contains a square $\neq 1$ and $\mu\left(p_{1} \cdots p_{t}\right)=(-1)^{t}$ if $p_{1}, \ldots, p_{t}$ are distinct primes.
3. $\Phi_{n p}(x)=\Phi_{n}\left(x^{p}\right)$ if the prime $p$ divides $n$.
4. $\Phi_{n p}(x)=\Phi_{n}\left(x^{p}\right) \Phi_{n}(x)^{-1}$ if the prime $p$ does not divide $n$.

Proof. (1) follows since both polynomials are monic and by Proposition III.1.1(3) have the same zeros.
(2) is a special case of the "Möbius inversion". This is the statement

$$
\begin{aligned}
& \text { If for all } n \geq 1 \text { the formula } f_{n}=\prod_{d \mid n} g_{d} \text { holds, then } \\
& \qquad g_{n}=\prod_{d \mid n} f_{d}^{\mu(n / d)} \text { holds for all } n \geq 1 .
\end{aligned}
$$

In the proof of this inversion formula one uses the easily deduced formulas $\sum_{d \mid n} \mu(d)=0$ for $n>1$ and $\sum_{d \mid n} \mu(d)=1$ for $n=1$.
The proofs of 3 . and 4. are left as an exercise.

## III. 2 An application to tangent values

Take $n \in \mathbb{Z}_{\geq 3}$ and consider $t_{n}:=\tan (\pi / n)$. With $\zeta:=\zeta_{4 n}=e^{2 \pi i / 4 n}$ we have $i=\zeta^{n}$ and

$$
t_{n}=\frac{\sin (\pi / n)}{\cos (\pi / n)}=\frac{\zeta^{2}-\zeta^{-2}}{\zeta^{n}\left(\zeta^{2}+\zeta^{-2}\right)} \in \mathbb{Q}(\zeta) .
$$

We will compute the degree $\left[\mathbb{Q}\left(t_{n}\right): \mathbb{Q}\right]$ by determining the subgroup of the Galois group $\operatorname{Gal}(\mathbb{Q}(\zeta) / \mathbb{Q}) \cong(\mathbb{Z} / 4 n \mathbb{Z})^{\times}$corresponding to the intermediate field $\mathbb{Q}\left(t_{n}\right) \subset \mathbb{Q}(\zeta)$.

To do so, first note that $\bar{a} \in(\mathbb{Z} / 4 n \mathbb{Z})^{\times}$yields the automorphism $\sigma_{a}$ of $\mathbb{Q}(\zeta)$ defined by $\zeta \mapsto \zeta^{a}$. Its effect on $t_{n}$ is given as

$$
t_{n} \mapsto \sigma_{a}\left(t_{n}\right)=\frac{\zeta^{2 a}-\zeta^{-2 a}}{\zeta^{a n}\left(\zeta^{2 a}+\zeta^{-2 a}\right)}=\left\{\begin{array}{rll}
\tan (\pi a / n) & \text { if } & a \equiv 1 \bmod 4 ; \\
-\tan (\pi a / n) & \text { if } & a \equiv 3 \bmod 4
\end{array}\right.
$$

So $\sigma_{a}$ is the identity when restricted to $\mathbb{Q}\left(t_{n}\right)$ precisely when

$$
\begin{aligned}
& \text { either } \bar{a} \in(\mathbb{Z} / 4 n \mathbb{Z})^{\times} \text {satisfies } a \equiv 1 \bmod 4 \text { and } \tan (\pi / n)=\tan (\pi a / n), \\
& \text { or } \bar{a} \in(\mathbb{Z} / 4 n \mathbb{Z})^{\times} \text {satisfies } a \equiv 3 \bmod 4 \text { and } \tan (\pi / n)=\tan (-\pi a / n)
\end{aligned}
$$

Solving these tangent equalities results in

$$
\left\{\begin{array} { l } 
{ a \equiv 1 \operatorname { m o d } n } \\
{ a \equiv 1 \operatorname { m o d } 4 }
\end{array} \text { or } \left\{\begin{array}{l}
a \equiv-1 \bmod n \\
a \equiv-1 \bmod 4
\end{array}\right.\right.
$$

The number of solutions $\bar{a} \in(\mathbb{Z} / 4 n \mathbb{Z})^{\times}$now depends on the parity of $n$ :
(i). If $n \equiv 0 \bmod 4$ then one finds the condition $a \equiv \pm 1 \bmod n$. This gives precisely 8 pairwise distinct values $\bar{a}$, namely

$$
\bar{a} \in\{\overline{1}, \overline{-1}, \overline{n+1}, \overline{n-1}, \overline{2 n+1}, \overline{2 n-1}, \overline{3 n+1}, \overline{3 n-1}\}
$$

They form a subgroup of order 8 in $(\mathbb{Z} / 4 n \mathbb{Z})^{\times} \cong \operatorname{Gal}(\mathbb{Q}(\zeta) / \mathbb{Q})$ which in fact by the Galois correspondence equals $\operatorname{Gal}\left(\mathbb{Q}(\zeta) / \mathbb{Q}\left(t_{n}\right)\right)$. So in this case by Corollary I.3.4 we have $\left[\mathbb{Q}(\zeta): \mathbb{Q}\left(t_{n}\right)\right]=8$ and therefore using Proposition III.1.1(c) we have

$$
\left[\mathbb{Q}\left(t_{n}\right): \mathbb{Q}\right]=[\mathbb{Q}(\zeta): \mathbb{Q}] / 8=\phi(4 n) / 8=\phi(n / 2) .
$$

(ii). If $n \equiv 2 \bmod 4$ then we obtain the condition $a \equiv \pm 1 \bmod 2 n$, and this yields exactly 4 distinct values $\bar{a}$, namely

$$
\bar{a} \in\{\overline{1}, \overline{-1}, \overline{2 n+1}, \overline{2 n-1}\} .
$$

Reasoning as in the previous case (note that now $n / 2$ is odd) one concludes

$$
\left[\mathbb{Q}\left(t_{n}\right): \mathbb{Q}\right]=\phi(4 n) / 4=\phi(n)=\phi(n / 2) .
$$

(iii). Finally, if $2 \nmid n$ then only $a \equiv \pm 1 \bmod 4 n$ solve the desired equations. So one obtains the subgroup $\{\overline{1}, \overline{-1}\} \subset(\mathbb{Z} / 4 n \mathbb{Z})^{\times}$. As a consequence one concludes for $n$ odd

$$
\left[\mathbb{Q}\left(t_{n}\right): \mathbb{Q}\right]=\phi(4 n) / 2=\phi(n) .
$$

Now consider the special case where $n=p$ is an odd prime number. The calculation presented above shows that the minimal polynomial of $t_{p}=\tan (\pi / p)$ over $\mathbb{Q}$ has degree $\left[\mathbb{Q}\left(t_{p}\right): \mathbb{Q}\right]=\phi(p)=p-1$. To construct this minimal polynomial, put $\alpha:=\pi / p$. One has $\left(e^{\alpha i}\right)^{p}=-1$, hence $(\cos (\alpha)+i \sin (\alpha))^{p}+1=0$. Dividing by $(\cos (\alpha))^{p}$ and taking the imaginary part of the resulting expression yields

$$
\sum_{j=0}^{(p-1) / 2}\binom{p}{2 j+1}(-1)^{j} t_{p}^{2 j+1}=0
$$

So $t_{p}$ is a zero of

$$
\sum_{j=0}^{(p-1) / 2}\binom{p}{2 j+1}(-1)^{j} x^{2 j} \in \mathbb{Z}[x] .
$$

As this has degree $p-1$ and leading coefficient $(-1)^{(p-1) / 2}$, it is up to at most a sign the minimal polynomial of $t_{p}$ over $\mathbb{Q}$.

Example: $t_{3}$ is a zero of $x^{2}-3$ and $t_{5}$ is a zero of $x^{4}-10 x^{2}+5$ and $t_{7}$ is a zero of $x^{6}-21 x^{4}+35 x^{2}-7$.

## III. 3 Quadratic reciprocity

In this section $p, q \in \mathbb{Z}_{>0}$ are prime numbers.
Quadratic reciprocity relates the problem whether $X^{2}-q$ has a zero in $\mathbb{F}_{p}$ to the 'reciprocal' problem whether $X^{2}-p$ has a zero in $\mathbb{F}_{q}$. In the proof presented here, no Galois theory will be used, but we do use roots of unity and certain so-called Gauss sums. We will briefly mention a property of such Gauss sums in terms of Galois theory.
III.3.1 Definition. Let $p$ be an odd prime and let $a \in \mathbb{Z}$.

The Legendre-symbol $\left(\frac{a}{p}\right) \in\{-1,0,1\}$ is defined by:

$$
\begin{aligned}
& \left(\frac{a}{p}\right)=0 \Leftrightarrow a \text { is divisible by } p \\
& \left(\frac{a}{p}\right)=1 \Leftrightarrow \exists x \in \mathbb{Z}: x \not \equiv 0 \bmod p \text { and } x^{2} \equiv a \bmod p \\
& \left(\frac{a}{p}\right)=-1 \Leftrightarrow \nexists x \in \mathbb{Z}: x^{2} \equiv a \bmod p
\end{aligned}
$$

So in particular: the polynomial $X^{2}-a$ has a zero in $\mathbb{F}_{p} \Longleftrightarrow\left(\frac{a}{p}\right) \in\{+1,0\}$.
III.3.2 Example. In $\mathbb{F}_{11}$ it holds that:

$$
( \pm 1)^{2}=1, \quad( \pm 2)^{2}=4, \quad( \pm 3)^{3}=9, \quad( \pm 4)^{2}=5, \quad( \pm 5)^{2}=3
$$

so by the definition of the Legendre-symbol one finds:

$$
\left(\frac{a}{11}\right)=1 \quad \text { if } a \equiv 1,3,4,5,9 \bmod 11
$$

The remaining $10 / 2=5$ elements of $\mathbb{F}_{11}^{\times}$are not squares:

$$
\left(\frac{a}{11}\right)=-1 \quad \text { for } a \equiv 2,6,7,8,10 \bmod 11
$$

III.3.3 Theorem. Let $a \in \mathbb{Z}$. Identifying $-1,0,1 \in \mathbb{Z}$ with $-\overline{1}, \overline{0}, \overline{1} \in \mathbb{F}_{p}$ it holds that

$$
\left(\frac{a}{p}\right)=\bar{a}^{\frac{p-1}{2}} \quad\left(\in \mathbb{F}_{p}\right)
$$

Moreover for all $n, m \in \mathbb{Z}$ one has

$$
\left(\frac{n m}{p}\right)=\left(\frac{n}{p}\right)\left(\frac{m}{p}\right)
$$

Proof. If $\bar{a}=\overline{0} \in \mathbb{F}_{p}$ then $\left(\frac{a}{p}\right)=0$ and $\bar{a}^{\frac{p-1}{2}}=\overline{0}$, showing the first claim in this case.
If $\bar{a} \neq \overline{0}$ then $\bar{a}^{p-1}=1$ since $\mathbb{F}_{p}^{\times}$is a group consisting of $p-1$ elements. Hence $\bar{a}^{\frac{p-1}{2}}$ is a zero of $X^{2}-1=0$ and therefore $\bar{a}^{\frac{p-1}{2}} \in\{\overline{1}, \overline{-1}\}$, as these are the only zeros of $X^{2}-1$ in the field $\mathbb{F}_{p}$. We can therefore define

$$
\epsilon: \mathbb{F}_{p}^{\times} \longrightarrow\{\overline{1}, \overline{-1}\} \quad x \mapsto x^{\frac{p-1}{2}}
$$

Note that $\epsilon$ is a homomorphism of (multiplicative) groups. Since $\mathbb{F}_{p}^{\times}$is a cyclic group is, $y \in \mathbb{F}_{p}^{\times}$exists with $y^{\frac{p-1}{2}} \neq \overline{1}$. So $\epsilon$ is surjective and $\# \operatorname{Ker}(\epsilon)=\frac{p-1}{2}$.

In case $\bar{a}=x^{2}$ for some $x \in \mathbb{F}_{p}, x \neq 0$, then $\bar{a}^{\frac{p-1}{2}}=x^{p-1}=1$, so $\bar{a} \in \operatorname{Ker}(\epsilon)$. As $\bar{a}^{2}=\bar{b}^{2} \Leftrightarrow \bar{a}= \pm \bar{b}$, one obtains in this way $\frac{p-1}{2}$ elements of $\operatorname{Ker}(\epsilon)$. Now $\# \operatorname{Ker}(\epsilon)=\frac{p-1}{2}$ and therefore

$$
\operatorname{Ker}(\epsilon)=\left\{\bar{a} \in \mathbb{F}_{p}^{\times}: \bar{a}=x^{2} \text { for some } x \in \mathbb{F}_{p}^{\times}\right\} .
$$

As a consequence $\epsilon(\bar{b})=\overline{-1}$ in case $b$ is not a square modulo $p$, so:

$$
\left(\frac{a}{p}\right)=\epsilon(\bar{a})
$$

for all $a \in \mathbb{Z}$ such that $\bar{a} \neq \overline{0} \in \mathbb{F}_{p}$. This proves the first assertion.
The second assertion is evident in case $\bar{n}=\overline{0}$ or $\bar{m}=\overline{0}$. In all other cases one uses that the Legendre-symbol agrees with the homomorphism $\epsilon$, as was observed before. This finishes the proof.
III.3.4 Example. In $\mathbb{F}_{11}$ it holds that

$$
2^{5}=32=-1, \quad 3^{5}=27 \cdot 9=5 \cdot 9=1, \quad 5^{5}=125 \cdot 25=4 \cdot 3=1,
$$

hence 2 is not a square modulo 11, while both 3 and 5 are squares modulo 11 (see also the previous Example III.3.2.
III.3.5 Remark. Given an odd prime $p$, it is not difficult to compute $(-1)^{\frac{p-1}{2}}$. Namely, write $p$ as $p=4 k+1$ or as $p=4 k+3$ for an integer $k$. Then

$$
(-1)^{\frac{p-1}{2}}=1 \quad \text { if } \quad p=4 k+1, \quad(-1)^{\frac{p-1}{2}}=-1 \quad \text { if } \quad p=4 k+3 .
$$

Hence the polynomial $X^{2}+1$ has a zero in $\mathbb{F}_{p}$, (here $p$ is an odd prime) if and only if $p \equiv 1 \bmod 4$. (Note that for $p=2$ one has $X^{2}+1=(X+1)^{2}$ in $\mathbb{F}_{2}$.)

Of course the same conclusions can be obtained by using the well-known fact that $\mathbb{F}_{p}^{\times}$is a cyclic group consisting of $p-1$ elements, so an element of order 4 (i.e., an element whose square is $-1 \neq 1$ ) exists, if and only if $4 \mid(p-1)$.
III.3.6 Definition. Given a field $K$ and a positive integer $n$, an element $\zeta \in K$ is called a primitive $n$-th root of unity if $\zeta$ has order $n$ in the multiplicative group $K^{\times}$. (In other words, $\zeta^{n}=1$ and $\zeta^{m} \neq 1$ for $1 \leq m<n$ ).

Note that using Proposition III.1.1 (1), the $n$-th cyclotomic polynomial $\Phi_{n}$ introduced in Section III.1 can be regarded as a polynomial in $K[x]$ for any field $K$. From Proposition III.1.3.(1) it then follows that any primitive $n$-th root of unity in $K$ is in fact a zero of the polynomial $\Phi_{n} \in K[x]$.

However, a zero in $K$ of $\Phi_{n}$ is not necessarily a primitive $n$-th root of unity. For example if $\operatorname{char}(K)=p>0$ then $x^{p}-1=(x-1)^{p}$ in $K[x]$, hence primitive $p$-th roots of unity do not exist in $K$.
III.3.7 Example. If $\operatorname{char}(K) \neq 2$ then -1 is the only primitive 2 -nd root of unity in $K$.

For $K=\mathbb{C}$ the primitive $p$-th roots of unity ( $p$ a prime number) are precisely the $p-1$ (pairwise distinct) complex numbers

$$
\zeta_{p}^{k} \quad(1 \leq k \leq p-1), \quad \text { with } \quad \zeta_{p}:=\cos \frac{2 \pi}{p}+i \sin \frac{2 \pi}{p} .
$$

In particular $\frac{-1 \pm i \sqrt{3}}{2}$ are the two primitive 3-rd roots of unity.
III.3.8 Example. The primitive 3 -rd roots of unity in $\mathbb{F}_{7}$ are 2 and 4, as follows from $2^{3}=8=1$ and $4^{3}=64=1$ in $\mathbb{F}_{7}$.

Note that $\mathbb{F}_{25} \cong \mathbb{F}_{5}[X] /\left(X^{2}-2\right)$, since $X^{2}-2$ has no zeros in $\mathbb{F}_{5}$. Therefore every element of $\mathbb{F}_{25}$ can be expressed uniquely as $a+b \alpha$ with $a, b \in \mathbb{F}_{5}$ and $\alpha^{2}=2$. The primitive 3 -rd roots of unity in $\mathbb{F}_{25}$ are therefore $3(-1 \pm \alpha)$, namely

$$
(3(-1+\alpha))^{3}=2\left(-1+3 \alpha-3 \alpha^{2}+\alpha^{3}\right)=2(-1+3 \alpha-1+2 \alpha)=1,
$$

and similarly for $3(-1-\alpha)$.

We will now study primitive $p$-th roots of unity, for $p$ an odd prime number, somewhat more extensively. In particular we introduce Gauss-sums and use them to obtain information concerning quadratic equations (over finite fields).
III.3.9 Definition. Let $p$ be an odd prime number and let $K$ be a field of characteristic $\operatorname{char}(K) \neq p$. Suppose $\zeta \in K$ is a primitive $p$-th root of unity.

The $\operatorname{map} \mathbb{Z} \rightarrow K^{\times}$given by $n \mapsto \zeta^{n}$ is a homomorphism of groups with kernel $p \mathbb{Z} \subset \mathbb{Z}$. Therefore one obtains a well-defined homomorphism $\mathbb{Z} / p \mathbb{Z}=\mathbb{F}_{p} \rightarrow K^{\times}$ given by

$$
x \mapsto \zeta^{x}:=\zeta^{n}, \quad \text { for } \quad x=\bar{n}=n+p \mathbb{Z} \in \mathbb{F}_{p}
$$

Similarly the Legendre symbol $\left(\frac{n}{p}\right) \in\{0,+1,-1\}$ can, as before, be regarded as an element of $K$. Since it depends only on $p$ and on $n \bmod p$, we obtain a map $\mathbb{F}_{p} \rightarrow K$ given by

$$
x \mapsto\left(\frac{x}{p}\right):=\left(\frac{n}{p}\right) \in K, \quad \text { for } \quad x=\bar{n}=n+p \mathbb{Z} \in \mathbb{F}_{p} .
$$

The Gauss-sum $\tau \in K$ is defined as

$$
\tau:=\sum_{x \in \mathbb{F}_{p}}\left(\frac{x}{p}\right) \cdot \zeta^{x} .
$$

III.3.10 Example. Take $p=3$. Considering $K=\mathbb{C}$ and $\zeta=\frac{-1+i \sqrt{3}}{2}$, one finds

$$
\tau=\left(\frac{0}{3}\right) \cdot \zeta^{0}+\left(\frac{1}{3}\right) \cdot \zeta^{1}+\left(\frac{2}{3}\right) \cdot \zeta^{2}=0+\zeta-\zeta^{2}=i \sqrt{3} .
$$

Note that with the choice $\zeta=\frac{-1-i \sqrt{3}}{2}$ one obtains $\tau=-i \sqrt{3}$, the complex conjugate of the Gauss-sum above. So the Gauss-sum depends on the choice of the primitive $p$-th root of unity $\zeta$, however for our application this is not important.

For $K=\mathbb{F}_{7}$ and 3-rd root of unity $\zeta=2 \in \mathbb{F}_{7}$ we have

$$
\tau=\left(\frac{0}{3}\right) \cdot 2^{0}+\left(\frac{1}{3}\right) \cdot 2^{1}+\left(\frac{2}{3}\right) \cdot 2^{2}=2-4=5 .
$$

Note that $\tau^{2}=5^{2}=-3$ in $\mathbb{F}_{7}$, analogous to the case $K=\mathbb{C}$.
III.3.11 Theorem. Suppose $p$ is an odd prime number and $K$ is a field with $\operatorname{char}(K) \neq p$ and $\zeta \in K$ is a primitive p-th root of unity. The associated Gauss-sum $\tau$ satisfies

$$
\tau^{2}=\left(\frac{-1}{p}\right) \cdot p
$$

Proof. One computes

$$
\begin{aligned}
\tau^{2} & =\left(\sum_{x \in \mathbb{F}_{p}}\left(\frac{x}{p}\right) \cdot \zeta^{x}\right) \cdot\left(\sum_{y \in \mathbb{F}_{p}}\left(\frac{y}{p}\right) \cdot \zeta^{y}\right) \\
& =\sum_{x, y \in \mathbb{F}_{p}}\left(\frac{x y}{p}\right) \cdot \zeta^{x+y} \\
& =\sum_{z \in \mathbb{F}_{p}}\left(\sum_{x \in \mathbb{F}_{p}}\left(\frac{x(z-x)}{p}\right)\right) \cdot \zeta^{z} \quad(\operatorname{met} z=x+y) .
\end{aligned}
$$

For $x=0$ one finds $\left(\frac{x(z-x)}{p}\right)=0$, and for $x \neq 0$ the following holds:

$$
\begin{aligned}
\left(\frac{x(z-x)}{p}\right) & =\left(\frac{-x^{2}}{p}\right) \cdot\left(\frac{1-z x^{-1}}{p}\right) \\
& =\left(\frac{-1}{p}\right)\left(\frac{x^{2}}{p}\right) \cdot\left(\frac{1-z x^{-1}}{p}\right) \\
& =\left(\frac{-1}{p}\right) \cdot\left(\frac{1-z x^{-1}}{p}\right) .
\end{aligned}
$$

Hence

$$
\tau^{2}=\left(\frac{-1}{p}\right) \cdot \sum_{z \in \mathbb{F}_{p}} c_{z} \zeta^{z} \quad \text { with } \quad c_{z}=\sum_{x \in \mathbb{F}_{p}^{\times}}\left(\frac{1-z x^{-1}}{p}\right) .
$$

The case $z=0$ leads to

$$
c_{0}=\sum_{x \in \mathbb{F}_{p}^{\times}}\left(\frac{1}{p}\right)=\sum_{x \in \mathbb{F}_{p}^{\times}} 1=p-1 .
$$

In case $z \neq 0$, and $x$ running over the elements of $\mathbb{F}_{p}^{\times}$, one has that $z x^{-1}$ runs over $\mathbb{F}_{p}^{\times}$as well, and $w:=1-z x^{-1}$ runs over $\mathbb{F}_{p}-\{1\}$. This implies

$$
\left.c_{z}=\left(\sum_{w \in \mathbb{F}_{p}}\left(\frac{w}{p}\right)\right)-\left(\frac{1}{p}\right)=0-1=-1 \quad \text { (with } \quad z \in \mathbb{F}_{p}^{\times}\right),
$$

where it is used that $\mathbb{F}_{p}$ contains as many elements $w$ such that $\left(\frac{w}{p}\right)=1$ as it contains elements with $\left(\frac{w}{p}\right)=-1$. Substituting the values of $c_{z}$ in the formula for $\tau^{2}$ now yields

$$
\begin{aligned}
\tau^{2} & =\left(\frac{-1}{p}\right) \cdot\left(p-1+\sum_{z \in \mathbb{F}_{p}^{\times}}(-1) \zeta^{z}\right) \\
& =\left(\frac{-1}{p}\right) \cdot\left(p-1-\sum_{z \in \mathbb{F}_{p}^{\times}} \zeta^{z}\right) \\
& =\left(\frac{-1}{p}\right) \cdot\left(p-\sum_{z=0}^{p-1} \zeta^{z}\right) \\
& =\left(\frac{-1}{p}\right) \cdot p,
\end{aligned}
$$

where we also used that

$$
\sum_{z=0}^{p-1} \zeta^{z}=\frac{\zeta^{p}-1}{\zeta-1}=0
$$

This proves Theorem III.3.11
III.3.12 Remark. We "explain" in terms of Galois theory why in some sense the formula defining a Gauss-sum is very natural. So let $p$ be an odd prime, and $\zeta=\zeta_{p}=e^{2 \pi i / p} \in \mathbb{C}$. The field $K:=\mathbb{Q}(\zeta)$ is a Galois extension of $\mathbb{Q}$ with Galois group $G \cong(\mathbb{Z} / p \mathbb{Z})^{\times}$(see Corollary III.1.2). The element $a \bmod n \in(\mathbb{Z} / p \mathbb{Z})^{\times}$corresponds to the automorphism of $K$ given by $\zeta \mapsto \zeta^{a}$.

The subgroup $H \subset G$ consisting of all squares has index 2 in $G$. Hence the field $M:=K^{H}$ which by the Galois correspondence is associated to the subgroup $H$, is a quadratic extension of $\mathbb{Q}$. We describe a generator of $M$. This means: we want an element $\alpha \in K$ with the properties $\sigma(\alpha)=\alpha$ for all $\sigma \in H$ (this guarantees that $\alpha \in M)$ and moreover $\tau(\alpha) \neq \alpha$ for some $\tau \in G \backslash H$ (this guarantees $\alpha \notin \mathbb{Q}$ ).

An obvious first choice is

$$
\alpha:=\sum_{x \in \mathbb{F}_{p}^{\times} \text {a square }} \zeta^{x} .
$$

Namely, if $\sigma \in H$, then $\sigma(\zeta)=\zeta^{y}$ for some square $y \in \mathbb{F}_{p}^{\times}$, hence

$$
\sigma(\alpha)=\sum_{x \in \mathbb{F}_{p}^{\times} \text {a square }} \zeta^{x y}=\alpha
$$

(since when $x$ runs over the squares in $\mathbb{F}_{p}^{\times}$, so does $x y$ ). Moreover when $v \in G \backslash H$ then $v(\zeta)=\zeta^{z}$ for some $z \in \mathbb{F}_{p}^{\times}$which is not a square. Then

$$
v(\alpha)=\sum_{x \in \mathbb{F}_{p}^{\times} \text {a square }} \zeta^{x z}=\sum_{w \in \mathbb{F}_{p}^{\times} \text {not a square }} \zeta^{w} .
$$

Now observe that

$$
\sum_{x \in \mathbb{F}_{p}^{\times}} \zeta^{x}+\sum_{w \in \mathbb{F}_{p}^{\times}} \sum_{\text {not a square }} \zeta^{w}=\sum_{j=1}^{p-1} \zeta^{j}=-1 .
$$

As a consequence $v(\alpha)=-1-\alpha$. To show that $\alpha \notin \mathbb{Q}$ we therefore need to establish that $\alpha \neq-1-\alpha$, in other words, that $2 \alpha \neq-1$.

A short argument for this (different and slightly more elementary reasonings are given later as well as in the exercises) is the following. Evaluating at $\zeta$ yields a surjective ring homomorphism $\mathbb{Z}[x] \rightarrow \mathbb{Z}[\zeta]$ with kernel $\mathbb{Z}[x] \cdot \Phi_{p}$. Hence we have $\mathbb{Z}[\zeta] \cong \mathbb{Z}[x] /\left(\Phi_{p}\right)$. Now let $\mathbb{F}_{q}$ be the splitting field of $X^{p}-1$ over $\mathbb{F}_{2}$, and take a primitive $p$-th root of unity $\bar{\zeta} \in \mathbb{F}_{q}$. Then via $\mathbb{Z}[x] \rightarrow \mathbb{F}_{q}, f(x) \mapsto f(\bar{\zeta})$ one obtains a ring homomorphism $\mathbb{Z}[\zeta] \rightarrow \mathbb{F}_{q}$. It sends $\alpha \in \mathbb{Z}[\zeta]$ to some element $\bar{\alpha} \in \mathbb{F}_{q}$ and then $2 \alpha$ to $2 \bar{\alpha}=0$, since $\operatorname{char}\left(\mathbb{F}_{q}\right)=2$. So $2 \alpha \neq-1$ because the image of -1 is not $0 \in \mathbb{F}_{q}$.

