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## Number Theoretic Symbols in $K$-theory and Motivic Homotopy Theory

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## Abstract

We start out by reviewing the theory of symbols over number fields, emphasizing how this notion relates to classical reciprocity laws and algebraic $K$-theory. Then we compute the second algebraic $K$-group of the fields $\mathbf{Q}(\sqrt{-1})$ and $\mathbf{Q}(\sqrt{-3})$ based on Tate's technique for $K_{2}(\mathbf{Q})$, and relate the result for $\mathbf{Q}(\sqrt{-1})$ to the law of biquadratic reciprocity.

We then move into the realm of motivic homotopy theory, aiming to explain how symbols in number theory and relations in $K$-theory and Witt theory can be described as certain operations in stable motivic homotopy theory. We discuss Hu and Kriz' proof of the fact that the Steinberg relation holds in the ring $\pi_{* \alpha} 1$ of stable motivic homotopy groups of the sphere spectrum 1. Based on this result, Morel identified the ring $\pi_{* \alpha} \mathbf{1}$ as the Milnor-Witt $K$-theory $K_{*}^{M W}(F)$ of the ground field $F$. Our last aim is to compute this ring in a few basic examples.

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## Introduction

## From reciprocity laws to $K$-theory

In number theory, one ubiquitous problem is that of solving equations over finite fields. In the case of quadratic equations, the problem led Legendre to introduce the symbols ( $-/ q$ ) representing solvability of quadratic equations over $\mathbf{F}_{q}$. The pinnacle of theory is the quadratic reciprocity law of Gauss (Theorem 1.1) relating Legendre symbols over different primes.

Since the time of Gauss, the Legendre symbols have been generalized to arbitrary number fields, and to symbols representing solvability over finite fields of equations of any degree. Most notable are perhaps the Hilbert symbols $\left(\frac{-,-}{v}\right)$, which we will introduce in Chapter 1. These symbols satisfy properties like antisymmetry, bimultiplicativity and the relation

$$
\begin{equation*}
\left(\frac{x, 1-x}{v}\right)=1 \tag{1}
\end{equation*}
$$

This last property turned out to be very important. Indeed, Steinberg's work on central extensions led him to study certain bimultiplicative maps with values in an abelian group satisfying the relation (1) above Tat76a. Such maps are now called Steinberg symbols, and the relation (1) is called the Steinberg relation. Hilbert symbols provide important examples of Steinberg symbols.

By Matsumoto's description of the second $K$-group $K_{2}(F)$ of a field $F$, it turns out that $K_{2}(F)$ is the universal object with respect to Steinberg symbols. By formulating the classical law of quadratic reciprocity in terms of Hilbert symbols on $\mathbf{Q}$, this insight led Tate to the discovery that the structure theorem for $K_{2}(\mathbf{Q})$ is actually equivalent to quadratic reciprocity.

## Classification problems and Milnor's conjecture

In every part of mathematics there are objects constructed for the purpose of classifying other mathematical objects. For example:

- In the machinery of algebraic $K$-theory, we have the group $K_{0}$ which classifies, e.g., isomorphism classes of finitely generated projective modules over a ring or classes of vector bundles on a scheme. The higher algebraic $K$-groups also take into account certain symmetry properties of the ground object of study.
- A classical problem that has been the subject of extensive study over the years is that of classifying quadratic forms over fields. The so-called Witt ring $W(F)$ of a field $F$ evolved out of this problem, whose elements represents equivalence classes of quadratic spaces over $F$.
- Parallel to the story of the Witt ring is that of the Galois cohomology $H_{\text {Gal }}^{*}(F ;-)$. It is known for example that $H_{\text {Gal }}^{2}\left(F ; F_{\text {sep }}^{\times}\right)$-where $F_{\text {sep }}$ is the separable closure of $F$-classifies equivalence classes of central simple algebras over the field $F$, while $H_{\text {Gal }}^{1}(F ; \mathbf{Z} / 2)$ is the set of square classes of elements in $F^{\times}$.

The algebraic $K$-theory, the Witt ring and the Galois cohomology of a field $F$ are clearly highly interesting invariants of $F$. What Milnor did in Mil70 was to define a ring $K_{*}^{M}(F)$, purely in terms of generators and relations, which would serve as a sort of "best approximation" to the Witt theory and Galois cohomology of $F$ [Dug04, 1.7]. The Milnor $K$-groups $K_{n}^{M}(F)$ agree with algebraic $K$-theory for $n$ equal to 0,1 and 2, but not generally in higher degrees. Letting $I(F)$ denote the fundamental ideal of the Witt ring $W(F)$, Milnor produced homomorphisms

and proved that these maps are isomorphisms in degrees 0,1 and 2 as well as always being surjective. But it required the machinery of motivic homotopy theory-invented during the 1990s - to prove Milnor's conjecture that the maps are actually isomorphisms in every degree. More on this in Chapter 2 .

The settlement of the Milnor conjecture shed a whole new light on the theory of quadratic forms over fields. For example, it made it possible to think about relations in both Witt theory and Milnor $K$-theory as operations in stable motivic homotopy theory. More precisely, further studies in motivic homotopy theory revealed that the stable motivic homotopy groups $\pi_{* \alpha} \mathbf{1}$ of the motivic sphere spectrum 1 could be described as the so-called Milnor-Witt $K$-theory $K_{*}^{M W}(F)$ of the ground field $F$. As the sphere spectrum is initial in the category of ring spectra, any homotopy group $\pi_{m+n \alpha} R$ of a ring spectrum $R$ will inherit the relations in the group $\pi_{m+n \alpha} \mathbf{1}$ via the map $\pi_{m+n \alpha} \mathbf{1} \rightarrow \pi_{m+n \alpha} R$ induced by the unique unit map $\mathbf{1} \rightarrow R$. In this regard, Milnor-Witt $K$-theory is a fundamental object in stable motivic homotopy theory.

## One prime at a time...

By the fundamental theorem of arithmetic, there is an isomorphism of abelian groups

$$
\begin{aligned}
\mathbf{Q}^{\times} & \cong \mathbf{Z} / 2 \oplus \bigoplus_{p \text { prime }} \mathbf{Z} \\
x & \longmapsto\left(\operatorname{sgn}(x), \operatorname{ord}_{2}(x), \operatorname{ord}_{3}(x), \ldots, \operatorname{ord}_{p}(x), \ldots\right)
\end{aligned}
$$

Here $\operatorname{sgn}(x)$ is the sign of the rational number $x$, and $\operatorname{ord}_{p}(x)$ is the $p$-adic valuation of $x$. This isomorphism suggests that in order to solve a problem concerning the field $\mathbf{Q}$, it might be helpful to consider the problem in each completion $\mathbf{Q}_{p}$ of $\mathbf{Q}$, as well as in $\mathbf{R}$. In other words, one should consider one prime at a time. Hasse-Minkowski's local-global principle MH73, Corollary 2.4, p.89] is a typical example of this, stating that a quadratic form defined over $\mathbf{Q}$ represents zero if and only if it represents zero in each completion $\mathbf{Q}_{p}$, as well as in $\mathbf{R}$. In fact, this local-global philosophy is used extensively in class field theory, as we will see examples of in Chapter 1.

Turning to the realm of $K$-theory, it was Tate who was the first one out to compute the group $K_{2}(\mathbf{Q})$. Using additive notation, the result is

$$
K_{2}(\mathbf{Q}) \cong \mathbf{Z} / 2 \oplus \bigoplus_{p \text { prime }} K_{1}\left(\mathbf{F}_{p}\right) .
$$

Tate's method of computation consisted of defining a filtration of $K_{2}(\mathbf{Q})$ by subgroups indexed over the primes, and then analyzing each of these subgroups separately in order to "conclude globally". By Hilbert's formulation of the known reciprocity laws by means of considering one reciprocity symbol for each prime (see Theorem 1.55), Tate's computation made it possible to relate $K$-groups of number fields to reciprocity laws. More on all this in Chapter 3 .

Tate's computation of $K_{2}(\mathbf{Q})$ inspired Milnor to give a new proof of the structure theorem for the Witt ring of $\mathbf{Q}$ [MH73, Theorem 2.1, p.88]:

$$
W(\mathbf{Q}) \cong \mathbf{Z} \oplus \bigoplus_{p \text { prime }} W\left(\mathbf{F}_{p}\right)
$$

Again, the proof consists of defining a filtration of $W(\mathbf{Q})$, one subgroup for each prime, and considering one of these at a time.

In Chapter 5, we will see that this method also carries over to the case of Milnor-Witt $K$-theory. More precisely, for $n \geq 2$ we compute $K_{n}^{M W}(\mathbf{Q})$; the result being

$$
K_{n}^{M W}(\mathbf{Q}) \cong \mathbf{Z} \oplus \bigoplus_{p \text { prime }} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)
$$

## Outline

Chapter 1 provides background material for understanding number theoretic aspects of the $K$ theory we consider in Chapter 3. This chapter has two main themes: First we give a more or less detailed treatment of the properties of the number fields $\mathbf{Q}(\sqrt{-1})$ and $\mathbf{Q}(\sqrt{-3})$, which is needed later on. We then move on to discuss symbols over number fields and their relation to classical reciprocity laws. Number theoretic symbols are perhaps the main connection between number theory and classical algebraic $K$-theory, so we make an effort to introduce this concept thoroughly. In particular, we take a detour to visit class field theory and the origins of tame symbols. It is of course possible to take the formula in Proposition 1.52 as the definition of tame symbols and then move directly onto $K$-theory, but we find it illuminating to consider a more unified treatment of all the different number theoretic symbols, which is achieved by realizing these symbols as offspring of the local reciprocity map of local class field theory. As a motivation, we provide some background material on the history of reciprocity laws in the beginning of Chapter 1 . We revisit classical reciprocity laws in new guises in Section 1.6, finishing the chapter by taking a closer look at biquadratic reciprocity.

Chapter 2 introduces some preliminary theory on classical algebraic $K$-theory. We also go through the definition and basic properties of Milnor $K$-theory and Witt theory. Finally, we briefly explain Milnor's conjecture relating Milnor $K$-theory, Witt theory and Galois cohomology.

Chapter 3 deals with some explicit calculations. Specifically, we compute the second $K$ group of the fields $\mathbf{Q}(\sqrt{-1})$ and $\mathbf{Q}(\sqrt{-3})$ based on Tate's technique for $K_{2}(\mathbf{Q})$. We hope to illustrate how the arithmetic of the given number field reveals itself in the $K$-group, e.g., by making use of previous results on units; the Euclidean algorithm; the structure of residue fields and the ramification and splitting behavior of primes in these number fields. In particular, we connect the result on $K_{2}(\mathbf{Q}(\sqrt{-1}))$ to the law of biquadratic reciprocity.

Chapter 4 is a brief introduction to motivic homotopy theory. The goal of this chapter is to provide background material and motivation for studying Milnor-Witt $K$-theory. The reader who wishes may therefore skip this chapter, as the definition of Milnor-Witt $K$ groups is purely algebraic and does not involve any motivic machinery. Note however that the defining generators and relations of Milnor-Witt $K$-theory have natural explanations as elements in stable motivic homotopy groups. Moreover, a few results from this chapter are used in the beginning of Chapter 5.

Chapter 5 aims to study how symbols in $K$-theory and Witt theory may be thought of as coming from operations in stable motivic homotopy theory. We start out by explaining Hu and Kriz' proof that the Steinberg relation holds in certain stable motivic homotopy groups, which was the starting point for Morel and Hopkins' discovery of Milnor-Witt $K$-theory. We take a brief look at how the ring spectrum map $\mathbf{1} \rightarrow \mathrm{KGL}$ behaves with respect to the well known symbols in $K$-theory. Then we proceed to study the properties of Milnor-Witt $K$-groups, and use similar methods as those of Chapter 3 to compute the Milnor-Witt $K$-theory of $\mathbf{Q}$ and $\mathbf{Q}(\sqrt{-1})$.

In Chapter 1 familiarity with some basic algebraic number theory is an advantage - for example a rough equivalent of the first two chapters of Jan96. However, we provide references to
all the results we are using. In Chapter 4 we leave out background material on model categories, simplicial sets, sites and abstract homotopy theory. The interested reader may consult $\left[\mathrm{DLO}^{+} 07\right]$ and [Hov99]. The remaining chapters should be fairly self-contained.

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We exploit the opportunity to address a formal apology to Paul Arne in the wake of the author's reckless behavior in the Easter holiday, during which he - without notice - went fishing instead of showing up at the gym. It shall not happen again.

## Conventions and notation

Throughout this text, the symbol $F$ will denote a field.
By saying that a number field $F$ is Euclidean we mean that the ring of integers $\mathcal{O}_{F}$ in $F$ is Euclidean with respect to the norm map Nm: $\mathcal{O}_{F} \rightarrow \mathbf{Z}$. By abuse of notation, we denote also by Nm the counting norm $\operatorname{Nm}(\mathfrak{a})$ defined on the set of integral ideals $\mathfrak{a}$ of $\mathcal{O}_{F}$, i.e., $\operatorname{Nm}(\mathfrak{a})$ is the integer $\#\left(\mathcal{O}_{F} / \mathfrak{a}\right)$. Recall that, if $\mathfrak{a}=(a)$ is a principal ideal of $\mathcal{O}_{F}$, then $\operatorname{Nm}(a)=\operatorname{Nm}(\mathfrak{a})$ (see [Jan96, pp.42-45]).

We use the word "place" for an equivalence class of nontrivial absolute values on a number field. Note that the word "prime" is often used for the same concept in the literature.

Following Gras [Gra03], the word "ramification" is-as opposed to classical literature - used only in connection with finite places (i.e., the places corresponding to prime ideals in the ring of integers); the phenomenon of real infinite places becoming complex in some extension of number fields will instead be referred to as complexification (however, we shall rarely need to make use of this term. The important thing to note is that infinite places are not included when we speak about ramification).

When discussing algebraic $K$-theory of rings, we implicitly assume that all rings in consideration are commutative and unital. This is done because we only apply the $K$-theoretic machinery to such rings, and furthermore in order to avoid overstating the word commutative. Note however that several of the constructions of Chapter 2 apply to any associative ring.

When we consider the algebraic $K$-groups $K_{n}(F)$ of fields in Chapter 2 and Chapter 33, we will use multiplicative notation in order to keep in line with the notations of class field theory. In the context of Milnor $K$-theory, however, we switch to additive notation (e.g., compare $\{x y, z\}=\{x, z\}\{y, z\} \in K_{2}(F)$ and $\left.\{x y, z\}=\{x, z\}+\{y, z\} \in K_{2}^{M}(F)\right)$.

As we shall see in Chapter 4, we will denote by $S^{\alpha}$ the motivic sphere $\mathbf{G}_{m}$. Note however that the symbol $S_{t}^{1}$ is also often used in the literature (the $t$ coming from the name "Tate circle"); this explains the notation $\Sigma_{t}:=-\wedge S^{\alpha}$ used in Section 5.2 of Chapter 5

## CHAPTER 1

## Results from Algebraic Number Theory

The history of reciprocity laws in number theory is impressively long and rich, to which several great names in mathematics have contributed. Below we start out by reviewing a tiny fraction of this story. From the time of Lagrange and Gauss, through the generalizations of Jacobi, Eisenstein, Hilbert and Artin we arrive at Tate and his discovery of the connection between reciprocity laws and $K$-theory. We then move on to study this connection in detail in the subsequent chapters.

### 1.1 Reciprocity laws

The famous law of quadratic reciprocity was stated in its complete form by Lagrange, but it was Gauss who first gave a full proof [Mil13b, vii]. Legendre introduced the symbols $(p / q)$ for distinct odd primes $p$ and $q$, demanding

$$
\left(\frac{p}{q}\right):=\left\{\begin{aligned}
1 & \text { if } p \text { is a square modulo } q \\
-1 & \text { otherwise. }
\end{aligned}\right.
$$

The Legendre symbols allow for a compact formulation of the quadratic reciprocity law, which by now is the most common phrasing of the theorem:

Theorem 1.1 (Quadratic reciprocity [Lem00, p.vi]). Let $p$ and $q$ be two distinct odd primes. Then

$$
\begin{equation*}
\left(\frac{p}{q}\right)=(-1)^{\frac{p-1}{2} \cdot \frac{q-1}{2}}\left(\frac{q}{p}\right) . \tag{1.1}
\end{equation*}
$$

Moreover, the following supplementary laws hold, for $p$ an odd prime:

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

In other words, $p$ is a square modulo $q$ if and only if $q$ is a square modulo $p$ unless both $p$ and $q$ are congruent to 3 modulo 4 , in which case $p$ is a square modulo $q$ if and only if $q$ is a nonsquare modulo $p$.

There are several possible reasons as to why Gauss considered the quadratic reciprocity law worthy of the name "aureum theorema" (golden theorem). First of all, the very existence of a connection between two arbitrarily chosen prime numbers is quite astonishing in its own right. Quadratic reciprocity strengthens the philosophy that in order to understand an object, we should try to understand how it relates to other objects. In this case the objects are integers, and the reciprocity law states that there is indeed a rather mysterious connection between two prime numbers emanating as we turn to the "relative case", i.e., look at their residue classes. This train of thought has of course been exploited to its full extent throughout modern mathematics. For example, it was the groundbreaking philosophy of Grothendieck that we should consider morphisms between objects rather than just the objects themselves, and the zeroth $K$-group in
algebraic $K$-theory is an example of how we can understand a ring by understanding modules over the ring.

Furthermore, the quadratic reciprocity law solves a certain inverse problem. Given a prime number $p$, determining the integers that are quadratic residues modulo $p$ is clearly a finite problem as one can simply check all elements in the finite group $\mathbf{Z} / p$. The inverse problem however, i.e., fixing an integer $n$ and asking for which primes $p$ is $n$ a quadratic residue modulo $p$, is a question of seemingly an entirely different nature. It is now the strength of Formula 1.1 becomes apparent: The "infinite" inverse problem can be transformed into the finite direct problem of determining the quadratic residues modulo $n$.

To study quadratic residues modulo primes is to study quadratic polynomials over finite fields, and motivation for the latter is readily given by the Hasse-Minkowski theorem on the local-global principle for quadratic forms. For instance, if $q$ is an odd prime number the polynomial $f(x):=x^{2}-q$ splits in $\mathbf{F}_{p}[x]$ if and only if $(q / p)=1$. Thus, the problem of determining the splitting of $f$ modulo different primes is an example of an "inverse problem" as described above: Each time one changes the prime $p$ and asks for the splitting behavior over $\mathbf{F}_{p}$ one has to compute a new Legendre symbol $(q / p)$. Using quadratic reciprocity we convert the problem into computing $(p / q)$ instead-a much more comfortable task, requiring only to check the $q$ residue classes once and for all.

Example 1.2 (Wym72, Example 1, p.573]). Let $f(x):=x^{2}-17$, and let us determine all primes $p$ such that $f$ splits over $\mathbf{F}_{p}$. A priori, we must compute $(17 / p)$ for all $p$. However, quadratic reciprocity yields

$$
\left(\frac{17}{p}\right)=\left(\frac{p}{17}\right)
$$

since $17 \equiv 1(\bmod 4)$. The quadratic residues modulo 17 are $1,2,4,8,9,13,15$ and 16 , hence $f$ splits over $\mathbf{F}_{p}$ if and only if $p$ is congruent modulo 17 to one of the listed numbers.

The splitting of quadratic polynomials over a number field $F$ modulo different primes is closely related to the knowledge of the extension of $F$ defined by the polynomial. In fact, knowing of the set of prime ideals that split in an extension $L / F$ of number fields is enough to determine $L$ Mil13b, p.2]. Quadratic reciprocity is the key to determine which primes that split in quadratic extensions of $\mathbf{Q}$, and should therefore be considered the first step toward understanding the splitting of prime ideals in arbitrary extensions of number fields. In other words-as Milne notes in [Mil13b, p.3]-quadratic reciprocity should be viewed as the first result in class field theory.

Based on the following description of the Legendre symbol one generalizes it to $n$-th power residue symbols (see Definition 1.56 below):

Proposition 1.3. The Legendre symbol satisfies $(p / q) \equiv p^{\frac{q-1}{2}}(\bmod q)$.
Proof. The sequence

$$
1 \longrightarrow \mu_{2} \longleftrightarrow \mathbf{F}_{q}^{\times} \xrightarrow{\alpha} \mathbf{F}_{q}^{\times} \xrightarrow{\beta} \mu_{2} \longrightarrow 1
$$

in which $\alpha(x):=x^{2}$ and $\beta(y):=y^{\frac{q-1}{2}}$ is clearly a complex by Fermat's little theorem. To show exactness at the second $\mathbf{F}_{q}^{\times}$, assume $x \in \mathbf{F}_{q}^{\times}$is such that $x^{\frac{q-1}{2}}=1$. Let $\zeta$ be a generator for the cyclic group $\mathbf{F}_{q}^{\times}$with $x=\zeta^{k}$, then

$$
x^{\frac{q-1}{2}}=\zeta^{k \frac{q-1}{2}}=1
$$

hence $2 \mid k$, hence $x$ is a square. Thus $p$ is a square modulo $q$ if and only if $p^{\frac{q-1}{2}} \equiv 1(\bmod q)$.
Gauss noticed that the ring of Gaussian integers $\mathbf{Z}[i]$, where $i^{2}=-1$, was the appropriate setting for a formulation of a biquadratic residue law [Lem00, p.vi]:

Definition 1.4. Let $\pi$ and $\tau$ be two distinct prime elements of $\mathbf{Z}[i]$ not dividing 2 . The biquadratic residue symbol $(\pi / \tau)_{4}$ is the unique element in $\mu_{4}=\{ \pm 1, \pm i\}$ satisfying

$$
\left(\frac{\pi}{\tau}\right)_{4} \equiv \pi^{\frac{\mathrm{Nm}(\tau)-1}{4}}(\bmod \tau) .
$$

We have $(\pi / \tau)_{4}=1$ if and only if $x^{4} \equiv \pi(\bmod \tau)$ has a solution in the Gaussian integers.
Definition 1.5. A Gaussian integer $\alpha$ is primary if $\alpha \equiv 1\left(\bmod (1+i)^{3}\right)$.
Theorem 1.6 (Biquadratic reciprocity Lem00, p.vi]). Let $\pi$ and $\tau$ be distinct primary Gaussian primes. Then

$$
\left(\frac{\pi}{\tau}\right)_{4}=(-1)^{\frac{\mathrm{Nm}(\pi)-1}{4} \cdot \frac{\mathrm{Nm}(\tau)-1}{4}}\left(\frac{\tau}{\pi}\right)_{4}
$$

Moreover, the following supplementary laws hold:

$$
\left(\frac{i}{\pi}\right)_{4}=i^{\frac{\mathrm{Nm}(\pi)-1}{4}}, \quad\left(\frac{1+i}{\pi}\right)_{4}=i^{\nu-\mu}
$$

where the integers $\mu$ and $\nu$ are determined as follows: By Proposition 1.23, the group $U_{1+i} /\left(U_{1+i}^{4}\right)$ of local units modulo fourth powers is isomorphic to $\mu_{4} \oplus\langle\overline{3+2 i}\rangle \oplus\langle 5\rangle$. Using this isomorphism, write $\pi \equiv i^{k}(3+2 i)^{\mu} 5^{\nu}\left(\bmod \left(U_{1+i}\right)^{4}\right)($ see Section 1.3.1 and Section 1.5 for explanation of the notation).

Theorem 1.6 will be proved at the end of this chapter using power reciprocity. Later we will also consider how biquadratic reciprocity is connected to the second $K$-group of $\mathbf{Q}(i)$.

The biquadratic reciprocity theorem settles a reciprocity law for the cyclotomic extension $\mathbf{Q}\left(\zeta_{4}\right)=\mathbf{Q}(i)$, and it is natural to ask whether there is a reciprocity law for each cyclotomic field $\mathbf{Q}\left(\zeta_{p}\right)$. Jacobi led the search for such a general law, but the failure of unique factorization in these fields 1 led to serious obstacles. It was not until Kummer's invention of ideal numbers that progress was made possible, and eventually it was Eisenstein who succeeded in finding a reciprocity law for all odd primes Lem00, p.vii]. In the following statement, we refer the reader to Definition 1.56 and the discussion following it for the meaning of the symbols $(-/ \alpha)_{n}$.

Theorem 1.7 (Reciprocity at odd primes Lem00, p.vii]). Let $p$ be an odd prime. We call an element $\alpha \in \mathbf{Z}\left[\zeta_{p}\right]=\mathcal{O}_{\mathbf{Q}\left(\zeta_{p}\right)}$ primary if $\alpha$ is congruent to a rational integer modulo $\left(1-\zeta_{p}\right)^{2}$.

If $\alpha \in \mathbf{Z}\left[\zeta_{p}\right]$ is primary, then

$$
\left(\frac{\alpha}{a}\right)_{p}=\left(\frac{a}{\alpha}\right)_{p}
$$

for all $a \in \mathbf{Z}$ prime to $p$.
The search for a complete generalization of quadratic reciprocity to arbitrary number fields is one of the main objectives of algebraic number theory. Indeed, this is the content of Hilbert's ninth problem. Hilbert discovered that the classical law of quadratic reciprocity could be phrased in terms of Hilbert symbols (see Definition 1.53) as a product formula

$$
\begin{equation*}
\prod_{v \in \operatorname{Pl}_{\mathbf{Q}}}\left(\frac{x, y}{v}\right)=1 \tag{1.2}
\end{equation*}
$$

for $x, y \in \mathbf{Q}^{\times}$, and conjectured that a similar formula should hold for any number field $F$ Lem00, p.viii]. This product formula was settled for abelian extensions of number fields by the breakthrough of class field theory and Artin's general reciprocity law, as we will see in Theorem 1.55 (in which the notations used in 1.2 are also explained). We shall also see explicitly that if we work over the number field $\mathbf{Q}(i)$, the corresponding product formula reduces to the biquadratic reciprocity law.

Artin's general reciprocity law will be stated in its idèlic form in Theorem 1.46 below, but for completeness we also mention the classical ideal group version:

[^0]Theorem 1.8 (Artin's reciprocity law [Gra03, II 4.4]). Let $F$ be a number field and $L / F$ a finite abelian extension. Let $T$ be the set of ramified places in $L / F$ and let $\mathfrak{m}$ be a modulus built from $T$. Then the Artin symbol $\alpha_{L / F}: I_{T} \rightarrow \operatorname{Gal}(L / F)$ induces a canonical isomorphism

$$
\frac{I_{T}}{P_{T, \mathfrak{m}, \operatorname{pos}} \operatorname{Nm}_{L / F}\left(I_{L, T}\right)} \xrightarrow[\alpha_{L / F}]{\cong} \operatorname{Gal}(L / F) .
$$

Here $I_{T}$ is the group of fractional ideals of $F$ prime to $T$, and $\mathfrak{m}$ is an integral ideal of the form $\mathfrak{m}:=\prod_{v \in T} \mathfrak{p}_{v}^{n_{v}}$ for some $n_{v} \geq 0$, where $\mathfrak{p}_{v}$ is the maximal ideal corresponding to the finite place $v$. The group $P_{T, \mathfrak{m}, \text { pos }}$ is the group of principal fractional ideals $(x)$ of $F$ prime to $T$, with $x$ totally positive and satisfying $\operatorname{ord}_{v}(x-1) \geq n_{v}$ for all $v \in T$, i.e., $x \equiv 1(\bmod \mathfrak{m})$.

Later Tate generalized the group isomorphism aspect of Artin's reciprocity law via the cohomological interpretation of class field theory. Specifically, Tate showed that-for $L / F$ a finite Galois extension - cup product with the fundamental class

$$
u_{L / F} \in \widehat{H}^{2}\left(\operatorname{Gal}(L / F), C_{L}\right)
$$

induces a canonical isomorphism

$$
x \mapsto x \smile u_{L / F}: \widehat{H}^{r}(\operatorname{Gal}(L / F) ; \mathbf{Z}) \xrightarrow{\cong} \widehat{H}^{r+2}\left(\operatorname{Gal}(L / F) ; C_{L}\right)
$$

for all $r \in \mathbf{Z}$; see Gra03, p.109]. Here the $\widehat{H}^{r}$ are Tate's modified cohomology groups (see, e.g., Wei94, Definition 6.2.4]), and $C_{L}$ is the idèle class group of $L$, defined in Definition 1.39 . Tate was also the first to discover the connection between the second algebraic $K$-group and reciprocity laws - more precisely, he computed the group $K_{2}(\mathbf{Q})$ and noticed that the method was very similar to the formal part of Gauss' first proof of quadratic reciprocity [Tat76a, p.318]. It involved induction over the primes and use of the Euclidean algorithm; we will return to this in Chapter 3. It is this $K$-theoretic view on classical number theory we shall keep in mind in the subsequent chapters.

As a final remark we mention that the story of the general reciprocity law is far from over: Although class field theory solves Hilbert's ninth problem in the case of abelian extensions of number fields, the general case remains unsolved to this date.

### 1.2 Preliminary results on quadratic fields

Let $d$ be a squarefree integer, and consider the field $F:=\mathbf{Q}(\sqrt{d})$. Recall the following result, whose proof can be found in, e.g., [IR90, Proposition 13.1.1].
Proposition 1.9. The ring of integers $\mathcal{O}_{F}$ in $F=\mathbf{Q}(\sqrt{d})$ equals

$$
\mathcal{O}_{F}=\left\{\begin{array}{cl}
\mathbf{Z}[\sqrt{d}] & \text { if } d \equiv 2,3(\bmod 4) \\
\mathbf{Z}\left[\frac{1+\sqrt{d}}{2}\right] & \text { if } d \equiv 1(\bmod 4)
\end{array}\right.
$$

Note that, by the fundamental identity Mil13a, Theorem 3.34], if a rational prime $p$ is ramified in a quadratic extension $F / \mathbf{Q}$, it is totally ramified, and similarly if $p$ splits, it is totally split.

Proposition 1.10. Consider the field $F=\mathbf{Q}(\sqrt{d})$ where $d$ is a squarefree integer. Let $p \in \mathbf{Z}$ be a rational prime. Assume first that $p>2$, then

- $p$ ramifies $\Longleftrightarrow p \mid d$;
- $p$ is inert $\Longleftrightarrow\left(\frac{d}{p}\right)=-1$;
- $p$ splits $\Longleftrightarrow\left(\frac{d}{p}\right)=1$.

For the prime 2, we have the following.

- 2 ramifies $\Longleftrightarrow d \equiv 2,3(\bmod 4)$;
- 2 is inert $\Longleftrightarrow d \equiv 5(\bmod 8)$;
- 2 splits $\Longleftrightarrow d \equiv 1(\bmod 8)$.

Proof. Write $\mathcal{O}_{F}=\mathbf{Z}[\theta]$. From Proposition 1.9 we see that the minimal polynomial $f$ of $\theta$ is either $x^{2}-d$ or $x^{2}-x-(d-1) / 4$ according as $d \equiv 2,3(\bmod 4)$ or $d \equiv 1(\bmod 4)$. Hence the discriminant $\delta_{F / \mathbf{Q}}$ is either $4 d$ or $d$, so the odd prime $p$ ramifies if and only if $p \mid d$ (cf. Jan96, I Theorem 7.3]). If $p \nmid d$, we have

$$
\frac{\mathcal{O}_{F}}{p \mathcal{O}_{F}} \cong \frac{\mathbf{Z}[x]}{(p, f)} \cong \frac{\mathbf{F}_{p}[x]}{(\bar{f})}
$$

where $\bar{f}$ is the reduction of $f$ modulo $p$. Thus, if $f$ splits modulo $p$ we have $\mathcal{O}_{F} / p \mathcal{O}_{F} \cong \mathbf{F}_{p} \oplus \mathbf{F}_{p}$ by the Chinese remainder theorem; if $f$ is irreducible modulo $p$ we have $\mathcal{O}_{F} / p \mathcal{O}_{F} \cong \mathbf{F}_{p^{2}}$-the finite field with $p^{2}$ elements. In the first case $p$ splits; in the second case, $p$ is inert. But $f$ splits modulo $p$ if and only if $(d / p)=1$, so this finishes the case when $p$ is odd.

We turn to the case $p=2$. We have seen that $2 \mid \delta_{F / \mathbf{Q}}$ if and only if $d \equiv 2,3(\bmod 4)$, in which case 2 ramifies. If $d \equiv 1(\bmod 4)$, the minimal polynomial $f(x)=x^{2}-x-(d-1) / 4$ of the extension $F / \mathbf{Q}$ is irreducible modulo 2 if and only if $(d-1) / 4 \equiv 1(\bmod 2)$, which is equivalent to $d \equiv 5(\bmod 8)$. Similarly, $f$ splits modulo 2 if and only if $(d-1) / 4 \equiv 0(\bmod 2)$, which happens if and only if $d \equiv 1(\bmod 8)$.

Note that these results also give some information about the structure of the residue fields at the different prime ideals. We record this as a proposition for future reference.

Proposition 1.11. If $F$ is a number field and $\mathfrak{p} \in \operatorname{Spec} \mathcal{O}_{F}$ is a nonzero prime ideal above the prime $p \in \mathbf{Z}$, then

$$
k(\mathfrak{p}) \cong \mathbf{F}_{\mathrm{Nm} \mathfrak{p}}
$$

where $k(\mathfrak{p}):=\mathcal{O}_{F} / \mathfrak{p}$ is the residue field of $\mathfrak{p}$.
In particular, if $F=\mathbf{Q}(\sqrt{d})$ is a quadratic number field, we have the following description of $k(\mathfrak{p})$ :

- If $p$ ramifies or splits, then $k(\mathfrak{p}) \cong \mathbf{F}_{p}$.
- If $p$ is inert, $k(\mathfrak{p}) \cong \mathbf{F}_{p^{2}}$.

Proof. We have $k(\mathfrak{p}) \cong \mathbf{F}_{\mathrm{Nm} \mathfrak{p}}$ by the definition of the counting norm Nmp as the number of elements in $k(\mathfrak{p})$.

In the case when $F$ is a quadratic number field, write $p \mathcal{O}_{F}=\mathfrak{P}_{1}^{e_{1}} \cdots \mathfrak{P}_{s}^{e_{s}}$ for $\mathfrak{P}_{j} \in \operatorname{Spec} \mathcal{O}_{F}$. By the fundamental identity we have

$$
\sum_{j=1}^{s} e_{j} f_{j}=[F: \mathbf{Q}]=2
$$

where $f_{j}:=\left[k\left(\mathfrak{P}_{j}\right): \mathbf{F}_{p}\right]$ are the local degrees. We have $\mathfrak{P}_{j}=\mathfrak{p}$ for some $j$ by assumption, and since the extension $F / \mathbf{Q}$ is Galois, all the local degrees are equal [Lan70, Corollary 2, p.26]. Thus we can apply the results above to deduce the following:

- If $p$ ramifies, then $s=1$ and $e_{1}=2$, thus $f_{1}=1$. Hence $k(\mathfrak{p})$ is a degree 1 extension of $\mathbf{F}_{p}$, in other words isomorphic to $\mathbf{F}_{p}$.
- If $p$ splits, we get $s=2$ and $f_{j}=e_{j}=1(j=1,2)$. Hence $k(\mathfrak{p}) \cong \mathbf{F}_{p}$.
- That $p$ is inert means that $s=1$ and $e_{1}=1$, thus $f_{1}=2$. In other words, $k(\mathfrak{p})$ is a degree 2 extension of $\mathbf{F}_{p}$. Hence $k(\mathfrak{p}) \cong \mathbf{F}_{p^{2}}$.

| $\delta_{\mathbf{Q}(\sqrt{d}) / \mathbf{Q}}$ | -3 | -4 | -7 | -8 | -11 | -15 |
| :---: | ---: | ---: | ---: | ---: | ---: | ---: |
| $M(\mathbf{Q}(\sqrt{d}))$ | $\frac{1}{3}$ | $\frac{1}{2}$ | $\frac{4}{7}$ | $\frac{3}{4}$ | $\frac{9}{11}$ | $\frac{16}{15}$ |

Table 1.1: Euclidean minima for the first six imaginary quadratic number fields.

The existence of a Euclidean algorithm on a number field simplifies certain $K$-theoretic questions; we shall make use of this in Chapter 3. In the case of imaginary quadratic number fields there are five instances for which there exists a Euclidean algorithm.

Proposition 1.12 (Lem95, Proposition 4.2]). Let $d<0$ be a squarefree integer and consider the imaginary quadratic number field $F:=\mathbf{Q}(\sqrt{d})$. Let $D$ denote the maximum distance from a point in $\mathbf{C}$ to a lattice point of $\mathcal{O}_{F}$ and let $M(F):=D^{2}$ be the Euclidean minimum (cf. Lem95, 3.1]). Then

$$
M(F)=\left\{\begin{array}{cl}
\frac{|d|+1}{4} & \text { if } d \equiv 2,3(\bmod 4) \\
\frac{(|d|+1)^{2}}{16|d|} & \text { if } d \equiv 1(\bmod 4)
\end{array}\right.
$$

Remark 1.13. In Chapter 3, Lemma 3.1, we will essentially prove this proposition in the case when $d \not \equiv 1(\bmod 4)$.

Table 1.1 shows the values of $M(\mathbf{Q}(\sqrt{d}))$ for the first six imaginary quadratic number fields, ordered by their discriminant.

Proposition 1.14 ([BT73, p.432]). A number field $F$ is Euclidean if $M(F)<1$. Hence, by Proposition 1.12, $\mathbf{Q}(\sqrt{-d})$ is Euclidean for

$$
d \in\{1,2,3,7,11\}
$$

### 1.3 The Gaussian integers

We consider the degree two extension $F:=\mathbf{Q}(i)$ of $\mathbf{Q}$, whose ring of integers $\mathcal{O}_{F}$ in $F$ is the Gaussian integers $\mathbf{Z}[i]$ (which follows from, e.g., Proposition 1.9).

Proposition 1.15. The ring of Gaussian integers is a principal ideal domain.
Proof. Note that the discriminant of $\mathbf{Q}(i) / \mathbf{Q}$ is $\delta_{\mathbf{Q}(i) / \mathbf{Q}}=-4$, and since there is exactly one pair of complex conjugate embeddings of $\mathbf{Q}(i)$ into $\mathbf{C}$, the Minkowski bound for the extension becomes

$$
M=\frac{2}{2^{2}} \cdot \frac{4}{\pi} \sqrt{\left|\delta_{\mathbf{Q}(i) / \mathbf{Q}}\right|}=\frac{4}{\pi} .
$$

By the Minkowski bound theorem ( Jan96, I Theorem 13.7]), any class in the ideal class group $\mathrm{Cl}(\mathbf{Z}[i])$ is represented by a nonzero integral ideal $\mathfrak{a}$ whose counting norm $\operatorname{Nm}(\mathfrak{a})$ is bounded by $M$. Since $M<2$, we must have that $\operatorname{Nm}(\mathfrak{a})=1$, i.e., $\mathfrak{a}=\mathbf{Z}[i]$, and thus the ideal class group is trivial.

Remark 1.16. The above result follows of course also from the fact that $\mathbf{Z}[i]$ is Euclidean.
Remark 1.17. The field $\mathbf{Q}(i)$ is one of the in total nine imaginary quadratic number fields with class number one; the other ones are $\mathbf{Q}(\sqrt{-d})$ for

$$
d \in\{2,3,7,11,19,43,67,167\}
$$

Proposition 1.18. We have $\mathbf{Z}[i]^{\times}=\mu(\mathbf{Q}(i))=\mu_{4}$.

Proof. This follows directly from Dirichlet's unit theorem ([Jan96, I Theorem 13.12]):

$$
\mathbf{Z}[i]^{\times} \cong \mu(\mathbf{Q}(i)) \times \mathbf{Z}^{r+c-1}=\mu(\mathbf{Q}(i))
$$

where $r$ and $c$ denotes respectively the number of real embeddings of $\mathbf{Q}(i)$, and the number of complex conjugate pairs of embeddings.

