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Argos seminar on intersections of modular correspondences
Astérisque, tome 312 (2007)
[http://www.numdam.org/item?id=AST_2007__312__R1_0](http://www.numdam.org/item?id=AST_2007__312__R1_0)
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## ARGOS SEMINAR <br> ON

 INTERSECTIONS OF MODULAR CORRESPONDENCES2000 Mathematics Subject Classification. - 11G18, 11E08, 11F03, 11F30, 11F32, 11G15, 11G18, 14B12, 14G35, 14K07, 14K22, 14L05.
Key words and phrases. - Modular correspondences, endomorphisms of formal groups, quadratic forms over the ring of $p$-adic integers, Siegel-Eisenstein series.

## ARGOS SEMINAR <br> ON <br> INTERSECTIONS OF MODULAR CORRESPONDENCES


#### Abstract

This volume contains the written account of the Bonn seminar on arithmetic geometry 2003/2004. It gives a coherent exposition of the theory of intersections of modular correspondences. The focus of the seminar is the formula for the intersection number of arithmetic modular correspondences due to Gross and Keating. Other topics treated are Hurwitz's theorem on the intersection of modular correspondences over the field of complex numbers, and the relation of the arithmetic intersection numbers to Fourier coefficients of Siegel-Eisenstein series.

Also included is background material on one-dimensional formal groups and their endomorphisms, and on quadratic forms over the ring of $p$-adic integers.


## Résumé (Séminaire ARGOS sur les intersections de correspondances modulaires)

Ce volume consiste des exposés faits dans le cadre du séminaire de géométrie arithmétique de Bonn en 2003/2004. Il donne une exposition systématique de la théorie des intersections de correspondances modulaires. Le but principal est la formule de Gross-Keating du nombre d'intersection de correspondances modulaires arithmétiques. Autres sujets traités sont le théorème de Hurwitz sur l'intersection de correspondances modulaires sur le corps des nombres complexes, et la relation des nombres d'intersection arithmétiques aux coefficients de Fourier des séries de Siegel-Eisenstein.

On a aussi inclus des rappels sur les groupes formels à un paramètre et leurs endomorphismes, et sur les formes quadratiques sur l'anneau des entiers $p$-adiques.

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## 1. FOREWORD

## 1. Motivation and main results

This book is based on the notes for the $\operatorname{AP\Gamma O} \Sigma^{(1)}$ seminar of the winter semester $2003 / 2004$ at the University of Bonn. Its aim was to go through the paper On the intersection of modular correspondences by Gross and Keating [GK], and understand it thoroughly. This subject was chosen for three reasons. First of all, it was felt that the mathematics contained in this paper (and the papers on which Gross and Keating base their article) is extremely interesting, and has become even more important recently, due to the use that $S$. Kudla and others have made of these results. Secondly, thanks to the elementary methods employed in the proofs of the main theorems, the seminar provided a rapid access, even to a novice in the field, to a deep and sophisticated topic in arithmetic algebraic geometry. Thirdly, it was felt from the start that the literature on the subject was not easy to penetrate and that therefore the effort made by all speakers to master this material should not be lost, and that a written account of the seminar should be made available.

The origin of the topics treated in the seminar goes back to the $19^{\text {th }}$ century. Let $j=j(\tau)$ be the elliptic modular function on the upper half plane. For $m \geq 1$ let $\varphi_{m}\left(j, j^{\prime}\right) \in \mathbb{Z}\left[j, j^{\prime}\right]$ be the classical modular polynomial, defined by

$$
\begin{equation*}
\varphi_{m}\left(j(\tau), j\left(\tau^{\prime}\right)\right)=\prod_{\substack{A \in \mathrm{M}_{2}(\mathbb{Z}) \\ \operatorname{det}(A)=m \\ \bmod S L_{2}(\mathbb{Z})}}\left(j(\tau)-j\left(A \tau^{\prime}\right)\right) \tag{1.1}
\end{equation*}
$$

Kronecker and Hurwitz established a number of important properties of these polynomials, as for instance their factorization into irreducible factors. They also proved degree formulas like

$$
\begin{equation*}
\operatorname{deg} f_{m}=\sum_{d d^{\prime}=m} \max \left(d, d^{\prime}\right) \tag{1.2}
\end{equation*}
$$

where $f_{m}(j)=\varphi_{m}(j, j)$, for $m$ not a square.