We now have that $M=K^{H}=\mathbb{Q}[\alpha]$, a quadratic extension of $\mathbb{Q}$ with basis (as $\mathbb{Q}$-vectorspace) $1, \alpha$ and with Galois group over $\mathbb{Q}$ generated by the restriction to $M$ of any $v \in G$ as above. On the basis $1, \alpha$ such $v$ is given by the matrix $\left(\begin{array}{ll}1 & -1 \\ 0 & -1\end{array}\right)$. So a basis for $M$ over $\mathbb{Q}$ consisting of eigenvectors of $v$ is $1,2 \alpha+1$. We have

$$
2 \alpha+1=2\left(\sum_{x \in \mathbb{F}_{p}^{\times} \text {a square }} \zeta^{x}\right)-\left(\sum_{j=1}^{p-1} \zeta^{j}\right)=\tau,
$$

precisely our Gauss-sum! It is an eigenvector of $v$ with eigenvalue -1 , and hence $v\left(\tau^{2}\right)=(v(\tau))^{2}=(-\tau)^{2}=\tau^{2}$, showing that $\tau^{2} \in \mathbb{Q}$. Theorem III.3.11 shows a little more, namely it determines $\tau^{2}$ as an element of $\mathbb{Q}$, namely $(2 \alpha+1)^{2}=\tau^{2}= \pm p$, providing another argument why $2 \alpha \neq-1$.

Theorem III.3.11 describes a zero $\tau \in K$, in case $\operatorname{char}(K) \neq p$ and $K$ contains a primitive $p$-th root of unity, of the polynomial $X^{2}-\left(\frac{-1}{p}\right) p$. For example, we can take as $K$ the field $\Omega_{\mathbb{F}_{q}}^{X^{p}-1}$, for $p \neq q$ two distinct primes (and $p \neq 2$ ). In the latter case a natural question is, whether $\tau$ is contained in the field $\mathbb{F}_{q}$. In other words, whether or not $\tau^{q}=\tau$.
III.3.13 Lemma. Suppose $p$ and $q$ are distinct odd primes. Let $K$ be a field with $\operatorname{char}(K)=q$, containing a primitive $p$-th root of unity. Then

$$
\tau^{q}=\left(\frac{q}{p}\right) \cdot \tau, \quad \text { and in particular } \quad \tau \in \mathbb{F}_{q} \Longleftrightarrow\left(\frac{q}{p}\right)=1
$$

Proof. By the definition of Gauss-sum and the fact that $x \mapsto x^{q}$ is a field homomorphism of $K$, it follows that

$$
\begin{aligned}
\tau^{q} & =\left(\sum_{x \in \mathbb{F}_{p}}\left(\frac{x}{p}\right) \cdot \zeta^{x}\right)^{q} \\
& =\sum_{x \in \mathbb{F}_{p}}\left(\frac{x}{p}\right) \cdot \zeta^{q x} \\
& =\sum_{y \in \mathbb{F}_{p}}\left(\frac{y q^{-1}}{p}\right) \cdot \zeta^{y} \quad\left(\text { with } \quad y=x \bar{q} \in \mathbb{F}_{p}\right) \\
& =\left(\frac{q^{-1}}{p}\right) \cdot \sum_{y \in \mathbb{F}_{p}}\left(\frac{y}{p}\right) \cdot \zeta^{y} \\
& =\left(\frac{q}{p}\right) \cdot \tau,
\end{aligned}
$$

here we used $\bar{q}=\bar{q}^{2} \cdot \bar{q}^{-1}$, hence $\bar{q}^{-1}$ is a square in $\mathbb{F}_{p}^{\times}$precisely when $\bar{q}$ is.
For the last assertion, note that $\tau \neq 0$ by Theorem III.3.11. This proves the lemma.
III.3.14 Theorem. (quadratic reciprocity) (Gauss, 1801).

Suppose $p$ and $q$ are distinct odd primes. Then

$$
\begin{aligned}
\left(\frac{q}{p}\right)=\left(\frac{p}{q}\right) & \text { if } p \equiv 1 \bmod 4 \text { or } q \equiv 1 \bmod 4, \\
\left(\frac{q}{p}\right)=-\left(\frac{p}{q}\right) & \text { if } p \equiv q \equiv 3 \bmod 4 .
\end{aligned}
$$

Formulated somewhat shorter:

$$
\left(\frac{p}{q}\right)\left(\frac{q}{p}\right)=(-1)^{\frac{p-1}{2} \cdot \frac{q-1}{2}}
$$

Moreover

$$
\left(\frac{-1}{p}\right)=(-1)^{\frac{p-1}{2}}, \quad \text { and } \quad\left(\frac{2}{p}\right)=(-1)^{\frac{p^{2}-1}{8}}
$$

Proof. First, note that $\left(\frac{-1}{p}\right)=(-1)^{\frac{p-1}{2}}$ is immediate from Theorem III.3.3.
We now show quadratic reciprocity. Take $K=\Omega_{\mathbb{F}_{q}}^{X^{p}-1}$ and let $\tau$ be the Gauss-sum as before. By Lemma III.3.13 we have

$$
\tau^{q-1}=\left(\frac{q}{p}\right)
$$

On the other hand from Theorem III.3.11 we obtain

$$
\tau^{q-1}=\left(\tau^{2}\right)^{\frac{q-1}{2}}=\left(\frac{-1}{p}\right)^{\frac{q-1}{2}} p^{\frac{q-1}{2}}
$$

Now Theorem III.3.3implies

$$
\left(\frac{-1}{p}\right)=(-1)^{\frac{p-1}{2}}, \quad p^{\frac{q-1}{2}}=\left(\frac{p}{q}\right) \quad\left(\in \mathbb{F}_{q}\right) .
$$

Combining the above equalities one concludes

$$
\left(\frac{q}{p}\right)=(-1)^{\frac{p-1}{2} \cdot \frac{q-1}{2}}\left(\frac{p}{q}\right),
$$

which is evidently equivalent to quadratic reciprocity.
Finally we will show for any odd prime $p$ that

$$
\left(\exists x \in \mathbb{F}_{p}: \quad x^{2}=2\right) \Longleftrightarrow p \equiv \pm 1 \bmod 8
$$

Given an odd prime $p$, the polynomial $X^{4}+1$ has a zero $\zeta$ in some finite extension $K$ of $\mathbb{F}_{p}$. So $\zeta^{4}=-1$ and $\zeta^{8}=1$. Put $x:=\zeta+\zeta^{7} \in K$, then

$$
x^{2}=\left(\zeta+\zeta^{7}\right)^{2}=\zeta^{2}+\zeta^{6}+2=2
$$

where we used $\zeta^{6}=\zeta^{4} \zeta^{2}=-\zeta^{2}$.
In case $p=8 k+1$ or $p=8 k+7$ then

$$
\begin{array}{ll}
\zeta^{p}=\zeta, & \zeta^{7 p}=\zeta^{7}
\end{array} \text { if } p \equiv 1 \bmod 81 子 \zeta^{p}=\left(\zeta+\zeta^{7}\right)^{p}=\zeta+\zeta^{7}=x,
$$

where we used that $x \mapsto x^{p}$ is a field homomorphism. From $x^{p}=x$ it follows that $x \in \mathbb{F}_{p}$. This shows that in case $p \equiv \pm 1 \bmod 8$ an $x \in \mathbb{F}_{p}$ exists with $x^{2}=2$, which means $\left(\frac{2}{p}\right)=1$.

In case $p=8 k+3$ or $p=8 k+5$ then (using $\zeta^{4}=-1$ ) one finds

$$
\left.\begin{array}{ll}
\zeta^{p}=\zeta^{3}=-\zeta^{7} & \text { if } p \equiv 3 \bmod 8 \\
\zeta^{p}=\zeta^{5}=-\zeta & \text { if } p \equiv 5 \bmod 8
\end{array}\right\} \Longrightarrow x^{p}=\left(\zeta+\zeta^{7}\right)^{p}=-\left(\zeta+\zeta^{7}\right)=-x
$$

Hence the zero $x$ of $X^{2}-2$ is not in $\mathbb{F}_{p}$. The other zero (which is $-x$ ) is then not in $\mathbb{F}_{p}$ either. Hence $\left(\frac{2}{p}\right)=-1$.

This completes the proof of Theorem III.3.14.
III.3.15 Example. We present a few simple applications of Theorem III.3.14

$$
\left(\frac{7}{11}\right)=-\left(\frac{11}{7}\right)=-\left(\frac{4}{7}\right)=-1,
$$

since $4=2^{2}$. So 7 is not a square modulo 11 . Next,

$$
\left(\frac{5}{11}\right)=\left(\frac{11}{5}\right)=\left(\frac{1}{5}\right)=1
$$

so 5 is a square modulo 11 . Indeed, $4^{2}=16=5$ in $\mathbb{F}_{11}$. Using that the Legendresymbol is multiplicative (see III.3.3), we find:

$$
\left(\frac{6}{11}\right)=\left(\frac{2}{11}\right)\left(\frac{3}{11}\right)=(-1) \cdot-\left(\frac{11}{3}\right)=(-1) \cdot-\left(\frac{2}{3}\right)=-1 .
$$

## III. 4 Exercises

1. Give a proof of Corollary III.1.2
2. Prove parts (3) and (4) of Proposition III.1.3
3. Let $\lambda \in \mathbb{Q}$. Prove that $\cos (2 \pi \lambda)$ is algebraic over $\mathbb{Q}$. Prove that $\mathbb{Q}(\cos (2 \pi \lambda))$ is a Galois extension. Determine its Galois group.
4. Determine $\Phi_{72}$ in a handy way.
5. Show for any odd prime $p$ that the minimal polynomial of $\tan (\pi / p)$ over $\mathbb{Q}$ is an Eisenstein polynomial.
6. Find all integers $n \geq 3$ such that $\tan (\pi / n) \in \mathbb{Q}$.
7. For each of the integers $n \geq 3$ such that $\tan (\pi / n)$ has degree 2 over $\mathbb{Q}$, find the corresponding minimal polynomial.
8. Consider $f:=x^{6}-21 x^{4}+35 x^{2}-7 \in \mathbb{Z}[x]$. As we saw in Section III.2 this is the minimal polynomial of $t_{7}=\tan (\pi / 7)$ over $\mathbb{Q}$, and $\mathbb{Q}\left(t_{7}\right) / \mathbb{Q}$ is a Galois extension.
(a) Describe all zeros of $f$ as values of the tangent function.
(b) Writing $f$ as a product of factors of degree 1 and comparing coefficients, what identities between tangent values do you obtain?
(c) Show using Galois theory that $\mathbb{Q}\left(t_{7}\right)$ contains a unique quadratic field $\mathbb{Q}(\sqrt{d})$. (in fact, it holds that $\sqrt{7} \in \mathbb{Q}\left(t_{7}\right)$. Factor $f$ over $\mathbb{Q}(\sqrt{7})$ and deduce more tangent indentities from this!)
9. Show for primes $p>3$ that: 3 is a square modulo $p \quad \Leftrightarrow \quad p \equiv 1,-1 \bmod 12$.
10. Show for $p>3$ prime that -3 is a square in $\mathbb{F}_{p}$ if and only if $p \equiv 1 \bmod 3$.
11. Let $p>2$ be prime and $n \in \mathbb{Z}$. Prove that if $p \mid\left(n^{2}-2\right)$, then $p \equiv \pm 1 \bmod 8$.
12. Compute the Legendre-symbols $\left(\frac{5}{101}\right),\left(\frac{6}{101}\right),\left(\frac{7}{101}\right),\left(\frac{11}{101}\right)$.
13. In Remark III.3.12, two arguments are provided showing that (with $p$ an odd prime and $\zeta$ a primitive $p$-th root of unity in $\mathbb{C}$ ) one has $\alpha:=\sum \bar{a} \in \mathbb{F}_{p}^{\times 2} \zeta^{a} \neq-1 / 2$. In this exercise we give two more proofs of this fact.
(a) Use that $[\mathbb{Q}(\zeta): \mathbb{Q}]=p-1$ and show that every subset of $\left\{1, \zeta, \zeta^{2}, \zeta^{3}, \ldots, \zeta^{p-1}\right\}$ consisting of $p-1$ elements, is a basis of $\mathbb{Q}(\zeta)$ as a vectorspace over $\mathbb{Q}$. Conclude from this that $\alpha \notin \mathbb{Q}$.
(b) Another proof: forgetting the multiplication in the ring $\mathbb{Z}[\zeta] \cong \mathbb{Z}[X] /\left(\Phi_{p}\right)$, one obtains the finitely generated abelian group ( $\mathbb{Z}[\zeta],+, 0)$. Use this to show that $(\mathbb{Z}[\alpha],+, 0)$ is finitely generated as well, and conclude from this that $\alpha \notin \mathbb{Q} \backslash \mathbb{Z}$.
14. Let $p$ be a prime number and let $\zeta \in \mathbb{C}^{\times}$be a primitive $p$-th root of unity. Put $K=\mathbb{Q}(\zeta)$. Show that a unique field $M$ exists with $\mathbb{Q} \subset M \subset K$ and $[M: \mathbb{Q}]=2$.
Present an example of a cyclotomic field $\mathbb{Q}\left(\zeta_{n}\right)$ containing more than one field $M$ with $[M: \mathbb{Q}]=2$.

## IV. 1 Definition and results

Let $R$ be a commutative ring with 1 , and $n$ an integer $\geq 1$.
IV.1.1 Definition. A polynomial $f \in R\left[X_{1}, X_{2}, \ldots, X_{n}\right]$ is called symmetric if $f$ is fixed under all permutations of $X_{1}, X_{2}, \ldots, X_{n}$.
IV.1.2 Examples. The polynomials

$$
\sum_{i=1}^{n} X_{i}, \quad \prod_{i=1}^{n} X_{i}, \quad \sum_{i=1}^{n} X_{i}^{k} \quad\left(\text { with } k \in \mathbb{Z}_{\geq 0}\right)
$$

are symmetric in $R\left[X_{1}, \ldots, X_{n}\right]$.
The polynomial $X_{1} X_{2}+X_{2} X_{3}+X_{3} X_{4}+X_{4} X_{1}$ is not symmetric: it is not fixed under interchanging $X_{1}$ and $X_{2}$ (here $n=4$ ).

If $Z$ is a new variable, then the polynomial

$$
\left(Z-X_{1}\right)\left(Z-X_{2}\right) \cdots\left(Z-X_{n}\right) \quad \in R\left[X_{1}, X_{2}, \ldots, X_{n}\right][Z]
$$

can be written as

$$
Z^{n}-\sigma_{1} Z^{n-1}+\sigma_{2} Z^{n-2}-\ldots+(-1)^{n-1} \sigma_{n-1} Z+(-1)^{n} \sigma_{n}
$$

with

$$
\begin{array}{ll}
\sigma_{1} & =X_{1}+X_{2}+\ldots+X_{n} \\
\sigma_{2} & =X_{1} X_{2}+X_{1} X_{3}+\ldots+X_{1} X_{n}+X_{2} X_{3}+\ldots+X_{n-1} X_{n} \\
\sigma_{3} & =X_{1} X_{2} X_{3}+\ldots=\sum_{1 \leq i<j<k \leq n} X_{i} X_{j} X_{k} \\
\vdots \\
\sigma_{t} & =\sum_{1 \leq i_{1}<i_{2}<\ldots<i_{t} \leq n} X_{i_{1}} X_{i_{2}} \cdots X_{i_{t}} \\
\vdots & \\
\sigma_{n} & =X_{1} X_{2} \cdots X_{n} .
\end{array}
$$

The coefficients $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}$ here are symmetric polynomials called the elementary symmetric polynomials. From $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}$ one obtains other symmetric polynomials by addition, multiplication, and multiplication by elements of $R$.
IV.1.3 Example. Take $n=2$, then

$$
\sigma_{1}=X_{1}+X_{2}, \quad \sigma_{2}=X_{1} X_{2}
$$

Other symmetric polynomials are

$$
\sigma_{1}^{2}=X_{1}^{2}+2 X_{1} X_{2}+X_{2}^{2}, \quad \sigma_{1}^{2}-2 \sigma_{2}=X_{1}^{2}+X_{2}^{2}, \quad \sigma_{1}^{3}-3 \sigma_{1} \sigma_{2}=X_{1}^{3}+X_{2}^{3}, \quad \text { etcetera. }
$$

One observes that every polynomial in $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}$ with coefficients from $R$, so every $g\left(\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}\right)$ for $g \in R\left[X_{1}, \ldots, X_{n}\right]$, is a symmetric polynomial. The converse of this holds as well:
IV.1.4 Theorem. (Main theorem of symmetric polynomials) Any symmetric polynomial $f \in R\left[X_{1}, X_{2}, \ldots, X_{n}\right]$ can be written as a polynomial in $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}$ with coefficients in $R$ (so $f\left(X_{1}, \ldots, X_{n}\right)=g\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ for some $g \in R\left[X_{1}, \ldots, X_{n}\right]$ ).

Moreover this way of writing $f$ is unique (so $g$ is uniquely determined by $f$ ).

Proof. Let $f \neq 0$ be symmetric. Order the nonzero terms $r X_{1}^{a_{1}} X_{2}^{a_{2}} \cdots X_{n}^{a_{n}}$ appearing in $f$ in such a way that a term $r \cdot X_{1}^{a_{1}} X_{2}^{a_{2}} \cdots X_{n}^{a_{n}}$ appears 'before' $r^{\prime} \cdot X_{1}^{b_{1}} X_{2}^{b_{2}} \cdots X_{n}^{b_{n}}$ if $a_{i}>b_{i}$ for the least $i$ with $a_{i} \neq b_{i}$ ('lexicographical ordering').

The 'leading term'

$$
r X_{1}^{c_{1}} X_{2}^{c_{2}} \cdots X_{n}^{c_{n}} \quad(r \in R, r \neq 0)
$$

of $f$ then has

$$
\begin{aligned}
& c_{1}=\left(\text { largest } a_{1} \text { appearing in } f \text { as exponent of } X_{1}\right), \\
& c_{2}=\left(\text { largest } a_{2} \text { apprearing with a given } a_{1}=c_{1}\right),
\end{aligned}
$$

et cetera. We call $r$ the leading coefficient of $f$.
As $f$ is symmetric, $c_{1} \geq c_{2} \geq \ldots \geq c_{n}$ since otherwise interchanging two $X_{i}$ 's can provide an 'earlier' term of $f$.

We claim that the symmetrical polynomial

$$
r \sigma_{1}^{c_{1}-c_{2}} \sigma_{2}^{c_{2}-c_{3}} \cdots \sigma_{n-1}^{c_{n-1}-c_{n}} \sigma_{n}^{c_{n}}
$$

also has leading term $r X_{1}^{c_{1}} X_{2}^{c_{2}} \cdots X_{n}^{c_{n}}$. Indeed,

$$
\begin{array}{ll}
\sigma_{1} \text { has leading term } & X_{1}, \\
\sigma_{2} \text { has leading term } & X_{1} X_{2} \\
\vdots & \\
\sigma_{n} \text { has leading term } & X_{1} X_{2} \cdots X_{n} .
\end{array}
$$

Now the rule

$$
\text { leading term }(g) \cdot \text { leading term }(h)=\text { leading term }(g \cdot h)
$$

(which holds for polynomials $g, h$ with leading coefficient 1 ), shows

$$
\begin{aligned}
\text { leading term }\left(\sigma_{1}^{c_{1}-c_{2}} \sigma_{2}^{c_{2}-c_{3}} \cdots \sigma_{n}^{c_{n}}\right) & =X_{1}^{c_{1}-c_{2}} \cdot\left(X_{1} X_{2}\right)^{c_{2}-c_{3}} \cdots\left(X_{1} X_{2} \cdots X_{n}\right)^{c_{n}} \\
& =X_{1}^{c_{1}} X_{2}^{c_{2}} \cdots X_{n}^{c_{n}},
\end{aligned}
$$

as claimed.
So one concludes that

$$
f_{1}:=f-r \sigma_{1}^{c_{1}-c_{2}} \cdots \sigma_{n}^{c_{n}}
$$

only contains terms which lexicographically come later than the leading term of $f$.
If $f_{1}=0$ then we found an expression for $f$ as desired:

$$
f=r \sigma_{1}^{c_{1}-c_{2}} \cdots \sigma_{n}^{c_{n}}
$$

If $f_{1} \neq 0$, note that $f_{1}$ is symmetric as well, so we can apply the same reasoning to $f_{1}$ as was done for $f$. This gives

$$
f_{2}=f_{1}-r^{\prime} \sigma_{1}^{c_{1}^{\prime}-c_{2}^{\prime}} \cdots \sigma_{n}^{c_{n}^{\prime}}
$$

and all terms appearing in $f_{2}$ come lexicographically later than the leading term $r^{\prime} X_{1}^{c_{1}^{\prime}} \cdots X_{n}^{c_{n}^{\prime}}$ of $f_{1}$. In case $f_{2}=0$ we are done:

$$
f=r \sigma_{1}^{c_{1}-c_{2}} \cdots \sigma_{n}^{c_{n}}+r^{\prime} \sigma_{1}^{c_{1}^{\prime}-c_{2}^{\prime}} \cdots \sigma_{n}^{c_{n}^{\prime}}
$$

otherwise we continue with $f_{2}$.
It remains to show that this process terminates, so in the sequence $f_{1}, f_{2}, f_{3}, \ldots$ one has $f_{k}=0$ for some $k$.

To this end one introduces the total degree $\operatorname{tdeg}(f)$ of $f$, which is the maximal $a_{1}+a_{2}+\cdots+a_{n}$ one obtains from the terms $r \cdot X_{1}^{a_{1}} \cdots X_{n}^{a_{n}}(\neq 0)$ of $f$. Note that $\operatorname{tdeg}\left(\sigma_{i}\right)=i$, and therefore

$$
\begin{aligned}
\operatorname{tdeg}\left(\sigma_{1}^{c_{1}-c_{2}} \cdots \sigma_{n}^{c_{n}}\right) & =1 \cdot\left(c_{1}-c_{2}\right)+2 \cdot\left(c_{2}-c_{3}\right)+\ldots+n \cdot c_{n} \\
& =c_{1}+c_{2}+\ldots+c_{n} \\
& \leq \operatorname{tdeg}(f) .
\end{aligned}
$$

It follows that

$$
\operatorname{tdeg}\left(f_{1}\right) \leq \operatorname{tdeg}(f)
$$

and more generally

$$
\ldots \leq \operatorname{tdeg}\left(f_{m}\right) \leq \operatorname{tdeg}\left(f_{m-1}\right) \leq \ldots \leq \operatorname{tdeg}(f)
$$

However, for a given total degree only finitely many products $X_{1}^{a_{1}} \cdots X_{n}^{a_{n}}$ are possible. In every step of the process at least one such product is erased and all remaining ones come lexicographically later. This shows that after finitely many steps we have $f_{k}=0$, completing the proof of the first assertion in Theorem IV.1.4

It remains to show: if $g_{1} \neq g_{2}$ are polynomals in $n$ variables over $R$, then $g_{1}\left(\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}\right) \neq g_{2}\left(\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}\right)$. Writing $g=g_{1}-g_{2}$ one concludes that it suffices to show:

$$
\text { if } g \in R\left[Y_{1}, \ldots, Y_{n}\right], g \neq 0, \quad \text { then } g\left(\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}\right) \neq 0
$$

Every term appearing in $g$ can be written as

$$
r Y_{1}^{a_{1}-a_{2}} \cdot Y_{2}^{a_{2}-a_{3}} \cdots Y_{n}^{a_{n}},
$$

with $r \in R, r \neq 0, a_{i} \in \mathbb{Z}_{\geq 0}$. Consider the term in $g$ such that the corresponding $a_{1}, a_{2}, \ldots, a_{n}$ (i.e., the product $X_{1}^{a_{1}} X_{2}^{a_{2}} \cdots X_{n}^{a_{n}}$ ) comes first in the lexicographical ordering. Substituting $\sigma_{i}$ for $Y_{i}$ one obtains a polynomial in $X_{1}, \ldots, X_{n}$ with leading term

$$
\text { (*) } \quad r X_{1}^{a_{1}} X_{2}^{a_{2}} \cdots X_{n}^{a_{n}}
$$

since the other terms $r^{\prime} \sigma^{a_{1}^{\prime}-a_{2}^{\prime}} \ldots \sigma_{n}^{a_{n}^{\prime}}$ give rise to polynomials in $X_{1}, \ldots, X_{n}$ with a later leading term. Hence ( $*$ ) does not cancel against other terms, and therefore $g\left(\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}\right) \neq 0$. This proves the theorem.
IV.1.5 Remark. In the case $R=K$ is a field, we discuss some relation between Theorem IV.1.4 and Galois theory. Namely, put $L:=K\left(X_{1}, \ldots, X_{n}\right)$, the field of rational functions over $K$ in the variables $X_{1}, \ldots X_{n}$. Then the elementary symmetric polynomials $\sigma_{j}$ are elements of $L$, so we obtain a subfield $M:=K\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ of $L$. Any permutation of the variables $X_{1}, \ldots, X_{n}$ extends to a $K$-linear automorphism of $L$ and these permutations in fact restrict to the identity map on the subfield $M$. Moreover $L$ is the splitting field over $M$ of the separable polynomial $\prod_{j=1}^{n}\left(X-X_{j}\right) \in M[X]$. As this polynomial has degree $n$, we have $[L: M] \leq n!$.

The Galois group $\operatorname{Gal}(L / M)$ contains the subgroup corresponding to all permutations of $X_{1}, \ldots, X_{n}$, so in particular $[L: M]=\# \operatorname{Gal}(L / M) \geq n!$.

The conclusion is that $[L: M]=n!$ and $G:=\operatorname{Gal}(L / M) \cong S_{n}$. In particular every permutation of the variables $X_{j}$ extends to an element of $G$, and all elements of $G$ are obtained in this way. Galois theory then tells us that $L^{G}=M$, so the rational functions in the $X_{j}$ 's that are symmetric, which means invariant under any permutation of the variables, are precisely the rational functions in the $\sigma_{j}$ 's.

TheoremIV.1.4, which we proved without any reference to Galois theory, states that this even holds when we replace 'rational functions' by 'polynomials', and even stronger, even when we replace the field $K$ by an arbitrary ring $R$.
IV.1.6 Example. Take $n=3$, and

$$
f=X_{1}^{3} X_{2}+X_{1}^{3} X_{3}+X_{1} X_{2}^{3}+X_{1} X_{3}^{3}+X_{2}^{3} X_{3}+X_{2} X_{3}^{3}
$$

The terms here are already lexicographically ordered, and the leading term $X_{1}^{3} X_{2}$ has $c_{1}=3,, c_{2}=1$, and $c_{3}=0$. Following the proof of Theorem IV.1.4 we must subtract from $f$ the symmetric expression

$$
\begin{gathered}
\sigma_{1}^{c_{1}-c_{2}} \sigma_{2}^{c_{2}-c_{3}} \sigma_{3}^{c_{3}}=\sigma_{1}^{2} \sigma_{2}=\left(X_{1}+X_{2}+X_{3}\right)^{2} \cdot\left(X_{1} X_{2}+X_{1} X_{3}+X_{2} X_{3}\right) \\
\quad=X_{1}^{3} X_{2}+X_{1}^{3} X_{3}+2 X_{1}^{2} X_{2}^{2}+5 X_{1}^{2} X_{2} X_{3}+2 X_{1}^{2} X_{3}^{2}+X_{1} X_{2}^{3} \\
+5 X_{1} X_{2}^{2} X_{3}+5 X_{1} X_{2} X_{3}^{2}+X_{1} X_{3}^{3}+X_{2}^{3} X_{3}+2 X_{2}^{2} X_{3}^{2}+X_{2} X_{3}^{3}
\end{gathered}
$$

The result is

$$
f_{1}=-2 X_{1}^{2} X_{2}^{2}-5 X_{1}^{2} X_{2} X_{3}-2 X_{1}^{2} X_{3}^{2}-5 X_{1} X_{2}^{2} X_{3}-5 X_{1} X_{2} X_{3}^{2}-2 X_{2}^{2} X_{3}^{2}
$$

From this one subtracts

$$
-2 \sigma_{2}^{2}=-2 X_{1}^{2} X_{2}^{2}-4 X_{1}^{2} X_{2} X_{3}-2 X_{1}^{2} X_{3}^{2}-4 X_{1} X_{2}^{2} X_{3}-4 X_{1} X_{2} X_{3}^{2}-2 X_{2}^{2} X_{3}^{2}
$$

which leads to

$$
f_{2}=f_{1}-\left(-2 \sigma_{2}^{2}\right)=-X_{1}^{2} X_{2} X_{3}-X_{1} X_{2}^{2} X_{3}-X_{1} X_{2} X_{3}^{2}
$$

Subtracting from this $-\sigma_{1} \sigma_{3}$ one obtains 0 , so

$$
f=\sigma_{1}^{2} \sigma_{2}-2 \sigma_{2}^{2}-\sigma_{1} \sigma_{3} .
$$

Hence the $g$ as mentioned in Theorem IV.1.4 is $X_{1}^{2} X_{2}-2 X_{2}^{2}-X_{1} X_{3}$.
Usually Theorem[IV.1.4 is applied to the following situation. Let $f \in R\left[X_{1}, \ldots, X_{n}\right]$ be symmetric, and take $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n} \in R$. By Theorem IV.1.4 $f$ can be expressed in $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}$, hence $f\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right)$ can be expressed in

$$
\begin{aligned}
& \sigma_{1}\left(\alpha_{1}, \ldots, \alpha_{n}\right)=\alpha_{1}+\ldots+\alpha_{n} \\
& \sigma_{2}\left(\alpha_{1}, \ldots, \alpha_{n}\right)=\alpha_{1} \alpha_{2}+\alpha_{1} \alpha_{3}+\ldots+\alpha_{n-1} \alpha_{n} \\
& \vdots \\
& \sigma_{n}\left(\alpha_{1}, \ldots, \alpha_{n}\right)=\alpha_{1} \alpha_{2} \ldots \alpha_{n}
\end{aligned}
$$

which are precisely $\pm$ the coefficients of

$$
\left(X-\alpha_{1}\right)\left(X-\alpha_{2}\right) \ldots\left(X-\alpha_{n}\right) .
$$

Roughly speaking this means that every symmetric expression in 'the $n$ zeros' of a monic polynomial of degree $n$ in one variable, can be expressed in the coefficients of this polynomial.