To introduce some geometry to the picture, let us take a look at the fibers over the closed points of $\operatorname{Spec} \mathbf{Z}$ of the morphism $X:=\operatorname{Spec} \mathbf{Z}[i] \rightarrow \operatorname{Spec} \mathbf{Z}$ induced by the inclusion $\mathbf{Z} \hookrightarrow \mathbf{Z}[i]$. Let $p$ be a rational prime, then the fiber

$$
\begin{aligned}
X_{(p)} & =\operatorname{Spec} \mathbf{Z}[i] \times_{\operatorname{Spec} \mathbf{Z}} \operatorname{Spec} \mathbf{F}_{p} \\
& =\operatorname{Spec}\left(\mathbf{Z}[x] /\left(x^{2}+1\right) \otimes_{\mathbf{Z}} \mathbf{F}_{p}\right) \\
& =\operatorname{Spec}\left(\mathbf{F}_{p}[x] /\left(x^{2}+1\right)\right) .
\end{aligned}
$$

Thus we see that determining the fibers $X_{(p)}$ is the same as determining the splitting of the prime $p$ in the extension $\mathbf{Q}(i) / \mathbf{Q}$. From Proposition 1.10 we have the following.

1. The prime 2 ramifies. Here we can also see at once that $\mathbf{F}_{2}[x] /\left(x^{2}+1\right)=\mathbf{F}_{2}[x] /(x+1)^{2}$, and by the correspondence $x \mapsto i$ we get that $2 \mathbf{Z}[i]=(2, i+1)^{2}=(i+1)^{2}$. Geometrically, the fiber is a nonreduced one-point scheme.
2. Assume $p \equiv 1(\bmod 4)$. By the first supplementary law of quadratic reciprocity we then obtain

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

Hence $p$ splits, and the fiber consists of two points.
3. If $p \equiv 3(\bmod 4)$ then $(-1 / p)=-1$, hence $p$ is inert in the extension and the fiber consists of one closed point.

Point 2 above illustrates a classical theorem of Fermat, namely that a rational prime $p$ is the sum of two squares if and only if $p \equiv 1(\bmod 4)$. Indeed, let $\mathrm{Nm}:=\mathrm{Nm}_{\mathbf{Q}(i) / \mathbf{Q}}: \mathbf{Z}[i] \rightarrow \mathbf{Z}$ be the norm map $a+i b \mapsto a^{2}+b^{2}$, so that $p$ is a sum of two squares if and only if $p$ lies in the image of the norm map. If $p$ splits, i.e., $p=\alpha \beta$ in $\mathbf{Z}[i]$, then $\operatorname{Nm}(p)=p^{2}=\operatorname{Nm}(\alpha) \operatorname{Nm}(\beta)$, and since we assume neither $\alpha$ nor $\beta$ is a unit we must have $p=\operatorname{Nm}(\alpha)=\operatorname{Nm}(\beta)$. Hence $p$ is a sum of two squares. But by the above discussion, $p$ splits if and only if $p \equiv 1(\bmod 4)$.

Let us summarize the results above.
Proposition 1.19. The nonzero prime ideals of $\mathbf{Z}[i]$ are of the following forms:

- $(1+i)$ lying above 2 .
- $(a+i b)$, which lies above the prime $p \equiv 1(\bmod 4)$ such that $p=a^{2}+b^{2}$.
- (p), for $p$ a rational prime congruent to 3 modulo 4.

By Proposition 1.11 we obtain the following result:
Proposition 1.20. Let $\mathfrak{p}$ be a nonzero prime ideal of $\mathbf{Z}[i]$ above the prime $p \in \mathbf{Z}$. Then the residue field $k(\mathfrak{p}):=\mathbf{Z}[i] / \mathfrak{p}$ is of the following form:

$$
k(\mathfrak{p}) \cong \mathbf{F}_{\mathrm{Nm} \mathfrak{p}}= \begin{cases}\mathbf{F}_{p} & \text { if } p=2 \text { or } p \equiv 1(\bmod 4) \\ \mathbf{F}_{p^{2}} \text { if } p \equiv 3(\bmod 4) .\end{cases}
$$

Figure 1.1 depicts the analogy with Riemann surfaces: One thinks of the morphism

$$
\operatorname{Spec} \mathbf{Z}[i] \longrightarrow \operatorname{Spec} \mathbf{Z}
$$

as a two sheeted covering, ramified only above (2). In general, given a number field $F \neq \mathbf{Q}$, let $\delta_{F}$ denote its discriminant. By Minkowski's theorem [Jan96, I Theorem 13.9], the extension


Figure 1.1: The closed points of $\operatorname{Spec} \mathbf{Z}[i]$, viewed as a two-sheeted covering of $\operatorname{Spec} \mathbf{Z}$.
$F / \mathbf{Q}$ is ramified, and the ramified prime ideals of $\mathcal{O}_{F}$ are exactly those lying above a rational prime $p$ dividing $\delta_{F}$. Let $S$ be the set consisting of the infinite places of $F$ together with the finite places corresponding to the ramified primes, and let

$$
\mathcal{O}_{F, S}:=\left\{x \in F: \operatorname{ord}_{v}(x) \geq 0 \text { for all } v \notin S\right\}
$$

denote the ring of $S$-integers. Then

$$
\operatorname{Spec} \mathcal{O}_{F, S}=\operatorname{Spec}\left(\mathcal{O}_{F}\left[1 / \delta_{F}\right]\right) \longrightarrow \operatorname{Spec} \mathbf{Z}
$$

is a finite étale covering. The fact that there are no unramified algebraic extensions of $\mathbf{Q}$ implies that $\operatorname{Spec} \mathbf{Z}$ is simply connected in the sense of étale fundamental groups: $\pi_{1}^{\text {ét }}(\operatorname{Spec} \mathbf{Z})=0$ Mor12b, p.36]. In this language, class field theory describes the maximal abelian quotient of the étale fundamental group of the arithmetic scheme $\operatorname{Spec} \mathcal{O}_{F}$ in terms of the arithmetic of the number field $F$ itself.

### 1.3.1 Local structure

In this subsection we will make use of a few lemmas whose statement we defer until Section 1.5.1 after general notation has been introduced. The reader may therefore skip this at the first reading and rather return when needed.

Later on we will need a few results on the structure of the local field $\mathbf{Q}(i)_{v}$ for $v$ a place of $\mathbf{Q}(i)$-in particular the place corresponding to the prime $1+i$. First of all, recall that if $v$ is a finite place of $\mathbf{Q}(i)$ lying above a rational prime $p$, then $\mathbf{Q}(i)_{v}=\mathbf{Q}_{p}(i)$ (see Gra03, p.13]). In particular, $\mathbf{Q}(i)_{(1+i)}=\mathbf{Q}_{2}(i)$ and thus $\mathbf{Z}[i]_{(1+i)}=\mathbf{Z}_{2}[i]$. Let $U_{1+i}:=\mathbf{Z}_{2}[i]^{\times}$, and let $U_{1+i}^{n}:=1+(1+i)^{n} \mathbf{Z}_{2}[i]$ denote the higher unit groups (cf. Definition 1.28).

Lemma 1.21. If $\alpha$ is a Gaussian integer relatively prime to $1+i$, then $\alpha$ has a primary associate. In other words, $\alpha$ becomes primary after multiplication by a suitable power of $i$.

Proof. Since we assume $\alpha$ is coprime to $1+i, \bar{\alpha}$ is a unit in $\mathbf{Z}[i] /(1+i)^{3}$. Now

$$
\#\left(\mathbf{Z}[i] /(1+i)^{3}\right)^{\times}=\varphi\left((1+i)^{3}\right)=\operatorname{Nm}(1+i)^{2}(\operatorname{Nm}(1+i)-1)=4
$$

where $\varphi$ is the Euler totient function on $\mathbf{Z}[i]$. Thus there are 4 units, and we can take $\mu_{4}$ as a system of representatives for $\left(\mathbf{Z}[i] /(1+i)^{3}\right)^{\times}$since the powers of $i$ are different modulo $(1+i)^{3}$. Hence $\bar{\alpha}=\overline{i^{k}}$ for some $k \in\{0,1,2,3\}$, and the statement follows.

Proposition 1.22. We have $\mathbf{Q}_{2}(i)^{\times} \cong(1+i)^{\mathbf{Z}} \oplus \mu_{4} \oplus U_{1+i}^{3}$. Consequently $U_{1+i} \cong \mu_{4} \oplus U_{1+i}^{3}$.
Proof. Let $v$ denote the place corresponding to $(1+i)$. Note that any element $\alpha \in U_{v}$ can be written as a power series

$$
\alpha=a_{0}+a_{1}(1+i)+a_{2}(1+i)^{2}+\cdots
$$

where $a_{n} \in\{0,1\}$ and $a_{0} \neq 0$. Hence $U_{v}=U_{v}^{1}$.
By Proposition 1.35 it is enough to show that $U_{v}^{1} \cong \mu_{4} \oplus U_{v}^{3}$. Since

$$
\frac{\mathbf{Z}_{2}[i]}{(1+i)^{3} \mathbf{Z}_{2}[i]} \cong \frac{\mathbf{Z}[i]}{(1+i)^{3}},
$$

the result of Lemma 1.21 holds also for $U_{v}^{1}$, hence any $\alpha \in U_{v}^{1}$ can be written as

$$
\alpha=i^{4-k}\left(i^{k} \alpha\right) \in \mu_{4} \oplus U_{v}^{3}
$$

for some $k$. Since the powers of $i$ are different modulo $(1+i)^{3}$ we have $U_{v}^{3} \cap \mu_{4}=1$, yielding $U_{v}^{1} \cong \mu_{4} \oplus U_{v}^{3}$.

Proposition 1.23. Let $\left(U_{1+i}\right)^{4}$ be the group of fourth powers in $\mathbf{Z}_{2}[i]^{\times}$. Then

1. the group $U_{1+i} /\left(U_{1+i}\right)^{4}$ has order 64 ;
2. $\left(U_{1+i}\right)^{4}=U_{1+i}^{7}$;
3. $U_{1+i} /\left(U_{1+i}\right)^{4} \cong \mu_{4} \oplus\langle\overline{3+2 i}\rangle \oplus\langle\overline{5}\rangle$.

Proof. To ease the notation let $v$ denote the place corresponding to $(1+i)$ and put $\pi:=1+i$.

1. Since $4=-\pi^{4}$ we have $\operatorname{ord}_{v}(4)=4$, so Lemma 1.34 yields

$$
\left(U_{v}:\left(U_{v}\right)^{4}\right)=\frac{\# \mu_{4}\left(\mathbf{Q}_{2}(i)\right)}{|4|_{v}}=\frac{4}{\mathrm{Nm}(\pi)^{-\operatorname{ord}_{v}(4)}}=2^{6}
$$

2. To show that $\left(U_{v}\right)^{4} \subseteq U_{v}^{7}$, take an $\alpha \in U_{v}$. Note that by the proof of Proposition 1.22 , we may assume $i^{k} \alpha$ is primary for some $k \in\{0,1,2,3\}$, say

$$
i^{k} \alpha=1+\pi^{3} \beta
$$

for some $\beta \in \mathbf{Z}_{2}[i]$. Then

$$
\begin{aligned}
\alpha^{4}=\left(i^{k} \alpha\right)^{4} & =1+4 \beta \pi^{3}+6 \beta^{2} \pi^{6}+4 \beta^{3} \pi^{9}+\beta^{4} \pi^{12} \\
& \equiv 1+4 \beta \pi^{3}+6 \beta^{2} \pi^{6}\left(\bmod \pi^{7}\right)
\end{aligned}
$$

Writing $4=-\pi^{4}$ and $2=(-i) \pi^{2}$ we get

$$
\alpha^{4} \equiv 1-\beta \pi^{7}-3 i \beta^{2} \pi^{8} \equiv 1\left(\bmod \pi^{7}\right)
$$

Hence $\left(U_{v}\right)^{4} \subseteq U_{v}^{7}$. Since $\left(U_{v}: U_{v}^{7}\right)=q_{v}^{6}\left(q_{v}-1\right)=2^{6}$ by Lemma 1.33 below (and the notations used there), we have by 1 that $\left(U_{v}: U_{v}^{7}\right)=\left(U_{v}:\left(U_{v}\right)^{4}\right)$, hence $\left(U_{v}\right)^{4}=U_{v}^{7}$.
3. Again by Lemma 1.33 we have

$$
\frac{U_{v}^{1}}{U_{v}^{7}} \cong\left(\frac{\mathbf{Z}_{2}[i]}{\pi^{7} \mathbf{Z}_{2}[i]}\right)^{\times} \cong\left(\frac{\mathbf{Z}[i]}{\pi^{7} \mathbf{Z}[i]}\right)^{\times}
$$

Choosing $3+2 i \in U_{v}^{3}, 5 \in U_{v}^{4}$ and $i \in U_{v}$ as generators for $\left(\mathbf{Z}[i] / \pi^{7}\right)^{\times}$yields the result.

### 1.4 The Eisenstein integers

Now we turn to the case $F=\mathbf{Q}(\sqrt{-3})$. Since $-3 \equiv 1(\bmod 4)$, the ring of integers $\mathcal{O}_{F}$ equals $\mathbf{Z}[\omega]$, where $\omega:=(-1+\sqrt{-3}) / 2$. The elements of $\mathbf{Z}[\omega]$ are called the Eisenstein integers. Note that $\omega$ is a primitive third root of unity, and that $\mathbf{Z}[\omega]=\mathbf{Z}\left[\zeta_{6}\right]$, where $\zeta_{6}$ is a primitive sixth root of unity. The discriminant of the extension $\mathbf{Q}(\sqrt{-3}) / \mathbf{Q}$ equals -3 , so the same reasoning as in Proposition 1.15 shows that $\mathbf{Z}[\omega]$ is a principal ideal domain.


Figure 1.2: Part of the lattice of integers in $\mathbf{Q}(\sqrt{-3})$.

Proposition 1.24. The ring of Eisenstein integers is a principal ideal domain, and the units are

$$
\mathbf{Z}[\omega]^{\times}=\mu_{6}=\left\{ \pm 1, \pm \omega, \pm \omega^{2}\right\}
$$

Proposition 1.25. The norm of an Eisenstein integer $a+b \omega \in \mathbf{Z}[\omega]$ is given by $\operatorname{Nm}(a+b \omega)=$ $a^{2}-a b+b^{2}$.

Proof. Let $\alpha:=a+b \omega \in \mathbf{Z}[\omega]$ be an Eisenstein integer. If $b=0$ we have $\operatorname{Nm}(\alpha)=\operatorname{Nm}(a)=a^{2}$, so the statement is true in this case. If $b \neq 0$, a straightforward computation shows that the minimal polynomial over $\mathbf{Q}$ of $\alpha$ is

$$
x^{2}-(2 a-b) x+\left(a^{2}-a b+b^{2}\right),
$$

and the norm $\operatorname{Nm}(\alpha)$ is the constant term of the minimal polynomial.
Proposition 1.26. Consider the extension $\mathbf{Q}(\sqrt{-3}) / \mathbf{Q}$. If $p \in \mathbf{Z}$ is a rational prime, then

- $p$ ramifies $\Longleftrightarrow p=3$;
- $p$ splits $\Longleftrightarrow p \equiv 1(\bmod 3)$;
- $p$ is inert $\Longleftrightarrow p \equiv 2(\bmod 3)$.

Thus, the nonzero prime ideals of $\mathbf{Z}[\omega]$ are of the following forms:

- $(2+\omega)$ lying above $3 \in \mathbf{Z}$;
- $(a+b \omega)$ lying above the rational prime $p \equiv 1(\bmod 3)$ such that $p=a^{2}-a b+b^{2}$;
- $(p)$, which lies above the prime $p \in \mathbf{Z}$ such that $p \equiv 2(\bmod 3)$.

Hence the Diophantine equation $x^{2}-x y+y^{2}=p$ has a solution if and only if $p=3$ or $p \equiv 1$ $(\bmod 3)$.

Proof. By Proposition 1.10 , only 3 ramifies. By Proposition 1.25 we have

$$
\operatorname{Nm}(2+\omega)=3
$$

so the ideal $(2+\omega)$ lies above $(3)$, hence $3 \mathbf{Z}[\omega]=(2+\omega)^{2}$.
Furthermore, we know by Proposition 1.10 that the prime 2 is inert. To show the rest of the statement for $p \neq 2$, all we need is quadratic reciprocity. By Proposition 1.10, a rational odd prime $p \neq 3$ splits in the extension if and only if $(-3 / p)=1$. By multiplicativity of Legendre symbols,

$$
\left(\frac{-3}{p}\right)=\left(\frac{-1}{p}\right)\left(\frac{3}{p}\right)=(-1)^{\frac{p-1}{2}}\left(\frac{3}{p}\right) .
$$

We consider two cases.


- If $p \equiv 1(\bmod 4),(3 / p)=(p / 3)$ by quadratic reciprocity. Hence

$$
\left(\frac{-3}{p}\right) \equiv(-1)^{\frac{p-1}{2}} p \equiv p \quad(\bmod 3)
$$

i.e., $(-3 / p)=1$ if and only if $p \equiv 1(\bmod 3)$.

- If $p \equiv 3(\bmod 4)$ we have $(3 / p)=-(p / 3)$, so

$$
\left(\frac{-3}{p}\right) \equiv-(-1)^{\frac{p-1}{2}} p \equiv p \quad(\bmod 3)
$$

so also in this case we get that $(-3 / p)=1$ if and only if $p \equiv 1(\bmod 3)$.
Thus $p$ splits if and only if $p \equiv 1(\bmod 3)$.
Proposition 1.27. Let $\mathfrak{p} \in \operatorname{Spec} \mathbf{Z}[\omega]$ be a nonzero prime ideal. Then the residue field $k(\mathfrak{p})$ is given by

$$
k(\mathfrak{p}) \cong \begin{cases}\mathbf{F}_{p} & \text { if } p=3 \text { or } p \equiv 1(\bmod 3) \\ \mathbf{F}_{p^{2}} & \text { if } p \equiv 2(\bmod 3)\end{cases}
$$

### 1.5 Class field theory

We state some of the main results of local and global class field theory. From here we introduce the notion of symbols over number fields, which serves as one of the first links between number theory and algebraic $K$-theory. This section will also help to fix notations, of which we mostly follow Gra03].

Let $F$ be a global number field, and let $\mathrm{Pl}_{F}$ denote the set of places of $F$. By Ostrowski's theorem Gra03, I Theorem 1.2], $\mathrm{Pl}_{F}$ decomposes as a disjoint union

$$
\mathrm{Pl}_{F}=\mathrm{Pl}_{0} \cup \mathrm{Pl}_{\infty}
$$

of finite and infinite places, respectively. Occasionally we will be interested in considering only noncomplex places of $F$, the set of which will be denoted by $\mathrm{Pl}_{F}^{\text {nc }}$.

Definition 1.28. Let $v$ be a place of $F$. We will denote by:

- $F_{v}$ the completion of $F$ at $v$, with valuation ring $\mathcal{O}_{v}$;
- $\pi_{v}$ a uniformizer of $F_{v}$, i.e., a generator of the maximal ideal $\mathfrak{p}_{v}$ of $\mathcal{O}_{v}$ (for $v$ finite);
- $k(v):=\mathcal{O}_{v} / \mathfrak{p}_{v}$ the residue field (for $v$ finite);
- $q_{v}:=\# k(v)$ (for $v$ finite);
- $U_{v}:=\mathcal{O}_{v}^{\times}$the unit group, where we set $U_{v}:=\mathbf{R}_{>0}$ when $v$ is infinite real;
- $U_{v}^{n}:=1+\mathfrak{p}_{v}^{n}(n \geq 1)$ the higher unit groups (for $v$ finite);
- $\mu(F), \mu\left(F_{v}\right)$ the groups of roots of unity of respectively $F, F_{v}$;
- $m:=\# \mu(F), m_{v}:=\# \mu\left(F_{v}\right)$ (for $v$ noncomplex);
- $i_{v}: F \hookrightarrow F_{v}$ the embedding of $F$ into its completion $F_{v}$ at $v$.

For each place $v \in \mathrm{Pl}_{F}$ we shall make use of the corresponding valuation on $F^{\times}$. It is clear that when $v$ is a finite place, the corresponding valuation is just the $\mathfrak{p}_{v}$-adic valuation. However, it is convenient to also have a similar notion when $v$ is infinite:

Definition 1.29. Let $v \in \mathrm{Pl}$ be a place of $F$. We denote the corresponding valuation by $\operatorname{ord}_{v}$, and it takes the following forms:

- If $v \in \mathrm{Pl}_{0}$ is a finite place, ord $\cos _{v}: F^{\times} \rightarrow \mathbf{Z}$ is the $\mathfrak{p}_{v}$-adic valuation.
- If $v$ is an infinite real place, the valuation $\operatorname{ord}_{v}: F^{\times} \rightarrow \mathbf{Z} / 2$ is defined by:

$$
\operatorname{ord}_{v}(x):= \begin{cases}0, & i_{v}(x)>0 \\ 1, & i_{v}(x)<0\end{cases}
$$

- If $v$ is an infinite complex place, we set $\operatorname{ord}_{v}:=0$.


### 1.5.1 On the higher unit groups

The higher unit groups $U_{v}^{n}$ of the local field $F_{v}$ form a filtration

$$
U_{v} \supseteq U_{v}^{1} \supseteq \cdots \supseteq U_{v}^{n} \supseteq \cdots
$$

of the unit group $U_{v}$, whose filtration coefficients $U_{v}^{n} / U_{v}^{n+1}$ are all finite. They allow us to $\operatorname{express} U_{v}$ as the limit $U_{v}=\lim _{\ddagger} U_{v} / U_{v}^{n}$, hence $U_{v}$ is a profinite group.

Proposition 1.30. The group $U_{v}$ contains the $\left(q_{v}-1\right)$-roots of unity $\mu_{q_{v}-1}$.
Proof. By the general Fermat's little theorem, the polynomial

$$
f(x):=x^{q_{v}-1}-1
$$

when reduced modulo $\mathfrak{p}_{v}$ has $q_{v}-1$ distinct roots in $k(v)^{\times} \cong \mu_{q_{v}-1}$. By Hensel's lemma, these roots lift uniquely to $q_{v}-1$ distinct roots of $f$ in $U_{v}$.

Remark 1.31. The map $k(v)^{\times} \hookrightarrow U_{v}$ which sends an element to its unique lift in $U_{v}$ provided by Hensel's lemma is called the Teichmüller lift.

Lemma 1.32 (Ser79, Proposition 6, p.66]). For any finite place $v \in \mathrm{Pl}_{F}$ we have a canonical isomorphism $U_{v} / U_{v}^{1} \cong \mu_{q_{v}-1} \cong k(v)^{\times}$, and for each $n \geq 1$ a noncanonical isomorphism $U_{v}^{n} / U_{v}^{n+1} \cong k(v)$.

Lemma 1.33. Given $v \in \mathrm{Pl}_{0}$ and any $n \geq 1$, let $\varphi$ be the generalized Euler function, where $\varphi\left(\mathfrak{p}_{v}^{n}\right)=q_{v}^{n-1}\left(q_{v}-1\right)$. Then $\left(U_{v}: U_{v}^{n}\right)=\varphi\left(\mathfrak{p}_{v}^{n}\right)$, and

$$
U_{v} / U_{v}^{n} \cong\left(\mathcal{O}_{v} / \mathfrak{p}_{v}^{n}\right)^{\times} .
$$

Proof. Note that we have exact sequences

$$
1 \longrightarrow U_{v}^{n} / U_{v}^{n+1} \longrightarrow U_{v} / U_{v}^{n+1} \longrightarrow U_{v} / U_{v}^{n} \longrightarrow 1
$$

for all $n$, i.e., $\left(U_{v}: U_{v}^{n+1}\right)=\left(U_{v}: U_{v}^{n}\right)\left(U_{v}^{n}: U_{v}^{n+1}\right)$. When $n=1$, this reads $\left(U_{v}: U_{v}^{2}\right)=$ $q_{v}\left(q_{v}-1\right)$ by Lemma 1.32 the general case follows by induction.

The second statement follows from the fact that the canonical surjective homomorphism

$$
\begin{aligned}
U_{v} & \longrightarrow\left(\mathcal{O}_{v} / \mathfrak{p}_{v}^{n}\right)^{\times} \\
u & \longmapsto u\left(\bmod \mathfrak{p}_{v}^{n}\right)
\end{aligned}
$$

has kernel $U_{v}^{n}$.

Lemma 1.34 ([Lan70, Proposition 6, p.47]). Let $v$ be a finite place of $F$. For any $n \geq 1$, let $\left(U_{v}\right)^{n}$ denote the group of $n$-th powers in $U_{v}$. Then

$$
\left(U_{v}:\left(U_{v}\right)^{n}\right)=\frac{\# \mu_{n}\left(F_{v}\right)}{|n|_{v}}
$$

where $\mu_{n}\left(F_{v}\right)$ is the group of $n$-th roots of unity contained in $F_{v}$, and $|\cdot|_{v}$ is the $v$-adic absolute value on $F_{v}$.

Proposition 1.35. For $v$ a finite place of $F$, there is a decomposition

$$
F_{v}^{\times} \cong \pi_{v}^{\mathbf{Z}} \oplus \mu_{q_{v}-1} \oplus U_{v}^{1}
$$

Proof. Any $x \in F_{v}^{\times}$can be written $x=\pi_{v}^{\operatorname{ord}_{v}(x)} u$ with $u \in U_{v}$. Since $U_{v} / U_{v}^{1} \cong \mu_{q_{v}-1}$ we have an exact sequence

$$
1 \longrightarrow U_{v}^{1} \longrightarrow U_{v} \longrightarrow \mu_{q_{v}-1} \longrightarrow 1
$$

which splits by the Teichmüller lift $\mu_{q_{v}-1} \cong k(v)^{\times} \hookrightarrow U_{v}$.
Example 1.36. Let us show that any element of $U_{p}^{1}=1+p \mathbf{Z}_{p}$ is in fact a square whenever $p$ is an odd rational prime. Indeed, for any $x \in 1+p \mathbf{Z}_{p}$, write $x=1+p y$ where $y \in \mathbf{Z}_{p}$. Then we have the series expansion

$$
x^{1 / 2}=(1+p y)^{1 / 2}=1+\frac{y}{2} p-\frac{y^{2}}{8} p^{2}+\frac{y^{3}}{16} p^{3}-\cdots
$$

which converges $p$-adically. The denominators are all powers of 2 , hence invertible in $\mathbf{Z}_{p}$, so the series converges to an element in $1+p \mathbf{Z}_{p}$. Thus $x$ has a square root in $1+p \mathbf{Z}_{p}$. From this and Proposition 1.35 above we can conclude for example that $\mathbf{Z}_{p}^{\times}$has only one subgroup of index 2 , namely $\mu_{p-1}^{2} \oplus\left(1+p \mathbf{Z}_{p}\right)=\left(\mathbf{Z}_{p}^{\times}\right)^{2}$.

### 1.5.2 Frobenius

Recall that if $L / F$ is an extension of number fields and $w$ a place of $L$ above a finite place $v$ of $F$, there is a canonical isomorphism $\operatorname{Gal}\left(L_{w} / F_{v}\right) \cong D_{w}$ [Mil13a, Proposition 8.10] where $D_{w}$ is the decomposition group of $w$ :

$$
D_{w}:=\{\sigma \in \operatorname{Gal}(L / F): \sigma w=w\}
$$

where the action of $\operatorname{Gal}(L / F)$ on $\mathrm{Pl}_{L}$ is defined by

$$
|x|_{\sigma w^{\prime}}:=\left|\sigma^{-1}(x)\right|_{w^{\prime}}
$$

for any $w^{\prime} \in \mathrm{Pl}_{L}$. There is a surjective homomorphism $D_{w} \rightarrow \operatorname{Gal}(k(w) / k(v))$ defined by $\sigma \mapsto \bar{\sigma}$, where

$$
\bar{\sigma}\left(y+\mathfrak{p}_{w}\right):=\sigma(y)+\mathfrak{p}_{w}
$$

Its kernel is called the inertia group of $w$, and is denoted by $I_{w}$. If $L_{w} / F_{v}$ is unramified, the inertia group is trivial and thus we have an isomorphism $\operatorname{Gal}\left(L_{w} / F_{v}\right) \cong \operatorname{Gal}(k(w) / k(v))$ in this case. Now $\operatorname{Gal}(k(w) / k(v))$ is a finite cyclic group, with a canonical generator which sends $\bar{y} \in k(w)$ to $\bar{y}^{q_{v}}$.

Definition 1.37. Suppose $L_{w} / F_{v}$ is an unramified extension of local fields. The image in $\operatorname{Gal}\left(L_{w} / F_{v}\right)$ of the canonical generator of $\operatorname{Gal}(k(w) / k(v))$ under the isomorphism $\operatorname{Gal}\left(L_{w} / F_{v}\right) \cong$ $\operatorname{Gal}(k(w) / k(v))$ is called the local Frobenius, and is denoted by

$$
\left(L_{w} / F_{v}\right)
$$

The image of $\left(L_{w} / F_{v}\right)$ in $\operatorname{Gal}(L / F)$ under the isomorphism $\operatorname{Gal}\left(L_{w} / F_{v}\right) \cong D_{w} \subseteq \operatorname{Gal}(L / F)$ is called the global Frobenius at $w$, and is denoted by

$$
\left(\frac{L / F}{w}\right)
$$

Remark 1.38. If $w$ and $w^{\prime}$ are two places above $v$, then the global Frobenius at $w$ and $w^{\prime}$ are conjugate to each other [Jan96, p.126]. Hence they are equal if $L / F$ is an abelian extension, i.e., if $\operatorname{Gal}(L / F)$ is abelian. In this case we will therefore denote by $\left(\frac{L / F}{v}\right)$ the Frobenius at any place above $v$.

### 1.5.3 Local and global class field theory

All the number theoretic symbols we shall consider in Section 1.6 are offspring of the local reciprocity map, which will be introduced shortly in the discussion on local class field theory (see Theorem 1.41). The treatment of global class field theory will be idèlic, i.e., we define the global reciprocity map on the idèle group of the global field in consideration, as opposed to the more classical ideal-theoretic definition. The idèle group, first introduced by Chevalley Mil13b, p.viii], is an object attached to a global field allowing for a unified treatment of the different embeddings of the field into its completion at all places, including the infinite ones. We start off by a brief introduction to the idèles before we move on to the local and global reciprocity maps.

The restricted product $\prod_{i \in I} G_{i}$ of a family of locally compact topological groups $\left\{G_{i}\right\}_{i \in I}$ is by definition the subset of $\prod_{i \in I} G_{i}$ consisting of all elements $\left(g_{i}\right)_{i \in I}$ for which $g_{i} \in K_{i}$ for almost all $i$, where $K_{i}$ is some fixed compact subgroup of $G_{i}$. The restricted product is then a locally compact group, and we provide it with a topology whose basis elements are of the form $\prod_{i} A_{i}$ where $A_{i} \subseteq G_{i}$ is open in $G_{i}$ for all $i$ and $A_{i}=K_{i}$ for almost all $i$.

Definition 1.39. The idèle group of the number field $F$ is the restricted product of the groups $F_{v}^{\times}$with respect to the compact subgroups $U_{v}$ :

$$
\mathbf{J}_{F}:=\coprod_{v \in \mathrm{Pl}_{F}} F_{v}^{\times}=\left\{\left(x_{v}\right)_{v \in \mathrm{Pl}_{F}}: x_{v} \in F_{v}^{\times}, x_{v} \in U_{v} \text { for almost all } v\right\} .
$$

An element $\mathbf{x}$ of $\mathbf{J}_{F}$ is called an idèle. We define the idèle class group $C_{F}$ as

$$
C_{F}:=\mathbf{J}_{F} / i\left(F^{\times}\right),
$$

where $i: F^{\times} \rightarrow \prod_{v \in \mathrm{Pl}} F_{v}^{\times}$is the injective map sending $x \in F^{\times}$to $\left(i_{v}(x)\right)_{v}$.
We sometimes abbreviate $\mathbf{J}_{F}$ and $C_{F}$ to $\mathbf{J}$ and $C$. By abuse of notation we may also write simply $F^{\times}$for the image $i\left(F^{\times}\right)$when it is clear from context that we consider the subset of $\mathbf{J}_{F}$ and not of $F$.

As in the general case mentioned above, the idèles $\mathbf{J}_{F}$ are endowed with a topology whose basis elements are of the form

$$
\prod_{v \in \mathrm{Pl}} V_{v},
$$

where $V_{v}$ is open in $F_{v}^{\times}$for all $v$ and $V_{v}=U_{v}$ for almost all $v$. With this topology, the image of $F^{\times}$in $\mathbf{J}$ under the diagonal embedding $i$ becomes discrete.

Remark 1.40. Alternatively, one may define $\mathbf{J}_{F}$ as the colimit

$$
\mathbf{J}_{F}=\underset{S \subseteq{\underset{\mathrm{P}}{F}}^{\lim }}{\mathbf{I}_{F}} \mathbf{J}_{F}(S)
$$

where $S$ varies over all finite subsets of $\mathrm{Pl}_{F}$ containing the infinite places, and

$$
\mathbf{J}_{F}(S):=\prod_{v \in S} F_{v}^{\times} \times \prod_{v \notin S} U_{v}
$$

Yet another definition is to put $\mathbf{J}_{F}:=\mathbf{A}_{F}^{\times}$, where $\mathbf{A}_{F}$ is the adèle ring of $F$. However, due to the multiplicative nature of class field theory, we shall not need to introduce adèles for our purposes.

For any idèle $\mathbf{x}=\left(x_{v}\right)_{v} \in \mathbf{J}$, define the volume of $\mathbf{x}$ as

$$
|\mathbf{x}|:=\prod_{v}\left|x_{v}\right|_{v}
$$

where $|\cdot|_{v}$ is the absolute value on $F_{v}$. The volume map $|\cdot|: \mathbf{J}_{F} \rightarrow \mathbf{R}_{>0}$ is then continuous. If we let $\mathbf{J}^{1}$ denote the kernel of the volume map, i.e., $\mathbf{J}^{1}:=\{\mathbf{x} \in \mathbf{J}:|\mathbf{x}|=1\}$, then $F^{\times}$is contained in $\mathbf{J}^{1}$ by the product formula for absolute values [Jan96, II Theorem 6.4] and the quotient $\mathbf{J}^{1} / F^{\times}$is compact Lan70, Theorem 4, p.142]. This compactness result is an equivalent formulation of the finiteness of the ideal class group and Dirichlet's unit theorem.

For the remainder of this section we assume that $L / F$ is an abelian extension of $F$, i.e., that $\operatorname{Gal}(L / F)$ is abelian. Let $v$ be a place of $F$, and $w \mid v$ a place of $L$ above $v$.
Theorem 1.41 (Local Reciprocity Law [Gra03, II Theorem 1.4]). There exists a unique canonical homomorphism

$$
\left(-, L_{w} / F_{v}\right): F_{v}^{\times} \longrightarrow \operatorname{Gal}\left(L_{w} / F_{v}\right)
$$

called the local reciprocity map, or the local Artin map, satisfying the following properties.

1. The kernel of $\left(-, L_{w} / F_{v}\right)$ is equal to the norm group

$$
\operatorname{Nm}_{L_{w} / F_{v}}\left(L_{w}^{\times}\right)=\left\{\operatorname{Nm}_{L_{w} / F_{v}}(x): x \in L_{w}^{\times}\right\}
$$

and we have an exact sequence

$$
1 \longrightarrow \operatorname{Nm}_{L_{w} / F_{v}}\left(L_{w}^{\times}\right) \longrightarrow F_{v}^{\times} \xrightarrow{\left(-, L_{w} / F_{v}\right)} \operatorname{Gal}\left(L_{w} / F_{v}\right) \longrightarrow 1
$$

2. If $L_{w} / F_{v}$ is unramified, then

$$
\left(x, L_{w} / F_{v}\right)=\left(L_{w} / F_{v}\right)^{\operatorname{ord}_{v}(x)}
$$

for all $x \in F_{v}^{\times}$.
3. The image of the unit group $U_{v}$ under $\left(-, L_{w} / F_{v}\right)$ is equal to the inertia group $I_{w}$, and we have the exact sequence

$$
1 \longrightarrow \mathrm{Nm}_{L_{w} / F_{v}}\left(U_{w}\right) \longrightarrow U_{v} \longrightarrow I_{w} \longrightarrow 1
$$

Theorem 1.42 ( Gra03, II Theorem 1.5]). For any subgroup $N$ of $F_{v}^{\times}$of finite index there exists a finite abelian extension $M$ of $F_{v}$ such that $\mathrm{Nm}_{M / F_{v}}\left(M^{\times}\right)=N$. The map $M \mapsto \operatorname{Nm}_{M / F_{v}}\left(M^{\times}\right)$ is a bijection between the set of finite abelian extensions of $F_{v}$ and subgroups of $F_{v}$ of finite index. This bijection has the following properties, where the field extensions $M_{1}, M_{2}$ correspond respectively to the groups $N_{1}$ and $N_{2}$ :

1. $M_{1} \subseteq M_{2}$ if and only if $N_{1} \supseteq N_{2}$;
2. $M_{1} M_{2}$ corresponds to $N_{1} \cap N_{2}$;
3. $M_{1} \cap M_{2}$ corresponds to $N_{1} N_{2}$;
4. if $M_{1} \subseteq M_{2}$ then $\operatorname{Gal}\left(M_{2} / M_{1}\right) \cong N_{1} / N_{2}$.

Example 1.43 (Mil13b, 1.6, p.22], Gra03, II Remark 1.4.6]). Suppose $F=\mathbf{Q}, L=\mathbf{Q}(i)$, and let $v$ be the infinite real place of $F$, so that $F_{v}=\mathbf{R}$. If $w$ is a place of $L$ above $v$, then $L_{w}=\mathbf{C}$. The local reciprocity map $\left(-, L_{w} / F_{v}\right)=(-, \mathbf{C} / \mathbf{R})$ is then given by

$$
(x, \mathbf{C} / \mathbf{R})=c^{\operatorname{ord}_{v}(x)}
$$

for $x \in \mathbf{R}^{\times}$, where $c \in \operatorname{Gal}(\mathbf{C} / \mathbf{R})$ is complex conjugation. Note that

$$
\mathrm{Nm}_{\mathbf{C} / \mathbf{R}}\left(\mathbf{C}^{\times}\right)=\mathbf{R}_{>0}=\operatorname{ker}(-, \mathbf{C} / \mathbf{R})
$$

(remember Definition 1.29). Finally, any subgroup $N$ of $\mathbf{R}^{\times}$of finite index is equal to either $\mathbf{R}^{\times}$or $\mathbf{R}_{>0}$. Indeed, $N$ contains $\mathbf{R}^{\times n}$ for some $n$, and by the intermediate value theorem one has $\mathbf{R}^{\times n}=\mathbf{R}^{\times}$if $n$ is odd; $\mathbf{R}^{\times n}=\mathbf{R}_{>0}$ if $n$ is even.

The local compactness of the idèle group makes it possible to define a global reciprocity map whose components are the local reciprocity maps:

Definition 1.44. The global reciprocity map-or global Artin map-is defined for the extension $L / F$ as the map

$$
\begin{aligned}
& \rho_{L / F}: \mathbf{J}_{F} \longrightarrow \operatorname{Gal}(L / F) \\
& \mathbf{x}=\left(x_{v}\right)_{v} \longmapsto \prod_{v \in \mathrm{Pl}_{F}}\left(\frac{x_{v}, L / F}{v}\right),
\end{aligned}
$$

where $\left(\frac{x_{v}, L / F}{v}\right)$ is the image of $\left(x_{v}, L_{v} / F_{v}\right) \in \operatorname{Gal}\left(L_{v} / F_{v}\right)$ in $\operatorname{Gal}(L / F)$ under the isomorphism $\operatorname{Gal}\left(L_{v} / F_{v}\right) \cong D_{v} \subseteq \operatorname{Gal}(L / F)$.

Note that for an idèle $\left(x_{v}\right)_{v} \in \mathbf{J}$, since $x_{v} \in U_{v}$ for almost all $v$ and only finitely many places $v$ are ramified in $L / F$, property 2 of Theorem 1.41 tells us that $\left(\frac{x_{v}, L / F}{v}\right)=1$ for almost all $v$, so $\rho_{L / F}$ is well defined.
Definition 1.45. Given any place $v \in \mathrm{Pl}_{F}$, the Hasse symbols $\left(\frac{-, L / F}{v}\right)$ are defined on $F^{\times}$as:

$$
\begin{aligned}
\left(\frac{-, L / F}{v}\right): F^{\times} & \longrightarrow \operatorname{Gal}(L / F) \\
x & \longmapsto\left(\frac{i_{v}(x), L / F}{v}\right) .
\end{aligned}
$$

The following theorem - which states in terms of Hasse symbols that the diagonal embedding $i\left(F^{\times}\right)$of $F^{\times}$in $\mathbf{J}_{F}$ is contained in the kernel of the global reciprocity map-is one of the main results of class field theory, and is considered as the generalization of the quadratic reciprocity law of Gauss to arbitrary abelian extensions of number fields. This constitutes a partial solution to Hilbert's ninth problem.
Theorem 1.46 (Gra03, II Theorem 3.4.1]). Let $L / F$ be a finite abelian extension of number fields. For any $x \in F^{\times}$, the following product formula holds:

$$
\prod_{v \in \mathrm{Pl}_{F}}\left(\frac{x, L / F}{v}\right)=1
$$

Example 1.47 (Gra03, II Example 3.4.2]). We calculate some Hasse symbols and illustrate how Gauss' quadratic reciprocity law is obtained from the product formula. Let $p$ be an odd prime number and take $F=\mathbf{Q}$. We consider the extension $L=\mathbf{Q}\left(\sqrt{p^{*}}\right)$ where $p^{*}:=$ $(-1)^{(p-1) / 2} p$. Write $\operatorname{Gal}(L / \mathbf{Q})=\{ \pm 1\}$. We shall compute the Hasse symbols

$$
\left(\frac{-1, L / \mathbf{Q}}{v}\right) \text { and }\left(\frac{q, L / \mathbf{Q}}{v}\right)
$$

where $q$ is a rational prime and $v$ runs through all places of $\mathbf{Q}$. Now there is exactly one archimedean place of $\mathbf{Q}$, which we denote by $\infty$, and otherwise all the nonarchimedean places correspond to prime numbers $\ell \in \mathbf{Q}$.

First of all, we have that the extension $L / \mathbf{Q}$ is ramified only at $p$. Indeed, since $p^{*} \equiv$ $1(\bmod 4)$, we know from the proof of Proposition 1.10 that the discriminant $\delta_{L / \mathbf{Q}}=p^{*}$. Thus, only $p$ ramifies in $L / \mathbf{Q}$.

By Theorem 1.41 (2) we know that if we are given an $x \in \mathbf{Q}^{\times}$and a place $v$ of $\mathbf{Q}$ unramified in $L / \mathbf{Q}$ which satisfies $\operatorname{ord}_{v}(x)=0$, then $\left(\frac{x, L / \mathbf{Q}}{v}\right)=1$.

- Consider first $\left(\frac{-1, L / \mathbf{Q}}{v}\right)$. We have $\operatorname{ord}_{v}(x)=0$ for any $v \neq \infty$, hence the only exceptional places are $p$ and $\infty$ in this case. The product formula yields

$$
\prod_{v}\left(\frac{-1, L / \mathbf{Q}}{v}\right)=\left(\frac{-1, L / \mathbf{Q}}{p}\right)\left(\frac{-1, L / \mathbf{Q}}{\infty}\right)=1
$$

i.e., $\left(\frac{-1, L / \mathbf{Q}}{p}\right)=\left(\frac{-1, L / \mathbf{Q}}{\infty}\right)$. We must examine the extension of local fields

$$
L_{\infty} / \mathbf{Q}_{\infty}=\mathbf{R}\left(\sqrt{p^{*}}\right) / \mathbf{R}
$$

If $p \equiv 1(\bmod 4)$, then $p^{*}=p$ and the extension is trivial. Thus $\left(\frac{-1, L / \mathbf{Q}}{\infty}\right)=1$ in this case. If $p \equiv 3(\bmod 4)$, then $L_{\infty}=\mathbf{C}$, and thus (again by property 2 of Theorem 1.41 :

$$
\left(\frac{-1, L / \mathbf{Q}}{\infty}\right)=(-1)^{\operatorname{ord}_{\infty}(-1)}=-1
$$

(here we write $\operatorname{Gal}(\mathbf{C} / \mathbf{R})=\{ \pm 1\})$. Combining these two cases we get

$$
\left(\frac{-1, L / \mathbf{Q}}{p}\right)=\left(\frac{-1, L / \mathbf{Q}}{\infty}\right)=(-1)^{\frac{p-1}{2}} .
$$

- Now consider $\left(\frac{q, L / \mathbf{Q}}{v}\right)$, which is equal to 1 except perhaps for $v \in\{p, q\}$. First, if $q=p$ we have $\left(\frac{p, L / \mathbf{Q}}{v}\right)=1$ for all $v \neq p$. But then the product formula is reduced to

$$
\left(\frac{p, L / \mathbf{Q}}{p}\right)=1
$$

We may therefore assume $q \neq p$ in the following.

- For $v=p$, to compute $\left(\frac{q, L / \mathbf{Q}}{p}\right)$ we must understand $\left(q, L_{p} / \mathbf{Q}_{p}\right)$. Now since $q \neq p$ we have $q \in U_{p}=\mathbf{Z}_{p}^{\times}$. Furthermore, the extension $L_{p} / \mathbf{Q}_{p}$ is ramified, so we can only conclude by property 1 in Theorem 1.41 that $\left(q, L_{p} / \mathbf{Q}_{p}\right)=1$ if and only if $q \in \mathrm{Nm}_{L_{p} / \mathbf{Q}_{p}}\left(L_{p}^{\times}\right)$. But we know that since $L_{p} / \mathbf{Q}_{p}$ is ramified, the inertia group $I_{p}$ has order 2, and hence by property 3 in Theorem 1.41 , the quotient group

$$
U_{p} /\left(U_{p} \cap \operatorname{Nm}_{L_{p} / \mathbf{Q}_{p}}\left(L_{p}^{\times}\right)\right) \cong I_{p}
$$

has order 2. Thus $U_{p} \cap \operatorname{Nm}_{L_{p} / \mathbf{Q}_{p}}\left(L_{p}^{\times}\right)$has index 2 in $U_{p}=\mathbf{Z}_{p}^{\times}$. Hence $\left(q, L_{p} / \mathbf{Q}_{p}\right)=1$ if and only if $q$ belongs to a subgroup of index 2 of $\mathbf{Z}_{p}^{\times}$. But by Example 1.36, the only subgroup of index 2 of $U_{p}=\mathbf{Z}_{p}^{\times}$is

$$
\mu_{p-1}^{2} \oplus U_{p}^{1}=\left(U_{p}\right)^{2}
$$

Therefore $\left(q, L_{p} / \mathbf{Q}_{p}\right)=1$ if and only if $q \in\left(U_{p}\right)^{2}$, which happens if and only if $\bar{q} \in \mathbf{F}_{p}^{\times 2}$, i.e., if and only if $q$ is a square modulo $p$. Thus we have that

$$
\left(\frac{q, L / \mathbf{Q}}{p}\right)=\left(\frac{q}{p}\right)
$$

- For $v=q$, since the extension $L_{q} / \mathbf{Q}_{q}$ is unramified, we can use property 2 of Theorem 1.41 to get

$$
\left(\frac{q, L / \mathbf{Q}}{q}\right)=\left(\frac{L / \mathbf{Q}}{q}\right) .
$$

This Frobenius is trivial if and only if the extension $L_{q} / \mathbf{Q}_{q}=\mathbf{Q}_{q}\left(\sqrt{p^{*}}\right) / \mathbf{Q}_{q}$ is trivial, which happens if and only if $p^{*}$ is a square modulo $q$. In other words,

$$
\left(\frac{q, L / \mathbf{Q}}{q}\right)=\left(\frac{p^{*}}{q}\right),
$$

where $\left(p^{*} / q\right)$ is the Kronecker symbol, which is the same as the Legendre symbol for $q \neq 2$, and for $q=2$ it is defined as

$$
\left(\frac{a}{2}\right)=\left\{\begin{aligned}
1 & \text { if } a \equiv 1(\bmod 8) \\
-1 & \text { otherwise }
\end{aligned}\right.
$$

| $p(\bmod 8)$ | $p^{*}$ | $\left(\frac{p^{*}}{2}\right)$ | $(-1)^{\frac{p^{2}-1}{8}}$ |
| :---: | ---: | ---: | ---: |
| 1 | $p$ | 1 | 1 |
| 3 | $-p$ | -1 | -1 |
| 5 | $p$ | -1 | -1 |
| 7 | $-p$ | 1 | 1 |

Table 1.2: $\left(\frac{2}{p}\right)=(-1)^{\frac{p^{2}-1}{8}}$.

Thus the product formula yields

$$
\left(\frac{q}{p}\right)=\left(\frac{p^{*}}{q}\right)
$$

Finally, for $p$ and $q$ odd primes, the formulation of quadratic reciprocity in Theorem 1.1 is obtained from the above by using the multiplicativity of the Legendre symbols:

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

yielding

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

For $q=2$ and $p$ odd we get $(2 / p)=\left(p^{*} / 2\right)$, but here we cannot use multiplicativity of the symbols as above, since e.g., $(3 / 2)=(5 / 2)=(15 / 2)=-1$. Instead we check all possible values for $p(\bmod 8)$ (see Table 1.2 ), and find

$$
\left(\frac{2}{p}\right)=(-1)^{\frac{p^{2}-1}{8}} .
$$

Thus we have also obtained the supplementary law of quadratic reciprocity.

### 1.6 Symbols over number fields

We now turn to the notion of symbols over number fields which we later will connect to algebraic $K$-theory. We mostly follow Gras' presentation [Gra03, II §7]. As before, $F$ denotes a number field and $v \in \mathrm{Pl}^{\text {nc }}$ a place of $F$.

In this section we will disregard all infinite complex places. There are several reasons for this exclusion. First, by our Definition 1.29 , all infinite complex places have trivial valuation. Secondly, we will repeatedly consider the group of roots of unity $\mu\left(F_{v}\right)$ of the local field $F_{v}$, which is finite cyclic whenever $v \in \mathrm{Pl}_{F}^{\mathrm{nc}}$, but infinite if $F_{v}=\mathbf{C}$. Finally, Moore has shown that any continuous Steinberg symbol (see Definition 2.16) on $\mathbf{C}$ is in fact trivial (cf. Mil71, Corollary A.2]).

Definition 1.48. The local Hilbert symbol at $v$ is the map

$$
(-,-)_{v}: F_{v}^{\times} \times F_{v}^{\times} \longrightarrow \mu\left(F_{v}\right)
$$

defined by

$$
(x, y)_{v}:=\frac{\left(y, F_{v}(\sqrt[m_{v}]{x}) / F_{v}\right)(\sqrt[m_{v}]{x})}{\sqrt[m_{v}]{x}}
$$

Here $\left(y, F_{v}(\sqrt[m_{v}]{x}) / F_{v}\right)$ is the local reciprocity symbol for the cyclic extension $F_{v}(\sqrt[m_{v}]{x}) / F_{v}$ evaluated at $y \in F_{v}^{\times}$. It is an element of the Galois $\operatorname{group} \operatorname{Gal}\left(F_{v}(\sqrt[m v]{x}) / F_{v}\right)$ and therefore permutes the different $m_{v}$-roots of $x$. Hence $\left(y, F_{v}(\sqrt[m_{n}]{x}) / F_{v}\right)(\sqrt[m_{v}]{x})$ and $\sqrt[m_{n}]{x}$ differ by a root of unity.

Proposition 1.49 ( $\boxed{\text { Gra03 }}$, II Proposition 7.1.1]). The local Hilbert symbol is Z-bilinear, and nondegenerate as a map on $F_{v}^{\times} / F_{v}^{\times m_{v}} \times F_{v}^{\times} / F_{v}^{\times m_{v}}$. Furthermore, it satisfies the following relations:

- $(x, 1-x)_{v}=1$ for all $x \in F_{v}^{\times}, x \neq 1 \quad$ (Steinberg relation);
- $(x,-x)_{v}=1$;
- $(x, y)_{v}=(y, x)_{v}^{-1}$, for all $x, y \in F_{v}^{\times}$.

Proof (following Gra03]). We only show the Steinberg relation. Antisymmetry and the relation $(x,-x)_{v}=1$ follow from the Steinberg relation; see Lemma 2.17 .

To this end, by Theorem 1.41 it is enough to show that for $x \in F^{\times}$, the element $1-x$ is a norm in the extension $M / F_{v}$, where $M:=F_{v}(\sqrt[m v]{x})$. Let $d:=\left[M: F_{v}\right]$, so that $d \mid m_{v}$. By Kummer theory (see e.g., Gra03, I §6]) there is a $t \in F_{v}^{\times}$such that $M=F_{v}(\sqrt[d]{t})$ and $x=t^{m_{v} / d}$. Note that for any $\xi \in \mu\left(F_{v}\right)$ we have

$$
\begin{aligned}
\operatorname{Nm}_{M / F_{v}}(1-\xi \sqrt[d]{t}) & =\prod_{\sigma \in \operatorname{Gal}\left(M / F_{v}\right)} \sigma(1-\xi \sqrt[d]{t}) \\
& =\prod_{j=1}^{d}\left(1-\zeta_{d}^{j} \xi \sqrt[d]{t}\right) \\
& =1-\xi^{d} t
\end{aligned}
$$

where $\zeta_{d}$ is a primitive $d$-th root of unity.
Now let $\zeta_{v}$ denote a generator for $\mu\left(F_{v}\right)$. We claim that the element

$$
\prod_{j=1}^{\frac{m_{v}}{d}}\left(1-\zeta_{v}^{j} \sqrt[d]{t}\right) \in M^{\times}
$$

has norm equal to $1-x$ :

$$
\begin{aligned}
\operatorname{Nm}_{M / F_{v}}\left(\prod_{j=1}^{\frac{m_{v}}{d}}\left(1-\zeta_{v}^{j} \sqrt[d]{t}\right)\right) & =\prod_{j=1}^{\frac{m_{v}}{d}} \operatorname{Nm}_{M / F_{v}}\left(1-\zeta_{v}^{j} \sqrt[d]{t}\right) \\
& =\prod_{j=1}^{\frac{m_{v}}{d}}\left(1-\zeta_{v}^{d j} t\right) \\
& =1-\zeta_{v}^{m_{v}} t^{\frac{m_{v}}{d}}=1-x
\end{aligned}
$$

as desired.
For a finite place $v \in \mathrm{Pl}_{0}$, let $\ell$ denote the characteristic of the residue field $k(v)=\mathcal{O}_{v} / \mathfrak{p}_{v}$. By Proposition 1.35 there is a decomposition $F_{v}^{\times} \cong \pi_{v}^{\mathbf{Z}} \oplus \mu_{q_{v}-1} \oplus U_{v}^{1}$. This yields a decomposition

$$
\mu\left(F_{v}\right)=\mu_{q_{v}-1} \oplus \mu\left(U_{v}^{1}\right)
$$

where $\mu\left(U_{v}^{1}\right)$ denotes the torsion part of $U_{v}^{1}$, whose order is a power of $\ell$. With this decomposition in mind we will from now on instead denote the group $\mu_{q_{v}-1}$ by $\mu\left(F_{v}\right)^{\text {reg }}$ and $\mu\left(U_{v}^{1}\right)$ by $\mu\left(F_{v}\right)^{1}$. If $v$ is a real place at infinity, we will set $\mu\left(F_{v}\right)^{\text {reg }}:=\mu_{2}$ and $\mu\left(F_{v}\right)^{1}:=1$. Here we do not consider complex places at infinity since $\operatorname{ord}_{v}(x)=0$ for all $x \in F_{v}$ for such a place. We write

$$
m_{v}^{1}:=\# \mu\left(F_{v}\right)^{1}
$$

There are only finitely many places $v$ for which $\mu\left(F_{v}\right)^{1} \neq 1$, and such places will be referred to as irregular places. The places $v$ with $\mu\left(F_{v}\right)=1$ are called regular.

## 1. Results from Algebraic Number Theory

Definition 1.50. Let $v$ be a noncomplex place of the number field $F$. The regular, or tame Hilbert symbol at $v,(-,-)_{v}^{\mathrm{reg}}$, is the map on $F_{v}^{\times} \times F_{v}^{\times}$defined by

$$
(x, y)_{v}^{\mathrm{reg}}:=\frac{\left(y, F_{v}(\sqrt[q_{v}-1]{x}) / F_{v}\right)(\sqrt[q_{v}-1]{x})}{\sqrt[q_{v}-1]{x}}
$$

Remark 1.51. The regular Hilbert symbol is nothing more than a power of the local Hilbert symbol at $v$ :

$$
(x, y)_{v}^{\mathrm{reg}}=(x, y)_{v}^{m_{v}^{1}}=\left(x^{m_{v}^{1}}, y\right)_{v} .
$$

Thus the results of Proposition 1.49 also hold for the regular Hilbert symbol.
Proposition 1.52 ( $(\boxed{G r a 03}, ~ I I ~ P r o p o s i t i o n ~ 7.1 .5]) . ~ F o r ~ a n y ~ x, y \in F_{v}^{\times}$the following formula for $(x, y)_{v}^{\text {reg }}$ holds:

$$
(x, y)_{v}^{\mathrm{reg}}= \begin{cases}(-1)^{\operatorname{ord}_{v}(x) \operatorname{ord}_{v}(y)} x^{\operatorname{ord}_{v}(y)} y^{-\operatorname{ord}_{v}(x)}+\mathfrak{p}_{v}, & v \text { finite } \\ (-1)^{\operatorname{ord}_{v}(x) \operatorname{ord}_{v}(y)}, & v \text { infinite real }\end{cases}
$$

where we in the finite case identify $\mu\left(F_{v}\right)^{\text {reg }}$ with $k(v)^{\times}$.
Proof (following [Gra03]). First assume $v$ is an infinite real place. Then, since $\mu\left(F_{v}^{\times}\right)=\mu_{2}$ we have

$$
(x, y)_{v}^{\mathrm{reg}}=\frac{\left(y, F_{v}(\sqrt{x}) / F_{v}\right)(\sqrt{x})}{\sqrt{x}}
$$

If $i_{v}(x)>0$, or in other words $\operatorname{ord}_{v}(x)=0$, we have $F_{v}(\sqrt{x})=F_{v}=\mathbf{R}$. Therefore $\left(y, F_{v}(\sqrt{x}) / F_{v}\right)=1$, and hence

$$
(x, y)_{v}^{\mathrm{reg}}=1=(-1)^{\operatorname{ord}_{v}(x) \operatorname{ord}_{v}(y)}
$$

If $i_{v}(x)<0$ so that $\operatorname{ord}_{v}(x)=1$ we have $F_{v}(\sqrt{x}) / F_{v}=\mathbf{C} / \mathbf{R}$, and thus, as we saw in Example 1.43 . $(y, \mathbf{C} / \mathbf{R})=c^{\operatorname{ord}_{v}(y)}$ where $c$ is complex conjugation. Hence

$$
(x, y)_{v}^{\mathrm{reg}}=(-1)^{\operatorname{ord}_{v}(y)}=(-1)^{\operatorname{ord}_{v}(x) \operatorname{ord}_{v}(y)} .
$$

Now assume $v$ is a finite place of $F$. For any $x \in F_{v}^{\times}$we may write $x=\pi_{v}^{n} u$ for some $u \in U_{v}$ and $n=\operatorname{ord}_{v}(x)$. By bilinearity and antisymmetry of the Hilbert symbol it is therefore enough to compute

$$
\left(u, u^{\prime}\right)_{v}^{\mathrm{reg}}, \quad(u, \pi)_{v}^{\mathrm{reg}}, \quad(\pi, \pi)_{v}^{\mathrm{reg}}
$$

where $u, u^{\prime} \in U_{v}$ and where $\pi_{v}$ is abbreviated by $\pi$.

- The extension $F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}$ is unramified, hence

$$
\left(u, u^{\prime}\right)_{v}^{\mathrm{reg}}=\frac{\left(u^{\prime}, F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)(\sqrt[q_{v}-1]{u})}{\sqrt[q_{v}-1 / \bar{u}]{\sqrt{2}}}=\left(F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)^{\operatorname{ord}_{v}\left(u^{\prime}\right)}=1
$$

since $\operatorname{ord}_{v}\left(u^{\prime}\right)=0$.

- For the symbol $(u, \pi)_{v}^{\mathrm{reg}}$ we again use that $F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}$ is unramified. We have

$$
\left(\pi, F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)=\left(F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)^{\operatorname{ord}_{v}(\pi)}=\left(F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)
$$

since $\operatorname{ord}_{v}(\pi)=1$. Let $\sigma$ denote the Frobenius $\left(F_{v}(\sqrt[q_{v}-1]{u}) / F_{v}\right)$. Then

$$
(u, \pi)_{v}^{\mathrm{reg}}=\frac{\sigma(\sqrt[q_{v}-1]{u})}{\sqrt[q_{v}-1]{u}} \equiv \frac{(\sqrt[q_{v}-1]{u})^{q_{v}}}{\sqrt[q_{v}-1]{u}} \equiv u\left(\bmod \mathfrak{p}_{v}\right)
$$

by the properties of the Frobenius automorphism. Thus

$$
(u, \pi)_{v}^{\operatorname{reg}} \equiv u=(-1)^{\operatorname{ord}_{v}(u) \operatorname{ord}_{v}(\pi)} u^{\operatorname{ord}_{v}(\pi)} \pi^{-\operatorname{ord}_{v}(u)}\left(\bmod \mathfrak{p}_{v}\right)
$$

- Finally, to compute $(\pi, \pi)_{v}^{\text {reg }}$ we make use of the properties in Proposition 1.49. We have

$$
(\pi,-\pi)_{v}^{\mathrm{reg}}=1=(\pi, \pi)_{v}^{\mathrm{reg}}(\pi,-1)_{v}^{\mathrm{reg}},
$$

hence $(\pi, \pi)_{v}^{\text {reg }}=(-1, \pi)_{v}^{\text {reg }}$ by antisymmetry. By the previous result with $u=-1$ we get

$$
(\pi, \pi)_{v}^{\mathrm{reg}} \equiv-1=(-1)^{\operatorname{ord}_{v}(\pi)^{2}} \pi^{\operatorname{ord}_{v}(\pi)} \pi^{-\operatorname{ord}_{v}(\pi)}\left(\bmod \mathfrak{p}_{v}\right)
$$

By bilinearity, this concludes the proof.
Definition 1.53 (Hilbert symbols). For $x, y \in F^{\times}$and $v$ a place of $F$, the Hilbert symbol at $v$ of order $m=\# \mu(F)$ is defined as

$$
\left(\frac{x, y}{v}\right):=\frac{\sigma_{v}(\sqrt[m]{x})}{\sqrt[m]{x}} \in \mu(F)
$$

where

$$
\sigma_{v}=\left(\frac{y, F(\sqrt[m]{x}) / F}{v}\right)
$$

is the Hasse symbol at $v$.
For any divisor $n$ of $m$, we also define the Hilbert symbol of order $n$ as

$$
\left(\frac{-,-}{v}\right)_{n}:=\left(\frac{-,-}{v}\right)^{\frac{m}{n}}
$$

Example 1.54. For $F=\mathbf{Q}$ and $v \in \mathrm{Pl}_{\mathbf{Q}}$, the $\operatorname{Hilbert} \operatorname{symbol}\left(\frac{-,-}{v}\right)$ is given by

$$
\left(\frac{a, b}{v}\right)=\left\{\begin{aligned}
1 & \text { if } a x^{2}+b y^{2}=z^{2} \text { has a nontrivial solution in } \mathbf{Q}_{v} \\
-1 & \text { otherwise }
\end{aligned}\right.
$$

for any $a, b \in \mathbf{Q}^{\times}$Dal06.
Theorem 1.55 (Product formula for Hilbert symbols). For any $x, y \in F^{\times}$and any $n \mid m=$ $\# \mu(F)$, the following product formula holds:

$$
\prod_{v \in \mathrm{Pl}_{F}}\left(\frac{x, y}{v}\right)_{n}=1
$$

Proof. The conclusion follows from the product formula for Hasse symbols (Theorem 1.46) and the fact that if $\left\{\sigma_{i}\right\}_{i=1}^{N}$ is a finite set of automorphisms in the Galois group of some abelian extension $L$ of $F$ containing an $n$-th root $z$ of an element $x \in F^{\times}$, we have

$$
\frac{\left(\prod_{i=1}^{N} \sigma_{i}\right)(z)}{z}=\prod_{i=1}^{N} \frac{\sigma_{i}(z)}{z} \in \mu(F)
$$

Indeed, let $\sigma, \tau \in \operatorname{Gal}(L / F)$. Then

$$
\begin{aligned}
\frac{(\sigma \tau)(z)}{z} & =\frac{\sigma(\tau(z))}{z} \cdot \frac{\tau(z)}{\tau(z)} \\
& =\frac{\tau(\sigma(z))}{\tau(z)} \cdot \frac{\tau(z)}{z} \\
& =\tau\left(\frac{\sigma(z)}{z}\right) \cdot \frac{\tau(z)}{z} \\
& =\frac{\sigma(z)}{z} \cdot \frac{\tau(z)}{z}
\end{aligned}
$$

where the second equality holds because the extension is abelian, and the last equality holds since $\tau$ acts trivially on $F$ and $\mu_{n} \subseteq F^{\times}$. The general case follows by induction.

Definition 1.56. Let $n>1$ be a natural number. For any $x \in F^{\times}$and any $v \in \mathrm{Pl}_{0}$ not ramified in $F(\sqrt[n]{x}) / F$, we define the $n$-th power residue symbol as

$$
\left(\frac{x}{v}\right)_{n}:=\frac{\left(\frac{F(\sqrt[n]{x}) / F}{v}\right)(\sqrt[n]{x})}{\sqrt[n]{x}}
$$

If $\mathfrak{p}$ is the prime ideal corresponding to $v$, we also write

$$
\left(\frac{x}{\mathfrak{p}}\right)_{n}:=\left(\frac{x}{v}\right)_{n}
$$

Finally, for $y \in F^{\times}$such that all $v \mid y$ is unramified in $F(\sqrt[n]{x}) / F$, we define

$$
\left(\frac{x}{y}\right)_{n}:=\prod_{v \mid y \infty}\left(\frac{x}{v}\right)_{n}^{\operatorname{ord}_{v}(y)}
$$

where the product runs over all places $v \mid y$ and all infinite real places.
Write $\sigma=\left(\frac{F(\sqrt[n]{x}) / F}{v}\right)$ for the global Frobenius at the finite place $v$, and let $\mathfrak{p}$ denote the prime ideal corresponding to $v$. Recall that, modulo $\mathfrak{p}, \sigma$ raises elements to the power of $q_{v}=\operatorname{Nm}(\mathfrak{p})$. Hence

$$
\left(\frac{x}{\mathfrak{p}}\right)_{n}=\frac{\sigma(\sqrt[n]{x})}{\sqrt[n]{x}} \equiv x^{\frac{q_{v}-1}{n}}(\bmod \mathfrak{p})
$$

so the $n$-th power reciprocity symbol is a direct generalization of the Legendre symbol $(p / q) \equiv$ $p^{\frac{q-1}{2}}(\bmod q)$.

Proposition 1.57 (Gra03, pp.204-205]). The n-th power residue symbols satisfy the following properties (for any place $v$ such that the given residue symbol is defined):

1. $\left(\frac{x}{v}\right)_{n}=1$ if and only if $i_{v}(x)$ is an $n$-th power in $F_{v}^{\times}$;
2. $\left(\frac{x}{v}\right)_{n}\left(\frac{y}{v}\right)_{v}=\left(\frac{x y}{v}\right)_{n}$;
3. $\left(\frac{x, y}{v}\right)_{n}=\left(\frac{x}{v}\right)_{n}^{\operatorname{ord}_{v}(y)}$.

In particular,

$$
\left(\frac{x, \pi_{v}}{v}\right)_{n}=\left(\frac{x}{v}\right)_{n}
$$

Theorem 1.58 (The $n$-th power reciprocity law [Gra03, II Theorem 7.4.4]). Let $n$ be a divisor of $m=\# \mu(F)$. Suppose $x, y \in F^{\times}$are such that $\operatorname{ord}_{v}(x) \operatorname{ord}_{v}(y)=0$ for any $v \in \mathrm{Pl}_{F}$ (i.e., $x$ and $y$ are coprime), and $\operatorname{ord}_{v}(x)=\operatorname{ord}_{v}(y)=0$ for all places $v$ dividing $n$. Then

$$
\left(\frac{y}{x}\right)_{n}\left(\frac{x}{y}\right)_{n}^{-1}=\prod_{v \mid n}\left(\frac{x, y}{v}\right)_{n}
$$

Proof (following Gra03]). By definition we have

$$
\left(\frac{y}{x}\right)_{n}\left(\frac{x}{y}\right)_{n}^{-1}=\prod_{v \nmid n}\left(\frac{y}{v}\right)_{n}^{\operatorname{ord}_{v}(x)} \prod_{v \nmid n}\left(\frac{x}{v}\right)_{n}^{-\operatorname{ord}_{v}(y)}
$$

since $\operatorname{ord}_{v}(x)=\operatorname{ord}_{v}(y)=0$ for all $v \mid n$.
Since

$$
\left(\frac{x}{v}\right)_{n}^{\operatorname{ord}_{v}(y)}=\left(\frac{x, y}{v}\right)_{n}
$$

we have

$$
\begin{aligned}
\left(\frac{y}{x}\right)_{n}\left(\frac{x}{y}\right)_{n}^{-1} & =\prod_{\substack{v \nmid n \\
\operatorname{ord}_{v}(x) \neq 0}}\left(\frac{y, x}{v}\right)_{n} \prod_{\substack{v \not n}}\left(\frac{x, y}{v}\right)_{n}^{-1} \\
& =\prod_{\substack{v \nmid n \\
\operatorname{ord}_{v}(y) \neq 0}}\left(\frac{y, x}{v}\right)_{n}=\prod_{v \nmid n}\left(\frac{y, x}{v}\right)_{n}
\end{aligned}
$$

where the last two equalities follow since by assumption $\operatorname{ord}_{v}(x y) \neq 0 \Longleftrightarrow \operatorname{ord}_{v}(x) \neq 0$ or $\operatorname{ord}_{v}(y) \neq 0$. By the product formula for Hilbert symbols this equals

$$
\prod_{v \nmid n}\left(\frac{y, x}{v}\right)_{n}=\prod_{v \mid n}\left(\frac{y, x}{v}\right)_{n}^{-1}=\prod_{v \mid n}\left(\frac{x, y}{v}\right)_{n},
$$

so this concludes the proof.
Let us return to the Gaussian integers. Recall the statement of biquadratic reciprocity (Theorem 1.6): For distinct primary Gaussian primes $\pi, \tau$, we have

$$
\left(\frac{\tau}{\pi}\right)_{4}=(-1)^{\frac{\operatorname{Nm}(\pi)-1}{4} \cdot \frac{\operatorname{Nm}(\tau)-1}{4}}\left(\frac{\pi}{\tau}\right)_{4}
$$

Proof of Theorem 1.6. The following argument is due to Lenstra and Stevenhagen.
Since $\mathbf{Q}(i)$ only has complex places at infinity, and since $(1+i)$ is the only place dividing $n=4$, power reciprocity states that

$$
\left(\frac{\tau}{\pi}\right)_{4}\left(\frac{\pi}{\tau}\right)_{4}^{-1}=\left(\frac{\pi, \tau}{(1+i)}\right)_{4},
$$

where the right hand side is the Hilbert symbol of order 4 at $(1+i)$. Thus we want to show

$$
\begin{equation*}
\left(\frac{\pi, \tau}{(1+i)}\right)_{4}=(-1)^{\frac{\mathrm{Nm}(\pi)-1}{4} \cdot \frac{\mathrm{Nm}(\tau)-1}{4}} \tag{1.3}
\end{equation*}
$$

Let $v$ denote the place corresponding to $(1+i)$. Note that for a primary Gaussian prime $\alpha$ we have $\alpha \in U_{v}^{3}$ by definition (here we identify $\alpha$ with its image $i_{v}(\alpha)$ in $U_{v}$ ). Thus $\alpha$ has component 1 on $\mu_{4}$ in the decomposition of $\mathbf{Z}_{2}[i]^{\times}$of Proposition 1.22 , and by Proposition 1.23 the same holds for the residue class of $\alpha$ in

$$
U_{v} /\left(U_{v}\right)^{4} \cong \mu_{4} \oplus\langle\overline{3+2 i}\rangle \oplus\langle\overline{5}\rangle
$$

We must therefore compute the Hilbert symbols $\left(\frac{\pi, \tau}{v}\right)_{4}$ for $\pi$ and $\tau$ in the congruence classes of $3+2 i$ and 5 modulo $\left(U_{v}\right)^{4}$. According to the right hand side of Equation 1.3 we must show that

$$
\left(\frac{\pi, \tau}{v}\right)_{4}=\left\{\begin{aligned}
1 & \text { if } \pi \text { or } \tau \equiv 5\left(\bmod \left(U_{v}\right)^{4}\right) \\
-1 & \text { if } \pi \equiv \tau \equiv 3+2 i\left(\bmod \left(U_{v}\right)^{4}\right)
\end{aligned}\right.
$$

If $\alpha, \beta \in \mathbf{Z}[i]$ satisfy $\operatorname{ord}_{v}(\alpha \beta) \geq 7$, then by Proposition 1.23 (3) we have $1-\alpha \beta \in\left(U_{v}\right)^{4}$. Using this we will show that $\left(\frac{1+\alpha, 1+\beta}{v}\right)_{4}=1$.

First, note that

$$
\left(\frac{1+\alpha,-\alpha}{v}\right)_{4}=1
$$

by the Steinberg relation (cf. Proposition 1.49). Hence, using bilinearity, we get

$$
\left(\frac{1+\alpha, 1+\beta}{v}\right)_{4}=\left(\frac{1+\alpha, 1+\beta}{v}\right)_{4}\left(\frac{1+\alpha,-\alpha}{v}\right)_{4}=\left(\frac{1+\alpha,-\alpha(1+\beta)}{v}\right)_{4}
$$

Since $1-\alpha \beta \in\left(U_{v}\right)^{4}$, we have

$$
\left(\frac{-\alpha(1+\beta), 1-\alpha \beta}{v}\right)_{4}=1=\left(\frac{(1-\alpha \beta)^{-1},-\alpha(1+\beta)}{v}\right)_{4}
$$

Thus

$$
\begin{aligned}
\left(\frac{1+\alpha,-\alpha(1+\beta)}{v}\right)_{4} & =\left(\frac{1+\alpha,-\alpha(1+\beta)}{v}\right)_{4}\left(\frac{(1-\alpha \beta)^{-1},-\alpha(1+\beta)}{v}\right)_{4} \\
& =\left(\frac{\frac{1+\alpha}{1-\alpha \beta},-\alpha(1+\beta)}{v}\right)_{4}
\end{aligned}
$$

Multiplying similarly in the second slot we obtain

$$
\begin{aligned}
\left(\frac{1+\alpha, 1+\beta}{v}\right)_{4} & =\left(\frac{1+\alpha,-\alpha(1+\beta)}{v}\right)_{4} \\
& =\left(\frac{\frac{1+\alpha}{1-\alpha \beta}, \frac{-\alpha(1+\beta)}{1-\alpha \beta}}{v}\right)_{4} \\
& =\left(\frac{\frac{1+\alpha}{1-\alpha \beta}, 1-\frac{1+\alpha}{1-\alpha \beta}}{v}\right)_{4}=1
\end{aligned}
$$

Since $\operatorname{ord}_{v}(4)=4$, this finishes the case where at least one of $\pi$ or $\tau$ is congruent to 5 modulo $\left(U_{v}\right)^{4}$. For if $\pi=5$ while $\tau$ equals 5 or $3+2 i$, then we can write $\pi=1+\alpha, \tau=1+\beta$ where $\operatorname{ord}_{v}(\alpha \beta) \geq 7$.

For the remaining case we shall show that

$$
\left(\frac{1+\alpha}{1+\beta}\right)_{4}\left(\frac{1+\beta}{1+\alpha}\right)_{4}^{-1}=-1
$$

for some $\alpha, \beta \in \mathbf{Z}[i]$ with $\operatorname{ord}_{v}(\alpha)=\operatorname{ord}_{v}(\beta)=3$, i.e., $1+\alpha, 1+\beta \in U_{v}^{3} \backslash U_{v}^{4}$. Put $\alpha=(-i)(1+i)^{3}$ and $\beta=-(1+i)^{3}$. Using that

$$
1+\alpha=3+2 i \equiv 1-i(\bmod 3-2 i)
$$

we get $\left(\frac{3+2 i}{3-2 i}\right)_{4}=\left(\frac{1-i}{3-2 i}\right)_{4}$. However,

$$
(1-i)^{\frac{\operatorname{Nm}(3-2 i)-1}{4}}=(1-i)^{3}=-2-2 i \equiv i(\bmod 3-2 i),
$$

hence

$$
\left(\frac{1+\alpha}{1+\beta}\right)_{4}=i
$$

Since $1+\beta$ is the complex conjugate of $1+\alpha$ we then have

$$
\left(\frac{1+\beta}{1+\alpha}\right)_{4}=\overline{\left(\frac{1+\alpha}{1+\beta}\right)_{4}}=-i
$$

where the bar denotes complex conjugation. Hence we end up with

$$
\left(\frac{1+\alpha}{1+\beta}\right)_{4}\left(\frac{1+\beta}{1+\alpha}\right)_{4}^{-1}=-1
$$

We now turn to the supplementary laws stated in Theorem 1.6 . The identity $(i / \pi)_{4}=$ $i^{(\operatorname{Nm}(\pi)-1) / 4}$ is true by definition, so we need only concern ourselves with the symbol $((1+i) / \pi)_{4}$. As noted in Theorem 1.6, we use the isomorphism

$$
\frac{U_{1+i}}{\left(U_{1+i}\right)^{4}} \cong \mu_{4} \oplus\langle\overline{3+2 i}\rangle \oplus\langle\overline{5}\rangle
$$

to write $\pi \equiv i^{k}(3+2 i)^{\mu} 5^{\nu}\left(\bmod \left(U_{1+i}\right)^{4}\right)$. Multiplicativity of residue symbols yields

$$
\left(\frac{1+i}{\pi}\right)_{4}=\left(\frac{1+i}{3+2 i}\right)_{4}^{\mu}\left(\frac{1+i}{2+i}\right)_{4}^{\nu}\left(\frac{1+i}{2-i}\right)_{4}^{\nu}
$$

We compute ${ }^{2}$

$$
\begin{aligned}
\left(\frac{1+i}{3+2 i}\right)_{4} & \equiv(1+i)^{\frac{\mathrm{Nm}(3+2 i)-1}{4}} \equiv(1+i)^{3} \equiv-i(\bmod 3+2 i) \\
\left(\frac{1+i}{2+i}\right)_{4} & \equiv 1+i \equiv-1(\bmod 2+i) \\
\left(\frac{1+i}{2-i}\right)_{4} & \equiv-i(\bmod 2-1)
\end{aligned}
$$

Hence

$$
\left(\frac{1+i}{\pi}\right)_{4}=(-i)^{\mu}(-1)^{\nu}(-i)^{\nu}=i^{\nu-\mu}
$$

as desired.