This assertion becomes particularly relevant of these $n$ zeros are not contained in the ring $R$ itself, but only in some extension $R^{\prime} \supset R$. We now first present an example and then continue with a general result.
IV.1.7 Example. Consider $h=X^{3}-X-1 \in \mathbb{Z}[X]$. In $\mathbb{Z}$, and even in $\mathbb{Q}$, no zero of $h$ exists (as follows using the lemma of Gauss). However, $\alpha_{1}, \alpha_{2}, \alpha_{3} \in \mathbb{C}$ exist such that

$$
X^{3}-X-1=\left(X-\alpha_{1}\right)\left(X-\alpha_{2}\right)\left(X-\alpha_{3}\right) .
$$

Comparing coefficients yields

$$
\begin{array}{llr}
\sigma_{1}\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right) & =\alpha_{1}+\alpha_{2}+\alpha_{3} & = \\
\sigma_{2}\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right) & =\alpha_{1} \alpha_{2}+\alpha_{1} \alpha_{3}+\alpha_{2} \alpha_{3} & = \\
\sigma_{3}\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right) & =\alpha_{1} \alpha_{2} \alpha_{3} & = \\
1 .
\end{array}
$$

Hence Theorem IV.1.4 implies:
is $f \in \mathbb{Z}\left[X_{1}, X_{2}, X_{3}\right]$ an arbitrary symmetric polynomial, then $f\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right) \in \mathbb{Z}$ (although $\alpha_{1}, \alpha_{2}, \alpha_{3} \notin \mathbb{Z}$ ).

Taking as $f$ the polynomial from ExampleIV.1.6, so

$$
f=X_{1}^{3} X_{2}+X_{1}^{3} X_{3}+\ldots+X_{2} X_{3}^{3}=\sigma_{1}^{2} \sigma_{2}-2 \sigma_{2}^{2}-\sigma_{1} \sigma_{3}
$$

one finds by substiting $X_{i}:=\alpha_{i}$ :

$$
\alpha_{1}^{3} \alpha_{2}+\alpha_{1}^{3} \alpha_{3}+\ldots+\alpha_{2} \alpha_{3}^{3}=0^{2} \cdot(-1)-2 \cdot(-1)^{2}-0.1=-2 .
$$

IV.1.8 Theorem. Let $R^{\prime}$ be a commutative ring (with 1) and $R \subset R$ a subring. Suppose $h \in R[X]$ has degree $n$ and suppose $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n} \in R^{\prime}$ exist such that

$$
h=\left(X-\alpha_{1}\right)\left(X-\alpha_{2}\right) \ldots\left(X-\alpha_{n}\right) .
$$

Then for every symmetric $f \in R\left[X_{1}, X_{2}, \ldots, X_{n}\right]$ it holds that

$$
f\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right) \in R .
$$

Proof. For all $i$ we have

$$
\sigma_{i}\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in R
$$

since this is up to sign the coefficient of $X^{i}$ in $h \in R[X]$. As the symmetric $f$ can be written as $f=g\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ for some $g \in R\left[X_{1}, \ldots, X_{n}\right]$ (see Theorem IV.1.4), substituting the $\alpha_{i}$ in $g\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ yields an element in $R$, as desired.
IV.1.9 Example. An important symmetric polynomial is

$$
D=\prod_{1 \leq i<j \leq n}\left(X_{i}-X_{j}\right)^{2} .
$$

The discriminant of a polynomial

$$
h=X^{n}+a_{1} X^{n-1}+\ldots+a_{n}=\left(X-\alpha_{1}\right)\left(X-\alpha_{2}\right) \ldots\left(X-\alpha_{n}\right)
$$

is defined as

$$
\Delta(h)=\prod_{1 \leq i<j \leq n}\left(\alpha_{i}-\alpha_{j}\right)^{2}=D\left(\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right) .
$$

This discriminant can be expressed in $a_{1}, a_{2}, \ldots, a_{n}$. For $n=2,3,4$ one obtains the following formulas:

$$
\begin{gathered}
\Delta\left(X^{2}+a X+b\right)=a^{2}-4 b \\
\Delta\left(X^{3}+a X^{2}+b X+c\right)=a^{2} b^{2}-4 b^{3}-4 a^{3} c-27 c^{2}+18 a b c \\
\Delta\left(X^{4}+a X^{3}+b X^{2}+c X+d\right)= \\
\frac{1}{27}\left\{4\left(b^{2}-3 a c+12 d\right)^{3}-\left(2 b^{3}-72 b d+27 a^{2} d-9 a b c+27 c^{2}\right)^{2}\right\}
\end{gathered}
$$

(when expanding the latter expression, it turns out that the 27 in the denominator cancels).

Given a (unitary) ring $R$, the discriminant of $h=X^{n}+a_{1} X^{n-1}+\ldots+a_{n} \in R[X]$ is defined as an expression in the $a_{j}$ 's, regardless whether or not $h$ can be factored in $R[X]$ as a product of $n$ factors $\left(X-\alpha_{i}\right)$.

The use of the discriminant relies on the observarion that in the case where the ring $R$ is an integral domain, we have

$$
\Delta(h)=0 \Longleftrightarrow \exists i, j \quad i \neq j: \quad \alpha_{i}=\alpha_{j},
$$

in other words, the discriminant is zero if and only if the polynomial has a multiple zero. The case $n=2$ is already discussed in high school. An other application of the discriminant is discussed in Exercise 5 on page 34

Is $K$ a field of characteristic $\operatorname{char}(K) \neq 3$ (so that $\frac{1}{3} \in K$ ), then the substitution $X:=X-\frac{1}{3} a$ brings a polynomial $f=X^{3}+a X^{2}+b X+c=\Pi\left(X-\alpha_{i}\right)$ of degree 3 into the form $g=X^{3}+p X+q$. Note that $\Delta(f)=\Delta(g)$, since the zeros of $g$ are $\beta_{i}:=\alpha_{i}+\frac{1}{3} a$ and $\alpha_{i}-\alpha_{j}=\beta_{i}-\beta_{j}$. The discriminant of $g$ is simply $\Delta(g)=-\left(4 p^{3}+27 q^{2}\right)$.
IV.1.10 Example. We will use symmetric polynomials to obtain expressions for the zeros of a polynomial of degree 3 .

Let $K$ be a field of characteristic $\operatorname{char}(K) \neq 2,3$, and let $f \in K[X]$ be a monic polynomial of degree 3 :

$$
f=X^{3}+a X^{2}+b X+c
$$

Let $\alpha_{1}, \alpha_{2}, \alpha_{3}$ be the zeros of $f$ (in a splitting field of $f$ over $K$, see Section I.3. Then

$$
\begin{aligned}
-a & =\alpha_{1}+\alpha_{2}+\alpha_{3} \\
b & =\alpha_{1} \alpha_{2}+\alpha_{1} \alpha_{3}+\alpha_{2} \alpha_{3} \\
-c & =\alpha_{1} \alpha_{2} \alpha_{3}
\end{aligned}
$$

Let $\omega$ be a primitive third root of unity (in an extension of the chosen splitting field of $f$ over $K$ ), so $\omega \neq 1, \omega^{3}=1$ ). Put

$$
\begin{aligned}
& A_{1}:=\alpha_{1}+\omega \alpha_{2}+\omega^{2} \alpha_{3}, \\
& A_{2}:=\alpha_{1}+\omega^{2} \alpha_{2}+\omega \alpha_{3}
\end{aligned}
$$

We consider what happens to $A_{i}$ when permuting the $\alpha_{i}$ 's. For $\rho:=(123) \in S_{3}$ we find

$$
\rho=(123): \begin{array}{ll} 
& A_{1} \mapsto \alpha_{2}+\omega \alpha_{3}+\omega^{2} \alpha_{1}=\omega^{2} A_{1}, \\
& A_{2} \mapsto \alpha_{2}+\omega^{2} \alpha_{3}+\omega \alpha_{1}=\omega A_{2} .
\end{array}
$$

Moreover,

$$
\tau=(23): A_{1} \mapsto A_{2}, \quad(23): A_{2} \mapsto A_{1} .
$$

Since the group $S_{3}$ is generated by $\rho$ and $\tau$ it follows that

$$
A_{1}^{3}+A_{2}^{3}, \quad A_{1} A_{2}
$$

are symmetric polynomials in $\alpha_{1}, \alpha_{2}, \alpha_{3}$.
Therefore by Theorem IV.1.4 they can be expressed in terms of the elementary symmetric polynomials $\sigma_{i}$ which are in our case, up to sign, equal to the coefficients of $f$. A short calculation reveals

$$
\begin{array}{rll}
2 B & := & A_{1}^{3}+A_{2}^{3} \\
A & =-2 a^{3}+9 a b-27 c, \\
& A_{1} A_{2} & =a^{2}-3 b .
\end{array}
$$

Since

$$
\left(T-A_{1}^{3}\right)\left(T-A_{2}^{3}\right)=T^{2}-2 B T+A^{3}
$$

we conclude that $A_{1}^{3}, A_{2}^{3}$ are given by

$$
A_{i}^{3}=B \pm \sqrt{B^{2}-A^{3}} .
$$

Hence

$$
A_{i}=\sqrt[3]{B \pm \sqrt{B^{2}-A^{3}}}
$$

(here one has three choices for the cube root ( $\sqrt[3]{ }$ '). Finally one determines $\alpha_{1}$ by noting that

$$
\begin{aligned}
3 \alpha_{1} & =\left(\alpha_{1}+\alpha_{2}+\alpha_{3}\right)+\left(\alpha_{1}+\omega \alpha_{2}+\omega^{2} \alpha_{3}\right)+\left(\alpha_{1}+\omega^{2} \alpha_{2}+\omega \alpha_{3}\right) \\
& =-a+A_{1}+A_{2},
\end{aligned}
$$

where it is used that $\omega^{2}+\omega+1=0$.
As an explicit example, take

$$
f=X^{3}+2 X^{2}-X-2 \quad \in \mathbb{Q}[X] .
$$

Here

$$
\begin{gathered}
B=10, \quad A=7 \\
A_{1}=\sqrt[3]{10+\sqrt{10^{2}-7^{3}}}=\sqrt[3]{10+9 i \sqrt{3}}
\end{gathered}
$$

(we chose $\mathrm{a}+$ sign, this turns out to be irrelevant). There are three solutions in $\mathbb{C}$ to $A_{1}^{3}=10+9 i \sqrt{3}$, namely

$$
A_{1}=-2+i \sqrt{3} \quad \text { and } \quad A_{1}=\frac{1}{2}(-1+3 i \sqrt{3}) \quad \text { and } \quad A_{1}=\frac{1}{2}(5+i \sqrt{3}) .
$$

Since $A_{1} A_{2}=A=7$ the corresponding $A_{2}$ 's are

$$
A_{2}=-2-i \sqrt{3}, \quad A_{2}=\frac{1}{2}(-1+3 i \sqrt{3}), \quad A_{2}=\frac{1}{2}(5-i \sqrt{3}) .
$$

We now obtain the three zeros $\alpha$ of $f$ using $\alpha=\frac{1}{3}\left(-a+A_{1}+A_{2}\right)$, they are

$$
-2, \quad-1, \quad 1 .
$$

The formulas for the roots of cubic equations which are obtained in this way, were in a different way found by Cardano en Tartaglia around 1540. They are called the Cardano formulas.
IV.1.11 Remark. We briefly return to the described solution method for solving cubic equations. We interpret the method in terms of Galois theory. Namely, let $K$ be a field (of characteristic $\neq 3, \neq 2$ ) and let $\sigma_{1}, \sigma_{2}, \sigma_{3}$ be variables. To be in exactly the situation discussed before, we will assume that $K$ contains a primitive third root of unity $\omega$. Over $K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ we consider the polynomial $f=X^{3}-\sigma_{1} X^{2}+\sigma_{2} X-\sigma_{3}$.

Let $L$ be a splitting field of $f$ over $K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ and let $X_{1}, X_{2}, X_{3}$ be the three zeros of $f$ in $L$. As we saw in Remark IV.1.5 $L=K\left(X_{1}, X_{2}, X_{3}\right)$ and the extension $K\left(X_{1}, X_{2}, X_{3}\right) \supset K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ is Galois, with Galois group $S_{3}$ acting by permutations on $X_{1}, X_{2}, X_{3}$.

The elements $A_{1}^{3}, A_{2}^{3} \in K\left(X_{1}, X_{2}, X_{3}\right)$ used in deriving the Cardano formulas are invariant under the subgroup $A_{3} \subset S_{3}$, while the 2-cycles in $S_{3}$ interchange $A_{1}^{3}$ and $A_{2}^{3}$. So $A_{1}^{3}, A_{2}^{3} \notin K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ but they are in the field $L^{A_{3}}$ of invariants under $A_{3}$. This is a quadratic extension of $K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$, and indeed we found a quadratic polynomial over $K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ with the $A_{j}^{3}$ 's as zeros: $T^{2}-2 B T+A^{3}$. The element $A_{1} \in L$ is not in this subfield since it is not fixed under the automorphisms in $A_{3}$. However $A_{1}$ satisfies the cubic equation $A_{i}^{3}=W$ for some given $W \in K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)\left(A_{1}^{3}\right)$.

So we found the zeros of $f$ by constructing a quadratic extension of $K\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ and then taking the cube root of a suitable element in that extension.

In Exercise 4 a similar idea is used to obtain the zeros of a polynomial $f$ of degree 4 , which we now briefly sketch. In this case, the 'general' Galois group is $S_{4}$. We first describe 3 combinations $C_{i}$ of the zeros of $f$ which are invariant under the normal subgroup

$$
H:=\{e,(12)(34),(13)(24),(14)(23)\} \subset S_{4} .
$$

Since $S_{4} / H \cong S_{3}$ the $C_{i}$ generate a Galois extension with Galois group $S_{3}$. The polynomial $\left(X-C_{1}\right)\left(X-C_{2}\right)\left(X-C_{3}\right)$ then has coefficients which are invariant under all of $S_{4}$, so they can be expressed in the coefficients of $f$. Since it turns out to be relatively easy to obtain the zeros of $f$ in terms of square roots of the $C_{i}$, this reduces the problem of finding expressions for the zeros of $f$ to the problem of finding such expressions for a cubic polynomial, and that was done before.

Unfortunately, for $n>4$ the group $A_{n}$ is the only nontrivial normal subgroup of $S_{n}$. For this reason, solving an equation of degree $n>4$ is in an essential way more difficult than the cases $n \leq 4$.

1. Express the symmetric polynomial $X_{1}^{3}+X_{2}^{3}+X_{3}^{3}$ (here $n=3$ ) in terms of $\sigma_{1}, \sigma_{2}, \sigma_{3}$.
2. In the proof of IV.1.4 we used the rule

$$
\text { leadingterm }(g) \cdot \text { leadingterm }(h)=\text { leadingterm }(g \cdot h)
$$

for polynomials $g, h$ with leading coefficient 1 . Show that this rule is false in general if we allow $g, h$ to have zero divisors as leading coefficents.
3. Let $\left(X-\alpha_{1}\right)\left(X-\alpha_{2}\right)\left(X-\alpha_{3}\right)=X^{3}-X-1$, with $\alpha_{1}, \alpha_{2}, \alpha_{3} \in \mathbb{C}$. Put $s_{k}=\alpha_{1}^{k}+\alpha_{2}^{k}+\alpha_{3}^{k}$ for $k \in \mathbb{Z}$. Show:

$$
\begin{array}{ll}
s_{-1}=-1, s_{0}=3, & s_{1}=0 \\
s_{k}=s_{k-2}+s_{k-3} & \text { for all } k \in \mathbb{Z}, \\
s_{k} \in \mathbb{Z} & \text { for all } k \in \mathbb{Z} \text { (also negative ones!). }
\end{array}
$$

4. Take $f=X^{4}+a X^{3}+b X^{2}+c X+d \in K[X]$ where $K$ is a field of characteristic $\operatorname{char}(K) \neq 2,3$. Write $\alpha_{1}, \ldots, \alpha_{4}$ for the zeros of $f$ in some extension of $K$, so $f=\Pi\left(X-\alpha_{j}\right)$.
(a) Define

$$
\begin{aligned}
& C_{1}=\left(\alpha_{1}+\alpha_{2}-\alpha_{3}-\alpha_{4}\right)^{2} \\
& C_{2}=\left(\alpha_{1}-\alpha_{2}+\alpha_{3}-\alpha_{4}\right)^{2} \\
& C_{3}=\left(\alpha_{1}-\alpha_{2}-\alpha_{3}+\alpha_{4}\right)^{2} .
\end{aligned}
$$

Express $\alpha_{1}$ in terms of $\sqrt{C_{i}}$ and the coefficient $a$ of $f$.
(b) Verify that the $S_{4}$ action (permuting the $\alpha_{i}$ 's) also permutes the $C_{i}$ 's. Check that the subgroup $H$ in Remark IV.1.11 leaves all $C_{i}$ invariant.
(c) Show that

$$
\begin{aligned}
C_{1}+C_{2}+C_{3} & =3 a^{2}-8 b \\
C_{1} C_{2}+C_{1} C_{3}+C_{2} C_{3} & =3 a^{4}-16 a^{2} b+16 b^{2}+16 a c-64 d \\
C_{1} C_{2} C_{3} & =\left(a^{3}-4 a b+8 c\right)^{2} .
\end{aligned}
$$

(d) Explain how the information above can be used to find an expression for the solutions of a general equation of degree 4 .
5. Suppose $f=X^{3}+a X^{2}+b X+c \in \mathbb{Q}[X]$ is an irreducible polynomial with zeros $\alpha_{1}, \alpha_{2}, \alpha_{3} \in \mathbb{C}$. The splitting field $\Omega_{\mathbb{Q}}^{f} \cong \mathbb{Q}\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)$, hence

$$
\sqrt{\triangle}:=\left(\alpha_{1}-\alpha_{2}\right)\left(\alpha_{1}-\alpha_{3}\right)\left(\alpha_{2}-\alpha_{3}\right) \in \Omega_{\mathbb{Q}}^{f}, \quad \text { and } \quad \Delta \in \mathbb{Q}
$$

where $\Delta$ is the discriminant of $f$ introduced in ExampleIV.1.9.
(a) Prove that $\left[\Omega_{\mathbb{Q}}^{f}: \mathbb{Q}\right]=3$ or 6 .
(b) Prove that $\sqrt{\triangle} \notin \mathbb{Q} \Rightarrow\left[\Omega_{\mathbb{Q}}^{f}: \mathbb{Q}\right]=6$.
(c) Now assume $\sqrt{\triangle} \in \mathbb{Q}$. Put $f=\left(X-\alpha_{1}\right)\left(X^{2}+r X+s\right) \in \mathbb{Q}\left(\alpha_{1}\right)[X]$. Show that $\alpha_{2} \in \mathbb{Q}\left(\alpha_{1}\right)$ by expressing $\alpha_{2}$ in terms of $\sqrt{\triangle}, a, b, c, r, s, \alpha_{1} \in \mathbb{Q}\left(\alpha_{1}\right)$. Conclude that $\left[\Omega_{\mathbb{Q}}^{f}: \mathbb{Q}\right]=3 \Leftrightarrow \sqrt{\Delta} \in \mathbb{Q}$.
(d) Describe $\operatorname{Gal}\left(\Omega_{\mathbb{Q}}^{f} / \mathbb{Q}\right)$ (the answer depends on $\Delta$ being a square or not).

## ALGEBRAICALLY CLOSED FIELDS

## V. 1 Definition and examples

V.1.1 Definition. A field $K$ is called algebraically closed if for every $f \in K[X]$ such that $f \notin K$, an $\alpha \in K$ exists such that $f(\alpha)=0$.

One of the statements in the next result is that if $K$ is algebraically closed, then every $f \in K[X], f \neq 0$, splits in $K[X]$ as a product of factors of degree 1 .
V.1.2 Theorem. Let $K$ be a field. The following statements are equivalent:
(a) $K$ is algebraically closed;
(b) every irreducible polynomial in $K[X]$ has degree 1 ;
(c) the only algebraic extension $L$ of $K$ is $L=K$;
(d) for every monic $f \in K[X]$ elements $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n} \in K$ exist with $f=\prod_{i=1}^{n}\left(X-\alpha_{i}\right)$.

Proof. (a) $\Rightarrow$ (b). If (a) holds, then a polynomial of degree $>1$ has a zero, hence is not irreducible.
(b) $\Rightarrow$ (c). Suppose $L \supset K$ is algebraic. Then any $\alpha \in L$ has a minimal polynomial $f_{K}^{\alpha}$ which is irreducible in $K[X]$, so $\operatorname{deg}\left(f_{K}^{\alpha}\right)=1$ using (b). Hence $\alpha \in K$, which implies $L=K$.
(c) $\Rightarrow$ (d). The splitting field $\Omega_{K}^{f}$ of $f$ over $K$ is algebraic over $K$, hence (c) implies $\Omega_{K}^{f}=K$, which is what we wanted to show.
(d) $\Rightarrow$ (a). This is evident, since every $\alpha_{i}$ is a zero of $f$.

This completes the proof of Theorem V.1.2.
V.1.3 Theorem. Every algebraically closed field $K$ is infinite.

Proof. If $K$ is finite, say $K=\left\{\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\}$, then

$$
f=1+\prod_{i=1}^{n}\left(X-\alpha_{i}\right)
$$

is not constant and has no zero in $K$. This proves V.1.3

The next result, originally proven in the doctoral dissertation (1799) of Carl Friedrich Gauss, is sometimes called the 'Fundamental Theorem of Algebra'.
V.1.4 Theorem. The field $\mathbb{C}$ of complex numbers is algebraically closed.

A simple analytic proof of this result, which relies on Liouville's theorem, is often discussed in introductory courses on complex function theory. Below we present a proof in which all ingredients except the following lemma are purely algebraic.
V.1.5 Lemma. Let $f \in \mathbb{R}[X]$, and assume that $\beta, \gamma \in \mathbb{R}$ exist with $f(\beta)>0$ and $f(\gamma)<0$.

Then $\alpha \in \mathbb{R}$ exists with $f(\alpha)=0$.
Proof. (of Lemma V.1.5. This is a special case of the intermediate value theorem from Real Analysis, since polynomials define continous functions.

In the proof of Theorem V.1.4 we will make use of the following lemmas.
V.1.6 Lemma. Any $f \in \mathbb{C}[X]$ of degree 2 has a zero in $\mathbb{C}$.

Proof. We may assume that $f$ is monic: $f=X^{2}+\beta X+\gamma$, with $\beta, \gamma \in \mathbb{C}$. Writing

$$
f=\left(X+\frac{1}{2} \beta\right)^{2}-\left(\frac{1}{4} \beta^{2}-\gamma\right)
$$

it suffices to show that the complex number $\frac{1}{4} \beta^{2}-\gamma$ has a square root in $\mathbb{C}$.
Write $\frac{1}{4} \beta^{2}-\gamma=a+b i$, for $a, b \in \mathbb{R}$. First consider the case $b=0$. If $a>0$ one obtains the root $\sqrt{a}$ in $\mathbb{R}$ by applying LemmaV.1.5 to $g=X^{2}-a$, observing that $g(0)<0, g(a+1)>0$. Is $a \leq 0$ then we have the square root $i \sqrt{|a|}$ in $\mathbb{C}$. This finishes the case $b=0$.

From now on we assume $b \neq 0$ and we want $c, d \in \mathbb{R}$ such that

$$
(c+d i)^{2}=\left(c^{2}-d^{2}\right)+2 c d i=a+b i .
$$

This means

$$
c^{2}-d^{2}=a, \quad 2 c d=b
$$

As $b \neq 0$ the desired $c, d$ should both be $\neq 0$ as well, and then

$$
c=\frac{b}{2 d}, \quad \text { so } \quad \frac{b^{2}}{4 d^{2}}-d^{2}=a
$$

Therefore the real number $d$ must be a zero of

$$
g=4 X^{4}+4 a X^{2}-b^{2} \in \mathbb{R}[X] .
$$

Since $g(0)<0$ and $g(x)>0$ for $x \in \mathbb{R}$ sufficiently large, Lemma V.1.5 implies that $d \in \mathbb{R}$ with $g(d)=0$ exists. Then taking $c=\frac{b}{2 d}$ shows that $a+b i$ has the square root $c+d i$ in $\mathbb{C}$. This proves Lemma V.1.6
V.1.7 Lemma. Any $f \in \mathbb{R}[X]$ of odd degree has a zero in $\mathbb{R}$.

Proof. We may assume that the leading coefficient of $f$ is positive. Then $f(x)>0$ for $x \in \mathbb{R}$ sufficiently large, and since $\operatorname{deg}(f)$ is odd we also have $f(x)<0$ for $x \in \mathbb{R}$ sufficiently negative. Hence Lemma V.1.5 implies that $f$ has a zero in $\mathbb{R}$.
V.1.8 Lemma. Suppose that every nonconstant $f \in \mathbb{R}[X]$ (so with real coefficients) has a zero in $\mathbb{C}$. Then $\mathbb{C}$ is algebraically closed.

Proof. Take $g=\sum_{i=0}^{n} a_{i} X^{i} \in \mathbb{C}[X], g \notin \mathbb{C}$. We have to show that $g$ has a zero in $\mathbb{C}$. Define

$$
\bar{g}:=\sum_{i=0}^{n} \bar{a}_{i} X^{i} \quad(\in \mathbb{C}[X])
$$

where $\bar{a}_{i}$ denotes the complex conjugate of $a_{i} \in \mathbb{C}$. Put

$$
f:=g \cdot \bar{g} \quad(\in \mathbb{C}[X]) .
$$

Clearly all coefficients of $f$ are real, so $f \in \mathbb{R}[X]$. Moreover $\operatorname{deg}(f)=2 \cdot \operatorname{deg}(g)$ hence $f$ is nonconstant. By the assumption of our lemma applied to $f$ some $\alpha \in \mathbb{C}$ exists with $f(\alpha)=0$. This means

$$
g(\alpha) \cdot \bar{g}(\alpha)=0 .
$$

Is $g(\alpha)=0$ then we are done. In case $g(\alpha) \neq 0$ one concludes $\bar{g}(\alpha)=0$, so

$$
\sum_{i=0}^{n} \bar{a}_{i} \alpha^{i}=0 .
$$

Taking the complex conjugate yields

$$
\sum_{i=0}^{n} a_{i} \bar{\alpha}^{i}=0
$$

so $g(\bar{\alpha})=0$ and again we found a zero in $\mathbb{C}$ of $g$. This proves Lemma V.1.8

We are now ready to prove the main result of this section.

Proof. (of TheoremV.1.4. By LemmaV.1.8 it suffices to show that any nonconstant $f \in \mathbb{R}[X]$ has a zero in $\mathbb{C}$. So take any such $f$. We may assume that $f$ is monic. Put $n=\operatorname{deg}(f)$. We write $n=2^{k} u$ with $k \in \mathbb{Z}_{\geq 0}$ and $u$ an odd positive integer. The proof will be done by mathematical induction w.r.t. $k$.