[^1]
## CHAPTER

## Algebraic $K$-theory

The zeroth algebraic $K$-group originated in Grothendieck's work on his formulation of the generalized Riemann-Roch theorem Dal06. Through his need to work on sheaves over a scheme $X$ there arose two definitions of algebraic $K$-groups, either of which is referred to as the Grothendieck group of $X$ : He defined $K_{0}(X)$ as a certain group built from isomorphism classes of locally free sheaves of finite rank over $X$, and similarly a group $G_{0}(X)$ of classes of coherent sheaves over $X$. If $X$ is a regular variety, the two definitions coincide. If $X$ is the spectrum of a ring $A$, we have an alternate definition of the $K$-group of $X$. We consider the monoid $M$ consisting of isomorphism classes of finitely generated projective modules over $A$, keeping in mind that such modules may be regarded as locally free $\mathcal{O}_{X}$-modules. Then $K_{0}(X)$ is defined as the universal object with respect to monoid homomorphisms of $M$ into various abelian groups.

The functor $K_{0}$ is the starting point for the definition of a series of functors attaching important invariants to a ring. As a starter, we shall see in Example 2.4 that the torsion part of $K_{0}(A)$ is precisely the ideal class group $\mathrm{Cl}(A)$ when $A$ is the ring of integers in a number field.

### 2.1 Classical $K$-theory of rings

There is a forgetful functor

$$
U: \mathrm{Ab} \longrightarrow \mathrm{CMon}
$$

from the category of abelian groups to the category of commutative monoids, which admits a left adjoint - a "group completion functor"

$$
(-)^{+}: \text {CMon } \longrightarrow \mathrm{Ab}
$$

For a given commutative monoid $M$, the group $M^{+}$is called the Grothendieck group, or the group completion of $M$. Thus, the Grothendieck group of $M$ is a group $M^{+}$together with a map $i: M \rightarrow M^{+}$(not necessarily an injection) such that the following universal property is satisfied. Given any abelian group $G$ with a monoid map $f: M \rightarrow G$, there is a unique map of abelian groups $\bar{f}: M^{+} \rightarrow G$ such that the following diagram commutes:


One constructs $M^{+}$by adding formal differences as follows Ros94, Theorem 1.1.3]. Define the equivalence relation $\sim$ on $M \times M$ by letting $\left(x_{1}, y_{1}\right) \sim\left(x_{2}, y_{2}\right)$ if and only if there is a $t \in M$ such that

$$
x_{1}+y_{2}+t=x_{2}+y_{1}+t
$$

Then $M^{+}=M \times M / \sim$.

Now let $(\mathscr{C}, \square)$ be a symmetric monoidal category (see [ML98, p.184] for a definition), and assume that the isomorphism classes of objects of $\mathscr{C}$ form a set, which we denote by $\mathscr{C}^{\text {iso }}$. The set $\mathscr{C}^{\text {iso }}$ is then a commutative monoid with addition given by

$$
[C]+\left[C^{\prime}\right]:=\left[C \square C^{\prime}\right]
$$

where we write $C \square C^{\prime}:=\square\left(C, C^{\prime}\right)$ for any $C, C^{\prime} \in \mathscr{C}$. The Grothendieck group of $\mathscr{C}$ is then defined as

$$
K_{0}^{\square}(\mathscr{C}):=\left(\mathscr{C}^{\text {iso }}\right)^{+},
$$

and is often abbreviated to $K_{0}(\mathscr{C})$ if no confusion is likely to arise (although confusion will occur in the case of $K_{0}$ of schemes-see below).

To define $K_{0}$ of a ring $A$, we take as our symmetric monoidal category the category $\mathscr{P}(A)$ of finitely generated projective $A$-modules, with direct sum as the associated bifunctor.
Definition 2.1. The zeroth algebraic $K$-group $K_{0}(A)$ of the ring $A$ is the Grothendieck group of $(\mathscr{P}(A), \oplus)$ :

$$
K_{0}(A):=K_{0}^{\oplus}(\mathscr{P}(A)) .
$$

Moreover, tensor product over $A$ induces a ring structure on $K_{0}(A)$.
Example 2.2. If $A$ is a principal ideal domain or a local ring, any finitely generated projective $A$-module is free ([Ros94, Theorem 1.3.1], Wei13, Lemma 2.2, p.11]). Since the rank is well defined up to isomorphism, the map $\mathscr{P}(A)^{\text {iso }} \rightarrow \mathbf{Z}_{\geq 0}$ defined by $\left[A^{n}\right] \mapsto n$ is an isomorphism. By group completing the two monoids we obtain an isomorphism $K_{0}(A) \cong \mathbf{Z}$.

To generalize the definition of $K_{0}$ of a ring to arbitrary schemes, we consider the category $\operatorname{VB}(X)$ of vector bundles over a scheme $X$. It turns out that the most satisfactory way of defining $K_{0}$ of $X$ is to set

$$
K_{0}(X):=K_{0}(\mathrm{VB}(X)),
$$

where the right hand side means $K_{0}$ of $\operatorname{VB}(X)$ viewed as an exact category, and not as a symmetric monoidal category. See Wei13, II Ch.7] for details on $K_{0}$ of an exact category. Since $\operatorname{VB}(X)$ is not always a split exact category ${ }^{1}$, we have $K_{0}(X) \not \not K_{0}^{\oplus}(X)$ in general. On the other hand, if $X=\operatorname{Spec} A$ is an affine scheme, there is an equivalence of categories $\mathscr{P}(A) \cong \mathrm{VB}(\operatorname{Spec} A)$ Wei13, p.51], hence $K_{0}(A)=K_{0}(\operatorname{Spec} A)$.

For any scheme $X$ there is a rank map

$$
\text { rank }: K_{0}(X) \longrightarrow H^{0}(X, \underline{\mathbf{Z}})
$$

induced by taking the rank of vector bundles, whose kernel $\widetilde{K}_{0}(X)$ is called the reduced $K$ theory of $X$. Here $\underline{\mathbf{Z}}$ is the constant sheaf on $X$ with group $\mathbf{Z}$, and $H^{0}(X, \underline{\mathbf{Z}})$ is viewed as a subring of $K_{0}(X)$. Similarly, there is a determinant map

$$
\operatorname{det}: K_{0}(X) \longrightarrow \operatorname{Pic}(X)
$$

induced by the determinant map on vector bundles sending a locally free rank $n$ bundle $\mathscr{F}$ to the line bundle $\operatorname{det} \mathscr{F}=\bigwedge^{n} \mathscr{F}$. The map $\operatorname{rank} \oplus \operatorname{det}$ is a surjective ring homomorphism Wei13, Theorem 8.1, p.158].

Theorem 2.3 ([Wei13, Proposition 8.2.1, p.159]). If $X$ is a 1-dimensional separated, regular noetherian scheme, the map rank $\oplus$ det yields an isomorphism

$$
K_{0}(X) \cong H^{0}(X, \underline{\mathbf{Z}}) \oplus \operatorname{Pic}(X)
$$

Example 2.4. Let $A$ be a Dedekind ring. Then the Picard group of $\operatorname{Spec} A$ is the ideal class $\operatorname{group} \mathrm{Cl}(A)$ Har77, II Example 6.3.2], and since $\operatorname{Spec} A$ is connected we have $H^{0}(\operatorname{Spec} A, \underline{\mathbf{Z}})=$ Z. Hence the theorem above yields

$$
K_{0}(A) \cong \mathbf{Z} \oplus \mathrm{Cl}(A)
$$

i.e., $\widetilde{K}_{0}(A) \cong \mathrm{Cl} A$.

[^2]
### 2.1.1 The functors $K_{1}$ and $K_{2}$

Let $A$ be a ring. Recall the definition of the general linear group $\mathrm{GL}(A)$ as the colimit of the directed system

$$
\mathrm{GL}_{1}(A) \hookrightarrow \mathrm{GL}_{2}(A) \hookrightarrow \cdots,
$$

where the injections of $\mathrm{GL}_{n}(A)$ into $\mathrm{GL}_{n+1}(A)$ are defined by $M \mapsto\left(\begin{array}{cc}M & 0 \\ 0 & 1\end{array}\right)$. Let $e_{i j}(a)$ denote the matrix which coincides with the identity matrix except having $a$ in the $(i, j)$-position, where $i \neq j$. Such a matrix is called an elementary matrix, and the elementary $n \times n$-matrices generate a subgroup $E_{n}(A)$ of $\mathrm{GL}_{n}(A)$. As in the above sequence we have maps $E_{n}(A) \rightarrow E_{n+1}(A)$ for each $n$, and the colimit of this system is denoted by $E(A)$.

Lemma 2.5 (Whitehead). The group $E(A)$ is the commutator subgroup of $\mathrm{GL}(A)$.
Whitehead's lemma tells us that $\mathrm{GL}(A) / E(A)$ is the maximal abelian quotient $\mathrm{GL}(A)^{\mathrm{ab}}$ of $\mathrm{GL}(A)$. This group was recognized by Bass to be the correct definition of $K_{1}(A)$ :

Definition 2.6. The first algebraic $K$-group of the ring $A$ is defined as

$$
K_{1}(A):=\mathrm{GL}(A) / E(A)
$$

A ring map $A \rightarrow B$ induces a map $\mathrm{GL}(A) \rightarrow \mathrm{GL}(B)$ preserving elementary matrices, and hence a map $K_{1}(A) \rightarrow K_{1}(B)$. Thus $K_{1}$ is a covariant functor on the category of rings.

There is a determinant map det: GL $(A) \rightarrow A^{\times}$, which induces a map $K_{1}(A) \rightarrow A^{\times}$. Denote by $\operatorname{SL}(A)$ the kernel of the determinant map, and let $S K_{1}(A):=\mathrm{SL}(A) / E(A)$. By realizing $A^{\times}$as $\mathrm{GL}_{1}(A)$, which then maps into $\mathrm{GL}(A)$, we obtain a split exact sequence

$$
1 \longrightarrow S K_{1}(A) \longrightarrow K_{1}(A) \longrightarrow A^{\times} \longrightarrow 1
$$

Example 2.7. The kernel $S K_{1}(A)$ is known to be trivial for example when $A$ is a local ring [Wei13, III Lemma 1.4] or the ring of integers in a number field [BMS67]. In other words, we get an isomorphism

$$
K_{1}(A) \cong A^{\times}
$$

in this case.
Definition 2.8. A central extension of a group $G$ by an abelian group $A$ is a short exact sequence

$$
1 \longrightarrow A \longrightarrow E \xrightarrow{\phi} G \longrightarrow 1,
$$

where $E$ is a group containing $A$ as a central subgroup. A central extension is denoted by $(E, \phi)$.

The central extensions of $G$ form a category where a morphism $\psi:(E, \phi) \rightarrow\left(E^{\prime}, \phi^{\prime}\right)$ is a homomorphism $\psi: E \rightarrow E^{\prime}$ over $G$, i.e., a commutative diagram


We may therefore define a central extension $(E, \phi)$ of $G$ to be universal if $(E, \phi)$ is initial in the category of central extensions of $G$.

Recall that a group $G$ is perfect if it equals its commutator subgroup $[G, G]$. In other words, a group is perfect if it has trivial abelianization, which is equivalent to the vanishing of the first group homology $H_{1}(G, \mathbf{Z})$.

Theorem 2.9 ([Ros94, Theorem 4.1.3]). A group admits a universal central extension if and only if it is perfect.

We turn back to a quick study of the elementary matrices over a ring $A$. There are a few evident relations satisfied by such matrices. First we have

$$
e_{i j}(a) e_{i j}(b)=e_{i j}(a+b)
$$

Secondly, the commutator of two elementary matrices satisfies

$$
\left[e_{i j}(a), e_{k l}(b)\right]= \begin{cases}1, & j \neq k, i \neq l \\ e_{i l}(a b), & j=k, i \neq l \\ e_{k j}(-a b), & j \neq k, i=l\end{cases}
$$

With this in mind, we define the Steinberg group as the group generated by symbols satisfying the analog of the relations given above:

Definition 2.10. Let $A$ be a ring and $n \geq 3$ a natural number. The Steinberg group $\operatorname{St}_{n}(A)$ of order $n$ over $A$ is defined as the free group generated by symbols $x_{i j}(a), 1 \leq i, j \leq n, i \neq j$, $a \in A$, modulo the relations

1. $x_{i j}(a) x_{i j}(b)=x_{i j}(a+b)$;
2. $\left[x_{i j}(a), x_{k l}(b)\right]=x_{i l}(a b) \quad$ if $i \neq l$;
3. $\left[x_{i j}(a), x_{k l}(b)\right]=1 \quad$ if $j \neq k, i \neq l$.

For every $n$ there are natural maps $\mathrm{St}_{n}(A) \rightarrow \mathrm{St}_{n+1}(A)$, and we let $\operatorname{St}(A)$ denote the colimit of the directed system. Given a ring homomorphism $f: A \rightarrow B$ we get an induced map on the free groups generated by the $\left\{x_{i j}(a): a \in A\right\}$ and $\left\{x_{i j}(b): b \in B\right\}$ by sending $x_{i j}(a)$ to $x_{i j}(f(a))$. This is compatible with the relations in $\operatorname{St}(A)$ and therefore factors through a map $\operatorname{St}(A) \rightarrow \operatorname{St}(B)$. Hence the formation of Steinberg groups is functorial. Additionally there are maps $\phi_{n}: \operatorname{St}_{n}(A) \rightarrow \mathrm{GL}_{n}(A)$ for any $n$, given by $x_{i j}(a) \mapsto e_{i j}(a)$; clearly its image is precisely the group $E_{n}(A)$ generated by elementary $n \times n$-matrices. Passing to the colimit, we obtain a $\operatorname{map} \phi: \operatorname{St}(A) \rightarrow E(A)$.

Definition 2.11. The second algebraic $K$-group $K_{2}(A)$ of the ring $A$ is defined as the kernel of the $\operatorname{map} \phi: \operatorname{St}(A) \rightarrow E(A)$.

That $K_{2}$ is a functor is a consequence of the functoriality of $\operatorname{St}(-)$ and $E(-)$. From the definition of $K_{2}$ we have the exact sequence

$$
1 \longrightarrow K_{2}(A) \longrightarrow \operatorname{St}(A) \xrightarrow{\phi} E(A) \longrightarrow 1
$$

The following theorem shows that this exact sequence is in fact a central extension of the group $E(A)$. In particular, $K_{2}(A)$ is an abelian group.
Theorem 2.12 (М씨기, Theorem 5.10]). The group $K_{2}(A)$ is the center of the Steinberg group $\operatorname{St}(A)$.

## $2.2 K_{2}$ of fields

The group $K_{2}(F)$ for $F$ a field has been extensively studied since the beginning of algebraic $K$-theory. If $F$ is a number field, $K_{2}(F)$ contains much information about the arithmetic of $F$. Later on we will see a few examples of this.

When it comes to computational questions on $K_{2}(F)$, one rarely operates with the original Definition 2.11. The first vast simplification on the description of $K_{2}(F)$ is due to Matsumoto (see Theorem 2.14 , which allows for the computation of $K_{2}$ of finite fields as well as $K_{2}$ of some number fields. In the computations in Chapter 3, we use only Matsumoto's theorem along with some number theory. To formulate Matsumoto's theorem, we first look at the general case of $K_{2}$ of a ring $A$ and later specialize to the case of a field.

In order to construct elements of $K_{2}(A)$ for a given ring $A$, consider two elements $a, b \in$ $E(A)$ such that the commutator $[a, b]=1$. Let $x, y$ denote representatives for $\phi^{-1}(a), \phi^{-1}(b)$, respectively. Then $[x, y]=x y x^{-1} y^{-1}$ lies in $\operatorname{St}(A)$, and by definition $\phi([x, y])=1$, so that $[x, y]$ is an element of $K_{2}(A)$. To see that the element $[x, y]$ is well defined, note that if $x^{\prime}$ is another representative for $\phi^{-1}(a)$, then $x^{\prime}=x z$ for some $z$ lying in the center of $\operatorname{St}(A)$. But then

$$
\begin{aligned}
{\left[x^{\prime}, y\right] } & =x^{\prime} y x^{\prime-1} y^{-1} \\
& =x z y z^{-1} x^{-1} y^{-1} \\
& =x y x^{-1} y^{-1} \\
& =[x, y],
\end{aligned}
$$

since $z$ commutes with every element of $\operatorname{St}(A)$. We therefore let $\left[\phi^{-1}(a), \phi^{-1}(b)\right.$ ] denote the element $[x, y]$ of $K_{2}(A)$ where $x$ and $y$ are any representatives for respectively $\phi^{-1}(a)$ and $\phi^{-1}(b)$.
Definition 2.13. Let $A$ be a ring and $u, v \in A^{\times}$. We define the Steinberg symbol $\{u, v\}$ to be the element $\left[\phi^{-1}\left(d_{12}(u)\right), \phi^{-1}\left(d_{13}(v)\right)\right] \in K_{2}(A)$, where

$$
d_{12}(u)=\left(\begin{array}{ccc}
u & 0 & 0 \\
0 & u^{-1} & 0 \\
0 & 0 & 1
\end{array}\right), \quad d_{13}(v)=\left(\begin{array}{ccc}
v & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & v^{-1}
\end{array}\right)
$$

The Steinberg symbols give rise to an antisymmetric bilinear map

$$
\{-,-\}: A^{\times} \times A^{\times} \longrightarrow K_{2}(A)
$$

satisfying the relation

$$
\{u, 1-u\}=1 \text { for all } u \in A^{\times} \text {with } 1-u \in A^{\times}
$$

Ros94, Lemma 4.2.14, Theorem 4.2.17]. When we restrict our attention to a field $F$, it turns out that the group $K_{2}(F)$ is generated by the symbols $\{u, v\}$ :

Theorem 2.14 (Matsumoto Mil71, §12]). The group $K_{2}(F)$ is the free abelian group on the generators $\{x, y\}$, where $x, y \in F^{\times}$, subject only to the following relations:

1. $\{x, 1-x\}=1$ for any $x \neq 1$;
2. $\left\{x_{1} x_{2}, y\right\}=\left\{x_{1}, y\right\}\left\{x_{2}, y\right\}$;
3. $\left\{x, y_{1} y_{2}\right\}=\left\{x, y_{1}\right\}\left\{x, y_{2}\right\}$.

Theorem 2.15 ([Mil71, Theorem 8.8]). There is an antisymmetric, bilinear pairing

$$
K_{1}(F) \times K_{1}(F) \longrightarrow K_{2}(F)
$$

mapping $(a, b) \in K_{1}(F) \times K_{1}(F)=F^{\times} \times F^{\times}$to $\{a, b\} \in K_{2}(F)$.
Let us proceed to describe the universal property of $K_{2}(F)$.
Definition 2.16. Let $F$ be a field and $G$ an abelian group written multiplicatively. A $G$-valued (Steinberg) symbol on $F$ is a bilinear map

$$
(-,-): F^{\times} \times F^{\times} \longrightarrow G
$$

such that $(x, 1-x)=1$ whenever $x \in F^{\times} \backslash\{1\}$.
By Matsumoto's theorem, any symbol $(-,-): F^{\times} \times F^{\times} \rightarrow G$ gives rise to a commutative diagram

$$
\begin{aligned}
& F^{\times} \times F^{\times} \xrightarrow{\{-,-\}} K_{2}(F) \\
& (-,-) \downarrow_{G}
\end{aligned}
$$

In other words, $K_{2}(F)$ is the universal object with respect to symbols on $F$ with values in abelian groups, and the corresponding Steinberg symbol $\{-,-\}$ is the universal symbol on $F$.

Lemma 2.17. Let $(-,-): F^{\times} \times F^{\times} \rightarrow G$ be a $G$-valued symbol on $F$. Then, for all $x, y \in F^{\times}$ we have
(i) $(x,-x)=1$;
(ii) $(x, y)^{-1}=(y, x)$;
(iii) $(x, x)=(x,-1)$.

Proof. First, note that $(1, x)=(1, x)(1, x)$ for any $x \in F^{\times}$by bilinearity, hence $(1, x)=1$. In particular, $(1,-1)=1$ so that property (i) holds for $x=1$. Therefore, assume that $x \neq 1$ and write $-x=(1-x)\left(1-x^{-1}\right)^{-1}$. We then have

$$
\begin{align*}
(x,-x) & =(x, 1-x)\left(x,\left(1-x^{-1}\right)^{-1}\right)  \tag{2.1}\\
& =\left(x,\left(1-x^{-1}\right)^{-1}\right)  \tag{2.2}\\
& =\left(x, 1-x^{-1}\right)^{-1}  \tag{2.3}\\
& =\left(x^{-1}, 1-x^{-1}\right)=1, \tag{2.4}
\end{align*}
$$

where the equality at $(2.3)$ holds since $\left(x, 1-x^{-1}\right)\left(x,\left(1-x^{-1}\right)^{-1}\right)=(x, 1)=1$. A similar computation shows that $\left(x, 1-x^{-1}\right)^{-1}=\left(x^{-1}, 1-x^{-1}\right)$.

The antisymmetry follows from (i):

$$
1=(x y,-x y)=(x,-x)(x, y)(y, x)(y,-y)=(x, y)(y, x) .
$$

Similarly we get property (iii):

$$
(x, x)=(x,-1)(x,-x)=(x,-1) .
$$

Theorem 2.18. For any finite field $F$, the group $K_{2}(F)$ is the trivial group.
Proof. Let $u$ denote a generator of the cyclic group $F^{\times}$. If $x, y \in F^{\times}$, write $x=u^{n}, y=u^{m}$. Then in $K_{2}(F)$ we have

$$
\{x, y\}=\left\{u^{n}, u^{m}\right\}=\{u, u\}^{n m}
$$

hence it suffices to show that $\{u, u\}=1$.
If the characteristic char $F=2$ we have $u=-u$, hence $\{u, u\}=\{u,-u\}=1$ by Lemma 2.17
(i). We may therefore assume char $F \neq 2$.

Suppose that $F$ has order $q$, where $q$ is a power of an odd prime. The set

$$
A:=F^{\times} \backslash F^{\times 2}
$$

of nonsquares of $F^{\times}$contains $(q-1) / 2$ elements and is in bijection with the set

$$
B:=\{1-v: v \in A\}
$$

Suppose that $A \cap B=\varnothing$. Since both $A$ and $B$ contain $(q-1) / 2$ elements we can then write $F^{\times}$ as the disjoint union of $A$ and $B$. But since 1 is a square, $1 \notin A$, and since $0 \notin A$ we have that $1 \notin B$. Hence $1 \notin A \cup B$, which contradicts the assumption that $F^{\times}=A \cup B$. We can therefore conclude that there are odd integers $i, j$ such that $u^{j}=1-u^{i}$, and thus $1=\left\{u^{i}, u^{j}\right\}=\{u, u\}^{i j}$. Writing $i j=2 k+1$ for $k$ an integer we therefore get that

$$
1=\{u, u\}^{2 k+1}=\{u, u\}^{2 k}\{u, u\}=\{u, u\}
$$

since $\{u, u\}^{2}=1$ by the antisymmetry.
We now return to the number theoretic setting. Let $F$ denote an algebraic number field and $v$ a place of $F$. Proposition 1.49 tells us that the local Hilbert symbol at $v$, defined in 1.48, is in fact a symbol in the sense of Definition 2.16. Thus the same holds for the regular Hilbert symbol of Definition 1.50 .

Definition 2.19. The global regular Hilbert symbol

$$
(-,-)^{\mathrm{reg}}: F^{\times} \times F^{\times} \longrightarrow \bigoplus_{v \in \mathrm{P}^{\mathrm{nc}}} \mu\left(F_{v}\right)^{\mathrm{reg}}
$$

is defined via the regular Hilbert symbols at the different places $v$ as

$$
(x, y)^{\mathrm{reg}}:=\left(\left(i_{v}(x), i_{v}(y)\right)_{v}^{\mathrm{reg}}\right)_{v \in \mathrm{P}^{\mathrm{nc}}}
$$

where $v$ runs over all noncomplex places of $F$.
Note that $(x, y)^{\text {reg }}$ indeed takes values in the direct sum, as $\left(i_{v}(x), i_{v}(y)\right)_{v}^{\mathrm{reg}}=1$ for almost all $v$.

With this definition $(-,-)^{\text {reg }}$ becomes a symbol on $F^{\times}$, and it therefore factors through $K_{2}(F)$. We denote by $h^{\text {reg }}$ the induced map $h^{\text {reg }}: K_{2}(F) \rightarrow \bigoplus_{v} \mu\left(F_{v}\right)^{\text {reg }}$ :

$$
\begin{aligned}
& F^{\times} \times F^{\times} \xrightarrow{\{-,-)^{\mathrm{reg}}}{ }_{(-,)^{2}} K_{2}(F) \\
& \bigoplus_{v} \mu\left(F_{v}\right)^{\mathrm{reg}}
\end{aligned}
$$

In a completely analogous manner we obtain the global Hilbert symbol,

$$
F^{\times} \times F^{\times} \ni(x, y) \longmapsto\left(\left(i_{v}(x), i_{v}(y)\right)_{v}\right)_{v \in \mathrm{P}^{\text {nc }}} \in \bigoplus_{v \in \mathrm{P}^{\text {nc }}} \mu\left(F_{v}\right)
$$

inducing a map

$$
h: K_{2}(F) \longrightarrow \bigoplus_{v \in \mathrm{P}^{\mathrm{nc}}} \mu\left(F_{v}\right)
$$

We aim to compare the maps $h$ and $h^{\text {reg }}$. To this end, we start out with a small observation before moving on to Moore's theorem.
Observation 2.20. The Hilbert symbol can be described as

$$
\left.\left(\frac{x, y}{v}\right)=i_{v}^{-1}\left(\left(i_{v}(x), i_{v}(y)\right)\right)^{\frac{m_{v}}{m}}\right)
$$

Proof. We have $\left(i_{v}(x), i_{v}(y)\right)_{v}^{m_{v} / m}=\left(i_{v}(x)^{m_{v} / m}, i_{v}(y)\right)_{v}$ by bilinearity, and

$$
\left(i_{v}(y), F_{v}\left(\sqrt[m_{v}]{i_{v}\left(x^{m_{v} / m}\right)}\right) / F_{v}\right)=\left(i_{v}(y), F_{v}\left(\sqrt[m]{i_{v}(x)}\right) / F_{v}\right)
$$

the right hand side being the image under the local reciprocity map in $L_{v} / F_{v}$, where $L:=$ $F\left(\sqrt[m]{i_{v}(x)}\right)$. Hence, letting $\sigma$ denote the Hasse symbol $\sigma:=\left(\frac{y, F(\sqrt[m]{x}) / F}{v}\right)$, we get

$$
i_{v}^{-1}\left(\left(i_{v}(x), i_{v}(y)\right)_{v}^{m_{v} / m}\right)=\frac{\sigma(\sqrt[m]{x})}{\sqrt[m]{x}}=\left(\frac{x, y}{v}\right)
$$

Now define a map

$$
\pi: \bigoplus_{v \in \mathrm{P}^{\mathrm{nc}}} \mu\left(F_{v}\right) \longrightarrow \mu(F)
$$

by

$$
\pi\left(\left(\zeta_{v}\right)_{v}\right):=\prod_{v} i_{v}^{-1}\left(\zeta_{v^{\frac{m v}{m}}}\right)
$$

The above observation together with the product formula for Hilbert symbols shows that the sequence

$$
K_{2}(F) \xrightarrow{h} \bigoplus_{v \in \mathrm{P}^{\text {nc }}} \mu\left(F_{v}\right) \xrightarrow{\pi} \mu(F) \longrightarrow 1
$$

is a complex. That the sequence is in fact exact is the statement of Moore's theorem, which is also known as the reciprocity uniqueness theorem. It says that the relation $\prod_{v}\left(\frac{x, y}{v}\right)=1$ is the only relation between the Hilbert symbols.

Theorem 2.21 (Moore Gra03, II Theorem 7.6]). The sequence

$$
K_{2}(F) \xrightarrow{h} \bigoplus_{v \in \mathrm{P}^{\text {nc }}} \mu\left(F_{v}\right) \xrightarrow{\pi} \mu(F) \longrightarrow 1
$$

is exact.
Definition 2.22 (Gra03, II 7.6.2]). The kernel of

$$
h: K_{2}(F) \rightarrow \bigoplus_{v \in \mathrm{P}^{\mathrm{nc}}} \mu\left(F_{v}\right)
$$

is called the wild kernel, and is denoted $W K_{2}(F)$. The kernel of $h^{\text {reg }}$ is called the regular, or tame kernel, and is denoted by $R_{2}^{\text {ord }}(F)$.

Clearly we have $W K_{2}(F) \subseteq R_{2}^{\text {ord }}(F)$. In the notations of Moore's theorem we have a commutative diagram Gra03, II Theorem 7.6]


Nontriviality of the regular kernel occurs for example when the prime 2 splits in the extension $F / \mathbf{Q}$. This is due to the existence of a so-called "wild" 2-adic Hilbert symbol on such a field; see Example 3.14. On the other hand, nontriviality of the wild kernel means that there are symbols on $F$ which do not come from Hilbert symbols-i.e., symbols that are detected by global, but not by local class field theory. This means that there are global symbols defined on $F$ which cannot be expressed in terms of the Hilbert symbols $\left(\frac{x, y}{v}\right)$ at different places $v$. Such symbols are called exotic symbols. As Gras notes in [Gra03, p.214], at the present time no one has managed to describe such a symbol explicitly. We will briefly revisit regular and wild kernels in Chapter 3 .

Remark 2.23. The regular kernel $R_{2}^{\text {ord }}(F)$ was identified by Quillen as $K_{2}\left(\mathcal{O}_{F}\right)$. The lower exact sequence in the diagram above is a special case of the localization sequence in higher algebraic $K$-theory Qui73, p.29]

$$
\cdots \longrightarrow K_{n+1}(F) \longrightarrow \bigoplus_{\mathfrak{p} \in \operatorname{Spec}(A) \backslash\{0\}} K_{n}(k(\mathfrak{p})) \longrightarrow K_{n}(A) \longrightarrow K_{n}(F) \longrightarrow \cdots
$$

for $A$ a Dedekind ring and $F$ the field of fractions of $A$.

### 2.3 Milnor $K$-groups and Witt theory

In Chapter 5 we will discuss the motivic homotopy groups $\pi_{* \alpha} \mathbf{1}$ of the sphere spectrum. These groups were identified by Morel in Mor04a as the so-called Milnor-Witt $K$-groups (defined in Definition 5.6), introduced by Morel in the same article. Milnor-Witt $K$-theory is in some sense a blend of Milnor $K$-theory and Witt theory, so let us here briefly review the definition and basic properties of these objects.

Let $F$ be a field. The Milnor $K$-groups were first defined in Milnor's seminal paper Mil70, inspired by Matsumoto's description of $K_{2}(F)$.

For any $n \geq 0$, let $\left(F^{\times}\right)^{\otimes n}:=F^{\times} \otimes_{\mathbf{z}} \cdots \otimes_{\mathbf{z}} F^{\times}$be the $n$-fold tensor product of $F^{\times}$, where we define $\left(F^{\times}\right)^{\otimes 0}:=\mathbf{Z}$.

Definition 2.24. Let $n \geq 0$. The $n$-th Milnor $K$-group $K_{n}^{M}(F)$ of $F$ is the group

$$
K_{n}^{M}(F):=\left(F^{\times}\right)^{\otimes n} /\left\langle a_{1} \otimes \cdots \otimes a_{n}: a_{j} \in F^{\times}, a_{j}+a_{j+1}=1 \text { for some } j\right\rangle .
$$

We write $\left\{a_{1}, \ldots, a_{n}\right\}$ for the image of $a_{1} \otimes \cdots \otimes a_{n} \in\left(F^{\times}\right)^{\otimes n}$ in $K_{n}^{M}(F)$.
The Milnor K-theory $K_{*}^{M}(F)$ is the graded ring $K_{*}^{M}(F):=\bigoplus_{n \geq 0} K_{n}^{M}(F)$, where $K_{n}^{M}(F)$ consists of the homogeneous elements of degree $n$.

Example 2.25. Theorem 2.18 shows that $K_{n}^{M}\left(\mathbf{F}_{q}\right)=0$ for any $n \geq 2$.
By Example 2.2 and Example 2.7. $K_{n}^{M}(F)$ agrees with the algebraic $K$-theory $K_{n}(F)$ for $n=0,1,2$.

In Mil70, Milnor establishes an analog of the regular map $h^{\text {reg }}: K_{2}(F) \rightarrow \bigoplus_{v \in \operatorname{Pin}_{F}^{\text {nc }}} k(v)^{\times}$ in Milnor $K$-theory (recall that $k(v)^{\times} \cong K_{1}(k(v))$ :
Theorem 2.26 (Mil70, Lemma 2.1]). Let $F$ be field. For each discrete valuation $v$ on $F$ and each $n \geq 1$ there exists a unique homomorphism

$$
\partial_{v}: K_{n}^{M}(F) \longrightarrow K_{n-1}^{M}(k(v))
$$

such that

$$
\partial_{v}\left(\left\{\pi_{v}, u_{2}, \ldots, u_{n}\right\}\right)=\left\{\bar{u}_{2}, \ldots, \bar{u}_{n}\right\} \in K_{n-1}^{M}(k(v))
$$

and

$$
\partial_{v}\left(\left\{u_{1}, u_{2}, \ldots, u_{n}\right\}\right)=0
$$

for all units $u_{i} \in \mathcal{O}_{v}^{\times}$.
Using these maps, Milnor showed in Mil70 that there are split exact sequences

$$
0 \longrightarrow K_{n}^{M}(F) \longrightarrow K_{n}^{M}(F(t)) \xrightarrow{\oplus_{\mathfrak{p}} \partial_{\mathfrak{p}}} \bigoplus_{\mathfrak{p} \in \operatorname{Spec}(F[t]) \backslash\{0\}} K_{n-1}^{M}(F[t] / \mathfrak{p}) \longrightarrow 0
$$

where the direct sum ranges over all nonzero prime ideals of $F[t]$. Again, the method of computation is attributed to Tate.

### 2.3.1 The Witt ring

The following is based on MH73, Sch85, Wei13. Let $A$ be a ring and let $M$ be an $A$-module. A symmetric bilinear form over $A$ is an $A$-bilinear map

$$
\beta: M \times M \longrightarrow A
$$

such that $\beta(x, y)=\beta(y, x)$ for all $x, y \in M$. The pair $(M, \beta)$ is called a symmetric inner product space over $A$ if in addition the following conditions are satisfied:

- $M$ is finitely generated and projective over $A$;
- the maps $M \rightarrow M^{*}=\operatorname{Hom}_{A}(M, R)$ given by $x \mapsto \beta(x,-)$ and $y \mapsto \beta(-, y)$ are bijective (i.e., the form is nondegenerate).

A map $f:(M, \beta) \rightarrow\left(M^{\prime}, \beta^{\prime}\right)$ of inner product spaces over $A$ is an $A$-module map $f: M \rightarrow M^{\prime}$ such that $\beta^{\prime}(f(x), f(y))=\beta(x, y)$. Thus we can form the category $\operatorname{SBil}(A)$ of inner product spaces over $A$. An isomorphism in $\operatorname{SBil}(A)$ is called an isometry.

The category $\operatorname{SBil}(F)$ becomes symmetric monoidal under orthogonal sum defined by

$$
(M, \beta) \oplus\left(M^{\prime}, \beta^{\prime}\right):=\left(M \oplus M^{\prime}, \beta \oplus \beta^{\prime}\right),
$$

where $\left(\beta \oplus \beta^{\prime}\right)\left(x \oplus x^{\prime}, y \oplus y^{\prime}\right):=\beta(x, y)+\beta^{\prime}\left(x^{\prime}, y^{\prime}\right)$. Moreover, tensor product over $A$ provides $\operatorname{SBil}(A)^{\text {iso }}$ with the structure of a semiring. Hence the Grothendieck group $K_{0} \operatorname{SBil}(A)$ becomes a ring, called the Grothendieck-Witt ring of $A$. We write

$$
G W(A):=K_{0} \operatorname{SBil}(A)
$$

If $M \in \mathscr{P}(A)$, the associated hyperbolic space is defined as

$$
\mathbf{H}(M):=\left(M \oplus M^{*}, \mathbf{h}_{M}\right),
$$

where

$$
\mathbf{h}_{M}((x, \alpha),(y, \beta)):=\beta(x)+\alpha(y) .
$$

The adjoint map $\widehat{\mathbf{h}}_{M}: M \oplus M^{*} \rightarrow M^{*} \oplus M$ of $\mathbf{h}$ is then given by $\left(\begin{array}{cc}0 & \mathrm{id}_{M^{*}} \\ \mathrm{id} & 0\end{array}\right)$.
Definition 2.27. The Witt ring $W(A)$ of $A$ is the quotient of $G W(A)$ by the ideal generated by hyperbolic spaces.

There is an augmentation $\epsilon: G W(A) \rightarrow \mathbf{Z}$; this is simply the rank map. We let $\widehat{I}$ denote the kernel of $\epsilon$. The augmentation map induces a map on the Witt ring

$$
\epsilon_{\mathbf{Z} / 2}: W(A) \longrightarrow \mathbf{Z} / 2
$$

its kernel $I(A)$ is called the fundamental ideal. It is the ideal of even-dimensional forms. We have $I(A) \cong \widehat{I}$ as abelian groups, and the following diagram summarizes the situation:


Example 2.28. The case when $A$ is a field $F$ of characteristic not 2 is perhaps the most studied. Any element $u \in F^{\times}$gives rise to a symmetric inner product space $\langle u\rangle$, where $V=F$ is the underlying vector space and $\beta(x, y):=u x y$ is the bilinear form. We have $\langle u\rangle \cong\left\langle u^{\prime}\right\rangle$ if and only if $u^{\prime}=a^{2} u$ for some $a \in F^{\times}$. We will write $\left\langle u, u^{\prime}\right\rangle$ for the sum $\langle u\rangle \oplus\left\langle u^{\prime}\right\rangle$, and we have $\langle u\rangle \otimes\left\langle u^{\prime}\right\rangle=\left\langle u u^{\prime}\right\rangle$.

The hyperbolic plane $\mathbf{H}$ is the inner product space $\left(F^{2}, \beta\right)$, where $\beta$ is given by the matrix $\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$. By diagonalizing we have $\mathbf{H}=\langle 1\rangle \oplus\langle-1\rangle$. The hyperbolic plane generates an ideal $\langle\mathbf{H}\rangle$, which is in fact equal to $\mathbf{Z} \cdot \mathbf{H}$. The ideal $\langle\mathbf{H}\rangle$ is the ideal of hyperbolic spaces, and we have $W(F)=G W(F) /\langle\mathbf{H}\rangle$.

The need to consider Witt rings of arbitrary rings becomes apparent shortly: We shall see in Example 2.35 and Theorem 2.36 that in order to describe the Witt ring of number fields, we are forced to consider the Witt ring of the ring of integers as well as $W\left(\mathbf{F}_{q}\right)$ for $q$ even.

The rings $G W(F)$ and $W(F)$ can be described by generators and relations as follows:
Theorem 2.29 ([Lam05, II Theorem 4.1]). The Grothendieck-Witt ring is the free abelian group generated by symbols $\langle u\rangle$ for $u \in F^{\times}$, subject to the following relations.