If $k=0$ then $\operatorname{deg}(f)$ is odd and Lemma V.1.7 shows that $f$ has a zero in $\mathbb{C}$ (even in $\mathbb{R}$ ).

Now assume $k \geq 1$, so $n$ is even. We use the field $L=\Omega_{\mathbb{C}}^{f} \supset \mathbb{C}$. In $L[X]$ we have

$$
f=\prod_{i=1}^{n}\left(X-\alpha_{i}\right), \quad \alpha_{i} \in L \quad(1 \leq i \leq n) .
$$

For any $c \in \mathbb{R}$ consider

$$
g_{c}=\prod_{1 \leq i<j \leq n}\left(X-\left(\alpha_{i}+\alpha_{j}+c \alpha_{i} \alpha_{j}\right)\right) \quad(\in L[X])
$$

The coefficients of $g_{c}$ are symmetric polynomials in $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}$, hence by TheoremIV.1.8 (applied to $R=\mathbb{R}, R^{\prime}=L$ ) we have $g_{c} \in \mathbb{R}[X]$.

The degree of $g_{c}$ equals the number of pairs $(i, j)$ with $1 \leq i<j \leq n$, which is $\frac{1}{2} n(n-1)=2^{k-1} \cdot u \cdot(n-1)$. Here $n-1$ is odd, so the number of factors 2 in $\operatorname{deg}\left(g_{c}\right)$ is $k-1$. Hence applying the induction hypothesis to $g_{c}$ we conclude that $g_{c}$ has a zero in $\mathbb{C}$. Now the zeros of $g_{c}$ are the $\frac{1}{2} n(n-1)$ expressions $\alpha_{i}+\alpha_{j}+c \alpha_{i} \alpha_{j}$. We conclude: for every $c \in \mathbb{R}$ there exist $i$ and $j$ with $1 \leq i<j \leq n$ and $\alpha_{i}+\alpha_{j}+c \alpha_{i} \alpha_{j} \in \mathbb{C}$.

Here $i$ and $j$ depend on $c$. However there are only finitely many possibilities for $i$ and $j$ while we have infinitely many $c \in \mathbb{R}$. So a pair of real numbers $c \neq c^{\prime}$ which give the same $i, j$. This means

$$
\alpha_{i}+\alpha_{j}+c \alpha_{i} \alpha_{j} \in \mathbb{C}, \quad \alpha_{i}+\alpha_{j}+c^{\prime} \alpha_{i} \alpha_{j} \in \mathbb{C} .
$$

Taking suitable linear combinations one finds that

$$
\beta=\alpha_{i}+\alpha_{j}, \quad \gamma=\alpha_{i} \alpha_{j} \in \mathbb{C} \quad \text { andhence } \quad\left(X-\alpha_{i}\right)\left(X-\alpha_{j}\right)=X^{2}-\beta X+\gamma \in \mathbb{C}[X] .
$$

Then Lemma V.1.6 tells us that this polynomial has a zero in $\mathbb{C}$, so $\alpha_{i} \in \mathbb{C}$ or $\alpha_{j} \in \mathbb{C}$. Hence $f$ has a zero in $\mathbb{C}$, as desired.

This proves the induction step and therefore Theorem V.1.4
V.1.9 Corollary. Every irreducible $f \in \mathbb{R}[X]$ has degree 1 or 2.

A polynomial $X^{2}+b X+c \in \mathbb{R}[X]$ is irreducible in $\mathbb{R}[X]$ if and only if $b^{2}-4 c<0$.
Proof. Let $f \in \mathbb{R}[X]$ be monic and irreducible, and let $\alpha \in \mathbb{C}$ be a zero of $f$. Then $f=f_{\mathbb{R}}^{\alpha}$, and

$$
\operatorname{deg}(f)=[\mathbb{R}(\alpha): \mathbb{R}] \leq[\mathbb{C}: \mathbb{R}]=2
$$

This proves the first assertion in V.1.9. The second one follows from

$$
X^{2}+b X+c=\left(X+\frac{1}{2} b\right)^{2}-\frac{1}{4}\left(b^{2}-4 c\right)
$$

and the fact that $X^{2}-a$ is irreducible in $\mathbb{R}[X]$ if and only if $a<0$. This shows V.1.9.

## V. 2 The algebraic closure

V.2.1 Definition. An algebraic closure of a field $K$ is a field extension $\bar{K} \supset K$ with the properties
i. $\bar{K}$ is algebraic over $K$;
ii. $\bar{K}$ is algebraically closed.
V.2.2 Example. From V.1.4 we know that $\mathbb{C}$ is an algebraic closure of of $\mathbb{R}$.
V.2.3 Theorem. The field $\mathbb{Q}$ of rational numbers has an algebraic closure.

Proof. Put

$$
\overline{\mathbb{Q}}=\{\alpha \in \mathbb{C}: \alpha \text { is algebraic over } \mathbb{Q}\} .
$$

From basic properties of the notion 'algebraic' we have that indeed $\overline{\mathbb{Q}}$ is a field. We verify condition (ii), which then implies that $\overline{\mathbb{Q}}$ is an algebraic closure of $\mathbb{Q}$.

Laat $f \in \overline{\mathbb{Q}}[X]$ be nonconstant. We need to show that $f$ has a zero in $\overline{\mathbb{Q}}$. Using Theorem V.1.4 one obtains a zero $\alpha$ of $f$ in $\mathbb{C}$. As $f \in \bar{Q}[X]$, this $\alpha$ is algebraic over $\overline{\mathbb{Q}}$. Moreover the field $\overline{\mathbb{Q}}$ is algebraic over $\mathbb{Q}$, and therefore $\overline{\mathbb{Q}}(\alpha)$ is an algebraic extension of $\mathbb{Q}$. In particular $\alpha$ is algebraic over $\mathbb{Q}$, so $\alpha \in \overline{\mathbb{Q}}$. This proves V.2.3.
V.2.4 Remark. The algebraic closure of $\mathbb{Q}$ certainly does not equal $\mathbb{C}$, since (uncountably many) transcendental complex numbers exist.

More generally we have:
V.2.5 Theorem. Every field $K$ has an algebraic closure $\bar{K}$.

Moreover $\bar{K}$ is unique up to $K$-isomorphism, meaning: if $\bar{K}_{1}$ and $\bar{K}_{2}$ are algebraic closures of $K$, then $\bar{K}_{1} \cong_{K} \bar{K}_{2}$.

Proofs of this for the general case necessarily depend on "Zorn's Lemma", the same axiom in the foundations of set theory that was also needed to prove the existence of a maximal ideal in an arbitrary ring. For completeness, we copy here three proofs taken from J.S. Milne's online lecture notes Fields and Galois Theory, see www.jmilne.org/math/CourseNotes/FTc.pdf(Version 4.52, March 17, 2017, pages 86-88).

## First proof of the existence of algebraic closures

(Bourbaki, Algèbre, Chap. V, §4.) An $F$-algebra is a ring containing $F$ as a subring. Let $\left(A_{i}\right)_{i \in I}$ be a family of commutative $F$-algebras, and define $\otimes_{F} A_{i}$ to be the quotient of the $F$-vector space with basis $\prod_{i \in I} A_{i}$ by the subspace generated by elements of the form:
$\left(x_{i}\right)+\left(y_{i}\right)-\left(z_{i}\right)$ with $x_{j}+y_{j}=z_{j}$ for one $j \in I$ and $x_{i}=y_{i}=z_{i}$ for all $i \neq j ;$ $\left(x_{i}\right)-a\left(y_{i}\right)$ with $x_{j}=a y_{j}$ for one $j \in I$ and $x_{i}=y_{i}$ for all $i \neq j$,
(ibid., Chap. II, 3.9). It can be made into a commutative $F$-algebra in an obvious fashion, and there are canonical homomorphisms $A_{i} \rightarrow \bigotimes_{F} A_{i}$ of $F$-algebras.

For each polynomial $f \in F[X]$, choose a splitting field $E_{f}$, and let $\Omega=\left(\otimes_{F} E_{f}\right) / M$ where $M$ is a maximal ideal in $\bigotimes_{F} E_{f}$ (whose existence is ensured by Zorn's lemma). Note that $F \subset \otimes_{F} E_{f}$ and $M \cap F=0$. As $\Omega$ has no ideals other than (0) and $\Omega$, it is a field (see 1.2). The composite of the $F$-homomorphisms $E_{f} \rightarrow \bigotimes_{F} E_{f} \rightarrow \Omega$, being a homomorphism of fields, is injective. Since $f$ splits in $E_{f}$, it must also split in the larger field $\Omega$. The algebraic closure of $F$ in $\Omega$ is therefore an algebraic closure of $F$.

## Second proof of the existence of algebraic closures

(Jacobson 1964, p144.) After (4.24) we may assume $F$ to be infinite. This implies that the cardinality of every field algebraic over $F$ is the same as that of $F$ (ibid. p143). Choose an uncountable set $\Xi$ of cardinality greater than that of $F$, and identify $F$ with a subset of $\Xi$. Let $S$ be the set of triples $(E,+, \cdot)$ with $E \subset \Xi$ and $(+, \cdot)$ a field structure on $E$ such that $(E,+, \cdot)$ contains $F$ as a subfield and is algebraic over it. Write $(E,+, \cdot) \leq\left(E^{\prime},+^{\prime},,^{\prime}\right)$ if the first is a subfield of the second. Apply Zorn's lemma to show that $S$ has maximal elements, and then show that a maximal element is algebraically closed.

## Third proof of the existence of algebraic closures

(Emil Artin.) Consider the polynomial ring $F\left[\ldots, x_{f}, \ldots\right]$ in a family of symbols $x_{f}$ indexed by the nonconstant monic polynomials $f \in F[X]$. If 1 lies in the ideal $I$ of $F\left[\ldots, x_{f}, \ldots\right]$ generated by the polynomials $f\left(x_{f}\right)$, then

$$
g_{1} f_{1}\left(x_{f_{1}}\right)+\cdots+g_{n} f_{n}\left(x_{f_{n}}\right)=1 \quad\left(\text { in } F\left[\ldots, x_{f}, \ldots\right]\right)
$$

for some $g_{i} \in F\left[\ldots, x_{f}, \ldots\right]$ and some nonconstant monic $f_{i} \in F[X]$. Let $E$ be an extension of $F$ such that each $f_{i}, i=1, \ldots, n$, has a root $\alpha_{i}$ in $E$. Under the $F$-homomorphism $F\left[\ldots, x_{f}, \ldots\right] \rightarrow E$ sending

$$
\left\{\begin{array}{l}
x_{f_{i}} \mapsto \alpha_{i} \\
x_{f} \mapsto 0, \quad f \notin\left\{f_{1}, \ldots, f_{n}\right\}
\end{array}\right.
$$

the above relation becomes $0=1$. From this contradiction, we deduce that 1 does not lie in $I$, and so Proposition 6.4 applied to $F\left[\ldots, x_{f}, \ldots\right] / I$ shows that $I$ is contained in a maximal ideal $M$ of $F\left[\ldots, x_{f}, \ldots\right]$. Let $\Omega=F\left[\ldots, x_{f}, \ldots\right] / M$. Then $\Omega$ is a field containing (a copy of) $F$ in which every nonconstant polynomial in $F[X]$ has at least one root. Repeat the process starting with $E_{1}$ instead of $F$ to obtain a field $E_{2}$. Continue in this fashion to obtain a sequence of fields

$$
F=E_{0} \subset E_{1} \subset E_{2} \subset \cdots,
$$

and let $E=\bigcup_{i} E_{i}$. Then $E$ is algebraically closed because the coefficients of any nonconstant polynomial $g$ in $E[X]$ lie in $E_{i}$ for some $i$, and so $g$ has a root in $E_{i+1}$. Therefore, the algebraic closure of $F$ in $E$ is an algebraic closure of $F$.

## (Non)uniqueness of algebraic closures

THEOREM $6.8\left({ }^{(*)}\right.$ Let $\Omega$ be an algebraic closure of $F$, and let $E$ be an algebraic extension of $F$. There exists an $F$-homomorphism $E \rightarrow \Omega$, and, if $E$ is also an algebraic closure of $F$, then every such homomorphism is an isomorphism.

Proof. Suppose first that $E$ is countably generated over $F$, i.e., $E=F\left[\alpha_{1}, \ldots, \alpha_{n}, \ldots\right]$. Then we can extend the inclusion map $F \rightarrow \Omega$ to $F\left[\alpha_{1}\right]$ (map $\alpha_{1}$ to any root of its minimal polynomial in $\Omega$ ), then to $F\left[\alpha_{1}, \alpha_{2}\right]$, and so on (see 2.2).

In the uncountable case, we use Zorn's lemma. Let $S$ be the set of pairs ( $M, \varphi_{M}$ ) with $M$ a field $F \subset M \subset E$ and $\varphi_{M}$ an $F$-homomorphism $M \rightarrow \Omega$. Write $\left(M, \varphi_{M}\right) \leq\left(N, \varphi_{N}\right)$ if $M \subset N$ and $\varphi_{N} \mid M=\varphi_{M}$. This makes $S$ into a partially ordered set. Let $T$ be a totally ordered subset of $S$. Then $M^{\prime}=\bigcup_{M \in T} M$ is a subfield of $E$, and we can define a homomorphism $\varphi^{\prime}: M^{\prime} \rightarrow \Omega$ by requiring that $\varphi^{\prime}(x)=\varphi_{M}(x)$ if $x \in M$. The pair ( $\left.M^{\prime}, \varphi^{\prime}\right)$ is an upper bound for $T$ in $S$. Hence Zorn's lemma gives us a maximal element $(M, \varphi)$ in $S$. Suppose that $M \neq E$. Then there exists an element $\alpha \in E, \alpha \notin M$. Since $\alpha$ is algebraic over $M$, we can apply (2.2) to extend $\varphi$ to $M[\alpha]$, contradicting the maximality of $M$. Hence $M=E$, and the proof of the first statement is complete.

If $E$ is algebraically closed, then every polynomial $f \in F[X]$ splits in $E[X]$ and hence in $\varphi(E)[X]$. Let $\alpha \in \Omega$, and let $f(X)$ be the minimum polynomial of $\alpha$. Then $X-\alpha$ is a factor of $f(X)$ in $\Omega[X]$, but, as we just observed, $f(X)$ splits in $\varphi(E)[X]$. Because of unique factorization, this implies that $\alpha \in \varphi(E)$.

## V. 3 Exercises

1. Suppose $K=\mathbb{F}_{q}$ is a finite field of charakteristic $p$. Let $f=1+\prod_{\alpha \in K}(X-\alpha)$ be the polynomial used in the proof of V.1.3. Put $L=\Omega_{K}^{f}$. Show that:
(a) $f=X^{q}-X+1$;
(b) for every $\alpha \in L$ such that $f(\alpha)=0$ it holds that

$$
\alpha^{q^{i}}=\alpha-\bar{i} \quad \text { for all } i \in \mathbb{Z}_{>0} ;
$$

here $\bar{i}=(i \bmod p) \in \mathbb{F}_{p} \subset K$. In particular, show that

$$
\alpha^{q^{p}}=\alpha ;
$$

(c) $L=\mathbb{F}_{q^{p}}$;
(d) every irreducible factor of $f$ in $K[X]$ has degree $p$.
2. Suppose $\bar{K}$ is an algebraic extension of the field $K$ such that for every monic $f \in K[X]$ the field $\bar{K}$ contains a splitting field of $f$ over $K$. Show that $\bar{K}$ is an algebraic closure of $K$.
3. Let $\overline{\mathbb{Q}}$ be the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$. Show that $[\overline{\mathbb{Q}}: \overline{\mathbb{Q}} \cap \mathbb{R}]=2$.

## VI. 1 Definitions

The analogue of a vectorspace over a field is a module over a ring.
VI.1.1 Definition. Let $R$ be a (unitary) ring. A left- $R$-module $M$ is an abelian group ( $M,+, 0$ ) with an action of the ring $R$, which means a map:

$$
R \times M \longrightarrow M, \quad(a, m) \mapsto a m
$$

such that for all $a, b \in R$ and all $m, n \in M$ it holds that:
(RM1) $a(m+n)=a m+a n$,
(RM2) $(a+b) m=a m+b m$,
(RM3) $a(b m)=(a b) m$,
(RM4) $1 m=m$.
A right- $R$-module is defined analogously, but in that case with an action $M \times R \rightarrow M$ et cetera.
VI.1.2 Remark. Unless explicitly stated otherwise, all modules considered here are left- $R$-modules and we simply write $R$-modules.
VI.1.3 Example. Let $R=K$ be a field. In this case the axioms for a left- $R$-module are exactly the axioms for a vectorspace over $K$.
VI.1.4 Example. Let $K$ be a field. The ring $R$ consisting of all $n \times n$ matrices with coefficients in $K$ acts on the additive group ( $\left.K^{n}\right)^{+}$of the vectorspace $K^{n}$ :

$$
R=M(n, K), \quad M=\left(K^{n}\right)^{+} \quad \text { with action } \quad(A, v) \mapsto A v,
$$

the usual product of a matrix and a vector. This turns $K^{n}$ into a left- $M(n, K)$ module.
VI.1.5 Example. Take $R=\mathbb{Z}$ and let $M=G$ be an abelian group. Define an action of $\mathbb{Z}$ on $G$ using $0 g:=e,(-1) g:=-g$ with $e \in G$ the unit element and $-g$ the inverse of $g$ in $G$. Next, define $n g:=g+g+\ldots+g(n$ times $)$ and $(-n) g=(-g)+(-g)+\ldots+(-g)$ ( $n$ times) whenever $n \geq 1$. Verify that in this way $G$ is a $\mathbb{Z}$-module.
VI.1.6 Example. Suppose $R$ is a ring and $I \subset R$ is an ideal. Then $I$ is an $R$-module. Moreover $R / I$ is an $R$-module as well, with action:

$$
R \times R / I \longrightarrow R / I \quad(r, a+I) \mapsto r a+I
$$

Note that in this way $r \bar{a}:=r a+I=\overline{r a}$.
VI.1.7 Remark. The axioms of a left- $R$-module imply:

$$
0 m=(0+0) m=0 m+0 m \quad \text { hence } \quad 0 m=0
$$

Moreover:

$$
0=(a+(-a)) m=a m+(-a) m \quad \text { so } \quad(-a) m=-(a m),
$$

in particular $(-1) m=-(1 m)=-m$, hence $a(-m)=a((-1) m)=(a(-1)) m=(-a) m$, which justifies the notation $-a m$ for $(-a) m=-(a m)=a(-m)$.
VI.1.8 Definition. A submodule $N$ of a left- $R$-module $M$ is a subgroup of $M$ which is closed under the action of $R$. This means that $N \subset M$ has the property that for all $a, b \in N$ and all $r \in R$ one has:
(SM1) $\quad 0 \in N, \quad a-b \in N$,
(SM2) $\quad r a \in N, \quad$ (one also writes $r N \subseteq N$ ).
A submodule $N$ of the $R$-module $M$ is itself also an $R$-module.
VI.1.9 Examples. 1. Take a field $R=K$, and let $M$ be a vectorspace over $K$. The submodules of $M$ are precisely the linear subspaces of $M$.
2. The only submodules of the $M(n, K)$-module $K^{n}$ are $\{0\}$ and $K^{n}$. Namely, if $x, y \in K^{n}-\{0\}$ then $A \in M(n, K)$ exists with $A x=y$ (verify this, and check how it implies the claim!).
3. The submodules of the $\mathbb{Z}$-module $G$, with $G$ any abelian group, are precisely the subgroups of $G$ (check for yourself!).
4. The ring $R$ is itself an $R$-module, and the $R$-submodules of $R$ are precisely the ideals of $R$ (as one easily verifies!).
5. Take $A \in M(n, K)$ where $K$ is a field, and let $R=K[A]=\left\{\sum_{i<\infty} a_{i} A^{i}: a_{i} \in K\right\}$. Suppose $v \in K^{n}$ is an eigenvector of $A$ with eigenvalue $\lambda \in K$. Then the (onedimensional) subspace $K v$ is a submodule of the $R$-module $K^{n}$. Namely, $K v$ is an additive group and since $A v=\lambda v \in K v$, also $\left(\sum_{i<\infty} a_{i} A^{i}\right)(v) \in K v$.

## VI. $2 \quad R$-module homomorphisms

The $R$-module homomorphisms generalise the linear maps:
VI.2.1 Definition. Let $R$ be a ring and let $M, N$ be (left)- $R$-modules. An $R$-module homomorphism is a map

$$
f: M \longrightarrow N
$$

which is a homomorphism of abelian groups and which is moreover $R$-linear. In other words, for all $x, y \in M$ and for all $r \in R$ the map $f$ satisfies:
(H1) $f(x+y)=f(x)+f(y)$,
(H2) $\quad f(r x)=r f(x)$.
In particular it follows that $f(0)=0$.
An $R$-module isomorphism is a bijective $R$-module homomorphism. (Note that the inverse of an $R$-module isomorphism is again a (bijective) $R$-module homomorphism and therefore it is again an $R$-module isomorphism.)
The kernel of an $R$-module homomorphism $f: M \rightarrow N$ is defined as

$$
\operatorname{Ker}(f):=\{m \in M: f(m)=0(\in N)\} .
$$

The image of an $R$-module homomorphism is defined as

$$
\operatorname{Im}(f)=f[M]=\{f(m) \in N: m \in M\} .
$$

VI.2.2 Examples. 1. Since $0=\{0\} \subset R$ is an $R$-module (it is an ideal in $R$ ), it follows that the obvious maps

$$
0 \longrightarrow M, \quad M \longrightarrow 0
$$

are $R$-module homomorphisms (the first one sends $0 \in R$ to $0 \in M$, the second one sends every $m \in M$ to $0 \in R$ ).
2. If $R$ is a commutative ring and $M$ is an $R$-module and $a \in R$, then

$$
\phi_{a}: M \longrightarrow M, \quad m \mapsto a m
$$

is an $R$-module homomorphism.
3. A homomorphism $f: G \rightarrow H$ of abelian groups is a $\mathbb{Z}$-module homomorphism, namely $f(2 g)=f(g+g)=f(g)+f(g)=2 f(g)$ et cetera.
4. Let $R$ be a ring and $I \subset R$ an ideal. Then $R$ and $R / I$ are $R$-modules. The canonical map (a ring homomorphism):

$$
\phi: R \longrightarrow R / I, \quad m \mapsto \bar{m}=m+I
$$

is an $R$-module homomorphism. Namely, condition (H1) is satisfied since $\phi$ is a ring homomorphism and therefore a homomorphism of additive groups. Moreover (see Example VI.1.6):

$$
\phi(r m):=r m+I=r \phi(m), \quad(r, m \in R)
$$

so also (H2) is satisfied.
VI.2.3 Theorem. Let $f: M \rightarrow N$ be an $R$-module homomorphism. Then:

1. $\operatorname{Ker}(f)$ is a submodule of $M$.
2. $\operatorname{Im}(f)$ is a submodule of $N$.

Proof. These assertions are immediate consequences of the definitions (check for yourself!).
VI.2.4 Theorem. If $f: M \rightarrow N$ and $g: N \rightarrow P$ are $R$-module homomorphisms, then the composition $g \circ f: M \rightarrow P$, so

$$
g \circ f: M \xrightarrow{f} N \xrightarrow{g} P
$$

is an $R$-module homomorphism as well.
Proof. This is immediate from the definition.
VI.2.5 Remark. One usually writes $g f$ for the composition $g \circ f$ of the $R$-module homomorphisms $f$ and $g$.

## VI. 3 Direct sums

All modules considered in this section are left-modules.
VI.3.1 Definition. Let $R$ be a ring and let $I$ be a nonempty set. For every $i \in I$ let $M_{i}$ be an $R$-module. The direct sum $M$ of the $M_{i}$ 's is defined as
$M=\bigoplus_{i \in I} M_{i}:=\left\{\left(x_{i}\right)_{i \in I}: x_{i} \in M_{i}\right.$ for all $i \in I$ and $x_{i} \neq 0$ for only finitely many $\left.i\right\}$.

This set is given the structure of an $R$-module by defining

$$
\begin{array}{rlr}
\left(x_{i}\right)_{i \in I}+\left(y_{i}\right)_{i \in I} & =\left(z_{i}\right)_{i \in I}, & \text { with } z_{i}=x_{i}+y_{i} \forall i \in I \\
0 & =(0)_{i \in I}, & \text { (so the } i \text {-th coordinate is } 0 \in M_{i}, \forall i \in I \\
r \cdot\left(x_{i}\right)_{i \in I} & =\left(z_{i}\right)_{i \in I} & \text { with } z_{i}=r x_{i} \forall i \in I .
\end{array}
$$

It is not hard to check that in this way indeed an $R$-module is defined.
VI.3.2 Definition. An $R$-module $F$ is called free (or, free $R$-module) if a nonempty set $I$ exists, and an isomorphism of $R$-modules

$$
F \stackrel{\cong}{\cong} \oplus_{i \in I} R .
$$

(So here we take all modules $M_{i}$ in VI.3.1 equal to $R$.)

Given a nonempty set $I$, the $R$-module $F:=\oplus_{i \in I} R$ is by definition free. For every $i \in I$, define

$$
e_{i}=\left(\ldots, x_{j}, \ldots\right)_{j \in I} \in F \quad \text { by } \quad x_{i}=1, \quad \text { and } x_{j}=0 \text { for all } j \neq i .
$$

Then every $x \in F$ can be written in a unique way as

$$
x=\sum_{i \in I} x_{i} e_{i}, \quad\left(x, e_{i} \in F, x_{i} \in R\right)
$$

(which is a finite sum since only finitely many $i \in I$ exist with $x_{i} \neq 0$ ).
As a special case, for $n \in \mathbb{Z} \geq 1$ one defines the free $R$-module

$$
R^{n}:=\bigoplus_{i \in\{1,2, \ldots n\}} R .
$$

VI.3.3 Examples. If $K$ is a field, $K^{n}$ is the familiar vector space over $K$.

The polynomial ring $R[X]$ is an $R$-module (for the usual addition and scalar multiplication). The map

$$
\phi: R[X] \longrightarrow F:=\bigoplus_{i \in \mathbb{Z}_{\geq 0}} R, \quad \sum_{i=0}^{n} a_{i} X^{i} \mapsto\left(a_{0}, a_{1}, \ldots, a_{n}, 0,0, \ldots\right),
$$

is an isomorphism of $R$-modules; $\phi$ is surjective since only finitely many $x_{i}$ differ from 0 in an element $\left(\ldots, x_{i}, \ldots\right) \in F$.

For $R$ non-commutative rings $R$ it turns out to be possible (see Exercise 7 on page 57) that $R \cong R^{2}$ (!). This is not possible in case $R$ is commutative:
VI.3.4 Theorem. If $R \neq(0)$ is a commutative (unitary) ring, then

$$
R^{m} \cong R^{n} \quad \Longrightarrow \quad m=n
$$

The number $m(=n)$ here is called the rank of the $R$-module $R^{m}$. More generally

$$
R^{m} \cong \oplus_{i \in I} R \quad \Longrightarrow \quad m=\# I
$$

Proof. A proof based on the well known property $\operatorname{det}(A B)=\operatorname{det}(A) \operatorname{det}(B)$ of $\operatorname{deter}$ minants of $n \times n$ matrices, is given in Section 3.4 of the book N. Jacobson, Basic Algebra I (San Fransisco: W.H. Freeman and Company, 1974). The advantage of that proof is that no use is made of Zorn's Lemma. The next argument uses the existence of a maximal ideal in a ring, which in general does require Zorn's Lemma.

Suppose $\phi: R^{m} \rightarrow R^{n}$ is an $R$-module isomorphism. Take $M$ a maximal ideal in $R$. Then $M^{m}=\oplus_{i=1}^{m} M$ is a submodule of $R^{m}$. The factor group $R^{m} / M^{m}$ is in a natural way a $K=R / M$-module, hence since $K$ is a field, $R^{m} / M^{m}$ is a vectorspace over $K$. One easily verifies that the map

$$
\left(r_{1}, \ldots, r_{m}\right)+M^{m} \mapsto\left(r_{1}+M, \ldots, r_{m}+M\right)
$$

yields a well-defined isomorphism between vectorspaces $R^{m} / M^{m}$ and $K^{m}$ over $K$. Hence $\operatorname{dim}_{K}\left(R^{m} / M^{m}\right)=m$.

The set $\phi\left(M^{m}\right)$ is a submodule of $R^{n}$. The factor group $R^{n} / \phi\left(M^{m}\right)$ therefore is an $R$-module, and even a $K$-module: namely, for $m \in M$ and $v \in R^{n}$ one has $m v=\phi \phi^{-1}(m v)=\phi\left(m \phi^{-1}(v)\right) \in \phi\left(M^{m}\right)$. Hence multiplying an element $r \in R$ by a class $v+\phi\left(M^{m}\right) \in R^{n} / \phi\left(M^{m}\right)$ only depends on the class of $r$ in $R / M=K$, which exactly means that $R^{n} / \phi\left(M^{m}\right)$ is a $K$-module. The isomorphism $\phi$ therefore induces an isomorphism $R^{m} / M^{m} \cong R^{n} / \phi\left(M^{n}\right)$ between vectorspaces over $K$, so in particular $\operatorname{dim}_{K}\left(R^{n} / \phi\left(M^{m}\right)\right)=m$.

Every $x \in R^{n}$ can be written in a unique way as

$$
\begin{aligned}
x & =\left(x_{1}, x_{2}, \ldots, x_{n}\right) \\
& =x_{1} \cdot(1,0 \ldots, 0)+x_{2} \cdot(0,1,0, \ldots, 0)+\ldots+x_{n} \cdot(0,0, \ldots, 1) \\
& =x_{1} e_{1}+x_{2} e_{2}+\ldots+x_{n} e_{n}
\end{aligned}
$$

Hence the vectorspace $R^{n} / \phi\left(M^{m}\right)$ over $K$ is spanned by $e_{1}+\phi\left(M^{m}\right), \ldots, e_{n}+\phi\left(M^{m}\right)$. We conclude that $n \geq \operatorname{dim}_{K}\left(R^{n} / \phi\left(M^{m}\right)\right)=m$.

Interchanging the roles of $R^{n}$ and $R^{m}$ in the argument above, and replacing $\phi$ by $\phi^{-1}$, it follows analogously that also $m \geq n$. Hence $n=m$.

The more general assertion in the theorem follows completely analogously. This shows Theorem VI.3.4.

We saw that modules over a field $K$ are in fact vectorspaces over $K$. In particular: if $K^{n}$ is the direct sum of two $K$-modules:

$$
K^{n} \cong V \oplus W, \quad \text { then } \quad V \cong K^{a} \quad \text { and } \quad W \cong K^{n-a}
$$

for some $a$, since $V$ and $W$ are (finite dimensional) vectorspaces over $K$ as well. The case of modules over a ring that is not a field is much more interesting. For a simple example, take the ring $R=\mathbb{Z} / 6 \mathbb{Z}$ and the $R$-modules $\mathbb{Z} / 2 \mathbb{Z}$ and $\mathbb{Z} / 3 \mathbb{Z}$, then

$$
\mathbb{Z} / 6 \mathbb{Z} \stackrel{\cong}{\rightrightarrows} \mathbb{Z} / 2 \mathbb{Z} \oplus \mathbb{Z} / 3 \mathbb{Z}, \quad n+6 \mathbb{Z} \mapsto(n+2 \mathbb{Z}, n+3 \mathbb{Z})
$$

Using the Chinese Remainder Theorem many more examples of this kind can be constructed. The next two examples reveal a slightly more surprising submodules of a free module over a ring.
VI.3.5 Example. Take $R=\mathbb{Z}[\sqrt{-5}]$. The map $\varphi: R \rightarrow \mathbb{F}_{2}$ defined by $\varphi(a+b \sqrt{5})=\bar{a}+\bar{b}$, with $a, b \in \mathbb{Z}$ and for any $n \in \mathbb{Z}$ using the notation $\bar{n}:=n \bmod 2$, is the surjective ring homomorphism. Hence $M:=\operatorname{Ker}(\varphi)$ is a maximal ideal in $R$, so in particular $M$ is an $R$-module. Note that

$$
\begin{aligned}
M & =\{a+b \sqrt{-5}: \exists k \in \mathbb{Z} \text { such that } a+b=2 k\} \\
& =\{2 k-b+b \sqrt{-5}: b, k \in \mathbb{Z}\} \\
& =\mathbb{Z} \cdot 2+\mathbb{Z} \cdot(-1+\sqrt{-5})
\end{aligned}
$$

and every element in $M$ can be written in a unique way as $n \cdot 2+m \cdot(-1+\sqrt{-5})$ with $n, m \in \mathbb{Z}$.

We claim that

$$
M \oplus M \cong R^{2}
$$

as $R$-modules. Indeed, the map $(\alpha, \beta) \mapsto(2 \alpha+(-1+\sqrt{-5}) \beta,(1+\sqrt{-5}) \alpha-2 \beta)$ for $\alpha, \beta \in R^{2}$ defines an $R$-module homomorphism from $R^{2}$ to $M \oplus M$. The "inverse" $R$-module homomorphism $M \oplus M \rightarrow R^{2}$ is given by

$$
\begin{gathered}
(a \cdot 2+b \cdot(-1+\sqrt{-5}), c \cdot 2+d \cdot(-1+\sqrt{-5})) \\
(-2 a+b+c+2 d+(d-b-c) \sqrt{-5},-a+3 b+2 c-d+(d-a) \sqrt{-5})
\end{gathered}
$$

for $a, b, c, d \in \mathbb{Z}$.
Next, we claim that $M$ itself is not free as an $R$-module. Indeed, if $M \cong R^{n}$ for some $n$, then $R^{2} \cong M \oplus M \cong R^{n} \oplus R^{n} \cong R^{2 n}$. So we necessarily have $n=1$. However, if $M \cong R$ as $R$-modules, then an isomorphism $f: R \rightarrow M$ would send $1 \in R$ to some element $e=f(1) \in M$, and then any $r \in R$ to $f(r)=f(r \cdot 1)=r \cdot f(1)=r e \in M$. As $f$ is supposed to be an isomorphism, this would imply $M=R e$, so $M$ is a principal ideal. In particular, the norm of any element of $M$ would be a multiple of the norm of $e$. Now $2 \in M$ has norm 4 and $-1+\sqrt{-5}$ has norm 6 , so the norm of $e$ divides $\operatorname{gcd}(4,6)=2$. Since $R=\mathbb{Z}[\sqrt{-5}]$ does not contain any element of norm 2 , we conclude that $e$ must have norm 1 , so $e \in R^{\times}$and therefore $M=R e=R$, which is absurd since, e.g., $1 \notin M$.

So in this example the module $M$ is not free, while $M \oplus M$ is free!
VI.3.6 Example. Consider the subring $R$ of periodic functions with period $2 \pi$ of the ring of all real $C^{\infty}$-functions on $\mathbb{R}$ :

$$
R:=\left\{f \in C^{\infty}(\mathbb{R}): f(x+2 \pi)=f(x) \forall x \in \mathbb{R}\right\} .
$$

Let $M$ be the $R$-module defined as

$$
M:=\left\{m \in C^{\infty}(\mathbb{R}): m(x+2 \pi)=-m(x) \forall x \in \mathbb{R}\right\} .
$$

Here addition and scalar multiplication on $M$ are given by

$$
\begin{array}{rlr}
(m+n)(x) & :=m(x)+n(x), & m, n \in M, x \in \mathbb{R} \\
(f m)(x) & :=f(x) m(x), & f \in R, m \in M, x \in \mathbb{R} ;
\end{array}
$$

note that indeed $f m \in M(!)$ and check for yourself that in this way $M$ is an $R$ module. We claim that like in ExampleVI.3.5 we have that

$$
M \nexists R, \quad \text { whereas } \quad M \oplus M \cong R^{2},
$$

and $M$ is not free.
Well known properties of the sine and cosine functions imply

$$
\forall a \in \mathbb{R}: C_{a}, S_{a} \in M, \quad \text { where } \quad \begin{cases}C_{a}: \mathbb{R} \rightarrow \mathbb{R}, & x \mapsto \cos \frac{x-a}{2}, \\ S_{a}: \mathbb{R} \rightarrow \mathbb{R}, & x \mapsto \sin \frac{x-a}{2} .\end{cases}
$$

Suppose an $R$-module isomorphism between $R$ and $M$ exists:

$$
\phi: R \longrightarrow M, \quad \text { then } \quad \phi(f)=f \cdot \phi(1)
$$

with $1 \in R$ the constant function 1 . Write $g:=\phi(1) \in M$. Since $g$ is continuous and $g(2 \pi)=-g(0)$ one concludes that $g$ has a zero $a \in[0,2 \pi]$.

As $\phi$ is supposed to be surjective, every $m \in M$ can be written as $m=\phi(f)$ for some $f \in R$, hence

$$
m=f g, \quad \text { implying } \quad m(a)=f(a) g(a)=0 .
$$

However, taking $m=C_{a} \in M$ one obtains $C_{a}(a)=\cos \frac{a-a}{2}=1 \neq 0$, a contradiction. We conclude that no isomorphism $\phi: R \rightarrow M$ exists, which shows the first assertion.

Next we show $R^{2} \cong M \oplus M$. Define

$$
\psi: R^{2} \longrightarrow M \oplus M, \quad(f, g) \mapsto\left(f C_{0}+g S_{0},-f S_{0}+g C_{0}\right)
$$

(this clearly is an $R$-module homomorphism). In terms of matrices (with coefficients in $M \subset C^{\infty}(\mathbb{R})$ ) this $\psi$ is given by

$$
A:=\left(\begin{array}{rr}
C_{0} & S_{0} \\
-S_{0} & C_{0}
\end{array}\right), \quad \text { so in particular } \quad A^{-1}=\left(\begin{array}{rr}
C_{0} & -S_{0} \\
S_{0} & C_{0}
\end{array}\right)
$$

where we used that $C_{0}^{2}+S_{0}^{2}=1$.
We claim that indeed the inverse of $\psi$ is given by

$$
\psi^{-1}: M \oplus M \longrightarrow R^{2}, \quad(m, n) \mapsto\left(C_{0} m-S_{0} n, S_{0} m+C_{0} n\right)
$$

To verify this, first note that the product of any two functions in $M$ is an element of $R$, hence $\psi^{-1}$ indeed maps $M \oplus M$ to $R^{2}$ ). Clearly $\psi^{-1}$ is an $R$-module homomorphism, and $\psi^{-1} \psi=\operatorname{id}_{R^{2}}$ and $\psi \psi^{-1}=\mathrm{id}_{M \oplus M}$. So $\psi^{-1}$ is the inverse of $\psi$.

Finally, we show that $M$ is not free. Exactly as in Example VI.3.5 above, if $M \cong \oplus_{i \in I} R$ for some set $I$, then $\# I=1$ and $R \cong M$. Since we already showed that $R \nexists M$ it follows that $M$ is not free.
VI.3.7 Remark. The module $M$ in Example VI.3.6 turns out to be isomorphic (as an $R$-module) to some (non-principal) ideal in $R$ (see Exercise 3 on page 56):

$$
M \cong I:=\operatorname{Ker}\left(\mathrm{ev}_{0}: R \longrightarrow \mathbb{R}\right), \quad \text { with } \quad \mathrm{ev}_{0}: f \mapsto f(0)
$$

The next result describes a "universal property" of direct sums: this means a property determining it up to isomorphisms of modules.
VI.3.8 Theorem. Let $R$ be a ring and let $I$ be a (nonempty) set. Suppose that for every $i \in I$ an $R$-module $M_{i}$ is given. Then $\oplus_{i \in I} M_{i}$ together with the $R$-module homomorphisms $\iota_{i}: M_{i} \rightarrow \oplus_{i \in I} M_{i}$ defined for $m \in M_{i}$ by $\iota_{i}(m)=\left(x_{j}\right)_{j \in I}$ (where $x_{j}=0 \in M_{j}$ if $j \neq i$ and $x_{i}=m$ ), have the following property.

For every $R$-module $M$, given $R$-module homomorphisms $f_{i}: M_{i} \rightarrow M$ for all $i \in I$, there exists a unique(!) $R$-module homomorphism $f: \oplus_{i \in I} M_{i} \rightarrow M$ such that $f_{i}=f \circ \iota_{i}$ for every $i \in I$.

Is $D$ any $R$-module equipped with homomorphisms $j_{i}: M_{i} \rightarrow D$ such that the same property holds:

For every $R$-module $M$, given $R$-module homomorphisms $f_{i}: M_{i} \rightarrow M$ for all $i \in I$, there exists a unique(!) $R$-module homomorphism $f: D \rightarrow M$ such that $f_{i}=f \circ j_{i}$ for every $i \in I$,
then $D \cong \oplus_{i \in I} M_{i}$.
Proof. Given maps $f_{i}: M_{i} \rightarrow M$, define $f: \oplus_{i \in I} M_{i} \rightarrow M$ by $f\left(\left(x_{i}\right)_{i \in I}\right)=\sum_{i} f_{i}\left(x_{i}\right)$. By the definition of 'direct sum' and the fact that any $R$-module homomorphism sends 0 to 0 , this is a finite sum. Evidently $f$ is an $R$-module homomorphism, and it has the property $f_{i}=f \circ \iota_{i}$ for all $i \in I$. Since any $\left(m_{i}\right)_{i \in I} \in \oplus_{i \in I} M_{i}$ can be written as a finite sum $\sum_{i \in I} \iota_{i}\left(m_{i}\right)$, it follows that any $R$-module homomorphism $\oplus_{i \in I} M_{i} \rightarrow M$ is completely determined by its restrictions to the $\iota_{i}\left(M_{i}\right), i \in I$. So $f$ is unique.

If $D$, with the maps $j_{i}$, has the same property, then we apply this property to the $R$-module $\oplus_{i \in I} M_{i}$ and the homomorphisms $\iota_{i}$. This gives us a unique $f: D \rightarrow \oplus_{i \in I} M_{i}$ such that $\iota_{i}=f \circ j_{i}$ for all $i$.

Reversing the roles of $D$ and $\oplus_{i \in I} M_{i}$ one obtains a unique $g: \oplus_{i \in I} M_{i} \rightarrow D$ with $j_{i}=g \circ \iota_{i}$. Combining the two conclusions yields $\iota_{i}=(f \circ g) \circ \iota_{i}$ which, by applying the 'universal property' of $\oplus_{i} M_{i}$ and the maps $\iota_{i}$, shows that $f \circ g$ is the identity map on $\oplus_{i} M_{i}$.

Finally, we also have $j_{i}=(g \circ f) \circ j_{i}$, which by the uniqueness in the property we assume for $D$ implies that $g \circ f$ is the identity on $D$. So indeed $D$ and $\oplus_{i} M_{i}$ are isomorphic, finishing the proof.

## VI. 4 Cyclic modules

In this section we discuss modules $M$ over a ring $R$, such that we have $M=R e$ for some $e \in M$. It turns out that this simple concept has some surprising applications.
VI.4.1 Definition. For a ring $R$, a (left) $R$-module $M$ is called cyclic if $e \in M$ exists with $M=R e$.
VI.4.2 Example. An abelian group, regarded as a module over $\mathbb{Z}$, is cyclic as a group if and only if it is cyclic as a $\mathbb{Z}$-module. Since any cyclic group is isomorphic to $\mathbb{Z} / n \mathbb{Z}$ for some $n \geq 0$ (the case $n=0$ is exactly the case of an infinite cyclic group!), we have that the cyclic $\mathbb{Z}$-modules are, up to isomorphisms, the modules $\mathbb{Z} / n \mathbb{Z}$ with $n \geq 0$. We will generalize this example in Theorem VI.4.4 below.
VI.4.3 Example. Suppose the ring $R$ is commutative and let $I \subset R$ be an ideal. Then $I$ is an $R$-module, and it is a cyclic $R$-module if and only if the ideal $I$ is a principal ideal.
VI.4.4 Theorem. If $R$ is a commutative ring (with 1) and $M$ is a cyclic $R$-module, then $I:=\{r \in R: r m=0 \forall m \in M\}$ is an ideal in $R$ and $M \cong R / I$ as $R$-modules.

We have that $M$ is free (and in that case of rank 1) precisely when $I=(0)$.
Proof. Write $M=R e$ for some $e \in M$. The map $R \rightarrow M$ given by $r \mapsto r e$ is an $R$ module homomorphism, with kernel exactly $I$. Hence $I$ is an $R$-module, and since $I \subset R$ this means it is an ideal in $R$. As the given map is surjective, one concludes $R / I \cong M$ as desired.

If $I=(0)$ then $R \cong M$ as $R$-modules, so indeed $M$ is free (of rank 1 ). And if $M$ is free, say $M \cong \oplus_{i \in S} R$ for some nonempty set $S$, then

$$
I=\{r \in R: r M=(0)\}=\left\{r \in R: \cdot\left(r \oplus_{i \in S} R\right)=(0)\right\}=(0) .
$$

This finishes the proof.
VI.4.5 Example. Let $K$ be a field and take $V$ a finite dimensional vectorspace over $K$. Suppose $\varphi: V \rightarrow V$ is a $K$-linear map. Evaluating at $\varphi$ defines a ring homomorphism from the polynomial ring $K[X]$ to the ring $\operatorname{End}(V)$ of all $K$-linear maps $V \rightarrow V$. The image of this map is denoted $R:=K[\varphi]$; it is a commutative subring of $\operatorname{End}(V)$ and it consists of all linear maps $\sum_{j=0}^{n} a_{j} \varphi^{j}$, for $n \in \mathbb{Z}_{\geq 0}$ and all $a_{j} \in K$, here $\varphi^{0}$ denotes the identity map on $V$.

The vectorspace $V$ obtains the structure of a $K[X]$-module by defining

$$
\left(\sum_{j=0}^{n} a_{j} X^{j}\right)(v)=\sum_{j=0}^{n} a_{j} \varphi^{j}(v) .
$$

Given any $v \in V$ we have the cyclic submodule $K[X] v \subset V$. By construction, it is the $K$-linear subspace of $V$ spanned by $v, \varphi(v), \varphi^{2}(v), \ldots, \varphi^{n}(v), \ldots$.

The vectorspace (cyclic $K[X]$-module) $K[X] v$ is called the Krylov subspace (corresponding to $\varphi$ and $v$ ), named after the Russian naval engineer and mathematician Aleksey Nikolaevich Krylov (1863-1945).

Using Theorem VI.4.4 we see that $K[X] v \cong K[X] / I$ with $I \in K[X]$ the ideal consisting of all $\sum_{j=0}^{n} a_{j} X^{j}$ such that $\sum_{j=0}^{n} a_{j} \varphi^{j}(v)=0$. As $K[X]$ is a principal ideal domain, $I=(g)$ for some polynomial $g$. Clearly $g \neq 0$ (the vectorspace $V$ was assumed to be finite dimensional, so there exist certainly linear dependencies between the vectors $\varphi^{j}(v)!$ ). Hence we may and will assume that $g$ is monic. Certainly the kernel of the evaluation map $K[X] \rightarrow K[\varphi]$ is contained in $I=(g)$. The latter kernel is the principal ideal generated by the minimal polynomial $m_{\varphi}$ of $\varphi$. So $\left(m_{\varphi}\right) \subset(g)$, which means that $g$ is a divisor of $m_{\varphi}$.

In particular

$$
K[x] v \cong K[x] /(g)
$$

and $\operatorname{dim}_{K}(K[x] v)=\operatorname{deg}(g) \leq \operatorname{deg}\left(m_{\varphi}\right)$. So for $V$ to be cyclic as a $K[X]$-module, a necessary condition is that $\operatorname{deg}\left(m_{\varphi}\right)=\operatorname{dim}_{K}(V)$. TheoremVI.5.7 will show that this condition also suffices.
VI.4.6 Example. (here I intend to discuss the proof of N. Katz for the fact that a $K(t)\left\langle\frac{d}{d t}\right\rangle$-module that is finite dimensional as $K(t)$-module, is cyclic. As a consequence, given an $n \times n$ system of first order linear differential equations we can find a corresponding scalar linear differential equation of order $n$...)

## VI. 5 An upper triangular form for matrices

Let $K$ be a field and let $V$ be a finite dimensional vectorspace over $K$, and let $\alpha: V \rightarrow V$ be a $K$-linear map. Under the condition that all eigenvalues of $\alpha$ are in $K$, we will construct a basis of $V$ such that the matrix $A=\left(a_{i j}\right)$ of $\alpha$ with respect to this basis is upper triangular (this means $a_{i j}=0$ for $i>j$ ).

With $\alpha$ as above, as in Example VI.4.5let

$$
K[\alpha]:=\operatorname{ev}_{\alpha}(K[X])=\left\{\sum_{i<\infty} a_{i} \alpha^{i}: a_{i} \in K\right\}
$$

be the image of the evaluation homomorphism

$$
\mathrm{ev}_{\alpha}: K[X] \longrightarrow \operatorname{End}_{K}(V), \quad f \mapsto f(\alpha)
$$

We have that $K[\alpha] \cong K[X] /\left(m_{\alpha}\right)$ with $m_{\alpha} \in K[X]$ the minimal polynomial of $\alpha$. Note that $m_{\alpha}$ is a divisor of the characteristic polynomial $P_{\alpha}=\operatorname{det}(\alpha-X I)$, and every irreducible factor of $P_{\alpha}$ is also an irreducible factor of $m_{\alpha}$. Nevertheless, there are many examples with $\operatorname{deg}\left(m_{\alpha}\right)<\operatorname{deg}\left(P_{\alpha}\right)$.

The vectorspace $V$ is a $K[\alpha]$-module with

$$
K[\alpha] \times V \longrightarrow V, \quad\left(\sum_{i=0}^{n} a_{i} \alpha^{i}, v\right) \mapsto \sum_{i=0}^{n} a_{i} \alpha^{i}(v)
$$

Since $K \subset K[\alpha]$, any $K[\alpha]$-submodule of $V$ is also a $K$-vectorspace. We will write $V$ as a direct sum of $K[\alpha]$-submodules of $V$. In each of these submodules the action of $\alpha$ will have a simple description.

The principal ideal domain $K[X]$ is a unique factorization domain, hence we can write

$$
m_{\alpha}=h_{1}^{n_{1}} h_{2}^{n_{2}} \ldots h_{k}^{n_{k}}
$$

with all $h_{i}$ monic and irreducible in $K[X]$ and $h_{i} \neq h_{j}$ for $i \neq j$. This factorization is unique up to permuting the indices $1,2, \ldots k$.

We will use this factorization to construct a direct sum decomposition of the $K[\alpha]$-module $V$.
VI.5.1 Theorem. Let $K$ be a field and let $V \neq(0)$ be a finite dimensional vectorspace over $K$ en let

$$
\alpha: V \longrightarrow V
$$

be a linear map, with minimal polynomial

$$
m_{\alpha}=h_{1}^{n_{1}} h_{2}^{n_{2}} \ldots h_{k}^{n_{k}}
$$

for irreducible monic $h_{i} \in K[X]$ and $h_{i} \neq h_{j}$ for $i \neq j$.
Put

$$
V_{i}:=\left\{v \in V: h_{i}^{n_{i}}(\alpha) v=0\right\} .
$$

Then each $V_{i}$ is a $K[\alpha]$-module. Moreover $V_{i} \neq 0$ for $i=1,2, \ldots, k$ and

$$
V \cong \bigoplus_{i \in\{1,2, \ldots, k\}} V_{i} .
$$

Proof. The fact that each $V_{i}$ is a $K[\alpha]$-module is immediate from the definitions. For the remaining assertions we use mathematical induction with respect to $k$. If $k=1$ then $h_{1}^{n_{1}}(\alpha)=m_{\alpha}(\alpha)=0$ hence $V=V_{1}$, from which the result follows for $k=1$.

Now let $k>1$. Define

$$
h:=h_{k}^{n_{k}}, \quad f:=h_{1}^{n_{1}} h_{2}^{n_{2}} \ldots h_{k-1}^{n_{k-1}}, \quad \text { and } \quad N:=\operatorname{Ker}(f(\alpha): V \rightarrow V) .
$$

Since $h$ and $f$ are coprime, there exist $g_{1}, g_{2} \in K[X]$ with

$$
g_{1} f+g_{2} h=1, \quad \text { so in particular } \quad g_{1}(\alpha) f(\alpha)+g_{2}(\alpha) h(\alpha)=1 .
$$

We claim that the $K[\alpha]$-module homomorphism

$$
\psi: N \oplus V_{k} \longrightarrow V \quad\left(n, v_{k}\right) \mapsto n+v_{k}
$$

is a $K[\alpha]$-module isomorphism.
Indeed, every $v \in V$ can be written as

$$
v=1 v=f(\alpha) g_{1}(\alpha) v+h(\alpha) g_{2}(\alpha) v .
$$

Now put $v_{k}=f(\alpha) g_{1}(\alpha) v$ and $n=h(\alpha) g_{2}(\alpha) v$, so $v=n+v_{k}$. As $m_{\alpha}(\alpha)=f(\alpha) h(\alpha)=0$ one finds

$$
\begin{aligned}
h(\alpha) v_{k} & =m_{\alpha}(\alpha) g_{1}(\alpha) v=0 & & \text { so } v_{k} \in V_{k}, \text { and } \\
f(\alpha) n & =m_{\alpha}(\alpha) g_{2}(\alpha) v=0 & & \text { so } n \in N .
\end{aligned}
$$

This shows that $\psi$ is surjective. Moreover for $n \in N$ en $v_{k} \in V_{k}$ we have

$$
n+v_{k}=0 \Longrightarrow x:=n=-v_{k} \in N \cap V_{k} .
$$

Therefore $f(\alpha) x=0$ as well as $h(\alpha) x=0$. From $x=1 x$ and $1=f(\alpha) g_{1}(\alpha)+h(\alpha) g_{2}(\alpha)$ now follows that $x=0$. $\operatorname{So} \operatorname{Ker}(\psi)=(0)$ which shows $\psi$ is injective, and we conclude that $V \cong N \oplus V_{k}$.

Let $\beta$ be the restriction of $\alpha$ to $N$. Then it is easy to verify that $\beta: N \rightarrow N$. We have (compare Exercise 9 on page 58 ) that $m_{\beta}=f$. The induction hypothesis therefore implies

$$
N \cong \oplus_{i=1}^{k-1} V_{i}, \quad \text { hence we conclude } \quad V \cong N \oplus V_{k} \cong V_{1} \oplus V_{2} \oplus \ldots V_{k}
$$

Finally $V_{i} \neq\{0\}$ since otherwise $h_{i}(\alpha)^{n_{i}}: V \rightarrow V$ would be injective, and therefore (since $\left.\operatorname{dim}_{K}(V)<\infty\right)$ invertible. Then $m_{\alpha}(\alpha)=0$ implies $m_{\alpha}(\alpha) h_{i}(\alpha)^{-n_{i}}=0$, hence $m_{\alpha} h_{i}^{-n_{i}}$ is a polynomial in the kernel of the evaluation homomorphism $\mathrm{ev}_{\alpha}$, and it has lower degree than $m_{\alpha}$. This contradiction finishes the proof.
VI.5.2 Remark. The $V_{i}$ in Theorem VI.5.1 are called the generalized eigenspaces of the linear map $\alpha$.
VI.5.3 Example. Suppose $\alpha: K^{n} \rightarrow K^{n}$ is a linear map having $n$ pairwise distinct eigenvalues $\lambda_{1}, \ldots, \lambda_{n} \in K$. In this case

$$
m_{\alpha}=\left(X-\lambda_{1}\right) \ldots\left(X-\lambda_{n}\right)
$$

and

$$
V_{i}=\operatorname{Ker}\left(\alpha-\lambda_{i}\right), \quad \text { so } \quad V_{i}:=\left\{v \in V: \alpha v=\lambda_{i} v\right\},
$$

which means $V_{i}$ is the eigenspace of $\alpha$ at the eigenvalue $\lambda_{i}$.
As each $V_{i} \neq\{0\}$ and $\sum_{i} \operatorname{dim}_{K} V_{i}=n$, every $V_{i}$ has dimension one. Choose for $i=1,2, \ldots, n$ an $f_{i} \in V_{i}-\{0\}$, so $f_{i}$ is an eigenvector of $\alpha$ with eigenvalue $\lambda_{i}$. Then $V_{i}=K \cdot f_{i}$ and

$$
V=\oplus_{i=1}^{n} V_{i}=K \cdot f_{1} \oplus \ldots \oplus K \cdot f_{n}
$$

implies that the $f_{i}$ are independent over $K$. Hence the $f_{i}$ form a basis of $K^{n}$. Since $\alpha f_{i}=\lambda_{i} f_{i}$ for every $i$, the matrix of $\alpha$ with respect to this basis is the diagonal matrix $\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}\right)$.