1. $\langle 1\rangle=1$,
2. $\langle u\rangle\langle v\rangle=\langle u v\rangle$,
3. $\langle u\rangle+\langle v\rangle=\langle u+v\rangle(1+\langle u v\rangle)$
whenever they make sense. Requiring the additional relation
4. $1+\langle-1\rangle=0$
one recovers the Witt ring.
The so-called Pfister forms are elements of the Witt ring of particular interest.
Definition 2.30. Let $a_{1}, \ldots, a_{n} \in F^{\times}$. The corresponding $n$-fold Pfister form is the $2^{n}$ dimensional form

$$
\left\langle\left\langle a_{1}, \ldots, a_{n}\right\rangle\right\rangle:=\prod_{i=1}^{n}\left\langle 1,-a_{i}\right\rangle=\left(\langle 1\rangle-\left\langle a_{1}\right\rangle\right) \cdots\left(\langle 1\rangle-\left\langle a_{n}\right\rangle\right) .
$$

| $q$ | $W\left(\mathbf{F}_{q}\right)$ | $I\left(\mathbf{F}_{q}\right)$ | $I^{2}\left(\mathbf{F}_{q}\right)$ |
| :---: | ---: | ---: | ---: |
| even | $\mathbf{Z} / 2$ | 0 | 0 |
| $\equiv 1(\bmod 4)$ | $\mathbf{Z} / 2 \oplus \mathbf{Z} / 2$ | $\mathbf{Z} / 2$ | 0 |
| $\equiv 3(\bmod 4)$ | $\mathbf{Z} / 4$ | $\mathbf{Z} / 2$ | 0 |

Table 2.1: The Witt ring and the fundamental ideal of finite fields.

Proposition 2.31 (【Lam05, X Proposition 1.2]). The n-th power $I^{n}(F)$ of the fundamental ideal of $F$ is generated as an abelian group by all $n$-fold Pfister forms.

There is a group homomorphism $F^{\times} / F^{\times 2} \rightarrow I(F) / I^{2}(F)$ induced by $u \mapsto\langle\langle u\rangle\rangle$; Pfister showed that this is an isomorphism:

Theorem 2.32 (Pfister MH73, Theorem 5.2]). There is a canonical isomorphism $F^{\times} / F^{\times 2} \cong$ $I(F) / I^{2}(F)$.

Example 2.33. The Witt ring of a finite field $\mathbf{F}_{q}$ is described in MH73, p.87]; see Table 2.1 . In particular, we have $I^{2}\left(\mathbf{F}_{q}\right)=0$ for any prime power $q$.
Example 2.34. Let us consider the case of ordered fields $F$. Recall the inertia theorem of Jacobi and Sylvester MH73, p.61], which says that any inner product space $X$ over an ordered field $F$ is isomorphic to an orthogonal sum

$$
X \cong X^{+} \oplus X^{-}
$$

where $X^{+}$is positive definite (i.e., $x \cdot x$ is positive for all nonzero $x \in X^{+}$), and $X^{-}$is negative definite. Moreover, the ranks of $X^{+}$and $X^{-}$are invariants up to isomorphism.

Suppose $P$ is an ordering in $F$. In light of the inertia theorem one defines the signature $\sigma_{P}(X)$ of the inner product space $X$ at $P$ as the integer

$$
\sigma_{P}(X):=\operatorname{rk}\left(X^{+}\right)-\operatorname{rk}\left(X^{-}\right)
$$

For the rank 1 inner product space $\langle u\rangle$-where $u \in F^{\times}$-we have $\sigma_{P}(\langle u\rangle)=\operatorname{sgn}(u)$, the sign of $u$. One can show that the signature induces a well defined homomorphism on the Witt ring, and that if $F$ is an ordered field in which all positive elements are squares, then $\sigma_{P}: W(F) \rightarrow \mathbf{Z}$ is an isomorphism (see [MH73, p.63]). In particular, we have $W(\mathbf{R}) \cong \mathbf{Z}$.

Example 2.35. We have the following description of the Witt ring of $\mathbf{Q}$ :

$$
W(\mathbf{Q}) \cong W(\mathbf{Z}) \oplus \bigoplus_{p \text { prime }} W\left(\mathbf{F}_{p}\right)
$$

The original proof uses the Hasse-Minkowski theorem for quadratic forms over Q. In MH73, pp.88-89], Milnor gives an alternative proof not making use of the Hasse-Minkowski principle. This proof is quite similar to Tate's computation of $K_{2}(\mathbf{Q})$, which we will come back to in Chapter 3

There are analogs of the tame symbols also in Witt theory: If $F$ is a number field, each finite place $v \in \mathrm{Pl}_{F}$ and each choice of uniformizer $\pi_{v}$ gives rise to an additive homomorphism

$$
\partial_{v}^{\pi_{v}}: W(F) \longrightarrow W(k(v))
$$

sending a generator $\left\langle u \pi_{v}^{n}\right\rangle$ of $W(F)$ to either $\langle\bar{u}\rangle$ or 0 according as $n \equiv 1(\bmod 2)$ or $n \equiv 0$ $(\bmod 2)$. The map $\partial_{v}^{\pi_{v}}$ is unique up to choice of uniformizer $\pi_{v}$.

Theorem 2.36 (MH73, IV Corollary 3.3]). If $F$ is a number field, there is an exact sequence

$$
0 \longrightarrow W\left(\mathcal{O}_{F}\right) \longrightarrow W(F) \xrightarrow{\oplus_{v} \partial_{v}^{\pi_{v}}} \bigoplus_{v \in \mathrm{Pl}_{0}} W(k(v))
$$

The extent to which $\partial:=\bigoplus_{v} \partial_{v}^{\pi_{v}}$ fails to be surjective is measured by the ideal class group:
Proposition 2.37 ([МН73, p.94]). There is an exact sequence

$$
W(F) \xrightarrow{\partial} \bigoplus_{v \in \mathrm{Pl}_{0}} W(k(v)) \longrightarrow \mathrm{Cl}(F) / \mathrm{Cl}(F)^{2} \longrightarrow 0,
$$

where the map $\bigoplus_{v} W(k(v)) \rightarrow \mathrm{Cl}(F) / \mathrm{Cl}(F)^{2}$ is defined by sending each generator $\langle\bar{u}\rangle$ of $W(k(v))$ to the ideal class of $\mathfrak{p}_{v}$ modulo $\mathrm{Cl}(F)^{2}$.
Example 2.38. If $F$ equals $\mathbf{Q}(i)$ or $\mathbf{Q}(\sqrt{-3})$ then $\mathrm{Cl}(F)=0$, hence we have exact sequences

$$
0 \longrightarrow W\left(\mathcal{O}_{F}\right) \longrightarrow W(F) \xrightarrow{\partial} \bigoplus_{v \in \mathrm{Pl}_{0}} W(k(v)) \longrightarrow 0
$$

In MH73, p.96], Milnor and Husemoller describe the groups $W\left(\mathcal{O}_{\mathbf{Q}(\sqrt{-d})}\right)$ for $d=1,2, \ldots, 15$. In our case, the result is

$$
W(\mathbf{Z}[i]) \cong \mathbf{Z} / 2 \oplus \mathbf{Z} / 2 ; \quad W(\mathbf{Z}[\omega]) \cong \mathbf{Z} / 4
$$

### 2.3.2 Relation to Milnor $K$-theory

Let $k_{n}^{M}(F)$ denote the mod 2 Milnor $K$-theory

$$
k_{n}^{M}(F):=K_{n}^{M}(F) / 2 K_{n}^{M}(F) .
$$

Thus we have $k_{0}^{M}(F)=\mathbf{Z} / 2$ and $k_{1}^{M}(F)=F^{\times} / F^{\times 2}$, the square classes of elements in $F^{\times}$. By definition of the fundamental ideal we have $W(F) / I(F) \cong k_{0}^{M}(F)$, and by Pfister's Theorem 2.32 we have $I(F) / I^{2}(F) \cong k_{1}^{M}(F)$.

In Mil70, Milnor establishes a unique homomorphism

$$
s_{n}: k_{n}^{M}(F) \longrightarrow I^{n}(F) / I^{n+1}(F)
$$

for each $n \geq 0$, induced by

$$
\left\{a_{1}, \ldots, a_{n}\right\} \longmapsto\left\langle\left\langle a_{1}, \ldots, a_{n}\right\rangle\right\rangle\left(\bmod I^{n+1}(F)\right) .
$$

He proved that $s_{n}$ is surjective for all values of $n$ and announced the famous Milnor conjecture, asking whether $s_{n}$ is always bijective.

Before Milnor's article [Mil70], Bass and Tate had discovered some connections between $K_{2}$ and Galois cohomology Tat76b]. Let $F$ be a field of characteristic not 2 , and let $F_{\text {sep }}$ denote its separable closure. We then have the Kummer sequence

$$
1 \longrightarrow \mu_{2} \longrightarrow F_{\mathrm{sep}}^{\times} \xrightarrow{(-)^{2}} F_{\mathrm{sep}}^{\times} \longrightarrow 1
$$

The start of the corresponding long exact sequence in Galois cohomology takes the form

$$
0 \longrightarrow \mathbf{Z} / 2 \longrightarrow F^{\times} \xrightarrow{.2} F^{\times} \xrightarrow{\delta} H_{\text {Gal }}^{1}(F ; \mathbf{Z} / 2) \longrightarrow 0 .
$$

Here $H_{\text {Gal }}^{n}(F ; M):=H^{n}\left(\operatorname{Gal}\left(F_{\text {sep }} / F\right) ; M\right)=H_{\text {ett }}^{n}(F ; M)$, and $H_{\text {Gal }}^{1}\left(F ; F_{\text {sep }}^{\times}\right)=0$ by the Hilbert-Speiser-Noether Theorem 90 [Ser79, Ch.X]. Hence there is an isomorphism

$$
k_{1}^{M}(F) \cong H_{\text {Gal }}^{1}(F ; \mathbf{Z} / 2)
$$

Theorem 2.39 ([Tat76b, Theorem 3.1]). The isomorphism

$$
\{a\} \mapsto \delta(a): k_{1}^{M}(F) \longrightarrow H_{\text {Gal }}^{1}(F ; \mathbf{Z} / 2)
$$

extends uniquely to a ring homomorphism

$$
h_{F}: k_{*}^{M}(F) \longrightarrow H_{\text {Gal }}^{*}(F ; \mathbf{Z} / 2)
$$

satisfying $h_{F}(\{a, b\})=\delta(a) \smile \delta(b)$.

In 1996-by the advent of motivic homotopy theory-Orlov-Vishik-Voevodsky brought clarity to all the connections mentioned above. Indeed, they settled the Milnor conjecture by showing that $s_{*}$ as well as $h_{*}$ are graded ring isomorphisms OVV07

where $\operatorname{Gr}_{I}(W(F)):=W(F) / I(F) \oplus I(F) / I^{2}(F) \oplus \cdots$ is the associated graded Witt ring. Every known proof of this fact uses motivic homotopy theory.

## CHAPTER 3

## On $K_{2}$ of Some Imaginary Quadratic Number Fields

In this chapter we look at the $K$-theory of the two imaginary quadratic number fields $\mathbf{Q}(\sqrt{-1})$ and $\mathbf{Q}(\sqrt{-3})$. Specifically, we compute $K_{2}$ of these fields and look at how the result for $\mathbf{Q}(\sqrt{-1})$ is connected with the law of biquadratic reciprocity.

In Tate's computation of $K_{2}(\mathbf{Q})$ he defined a filtration

$$
L_{2} \subseteq L_{3} \subseteq L_{5} \subseteq \cdots \subseteq L_{p} \subseteq \cdots \subseteq K_{2}(\mathbf{Q})
$$

of $K_{2}(\mathbf{Q})$, where for each prime $p$,

$$
L_{p}:=\langle\{x, y\}: x, y \in \mathbf{Z} \backslash\{0\}, 1 \leq| x|,|y| \leq p\rangle .
$$

Tate used induction and the Euclidean algorithm to show that $L_{p} \cong \mu_{2} \bigoplus_{2<\ell \leq p} \mathbf{F}_{\ell}^{\times}$, yielding

$$
K_{2}(\mathbf{Q}) \cong \mu_{2} \oplus \bigoplus_{p \neq 2} \mathbf{F}_{p}^{\times}
$$

We aim to follow the same technique, only with slight modifications. In particular, the filtration has to be defined differently due to the fact that there may be two primes with the same norm in these number fields. Since we want to consider each prime separately, the norm is not sufficient to distinguish between them.

Now the groups $K_{2}(\mathbf{Q}(\sqrt{-1}))$ and $K_{2}(\mathbf{Q}(\sqrt{-3}))$ are of course known, e.g., from BT73, p.429] where Tate computes the second $K$-group of the first six imaginary quadratic number fields (ordered by their discriminants). Tate's method of computation is based on a few general results; here we will instead - as mentioned-give a direct computation based on the calculation of $K_{2} \mathbf{Q}$ mentioned above. This approach should work for all Euclidean imaginary quadratic number fields, but not in general because of the need of an Euclidean algorithm. However, the computation exemplifies how the arithmetic of the number field reveals itself in the $K$-group.

### 3.1 The group $K_{2}(\mathbf{Q}(\sqrt{-1}))$

Recall the situation from Chapter 2 for the global regular Hilbert symbol:


Here we specialize to the case when $F=\mathbf{Q}(i)$, where $i^{2}=-1$. Thus we have $\mathcal{O}_{\mathbf{Q}(i)}=\mathbf{Z}[i]$, the Gaussian integers. The map $h^{\text {reg }}: K_{2}(\mathbf{Q}(i)) \rightarrow \bigoplus_{v} \mu\left(\mathbf{Q}(i)_{v}\right)^{\text {reg }}$ will be our primary object of study.

## 3. On $K_{2}$ of Some Imaginary Quadratic Number Fields



Figure 3.1: Gaussian integers lying in the same relative position in the squares defined by $\pi$ are congruent modulo $\pi$.

First recall-as mentioned in Chapter 1 - that if $v$ is a finite place of $\mathbf{Q}(i)$ above a prime $p \in \mathbf{Z}$, then $\mathbf{Q}(i)_{v}=\mathbf{Q}_{p}(i)$. Secondly, recall that we in Proposition 1.20 have described the $\operatorname{groups}\left(\mathbf{Z}[i] / \mathfrak{p}_{v}\right)^{\times} \cong k(v)^{\times} \cong \mu\left(\mathbf{Q}(i)_{v}\right)^{\text {reg }}$, so that we have a map

$$
h^{\mathrm{reg}}: K_{2}(\mathbf{Q}(i)) \longrightarrow \bigoplus_{\substack{p \text { prime } \\ p \equiv 1(\bmod 4)}}\left(\mathbf{F}_{p}^{\times}\right)^{2} \oplus \bigoplus_{\substack{p \text { prime } \\ p \equiv 3(\bmod 4)}} \mathbf{F}_{p^{2}}^{\times}
$$

induced by the global regular Hilbert symbol. We aim to show that this is an isomorphism. To begin with we need a few preliminary lemmas.

Lemma 3.1. For any finite place $v$ of $\mathbf{Q}(i)$, there is a complete system of representatives $\left\{\alpha_{j}\right\}_{j=1, \ldots, \mathrm{Nm}\left(\pi_{v}\right)}$ for $k(v)$ such that $\left|\alpha_{j}\right| \leq\left|\pi_{v}\right| / \sqrt{2}$ for each $j$ (where $|\alpha|$ is the standard absolute value of $\alpha$, satisfying $\left.|\alpha|^{2}=\operatorname{Nm}(\alpha)\right)$.

Proof. This is the Euclidean algorithm on the Gaussian integers and follows from Table 1.1 in Chapter 1, but we write out the details for convenience. Geometrically we can picture the situation as follows (see Figure 3.1): Given $\alpha \in \mathbf{Z}[i] \backslash\{0\}$ and $\pi$ a prime element, consider the lattice spanned by unit multiples of $\pi$. Then $\alpha$ lies in one of the squares in the lattice, with distance at most $|\pi| / \sqrt{2}$ from a corner $\gamma \pi$. If $\beta:=\alpha-\gamma \pi$ then $\beta \equiv \alpha(\bmod \pi)$ and $|\beta| \leq|\pi| / \sqrt{2}$.

Lemma 3.2. The group $K_{2}(\mathbf{Q}(i))$ is generated by elements $\{\alpha, \beta\}$ for which $\alpha, \beta \in \mathbf{Z}[i] \backslash\{0\}$.
Proof. Let $\{x, y\} \in K_{2}(\mathbf{Q}(i))$, then we may write $x=\alpha / \beta, y=\alpha^{\prime} / \beta^{\prime}$ for $\alpha, \alpha^{\prime}, \beta, \beta^{\prime} \in \mathbf{Z}[i]$ nonzero. Using the relations in Lemma 2.17 and bilinearity we then have

$$
\begin{aligned}
\{x, y\} & =\left\{\alpha \beta^{-1}, \alpha^{\prime} \beta^{\prime-1}\right\} \\
& =\left\{\alpha, \alpha^{\prime}\right\}\left\{\alpha, \beta^{\prime-1}\right\}\left\{\beta^{-1}, \alpha^{\prime}\right\}\left\{\beta^{-1}, \beta^{\prime-1}\right\} \\
& =\left\{\alpha, \alpha^{\prime}\right\}\left\{\beta^{\prime}, \alpha\right\}\left\{\alpha^{\prime}, \beta\right\}\left\{\beta, \beta^{\prime}\right\},
\end{aligned}
$$

where we have used that $\left\{\gamma^{-1}, \delta\right\}=\{\gamma, \delta\}^{-1}=\{\delta, \gamma\}$.
Remark 3.3. The same proof holds for $K_{2}(F)$ whenever $F$ is the quotient field of some integral domain.

Definition 3.4. Let $\Lambda_{\infty}$ denote the following subgroup of $K_{2}(\mathbf{Q}(i))$ :

$$
\Lambda_{\infty}:=\left\langle\{u, v\}: u, v \in \mathbf{Z}[i]^{\times}=\mu_{4}\right\rangle .
$$

Lemma 3.5. The group $\Lambda_{\infty}$ is trivial.
Proof. The generators of $\Lambda_{\infty}$ are of the form $\{u, v\}$ with $u, v \in \mathbf{Z}[i]^{\times}$. First, note that if either $u=1$ or $v=1$, then $\{u, v\}=1$ by the bilinearity. We consider the remaining cases:

- We already know that $\{i,-i\}=1$ by the symbol properties. Therefore also $\{-i, i\}=$ $\{i,-i\}^{-1}=1$.
- Again by the relations in 2.17 we have $\{i,-1\}=\{i, i\}$. But $\{i, i\}\{i, i\}=\{i,-1\}=\{i, i\}$, hence $\{i, i\}=1=\{i,-1\}$ and also $\{-1, i\}=1$.
- For the element $\{-1,-1\}$ we now get $\{-1,-1\}=\left\{i^{2},-1\right\}=\{i,-1\}\{i,-1\}=1$.
- Similarly as above we have that $\{-i,-1\}=\{-i,-i\}$ and

$$
\{-i,-i\}\{-i,-i\}=\{-i,-1\}=\{-i,-i\}
$$

hence $\{-i,-i\}=1$ and therefore also $\{-i,-1\}=\{-1,-i\}=1$.
We aim to define a filtration of $K_{2}(\mathbf{Q}(i))$ indexed by the Gaussian primes. To start out, let us list the finite places of $\mathbf{Q}(i)$ by increasing norm:

$$
\mathrm{Pl}_{0}=\left\{v_{1}, v_{2}, \ldots \mid \mathrm{Nm}\left(\pi_{v_{n}}\right) \leq \operatorname{Nm}\left(\pi_{v_{n+1}}\right) \quad \forall n\right\} .
$$

Definition 3.6. Let $\mathrm{Pl}_{0}$ be enumerated as above. For all $n \geq 1$, let

$$
S_{n}:=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\} \subseteq \mathrm{Pl}_{0}
$$

Following the notation of Tate in [BT73], let $K_{2}^{S_{n}} \mathbf{Q}(i)$ denote the following subgroup of $K_{2} \mathbf{Q}(i)$ :

$$
K_{2}^{S_{n}} \mathbf{Q}(i):=\left\langle\{\alpha, \beta\} \in K_{2} \mathbf{Q}(i): \alpha, \beta \in \mathbf{Z}[i]_{S_{n}}^{\times}\right\rangle
$$

where

$$
\mathbf{Z}[i]_{S_{n}}=\mathbf{Z}\left[i, \frac{1}{\pi_{v_{1}}}, \ldots, \frac{1}{\pi_{v_{n}}}\right]
$$

is the ring of Gaussian $S_{n}$-integers.
The groups $K_{2}^{S_{n}} \mathbf{Q}(i)$ then form a filtered system and thus we have $K_{2} \mathbf{Q}(i)=\varliminf_{n} K_{2}^{S_{n}} \mathbf{Q}(i)$.
Theorem 3.7. The global regular Hilbert symbol induces an isomorphism

$$
K_{2}(\mathbf{Q}(i)) \cong \bigoplus_{v \in \mathrm{P}_{\mathbf{Q}(i)}^{\mathrm{ncc}}} k(v)^{\times} \cong \bigoplus_{\substack{p \text { prime } \\ p \equiv 1(\bmod 4)}}\left(\mathbf{F}_{p}^{\times}\right)^{2} \oplus \bigoplus_{\substack{p \text { prime } \\ p \equiv 3(\bmod 4)}} \mathbf{F}_{p^{2}}^{\times}
$$

Proof. We will use induction to show that for any $n \geq 1$, the restriction of $h^{\text {reg }}$ to $K_{n}^{S_{n}} \mathbf{Q}(i)$ yields an isomorphism $K_{n}^{S_{n}} \mathbf{Q}(i) \cong \bigoplus_{j=1}^{n} k\left(v_{j}\right)^{\times}$. To begin with we must show that $K_{2}^{S_{1}} \mathbf{Q}(i)$ is trivial. Note that $S_{1}=\left\{v_{1}\right\}$, where $v_{1}$ is the place above $2 \in \mathbf{Z}$, and that the Gaussian $S_{1}$-units consist of the torsion part $\mu_{4}$, and a free part generated by $1+i$.

Using bilinearity, Lemma 3.5 and the relation $\{x,-1\}=\{x, x\}$, we need only consider the generating symbols $\{\alpha, \beta\} \in \overline{K_{2}^{S_{1}}} \mathbf{Q}(i)$ where $\alpha, \beta=i, 1+i$. From the proof of Lemma 3.5 we already know that $\{i, i\}=1$, thus it remains to check the elements $\{1+i, i\},\{i, 1+i\}$ and $\{1+i, 1+i\}$. But by the Steinberg relation we have

$$
\{-i, 1+i\}=1
$$

Hence

$$
1=\{-i, 1+i\}=\{-1,1+i\}\{i, 1+i\}
$$

and therefore

$$
\{1+i, i\}=\{-1,1+i\}=\{1+i, 1+i\}
$$

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But then $\{1+i, i\}=\{i, 1+i\}$ and

$$
\{i, 1+i\}=\{1+i, i\}=\{-1,1+i\}=\{i, 1+i\}\{i, 1+i\}
$$

hence $\{i, 1+i\}=1$. In total we have

$$
1=\{i, 1+i\}=\{1+i, i\}=\{1+i, 1+i\}
$$

so $K_{2}^{S_{1}} \mathbf{Q}(i)$ is trivial-in other words $K_{2}^{S_{1}} \mathbf{Q}(i) \cong k\left(v_{1}\right)^{\times}$.

Lemma 3.8. For any $n \geq 1$, the quotient group $K_{2}^{S_{n+1}} \mathbf{Q}(i) / K_{2}^{S_{n}} \mathbf{Q}(i)$ is isomorphic to $k\left(v_{n+1}\right)^{\times}$
Proof. Define a map

$$
\phi: k\left(v_{n+1}\right)^{\times} \longrightarrow K_{2}^{S_{n+1}} \mathbf{Q}(i) / K_{2}^{S_{n}} \mathbf{Q}(i)
$$

by

$$
\bar{\alpha} \longmapsto\{\alpha, \pi\} \quad\left(\bmod K_{2}^{S_{n}} \mathbf{Q}(i)\right),
$$

where $\pi:=\pi_{v_{n+1}}$, and where we may assume $\operatorname{Nm}(\alpha) \leq \operatorname{Nm}(\pi) / 2$ by Lemma 3.1. To show that $\phi$ is a well defined homomorphism, assume $\alpha \beta \equiv \gamma(\bmod \pi)$, say $\alpha \beta=\gamma+\delta \pi$, where $|\alpha|,|\beta|,|\gamma| \leq|\pi| / \sqrt{2}$. Then

$$
|\delta \pi| \leq|\alpha \beta|+|\gamma| \leq \frac{|\pi|^{2}}{2}+\frac{|\pi|}{\sqrt{2}}
$$

thus $|\delta| \leq|\pi| / 2+1 / \sqrt{2}$, which is less than $|\pi|$ since $\operatorname{Nm}(\pi)>2$. Hence $\operatorname{Nm}(\delta)<\operatorname{Nm}(\pi)$, so that $\bar{\delta} \in k\left(v_{n+1}\right)^{\times}$. Now compute

$$
\begin{aligned}
1 & =\left\{\frac{\gamma}{\alpha \beta}, 1-\frac{\gamma}{\alpha \beta}\right\} \\
& =\left\{\frac{\gamma}{\alpha \beta}, \frac{\delta \pi}{\alpha \beta}\right\} \\
& =\{\gamma, \delta\}\{\gamma, \pi\}\{\alpha \beta, \gamma\}\{\delta, \alpha \beta\}\{\pi, \alpha \beta\}\{\alpha \beta, \alpha \beta\} \\
& \equiv\{\gamma, \pi\}\{\pi, \alpha \beta\} \quad\left(\bmod K_{2}^{S_{n}} \mathbf{Q}(i)\right)
\end{aligned}
$$

where the last congruence holds since $\alpha, \beta, \gamma$ and $\delta$ all have norms bounded by $\operatorname{Nm}(\pi)$. Hence, modulo $K_{2}^{S_{n}} \mathbf{Q}(i)$ we have $\{\alpha \beta, \pi\}=\{\gamma, \pi\}$, which means that $\phi$ is multiplicative, and also well defined by taking $\beta=1$.

Furthermore we have that $\phi$ is surjective. Indeed, the group $K_{2}^{S_{n+1}} \mathbf{Q}(i)$ is generated by the elements of $K_{2}^{S_{n}} \mathbf{Q}(i)$ in addition to the symbols $\{\alpha, \pi\}$ where $\alpha$ is a Gaussian $S_{n}$-unit (remember that $\{\pi, \pi\}=\{-1, \pi\}$ ). This means that $\bar{\alpha} \in k\left(v_{n+1}\right)^{\times}$and $\phi(\bar{\alpha})=\{\alpha, \pi\}$, hence $\phi$ is surjective and $\#\left(K_{2}^{S_{n+1}} \mathbf{Q}(i) / K_{2}^{S_{n}} \mathbf{Q}(i)\right) \leq \operatorname{Nm}(\pi)-1$.

Now let $\bar{\zeta}$ denote a generator for the cyclic group $k\left(v_{n+1}\right)^{\times}$. By Proposition 1.52 we then have that

$$
(\zeta, \pi)_{v_{n+1}}^{\mathrm{reg}}=\bar{\zeta} \in k\left(v_{n+1}\right)^{\times} .
$$

But

$$
(\zeta, \pi)^{\mathrm{reg}}=h^{\mathrm{reg}}(\{\zeta, \pi\}),
$$

hence $\{\zeta, \pi\}$ has order $\operatorname{Nm}(\pi)-1$ in $K_{2}^{S_{n+1}} \mathbf{Q}(i) / K_{2}^{S_{n}} \mathbf{Q}(i)$. Therefore the group $K_{2}^{S_{n+1}} \mathbf{Q}(i) / K_{2}^{S_{n}} \mathbf{Q}(i)$ also has order at least $\operatorname{Nm}(\pi)-1$. By the above inequality, its order must be exactly $\operatorname{Nm}(\pi)-1$, hence $\phi$ is an isomorphism.

To conclude the proof, let $n \geq 1$ and assume by induction that $K_{2}^{S_{n}} \mathbf{Q}(i) \cong \bigoplus_{j=1}^{n} k\left(v_{j}\right)^{\times}$ via $h^{\text {reg }}$. We must show that $h^{\text {reg }}$ yields an isomorphism

$$
K_{2}^{S_{n+1}} \mathbf{Q}(i) \cong \bigoplus_{j=1}^{n+1} k\left(v_{j}\right)^{\times}
$$

If an element $x \in K_{2}^{S_{n+1}} \mathbf{Q}(i)$ maps to 1 under $h^{\text {reg }}$, Lemma 3.8 states that $x \in K_{2}^{S_{n}} \mathbf{Q}(i)$, and so $x=1$ by induction. Similarly, if $u=\left(u_{j}\right)_{j=1}^{n+1} \in \bigoplus_{j=1}^{n+1} k\left(v_{j}\right)^{\times}$, we can find an element in $K_{2}^{S_{n+1}} \mathbf{Q}(i)$ mapping to $u_{n+1}$ by Lemma 3.8. By the induction hypothesis there is an element in $K_{2}^{S_{n}} \mathbf{Q}(i)$ mapping to $\left(u_{j}\right)_{j=1}^{n}$, and so the product of these elements constitutes a suitable preimage of $u$.

### 3.1.1 Connection with biquadratic reciprocity

Let us show how the structure of $K_{2}(\mathbf{Q}(i))$ is related to the law of biquadratic reciprocity. By Theorem 3.7 there exist maps $\psi_{v}: k(v)^{\times} \rightarrow \mu_{4}$ for all $v \nmid 2$ such that the Hilbert symbol at the place $(1+i)$ factors through $\bigoplus_{v \nmid 2} k(v)^{\times}$via these maps:

$$
\begin{aligned}
& \mathbf{Q}(i)^{\times} \times \mathbf{Q}(i)^{\times} \stackrel{\stackrel{(-,-)^{\mathrm{reg}}}{\longrightarrow}}{\bigoplus_{v \nmid 2}} k(v)^{\times} \\
& \left(\frac{-,-}{(1+i)}\right)_{{ }_{4}} \downarrow_{\mu_{4}}
\end{aligned}
$$

In other words, for all $x, y \in \mathbf{Q}(i)^{\times}$we have

$$
\left(\frac{x, y}{(1+i)}\right)_{4}=\prod_{v \nmid 2} \psi_{v}\left((x, y)_{v}^{\mathrm{reg}}\right)
$$

Lemma 3.9. Let $v$ be a finite place of $\mathbf{Q}(i)$ not lying above 2. If $f$ is a map $f: k(v)^{\times} \rightarrow \mu_{4}$, then $f$ satisfies

$$
f(x) \equiv x^{\frac{q_{v}-1}{4} \delta}\left(\bmod \mathfrak{p}_{v}\right)
$$

where $\delta \in\{0,1,2,3\}$ and $q_{v}=\operatorname{Nm}\left(\mathfrak{p}_{v}\right)$.
Proof. Write $q:=q_{v}$. In the following, we abuse notation by identifying elements with their image in $k(v)^{\times}$.

The map $f$ is determined by where a generator of the cyclic group $k(v)^{\times}$is sent. If $\zeta$ denotes a generator for $k(v)^{\times}$, then clearly $\zeta^{\frac{q-1}{4}}$ equals $i$ or $-i$, for if $\zeta^{\frac{q-1}{4}} \in \mu_{2}$ then $\zeta$ is a square or a fourth power.

Assume first that $f(\zeta)=i$. We consider two cases.

1. If $\zeta^{\frac{q-1}{4}}=i$, then $f(x)=x^{\frac{q-1}{4}}$ since their action on the generator coincide.
2. If $\zeta^{\frac{q-1}{4}}=-i$, then $\zeta^{3 \frac{q_{v}-1}{4}}=i$, hence $f(x)=x^{3 \frac{q-1}{4}}$.

Now if $f(\zeta)=i^{\delta}$, where $\delta \in\{0,1,2,3\}$, then in case 1 above we have $f(x)=x^{\frac{q-1}{4} \delta}$ for any $x \in k(v)^{\times}$, and in case 2 we get

$$
f(x)=\left(x^{\frac{q-1}{4}}\right)^{3 \delta}=x^{\frac{q-1}{4} \delta^{\prime}}
$$

for $\delta^{\prime} \in\{0,1,2,3\}$ with $\delta^{\prime} \equiv \delta(\bmod 4)$.
Note that, by Proposition 1.52 and Observation 2.20

$$
\left(\frac{a, b}{v}\right)_{n} \equiv\left((a, b)_{v}^{\mathrm{reg}}\right)^{\frac{q_{v}-1}{n}}\left(\bmod \mathfrak{p}_{v}\right)
$$

Thus, by applying the above lemma to the maps $\psi_{v}$ above we get the relation of Hilbert symbols

$$
\left(\frac{x, y}{(1+i)}\right)_{4}=\prod_{v \nmid 2}\left((x, y)_{v}^{\mathrm{reg}}\right)^{\frac{q_{v}-1}{4} \delta_{v}}=\prod_{v \nmid 2}\left(\frac{x, y}{v}\right)_{4}^{\delta_{v}}
$$

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Figure 3.2: Lattice spanned by unit multiples of an Eisenstein prime.
where $\delta_{v} \in\{0,1,2,3\}$ for all $v \nmid 2$. In other words, the structure of the group $K_{2}(\mathbf{Q}(i))$ gives rise to a relation between Hilbert symbols of order 4. But Moore's theorem says that there is only one relation of Hilbert symbols, hence we must have $\delta_{v}=1$ for all $v \nmid 2$. Since $\left(\frac{x, y}{(1+i)}\right)_{4}$ takes values in $\{ \pm 1\}$, we have thus obtained the product formula for Hilbert symbols

$$
\prod_{v \in \operatorname{Pl}_{\mathbf{Q}(i)}^{\mathrm{nc}}}\left(\frac{x, y}{v}\right)_{4}=1
$$

which implies the biquadratic reciprocity law, as we have seen in Chapter 1

### 3.2 The group $K_{2}(\mathbf{Q}(\sqrt{-3}))$

Now we turn to the case $F=\mathbf{Q}(\sqrt{-3})=\mathbf{Q}(\omega)$, so that $\mathcal{O}_{F}=\mathbf{Z}[\omega]$, the Eisenstein integers. Here $\omega=\zeta_{3}=(-1+\sqrt{-3}) / 2$ is a primitive third root of unity. Again we consider the map

$$
h^{\mathrm{reg}}: K_{2} \mathbf{Q}(\omega) \longrightarrow \bigoplus_{v \in \mathrm{P}_{\mathbf{Q}(\omega)}^{\mathrm{nc}}} k(v)^{\times}
$$

induced by the global regular Hilbert symbol.
Definition 3.10. Let $\Lambda_{\infty}$ denote the subgroup of $K_{2} \mathbf{Q}(\omega)$ generated by the elements $\{u, v\}$ for which $u, v \in \mathbf{Z}[\omega]^{\times}=\mu_{6}=\left\{ \pm 1, \pm \omega, \pm \omega^{2}\right\}$.

Lemma 3.11. We have $\Lambda_{\infty}=1$.
Proof. Note that, by using bilinearity and the relations for Steinberg symbols, any element $\{u, v\}$ for which $u, v \in\left\{-1, \pm \omega, \pm \omega^{2}\right\}$ is a product of the symbols $\{\omega, \omega\},\{-1,-1\}$. We check that these are trivial. Recall that $\omega^{-1}=\omega^{2}$ and that $1+\omega+\omega^{2}=0$ in the following.

- We have

$$
\left\{\omega, \omega^{2}\right\}=\{\omega, \omega\}\{\omega, \omega\},
$$

but also

$$
\left\{\omega, \omega^{2}\right\}=\left\{\omega, \omega^{-1}\right\}=\{\omega, \omega\}^{-1}=\{\omega, \omega\},
$$

hence $\{\omega, \omega\}=1$.

- Note that

$$
1=\{-\omega, 1+\omega\}=\{-1,1+\omega\}\{\omega, 1+\omega\}
$$

by the Steinberg relation, hence

$$
\{1+\omega, 1+\omega\}=\{-1,1+\omega\}=\{1+\omega, \omega\} .
$$

Using that $1+\omega=-\omega^{2}$, this yields

$$
\{1+\omega, \omega\}=\{-1, \omega\}\left\{\omega^{2}, \omega\right\}=\{\omega, \omega\}\{\omega, \omega\}^{2}=1
$$

Hence we get

$$
\begin{aligned}
\{-1,-1\} & =\{\omega(1+\omega), \omega(1+\omega)\} \\
& =\{\omega, \omega\}\{\omega, 1+\omega\}\{1+\omega, \omega\}\{1+\omega, 1+\omega\} \\
& =1
\end{aligned}
$$

which concludes the proof.
As in Section 3.1 we now order the noncomplex places of $\mathbf{Q}(\omega)$ by increasing norm,

$$
\mathrm{Pl}_{\mathbf{Q}(\omega)}^{\mathrm{nc}}=\mathrm{Pl}_{0}=\left\{v_{1}, v_{2}, \ldots \mid \mathrm{Nm}\left(\pi_{v_{n}}\right) \leq \mathrm{Nm}\left(\pi_{v_{n+1}}\right) \quad \forall n\right\}
$$

and we let $K_{2}^{S_{n}} \mathbf{Q}(\omega)$ denote the subgroup of $K_{2} \mathbf{Q}(\omega)$ generated by elements $\{u, v\}$ for $u$, $v \in \mathbf{Z}[\omega]_{S_{n}}^{\times}$. Here $S_{n}$ denotes as usual the set $S_{n}=\left\{v_{1}, \ldots, v_{n}\right\}$ of the first $n$ places of $\mathbf{Q}(\omega)$ of smallest norm.

By Proposition 1.26, $\mathbf{Q}(\omega) / \mathbf{Q}$ is ramified only at 3 , hence the smallest value for $\operatorname{Nm} \pi_{v}$ is 3 . Thus $v_{1}$ is the place corresponding to the prime ideal $(2+\omega)$, and by Proposition 1.27 we have $k\left(v_{1}\right) \cong \mathbf{F}_{3}$.
Lemma 3.12. We have $K_{2}^{S_{1}} \mathbf{Q}(\omega) \cong k\left(v_{1}\right)^{\times} \cong \mathbf{F}_{3}^{\times}$.
Proof. We claim that the element $x:=\{-1,2+\omega\}=\{2+\omega, 2+\omega\}$ is the only nontrivial element in $K_{2}^{S_{1}} \mathbf{Q}(\omega)$. Note that

$$
x^{2}=\{-1,2+\omega\}^{2}=\{1,2+\omega\}=1
$$

so this element has order 2 . Computing the regular symbol at $v_{1}$ yields

$$
(-1,2+\omega)_{v_{1}}^{\mathrm{reg}} \equiv-1\left(\bmod \mathfrak{p}_{v_{1}}\right)
$$

by Proposition 1.52 Since the regular symbol factors through $K_{2} \mathbf{Q}(\omega)$, it follows that $x$ is nontrivial.

To conclude we must show that the rest of the symbols in $K_{2}^{S_{1}} \mathbf{Q}(\omega)$ are trivial. Using Lemma 3.11 along with bilinearity and antisymmetry, we see that it is enough to show that $\{2+\omega, \omega\}=1$. To this end, note that since $2+\omega=1-\omega^{2}$, the Steinberg relation yields

$$
\left\{\omega^{2}, 2+\omega\right\}=1
$$

Using that $\omega=\omega^{-2}$, we thus get

$$
\begin{aligned}
\{2+\omega, \omega\} & =\left\{2+\omega, \omega^{-2}\right\} \\
& =\left\{2+\omega, \omega^{2}\right\}^{-1} \\
& =\left\{\omega^{2}, 2+\omega\right\}=1
\end{aligned}
$$

Theorem 3.13. The global regular Hilbert symbol induces an isomorphism

$$
K_{2}(\mathbf{Q}(\omega)) \cong \bigoplus_{v \in \mathrm{P}_{\mathbf{P}(\omega)}^{\mathrm{nc}}(v)^{\times}} k(v)^{\times} \cong \mathbf{F}_{3}^{\times} \oplus \bigoplus_{p \equiv 1(\bmod 3)}\left(\mathbf{F}_{p}^{\times}\right)^{2} \oplus \bigoplus_{p \equiv 2}^{(\bmod 3)} \mathbf{F}_{p^{2}}^{\times}
$$

Proof. We proceed in the same manner as in Theorem 3.7, i.e., we use induction to prove that $K_{2}^{S_{n}} \mathbf{Q}(\omega) \cong \bigoplus_{j=1}^{n} k\left(v_{j}\right)^{\times}$. Lemma 3.12 furnishes the base case of the induction. For the induction step, we will again show that the map

$$
\phi: k\left(v_{n+1}\right)^{\times} \longrightarrow K_{2}^{S_{n+1}} \mathbf{Q}(\omega) / K_{2}^{S_{n}} \mathbf{Q}(\omega)
$$

defined by

$$
\phi(\bar{\alpha}):=\{\alpha, \pi\}\left(\bmod K_{2}^{S_{n}} \mathbf{Q}(\omega)\right)
$$

is a well defined isomorphism, where $\pi:=\pi_{v_{n+1}}$. Here $\alpha$ is to vary between all Eisenstein integers of norm bounded by $\operatorname{Nm}(\pi) / 3$. This assumption can be made thanks to the Euclidean algorithm on $\mathbf{Z}[\omega]$ and Table 1.1, just as in Lemma 3.1.