In the remainder of this section the special case where

$$
m_{\alpha}=\left(X-\lambda_{1}\right)^{n_{1}}\left(X-\lambda_{2}\right)^{n_{2}} \ldots\left(X-\lambda_{k}\right)^{n_{k}},
$$

with the $\lambda_{i} \in K$ pairwise distinct, is considered.
In particular the $\lambda_{i}$ are exactly the eigenvalues of $\alpha$. By Theorem VI.5.1 we have $V=\oplus_{i=1}^{k} V_{\lambda_{i}}$ as $K[\alpha]$-modules, with

$$
V_{\lambda_{i}}:=\operatorname{Ker}\left(\left(\alpha-\lambda_{i}\right)^{n_{i}}\right) \neq\{0\} .
$$

The fact that $V_{\lambda_{i}}$ is a $K[\lambda]$-module implies in particular that if $v_{i} \in V_{\lambda_{i}}$ then also $\alpha\left(v_{i}\right) \in V_{\lambda_{i}}$. This means that a matrix of $\alpha$ with respect to a suitable basis will consist of blocks:

$$
V=\oplus_{i=1}^{k} V_{\lambda_{i}} \xrightarrow{\alpha} \oplus_{i=1}^{k} V_{\lambda_{i}}, \quad \alpha\left(V_{\lambda_{i}}\right) \subset V_{\lambda_{i}} .
$$

We consider each of these blocks separately, in other words we restrict $\alpha$ to $V_{\lambda_{i}}$. Note that the restriction of $\alpha$ to $V_{\lambda_{i}}$ has minimal polynomial ( $\left.X-\lambda_{i}\right)^{n_{i}}$.
VI.5.4 Example. Take $V=\mathbb{R}^{3}$ and

$$
A:=\left(\begin{array}{lll}
4 & -4 & 4 \\
1 & -1 & 4 \\
0 & -1 & 4
\end{array}\right)
$$

We determine the minimal polynomial of $A$ and we find a basis of $V$ on which $A$ is given by blocks. First compute the characteristic polynomial of $A$, so

$$
\operatorname{det}(A-X I)=-\left(X^{3}-7 X^{2}+16 X-12\right)=-(X-2)^{2}(X-3) .
$$

The minimal polynomial of $A$ is therefore either $(X-2)(X-3)$ or $(X-2)^{2}(X-3)$. One checks that $(A-2)(A-3) \neq 0$, and hence

$$
m_{A}=(X-2)^{2}(X-3)
$$

The generalized eigenspaces are the kernels of

$$
(A-2)^{2}=\left(\begin{array}{rrr}
0 & 0 & 0 \\
-1 & 1 & 0 \\
-1 & 1 & 0
\end{array}\right) \text { and of } A-3=\left(\begin{array}{lll}
1 & -4 & 4 \\
1 & -4 & 4 \\
0 & -1 & 1
\end{array}\right) .
$$

Hence (writing $V_{\lambda}$ for the generalized eigenspace at the eigenvalue $\lambda$ ) one computes

$$
V_{2}=K \cdot(1,1,0)+K \cdot(1,1,1), \quad \text { and } \quad V_{3}=K \cdot(0,1,1) .
$$

Since

$$
\left.\begin{array}{l}
A(1,1,0)=(0,0,-1)=1 \cdot(1,1,0)+(-1) \cdot(1,1,1) \\
A(1,1,1)=(4,4,3)=1 \cdot(1,1,0)+3 \cdot(1,1,1)
\end{array}\right\},
$$

we find that $\left(\begin{array}{cc}1 & 1 \\ -1 & 3\end{array}\right)$ is the matrix of the restriction of $A$ to $V_{2}$ (with respect to the basis $(1,1,0),(1,1,1)$ of $\left.V_{2}\right)$. Restricting $A$ to $V_{3}$ yields the $1 \times 1$ matrix (3). The matrix of $A$ in terms of the basis of $V$ given by $\{(1,1,0),(1,1,1),(0,1,1)\}$ is therefore

$$
\left(\begin{array}{rrr}
1 & 1 & 0 \\
-1 & 3 & 0 \\
0 & 0 & 3
\end{array}\right)
$$

VI.5.5 Theorem. Let $W$ be a finite dimensional vectorspace over a field $K$ and let $B: W \longrightarrow W$ be a linear map with minimal polynomial $m_{B}=(X-\lambda)^{m}$.

Then

$$
B=\lambda I+N, \quad \text { with } \quad N^{m}=0
$$

(here I denotes the identity map).
There exists a basis of $W$ such that the matrix of $B$ with respect to that basis has the form

$$
\left(\begin{array}{ccccc}
\lambda & * & \ldots & \ldots & * \\
0 & \lambda & * & \ldots & * \\
0 & 0 & & & * \\
0 & 0 & 0 & 0 & \lambda
\end{array}\right)
$$

Proof. The first statement follows by writing $N:=B-\lambda I$ and noting that $m_{B}(B)=0$. Next, for $i=1,2, \ldots m$ define the linear subspace $W_{i} \subset W$ by

$$
W_{1}=\operatorname{Ker}(N), \ldots, W_{i}=\operatorname{Ker}\left(N^{i}\right), \ldots, W_{m}=\operatorname{Ker}\left(N^{m}\right)=\operatorname{Ker}(0)=W
$$

Note that $W_{i} \subset W_{i+1}$, since if $N^{i} v=0$ then also $N^{i+1} v=0$. Now take a basis of $W$ starting from a basis of $W_{1}$, extending this to a basis of $W_{2}$, and continue in this way by extending a basis of $W_{i}$ to a basis of $W_{i+1}$. At last this yields a basis $\left\{f_{i}\right\}$ of $W=W_{m}$.

Note that

$$
N W_{i} \subset W_{i-1}, \quad \text { because if } \quad N^{i} v=0 \text { then } N^{i-1}(N v)=0 \text {, so } N v \in W_{i-1}
$$

Now let $f_{j}$ be any of the chosen basis vectors. If $f_{j} \in W_{1}$ then $N\left(f_{j}\right)=0$. If $f_{j} \notin W_{1}$, then $i$ exists such that

$$
f_{j} \in W_{i} \backslash W_{i-1}
$$

Since $N W_{i} \subset W_{i-1}$ we have by the construction of the basis that

$$
N\left(f_{j}\right)=\sum_{k} x_{k j} f_{k} \quad \text { with } \quad x_{k j}=0 \text { for } k \geq j .
$$

Hence the matrix of $N$ with respect to the basis $\left\{f_{k}\right\}$ is an upper triangular matrix with zeros on the diagonal. Hence $B=\lambda I+N$ is on the same basis given by the asserted form. This completes the proof.
VI.5.6 Example. Consider the $\mathbb{R}$-linear map

$$
B: \mathbb{R}^{2} \longrightarrow \mathbb{R}^{2} \quad \text { with matrix } \quad\left(\begin{array}{rr}
1 & 1 \\
-1 & 3
\end{array}\right)
$$

The minimal polynomial $m_{B}$ is a divisor of the characteristic polynomial

$$
\operatorname{det}(B-X I)=X^{2}-4 X+4=(X-2)^{2}
$$

of $B$. Since $B-2 I \neq 0$, it follows that

$$
m_{B}=(X-2)^{2}, \quad \text { and } \quad N:=B-2 I=\left(\begin{array}{cc}
-1 & 1 \\
-1 & 1
\end{array}\right) .
$$

Check, as predicted by TheoremVI.5.5 that indeed $N^{2}=0$. We determine the $W_{i}$ 's:

$$
W_{1}=\operatorname{Ker}(N)=\mathbb{R} \cdot(1,1), \quad W_{2}=\mathbb{R}^{2} .
$$

As basis of $\mathbb{R}^{2}$ we take $f_{1}:=(1,1) \in W_{1}$, supplemented with $f_{2}:=(0,1) \in W_{2}$. Then

$$
N f_{1}=0, \quad N f_{2}=\binom{1}{1}=f_{1} .
$$

Therefore $\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right)$ is the matrix of $N$ with respect to the basis $f_{1}, f_{2}$. On the same basis the matrix van $B$ has the upper triangular form

$$
\left(\begin{array}{ll}
2 & 1 \\
0 & 2
\end{array}\right) .
$$

The arguments used in this section also make it possible to give a criterion for a $K$-vectorspace to be cyclic as a $K[\varphi]$-module, where $\varphi: V \rightarrow V$ is a linear map. This extends the discussion we started in ExampleVI.4.5. The result is the following.
VI.5.7 Theorem. Let $V \neq(0)$ be a finite dimensional vectorspace over a field $K$ and let $\varphi: V \rightarrow V$ be a linear map with minimal polynomial $m_{\varphi}$.

Then $V$ is cyclic as a $K[\varphi]$-module if and only if $\operatorname{deg}\left(m_{\varphi}\right)=\operatorname{dim}_{K}(V)$.
Proof. If $V=K[\varphi] \cdot v$ for some $v \in V$, then (as also discussed in Example VI.4.5) by Theorem VI.4.4 we have $V \cong K[\varphi] / I$ as $K[\varphi]$-modules, for some ideal $I \subset K[\varphi]$. Since as rings $K[\varphi] \cong K[X] /\left(m_{\varphi}\right)$, the ideal $I$ corresponds to an ideal $J \subset K[X]$ with $\left(m_{\varphi}\right) \subset J$ and $K[X] / J \cong K[\varphi] / I \cong V$. As $K[X]$ is a principal ideal domain, $J=(g)$ for some $g \in K[X]$ and the property $\left(m_{\varphi}\right) \subset J$ means that $g \mid m_{\varphi}$. Now

$$
\operatorname{dim}_{K}(V)=\operatorname{dim}_{K}(K[X] /(g))=\operatorname{deg}(g)
$$

implies that $\operatorname{deg}\left(m_{\varphi}\right) \geq \operatorname{dim}_{K}(V)$. Since $m_{\varphi}$ divides the characteristic polynomial of $\varphi$ which is of degree $\operatorname{dim}_{K}(V)$, we also have $\operatorname{deg}\left(m_{\varphi}\right) \leq \operatorname{dim}_{K}(V)$, hence it follows that $\operatorname{deg}\left(m_{\varphi}\right)=\operatorname{dim}_{K}(V)$.

For the converse, assume that $\operatorname{deg}\left(m_{\varphi}\right)=\operatorname{dim}_{K}(V)$. Factor

$$
m_{\varphi}=h_{1}^{n_{1}} h_{2}^{n_{2}} \cdot \ldots \cdot h_{k}^{n_{k}}
$$

for irreductible monic and pairwise distinct $h_{i} \in K[X]$, and all $n_{i}>0$. Then by TheoremVI.5.1 $V \cong \oplus_{i} V_{i}$ as $K[\varphi]$-modules, with $V_{i}=\operatorname{Ker}\left(h_{i}^{n_{i}}(\varphi)\right)$.

From the proof of VI.5.1 we know that the restriction of $\varphi$ to $V_{i}$ has minimal polynomial $h_{i}^{n_{i}}$. Hence $v_{i} \in V_{i}$ exists with $h_{i}^{n_{i}-1}\left(v_{i}\right) \neq 0$. We claim that

$$
v:=v_{1}+v_{2}+\ldots+v_{k} \in V
$$

has the property $K[\varphi] \cdot v=V$. Indeed, suppose $f(\varphi)(v)=0$ for some $f \in K[X]$. Then $f(\varphi)\left(v_{1}\right)=-f(\varphi)\left(v_{2}+\ldots+v_{k}\right)$. As the left-hand side is in $V_{1}$ and the right-hand side in $V_{2} \oplus \ldots \oplus V_{k}$, the fact that $V$ is the direct sum of the subspaces $V_{i}$ implies $f(\varphi)\left(v_{1}\right)=0$ and $f(\varphi)\left(v_{2}+\ldots+v_{k}\right)=0$. Continuing inductively we find $f(\varphi)\left(v_{i}\right)=0$ for all $i$. This implies that $h_{i}^{n_{i}} \mid f$, and therefore $m_{\varphi} \mid f$.

As $m_{\varphi}(\varphi)=0$ one concludes that the ideal $I \subset K[\varphi]$ corresponding via TheoremVI.4.4 to the cyclic module $K[\varphi] \cdot v$, is the ideal (0). Hence $K[\varphi] \cdot v \cong K[\varphi]$. The latter $K$-vectorspace has dimension $\operatorname{dim}_{K}(V)$ by assumption, so we conclude that $K[\varphi] \cdot v \subset V$ has dimension $\operatorname{dim}_{K}(V)$ as well. This implies $K[\varphi] \cdot v=V$ hence $V$ is cyclic as a $K[\varphi]$-module, finishing the proof.

## V1. 6 Exercises

1. Take $R=\mathbb{R}[X]$. For $a \in \mathbb{R}$ define

$$
I_{a}:=(X-a)=\mathbb{R}[X] \cdot(X-a)
$$

which is an ideal of $R$ (and therefore an $R$-module).
(a) Another $R$-module is the field of rational functions in the variable $X$ over $R$. For $a, b \in \mathbb{R}$ define

$$
\phi: I_{a} \longrightarrow \mathbb{R}(X) \quad \text { by } \quad f \mapsto f \cdot \frac{X-b}{X-a}
$$

Verify that $\phi$ is an injective $R$-module homomorphism and that $\operatorname{Im}(\phi)=I_{b}$. Conclude that $I_{a}$ and $I_{b}$ are isomorphic as $R$-modules.
(b) Show that the $R$-modules $\mathbb{R}_{a}:=R / I_{a}$ and $\mathbb{R}_{b}:=R / I_{b}$ are not isomorphic in case $a \neq b$.
2. Find using the methods of this chapter an upper triangular form for the following matrices:

$$
A=\left(\begin{array}{lll}
3 & -1 & 0 \\
4 & -2 & 1 \\
4 & -4 & 3
\end{array}\right), \quad B=\left(\begin{array}{rrr}
1 & 0 & 0 \\
-1 & -2 & 3 \\
-1 & -3 & 4
\end{array}\right)
$$

3 . Let $R$ be the ring of $C^{\infty}$ functions $\mathbb{R} \rightarrow \mathbb{R}$ with period $2 \pi$, and

$$
M:=\left\{g \in C^{\infty}(\mathbb{R}): g(x+2 \pi)=-g(x)\right\}, \quad \text { and } \quad I:=\{f \in R: f(0)=0\}
$$

Then both $M$ and $I$ are $R$-modules (for $M$ see Example VI.3.6, and $I \subset R$ is an ideal).
(a) Show that

$$
\phi: M \longrightarrow I, \quad g \mapsto g S_{0}
$$

with $S_{0}(x):=\sin \frac{x}{2}$, is an injective $R$-module homomorphism.
(b) Take $f \in I$. Show that the function $g$ defined by

$$
g(x):= \begin{cases}f(x) / S_{0}(x), & \text { if } \quad x \not \equiv 0 \bmod 2 \pi \\ 2 f^{\prime}(x) & \text { if } \quad x \equiv 0 \bmod 2 \pi\end{cases}
$$

is a $C^{\infty}$ function. Hint: check that $g \in M$ using

$$
\left.f(x)=\int_{0}^{1} \frac{\partial f}{\partial t}(t x) d t=x \int_{0}^{1} f^{\prime}(t x) d t\right)
$$

(c) Prove that the $R$-modules $M$ and $I$ are isomorphic.
4. Let $R$ be a commutative ring and let $I, J \subset R$ be ideals such that

$$
I+J=R \quad \text { and let } \quad i_{1} \in I, j_{1} \in J \quad \text { satisfy } \quad i_{1}+j_{1}=1
$$

(a) Show that

$$
\phi: I \oplus J \longrightarrow R, \quad(i, j) \mapsto i+j
$$

is a surjective $R$-module homomorphism with $\operatorname{kernel} \operatorname{Ker}(\phi)$ as an $R$-module isomorphic to $I \cap J=I J$.
(b) Show that

$$
\psi: I \oplus J \longrightarrow R \oplus I J, \quad(i, j) \mapsto\left(i+j, i j_{1}-j i_{1}\right)
$$

is an $R$-module isomorphism.
5. Let $R$ be the ring of polynomial functions on the circle:

$$
R=\mathbb{R}[X, Y] /\left(X^{2}+Y^{2}-1\right)
$$

and define

$$
x:=X+\left(X^{2}+Y^{2}-1\right), y:=Y+\left(X^{2}+Y^{2}-1\right) \in R .
$$

Since $X^{2}+Y^{2}-1$ is irreducible (Eisenstein polynomial at $p=Y-1$ in the unique factorization domain $(\mathbb{R}[X])[Y])$, the ideal $\left(X^{2}+Y^{2}-1\right) \subset \mathbb{R}[X, Y]$ is a prime ideal and hence $R$ is an integral domain. The field of fractions of $R$ we denote by $Q(R)$. In $R$ we define the ideals

$$
I:=(x-1, y) \quad \text { and } \quad J:=(x, y-1),
$$

it can be verified that $I$ is not a principal ideal (this is an exercise in the Algebraic Structures course).
(a) Prove that in $R$ it holds that

$$
(x+y-1)^{2}=2(x-1)(y-1), \quad(x+y-1)(x-y+1)=-2 y(y-1)
$$

(b) Define

$$
\psi: I \longrightarrow Q(R), \quad i \mapsto i \cdot \frac{2(y-1)}{x+y-1}
$$

Verify that $\psi$ is an $R$-module homomorphism, that $\psi$ is injective, and that $\operatorname{Im}(\psi)=J \subset R \subset Q(R)$.
(c) Prove that $I+J=R$ and that $I J=(x+y-1) R \cong R$ (isomorphic as $R$ modules).
(d) Conclude (using Exercise 4 that $I \neq R$, but $I \oplus I \cong R^{2}$.
6. Take $R=\mathbb{Z}[\sqrt{-5}]$ and $I:=(2,1+\sqrt{-5})$ and $J:=(3,1-\sqrt{-5})$ (so $I$ and $J$ are ideals in $R$ ).
(a) Verify that

$$
\psi: I \longrightarrow \mathbb{Q}[\sqrt{-5}], \quad i \mapsto i \cdot \frac{3}{1+\sqrt{-5}}
$$

is an injective $R$-module homomorphism with $\operatorname{Im}(\psi)=J \subset R$.
(b) Show that $I+J=R$ and $I J=(1-\sqrt{-5}) R$.
(c) Show using Exercise 4 that $I \nexists R$ but $I \oplus I \cong R^{2}$ (for a different argument compare Example VI.3.5.
7. This exercise provides an example of a ring $R$ such that $R \cong R^{2}$ as $R$-modules. Let $R$ be the ring of row-finite matrices with coefficients in the field $K$. So any $r \in R$ is an infinite matrix $r=\left(r_{i j}\right)_{i, j \in \mathbb{Z}_{\geq 1}}$, and for every $i$ we have $r_{i j} \neq 0$ for only finitely many $j$.
Addition and multiplication are analogous to the usual matrix operations:

$$
r+s=t \quad \text { with } t_{i j}:=r_{i j}+s_{i j}, \quad \text { and } \quad r s=u \quad \text { with } u_{i j}:=\sum_{k=0}^{\infty} r_{i k} s_{k j}
$$

(note that the latter sum is a finite sum since for any $i$ only finitely many $r_{i k}$ 's are nonzero).
(a) Verify that $R$ is a ring.
(b) Define $b, c \in R$ by

$$
\begin{aligned}
b & :=\left(b_{i j}\right) \quad \text { and } \quad b_{i j}=1 \text { for } j=2(i-1)+1 \text { and } b_{i j}=0 \text { otherwise }, \\
c & :=\left(c_{i j}\right), \quad \text { and } \quad c_{i j}=1 \text { for } j=2(i-1)+2 \text { and } c_{i j}=0 \text { otherwise } .
\end{aligned}
$$

Prove that

$$
R \cong R b \oplus R c \cong R \oplus R
$$

as $R$-modules.
8. An $R$-module $M \neq 0$ is called simple if $\{0\}$ and $M$ are the only $R$-submodules of $M$. Suppose $M$ is a simple $R$-module.
Prove that for every $R$-module homomorphism $f: M \rightarrow M$ either $f=0$ or $f$ is an $R$-module isomorphism. Conclude that $\operatorname{End}_{R}(M)$ is a division ring (so every nonzero element is a unit).
Determine $\operatorname{End}_{R}(M)$ in the case $R=M(n, K)$ and $M=K^{n}$, where $K$ is a field (see ExampleVI.1.3).
9. Given are two finite dimensional vectorspaces $V, W$ over a field $K$. Let $\alpha: V \rightarrow V$ and $\beta: W \rightarrow W$ be linear maps over $K$, with minimal polynomials $m_{\alpha}$ and $m_{\beta}$. Define

$$
\alpha \oplus \beta: V \oplus W \rightarrow V \oplus W \quad \text { by } \quad(v, w) \mapsto(\alpha(v), \beta(w)) .
$$

Prove that if $\operatorname{gcd}\left(m_{\alpha}, m_{\beta}\right)=1$ then $m_{\alpha \oplus \beta}=m_{\alpha} m_{\beta}$.

# VII QUOTIENTS, EXACTNESS, TENSOR PRODUCTS, AND PROJECTIVE MODULES 

## VII. 1 Quotients of modules

Recall that any subgroup $N$ of an abelian group $M$ is a normal subgroup, so the factor group $M / N$ exists. If moreover $M$ is an $R$-module for some ring $R$, and $N \subset M$ is an $R$-submodule, then a natural structure of $R$-module exists on the factor group $M / N$. This is explained below.
VII.1.1 Definition. Let $R$ be a ring and let $N$ be a submodule of an $R$-module $M$. Then

$$
M / N=\{\bar{m}=m+N \subset M: m \in M\}
$$

obtains the structure of an $R$-module by defining

$$
R \times M / N \longrightarrow M / N, \quad(r, m+N) \mapsto r m+N .
$$

This $R$-module is simply denoted $M / N$ and it is called the quotient of $M$ by $N$.
VII.1.2 Remark. Of course one needs to check that the map $R \times M / N \rightarrow M / N$ used here is well-defined, in other words:

$$
\text { if } \quad m_{1}+N=m_{2}+N \quad \text { then } \quad r m_{1}+N=r m_{2}+N
$$

This property indeed holds, since

$$
m_{1}+N=m_{2}+N \Longrightarrow m_{1}-m_{2} \in N, \quad \text { and } \quad r N \subset N
$$

(because $N$ is an $R$-module). Now $m_{1}-m_{2} \in N$ implies $r\left(m_{1}-m_{2}\right)=r m_{1}-r m_{2} \in N$, which exactly means $r m_{1}+N=r m_{2}+N$. So $m_{1}+N=m_{2}+N \Rightarrow r m_{1}+N=r m_{2}+N$.

One easily verifies that in this way the abelian group obtains the structure of an $R$-module, as asserted in the definition.

## VII.1.3 Examples.

1. Let $R=K$ be a field and let $V \subset K^{n}$ be a linear subspace of $K^{n}$. Then the quotient $K^{n} / V$ is a vectorspace over $K$. To make this more explicit, take a basis $f_{1}, \ldots, f_{n}$ of $K^{n}$ by extending a basis $f_{1}, f_{2}, \ldots, f_{k}$ of $V$. So

$$
V=\left\{\sum_{i=1}^{n} x_{i} f_{i} \in K^{n}: x_{i}=0 \text { for all } i>k\right\} .
$$

Define

$$
W:=\sum_{j=k+1}^{n} K \cdot f_{j} .
$$

Then $K^{n}=V \oplus W$, which means every $x \in K^{n}$ can be written in a unique way as $x=v+w$ with $v \in V, w \in W$. In particular $x \in V$ if and only if $w=0$. So if $x_{1}=v_{1}+w_{1}$ and $x_{2}=v_{2}+w_{2}$ for some $v_{i} \in V$ and $w_{i} \in W$ then

$$
x_{1}+V=x_{2}+V \Longleftrightarrow x_{1}-x_{2} \in V \Longleftrightarrow w_{1}=w_{2}
$$

As a consequence

$$
K^{n} / V=\left\{w+V \subset K^{n}: w \in W\right\}
$$

and

$$
w_{1}+V=w_{2}+V \Leftrightarrow w_{1}=w_{2}
$$

This means that any residue class $x+V$, by writing $x=v+w$ with $v \in V$ and $w \in W$, can be written in a unique way as $x+V=w+V$ with $w \in W$. In this way the vectorspace $K^{n} / V$ is identified with the vectorspace $W$. More precisely, the map

$$
f: W \longrightarrow K^{n} / V, \quad w \mapsto w+V
$$

is an isomorphism of vectorspaces over $K$. Note that there are many different choices for $W$, but all of them are $K$-vectorspaces isomorphic to $K^{n-k}$.
2. If $G$ is an abelian group and $H \subset G$ a subgroup, then $G / H$ is an abelian group as well and therefore it is a $\mathbb{Z}$-module. The action of $\mathbb{Z}$ we defined in Example VI.1.5 on any abelian group and hence also on $G / H$, coincides with the action of $\mathbb{Z}$ defined here on the quotient module $G / H$.
3. If $R$ is a ring and $I \subset R$ is an ideal, then $I$ is a submodule of the $R$-module $R$. The $R$-module structure on the quotient $R / I$ is the same as the one defined in ExampleVI.1.3 (verify this yourself!).
VII.1.4 Theorem. Let $f: M \rightarrow N$ be an $R$-module homomorphism. Then:
(a) The $R$-module homomorphism $f$ induces an $R$-module isomorphism

$$
M / \operatorname{Ker}(f) \xrightarrow{\cong} \operatorname{Im}(f), \quad m+\operatorname{Ker}(f) \mapsto f(m)
$$

(b) If $K \subset M$ is a submodule, then the canonical map

$$
\phi: M \longrightarrow M / K, \quad m \mapsto m+K
$$

is a (surjective) $R$-module homomorphism with kernel $\operatorname{Ker}(\phi)=K$.
Proof. This is completely analogous to similar homomorphism theorems in Group Theory and in the theory of rings. For this reason the details are left as an exercise for the reader.

## VII. 2 Hom and exactness

The following terminology for $R$-module homomorphisms is very common not only in algebra, but also in, for example, topology (in particular homology and homotopy theory) and in differential geometry.
VII.2.1 Definition. A sequence of $R$-module homomorphisms

$$
\ldots \longrightarrow M \xrightarrow{f} N \xrightarrow{g} P \longrightarrow \ldots
$$

is called exact in $N$ if

$$
\operatorname{Im}(f)=\operatorname{Ker}(g) .
$$

The sequence is called exact if it is exact in all modules appearing in it.
VII.2.2 Remark. A special case of a sequence of $R$-module homomorphisms is

$$
0 \longrightarrow N \xrightarrow{g} P .
$$

This sequence is exact in $N$ precisely when $g$ is injective: namely, the image of the map on the left is $\{0\} \subset N$, so exactness (in $N$ ) means that $\{0\}=\operatorname{Ker}(g)$, in other words, $g$ is injective.

Similarly, the sequence

$$
M \xrightarrow{f} N \longrightarrow 0
$$

is exact in $N$ precisely when $f$ is surjective. Indeed, the kernel of the map on the right is all of $N$, so exactness (in $N$ ) means $\operatorname{Im}(f)=N$.

Finally, consider the sequence

$$
0 \longrightarrow M \xrightarrow{f} N \xrightarrow{g} P \longrightarrow 0 .
$$

This sequence being exact means first of all that $f$ is injective, so $f: M \rightarrow f(M)$ is an $R$-module isomorphism. Next, $g$ is surjective and therefore $\operatorname{Ker}(g)=\operatorname{Im}(f) \cong M$, hence $g$ induces an $R$-module isomorphism $N / \operatorname{Ker}(g) \cong P$. Up to identifications (which means $R$-module isomorphisms), this means that an exact sequence as above has the form

$$
0 \longrightarrow \operatorname{Ker}(g) \hookrightarrow N \stackrel{g}{\longrightarrow} N / \operatorname{Ker}(g) \longrightarrow 0 .
$$

If a sequence

$$
\ldots \longrightarrow M \xrightarrow{f} N \xrightarrow{g} P \longrightarrow \ldots
$$

is exact, then the property $\operatorname{Ker}(g)=\operatorname{Im}(f)$ implies in particular that $g \circ f=0$. In other words, the composition of any two consecutive maps in an exact sequence, is the zero map. On the other hand, if in a sequence $M \xrightarrow{f} N \xrightarrow{g} P$ we have $g \circ f=0$ then this implies $\operatorname{Im}(f) \subset \operatorname{Ker}(g)$, which is evidently weaker than the assertion that the sequence is exact at $N$.
VII.2.3 Definition. Given a ring $R$, a diagram of $R$-modules is a simple, directed graph in which the vertices represent $R$-modules $M_{i}$, and the (directed) edges represent $R$-module homomorphisms $M_{i} \rightarrow M_{j}$.