To show $\phi$ is a well defined homomorphism, we proceed as before: If $\alpha \beta=\gamma+\pi \delta$ with $|\alpha|$, $|\beta|,|\gamma| \leq|\pi| / \sqrt{3}$, then

$$
|\pi \delta| \leq|\alpha||\beta|+|\gamma| \leq \frac{|\pi|^{2}}{3}+\frac{|\pi|}{\sqrt{3}}
$$

hence

$$
|\delta| \leq \frac{|\pi|}{3}+\frac{1}{\sqrt{3}}<|\pi|
$$

An identical computation as the one done in the proof of Theorem 3.7 yields that $\phi$ is well defined and multiplicative. The proof of the bijectivity also follows from the same argument as we have seen.

### 3.3 Tame and wild kernels, exotic symbols

In Section 2.2 we defined regular-or tame-kernels as well as the wild kernel (see Definition 2.22). The wild kernel $W K_{2}(F)$ was shown by Bass and Tate to always be finitely generated [Tat71]. Here $F$ can be any field. A deep theorem of Garland established that in the case when $F$ is a number field, $W K_{2}(F)$ is in fact finite Gar71.

The tame and wild kernels encode deep arithmetic information about the number field $F$. For example, as conjectured by Birch-Tate and proved by Mazur-Wiles-Kolyvagin, if $F / \mathbf{Q}$ is a totally real abelian extension we have the expression Gra03, II Remark 7.8.2]

$$
\# R_{2}^{\mathrm{ord}}(F)=\frac{w_{2}}{2^{[K: \mathbf{Q}]}}\left|\zeta_{F}(-1)\right|
$$

Here $\zeta_{F}$ is the Dedekind zeta-function of $F$, and $w_{2}$ is the largest integer $n$ such that $\operatorname{Gal}\left(F\left(\mu_{n}\right) / F\right)$ is killed by 2 [Wei13, p.516].

The computation in the previous sections is of course intimately linked with the study of the groups $R_{2}^{\text {ord }}(F)$ and $W K_{2}(F)$. Indeed, the above computations show that the regular kernel $R_{2}^{\text {ord }}(F)$ is trivial in the case when $F=\mathbf{Q}(\sqrt{-1})$ or $F=\mathbf{Q}(\sqrt{-3})$. In other words, $K_{2}(\mathbf{Z}[i])=1=K_{2}(\mathbf{Z}[\omega])$ by Remark 2.23 .
Example 3.14. In the case of a quadratic number field $\mathbf{Q}(\sqrt{d})$, it is quite easy to give examples for which there is a nontrivial regular kernel. For example, let $F$ be the field $\mathbf{Q}(\sqrt{d})$ where $d<0$ is a squarefree integer such that $d \equiv 1(\bmod 8)$. By Proposition 1.10, 2 splits in the extension $F / \mathbf{Q}$. Now let $v$ be a place of $F$ above 2 . The fact that $d \equiv 1(\bmod 8)$ means that $d$ has a square root in $\mathbf{Q}_{2}$ by Hensel's lemma, hence $F_{v}=\mathbf{Q}_{2}(\sqrt{d})=\mathbf{Q}_{2}$. Thus

$$
\left(\frac{-1,-1}{v}\right)=\left(\frac{-1,-1}{2}\right)=-1
$$

by Example 1.54 This means that the element $\{-1,-1\}$ is nontrivial in $K_{2}(F)$. But the formula for tame symbols (Proposition 1.52) along with the fact that

$$
\mu\left(F_{v}\right)^{\mathrm{reg}}=\mu\left(\mathbf{Q}_{2}\right)^{\mathrm{reg}}=1
$$

shows that $\{-1,-1\} \in R_{2}^{\text {ord }}(F)$. On the other hand we clearly have $\{-1,-1\} \notin W K_{2}(F)$ since $h(\{-1,-1\})$ is nontrivial on the component $\mu\left(F_{v}\right)=\mu\left(\mathbf{Q}_{2}\right)=\mu_{2}$.

Among the imaginary quadratic fields, the first instance of a nontrivial regular kernel occurs for $\mathbf{Q}(\sqrt{-7})$.

Tate shows in [BT73, p.435] that the field $F:=\mathbf{Q}(\sqrt{-35})$ has $R_{2}^{\text {ord }}(F)$ nontrivial, yet there are no local Hilbert symbols showing this. This is an example of an imaginary quadratic field on which there is a symbol that cannot be expressed in terms of Hilbert symbols. Equivalently, such a symbol is nontrivial on the wild kernel. These symbols are called exotic. Thus the wild kernel is in some sense a measure of how many exotic symbols one can have on a number field.

Example 3.15. In Keu97, Keune shows that if $d<0$ is a squarefree integer with $d \equiv$ $2(\bmod 16)$, the element $\{-1,-1\} \in K_{2}(\mathbf{Q}(\sqrt{d}))$ lies in the wild kernel.

Example 3.16. Hutchinson shows in Hut04 that the imaginary quadratic fields $\mathbf{Q}(\sqrt{d})$, where $d \equiv 2(\bmod 16)$ is negative and squarefree, possess exotic symbols. Specifically, he shows that there is a Steinberg symbol $\lambda$ defined on $\mathbf{Q}(\sqrt{d})$ with values in $\mu_{8}$ such that $\lambda(\{-1,-1\})=-1$.

## CHAPTER 4

## Motivic Homotopy Theory

Motivic-, or $\mathbf{A}^{1}$-homotopy theory is a homotopy theory for smooth schemes, developed by Morel and Voevodsky MV99]. Letting $\mathrm{Sm}_{F}$ denote the category of smooth schemes over a field $F$, the basic goal is to create suitably nice unstable and stable homotopy categories from $\mathrm{Sm}_{F}$, in which the affine line $\mathbf{A}^{1}:=\mathbf{A}_{F}^{1}$ plays the role of the unit interval. Now the naive approach toward such a construction would be simply to directly lift the definition of a homotopy in topology by saying two morphisms

$$
f, g: X \longrightarrow Y
$$

in $\mathrm{Sm}_{F}$ are $\mathbf{A}^{1}$-homotopic if there is a morphism

$$
H: X \times_{F} \mathbf{A}^{1} \longrightarrow Y
$$

such that $\left.H\right|_{X \times\{0\}}=f$ and $\left.H\right|_{X \times\{1\}}=g$. This approach is however unsatisfactory. Indeed, the resulting relation of " $\mathbf{A}^{1}$-homotopy" is not an equivalence relation, as it fails to be transitive. And even if we consider the equivalence relation generated by " $\mathbf{A}^{1}$-homotopy" there is no reasonable way in general to put a group structure on the resulting set of $\mathbf{A}^{1}$-homotopy classes of morphisms that we would like to call " $\mathbf{A}^{1}$-homotopy groups" Mor04a, 2.2].

Instead of the above approach, one could try to invoke the machinery of model categories in order to create the desired homotopy category. But then the lack of certain categorical properties of $\mathrm{Sm}_{F}$ becomes apparent: For example, one of the first properties demanded of the category in order to put a model structure on it is bicompleteness, i.e., all small limits and colimits exist. This is not fulfilled by $\mathrm{Sm}_{F}$ since this category does not have all quotients. To fix this problem, one first enlarges the category $\mathrm{Sm}_{F}$ by considering instead the category

$$
\operatorname{Pre}\left(\operatorname{Sm}_{F}\right):=\left[\mathrm{Sm}_{F}^{\mathrm{op}}, \operatorname{Set}\right]
$$

of presheaves on $\operatorname{Sm}_{F}$. Thus $\operatorname{Pre}\left(\operatorname{Sm}_{F}\right)$ is the category whose objects are functors $\operatorname{Sm}_{F}^{\text {op }} \rightarrow$ Set, and whose morphisms are natural transformations. Via the Yoneda embedding

$$
X \mapsto \operatorname{Hom}_{\operatorname{Sm}_{F}}(-, X): \operatorname{Sm}_{F} \rightarrow \operatorname{Pre}\left(\operatorname{Sm}_{F}\right)
$$

we identify $\mathrm{Sm}_{F}$ with the full subcategory of $\operatorname{Pre}\left(\mathrm{Sm}_{F}\right)$ consisting of representable functors. Moreover, the category $\operatorname{Pre}\left(\mathrm{Sm}_{F}\right)$ is bicomplete since Set is bicomplete. But in order to put a model structure on a functor category it is convenient to have a model structure on the target category. So we replace Set by the model category $\mathscr{S}:=\left[\Delta^{\mathrm{op}}\right.$, Set $]$ of simplicial sets (see Hov99, $\mathrm{DL}^{+} 07$ ). Thus we further embed $\operatorname{Pre}\left(\mathrm{Sm}_{F}\right)$ into the category

$$
\Delta^{\mathrm{op}} \operatorname{Pre}\left(\operatorname{Sm}_{F}\right):=\left[\Delta^{\mathrm{op}}, \operatorname{Pre}\left(\operatorname{Sm}_{F}\right)\right]=\left[\mathrm{Sm}_{F}^{\mathrm{op}}, \mathscr{S}\right]=\left[\Delta^{\mathrm{op}} \times \operatorname{Sm}_{F}^{\mathrm{op}}, \operatorname{Set}\right]
$$

of simplicial presheaves. The embedding $\operatorname{Pre}\left(\operatorname{Sm}_{F}\right) \hookrightarrow \Delta^{\mathrm{op}} \operatorname{Pre}\left(\mathrm{Sm}_{F}\right)$ is given by considering a presheaf $X$ as a constant simplicial presheaf $[n] \mapsto X$. Note that the category of simplicial sets also embeds in $\Delta^{\mathrm{op}} \operatorname{Pre}\left(\operatorname{Sm}_{F}\right)$ by considering a simplicial set as constant on $\operatorname{Sm}_{F}$, i.e., for $K \in \mathscr{S}$, define

$$
K([n], X):=K_{n} \in \operatorname{Set}
$$

for all $X \in \operatorname{Sm}_{F}$.

The category $\Delta^{\mathrm{op}} \operatorname{Pre}\left(\mathrm{Sm}_{F}\right)$ has all the desired properties one can ask for, and is therefore called the category of motivic spaces, denoted by $\mathcal{M S}(F)$. The category of motivic spaces is then provided with a model structure which generalizes the naive point of view, yielding the notion of $\mathbf{A}^{1}$-weak equivalences.

Notation 4.1. We will identify a scheme $X \in \operatorname{Sm}_{F}$ with its corresponding representable presheaf $\operatorname{Hom}_{\operatorname{Sm}_{F}}(-, X) \in \operatorname{Pre}\left(\operatorname{Sm}_{F}\right)$. For a smooth subscheme $Y \hookrightarrow X$ we write $X / Y$ for the pointed sheaf associated to the presheaf $U \mapsto \operatorname{Hom}(U, Y) / \operatorname{Hom}(U, X)$.

We will denote by $\mathcal{M S} .(F)$ the category of pointed motivic spaces. The category $\mathcal{M S}(F)$ embeds in $\mathcal{M S} \bullet(F)$ via $X \mapsto X_{+}:=X \amalg \operatorname{Spec}(F)$, where $X_{+}$is pointed by $\operatorname{Spec}(F)$.

Below we will briefly discuss the model structure on $\mathcal{M S} \bullet(F)$ and the construction of the motivic unstable and stable homotopy category $\mathcal{H}_{\bullet}(F)$ and $\mathcal{S H}(F)$. However, background material on, e.g., model categories is left out, as this is not the main objective of this thesis. The curious reader may consult Hov99.

On the homological side we also mention below the theory of motives and motivic cohomology; in particular we will discuss Voevodsky's construction of the derived category of geometric motives.

### 4.1 The motivic unstable and stable homotopy category

The model structure on $\mathcal{M S}(F)$ is constructed to take into account the Nisnevich topology on $\mathrm{Sm}_{F}$. This topology is finer than the Zariski topology, but coarser than the étale topology, and is precisely what is needed to prove, e.g., descent theorems for algebraic $K$-theory and the homotopy purity theorem.

Definition 4.2. A map $f: U \rightarrow X$ of schemes is completely decomposed at $x \in X$ if there is $u \in U$ such that $f(u)=x$ and the induced map on residue fields $k(x) \rightarrow k(u)$ is an isomorphism.

Definition 4.3. Let $X \in \operatorname{Sm}_{F}$. A covering $\left\{f_{i}: U_{i} \rightarrow X\right\}_{i \in I}$ is called a Nisnevich covering if $I$ is finite, all the $f_{i}$ 's are étale and for all $x \in X$ there is an $i \in I$ such that $f_{i}$ is completely decomposed at $x$.

The Nisnevich coverings provide a basis for a Grothendieck topology on $\mathrm{Sm}_{F}$, called the Nisnevich topology. The resulting Nisnevich site is denoted $\left(\mathrm{Sm}_{F}\right)_{\text {Nis }}$.

Definition 4.4. A cartesian square

is called an elementary distinguished square in the Nisnevich topology if $i$ is an open embedding, $p$ is étale and $p^{-1}(X \backslash U) \xrightarrow{p} X \backslash U$ is an isomorphism of reduced schemes.

Proposition 4.5 ([Sev06, Proposition 6.1.11]). Given an elementary distinguished square ( $i, p$ ) as in Definition 4.4 above, the set $\{U \xrightarrow{i} X, V \xrightarrow{p} X\}$ forms a Nisnevich covering of $X$.

Moreover, a presheaf $\mathscr{F}$ on the Nisnevich site $\left(\operatorname{Sm}_{F}\right)_{\text {Nis }}$ is a sheaf if and only if $\mathscr{F}(\varnothing)=*$ and for all elementary distinguished squares $(i, p)$, the diagram

is a pullback square in Set.
We let $\operatorname{Shv}_{\mathrm{Nis}}\left(\mathrm{Sm}_{F}\right)$ denote the category of sheaves on the Nisnevich site $\left(\mathrm{Sm}_{F}\right)_{\mathrm{Nis}}$.

Definition 4.6. The projective objectwise model structure on $\mathcal{M} \mathcal{S}_{\bullet}(F)$ is defined as follows. Weak equivalences and cofibrations are defined objectwise - i.e., $\mathscr{X} \rightarrow \mathscr{Y}$ is a weak equivalence in $\mathcal{M S} .(F)$ if and only if $\mathscr{X}(X) \rightarrow \mathscr{Y}(X)$ is a weak equivalence in $\mathscr{S}$ for all $X \in \operatorname{Sm}_{F}$, and similarly for cofibrations. Fibrations are the maps having the right lifting property with respect to acyclic cofibrations.

Now let $Q$ denote a Nisnevich elementary distinguished square as in Definition 4.4 and let $Q^{\text {hp }}$ denote the homotopy pushout of $U \leftarrow U \times_{X} V \rightarrow V$. Form the class of morphisms

$$
S_{\mathrm{Nis}}^{\mathrm{hp}}:=\left\{Q^{\mathrm{hp}} \rightarrow X\right\}_{Q} \cup\left\{\varnothing \rightarrow h_{\varnothing}\right\}
$$

where $Q$ ranges over all Nisnevich elementary distinguished squares. Here $\varnothing$ means the empty simplicial set, while $h_{\varnothing}$ is the presheaf represented by the empty scheme. The set $\left\{\varnothing \rightarrow h_{\varnothing}\right\}$ is included only for technical reasons.
Definition 4.7. The local projective model structure on $\mathcal{M S} \bullet(F)$ is the $S_{\mathrm{Nis}}^{\mathrm{hp}}$-localization of the projective model structure on $\mathcal{M S} \bullet(F)$.

The local projective model structure makes $\mathcal{M S} \mathcal{S}_{\bullet}(F)$ into a proper model category [Sev06, Theorem 6.2.7]. This model structure takes into account certain relations that is needed to prove for example the purity theorem for Thom spaces. However, one can show that there is a full embedding of $\mathrm{Sm}_{F}$ into the homotopy category of $\mathcal{M S} \bullet(F)$ with the local projetive model structure [Sev06, Proposition 6.2.4]. In other words, no schemes are identified in this homotopy category. Therefore we need to localize once more if we want the affine line to be contractible.

Definition 4.8. The motivic projective model structure on $\mathcal{M S}_{\bullet}(F)$ is the left Bousfield localization of the local projective model structure with respect to $S:=\left\{X \times{ }_{F} \mathbf{A}^{1} \rightarrow X\right\}_{X \in \operatorname{Sm}_{F}}$. The resulting homotopy category is the unstable motivic homotopy category $\mathcal{H}_{\bullet}(F)$.

There are two kinds of circles in $\mathcal{H}_{\bullet}(F)$ : Familiar from topology we have the simplicial circle $S^{1}:=\Delta^{1} / \partial \Delta^{1}$, which is pointed by $\Delta^{0}$. On the geometric side we have the Tate circle $S^{\alpha}:=\mathbf{A}^{1} \backslash\{0\}$, pointed by 1 . We will often identify $S^{\alpha}$ with the group scheme $\mathbf{G}_{m}$. There is a smash product on $\mathcal{H}_{\bullet}(F)$, defined pointwise, which allows us to talk about the mixed spheres

$$
S^{m+n \alpha}:=\left(S^{1}\right)^{\wedge m} \wedge\left(S^{\alpha}\right)^{\wedge n}
$$

Example 4.9. By taking complex realization we obtain $S^{m+n \alpha}(\mathbf{C})=S^{m+n}$, where $S^{k}$ is the usual topological sphere. If we consider the real realization we get $S^{m+n \alpha}(\mathbf{R})=S^{m}$.

In Section 4.2 we will see that there is a canonical isomorphism $S^{1+\alpha} \cong \mathbf{P}^{1}$ in $\mathcal{H}_{\bullet}(F)$. We let $\Sigma_{\mathbf{P}^{1}}$ denote the $\mathbf{P}^{1}$-suspension functor $X \mapsto \mathbf{P}^{1} \wedge X$. The process of stabilizing with respect to $\mathbf{P}^{1}$-suspension results in the stable motivic homotopy category $\mathcal{S H}(F)$, which we elaborate on next.

Definition 4.10. A $\mathbf{P}^{1}$-spectrum $X=\left(X_{n}\right)_{n \geq 0}$ is a sequence of pointed motivic spaces $X_{n} \in$ $\mathcal{M S} .(F)$ together with structure maps $\sigma: \Sigma_{\mathbf{P}^{1}} X_{n} \rightarrow X_{n+1}$ for all $n$. A map $f: X \rightarrow Y$ of spectra is a sequence of maps $f_{n}: X_{n} \rightarrow Y_{n}$ compatible with the structure maps, i.e., the diagram

commutes for each $n \geq 0$. We let $\operatorname{Spt}(F)$ denote the category of $\mathbf{P}^{1}$-spectra. There is an embedding $\Sigma_{\mathbf{P}^{1}}^{\infty}: \mathcal{M S} \bullet(F) \rightarrow \operatorname{Spt}(F)$, where $\Sigma_{\mathbf{P}^{1}}^{\infty} X$ is the spectrum with constituent spaces $\left(\Sigma_{\mathbf{P}^{1}}^{\infty} X\right)_{n}:=\Sigma_{\mathbf{P}^{1}}^{n} X$ and identity structure maps.
Definition 4.11. The sphere spectrum is

$$
1:=\Sigma_{\mathbf{P}^{1}}^{\infty} \operatorname{Spec}(F)_{+} \in \operatorname{Spt}(F)
$$

This spectrum is initial in the category of ring spectra.

We do not intend to go into details of the construction of the stable model structure on $\operatorname{Spt}(F)$; this can be found in, e.g., Jar00. The rough idea is however as follows. A map $f: X \rightarrow Y$ of spectra is called a levelwise weak equivalence (resp. (co)fibration) if the component maps $f_{n}: X_{n} \rightarrow Y_{n}$ are weak equivalences (resp. (co)fibrations) for all $n$. Thus the category $\operatorname{Spt}(F)$ inherits a model structure from $\mathcal{M} \mathcal{S}_{\bullet}(F)$. The projective model structure on $\operatorname{Spt}(F)$ is defined by these levelwise fibrations and cofibrations. From this model structure one defines a certain localizing set $S$ of morphisms in $\operatorname{Spt}(F)$, and weak equivalences are defined as the $S$-local weak equivalences. The resulting model structure is the stable model structure. The homotopy category of $\operatorname{Spt}(F)$ with the stable model structure is the motivic stable homotopy category $\mathcal{S H}(F)$.

Theorem 4.12 (Voe98). We list a few properties of the category $\mathcal{S H}(F)$.

- $\mathcal{S H}(F)$ is a triangulated category, with simplicial suspension $\Sigma_{s}:=-\wedge S^{1}$ as shift functor.
- There are realization functors


Here $\mathcal{S H}$ denotes the topological stable homotopy category, and $\mathcal{S H}{ }^{\mathbf{Z} / 2}$ is the $\mathbf{Z} / 2$-equivariant stable homotopy category. The realization of the motivic sphere spectrum is the topological sphere spectrum.

- $\mathcal{S H}(F)$ satisfies Grothendieck's six functor formalism [CD12].

In the category $\mathcal{S H}(F)$ we can define stable motivic homotopy groups:
Definition 4.13. Let $m, n \in \mathbf{Z}$ and $X \in \operatorname{Spt}(F)$. We write $\underline{\pi}_{m+n \alpha}(X)$ for the sheaf associated to the presheaf of stable homotopy groups

$$
U \longmapsto \operatorname{colim}_{k} \operatorname{Hom}_{\mathcal{S H}(F)}\left(S^{(m+k)+(n+k) \alpha} \wedge \Sigma_{\mathbf{P}^{1}}^{\infty} U_{+}, X\right)
$$

where $U \in \operatorname{Sm}_{F}$. The bigraded stable homotopy groups of $X$ is $\pi_{m+n \alpha}(X):=\underline{\pi}_{m+n \alpha}(X)(\operatorname{Spec} F)$.
There is a description of the stable weak equivalences in the stable model structure on $\operatorname{Spt}(F)$ that is more in line with classical topology:

Theorem 4.14 (Jar00, Lemma 3.7]). A map $f: X \rightarrow Y$ in $\operatorname{Spt}(F)$ is a stable weak equivalence if and only if $f$ induces an isomorphism of presheaves $\underline{\pi}_{n+m \alpha} X \rightarrow \underline{\pi}_{n+m \alpha} Y$ for all $m, n \in \mathbf{Z}$.

Notation 4.15. The notation $\pi_{p, q}$ and $S^{p, q}$ for the motivic homotopy groups and mixed spheres is also often used in the literature. To transition between the two, use $S^{p, q}=S^{(p-q)+q \alpha}$.

We will use a star $\star$ to denote bigrading, i.e., all indices $(m, n) \in \mathbf{Z}^{2}$, and an asterisk $*$ will denote monograding (e.g., compare $\pi_{\star}$ and $\pi_{* \alpha}$ ).

The most fundamental objects in stable motivic homotopy theory are perhaps the homotopy groups $\pi_{m+n \alpha} \mathbf{1}$ of the sphere spectrum. Very few of these groups are known; indeed, since the realization of $\mathbf{1}$ is the topological sphere spectrum, the groups $\pi_{m+n \alpha} \mathbf{1}$ are at least as hard to compute as the topological stable stems.

Morel's $\mathbf{A}^{1}$-connectivity theorem Mor05 shows that $\pi_{m+n \alpha} \mathbf{1}=0$ whenever $m<0$, and in Mor04a the homotopy groups are computed for $m=0$ (see Chapter 5. Theorem 5.1). Chapter 5 is essentially devoted to the study of these groups for $m=0$, which Morel identifies as the Milnor-Witt K-theory of $F$ (see Definition 5.6 in Chapter 5 for the definition). In Chapter 5 we will see a few computations of the Milnor-Witt $K$-theory of different fields-in other words, we compute $\pi_{n \alpha} \mathbf{1}$ for different base fields $F$.

The groups $\pi_{1+n \alpha} \mathbf{1}$ have recently been computed in [RSØ16], but the situation for higher values of $m$ is in general unknown to this date.

### 4.2 Motivic spheres

In this section we list a few useful isomorphisms in $\mathcal{H}_{\bullet}(F)$ Voe98, Lemma 4.1].
Proposition 4.16. The Tate object $T:=\mathbf{A}^{1} /\left(\mathbf{A}^{1} \backslash\{0\}\right)$ is weakly equivalent to the sphere $S^{1+\alpha}$.

Proof. The homotopy colimit of the diagram $* \leftarrow \mathbf{A}^{1} \backslash\{0\} \rightarrow *$ equals the simplicial suspension of $\mathbf{A}^{1} \backslash\{0\}$ [Str11, p.153]; identifying $\mathbf{A}^{1} \backslash\{0\}$ with $\mathbf{G}_{m}$, this yields

$$
\operatorname{hocolim}\left(* \longleftarrow \mathbf{A}^{1} \backslash\{0\} \longrightarrow *\right)=\Sigma_{s}\left(\mathbf{A}^{1} \backslash\{0\}\right)=S^{1} \wedge \mathbf{G}_{m}=S^{1+\alpha}
$$

On the other hand, the following diagram is a pushout by definition:


Since the upper arrow in the diagram above is a cofibration and all $X \in \operatorname{Sm}_{F}$ are $\mathbf{A}^{1}$-local fibrant [Sev06, Lemma 6.2.16], the colimit of the above diagram coincides with the homotopy colimit Lur09, Proposition A.2.4.4]. Since $\mathbf{A}^{1} \simeq *$ this yields $\mathbf{A}^{1} /\left(\mathbf{A}^{1} \backslash\{0\}\right) \simeq S^{1+\alpha}$.

Proposition 4.17. In the pointed motivic homotopy category $\mathcal{H}_{\bullet}(F)$ we have a canonical isomorphism

$$
\mathbf{A}^{1} /\left(\mathbf{A}^{1} \backslash\{0\}\right) \cong \mathbf{P}^{1}
$$

and more generally,

$$
\mathbf{A}^{n} /\left(\mathbf{A}^{n} \backslash\{0\}\right) \cong \mathbf{P}^{n} / \mathbf{P}^{n-1}
$$

for any $n \geq 1$.

Proof. We prove the assertion for $n=1$ by computing the homotopy pushout of

$$
* \longleftarrow \mathbf{A}^{1} /\left(\mathbf{A}^{1} \backslash\{0\}\right) \longrightarrow *
$$

in a third way. Consider the diagram

where the cofibrations are given by $x \mapsto x$ and $x \mapsto x^{-1}$. By gluing we obtain $\mathbf{P}^{1}$ as the pushout of this diagram, hence by the results above we obtain the following isomorphisms of homotopy pushouts:

$$
\mathbf{P}^{1} \cong \mathbf{A}^{1} /\left(\mathbf{A}^{1} \backslash\{0\}\right) \cong S^{1+\alpha} .
$$

Proposition 4.18 ( $\overline{\mathrm{DI} 05}$, Example 2.11]). In $\mathcal{H}_{\bullet}(F)$, we have a canonical isomorphism

$$
\mathbf{A}^{n} \backslash\{0\} \cong S^{(n-1)+n \alpha}
$$

for any $n \geq 1$.

Proof (following DI05]). We use induction on $n$. For $n=1$ this follows by definition: $\mathbf{A}^{1} \backslash\{0\}=$ $\mathbf{G}_{m}=S^{\alpha}$. For $n=2$, we use the open covering $\mathbf{A}^{2} \backslash\{0\}=U \cup V$, where

$$
\begin{aligned}
U & :=\left(\mathbf{A}^{1} \backslash\{0\}\right) \times \mathbf{A}^{1}, \\
V & :=\mathbf{A}^{1} \times\left(\mathbf{A}^{1} \backslash\{0\}\right) .
\end{aligned}
$$

Then $U \cap V=\left(\mathbf{A}^{1} \backslash\{0\}\right) \times\left(\mathbf{A}^{1} \backslash\{0\}\right)$, and $\mathbf{A}^{2} \backslash\{0\}$ is the homotopy pushout of $U \leftarrow U \cap V \rightarrow V$ in the category of presheaves. Using that the projections from $\left(\mathbf{A}^{1} \backslash\{0\}\right) \times \mathbf{A}^{1}$ and $\mathbf{A}^{1} \times\left(\mathbf{A}^{1} \backslash\{0\}\right)$ onto $\mathbf{A}^{1} \backslash\{0\}$ are weak equivalences, we get the diagram

where $P$ is the pushout of the outer diagram

$$
\mathbf{A}^{1} \backslash\{0\} \longleftarrow\left(\mathbf{A}^{1} \backslash\{0\}\right) \times\left(\mathbf{A}^{1} \backslash\{0\}\right) \longrightarrow \mathbf{A}^{1} \backslash\{0\}
$$

Since the arrows in the upper left square are cofibrations and $\mathcal{M S}_{\bullet}(F)$ is a left proper model category we have $\mathbf{A}^{2} \backslash\{0\} \simeq P$, so we proceed to compute the pushout $P$ instead. Consider the diagram

where the middle and lower horizontal maps are the projections onto the factors. The homotopy colimits of the rows in this diagram are respectively $*, *$ and $P$, going from the top row to the bottom row. We thus obtain the diagram $* \leftarrow * \rightarrow P$, whose homotopy pushout is $P$. Here we have used that

$$
\operatorname{hocolim}\left(X \stackrel{\mathrm{pr}_{1}}{\longleftarrow} X \vee Y \xrightarrow{\mathrm{pr}_{2}} Y\right)=* .
$$

On the other hand, if we identify $\mathbf{A}^{1} \backslash\{0\}=\mathbf{G}_{m}$ and compute the homotopy pushouts of the columns, we obtain the diagram

$$
* \longleftarrow \mathbf{G}_{m} \wedge \mathbf{G}_{m} \longrightarrow *
$$

whose homotopy pushout is $\Sigma_{s}\left(\mathbf{G}_{m} \wedge \mathbf{G}_{m}\right)=S^{1+2 \alpha}$. This yields $P=S^{1+2 \alpha}$, as we must obtain the same homotopy type no matter which direction we start computing the homotopy pushouts.

For $n>2$, one proceeds in the same manner as above: Cover $\mathbf{A}^{n}$ by open sets

$$
\begin{aligned}
U & =\left(\mathbf{A}^{n-1} \backslash\{0\}\right) \times \mathbf{A}^{1} \\
V & =\mathbf{A}^{1} \times\left(\mathbf{A}^{n-1} \backslash\{0\}\right)
\end{aligned}
$$

and perform a similar computation as above to obtain

$$
\mathbf{A}^{n} \backslash\{0\} \simeq \Sigma_{s}\left(\left(\mathbf{A}^{n-1} \backslash\{0\}\right) \wedge\left(\mathbf{A}^{1} \backslash\{0\}\right)\right),
$$

from which the result follows by induction.

### 4.3 The algebraic $K$-theory spectrum

We will concern ourselves with the motivic $\mathbf{P}^{1}$-spectrum KGL representing algebraic $K$-theory. Voevodsky defines this spectrum in Voe98; we start out by briefly recalling this construction.

For any $N \geq n$, consider the Grassmannian $\operatorname{Gr}_{n}\left(\mathbf{A}^{N}\right)$. There are canonical inclusions $\operatorname{Gr}_{n}\left(\mathbf{A}^{N}\right) \hookrightarrow \operatorname{Gr}_{n}\left(\mathbf{A}^{N+1}\right)$; let $\mathrm{BGL}_{n}$ denote the union

$$
\mathrm{BGL}_{n}:=\underset{m}{\operatorname{colim}} \mathrm{Gr}_{n}\left(\mathbf{A}^{n+m}\right)
$$

Furthermore, there are canonical inclusions

$$
\operatorname{Gr}_{n}\left(\mathbf{A}^{N}\right) \longleftrightarrow \operatorname{Gr}_{n+1}\left(\mathbf{A}^{N+1}\right)
$$

given by sending a linear subspace $L$ of $\mathbf{A}^{N}$ to $L \oplus\{0\} \subseteq \mathbf{A}^{N+1}$. Hence we can define

$$
\mathrm{BGL}:=\underset{n}{\operatorname{colim}} \mathrm{BGL}_{n}=\underset{n}{\operatorname{colim}} \operatorname{colim}_{m} \mathrm{Gr}_{n}\left(\mathbf{A}^{n+m}\right) .
$$

The constituent spaces of the spectrum KGL will be the spaces KGL, where KGL is the fibrant replacement of

$$
\mathbf{Z} \times \mathrm{BGL}:=\coprod_{\mathbf{Z}} \mathrm{BGL} .
$$

The reason for taking a fibrant replacement lies in the definition of the structure maps $\mathbf{P}^{1} \wedge$ KGL $\rightarrow$ KGL of KGL, which again relies on the isomorphism

$$
\operatorname{Hom}_{\mathcal{H} \cdot(F)}\left(\mathbf{P}^{1} \wedge(\mathbf{Z} \times \mathrm{BGL}), \mathbf{Z} \times \mathrm{BGL}\right) \cong \operatorname{Hom}_{\mathcal{H} \cdot(F)}(\mathbf{Z} \times \mathrm{BGL}, \mathbf{Z} \times \mathrm{BGL})
$$

of Voe98, p.600]. Let

$$
\bar{e}: \mathbf{P}^{1} \wedge(\mathbf{Z} \times \mathrm{BGL}) \longrightarrow \mathbf{Z} \times \mathrm{BGL}
$$

be the map corresponding to the identity morphism of $\mathbf{Z} \times$ BGL under the above isomorphism. Since KGL is fibrant, $\bar{e}$ lifts to a map

$$
e: \mathbf{P}^{1} \wedge \mathrm{KGL} \longrightarrow \mathrm{KGL}
$$

in $\operatorname{Spc}_{.}(F)$.
Definition 4.19. The algebraic $K$-theory spectrum KGL has constituent spaces $\mathrm{KGL}_{n}=\mathrm{KGL}$ and structure maps given by $e: \mathbf{P}^{1} \wedge$ KGL $\rightarrow$ KGL.

The spectrum KGL represents algebraic $K$-theory:
Theorem 4.20 ( $\overline{\text { Voe98, }}$, Theorem 6.9]). For any $X \in \operatorname{Sm}_{F}$ there are canonical isomorphisms

$$
K_{m-n}(X) \cong \operatorname{Hom}_{\mathcal{S H}(F)}\left(S^{m+n \alpha} \wedge X_{+}, \mathrm{KGL}\right)
$$

In particular, for $X=\operatorname{Spec} F$ we have $K_{m-n}(F) \cong \pi_{m+n \alpha} \mathrm{KGL}$.
In PPR09, it is shown that this spectrum is essentially unique, and a product

$$
\mu_{\mathrm{KGL}}: \mathrm{KGL} \wedge \mathrm{KGL} \longrightarrow \mathrm{KGL}
$$

is constructed, providing KGL with the structure of a ring spectrum PPR09, Theorem 2.2.1]. This product is the unique map in $\operatorname{Hom}_{\mathcal{S H}(F)}(\mathrm{KGL} \wedge \mathrm{KGL}, \mathrm{KGL})$ making the following diagram commutative:

$$
\begin{aligned}
& K_{0}(X) \times K_{0}(X) \\
& \cong \downarrow \\
& \pi_{m(1+\alpha)} \mathrm{KGL}(X) \otimes \pi_{n(1+\alpha)} \mathrm{KGL}(X) \xrightarrow[\mu_{\mathrm{KLL}}]{ } \pi_{(m+n)(1+\alpha)} \mathrm{KGL}(X)
\end{aligned}
$$

### 4.4 Motives

Voevodsky's construction of the derived category of motives - which predates the birth of motivic homotopy theory-allows for a definition of the long sought motivic cohomology theory. Originally it was Grothendieck and his collaborators who speculated upon the existence of such a universal cohomology theory. This section is therefore meant as a short "history lesson", although we conclude with a discussion on the connection between motivic cohomology-, and homotopy theory. In particular, we mention the "Hurewicz functor" $H: \mathcal{H}_{\bullet}(F) \rightarrow \mathrm{DM}_{-}^{\text {eff }}(F)$, relating homotopy on one side with homology on the other. We will also see in Theorem 4.30 that familiar objects like Milnor $K$-theory are found among the motivic cohomology groups.

The notion of a motive was introduced by Grothendieck in 1964 Mil14. The basic idea was that motives should define a category through which all Weil cohomology theories should factor. Grothendieck called the objects of this hypothetical category pure motives. Until now, the socalled mixed motives have been defined. However, the existence of pure motives is equivalent to the standard conjecture on algebraic cycles - a problem which is still unsolved.

Below we briefly illustrate the construction of Voevodsky's derived categories of motives $\mathrm{DM}_{\mathrm{gm}}(F)$ and $\mathrm{DM}_{-}^{\mathrm{eff}}(F)$. Roughly speaking, we start out with $\mathrm{Sm}_{F}$, i.e., smooth schemes over $F$, but we want this category to be Ab-enriched. So $\mathrm{Sm}_{F}$ is replaced with the category Cor $_{F}$ of finite correspondences, the definition of which relies on Suslin-Voevodsky's theory of relative cycles. Through a process of localization and idempotent completion we arrive at the so-called category of effective geometric motives $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$. Analogously to the situation in stable homotopy theory, the category $\mathrm{DM}_{\mathrm{gm}}(F)$ of geometric motives is then a stabilized version of $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ CD12.

Definition 4.21. Let $S$ be a regular scheme and let $X \in \operatorname{Sch}_{S}$ be an $S$-scheme. Let $Z(X)$ denote the free abelian group on the prime cycles of $X$ (see [Ful98]). The group $C_{\text {equi }}(X / S, 0)$ of relative cycles of dimension 0 is the subgroup of $Z(X)$ generated by the points $x \in X$ such that $\overline{\{x\}} \rightarrow S$ is finite and surjective (where $\overline{\{x\}}$ is given the reduced scheme structure).

The groups of relative cycles will form the Hom-sets in the category Cor $_{F}$ :
Definition 4.22. Let $X, Y \in \operatorname{Sm}_{F}$. The group of finite correspondences from $X$ to $Y$ is

$$
\operatorname{Cor}_{F}(X, Y):=C_{\text {equi }}\left(X \times_{F} Y / X, 0\right)
$$

A generator in $C_{\text {equi }}(X \times Y / X, 0)$ is called an elementary correspondence from $X$ to $Y$.
We can compose correspondences: The composition of two elementary correspondences $V \in \operatorname{Cor}_{F}(X, Y), W \in \operatorname{Cor}_{F}(Y, Z)$ is defined as the pushforward along the projection

$$
p: X \times Y \times Z \rightarrow X \times Z
$$

of the intersection product $[T]:=(V \times Z) \cdot(X \times W)$; see MVW06, p.4].
Definition 4.23. Let $\operatorname{Cor}_{F}$ be the category whose objects are the objects of $\mathrm{Sm}_{F}$, and whose morphisms are correspondences: $\operatorname{Hom}_{\operatorname{Cor}_{F}(X, Y)}:=C_{\text {equi }}(X \times Y / X, 0)$.

Definition 4.24. Let

$$
\operatorname{PST}(F):=\left[\operatorname{Cor}_{F}^{\mathrm{op}}, \mathrm{Ab}\right]
$$

denote category of presheaves with transfers. For $X$ a smooth $F$-scheme we let $\mathbf{Z}_{\text {tr }}(X)$ denote the representable presheaf with transfers $\mathbf{Z}_{\mathrm{tr}}(X):=\operatorname{Hom}_{\operatorname{Cor}_{F}}(-, X)$.

Defining $X \oplus Y:=X \amalg Y$, it follows that $\operatorname{Cor}_{F}$ is an additive category with the empty scheme as zero-object. The category $\mathrm{Cor}_{F}$ is also symmetric monoidal-see [MVW06, p.6] for the definition of the tensor product $\otimes^{\operatorname{tr}}$ on $\operatorname{Cor}_{F}$. The homotopy category $K^{b}\left(\operatorname{Cor}_{F}\right)$ of bounded chain complexes on $\mathrm{Cor}_{F}$ then becomes a tensor-triangulated category. We aim to localize the category $K^{b}\left(\operatorname{Cor}_{F}\right)$ with respect to a certain subcategory:

Definition 4.25. Let $\mathscr{B}$ be the localizing subcategory of $K^{b}\left(\operatorname{Cor}_{F}\right)$ generated by:

Homotopy invariance: For all $X \in \operatorname{Sm}_{F},\left[X \times \mathbf{A}^{1}\right] \rightarrow[X]$ is in $\mathscr{B}$.
Mayer-Vietoris: For each Zariski open covering $X=U \cup V$ of $X$, the map

$$
[U \cap V] \longrightarrow[U] \oplus[V] \longrightarrow[X]
$$

is in $\mathscr{B}$.
We then define the category $\widehat{\mathrm{DM}}_{\mathrm{gm}}^{\text {eff }}(F)$ as the localization of $K^{b}\left(\operatorname{Cor}_{F}\right)$ with respect to $\mathscr{B}$ :

$$
\widehat{\mathrm{DM}}_{\mathrm{gm}}^{\mathrm{eff}}(F):=K^{b}\left(\operatorname{Cor}_{F}\right) / \mathscr{B} .
$$

Finally, the category of effective geometric motives $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ is the idempotent completion (see Wei13, p.143]) of $\widehat{\mathrm{DM}}_{\mathrm{gm}}^{\mathrm{eff}}(F)$.

We have a functor $M_{\mathrm{gm}}: \mathrm{Sm}_{F} \rightarrow \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ defined by the composition

$$
\mathrm{Sm}_{F} \longrightarrow \operatorname{Cor}_{F} \longrightarrow K^{b}\left(\operatorname{Cor}_{F}\right) \longrightarrow \widehat{\mathrm{DM}}_{\mathrm{gm}}^{\mathrm{eff}}(F) \longrightarrow \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)
$$

Here $K^{b}\left(\operatorname{Cor}_{F}\right) \rightarrow \widehat{\mathrm{DM}}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ is the localization functor, and $\widehat{\mathrm{DM}}_{\mathrm{gm}}^{\mathrm{eff}}(F) \rightarrow \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ is the full embedding into the idempotent completion, defined by $X \mapsto\left(X, \mathrm{id}_{X}\right)$. For $X \in \operatorname{Sm}_{F}$, we call $M_{\mathrm{gm}}(X)$ the effective geometric motive of $X$. The functor $M_{\mathrm{gm}}$ is monoidal, satisfying $M_{\mathrm{gm}}(X \times Y)=M_{\mathrm{gm}}(X) \otimes M_{\mathrm{gm}}(Y)$. We write $\mathbf{Z}:=M_{\mathrm{gm}}(\operatorname{Spec} F)$, as this object is the unit for the tensor product on $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$.

The reduced effective geometric motive $\widetilde{M}_{\mathrm{gm}}(X)$ of $X \in \operatorname{Sm}_{F}$ is defined by the distinguished triangle

$$
\widetilde{M}_{\mathrm{gm}}(X) \longrightarrow M_{\mathrm{gm}}(X) \longrightarrow \mathbf{Z} \longrightarrow \widetilde{M}_{\mathrm{gm}}(X)[1]
$$

Definition 4.26. The Lefschetz motive is

$$
\mathbf{Z}(1):=\widetilde{M}_{\mathrm{gm}}\left(\mathbf{P}^{1}\right)[-2] .
$$

More generally, $\mathbf{Z}(n):=\mathbf{Z}(1)^{\otimes n}$ for any $n \geq 0$. Moreover, for any effective geometric motive $\mathscr{M}$ in $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ we define its twist by $n$ as $\mathscr{M}(n):=\mathscr{M} \otimes \mathbf{Z}(n)$.

To construct the category $\mathrm{DM}_{\mathrm{gm}}(F)$ of geometric motives we invert the Lefschetz motive:
Definition 4.27. The category of geometric motives $\mathrm{DM}_{\mathrm{gm}}(F)$ is defined as having as its objects pairs $(\mathscr{M}, m)$, where $\mathscr{M} \in \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ is an effective geometric motive and $m \in \mathbf{Z}$ is an integer. The morphisms are defined as

$$
\mathrm{DM}_{\mathrm{gm}}(F)((\mathscr{M}, m),(\mathscr{N}, n)):=\underset{k \geq-m,-n}{\operatorname{colim}_{\mathrm{gm}}} \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)(\mathscr{M}(k+m), \mathscr{N}(k+n)) .
$$

In the category $\mathrm{DM}_{\mathrm{gm}}(F)$, tensor product with the Lefschetz motive $\mathbf{Z}(1)$ is invertible. The process of inverting an object of $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ to obtain the category $\mathrm{DM}_{\mathrm{gm}}(F)$ is an analog of the Spanier-Whitehead construction in topology. In fact, one can carry out this process in any symmetric monoidal category [Voe98, p.588].

The Hom-groups in $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ are of great interest. In fact, they are one possible definition of motivic cohomology:

Definition 4.28. For $X \in \operatorname{Sm}_{F}$, let

$$
H^{p, q}(X, \mathbf{Z}):=\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)\left(M_{\mathrm{gm}}(X), \mathbf{Z}(q)[p]\right)
$$

be the motivic cohomology of $X$ with coefficients in $\mathbf{Z}$. The index $p$ is called the cohomological degree, and $q$ is called the weight.

Example 4.29. Suppose $X \in \mathrm{Sm}_{F}$ is a connected scheme. By [MVW06, Lecture 4] there is a quasi-isomorphism of motivic complexes $\mathbf{Z}(1) \xrightarrow{\sim} \mathcal{O}_{X}^{\times}[-1]$, where $\mathcal{O}_{X}^{\times}$is the sheaf of units in the structure sheaf on $X$. This theorem can be used to describe some motivic cohomology groups in weight one. The result is:

$$
H^{p, q}(X, \mathbf{Z})= \begin{cases}0, & q \leq 1,(p, q) \neq(0,0),(1,1),(2,1) \\ \mathbf{Z}, & (p, q)=(0,0) \\ \mathcal{O}_{X}^{\times}(X), & (p, q)=(1,1) \\ \operatorname{Pic}(X), & (p, q)=(2,1)\end{cases}
$$

Milnor $K$-theory is also found along the diagonal motivic cohomology groups:
Theorem 4.30 ([MVW06, Theorem 5.1]). For any integer $n \geq 0$ there is a natural isomorphism $H^{n, n}(\operatorname{Spec} F, \mathbf{Z}) \cong K_{n}^{M}(F)$.

The effective geometric motives can be realized as a thick subcategory of the category of effective motives:

Definition 4.31 (MVW06, p.109]). Let $\mathrm{DM}_{-}^{\text {eff }}(F)$ denote the localization

$$
\mathrm{DM}_{-}^{\mathrm{eff}}(F):=D^{-}\left(\operatorname{Shv}_{\mathrm{Nis}}\left(\operatorname{Cor}_{F}\right)\right) / \mathscr{A}
$$

where $\mathscr{A}$ is the localizing subcategory generated by complexes

$$
\mathbf{Z}_{\mathrm{tr}}\left(X \times \mathbf{A}^{1}\right) \xrightarrow{\mathbf{Z}_{\mathrm{tr}}(p)} \mathbf{Z}_{\mathrm{tr}}(X)
$$

for any $X \in \operatorname{Sm}_{F}$ with $p: X \times \mathbf{A}^{1} \rightarrow X$ being the canonical projection.
Theorem 4.32 (The embedding theorem [MVW06, p.110]). The category $\mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F)$ embeds as a full subcategory of $\mathrm{DM}_{-}^{\mathrm{eff}}(F)$.

Theorem 4.33 (MVW06, pp.110-111]). We list a few properties of the category $\mathrm{DM}_{-}^{\mathrm{eff}}(F)$.

- (Mayer-Vietoris) Given a Zariski open covering $X=U \cup V$ of $X \in \operatorname{Sm}_{F}$, there is a Mayer-Vietoris distinguished triangle

$$
M_{\mathrm{gm}}(U \cap V) \longrightarrow M_{\mathrm{gm}}(U) \oplus M_{\mathrm{gm}}(V) \longrightarrow M_{\mathrm{gm}}(X) \longrightarrow M_{\mathrm{gm}}(U \cap V)[1]
$$

in $\mathrm{DM}_{-}^{\mathrm{eff}}(F)$.

- If $\mathscr{E} \rightarrow X$ is a vector bundle, then the induced map $M_{\mathrm{gm}}(\mathscr{E}) \rightarrow M_{\mathrm{gm}}(X)$ is an isomorphism.
- (Projective bundle formula) Let $\mathscr{E} \rightarrow X$ be a vector bundle of rank d +1 , giving rise to the projective bundle $p: \mathbf{P}(\mathscr{E}) \rightarrow X$. The map $p$ induces an isomorphism

$$
\bigoplus_{n=0}^{d} M_{\mathrm{gm}}(X)(n)[2 n] \stackrel{\cong}{\cong} M_{\mathrm{gm}}(\mathbf{P}(\mathscr{E})) .
$$

- (Blowup triangle) If $Z$ is a smooth closed subscheme of $X \in \operatorname{Sm}_{F}$, let $p: \mathrm{Bl}_{Z}(X) \rightarrow X$ denote the blowup morphism. Then there is a distinguished triangle

$$
M_{\mathrm{gm}}\left(p^{-1}(Z)\right) \longrightarrow M_{\mathrm{gm}}(Z) \oplus M_{\mathrm{gm}}\left(\mathrm{Bl}_{Z}(X)\right) \longrightarrow M_{\mathrm{gm}}(X) \longrightarrow M_{\mathrm{gm}}\left(p^{-1}(Z)\right)[1]
$$

in $\mathrm{DM}_{-}^{\mathrm{eff}}(F)$.

### 4.5 Overview

The following diagram is meant as a short overview of all the categories constructed above and how they relate.


Here $i: \mathrm{DM}_{\mathrm{gm}}^{\mathrm{eff}}(F) \rightarrow \mathrm{DM}_{-}^{\mathrm{eff}}(F)$ is the full embedding from Theorem 4.32, and the "Hurewicz functor" $H: \mathcal{H}_{\bullet}(F) \rightarrow \mathrm{DM}_{-}^{\text {eff }}(F)$ is defined as follows. By Nisnevich sheafification we obtain a functor $\mathcal{M S} \mathcal{S}_{\bullet}(F) \rightarrow D^{-}\left(\operatorname{Shv}_{\mathrm{Nis}}\left(\operatorname{Cor}_{F}\right)\right)$. This functor sends $X \times \mathbf{A}^{1} \rightarrow X$ to $\mathbf{Z}_{\mathrm{tr}}(X \times$ $\left.\mathbf{A}^{1}\right) \rightarrow \mathbf{Z}_{\mathrm{tr}}(X)$, hence $\mathbf{A}^{1}$-weak equivalences of spaces are mapped to $\mathbf{A}^{1}$-weak equivalences of complexes. It follows that this functor preserves localization, i.e., it induces the functor $H$. By Wei04, Lemma 2], $H$ extends to a functor $\mathcal{S H}(F) \rightarrow \mathrm{DM}_{-}^{\text {eff }}(F)$. We refer the reader to Wei04 for further details.

## CHAPTER 5

## Symbols in Motivic Homotopy Theory

In Chapter 1 and Chapter 3 we have seen the importance of Steinberg symbols in number theory. The aim of this chapter is to illustrate how elements in certain motivic homotopy groups give rise to Steinberg symbols as well as Witt classes. These homotopy groups may perhaps be thought of as more fundamental, since the relations of both Milnor $K$-theory and Witt theory can be described by operations in stable motivic homotopy theory.

### 5.1 The Steinberg relation

Let $a$ denote a rational point $a: \operatorname{Spec}(F) \rightarrow \mathbf{G}_{m}$ of $\mathbf{G}_{m}$. Then $a$ induces a map of sheaves

$$
[a]: S^{0} \longrightarrow \mathbf{G}_{m}
$$

by sending the nonbase point to $a \in \mathbf{G}_{m}(F)$ and preserving the base point, and hence defines an element $[a] \in \pi_{-\alpha} \mathbf{1}$. If $b$ is another rational point of $\mathbf{G}_{m}$, then the smash product yields a map

$$
[a] \cdot[b] \in \operatorname{Hom}_{\mathcal{S H}(F)}\left(S^{0}, \mathbf{G}_{m} \wedge \mathbf{G}_{m}\right)
$$

that is, $[a] \cdot[b] \in \pi_{-2 \alpha} \mathbf{1}$. In HK01, Hu and Kriz show that these elements satisfy the Steinberg relation, i.e.,

$$
[a] \cdot[1-a]=0 \in \pi_{-2 \alpha} \mathbf{1} .
$$

Based on this result, Morel gives in Mor04a, Ch.6] a full description of the homotopy groups $\pi_{* \alpha} \mathbf{1}$ in terms of the graded Milnor-Witt ring $K_{*}^{M W}(F)$ (defined in Section 5.3):
Theorem 5.1 ([Mor04a, p.40]). If $F$ is a perfect field of characteristic not 2, then for each $n \in \mathbf{Z}$ there is an isomorphism

$$
\pi_{n \alpha} \mathbf{1} \xrightarrow{\cong} K_{-n}^{M W}(F)
$$

Observe that this isomorphism switches degrees.
In addition to the elements $[u] \in \pi_{\star} \mathbf{1}$, the following maps are of great interest in stable motivic homotopy theory:

- By performing the Hopf construction [DI13] on the multiplication map $\mathbf{G}_{m} \times \mathbf{G}_{m} \rightarrow \mathbf{G}_{m}$ one obtains a morphism

$$
\eta: \mathbf{A}^{2} \backslash\{0\}=S^{1+2 \alpha} \longrightarrow S^{1+\alpha}=\mathbf{P}^{1}
$$

given by $\mathbf{A}^{2} \backslash\{0\} \ni(x, y) \mapsto(x: y) \in \mathbf{P}^{1}$. This map represents an element $\eta \in \pi_{\alpha} \mathbf{1}$; it is the motivic analog of the Hopf fibration $\eta_{\mathrm{top}}: S^{3} \rightarrow S^{2}$ in topology. Indeed, we obtain $\eta_{\text {top }}$ from $\eta$ by taking complex realization:

$$
\eta_{\text {top }}: S^{1+2 \alpha}(\mathbf{C})=S^{3} \rightarrow S^{1+\alpha}(\mathbf{C})=S^{2}
$$

In Section 5.3 we will mention that multiplication with the image of $\eta$ in Milnor-Witt $K$-theory is an isomorphism. It follows that $\eta$ is not nilpotent. This is a major difference from the situation in topology, where $\eta_{\mathrm{top}}^{4}=0$.


Figure 5.1: Suggested picture of the real points of $\mathbf{G}_{m} \times \mathbf{G}_{m}$-where $\mathbf{G}_{m}$ is pointed by 1-and of $\mathbf{G}_{m} \wedge \mathbf{G}_{m}=\left(\mathbf{G}_{m} \times \mathbf{G}_{m}\right) /\left(\mathbf{G}_{m} \vee \mathbf{G}_{m}\right)$. The thick line is the graph of the Steinberg morphism $\mathbf{A}^{1} \backslash\{0,1\} \rightarrow \mathbf{G}_{m} \times \mathbf{G}_{m}, a \mapsto(a, 1-a)$. To the right we see that the image of this graph becomes connected in $\mathbf{G}_{m} \wedge \mathbf{G}_{m}$, suggesting that the Steinberg morphism should extend to $\mathbf{A}^{1}$.

- The twist morphism

$$
\epsilon: \mathbf{G}_{m} \wedge \mathbf{G}_{m} \longrightarrow \mathbf{G}_{m} \wedge \mathbf{G}_{m}
$$

defines a degree-zero element $\epsilon \in \pi_{0} \mathbf{1}$. Morel has shown that the elements $\eta$ and $\epsilon$ satisfy the relation

$$
\epsilon \eta=\eta
$$

and that the graded ring $\pi_{* \alpha} \mathbf{1}$ is graded $\epsilon$-commutative Mor04a.

- For $u \in F^{\times}$we let $\langle u\rangle$ denote the degree-zero element

$$
\langle u\rangle:=1+\eta[u] \in \pi_{0} \mathbf{1} .
$$

Then $\epsilon$ is given by multiplication by $-\langle-1\rangle$ Mor04a.
We will revisit all the elements considered above when we discuss Milnor-Witt $K$-theory in Section 5.3, although in a purely algebraic setting. But first we start out by explaining Hu and Kriz' proof HK01 of the Steinberg relation in motivic homotopy theory.

The idea is to show that, after a single simplicial suspension (i.e., smashing with the simplicial circle $S^{1}$ ), the "Steinberg morphism"

$$
\begin{equation*}
a \longmapsto(a, 1-a): \mathbf{A}^{1} \backslash\{0,1\} \longrightarrow \mathbf{G}_{m} \times \mathbf{G}_{m} \tag{5.1}
\end{equation*}
$$

becomes nullhomotopic when $\mathbf{G}_{m} \times \mathbf{G}_{m}$ is replaced by $\mathbf{G}_{m} \wedge \mathbf{G}_{m}$. In other words, we want to show that the composition

$$
\mathbf{A}^{1} \backslash\{0,1\} \longrightarrow \mathbf{G}_{m} \times \mathbf{G}_{m} \longrightarrow \mathbf{G}_{m} \wedge \mathbf{G}_{m}
$$

extends to $\mathbf{A}^{1}$ after simplicial suspension (see Figure 5.1).
To begin with, we need a lemma on Zariski excision. Let $X \in \mathrm{Sm}_{F}$ be a scheme and $X=U \cup V$ a Zariski open covering of $X$. Thus, in the category of motivic spaces there is a pushout square


Given such a covering of $X$, we have an analog in the motivic setting to excision in topology:

Lemma 5.2 (Zariski excision MP02). Let $X=U \cup V$ be a Zariski open covering of $X \in \operatorname{Sm}_{F}$. Then we have the following isomorphisms in the category $\mathcal{H}_{\bullet}(F)$ :

1. $U /(U \cap V) \cong X / V$;
2. $(X / U) \vee(X / V) \cong X /(U \cap V)$.

Proof. In the category of motivic spaces we have the diagram

in which the two small squares are pushouts, hence $X / V$ is the pushout of the large rectangle. But the pushout of the rectangle is $U /(U \cap V)$ by definition, hence $X / V \simeq U /(U \cap V)$.

For the second claim, we make use of the first result along with the following pushout diagram:


This yields an isomorphism

$$
\frac{U}{U \cap V} \vee \frac{V}{U \cap V} \xrightarrow{\sim} \frac{X}{U \cap V}
$$

By 1 we have $U /(U \cap V) \simeq X / V$ and $V /(U \cap V) \simeq X / U$, hence the result follows.
Now consider the space

$$
X_{1}:=\left(\mathbf{G}_{m} \times \mathbf{G}_{m}\right) \underset{\mathbf{G}_{m} \times\{1\}}{\amalg}\left(\mathbf{A}^{1} \times\{1\}\right) .
$$

That is, $X_{1}$ is the pushout of $\mathbf{G}_{m} \times \mathbf{G}_{m}$ and $\mathbf{A}^{1} \times\{1\}$ along $\mathbf{G}_{m} \times\{1\}$, with attaching maps the inclusions. For simplicity, we will in the following write $V:=[y-1=x z, y \neq 0]$ for the affine scheme

$$
V:=\operatorname{Spec}\left(F\left[x, y, y^{-1}, z\right] /(y-1-x z)\right)
$$

Let $\phi$ be the natural map

$$
\begin{aligned}
\phi: X_{1} & \longrightarrow V \\
(a, b) & \longmapsto\left(a, b, \frac{b-1}{a}\right) .
\end{aligned}
$$

Lemma 5.3 (HK01, Lemma 2]). The simplicial suspension $\Sigma_{s} \phi$ is a weak equivalence in $\mathcal{M S} \boldsymbol{\bullet}^{(F)}$.

Proof. We consider the diagram


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i.e., diagram (5) in the article [HK01]. Here the two vertical columns are cofiber sequences, where

$$
\begin{aligned}
C & :=\frac{\left(\mathbf{G}_{m} \times \mathbf{A}^{1}\right) \underset{\mathbf{G}_{m} \times\{1\}}{\amalg}\left(\mathbf{A}^{1} \times\{1\}\right)}{\left(\mathbf{G}_{m} \times \mathbf{G}_{m}\right) \underset{\mathbf{G}_{m} \times\{1\}}{\amalg}\left(\mathbf{A}^{1} \times\{1\}\right)} ; \\
D & :=\frac{[y-1=x z]}{[y-1=x z, y \neq 0]}
\end{aligned}
$$

Here the middle horizontal map is defined similarly as $\phi$, and $\psi$ is the induced map on cofibers.
Now we claim that both $\psi$ and the middle horizontal map are equivalences. If so, the statement follows, because we then have a diagram of the form

hence the 5-lemma for triangulated categories Wei94, Exercise 10.2.2] implies that $\Sigma_{s} \phi$ is an equivalence.

To prove the claim, we start out by showing that both sides of the middle horizontal map are $\mathbf{A}^{1}$-contractible. For the space to the right, note that the assignment

$$
\begin{aligned}
F[x, y, z] & \longrightarrow F[s, t], \\
x & \longmapsto s ; \\
y & \longmapsto 1+s t ; \\
z & \longmapsto t
\end{aligned}
$$

identifies

$$
F[x, y, z] /(y-1-x z) \cong F[s, t]
$$

hence

$$
[y-1=x z]=\operatorname{Spec}(F[x, y, z] /(y-1-x z)) \cong \mathbf{A}^{2},
$$

which is contractible in the motivic homotopy category. Turning the the space on the left, we use again that $\mathbf{A}^{1}$ is contractible to obtain:
$\operatorname{hocolim}\left(\mathbf{G}_{m} \times \mathbf{A}^{1} \longleftarrow \mathbf{G}_{m} \times\{1\} \longrightarrow \mathbf{A}^{1} \times\{1\}\right)=\operatorname{hocolim}\left(\mathbf{G}_{m} \longleftarrow \mathbf{G}_{m} \longrightarrow *\right)=*$.
Thus $\left(\mathbf{G}_{m} \times \mathbf{A}^{1}\right) \amalg_{\mathbf{G}_{m} \times\{1\}}\left(\mathbf{A}^{1} \times\{1\}\right) \simeq *$, and hence the middle horizontal map is an equivalence.
The next step is to show that $\psi$ is an equivalence. In fact, we show that $\psi$ is an isomorphism of sheaves. Note that, by Zariski excision (Lemma 5.2) we have

$$
D \cong \frac{[y-1=x z, y \neq 1]}{[y-1=x z, y \neq 0, y \neq 1]}
$$

We claim that the right hand side is isomorphic to

$$
C^{\prime}:=\frac{\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{1\}\right)}{\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{0,1\}\right)}
$$

Indeed,

$$
\begin{aligned}
\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{1\}\right) & \cong \operatorname{Spec}\left(F\left[t, t^{-1}\right] \otimes_{F} F\left[s,(s-1)^{-1}\right]\right) \\
& \cong \operatorname{Spec}\left(F\left[t, t^{-1}, s,(s-1)^{-1}\right]\right) \\
& \cong \operatorname{Spec}\left(F\left[u, u^{-1}, v, v^{-1}\right]\right)
\end{aligned}
$$

where the last isomorphism is induced by the map $F[u, v] \rightarrow F[s, t]$ sending $u$ to $s$ and $v$ to $s-1$. Furthermore, we have

$$
\begin{aligned}
{[y-1=x z, y \neq 1] } & =\operatorname{Spec}\left(F\left[x, y, z,(y-1)^{-1}\right] /(y-1-x z)\right) \\
& \cong \operatorname{Spec}\left(F\left[x, y,(x y)^{-1}\right]\right) \\
& =\operatorname{Spec}\left(F\left[x, y, x^{-1}, y^{-1}\right]\right),
\end{aligned}
$$

so that $\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{1\}\right) \cong[y-1=x z, y \neq 1]$, hence $D \cong C^{\prime}$.
Now consider the diagram


The square (1) is a pushout because the spaces $\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{0\}\right)$ and $\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{1\}\right)$ form an open cover of $\mathbf{G}_{m} \times \mathbf{A}^{1}$, with intersection $\mathbf{G}_{m} \times\left(\mathbf{A}^{1} \backslash\{0,1\}\right)$. The square (2) is a pushout by definition, and the space in the lower right corner of (3) is the pushout of the large rectangle made by (2) and (3), which must coincide with the pushout of (3). Finally, note that $C^{\prime}$ is the pushout of the large rectangle consisting of (1), (3) and (4). Since this pushout must coincide with the pushout of (4), we have $C \cong C^{\prime}$, thus $C \cong D$.

Consider the composition

$$
\mathbf{A}^{1} \backslash\{0,1\} \longrightarrow \mathbf{G}_{m} \times \mathbf{G}_{m} \longrightarrow X_{1} \xrightarrow{\phi} V
$$

where the first map is the Steinberg morphism 5.1 and the second map is the induced map on the pushout. This composition maps $a \in \mathbf{A}^{1} \backslash\{0,1\}$ to $(a, 1-a,-1) \in V$, and extends to a map

$$
\mathbf{A}^{1} \backslash\{1\} \longrightarrow V
$$

by letting $a \mapsto(a, 1-a,-1)$.
Symmetrically, define

$$
\begin{aligned}
X_{2} & :=\left(\mathbf{G}_{m} \times \mathbf{G}_{m}\right) \underset{\{1\} \times \mathbf{G}_{m}}{\amalg}\left(\{1\} \times \mathbf{A}^{1}\right), \\
W & :=[x-1=y z, x \neq 0]
\end{aligned}
$$

and a map

$$
\begin{aligned}
\phi^{\prime}: X_{2} & \longrightarrow W \\
(a, b) & \longmapsto\left(a, b, \frac{a-1}{b}\right) .
\end{aligned}
$$

Similarly as above, 5.1 extends to a map

$$
\begin{aligned}
\mathbf{A}^{1} \backslash\{0\} & \longrightarrow W \\
a & \longmapsto(a, 1-a,-1)
\end{aligned}
$$

Hence 5.1 also extends to a map

$$
\mathbf{A}^{1} \longrightarrow V \underset{\mathbf{G}_{m} \times \mathbf{G}_{m}}{\amalg} W
$$

where $V \amalg_{\mathbf{G}_{m} \times \mathbf{G}_{m}} W$ means the pushout of $V \stackrel{\phi}{\leftarrow} \mathbf{G}_{m} \times \mathbf{G}_{m} \xrightarrow{\phi^{\prime}} W$.
By a similar argument as the one given in Lemma 5.3, $\Sigma_{s} \phi^{\prime}$ is an equivalence.

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Theorem 5.4. The maps $\phi$ and $\phi^{\prime}$ induce an equivalence

$$
\Sigma_{s}\left(\mathbf{G}_{m} \wedge \mathbf{G}_{m}\right) \simeq \Sigma_{s}\left(V \underset{\mathbf{G}_{m} \times \mathbf{G}_{m}}{\amalg} W\right)
$$

Hence, after simplicial suspension, the Steinberg morphism 5.1 extends to

$$
\Sigma_{s} \mathbf{A}^{1} \longrightarrow \Sigma_{s}\left(\mathbf{G}_{m} \wedge \mathbf{G}_{m}\right)
$$

Proof. First note that the colimit of the diagram

is $\left(\mathbf{G}_{m} \times \mathbf{G}_{m}\right) /\left(\mathbf{G}_{m} \vee \mathbf{G}_{m}\right)=\mathbf{G}_{m} \wedge \mathbf{G}_{m}$. Now consider the diagram

where $P$ is the pushout of the lower right square. This diagram is an extended version of (11) in HK01. Since the arrows in the upper left half-square are cofibrations, $P$ coincides with the homotopy colimit of the diagram

which is $\mathbf{G}_{m} \wedge \mathbf{G}_{m}$ by the discussion above since $\mathbf{A}^{1} \simeq *$.
Finally, note that $\phi$ and $\phi^{\prime}$ induce maps from $P=\mathbf{G}_{m} \wedge \mathbf{G}_{m}$ to $V \amalg_{\mathbf{G}_{m} \times \mathbf{G}_{m}} W$. Since $\Sigma_{s} \phi$ and $\Sigma_{s} \phi^{\prime}$ are equivalences, the induced map on pushouts yields an equivalence

$$
\Sigma_{s}\left(\mathbf{G}_{m} \wedge \mathbf{G}_{m}\right) \simeq \Sigma_{s}\left(V \underset{\mathbf{G}_{m} \times \mathbf{G}_{m}}{\amalg} W\right),
$$

as desired.

### 5.2 On the ring spectrum map $1 \rightarrow \mathrm{KGL}$

Since the sphere spectrum $\mathbf{1}$ is initial in the category of ring spectra, there is a unique map from 1 to the algebraic $K$-theory spectrum KGL. This morphism induces a ring map

$$
u: \pi_{\star} \mathbf{1} \longrightarrow \pi_{\star} \mathrm{KGL}
$$

The unit map $\mathbf{1} \rightarrow \mathrm{KGL}$ corresponds to a choice of generator for $K_{0}(F) \cong \mathbf{Z}$; indeed, we have

$$
\begin{aligned}
\operatorname{Hom}_{\mathcal{S H}(F)}(\mathbf{1}, \mathrm{KGL}) & =\pi_{0} \mathrm{KGL}(F) \\
& \cong K_{0}(F) \\
& \cong \mathbf{Z}\{[F]\}
\end{aligned}
$$

by Chapter 2, Example 2.2. By composing with the unit map $\mathbf{1} \rightarrow \mathrm{KGL}$ we obtain morphisms

$$
\pi_{* \alpha} 1 \longrightarrow \pi_{* \alpha} \mathrm{KGL} .
$$

Below we will consider the map

$$
\pi_{-\alpha} \mathbf{1}=\operatorname{Hom}_{\mathcal{S H}(F)}\left(\Sigma_{t}^{-1} \mathbf{1}, \mathbf{1}\right) \longrightarrow \operatorname{Hom}_{\mathcal{S H}(F)}\left(\Sigma_{t}^{-1} \mathbf{1}, \mathrm{KGL}\right)=\pi_{-\alpha} \mathrm{KGL}
$$

where $\Sigma_{t}=-\wedge \mathbf{G}_{m}$ denotes $\mathbf{G}_{m}$-suspension.
Proposition 5.5. The map

$$
u: \pi_{-2 \alpha} \mathbf{1} \longrightarrow \pi_{-2 \alpha} \mathrm{KGL}
$$

sends the element $[a] \cdot[b] \in \pi_{-2 \alpha} \mathbf{1}$ to $\{a, b\} \in K_{2}(F)=\pi_{-2 \alpha}$ KGL.
Sketch of proof. By the identification $\pi_{-\alpha} \mathrm{KGL} \cong K_{1}(F)$, it follows that the map $u$ sends $[a] \in$ $\pi_{-\alpha} \mathbf{1}$ to $\{a\} \in K_{1}(F)$. Since $u$ is a ring map we obtain a commutative diagram


The bottom horizontal map is the product on $K_{1}(F)$, which sends $(\{a\},\{b\}) \in K_{1}(F) \times K_{1}(F)$ to $\{a, b\} \in K_{2}(F)$ (cf. Theorem 2.15). The top horizontal morphism is the product in $\pi_{-* \alpha} \mathbf{1}$, mapping $([a],[b])$ to $[a][b] \in \pi_{-2 \alpha} 1$. Thus $u([a][b])=\{a, b\} \in K_{2}(F)$.

Proposition 5.5 tells us that many of the relations in $K$-theory comes from relations in motivic homotopy groups of spheres. We move on to study these fundamental objects nextalbeit in the form of Milnor-Witt $K$-theory.

### 5.3 The Milnor-Witt ring

As mentioned at the beginning of this chapter, the so-called Milnor-Witt $K$-theory gives an exact description of certain operations in stable motivic homotopy theory. In this section, we proceed to study the Milnor-Witt $K$-groups in their own right.

Definition 5.6 (Hopkins-Morel [Mor04a, p.40]). The Milnor-Witt K-theory $K_{*}^{M W}(F)$ of $F$ is the $\mathbf{Z}$-graded associative ring freely generated by symbols $[x]\left(x \in F^{\times}\right)$in degree +1 , and one symbol $\eta$ of degree -1 , subject to the following relations:

1. $[x][1-x]=0$ for all $x \in F^{\times} \backslash\{1\}$,
2. $[x y]=[x]+[y]+\eta[x][y]$,
3. $\eta[x]=[x] \eta$,
4. $\eta h=0$, where $h:=2+\eta[-1] \in K_{0}^{M W}(F)$
for $x, y \in F^{\times}$.

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There is a natural map

$$
K_{*}^{M W}(F) \longrightarrow K_{*}^{M}(F)
$$

sending $[x]$ to $\{x\}$ and $\eta$ to 0 , which yields a canonical isomorphism

$$
K_{*}^{M W}(F) / \eta \cong K_{*}^{M}(F)
$$

Indeed, the term $\eta[x][y]$ in the second relation of Milnor-Witt $K$-theory is the obstruction for additivity in $K_{*}^{M W}(F)$.

Definition 5.7. For $u \in F^{\times}$, let $\langle u\rangle$ denote the degree-zero element

$$
\langle u\rangle:=1+\eta[u] \in K_{0}^{M W}(F) .
$$

In particular, let $\epsilon:=-\langle-1\rangle$.
Note the relation

$$
\epsilon \eta=\eta
$$

Moreover, we have that

$$
\eta(1+\langle-1\rangle)=0
$$

which is reminiscent of the defining relation of the Witt ring. In fact, the map $F^{\times} \rightarrow K_{0}^{M W}(F)$ defined by $u \mapsto\langle u\rangle$ induces an isomorphism $G W(F) \stackrel{\cong}{\leftrightarrows} K_{0}^{M W}(F)$ [Mor12a, p.65], and for $n>0$, multiplication with $\eta^{n}$ induces an isomorphism $W(F) \xrightarrow{\cong} K_{-n}^{M W}(F)$ Mor12a, Lemma 2.10, p.65].

In addition to the maps $K_{n}^{M W}(F) \rightarrow K_{n}^{M}(F)$ there are, for $n \geq 1$, canonical homomorphisms

$$
K_{n}^{M W}(F) \longrightarrow I^{n}(F),
$$

which maps an element $\left[a_{1}\right] \cdots\left[a_{n}\right]$ of $K_{n}^{M W}(F)$ to the class of the Pfister form $\left\langle\left\langle a_{1} \ldots, a_{n}\right\rangle\right\rangle \in$ $I^{n}(F)$ (defined in Chapter 2. Definition 2.30. A theorem of Morel shows how these maps relate:

Theorem 5.8 ([Mor04b]). The diagram

is a pullback square.
Example 5.9. Since $I^{2}\left(\mathbf{F}_{q}\right)=0$ (cf. Table 2.1 in Chapter 2) for any finite field $\mathbf{F}_{q}$, Theorem 5.8 implies that $K_{n}^{M W}\left(\mathbf{F}_{q}\right) \cong K_{n}^{M}\left(\mathbf{F}_{q}\right)$ for all $n \geq 1$.

Lemma 5.10 (Mor12a, Lemma 2.5, p. 62 and Lemma 2.7, p.63]). For any $x, y \in F^{\times}$, the following relations hold in $K_{*}^{M W}(F)$ :

1. $[1]=0$.
2. $[x][-x]=0$.
3. $[x][x]=[x][-1]=[-1][x]$.
4. $[x][y]=\epsilon[y][x]$.
5. $\left\langle x^{2}\right\rangle=1$.

Lemma 5.11. The element $\langle x\rangle=1+\eta[x]$ is a unit in $K_{0}^{M W}(F)$, whose inverse is $\left\langle x^{-1}\right\rangle$. Moreover, for any $x, y \in F^{\times}$, the following relations hold in $K_{*}^{M W}(F)$ :

1. $\left[\frac{x}{y}\right]=[x]-\left\langle\frac{x}{y}\right\rangle[y]$, hence $\left[x^{-1}\right]=-\left\langle x^{-1}\right\rangle[x]$.
2. $\langle x y\rangle=\langle x\rangle\langle y\rangle$.
3. $[x][y]=\left[y^{-1}\right][x]$.
4. $[x y]=[x]+\langle x\rangle[y]$.
5. $\left[x^{2}\right][y]=[x][y]-[y][x]$.

Proof. We only prove the last relation. Proof of the first assertion along with Item 1 and Item 2 is found in Mor12a, Lemma 2.5, p.62]. Moreover, Item3 is proved in [HT08, Lemma 2.7], and Item 4 follows from the definitions.

By definition of Milnor-Witt $K$-theory we have

$$
\left[x^{2}\right]=2[x]+\eta[x]^{2}
$$

Multiplying from the right by $[y]$ then yields

$$
\left[x^{2}\right][y]=2[x][y]+\eta[x]^{2}[y] .
$$

On the other hand, using the $\epsilon$-commutativity $[x][y]=\epsilon[y][x]$ we have

$$
\begin{aligned}
{[x][y]-[y][x] } & =[x][y]-\epsilon[x][y] \\
& =[x][y]+\langle-1\rangle[x][y] \\
& =[x][y]+(1+\eta[-1])[x][y] \\
& =2[x][y]+\eta[-1][x][y] \\
& =2[x][y]+\eta[x]^{2}[y]=\left[x^{2}\right][y],
\end{aligned}
$$

where we have used that $[-1][x]=[x]^{2}$.
Proposition 5.12 (Mor12a, Proposition 2.13, p.67]). If any unit in $F$ is a square, then for each $n \geq 0$ the canonical surjection

$$
K_{n}^{M W}(F) \longrightarrow K_{n}^{M}(F)
$$

is an isomorphism.
Proof. By the second and fourth relation of Milnor-Witt $K$-theory along with the relation $[x][x]=[x][-1]$ we have

$$
\eta\left[x^{2}\right]=\eta(2[x]+\eta[x][-1])=\eta[x](2+\eta[-1])=0
$$

Since any unit in $F$ is a square, this says that $\eta[z]=0$ for all $z \in F^{\times}$. This implies additivity in $K_{n}^{M W}(F)$, i.e.,

$$
[x y]=[x]+[y]+\eta[x][y]=[x]+[y] .
$$

Therefore the map $K_{n}^{M}(F) \rightarrow K_{n}^{M W}(F)$ which sends $\left\{x_{1}, \ldots, x_{n}\right\}$ to $\left[x_{1}\right] \cdots\left[x_{n}\right]$ is well defined. Composing with the canonical map $K_{n}^{M W}(F) \rightarrow K_{n}^{M}(F)$, which sends $\eta$ to 0 , yields the respective identities, so this map is an isomorphism.

Lemma 5.13 (Mor12a, Lemma 2.6, p.63]). For all $n \geq 1$, the group $K_{n}^{M W}(F)$ is generated by the products of the form $\left[u_{1}\right]\left[u_{2}\right] \cdots\left[u_{n}\right]$, with $\left[u_{j}\right] \in F^{\times}$.

Proof. Clearly any generator of $K_{n}^{M W}(F)$ is of the form $\eta^{m}\left[u_{1}\right] \cdots\left[u_{\ell}\right]$ where $m, \ell \geq 0$ and $\ell-m=n$. The result follows from the relation $\eta[x][y]=[x y]-[x]-[y]$.