A diagram of $R$-modules is called a commutative diagram if (for all $i, j$ ) every path from $M_{i}$ to $M_{j}$, seen as a composition of $R$-module homomorphisms, is the same map.
VII.2.4 Examples. A triangle of $R$-modules and maps

being commutative simply means that $h \circ f=g$.
Similarly, commutativity of a square

| $K$ | $\xrightarrow{e}$ | $L$ |
| :--- | :--- | :--- |
| $\downarrow f$ |  | $\downarrow g$ |
| $M$ | $\xrightarrow{h}$ | $N$ |

means that $h \circ f=g \circ e$.

A famous result featuring commutative diagrams as well as exact sequences is the so-called snake lemma. It appeared in the movie It's my turn (1980) where it was proven by actress Jill Clayburgh (1944-2010). We present it here, using the same notations as she did in the movie. To this end we first introduce one more notation.
VII.2.5 Notation. If $R$ is a ring and $f: M \rightarrow N$ is an $R$-module homomorphism, then we write

$$
\operatorname{Coker}(f):=N / f(M)
$$

for the $R$-module obtained by taking the quotient of $N$ by the submodule $f(M)$.
VII.2.6 Theorem. Suppose the diagram of $R$-modules

$$
\begin{array}{lllllll}
0 \rightarrow & A & \xrightarrow{f} & B & \xrightarrow{g} & C & \rightarrow 0 \\
& \downarrow \alpha & & \downarrow \beta & & \downarrow \gamma & \\
0 \rightarrow & A^{\prime} & \xrightarrow{f^{\prime}} & B^{\prime} & \xrightarrow{g^{\prime}} & C^{\prime} & \rightarrow 0
\end{array}
$$

is commutative, and its two rows

$$
0 \rightarrow A \xrightarrow{f} B \stackrel{g}{\rightarrow} C \rightarrow 0 \quad \text { and } \quad 0 \rightarrow A^{\prime} \xrightarrow{f^{\prime}} B^{\prime} \stackrel{g^{\prime}}{\rightarrow} C^{\prime} \rightarrow 0
$$

are exact.
Then this induces an exact sequence of $R$-modules

$$
0 \rightarrow \operatorname{Ker}(\alpha) \xrightarrow{f} \operatorname{Ker}(\beta) \xrightarrow{g} \operatorname{Ker}(\gamma) \xrightarrow{\delta} \operatorname{Coker}(\alpha) \xrightarrow{\overline{f^{\prime}}} \operatorname{Coker}(\beta) \stackrel{\overline{g^{\prime}}}{\rightarrow} \operatorname{Coker}(\gamma) \rightarrow 0
$$

Here $f$ and $g$ denote the restrictions to $\operatorname{Ker}(\alpha) \subset A$ and $\operatorname{Ker}(\beta) \subset B$ of $f: A \rightarrow B$ resp. $g: B \rightarrow C$. Similarly, $\overline{f^{\prime}}$ and $\overline{g^{\prime}}$ are the maps induced by $f^{\prime}: A^{\prime} \rightarrow B^{\prime}$ resp. $g^{\prime}: B^{\prime} \rightarrow C^{\prime}$. Finally, $\delta$ is defined as follows. Given any $c \in \operatorname{Ker}(\gamma) \subset C$, since $g: B \rightarrow C$ is surjective, $b \in B$ exists with $g(b)=c$. Then $g^{\prime}(\beta(b))=\gamma(g(b))=\gamma(c)=0$ by the commutativity of the diagram, hence $\beta(b) \in \operatorname{Ker}\left(g^{\prime}\right)=\operatorname{Im}\left(f^{\prime}\right)$, so $a^{\prime} \in A^{\prime}$ exists with $f^{\prime}\left(a^{\prime}\right)=\beta(b)$. Now put $\delta(c):=a^{\prime}+\alpha(A) \in \operatorname{Coker}(\alpha)$.

Proof. We first show that the maps in the sequence are well-defined and indeed map to the indicated $R$-modules.
For (the restriction of) $f$, if $a \in \operatorname{Ker}(\alpha)$ then $\beta(f(\alpha))=f^{\prime}(\alpha(\alpha))=f^{\prime}(0)=0$ where we used that the diagram is commutative. So indeed $f(a) \in \operatorname{Ker}(\beta)$.
The same argument works for the restriction of $g$.
In the construction of $\delta$, there are two places where potentially a choice is made. One of them is that given $\beta(b) \in \operatorname{Im}\left(f^{\prime}\right)$ we choose $a^{\prime} \in A^{\prime}$ with $f^{\prime}\left(a^{\prime}\right)=\beta(b)$. However, $f^{\prime}$ is injective (by the exactness of the lower row in the diagram), so in fact we have a unique $\alpha^{\prime}$ here. The other choice is the $b \in B$ such that $g(b)=c$. If a different choice $b_{1} \in B$ was used, then $g(b)=c=g\left(b_{1}\right)$ hence $b_{1}-b \in \operatorname{Ker}(g)=\operatorname{Im}(f)$. This means $a \in A$ exists with $b_{1}=b+f(a)$. Then $\beta\left(b_{1}\right)=\beta(b)+\beta(f(a))=\beta(b)+f^{\prime}(\alpha(a))$. Hence $\beta\left(b_{1}\right)$ is the image (under $f^{\prime}$ ) of the unique element $a^{\prime}+\alpha(a) \in A^{\prime}$. In particular the classes of $a^{\prime}$ and of $a^{\prime}+\alpha(\alpha)$ in $\operatorname{Coker}(\alpha)=A^{\prime} / f(A)$ are equal, showing that $\delta$ is well-defined.
For the map $\overline{f^{\prime}}$, first consider the composition $\pi \circ f^{\prime}: A^{\prime} \xrightarrow{f^{\prime}} B^{\prime} \xrightarrow{\pi} B^{\prime} / \beta(B)$ where $\pi$ is the canonical map. If $a^{\prime} \in \alpha(A)$, then write $\alpha^{\prime}=\alpha(\alpha)$ with $a \in A$. We have, using the commutativity of the diagram that $\pi\left(f^{\prime}(\alpha(a))\right)=\pi(\beta(f(a))=0$, since $\beta(f(a)) \in \beta(B)$. Hence $\alpha(A) \subset \operatorname{Ker}\left(\pi \circ f^{\prime}\right)$, which implies that $\overline{f^{\prime}}: \operatorname{Coker}(\alpha) \rightarrow \operatorname{Coker}(\beta)$ given by $a^{\prime}+\alpha(A) \mapsto f^{\prime}(a)+\beta(B)$ is well-defined.
The same argument shows that $\overline{g^{\prime}}$ is well-defined.

The fact that all maps in the sequence are $R$-module homomorphisms, is easily checked from the definitions. So it remains to show that

$$
0 \rightarrow \operatorname{Ker}(\alpha) \xrightarrow{f} \operatorname{Ker}(\beta) \stackrel{g}{\xrightarrow[g]{n}} \operatorname{Ker}(\gamma) \xrightarrow{\delta} \operatorname{Coker}(\alpha) \xrightarrow{\overline{f^{\prime}}} \operatorname{Coker}(\beta) \stackrel{\overline{g^{\prime}}}{\rightarrow} \operatorname{Coker}(\gamma) \rightarrow 0
$$

is exact. This is checked at each of the modules in the sequence, as follows.
$\operatorname{Ker}(\alpha)$ : since $f: A \rightarrow B$ is injective, so is its restriction to $\operatorname{Ker}(\alpha)$.
$\operatorname{Ker}(\beta)$ : the exactness of $A \xrightarrow{f} B \xrightarrow{g} C$ says $f(A)=\operatorname{Ker}(g)$. So restricting $g$ to $\operatorname{Ker}(\beta)$ one obtains as kernel $\operatorname{Ker}(\beta) \cap f(A)$. Since $g \circ f=0$ this certainly contains $f(\operatorname{Ker}(\alpha))$. On the other hand, if $b \in \operatorname{Ker}(\beta) \cap f(A)$ then $b=f(a)$ for some $a \in A$ and $\beta(f(a))=0$. As $\beta \circ f=f^{\prime} \circ \alpha$ and $f^{\prime}$ is injective, this implies $\alpha(\alpha)=0$. So $a \in \operatorname{Ker}(\alpha)$ and then one concludes $b=f(a) \in f(\operatorname{Ker}(\alpha))$. This shows exactness at $\operatorname{Ker}(\beta)$.
$\operatorname{Ker}(\gamma)$ : take $c \in \operatorname{Ker}(\gamma)$ such that $c \in \operatorname{Im}(g)$. This means we have $b \in \operatorname{Ker}(\beta)$ with $g(b)=c$. The fact that $b \in \operatorname{Ker}(\beta)$, says that $\beta(b)=0$. And $0 \in B^{\prime}$ is the image of $0 \in A^{\prime}$ under the map $f^{\prime}$, so by definition $\delta(c)=0+\alpha(A) \in \operatorname{Coker}(\alpha)$, in other words $c \in \operatorname{Ker}(\delta)$. This shows $\operatorname{Im}(g) \subset \operatorname{Ker}(\delta)$. On the other hand, is $\delta(c)=\overline{0}$ for some $c \in \operatorname{Ker}(\gamma)$, then for a $b \in B$ with $g(b)=c$ we have that $\beta(b)=f^{\prime}\left(a^{\prime}\right)$ with $a^{\prime} \in \alpha(A)$. So $a^{\prime}=\alpha(a)$ for some $a \in A$, and then $\beta(b)=f^{\prime}(\alpha(a))=\beta(f(a))$. As a result $b-f(a) \in \operatorname{Ker}(\beta)$, and we have $g(b-f(a))=g(b)-g(f(a))=c-0=c$. This shows $\operatorname{Ker}(\delta) \subset \operatorname{Im}(g)$ and finishes the proof of exactness at $\operatorname{Ker}(\gamma)$.
$\operatorname{Coker}(\alpha)$ : first, we show that $\operatorname{Im}(\delta) \subset \operatorname{Ker}\left(\overline{f^{\prime}}\right)$. So take $\overline{a^{\prime}} \in \operatorname{Coker}(\alpha)$ with $\overline{a^{\prime}}=\delta(c)$ where $c \in C$ satisfies $\gamma(c)=0$. By definition $\overline{f^{\prime}}\left(\overline{a^{\prime}}\right)=f^{\prime}\left(a^{\prime}\right)+\beta(B)$. We have to show that this is zero, in other words, that $f^{\prime}\left(a^{\prime}\right) \in \beta(B)$. To this end, pick $b \in B$ with $g(b)=c$ and $a^{\prime \prime} \in A^{\prime}$ with $f^{\prime}\left(a^{\prime \prime}\right)=\beta(b)$. By the construction of the map $\delta$, then $\overline{a^{\prime}}=\delta(c)=\overline{a^{\prime \prime}}$. Hence $a^{\prime}=a^{\prime \prime}+\alpha(a)$ for some $a \in A$. As a result, $f^{\prime}\left(a^{\prime}\right)=f^{\prime}\left(a^{\prime \prime}\right)+f^{\prime}(\alpha(a))$, and this equals $\beta(b)+\beta(f(a))=\beta(b+f(a))$. So indeed $f^{\prime}\left(a^{\prime}\right) \in \beta(B)$.
Next we show $\operatorname{Ker}\left(\overline{f^{\prime}}\right) \subset \operatorname{Im}(\delta)$. Take $\overline{a^{\prime}} \in \operatorname{Coker}(\alpha)$ with $\overline{f^{\prime}}\left(\overline{a^{\prime}}\right)=0$. By definition this means that $b \in B$ exists with $f^{\prime}\left(a^{\prime}\right)=\beta(b)$. Now put $c:=g(b) \in C$. Then $c \in \operatorname{Ker}(\gamma)$ since $\underline{\gamma(c)}=\gamma(g(b))=g^{\prime}(\beta(b))=g^{\prime}\left(f^{\prime}\left(a^{\prime}\right)\right)=0$. By the construction of $\delta$ we have $\delta(c)=\overline{a^{\prime}}$. This finished proving exactness at $\operatorname{Coker}(\alpha)$.
$\operatorname{Coker}(\beta)$ : Given $\overline{b^{\prime}} \in \operatorname{Coker}(\beta)$, if $\overline{g^{\prime}}\left(\overline{b^{\prime}}\right)=0$ then $c \in C$ exists with $g^{\prime}\left(b^{\prime}\right)=\gamma(c)$. Now write $c=g(b)$ for some $b \in B$, then $g^{\prime}\left(b^{\prime}-\beta(b)\right)=\gamma(c)-\gamma(g(b))=0$ hence $a^{\prime} \in A^{\prime}$ exists with $f^{\prime}\left(a^{\prime}\right)=b^{\prime}-\beta(b)$. This implies that $\overline{b^{\prime}} \in \operatorname{Im}\left(\overline{f^{\prime}}\right)$.
Vice versa, is $\overline{b^{\prime}} \in \operatorname{Im}\left(\overline{f^{\prime}}\right)$ then $a^{\prime} \in A^{\prime}$ and $a \in A$ exist such that $f^{\prime}\left(a^{\prime}\right)-b^{\prime}=\alpha(a)$. Hence $g^{\prime}\left(b^{\prime}\right)=g^{\prime}\left(f^{\prime}\left(a^{\prime}\right)-g^{\prime}(\alpha(a))=0-\beta(g(a)) \in \beta(B)\right.$. This means $\overline{b^{\prime}} \in \operatorname{Ker}\left(\overline{g^{\prime}}\right)$.
$\operatorname{Coker}(\gamma)$ : here we have to show that $\overline{g^{\prime}}$ is surjective. So take $\overline{c^{\prime}} \in \operatorname{Coker}(\gamma)$. Since $g^{\prime}$ is surjective, $b^{\prime} \in B^{\prime}$ exists with $g^{\prime}\left(b^{\prime}\right)=c^{\prime}$. Then $\overline{g^{\prime}}\left(\overline{b^{\prime}}\right)=\overline{c^{\prime}}$.
This finishes the proof of the snake lemma.
VII.2.7 Definition. Given a ring $R$ and (left) $R$-modules $M, N$, the set of all $R$ module homomorphisms $M \rightarrow N$ is denoted

$$
\operatorname{Hom}(M, N)=\{f: M \rightarrow N: f \text { is an } R-\text { module homomorphism }\} .
$$

If one wants to emphasize that the homomorphisms are considered over the ring $R$, then the notation $\operatorname{Hom}_{R}(M, N)$ is used.

In case $N=M$ one writes $\operatorname{End}_{R}(M):=\operatorname{Hom}_{R}(M, M)$, the $R$-module endomorphisms of $M$.

The set $\operatorname{Hom}(M, N)$ is in a natural way an abelian group, with

$$
(f+g)(m):=f(m)+g(m) \quad(f, g \in \operatorname{Hom}(M, N), m \in M) .
$$

In the case $M=N$, i.e., considering $\operatorname{End}(M)$, the product given by composing maps: $f g=f \circ g$, provides $\operatorname{End}(M)$ with a ring structure. In fact this is a subring of the ring consisting of all group homomorphisms $M \rightarrow M$.

In case $R$ is commutative, the abelian group $\operatorname{Hom}(M, N)$ obtains the structure of an $R$-module by defining

$$
(r f)(m):=r f(m), \quad(f \in \operatorname{Hom}(M, N), m \in M, r \in R)
$$

Commutativity of $R$ is essential, since we want $r f \in \operatorname{Hom}(M, N)$. This means that $r f$ should be $R$-linear. Now

$$
(r f)(s m):=r \cdot(f(s m))=(r s) \cdot f(m)
$$

and this equals $(s r) \cdot f(m)$ in case $R$ is commutative, but not necessarily otherwise. One readily checks that indeed in this way $\operatorname{Hom}(M, N)$ is an $R$-module in case $R$ is commutative.

In de remainder of this section we will assume the ring $R$ to be commutative.
(Most of the results presented here also hold for non-commutative rings, provided $\operatorname{Hom}(M, N)$ is considered as an abelian group only.)
VII.2.8 Examples. 1. Let $K$ be a field and take $M=N=K^{n}$. Then

$$
\operatorname{End}_{K}\left(K^{n}\right):=\operatorname{Hom}_{K}\left(K^{n}, K^{n}\right) \cong M(n, K),
$$

with $M(n, K)$ the $K$-module consisting of $n \times n$ matrices with coefficients in $K$, and the usual addition and scalar multiplication of matrices. Namely, any $\alpha \in \operatorname{Hom}_{K}\left(K^{n}, K^{n}\right)$ is by definition a $K$-linear map.
2 . Let $R$ be a ring and $M$ an $R$-module. Then

$$
\mathrm{ev}_{1}: \operatorname{Hom}_{R}(R, M) \stackrel{ }{\cong} M, \quad f \mapsto f(1),
$$

defines an isomorphism of $R$-modules. Injectivity of $\mathrm{ev}_{1}$ follows from $f(r)=r f(1)$. Namely, is $\mathrm{ev}_{1}(f):=f(1)=0$ then $f(r)=r f(1)=r \cdot 0=0$ for all $r \in R$, showing that $f=0$. To show that $\mathrm{ev}_{1}$ is surjective, note that for any $m \in M$ the map

$$
f_{m}: R \longrightarrow M \quad f_{m}(r):=r m
$$

defines an $R$-module homomorphism. It satisfies $\operatorname{ev}_{1}\left(f_{m}\right)=m$. It remains to show that $\mathrm{ev}_{1}$ is an $R$-module homomorphism. This is immediate:

$$
\left.\begin{array}{rll}
\mathrm{ev}_{1}(f+g) & := & (f+g)(1)
\end{array}\right):=f(1)+g(1)=\mathrm{ev}_{1}(f)+\mathrm{ev}_{1}(g), ~ 子=r \mathrm{ev}_{1}(f) .
$$

So indeed $\mathrm{ev}_{1}$ defines an isomorphism of modules.
3. For every $n \in \mathbb{Z}_{\geq 1}$ it holds that

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

Indeed, if $f: \mathbb{Z} / n \mathbb{Z} \rightarrow \mathbb{Z}$ is $\mathbb{Z}$-linear, then

$$
0=f(\overline{0})=f(\bar{n})=f(n \cdot \overline{1})=n f(\overline{1}),
$$

and since $\mathbb{Z}$ is an integral domain and $n \neq 0$, this implies $f(\overline{1})=0$. Hence $f(\bar{a})=a f(\overline{1})=a \cdot 0=0$ for all $\bar{a} \in \mathbb{Z} / n \mathbb{Z}$ which shows that $f=0$. In words: every element in $\mathbb{Z} / n \mathbb{Z}$ has finite order and its image under $f$ therefore has finite order as well. As a consequence, this image is 0 since 0 is the only element of finite order in $\mathbb{Z}$.

We will now consider $\operatorname{Hom}_{R}(A,-)$; in other words we fix the $R$-module $A$ and consider the recipe that on input any $R$-module $M$, outputs the $R$-module $\operatorname{Hom}_{R}(A, M)$. As before, all rings $R$ considered here will be assumed commutative.
VII.2.9 Definition. If $A$ is an $R$-module, and $f \in \operatorname{Hom}(M, N)$, then composing any $R$-module homomorphism $\phi: A \rightarrow M$ with $f: M \rightarrow N$ results in an $R$-module homomorphism denoted $f_{*}(\phi):=f \circ \phi: A \rightarrow N$ :


One also uses the somewhat suggestive notation $\operatorname{Hom}(A, f)$ for $f_{*}$.
Using the definitions it is not hard to check that indeed $f_{*}$ is an $R$-module homomorphism, which in this case means that

$$
f_{*}(\phi+\psi)=f_{*}(\phi)+f_{*}(\psi), \quad \text { and } \quad f_{*}(r \phi)=r f_{*}(\phi)
$$

for all $\phi, \psi \in \operatorname{Hom}(A, M)$ and all $r \in R$.
We will now consider the following problem. Suppose $f: M \rightarrow N$ is surjective. Does it follow that $f_{*}: \operatorname{Hom}(A, M) \rightarrow \operatorname{Hom}(A, N)$ is surjective as well? $R$-modules $A$ for which the answer is affirmative, so with the property

$$
f: M \longrightarrow N \quad \text { surjective } \quad \Longrightarrow \quad f_{*}: \operatorname{Hom}(A, M) \longrightarrow \operatorname{Hom}(A, N) \quad \text { surjective }
$$ will be studied in more detail in Section VII. 4

(We will see in Theorem VII.2.12 that if one replaces 'surjective' by 'injective'in the question above, then for all $R$-modules $A$ the answer is affirmative. Note also that if $f$ is not surjective, then Exercise 1 on page 74 shows that already for $A=R$ the map $f_{*}$ is also not surjective.
VII.2.10 Example. This example shows that $f_{*}$ is not necessarily surjective even if $f$ is surjective. Consider $n \in \mathbb{Z}_{\geq 2}$ and


If $\pi_{*}$ were surjective, then $\tilde{h}: \mathbb{Z} / n \mathbb{Z} \rightarrow \mathbb{Z}$ exists with $\pi \circ \tilde{h}=\mathrm{id}_{\mathbb{Z} / n \mathbb{Z}}$.
By Example VII.2.8 3 ) we know $\operatorname{Hom}(\mathbb{Z} / n \mathbb{Z}, \mathbb{Z})=0$, hence $\tilde{h}=0$. But then $\mathrm{id}_{\mathbb{Z} / n \mathbb{Z}}=\pi \circ \tilde{h}=0$, a contradiction. So $\pi_{*}$ is not surjective.
VII.2.11 Example. We now show: if $A=R$ and $f: M \rightarrow N$ is a surjective $R$-module homomorphism, then $f_{*}: \operatorname{Hom}(R, M) \rightarrow \operatorname{Hom}(R, N)$ is surjective as well.


To show this, take any $h \in \operatorname{Hom}(R, N)$. We want to find $\tilde{h} \in \operatorname{Hom}(R, M)$ such that $f_{*}(h)=f \circ \tilde{h}=h$. Since $f$ is surjective, $m \in M$ exists with $f(m)=h(1)$. Now define

$$
\tilde{h}: R \rightarrow M \quad \text { by } \quad \tilde{h}(r):=r m .
$$

Then $\tilde{h} \in \operatorname{Hom}(R, M)$ (see ExampleVII.2.8, 2)) and moreover

$$
\left(f_{*} \tilde{h}\right)(r):=f \tilde{h}(r)=f(r m)=r f(m)=r h(1)=h(r)
$$

for all $r \in R$. So indeed $f_{*} \tilde{h}=h$, showing that $f_{*}$ is surjective in this case.
VII.2.12 Theorem. Let $R$ be a commutative ring and suppose

$$
0 \longrightarrow M \xrightarrow{f} N \xrightarrow{g} P
$$

is an exact sequence of $R$-modules.
For any $R$-module $A$ it holds that

$$
0 \longrightarrow \operatorname{Hom}(A, M) \xrightarrow{f_{*}} \operatorname{Hom}(A, N) \xrightarrow{g_{*}} \operatorname{Hom}(A, P)
$$

is an exact sequence of $R$-modules.
In particular: if $f: M \rightarrow N$ is injective, then so is $f_{*}: \operatorname{Hom}(A, M) \rightarrow \operatorname{Hom}(A, N)$.
VII.2.13 Remark. One usually expresses Theorem VII.2.12 by asserting that the map $\operatorname{Hom}(A,-)$ is left-exact.

Proof. (of Theorem VII.2.12) Given are an exact sequence as above and an $R$ module $A$.

Exactness at $\operatorname{Hom}(A, M)$ : suppose $\phi \in \operatorname{Hom}(A, M)$ satisfies $\phi \neq 0$. Then $a \in A$ exists with $\phi(a) \neq 0$. As $\phi(a) \in M$ and $f: M \rightarrow N$ is injective, $\left(f_{*} \phi\right)(a):=f(\phi(a)) \neq 0$ so $f_{*} \phi \neq 0$. Hence $f_{*}$ is injective.

Exactness at $\operatorname{Hom}(A, N)$ : we must show $\operatorname{Im}\left(f_{*}\right)=\operatorname{Ker}\left(g_{*}\right)$.
' $\subset$ ': take an arbitrary $\psi \in \operatorname{Im}\left(f_{*}\right)$. We may write $\psi=f_{*}(\phi)$ for some $\phi \in \operatorname{Hom}(A, M)$. For any $a \in A$ now $\psi(a)=f(\phi(a)) \in \operatorname{Im}(f)$. By assumption $\operatorname{Im}(f)=\operatorname{Ker}(g)$, hence $g \psi(a)=0$ for all $a \in A$. This shows $g_{*}(\psi)=0$, proving that $\operatorname{Im}\left(f_{*}\right) \subset \operatorname{Ker}\left(g_{*}\right)$.
' $\supset$ ': let $\psi \in \operatorname{Hom}(A, N)$ satisfy $g_{*}(\psi)=0$. We will construct $\phi \in \operatorname{Hom}(A, M)$ such that $f_{*}(\phi)=\psi$. The condition $g_{*}(\psi)=0$ means that $g \psi(a)=0$ for all $a \in A$, i.e., $\psi(a) \in \operatorname{Ker}(g)$ for all $a \in A$. Moreover by assumption $\operatorname{Ker}(g)=\operatorname{Im}(f)$. We also have that $f$ is injective, so $f: M \stackrel{\cong}{\rightrightarrows} \operatorname{Im}(f)$. As a result, an $R$-module isomorphism

$$
M \stackrel{h}{\longleftarrow} \operatorname{Im}(f) \quad(\subset N), \quad \text { with } \quad f \circ h=i d_{\operatorname{Im}(f)}
$$

exists. Since $\operatorname{im}(\psi) \subset i m(f)$, the composition $\phi:=h \circ \psi: A \rightarrow M$ is a well defined $R$-module homomorphism. One finds for all $a \in A$ that

$$
\left(f_{*} \phi\right)(a):=f(\phi(a))=f(h(\psi(a)))=\psi(a), \quad \text { hence } \quad f_{*}(\phi)=\psi
$$

(here we used that $\psi(a) \in \operatorname{Ker}(g)=\operatorname{Im}(f)$ and $f \circ h=i d_{\operatorname{Im}(h)}$. So $\operatorname{Ker}\left(g_{*}\right) \subset \operatorname{Im}\left(f_{*}\right)$. This concludes the proof of the theorem.

## VII. 3 Tensor products

In linear algebra, and in various applications of linear algebra such as coding theory, inner products and generalizations of inner products are an important notion. Tensor products, which will be introduced in this section, provide a general framework for this. They are in particular used in representation theory and in differential geometry, and in applications of these areas in theoretical physics.
VII.3.1 Definition. If $R$ is a ring and $M, N, T$ are $R$-modules, then a map

$$
b: M \times N \longrightarrow T
$$

is called bilinear (or $R$-bilinear) if for every $m \in M$ and every $n \in N$ the maps $M \rightarrow T$ and $N \rightarrow T$ given by $x \mapsto b(x, n)$ resp. $y \mapsto b(m, y)$ are $R$-module homomorphisms.
VII.3.2 Example. Take $n \in \mathbb{Z}_{\geq 1}$. Let $R$ be a commutative ring and let $M=N=R^{n}$ be the free $R$-module of rank $n$. Then

$$
b: R^{n} \times R^{n} \rightarrow R \quad \text { given by } \quad b\left(\left(x_{1}, \ldots, x_{n}\right),\left(y_{1}, \ldots, y_{n}\right)\right)=\sum x_{j} y_{j}
$$

is bilinear.
VII.3.3 Definition. If $R$ is a ring and $M, N$ are $R$-modules, then a tensor product of $M$ and $N$ is a pair ( $T, \beta$ ) in which $T$ is an $R$-module and $\beta: M \times N \rightarrow T$ is a bilinear map, such that the following holds:
given any $R$-bilinear map $b: M \times N \rightarrow S$ for some $R$-module $S$, there exists a unique $R$-module homomorphism $f: T \rightarrow S$ such that $b=f \circ \beta$.