Lemma 5.14. Suppose $v$ is a discrete valuation on $F$. If $\pi$ and $\pi^{\prime}$ are two uniformizers, then

$$
[\pi]\left[\pi^{\prime}\right]=[\pi][u]
$$

for a unit $u \in \mathcal{O}_{v}^{\times}$.

## 5. Symbols in Motivic Homotopy Theory

Proof. We may write $\pi^{\prime}=-u \pi$ for $u$ a unit of $\mathcal{O}_{v}$. Then

$$
\begin{aligned}
{[\pi][u(-\pi)] } & =[\pi]([u]+[-\pi]+\eta[u][-\pi]) \\
& =[\pi][u]+[\pi][-\pi]+\eta \epsilon[u][\pi][-\pi] \\
& =[\pi][u]
\end{aligned}
$$

since $[\pi][-\pi]=0$.
Lemma 5.15 (Mor12a, Lemma 2.14, p.67]). For any $x \in F^{\times}$and any $n \in \mathbf{Z}$, we have the following expression for $\left[x^{n}\right] \in K_{1}^{M W}(F)$ :

$$
\left[x^{n}\right]=n_{\epsilon}[x],
$$

where $n_{\epsilon} \in K_{0}^{M W}(F)$ is defined as follows. If $n \geq 0$, we set

$$
n_{\epsilon}:=\sum_{j=1}^{n}\left\langle(-1)^{j-1}\right\rangle,
$$

and for $n<0, n_{\epsilon}:=-\langle-1\rangle(-n)_{\epsilon}=\epsilon(-n)_{\epsilon}$.
Lemma 5.16 (Mor12a, Corollary 2.8, p.63]). Given $\alpha \in K_{n}^{M W}(F)$ and $\beta \in K_{m}^{M W}(F)$, we have

$$
\alpha \beta=\epsilon^{n m} \beta \alpha .
$$

Thus the graded $K_{0}^{M W}(F)$-algebra $K_{*}^{M W}(F)$ is graded $\epsilon$-commutative.
In Mor12a, Morel establishes the existence of an analog in Milnor-Witt $K$-theory to the well known regular symbols in Milnor $K$-theory:

Theorem 5.17 (Mor12a, Section 2.2]). For each discrete valuation $v$ on $F$ and each choice of uniformizer $\pi$, there exists a unique homomorphism

$$
\partial_{v}^{\pi}: K_{*}^{M W}(F) \longrightarrow K_{*}^{M W}(k(v))
$$

commuting with product by $\eta$ and satisfying

$$
\partial_{v}^{\pi}\left([\pi]\left[u_{2}\right] \cdots\left[u_{n}\right]\right)=\left[\bar{u}_{2}\right] \cdots\left[\bar{u}_{n}\right]
$$

and

$$
\partial_{v}^{\pi}\left(\left[u_{1}\right]\left[u_{2}\right] \cdots\left[u_{n}\right]\right)=0
$$

for all units $u_{1}, \ldots, u_{n} \in \mathcal{O}_{v}^{\times}$.
Remark 5.18. In Milnor-Witt $K$-theory, the maps $\partial_{v}^{\pi}$ depends on the choice of uniformizer. For if $\pi^{\prime}=u \pi$, where $u \in \mathcal{O}_{v}^{\times}$, then $\left[\pi^{\prime}\right]=[u]+[\pi]+\eta[u][\pi]$. This differs from the situation in Milnor $K$-theory, where $\left\{\pi^{\prime}\right\}=\{u\}+\{\pi\}$, hence $\partial_{v}\left(\pi^{\prime}\right)=\partial_{v}(\pi)$ in this case.

Using the technique of Tate, Morel establishes in Mor12a, Theorem 2.24, p.73] a split exact sequence

$$
0 \longrightarrow K_{n}^{M W}(F) \longrightarrow K_{n}^{M W}(F(t)) \xrightarrow{\oplus_{\mathbf{p}} \partial_{\mathrm{p}}^{\boldsymbol{q}}} \underset{\mathfrak{p} \in \mathrm{Spec}(F[t]) \backslash\{0\}}{ } K_{n-1}^{M W}(F[t] / \mathfrak{p}) \longrightarrow 0,
$$

similarly to the sequence in Milnor $K$-theory.
Theorem 5.17 also allows us to define sheaves of unramified Milnor-Witt $K$-groups CF14, p.5]. Let $X \in \operatorname{Sm}_{F}$ be an integral scheme, and let $X^{(n)}$ denote the set of points $x \in X$ of codimension $n$. We let $k(X)$ denote the field of rational functions on $X$.

Definition 5.19. For any $X \in \operatorname{Sm}_{F}$, choose for each $x \in X^{(1)}$ a corresponding uniformizer $\pi_{x}$. Then, for each $n \in \mathbf{Z}$ we define

$$
K_{n}^{M W}(X):=\operatorname{ker}\left(K_{n}^{M W}(k(X)) \longrightarrow \bigoplus_{x \in X^{(1)}} K_{n-1}^{M W}(k(x))\right)
$$

where the map is induced by the morphisms $\partial_{v_{x}}^{\pi_{x}}$.
Any morphism $f: X \rightarrow Y$ in $\operatorname{Sm}_{F}$ induces a homomorphism $h^{*}: K_{n}^{M W}(Y) \rightarrow K_{n}^{M W}(X)$ CF14, p.6], so that $X \mapsto K_{n}^{M W}(X)$ is a presheaf on $\mathrm{Sm}_{F}$. This presheaf is in fact a Nisnevich sheaf Mor12a, which is denoted by $\mathbf{K}_{n}^{M W}$.

### 5.4 Milnor-Witt $K$-theory of Q

In this section we aim to compute the groups $K_{n}^{M W}(\mathbf{Q})$ for all $n \geq 2$.
Definition 5.20. Let $\Lambda_{\infty} \subseteq K_{*}^{M W} \mathbf{Q}$ be the graded subring of $K_{*}^{M W} \mathbf{Q}$ generated by the elements $[-1]^{n} \in K_{n}^{M W} \mathbf{Q}$ for $n \geq 0$, and let $\Lambda_{\infty}^{n}$ denote the $n$-th graded piece of $\Lambda_{\infty}$. Thus $\Lambda_{\infty}^{n}$ is the subgroup of $K_{n}^{M W} \mathbf{Q}$ additively generated by $[-1]^{n}$.

Moreover, let $\Lambda_{0}$ be the graded subring of $K_{*}^{M W} \mathbf{Q}$ generated by $\eta$ and the elements $[a]$ with $a$ a positive rational number:

$$
\Lambda_{0}:=\left\langle\eta,[a]: a \in \mathbf{Q}_{>0}\right\rangle
$$

For any $n \in \mathbf{Z}$ we let $\Lambda_{0}^{n}$ denote the $n$-th graded piece of $\Lambda_{0}$.
Lemma 5.21. For any $n \geq 2$ we have $K_{n}^{M W} \mathbf{Q}=\Lambda_{\infty}^{n} \oplus \Lambda_{0}^{n}$.
Proof. The relation $[-1][a]=[a][a]$ implies that

$$
([-1]+x)([-1]+y)=[-1]^{2}+[-1] x+[-1] y+x y \in \Lambda_{\infty} \oplus \Lambda_{0}
$$

for any $x, y \in \Lambda_{0}$.
Now, fix an $n \geq 2$ and recall from Lemma 5.13 that $K_{n}^{M W} \mathbf{Q}$ is generated by the symbols $\left[a_{1}\right] \cdots\left[a_{n}\right]$ for which $a_{i} \in \mathbf{Q}^{\times}$. If $a_{1}, \ldots, a_{n} \in \mathbf{Q}_{>0}$, then the observation above along with the relation $[-a]=[-1]+[a]+\eta[a]^{2}$ implies that any generator $\left[ \pm a_{1}\right] \cdots\left[ \pm a_{n}\right]$ of $K_{n}^{M W}(\mathbf{Q})$ lies in $\Lambda_{\infty}^{n} \oplus \Lambda_{0}^{n}$.

To conclude we must show that $\Lambda_{\infty}^{n} \cap \Lambda_{0}^{n}=0$. Recall again that $\Lambda_{0}^{n}$ is generated by elements of the form $\left[a_{1}\right] \cdots\left[a_{n}\right]$ for $a_{i} \in \mathbf{Q}_{>0}$. So suppose that $k[-1]^{n} \in \Lambda_{0}^{n}$ for some $k \neq 0$, say

$$
k[-1]^{n}=\sum_{i} k_{i}\left[a_{1 i}\right] \cdots\left[a_{n i}\right]
$$

with all the $a_{j i}$ 's positive rationals and $k_{i} \in \mathbf{Z}$. Recall that the canonical map

$$
K_{n}^{M W}(\mathbf{Q}) \longrightarrow I^{n}(\mathbf{Q})
$$

maps a generator $\left[a_{1}\right] \cdots\left[a_{n}\right] \in K_{n}^{M W} \mathbf{Q}$ to the Pfister form

$$
\left\langle\left\langle a_{1}, \ldots, a_{n}\right\rangle\right\rangle=\left(\langle 1\rangle-\left\langle a_{1}\right\rangle\right) \cdots\left(\langle 1\rangle-\left\langle a_{n}\right\rangle\right) \in I^{n}(\mathbf{Q}) .
$$

Hence, passing to the Witt ring $W(\mathbf{Q})$ we obtain the relation of Pfister forms

$$
\begin{equation*}
k\langle\langle-1, \ldots,-1\rangle\rangle=\sum_{i} k_{i}\left\langle\left\langle a_{1 i}, \ldots, a_{n i}\right\rangle\right\rangle \in I^{n}(\mathbf{Q}) \tag{5.2}
\end{equation*}
$$

We claim that such a relation is impossible. Indeed, consider the natural map $W(\mathbf{Q}) \rightarrow W(\mathbf{R})$ induced by the inclusion $\mathbf{Q} \hookrightarrow \mathbf{R}$, giving rise to the same relation (5.2) in $W(\mathbf{R})$. Let $P$ denote the standard ordering on $\mathbf{R}$. Note that, since the $a_{j i}$ 's are positive we have $\sigma_{P}\left(\left\langle a_{j i}\right\rangle\right)=+1$ for all $i, j$ by Example 2.34. Hence

$$
\sigma_{P}\left(\left\langle\left\langle a_{1 i}, \ldots, a_{n i}\right\rangle\right\rangle\right)=\sigma_{P}\left(\left(\langle 1\rangle-\left\langle a_{1 i}\right\rangle\right) \cdots\left(\langle 1\rangle-\left\langle a_{n i}\right\rangle\right)\right)=\prod_{j}\left(\sigma_{P}\langle 1\rangle-\sigma_{P}\left\langle a_{j i}\right\rangle\right)=0
$$

On the other hand, a similar computation shows that $\sigma_{P}(\langle\langle-1, \ldots,-1\rangle\rangle)=2^{n}$, hence the relation 5.2 is impossible whenever $k \neq 0$.

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Lemma 5.22. For any $n \geq 1$, the group $\Lambda_{\infty}^{n}$ is the free additive group generated by $[-1]^{n}$.
Proof. Note that, since $\langle-a\rangle=-\langle a\rangle$ in the Witt ring, it follows that the canonical morphism $K_{n}^{M W} \mathbf{Q} \rightarrow I^{n}(\mathbf{Q})$ maps $k[-1]^{n}$ to

$$
k\langle\langle-1, \ldots,-1\rangle\rangle=k(\langle 1\rangle-\langle-1\rangle)^{n}=k(2\langle 1\rangle)^{n}=k \cdot 2^{n}\langle 1\rangle .
$$

By [MH73, Theorem 4.5, p.75], $m\langle 1\rangle \neq 0 \in W(\mathbf{Q})$ for any nonzero integer $m$. Hence $k[-1]^{n} \neq$ $\ell[-1]^{n}$ whenever $k \neq \ell$.

Now we proceed in the same manner as for Milnor $K$-theory to show that

$$
\Lambda_{0}^{n} \cong \bigoplus_{p \geq 2} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)
$$

for any $n \geq 2$. For each prime $p$ and each $n \geq 2$, Theorem 5.17 establishes a map

$$
\partial:=\bigoplus_{p \geq 2 \text { prime }} \partial_{p}: K_{n}^{M W}(\mathbf{Q}) \longrightarrow \bigoplus_{p \geq 2 \text { prime }} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right),
$$

where we write $\partial_{p}:=\partial_{(p)}^{p} ; p$ being the canonical choice of uniformizer for the place $(p)$. Note that $\Lambda_{\infty}^{n} \subseteq \operatorname{ker} \partial$ by Theorem 5.17.

Definition 5.23. Define a filtration

$$
L_{2} \subseteq L_{3} \subseteq \cdots \subseteq L_{p} \subseteq \cdots \subseteq \Lambda_{0}
$$

of $\Lambda_{0} \subseteq K_{*}^{M W} \mathbf{Q}$ by letting $L_{p}$ be the subring

$$
L_{p}:=\langle\eta,[a]: a \in \mathbf{Z}, 1 \leq a \leq p\rangle
$$

For each $n \in \mathbf{Z}$, we let $L_{p}^{n}$ denote the $n$-th graded piece of $L_{p}$.
By Lemma 5.11, if $[a]$ and $[b]$ are elements of $L_{p}$ then so are $[a b]$ and $[a / b]$. Thus, by Lemma 5.13 this filtration is exhaustive, i.e., $\lim _{p} L_{p}=\Lambda_{0}$.

Lemma 5.24. For each $n \geq 2$, the group $L_{2}^{n}$ is trivial.
Proof. We have $[2][2]=[2][-1]=0$ by the Steinberg relation. Thus any generator $\left[a_{1}\right] \cdots\left[a_{n}\right]$ with $a_{i} \in\{1,2\}$ is trivial.

Lemma 5.25. Let $n \geq 1$. For any prime $p>2$, let $q$ be the largest prime $<p$. Then $L_{p}^{n}$ is generated by $L_{q}^{n}$ along with elements of the form

$$
\eta^{m}[p]\left[a_{2}\right]\left[a_{3}\right] \cdots\left[a_{n+m}\right]
$$

where the $a_{j} \in \mathbf{Z}$ are units modulo $p$ for all $j \geq 2$.
Proof. Let $x=\eta^{m}\left[a_{1}\right] \cdots\left[a_{m+n}\right]$ be a generator of $L_{p}^{n}$, with $1 \leq a_{i} \leq p$. Using $\epsilon$-commutativity we may assume $x$ is of the form $x=\eta^{m}[p]^{k}\left[a_{k+1}\right] \cdots\left[a_{n+m}\right]$ for some $k \leq n+m$ and $a_{i}$ units modulo $p$. But then the result follows from the relation $[p][p]=[-1][p]=[p][-1]$.

Lemma 5.26. Let $p$ be an odd prime and let $q$ be the biggest prime $<p$. Suppose $i \geq 2$ and $a_{j} \in \mathbf{Z}$ such that $a_{1} \cdots a_{i} \equiv c(\bmod p)$, say $\prod_{j} a_{j}=c+p d$, where $|c|<p$ and $\left|a_{j}\right|<p$ for all $j$. Then

$$
[p]\left[a_{1} \cdots a_{i}\right] \equiv[p][c]\left(\bmod L_{q}\right)
$$

Proof (following [Mor12a, Lemma 2.25, p.74]). We proceed by induction. For the case $i=2$, assume $a b=c+p d$, with $|a|,|b|,|c|<p$. Note that

$$
|p d| \leq|a b|+|c| \leq(p-1)^{2}+p-1
$$

hence $|d|<p$. Since

$$
\frac{p d}{a b}+\frac{c}{a b}=1
$$

we have

$$
\left[\frac{p d}{a b}\right]\left[\frac{c}{a b}\right]=0
$$

Moreover,

$$
\left[\frac{p d}{a b}\right]=\left[\frac{d}{a b}\right]+\left\langle\frac{d}{a b}\right\rangle[p]
$$

by Item 4 of Lemma 5.11. Thus

$$
\begin{aligned}
0 & =\left[\frac{p d}{a b}\right]\left[\frac{c}{a b}\right] \\
& =\left(\left[\frac{d}{a b}\right]+\left\langle\frac{d}{a b}\right\rangle[p]\right)\left[\frac{c}{a b}\right]
\end{aligned}
$$

i.e.,

$$
\left\langle\frac{d}{a b}\right\rangle[p]\left[\frac{c}{a b}\right]=-\left[\frac{d}{a b}\right]\left[\frac{c}{a b}\right]
$$

Multiplying from the left by $\langle d /(a b)\rangle^{-1}=\langle a b / d\rangle$ we obtain

$$
[p]\left[\frac{c}{a b}\right]=-\left\langle\frac{a b}{d}\right\rangle\left[\frac{d}{a b}\right]\left[\frac{c}{a b}\right] \in L_{q}
$$

Note that

$$
\begin{aligned}
{\left[\frac{x}{y}\right] } & =[x]-\left\langle\frac{x}{y}\right\rangle[y] \\
& =[x]-[y]-\eta\left[\frac{x}{y}\right][y]
\end{aligned}
$$

by Lemma 5.11. Using this relation on $[c /(a b)]$ and multiplying from the left by $[p]$ get that

$$
\begin{aligned}
{[p]\left[\frac{c}{a b}\right] } & =[p][c]-[p][a b]-\eta[p][a b]\left[\frac{c}{a b}\right] \\
& =[p][c]-[p][a b]-\eta \epsilon[a b][p]\left[\frac{c}{a b}\right] .
\end{aligned}
$$

Since

$$
\begin{aligned}
-\eta \epsilon[a b][p]\left[\frac{c}{a b}\right] & =\eta\langle-1\rangle[a b][p]\left[\frac{c}{a b}\right] \\
& =\eta(1+\eta[-1])[a b][p]\left[\frac{c}{a b}\right] \\
& =\eta[a b][p]\left[\frac{c}{a b}\right]+\eta^{2}[a b]^{2}[p]\left[\frac{c}{a b}\right] \in L_{q}
\end{aligned}
$$

we thus have

$$
[p][c] \equiv[p][a b]\left(\bmod L_{q}\right)
$$

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For the general case $i>2$, write $a_{2} \cdots a_{i} \equiv c^{\prime}(\bmod p)$. Then $c \equiv a_{1} \cdots a_{i} \equiv a_{1} c^{\prime}(\bmod p)$. By the second relation of Milnor-Witt $K$-theory and the induction hypothesis we have

$$
\begin{aligned}
{[p]\left[a_{1} \cdots a_{i}\right] } & =[p]\left[a_{1}\right]+[p]\left[a_{2} \cdots a_{i}\right]+\eta[p]\left[a_{1}\right]\left[a_{2} \cdots a_{i}\right] \\
& \equiv[p]\left[a_{1}\right]+[p]\left[c^{\prime}\right]+\eta[p]\left[a_{1}\right]\left[c^{\prime}\right]\left(\bmod L_{q}\right) .
\end{aligned}
$$

Now $[p]\left[a_{1}\right]+[p]\left[c^{\prime}\right]+\eta[p]\left[a_{1}\right]\left[c^{\prime}\right]=[p]\left[a_{1} c^{\prime}\right]$, which we know from the case $i=2$ equals $[p][c]$ modulo $L_{q}$. Hence $[p]\left[a_{1} \cdots a_{i}\right] \equiv[p][c]\left(\bmod L_{q}\right)$.

Theorem 5.27. For each $n \geq 2$, the map $\partial=\bigoplus_{p} \partial_{p}$ induces an isomorphism

$$
K_{n}^{M W}(\mathbf{Q}) \cong \mathbf{Z} \oplus \bigoplus_{p \geq 2} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right) \cong \mathbf{Z} \oplus \bigoplus_{p \geq 2} K_{n-1}^{M}\left(\mathbf{F}_{p}\right)
$$

Proof. The proof is similar to [Mor12a, p.76]. We fix an $n \geq 2$. By Lemma 5.21 and Lemma 5.22 it remains to show that

$$
\Lambda_{0}^{n} \cong \bigoplus_{p \geq 2} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)
$$

As before, we use induction to prove that for each prime $p, L_{p}^{n} \cong \bigoplus_{\ell=2}^{p} K_{n-1}^{M W}\left(\mathbf{F}_{\ell}\right)$; the statement then follows.

To start the induction, note that $K_{m}^{M W}\left(\mathbf{F}_{2}\right) \cong K_{m}^{M}\left(\mathbf{F}_{2}\right)$ for each $m \geq 0$ by Proposition 5.12 . But $K_{m}^{M}\left(\mathbf{F}_{2}\right)=0$ for $m \geq 1$, hence Lemma 5.24 furnishes the base case of the induction.

Now let $p$ be an odd prime. As in Chapter 3, it suffices to show that there is an isomorphism

$$
\phi: K_{n-1}^{M W}\left(\mathbf{F}_{p}\right) \xrightarrow{\cong} L_{p}^{n} / L_{q}^{n}
$$

where $q$ is the greatest prime less than $p$.
We first establish the map $\phi$. Let $\mu(x)$ denote the endomorphism of the quotient of graded Z-modules $L_{p} / L_{q}$ defined by multiplication from the left by an element $x \in L_{p}$. Thus, if for example $x=\epsilon[u]$, then $\mu(\epsilon[u])$ is the degree- 1 homomorphism

$$
\begin{aligned}
\mu(\epsilon[u]): L_{p} / L_{q} & \longrightarrow L_{p} / L_{q} \\
\eta^{m}[p]\left[u_{1}\right] \cdots\left[u_{n+m}\right] & \longmapsto \epsilon[u] \eta^{m}[p]\left[u_{1}\right] \cdots\left[u_{n+m}\right]=\eta^{m}[p][u]\left[u_{1}\right] \cdots\left[u_{n+m}\right]
\end{aligned}
$$

Let

$$
\mathscr{E}_{p}:=\operatorname{Hom}\left(L_{p} / L_{q}, L_{p} / L_{q}\right)
$$

denote the graded associative ring of graded endomorphisms of $L_{p} / L_{q}$. We shall prove that the maps

$$
\bar{u} \mapsto \mu(\epsilon[u]): \mathbf{F}_{p}^{\times} \longrightarrow\left(\mathscr{E}_{p}\right)_{1}
$$

(where $|u|<p$ ) along with $\mu(\eta) \in\left(\mathscr{E}_{p}\right)_{-1}$ satisfy the relations of Milnor-Witt $K$-theory.

1. For the Steinberg relation, take $u \not \equiv 1(\bmod p)$. Then

$$
\mu(\epsilon[u]) \circ \mu(\epsilon[1-u])=\mu\left(\epsilon^{2}[u][1-u]\right)=0 .
$$

2. Suppose $u$ and $v$ are units modulo $p$, with $u v \equiv w(\bmod p)$. By Lemma 5.26 we then have

$$
\begin{aligned}
\mu(\epsilon[w])\left(\eta^{m}[p]\left[u_{1}\right] \cdots\left[u_{n+m}\right]\right) & =\eta^{m}[p][w]\left[u_{1}\right] \cdots\left[u_{n+m}\right] \\
& =\eta^{m}[p][u v]\left[u_{1}\right] \cdots\left[u_{n+m}\right]
\end{aligned}
$$

in $L_{p} / L_{q}$. Writing out $[u v]=[u]+[v]+\eta[u][v]$ yields

$$
\mu(\epsilon[u v])=\mu(\epsilon[u])+\mu(\epsilon[v])+\mu(\eta) \circ \mu(\epsilon[u]) \circ \mu(\epsilon[v]) .
$$

The last two relations $\mu(\eta) \circ \mu([u])=\mu([u]) \circ \mu(\eta)$ and $\mu(\eta) \circ \mu(1+\langle-1\rangle)=0$ follows similarly from the relations in $L_{p} / L_{q}$.

Thus we obtain a graded homomorphism $f: K_{*}^{M W}\left(\mathbf{F}_{p}\right) \rightarrow \mathscr{E}_{p}$, with $[\bar{u}] \mapsto \mu(\epsilon[u])$. Under the adjunction

$$
\operatorname{Hom}\left(K_{*}^{M W}\left(\mathbf{F}_{p}\right), \mathscr{E}_{p}\right) \cong \operatorname{Hom}\left(K_{*}^{M W}\left(\mathbf{F}_{p}\right) \otimes L_{p} / L_{q}, L_{p} / L_{q}\right)
$$

the morphism $f \in \operatorname{Hom}\left(K_{*}^{M W}\left(\mathbf{F}_{p}\right), \mathscr{E}_{p}\right)$ corresponds to the map

$$
\begin{aligned}
\tilde{f}: K_{*}^{M W}\left(\mathbf{F}_{p}\right) \otimes L_{p} / L_{q} & \longrightarrow L_{p} / L_{q} \\
{[\bar{u}] \otimes y } & \longmapsto f([\bar{u}])(y)=\mu(\epsilon[u])(y)=\epsilon[u] y .
\end{aligned}
$$

By restricting $\tilde{f}$ to $K_{*}^{M W}\left(\mathbf{F}_{p}\right) \otimes[p]$ we establish the desired map

$$
\begin{aligned}
\phi: K_{n-1}^{M W}\left(\mathbf{F}_{p}\right) & \longrightarrow L_{p}^{n} / L_{q}^{n} \\
\eta^{m}\left[\bar{u}_{1}\right] \cdots\left[\bar{u}_{n+m-1}\right] & \longmapsto \eta^{m}[p]\left[u_{1}\right] \cdots\left[u_{n+m-1}\right] .
\end{aligned}
$$

By Lemma 5.25, $\phi$ is surjective. By Theorem 5.17, $\partial_{p}$ vanishes on $L_{q}^{n}$ and thus induces a map $\partial_{p}: L_{p}^{n} / L_{q}^{n} \rightarrow K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)$. Furthermore, the compostition

$$
K_{n-1}^{M W}\left(\mathbf{F}_{p}\right) \xrightarrow{\phi} L_{p}^{n} / L_{q}^{n} \xrightarrow{\partial_{p}} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)
$$

is the identity; indeed,

$$
\partial_{p}\left(\phi\left(\left[\bar{u}_{1}\right] \cdots\left[\bar{u}_{n-1}\right]\right)\right)=\partial_{p}\left([p]\left[u_{1}\right] \cdots\left[u_{n-1}\right]\right)=\left[\bar{u}_{1}\right] \cdots\left[\bar{u}_{n-1}\right]
$$

by Theorem 5.17, By Lemma $5.13, K_{n-1}^{M W}\left(\mathbf{F}_{p}\right)$ is generated by the elements $\left[\bar{u}_{1}\right] \cdots\left[\bar{u}_{n-1}\right]$. Thus $\phi$ is an isomorphism.

Example 5.28. In the case $n=2$, note that $2 k[-1]^{2} \in \Lambda_{\infty}^{2}$ is mapped to

$$
2 k\{-1,-1\}=k\left\{(-1)^{2},-1\right\}=0 \in K_{2}^{M}(\mathbf{Q})
$$

under the canonical map $K_{2}^{M W}(\mathbf{Q}) \rightarrow K_{2}^{M}(\mathbf{Q})$. On the other hand, if $k$ is odd, then

$$
k[-1]^{2} \mapsto\left\{(-1)^{k},-1\right\}=\{-1,-1\} \neq 0 \in K_{2}^{M}(\mathbf{Q})
$$

Hence the kernel of the restriction of the canonical map $K_{2}^{M W}(\mathbf{Q}) \rightarrow K_{2}^{M}(\mathbf{Q})$ to $\Lambda_{\infty}^{2}$ is $2 \Lambda_{\infty}^{2}$.
By [MH73, Corollary 2.5, p.90], $I^{3}(\mathbf{Q})$ is the free additive group generated by $8\langle 1\rangle$. Since

$$
K_{2}^{M W}(\mathbf{Q}) \ni 2 k[-1]^{2} \longmapsto k \cdot 8\langle 1\rangle \in I^{3}(\mathbf{Q})
$$

the canonical map $K_{2}^{M W}(\mathbf{Q}) \rightarrow I^{2}(\mathbf{Q})$ sends $2 \Lambda_{\infty}^{2}$ onto $I^{3}(\mathbf{Q})=\mathbf{Z} \cdot 8\langle 1\rangle$. Moreover, for each $n \geq 0$ there is an inclusion

$$
I^{n+1}(F) \longrightarrow K_{n}^{M W}(F)
$$

given by

$$
\left\langle\left\langle x_{1}, \ldots, x_{n+1}\right\rangle\right\rangle \longmapsto \eta\left[x_{1}\right] \cdots\left[x_{n+1}\right] .
$$

The element

$$
k \cdot 8\langle 1\rangle=k\langle\langle-1,-1,-1\rangle\rangle \in I^{3}(\mathbf{Q})
$$

will thus be mapped to

$$
k \eta[-1]^{3}=-2 k[-1]^{2} \in 2 \Lambda_{\infty}^{2}
$$

We may thus identify $I^{3}(\mathbf{Q})$ with the kernel of the map $K_{2}^{M W}(\mathbf{Q}) \rightarrow K_{2}^{M}(\mathbf{Q})$. This illustrates the exact sequence

$$
0 \longrightarrow I^{3}(\mathbf{Q}) \longrightarrow K_{2}^{M W}(\mathbf{Q}) \longrightarrow K_{2}^{M}(\mathbf{Q}) \longrightarrow 0
$$

that Theorem 5.8 predicts.

## 5. Symbols in Motivic Homotopy Theory

Example 5.29. Theorem 5.27 yields $K_{n}^{M W}(\mathbf{Q}) \cong \mathbf{Z}$ for any $n \geq 3$. On the other hand, for any $n \geq 3$ we have $K_{n}^{M}(\mathbf{Q}) \cong \mathbf{Z} / 2$ by Mil70, Example 1.8], and $I^{n}(\mathbf{Q}) \cong \mathbf{Z}$ by MH73. Theorem 4.5 , p.75]. The exact sequence

$$
0 \longrightarrow I^{n+1}(\mathbf{Q}) \longrightarrow K_{n}^{M W}(\mathbf{Q}) \longrightarrow K_{n}^{M W}(\mathbf{Q}) \longrightarrow 0
$$

takes the form $0 \rightarrow \mathbf{Z} \xrightarrow{\cdot 2} \mathbf{Z} \rightarrow \mathbf{Z} / 2 \rightarrow 0$ whenever $n \geq 3$.
Example 5.30. In light of Definition 5.19, we can define $K_{n}^{M W}(\mathbf{Z}):=K_{n}^{M W}$ (Spec Z). Theorem 5.27 then shows that $K_{n}^{M W}(\mathbf{Z})=\overline{\mathbf{Z}}$ for $n \geq 2$, and also that the sequence

$$
0 \longrightarrow K_{n}^{M W}(\mathbf{Z}) \longrightarrow K_{n}^{M W}(\mathbf{Q}) \xrightarrow{\partial} \bigoplus_{p \geq 2} K_{n-1}^{M W}\left(\mathbf{F}_{p}\right) \longrightarrow 0
$$

is split exact. We do not expect this split exactness to hold for general number fields; indeed, it is not true in general for the corresponding sequence in neither Milnor $K$-theory nor Witt theory.

### 5.5 Milnor-Witt $K$-theory of $\mathbf{Q}(\sqrt{-1})$

Let us consider briefly the case of the Gaussian numbers. By [MH73, p.81], $I^{3}(F)=0$ for any totally imaginary number field $F$. This simplifies matters greatly, because then Theorem 5.8 implies that the map $K_{n}^{M W}(F) \rightarrow K_{n}^{M}(F)$ is an isomorphism in degrees $n \geq 2$. However, it is possible to obtain the isomorphism directly by the same method as we have seen. Let us illustrate the case $n=2$ :

Definition 5.31. Define a filtration

$$
L_{1} \subseteq \cdots \subseteq L_{n} \subseteq \cdots \subseteq K_{2}^{M W}(\mathbf{Q}(i))
$$

of $K_{2}^{M W} \mathbf{Q}(i)$ by letting

$$
L_{n}:=\left\langle\eta,[\alpha]: \alpha \in \mathbf{Z}[i]_{S_{n}}^{\times}\right\rangle \cap K_{2}^{M W}(\mathbf{Q}(i)),
$$

where the $S_{n}$ are defined as in Definition 3.6 of Chapter 3

$$
\text { As before, we have } \lim _{\longrightarrow} L_{n}=K_{2}^{M W} \mathbf{Q}(i) .
$$

Lemma 5.32. The group $L_{1}$ is trivial.
Proof. By Lemma 5.25 , it is enough to show that the elements $[u]\left[u^{\prime}\right]$, for $u, u^{\prime} \in \mathbf{Z}[i]_{S_{1}}^{\times}$, are trivial. Note that $[i][-i]=0$ by Lemma 5.10 .

- By Item 5 of Lemma 5.11, we have:

$$
[i][i]=[-1][i]=\left[i^{2}\right][i]=[i][i]-[i][i]=0 .
$$

Moreover,

$$
[-1][-1]=\left[i^{2}\right][-1]=[i][-1]-[-1][i]=[i][i]-[i][i]=0 .
$$

- Note that $[1+i][-i]=0=[1-i][i]$ by the Steinberg relation. Using this we have:

$$
[1-i]=[(1+i)(-i)]=[1+i]+[-i]+\eta[1+i][-i]=[1+i]+[-i]
$$

Multiplying by $[i]$ from the right we obtain

$$
0=[1-i][i]=[1+i][i]+[-i][i]
$$

Since $[-i][i]=0$, this yields $[1+i][i]=0$.

- The second relation of Milnor-Witt $K$-theory now gives $[ \pm 1 \pm i][u]=0$ for any $u \in \mu_{4}$; indeed, if $u^{\prime} \in \mu_{4}$ we have

$$
\left[u^{\prime}(1+i)\right][u]=\left[u^{\prime}\right][u]+[1+i][u]+\eta\left[u^{\prime}\right][1+i][u]=0 .
$$

- By Item 3 of Lemma 5.11, $[u]\left[(1+i)^{-1}\right]=[1+i][u]=0$ for any $u \in\{ \pm 1, \pm i, 1+i\}$.

By the cases considered above, Lemma 5.14 assures that $[1+i][ \pm 1 \pm i]=0$. Using the $\epsilon$ commutativity of Milnor-Witt $K$-groups we establish also the symmetric counterparts of all cases above. Finally, we consider the general case $\left[u(1+i)^{n}\right]\left[u^{\prime}(1+i)^{m}\right]$ for $n, m \in \mathbf{Z}, u, u^{\prime} \in \mu_{4}$. By the second relation of Milnor-Witt $K$-theory along with Lemma 5.15 and Lemma 5.16 we have

$$
\begin{aligned}
{\left[u(1+i)^{n}\right]\left[u^{\prime}(1+i)^{m}\right]=} & \left([u]+n_{\epsilon}[1+i]+\eta[u] n_{\epsilon}[1+i]\right)\left(\left[u^{\prime}\right]+m_{\epsilon}[1+i]+\eta\left[u^{\prime}\right] m_{\epsilon}[1+i]\right) \\
= & {[u]\left[u^{\prime}\right]+m_{\epsilon}[u][1+i]+\eta n_{\epsilon}[u]\left[u^{\prime}\right][1+i] } \\
& +n_{\epsilon}[1+i]\left[u^{\prime}\right]+n_{\epsilon} m_{\epsilon}[1+i][1+i]+\eta n_{\epsilon} m_{\epsilon}[1+i]\left[u^{\prime}\right][1+i] \\
& +\eta n_{\epsilon}[u][1+i]\left[u^{\prime}\right]+\eta m_{\epsilon}[u][1+i]+\eta^{2} n_{\epsilon} m_{\epsilon}[u][1+i]\left[u^{\prime}\right][1+i] .
\end{aligned}
$$

In the last expression, all terms vanish by the cases considered above, so this concludes the proof.

From here on we proceed identically as in Theorem 5.27 to obtain the isomorphism:
Theorem 5.33. For each finite place $v$ of $\mathbf{Q}(i)$, choose a uniformizer $\pi_{v}$. Then the map $\partial:=\bigoplus_{v} \partial_{v}^{\pi_{v}}$ induces an isomorphism

$$
K_{2}^{M W}(\mathbf{Q}(i)) \cong \bigoplus_{v \in \mathrm{P}_{\mathbf{Q}(i)}^{\mathrm{nc}}} K_{1}^{M W}(k(v)) \cong \bigoplus_{v \in \mathrm{P}_{\mathbf{Q}_{(i)}}^{\mathrm{nc}}} K_{1}^{M}(k(v))
$$

## Final Remarks

The theory of motivic homotopy is quite new, and the discovery of Milnor-Witt $K$-theory is even newer. Thus there are yet many unanswered questions. One possibility for further studies is to try to see if there exist analogs of Hilbert symbols in Milnor-Witt $K$-theory-at least for $K_{2}^{M W}$. Recall that in algebraic $K$-theory we have Moore's theorem, stating that there is an exact sequence

$$
K_{2}(F) \xrightarrow{h} \bigoplus_{v \in \mathrm{Pl}^{\mathrm{nc}}} \mu\left(F_{v}\right) \xrightarrow{\pi} \mu(F) \longrightarrow 1
$$

Here the map $h$ is induced by the Hilbert symbols defined on the number field $F$. As we have seen in Chapter 2, the above exact sequence should be considered as a uniqueness theorem for Hilbert reciprocity.

On the other hand, in Lam05, p.159], Lam defines Hilbert symbols in Witt theory and uses these to show that there is an exact sequence

$$
0 \longrightarrow I^{2}(\mathbf{Q}) / I^{3}(\mathbf{Q}) \longrightarrow \bigoplus_{v \in \mathrm{Pl}_{\mathbf{Q}}} I^{2}\left(\mathbf{Q}_{v}\right) / I^{3}\left(\mathbf{Q}_{v}\right) \longrightarrow \mathbf{Z} / 2 \longrightarrow 0
$$

which is interpreted as a uniqueness theorem for Hilbert reciprocity in Witt theory. This result was originally proved for any global field $F$ by Bass and Tate in the equivalent form Mil70, Lemma A.1]

$$
0 \longrightarrow k_{2}^{M}(F) \longrightarrow \bigoplus_{v \in \mathrm{Pl}_{F}^{\text {nc }}} k_{2}^{M}\left(F_{v}\right) \longrightarrow \mathbf{Z} / 2 \longrightarrow 0
$$

(remember that $I^{2}(F) / I^{3}(F) \cong K_{2}^{M}(F) / 2 K_{2}^{M}(F)=k_{2}^{M}(F)$ ). As there is a uniqueness of reciprocity laws-theorem in both Milnor $K$-theory and Witt theory, this suggests that there should also be a corresponding result in Milnor-Witt $K$-theory that specializes to the two mentioned sequences. The first step toward such a result would be to define Hilbert symbols on Milnor-Witt $K$-theory.

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[^0]:    ${ }^{1}$ The class group of $\mathbf{Q}\left(\zeta_{p}\right)$ is trivial for $2 \leq p \leq 19$; the first instance of a cyclotomic field with nontrivial class group is $\mathbf{Q}\left(\zeta_{23}\right)$, where $\operatorname{Cl}\left(\mathbf{Q}\left(\zeta_{23}\right)\right) \cong \mathbf{Z} / 3$ Was82, p.353]. As $p$ increases, the class number of $\mathbf{Q}\left(\zeta_{p}\right)$ grows exponentially Was82, Theorem 4.20].

[^1]:    ${ }^{2}$ To solve congruences of Gaussian integers it can often be helpful to draw the lattice spanned by unit multiples of the given Gaussian prime, and look at Gaussian integers relative to this lattice; see Figure 3.1 in Chapter 3

[^2]:    ${ }^{1}$ E.g., taking $X=\mathbf{P}_{\mathbf{C}}^{1}$, the exact sequence $0 \rightarrow \mathcal{O}_{X}(-2) \rightarrow \mathcal{O}_{X}(-1)^{\oplus 2} \rightarrow \mathcal{O}_{X} \rightarrow 0$ does not split.