One can visualize this definition by means of diagrams:


The first thing we will see concerning a tensor product, is that if it exists then it is unique (in fact up to a unique(!) isomorphism). This should be compared with Theorem VI.3.8, where a similar argument for direct sums rather than tensor products is given.
VII.3.4 Theorem. Suppose $M, N$ are $R$-modules and $\left(T_{1}, \beta_{1}\right)$ and $\left(T_{2}, \beta_{2}\right)$ are tensor products of $M$ and $N$. Then there is a unique $R$-module isomorphism $f: T_{1} \rightarrow T_{2}$ such that $\beta_{2}=f \circ \beta_{1}$.

Proof. Using that $\left(T_{1}, \beta_{1}\right)$ is a tensor product of $M$ and $N$, the diagram

yields a unique $f_{1}: T_{1} \rightarrow T_{2}$ with $\beta_{2}=f_{1} \circ \beta_{1}$. Interchanging the roles of $T_{1}$ and $T_{2}$ one obtains $f_{2}: T_{2} \rightarrow T_{1}$ with $\beta_{1}=f_{2} \circ \beta_{2}$. As a result, $\beta_{2}=\left(f_{1} \circ f_{2}\right) \circ \beta_{2}$. However, starting from

$$
\begin{array}{cll}
M \times N & \xrightarrow{\beta_{2}} & T_{2} \\
\downarrow \beta_{2} & & \\
T_{2} & &
\end{array}
$$

and the assumption that ( $T_{2}, \beta_{2}$ ) is a tensor product, the unique arrow $T_{2} \rightarrow T_{2}$ making the above diagram commutative is the identity map. So $f_{1} \circ f_{2}=\mathrm{id}_{T_{2}}$. Again interchanging the roles of $T_{1}$ and $T_{2}$, one finds in the same way that $f_{2} \circ f_{1}=\operatorname{id}_{T_{1}}$. So $f_{1}$ and $f_{2}$ are isomorphisms as desired.

To see that $f_{1}$ is unique, suppose that also $f_{1}^{\prime}: T_{1} \rightarrow T_{2}$ is an isomorphism with $\beta_{2}=f_{1}^{\prime} \circ \beta_{1}$. Then $\beta_{1}=f_{1}^{\prime-1} \circ \beta_{2}$ and hence $\beta_{2}=\left(f_{1} \circ f_{1}^{\prime-1}\right) \circ \beta_{2}$. As above, since $\left(T_{2}, \beta_{2}\right)$ is a tensor product this implies $f_{1} \circ f_{1}^{\prime-1}=\mathrm{id}_{T_{2}}$ and hence $f_{1}=f_{1}^{\prime}$.
VII.3.5 Notation. The (if it exists at all) by Theorem VII.3.4 unique pair ( $T, \beta$ ) which is a tensor product of the $R$-modules $M$ and $N$, we will from now on denote by $M \otimes_{R} N$ (or simply $M \otimes N$ if it is clear from the context which ring $R$ is considered).

The bilinear map $\beta: M \times N$ will be written as $(m, n) \mapsto m \otimes n$.
It remains to show existence of a tensor product. The next result claims exactly that.
VII.3.6 Theorem. For any $R$-modules $M$ and $N$ a tensor product $(T, \beta)$ exists.

Proof. One constructs an abelian group ( $T,+, 0$ ) as follows. Start with the free abelian group $F$ with (independent) generators denoted $m \times n$ for $m \in M, n \in N$. So

$$
F \cong \bigoplus_{(m, n) \in M \times N} \mathbb{Z}
$$

Next, take the subgroup $S \subset F$ generated by all elements

$$
\begin{gathered}
(m \times n)+\left(m^{\prime} \times n\right)-\left(\left(m+m^{\prime}\right) \times n\right), \quad(m \times n)+\left(m \times n^{\prime}\right)-\left(m \times\left(n+n^{\prime}\right)\right), \\
((r m) \times n)-(m \times(r n)), \quad(0 \times n)
\end{gathered}
$$

with $m, m^{\prime}, 0 \in M$ and $n, n^{\prime} \in N$ and $r \in R$. Let $T$ be the factor group: $T:=F / S$. The class $(m \times n) \bmod S \in T$ is denoted $m \otimes n$.

One obtains a scalar multiplication by elements of $R$ on $F$ and on $S$ (and therefore on $T$ ) by defining

$$
r \cdot\left(\sum_{j} m_{j} \times n_{j}\right):=\sum_{j}\left(r \cdot m_{j}\right) \times n_{j} .
$$

This makes $T$ into an $R$-module. The map

$$
\beta: M \times N \rightarrow T \quad \text { given by } \quad \beta(m, n):=m \otimes n=(m \times n) \bmod S
$$

is by the definition of $S$ bilinear. The fact that indeed $(T, \beta)$ defines a tensor product is left as an exercise to the reader.

The construction of a tensor product as given in the proof of Theorem VII.3.6 is rarely needed in showing properties of the tensor product. Instead, one shows that a given module is the tensor product of $M$ and $N$ by showing that the module satisfies Definition VII.3.3 Then the unicity of the tensor product (Theorem VII.3.4) shows that indeed the given module is the tensor product. We illustrate this by some simple results and examples.
VII.3.7 Proposition. Let $R$ be a ring and let $M$ and $N$ be $R$-modules. Then

$$
M \otimes_{R} N \cong N \otimes_{R} M
$$

Proof. It suffices to show that $M \otimes_{R} N$ has the defining property of a tensor product of $N$ and $M$.

Firstly, $\beta: N \times M \rightarrow M \otimes_{R} N$ given by $\beta(n, m)=m \otimes n$ is bilinear. Now suppose $b: N \times M \rightarrow T$ is any bilinear map. Then $\tilde{b}: M \times N \rightarrow T$ defined as $\tilde{b}(m, n)=b(n, m)$ is bilinear as well, so because $M \otimes_{R} N$ is a tensor product of $M$ and $N$, a unique $R$-homomorphism $f: M \otimes_{R} N \rightarrow T$ exists with $\tilde{b}=f \circ \beta$. Then also $b=f \circ \beta$, and the uniqueness of an $f$ with the latter property is obvious.
VII.3.8 Proposition. Let $R$ be a ring and let $f: M_{1} \rightarrow M_{2}$ and $g: N_{1} \rightarrow N_{2}$ be two $R$-module homomorphisms. Then a unique $R$-module homomorphism

$$
M_{1} \otimes_{R} N_{1} \rightarrow M_{2} \otimes_{R} N_{2}
$$

exists with $m_{1} \otimes n_{1} \mapsto f\left(m_{1}\right) \otimes g\left(n_{1}\right)$ for all $m_{1} \in M_{1}$ and $n_{1} \in N_{1}$.
Proof. Define $b: M_{1} \times N_{1} \rightarrow M_{2} \otimes_{R} N_{2}$ by $b\left(m_{1}, n_{1}\right)=f\left(m_{1}\right) \otimes g\left(n_{1}\right)$. Then $b$ is bilinear, hence by the definition of tensor product a unique $R$-module homomorphism as desired exists.
VII.3.9 Proposition. If $R$ is a unitary ring and $M$ is an $R$-module, then $R \otimes_{R} M \cong M$.

Proof. Define $\beta: R \times M \rightarrow M$ by $\beta(r, m)=r m$. Clearly $\beta$ is bilinear. We claim that $(M, \beta)$ satisfies the definition of a tensor product for $R$ and $M$. Indeed, take any $R$ module $N$ and a bilinear map $b: R \times M \rightarrow N$. Define $f: M \rightarrow N$ by $f(m)=b(1, m)$. Then for any $m \in M$ and any $r \in R$ one has $f(\beta(r, m))=f(r m)=b(1, r m)=b(r, m)$ since $b$ is bilinear. So $f \circ \beta=b$. If also $f^{\prime} \circ \beta=b$, then taking any $m \in M$ we have $f(m)=b(1, m)=f^{\prime}(\beta(1, m))=f^{\prime}(m)$, so $f=f^{\prime}$. Hence $f$ is unique, proving the claim.
VII.3.10 Remark. The proof of Proposition VII.3.9 in fact shows that $r \otimes m \rightarrow r m$ defines an isomorphism of $R$-modules $R \otimes_{R} M \rightarrow M$. This is used in the proof of the next result.
VII.3.11 Proposition. Let $K$ be a field and let $V$, $W$ be vectorspaces over $K$ with bases $\left\{e_{i}: i \in I\right\}$ and $\left\{f_{j}: j \in J\right\}$, respectively. Then the tensor product $V \otimes_{K} W$ is a $K$-vectorspace with basis $\left\{e_{i} \otimes f_{j}: i \in I, j \in J\right\}$.

Proof. By assumption $V \cong \oplus_{i \in I} K$ and $W \cong \oplus_{j \in J} K$. Hence

$$
V \otimes_{K} W \cong\left(\oplus_{i \in I} K\right) \otimes_{K}\left(\oplus_{j \in J} K\right) .
$$

Now define

$$
b:\left(\oplus_{i \in I} K\right) \times\left(\oplus_{j \in J} K\right) \longrightarrow \oplus_{(i, j) \in I \times J} K \otimes_{K} K \cong \oplus_{(i, j) \in I \times J} K
$$

by $b\left(\left(x_{i}\right)_{i \in I},\left(y_{j}\right)_{j \in J}\right)=\left(x_{i} \otimes y_{j}\right)_{(i, j) \in I \times J}$. Since $b$ is bilinear, a (unique)

$$
f:\left(\oplus_{i \in I} K\right) \otimes_{K}\left(\oplus_{j \in J} K\right) \rightarrow \oplus_{(i, j) \in I \times J} K \otimes_{K} K
$$

exists with $b\left(\left(x_{i}\right)_{i \in I},\left(y_{j}\right)_{j \in J}\right)=f\left(\left(x_{i}\right)_{i \in I} \otimes\left(y_{j}\right)_{j \in J}\right)$. In fact $f$ defines an isomorphism. Since as a special case of PropositionVII.3.9 we have $K \otimes_{K} K \cong K$, the result follows.
VII.3.12 Remark. In the special case of Proposition VII.3.11 that $\operatorname{dim}_{K}(V)=n<\infty$ and $\operatorname{dim}_{K}(W)=m<\infty$, the tensor product $V \otimes_{K} W$ is a vectorspace of dimension $n \cdot m$.

If in this case $f: V \rightarrow V$ and $g: W \rightarrow W$ are $K$-linear maps, then Proposition VII.3.8 combines these maps into a $K$-linear map $V \otimes_{K} W \rightarrow V \otimes_{K} W$. The latter map is usually denoted $f \otimes g$. If $f$ is given by an $n \times n$ matrix and $g$ by an $m \times m$ matrix, then $f \otimes g$ is given by an $n m \times n m$ matrix. With respect to the bases as described in Proposition VII.3.11 the latter matrix is called the Kronecker product of the matrices for $f$ and $g$.

The final result we will present here concerning tensor products, shows a relation between ( $T \otimes_{R}$ ), so the map assigning to an $R$-module $N$ the $R$-module $T \otimes_{R} N$, and $\operatorname{Hom}_{R}(T,-)$. Recall that for the latter to be a map from $R$-modules to $R$-modules, we require that $R$ is commutative.

The result we are about to describe is usually referred to as the statement that ( $T \otimes_{R^{-}}$) and $\operatorname{Hom}_{R}(T,-)$ are adjoint. To understand this terminology, suppose that $V$ is a vector space over $\mathbb{R}$, equipped with an inner product $\langle\cdot, \cdot\rangle$. Then two linear maps $\varphi, \varphi^{*}: V \rightarrow V$ are called adjoint if $\langle\varphi(v), w\rangle=\left\langle v, \varphi^{*}(w)\right\rangle$ for all $v, w \in V$. Now replace "vectors in $V$ " by " $R$-modules", and "taking the inner product" by "taking $\operatorname{Hom}_{R}(-,-)$ ", and the maps $\varphi, \varphi^{*}$ by ( $T \otimes_{R}-$ ) respectively $\operatorname{Hom}_{R}(T,-)$. In this way the following result obtained the name adjointness of Hom and tensor.
VII.3.13 Theorem. Let $R$ be a commutative ring and let $T$ be an $R$-module. Then for any pair $M, N$ of $R$-modules one has

$$
\operatorname{Hom}_{R}\left(T \otimes_{R} M, N\right) \cong \operatorname{Hom}_{R}\left(M, \operatorname{Hom}_{R}(T, N)\right)
$$

as $R$-modules.
Proof. Write $\operatorname{Bilin}_{R}(T \times M, N)$ for the set of all bilinear maps $T \times M \rightarrow N$. In fact this is an $R$-module with respect to pointwise addition and multiplication by elements of $R$, so

$$
\left(b+b^{\prime}\right)(t, m):=b(t, m)+b^{\prime}(t, m) \quad \text { and } \quad(r b)(t, m)=r \cdot b(t, m)
$$

The definition of a tensor product implies that every element of $\operatorname{Bilin}_{R}(T \times M, N)$ corresponds to an element of $\operatorname{Hom}_{R}\left(T \otimes_{R} M, N\right)$, and in this way

$$
\operatorname{Bilin}_{R}(T \times M, N) \cong \operatorname{Hom}_{R}\left(T \otimes_{R} M, N\right)
$$

as $R$-modules. Hence it suffices to show that

$$
\operatorname{Bilin}_{R}(T \times M, N) \cong \operatorname{Hom}_{R}\left(M, \operatorname{Hom}_{R}(T, N)\right),
$$

which is done as follows. Given $b \in \operatorname{Bilin}_{R}(T \times M, N)$, which means a bilinear map $b: T \times M \rightarrow N$, the map $b(-, m): T \rightarrow N$ is in $\operatorname{Hom}_{R}(T, N)$. Vice versa, is $f \in \operatorname{Hom}_{R}\left(M, \operatorname{Hom}_{R}(T, N)\right)$ then $(t, m) \mapsto f(m)(t)$ defines a bilinear map $T \times M \rightarrow N$.

A straightforward verification shows that

$$
b \mapsto[m \mapsto b(-, m)]
$$

and

$$
f \mapsto[(t, m) \mapsto f(m)(t)]
$$

are each other's inverse, and moreover these maps are $R$-module homomorphisms. This proves the theorem.

## VII. 4 Projective modules

In ExampleVII.2.10 we saw that $\operatorname{Hom}(A,-)$ is in general not (right-)exact, in other words an exact sequence $M \rightarrow N \rightarrow 0$ does not necessarily give rise to an exact sequence $\operatorname{Hom}(A, M) \rightarrow \operatorname{Hom}(A, N) \rightarrow 0$. And in Section VI.3 we saw examples of modules $M, N$ that are not free, but $M \oplus N$ is a free module (the examples VI.3.5 andVI.3.6 even have $M=N$ ).

In this section the notion 'projective module' (Definition VII.4.1) is introduced, and it is shown that $\operatorname{Hom}(P,-)$ is exact precisely when $P$ is projective (see Theorem VII.4.5. Moreover it will be shown that $P$ being projective if and only if a module $Q$ exists such that $P \oplus Q$ is a free module (see Theorem VII.4.6.
VII.4.1 Definition. An $R$ module $P$ is called projective if for every surjective $R$ module homomorphism $M \xrightarrow{f} N$ and for every $R$-module homomorphism $h: P \rightarrow N$, so

there exists an $R$-module homomorphism $\tilde{h}: P \rightarrow M$, so


The next three lemmas, especially VII.4.2 and VII.4.4, assist in appreciating this definition.
VII.4.2 Lemma. Any free $R$-module $F$ is projective.

Proof. By definition $F \cong \oplus_{i \in I} R$ for some non-empty set $I$. Now given and exact sequence $M \xrightarrow{f} N \rightarrow 0$ and $h: F \rightarrow N$, one constructs $\tilde{h}$ as follows. Put

$$
n_{i}:=h\left(e_{i}\right), \quad \text { and choose } \quad m_{i} \in M \quad \text { with } \quad f\left(m_{i}\right)=n_{i},
$$

such $m_{i}$ exist since $f$ is surjective. Since $h$ is an $R$-module homomorphism, we have $h\left(\sum x_{i} e_{i}\right)=\sum x_{i} h\left(e_{i}\right)=\sum x_{i} n_{i}$. Next, define

$$
\tilde{h}: F \longrightarrow M, \quad \text { by } \quad \sum_{i \in I} x_{i} e_{i} \mapsto \sum_{i \in I} x_{i} m_{i} .
$$

Then $\tilde{h}$ is an $R$-module homomorphism and $f \tilde{h}=h$. This proves the lemma.
VII.4.3 Lemma. Let $R$ be a ring and let $M, N$ be $R$-modules. Suppose given two $R$-module homomorphisms

$$
M \xrightarrow[f]{\longrightarrow} N, \quad \text { and } \quad M \stackrel{g}{\leftrightarrows} N \quad \text { such that } f \circ g=\mathrm{id}_{N} .
$$

Then $f$ is surjective (one says that $g$ splits the exact sequence $M \xrightarrow{f} N \rightarrow 0$ ). We have

$$
M \cong \operatorname{Im}(g) \oplus \operatorname{Ker}(f), \quad \text { and } \quad \operatorname{Im}(g) \cong N
$$

Proof. We claim that $\operatorname{Im}(g) \cap \operatorname{Ker}(f)=\{0\}$ and that $\operatorname{Im}(g)+\operatorname{Ker}(f)=M$. If this holds, then the map

$$
\operatorname{Im}(g) \oplus \operatorname{Ker}(f) \longrightarrow M, \quad(a, b) \mapsto a+b
$$

is an isomorphism of $R$-modules (as the reader should verify!).
We first show $\operatorname{Im}(g) \cap \operatorname{Ker}(f)=\{0\}$. Take $x \in \operatorname{Im}(g) \cap \operatorname{Ker}(f)$, then

$$
\left.\begin{array}{rl}
\exists n \in N: & g(n)=x \\
f(x) & =0
\end{array}\right\} \Longrightarrow n=f(g(n))=f(x)=0,
$$

and therefore $x=g(n)=g(0)=0$.
Now take $m \in M$, then

$$
g f(m) \in \operatorname{Im}(g) \quad \text { and } \quad m=g f(m)+(m-g f(m))
$$

moreover it holds that

$$
f(m-g f(m))=f(m)-f g f(m)=f(m)-f(m)=0,
$$

and therefore $m-g f(m) \in \operatorname{Ker}(f)$. One concludes that every $m \in M$ can be written as a sum of the element $g f(m) \in \operatorname{Im}(g)$ and the element $m-g f(m) \in \operatorname{Ker}(f)$. This shows the lemma.
VII.4.4 Lemma. If $P$ is a projective $R$-module and $M \xrightarrow{f} P \longrightarrow 0$ is a surjective $R$-module homomorphism, then $M \cong P \oplus \operatorname{Ker}(f)$.

Proof. Use the definition of projective with $N=P$ and $h=\operatorname{id}_{P}: P \rightarrow P=N$ to obtain $\tilde{h}: P \rightarrow M$

$$
\begin{array}{ccccc}
M & \xrightarrow{f} \quad P \quad \rightarrow 0 \quad \text { with } \quad f \tilde{h}=h=\operatorname{id}_{P} . \\
\tilde{h} \uparrow & \nearrow \cong & & & \\
P
\end{array}
$$

In particular one concludes that $\tilde{h}$ splits the exact sequence $N \stackrel{f}{\rightarrow} P \rightarrow 0$. The lemma now follows from Lemma VII.4.3
VII.4.5 Theorem. If $P$ is a projective $R$-module, then $\operatorname{Hom}(P,-)$ is right exact. This means the following. If

$$
A \xrightarrow{f} B \xrightarrow{g} C \longrightarrow 0
$$

is an exact sequence of $R$-modules, then

$$
\operatorname{Hom}(P, A) \xrightarrow{f_{*}} \operatorname{Hom}(P, B) \xrightarrow{g_{*}} \operatorname{Hom}(P, C) \longrightarrow 0
$$

is also an exact sequence of $R$ modules.
Moreover, if an $R$-module $A$ is not projective, then $\operatorname{Hom}(A,-)$ is not right exact.
In conclusion, an $R$-module $P$ is projective precisely when $\operatorname{Hom}(P,-)$ is right exact.
Proof. Given are an exact sequence and a projective module $P$ as in the statement of the theorem.

Exactness at $\operatorname{Hom}(P, C)$ : this means we have to verify that $g_{*}$ is surjective. This follows immediately from the definition of projective: indeed, take $h \in \operatorname{Hom}(P, C)$. Then $\tilde{h} \in \operatorname{Hom}(P, C)$ exists with $f \tilde{h}=h$. This says $f_{*}(\tilde{h})=h$, proving surjectivity.

Exactness at $\operatorname{Hom}(P, B)$ : we have to show $\operatorname{Im}\left(f_{*}\right)=\operatorname{Ker}\left(g_{*}\right)$.
' $\subset$ ': let $\psi \in \operatorname{Im}\left(f_{*}\right)$, then one may write $\psi=f_{*}(\phi)$ for some $\phi \in \operatorname{Hom}(P, A)$. By definition, for any $x \in P$ one has $\psi(x)=f(\phi(x)) \in \operatorname{Im}(f)$. By assumption $\operatorname{Im}(f)=\operatorname{Ker}(g)$, hence $g(\psi(x))=0$. This implies $g_{*}(\psi)=0$, and therefore $\operatorname{Im}\left(f_{*}\right) \subset \operatorname{Ker}\left(g_{*}\right)$. ' $\supset$ ': let $\psi \in \operatorname{Hom}(P, B)$ satisfy $g_{*}(\psi)=0$. We will construct $\phi \in \operatorname{Hom}(P, A)$ with $f_{*}(\phi)=\psi$. The condition $g_{*} \psi=0$ means

$$
\operatorname{Im}(\psi) \subset \operatorname{Ker}(g)=\operatorname{Im}(f) .
$$

Furthermore $A \xrightarrow{f} \operatorname{Im}(f) \longrightarrow 0$ is an exact sequence. Considering $\psi$ as an $R$ module homomorphism $\psi: P \rightarrow \operatorname{Im}(f)$, the definition of projective (with $h:=\psi$, and $M:=A, N:=\operatorname{Im}(f))$ implies that $\phi: P \rightarrow A$ exists with $f \phi=\psi$. This means $f_{*}(\phi)=\psi$. As a consequence, $\operatorname{Ker}\left(g_{*}\right) \subset \operatorname{Im}\left(f_{*}\right)$.

Finally we prove the remaining assertion of the theorem. If $A$ is not projective, then an exact sequence $M \xrightarrow{f} N \rightarrow 0$ and an $h \in \operatorname{Hom}(A, M)$ exist, for which there is no $\tilde{h} \in \operatorname{Hom}(A, N)$ such that $f \tilde{h}=h$. This means that $\operatorname{Hom}(A,-)$, applied to the exact sequence

$$
\operatorname{Ker}(f) \hookrightarrow M \xrightarrow{f} N \longrightarrow 0,
$$

yields a sequence in which $f_{*}: \operatorname{Hom}(A, M) \rightarrow \operatorname{Hom}(A, N)$ is not surjective. Hence $\operatorname{Hom}(A,-)$ is not right exact. This finishes the proof.
VII.4.6 Theorem. Let $R$ be a ring and let $P$ be an $R$-module. Then:
$P$ is projective if and only if an $R$-module $Q$ exists with $P \oplus Q=F \quad$ for some free $R$-module $F$.

Proof. ' $\Leftarrow$ ': Suppose $P \oplus Q=F$ is a free $R$-module. Given a diagram

$$
\begin{aligned}
& M \underset{\rightarrow h}{\underset{\nearrow}{f}} N \quad N \quad 0, \\
& P
\end{aligned}
$$

in which $f$ is surjective, we must construct $\tilde{h}$ such that $f \tilde{h}=h$. Define

$$
g: F=P \oplus Q \longrightarrow N, \quad g(p, q):=h(p) \quad(p \in P, q \in Q) .
$$

Then $g \upharpoonright P=h$ (the restriction of $g$ to $P \cong P \oplus(0) \subset F$ ). By Lemma VII.4.2 $F$ is projective, hence $\tilde{g}$ exists with

$$
\begin{aligned}
& M \xrightarrow{f} \quad N \quad \rightarrow 0 \\
& \tilde{g} \uparrow \nearrow g \text { such that } f \tilde{g}=g \text {. } \\
& P \oplus Q=F
\end{aligned}
$$

Now define $\tilde{h}:=\tilde{g} \upharpoonright P$. Then $f \tilde{h}=h$ follows from $f \tilde{g}=g$ by restricting to $P$. ' $\Rightarrow$ ': Let $P$ be a projective $R$-module. We construct a free $R$-module $F$ by using $P$ as index set:

$$
F:=\bigoplus_{p \in P} R .
$$

Now define

$$
f: F \rightarrow P, \quad\left(x_{p}\right)_{p \in P} \mapsto \sum_{p \in P} x_{p} \cdot p .
$$

Clearly $f$ is an $R$-module homomorphism. Moreover $f$ is surjective since $f\left(\left(x_{q}\right)_{q \in P}\right)=p$ if one takes $x_{q}=0$ for all $q \neq p$ and $x_{p}=1$. LemmaVII.4.4 then implies $F=P \oplus \operatorname{Ker}(f)$, hence one can take $Q=\operatorname{Ker}(f)$. This proves Theorem|VII.4.6.

## VII. 5 Exercises

1. Let $R$ be a commutative (unitary) ring and suppose $f: M \rightarrow N$ is a non-surjective $R$-module homomorphism.
Show that $f_{*}: \operatorname{Hom}(R, M) \longrightarrow \operatorname{Hom}(R, N)$ is not surjective.
2. Let $R$ be a commutative ring and consider $R$-modules $M, N, P, Q$. Construct $R$-module isomorphisms

$$
\begin{aligned}
& \operatorname{Hom}(M \oplus N, P) \cong \operatorname{Hom}(M, P) \oplus \operatorname{Hom}(N, P), \\
& \operatorname{Hom}(M, P \oplus Q) \cong \operatorname{Hom}(M, P) \oplus \operatorname{Hom}(M, Q)
\end{aligned}
$$

3. Let $R$ be a commutative ring and let $I, J$ be ideals in $R$ such that $I+J=R$. Prove that $\operatorname{Hom}(R / I, R / J)=0$.
Conclude that the $\mathbb{R}[X]$-modules $\mathbb{R}[X] /(X-a)$ and $R[X] /(X-b)$ are not isomorphic if $a \neq b$.
4. Prove the " 5 -lemma": this is the following assertion. Suppose $R$ is a ring, and given is a commutative diagram of $R$-modules

in which the two rows are exact. Show that if $\ell$ is surjective and $q$ is injective and both $m$ and $p$ are isomorphisms, then $n$ is an isomorphism.
5. Given a ring $R$ and $R$-modules $M_{1}, M_{2}, N$, show that

$$
\left(M_{1} \oplus M_{2}\right) \otimes_{R} N \cong\left(M_{1} \otimes_{R} N\right) \oplus\left(M_{2} \otimes_{R} N\right)
$$

6. Prove that the tensor product is right exact. This means the following. Let $R$ be a ring and let $T$ be an $R$-module. Suppose

$$
M \rightarrow N \rightarrow P \rightarrow 0
$$

is an exact sequence of $R$-modules. Show that this results in a sequence

$$
T \otimes_{R} M \rightarrow T \otimes_{R} N \rightarrow T \otimes_{R} P \rightarrow 0
$$

which is also exact.
7. Given the linear maps $\mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ with matrices $\left(\begin{array}{ll}1 & 2 \\ 3 & 4\end{array}\right)$ and $\left(\begin{array}{ll}5 & 6 \\ 7 & 8\end{array}\right)$, respectively (with respect to the standard basis of $\mathbb{R}^{2}$ ).
Determine the Kronecker product of these linear maps (see Remark VII.3.12).
8. Let $R$ be a commutative ring and let $I \subset R$ be an ideal.

Prove: $R / I$ is a projective $R$-module if and only if
an ideal $J \subset R$ exists such that the canonical map $R \rightarrow R / I \times R / J$ given by $r \mapsto(r \bmod I, r \bmod J)$ defines an isomorphism of rings.
(Hint: if $g$ splits the canonical map $R \rightarrow R / I$, then consider $g(R / I) \subset R$.)
9. Show that the ideal $I=(X, Y) \subset \mathbb{R}[X, Y]$ is not projective as an $R$-module. (Hint: if $\phi$ splits the surjection

$$
R \oplus R \longrightarrow I \quad \text { given by } \quad(f, g) \mapsto f X+g Y
$$

then consider $\phi(X Y) \in R \oplus R$.
10. Suppose $M$ is a projective module over the commutative ring $R$.

Prove that $\operatorname{Hom}(M, R)$ is a projective $R$-module as well. (This $R$-module is usually called the dual of the $R$-module $M$ ).