[^0]From the point of view of the seminar, the interest in these results lies in the fact that they can be interpreted as giving intersection numbers on the complex surface $S_{\mathbb{C}}=\operatorname{Spec} \mathbb{C}\left[j, j^{\prime}\right]$. Let $T_{m, \mathbb{C}} \subseteq S_{\mathbb{C}}$ be the divisor defined by $\varphi_{m}=0$. Then (1.2) can be interpreted as the intersection formula

$$
\begin{equation*}
\left(T_{m, \mathbb{C}} \cdot T_{1, \mathbb{C}}\right)=\sum_{d d^{\prime}=m} \max \left(d, d^{\prime}\right) \tag{1.3}
\end{equation*}
$$

if $m$ is not a square. Here $\left(T_{m_{1}, \mathbb{C}} \cdot T_{m_{2}, \mathbb{C}}\right)$ is defined by

$$
\begin{equation*}
\left(T_{m_{1}, \mathbb{C}} \cdot T_{m_{2}, \mathbb{C}}\right)=\operatorname{dim}_{\mathbb{C}} \mathbb{C}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}\right) \tag{1.4}
\end{equation*}
$$

More generally, Hurwitz showed that the divisors $T_{m_{1}, \mathbb{C}}$ and $T_{m_{2}, \mathbb{C}}$ intersect properly on $S_{\mathbb{C}}$ if and only if $m_{1} m_{2}$ is not a perfect square and gave an explicit expression for the intersection number $\left(T_{m_{1}, \mathbb{C}} \cdot T_{m_{2}, \mathbb{C}}\right)$. This in turn leads to the famous class number relations of Kronecker and Hurwitz.

Gross and Keating took up this classical subject by adding an arithmetic dimension to it. Instead of usual intersection numbers they consider arithmetic intersection numbers. Let $S=\operatorname{Spec} \mathbb{Z}\left[j, j^{\prime}\right]$, which we consider as an arithmetic threefold. Let $T_{m}$ be the arithmetic divisor defined by $\varphi_{m}=0$. The arithmetic intersection number is defined for any triple of positive integers $m_{1}, m_{2}, m_{3}$ by

$$
\begin{equation*}
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)=\log \# \mathbb{Z}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right) \tag{1.5}
\end{equation*}
$$

Gross and Keating derive a criterion for when this number is finite and give in this case an explicit expression for it (see below). This result is the main focus of the present book. Let us state it from the point of view adopted in these notes. Let $\mathcal{M}$ be the moduli space of elliptic curves over Spec $\mathbb{Z}$ (since we impose no level structure, $\mathcal{M}$ is not a scheme, but a Deligne-Mumford stack). Put $\mathcal{S}=\mathcal{M} \times_{\text {Spec } \mathbb{Z}} \mathcal{M}$. For a positive integer $m$, let $\mathcal{T}_{m}$ be the moduli space of isogenies of elliptic curves $E \rightarrow E^{\prime}$ of degree $m$. Then $\mathcal{I}_{m}$ maps by a finite unramified morphism to $\mathcal{M} \times \mathcal{M}$. From this point of view, the intersection number above should be interpreted as

$$
\sum_{p} \log (p) \cdot \sum_{x \in \mathcal{X}\left(\mathbb{F}_{p}\right)} \frac{1}{\# \operatorname{Aut}(x)} \lg \widehat{\mathcal{O}}_{\mathcal{X}, x},
$$

where we denote by $\mathcal{X}$ the triple fiber product of $\mathcal{T}_{m_{1}}, \mathcal{T}_{m_{2}}$, and $\mathcal{T}_{m_{3}}$ over $\mathcal{M} \times \mathcal{M}$. Here the weighting factor $\frac{1}{\# \operatorname{Aut}(x)}$ is due to the fact that $\mathcal{X}$ is a stack.

We now state the main results contained in this volume.
We denote by $S_{\mathbb{C}}$ resp. by $T_{m, \mathbb{C}}$ the base change of $S$ resp. $T_{m}$ to Spec $\mathbb{C}$. The first result is Hurwitz's theorem.

Theorem 1.1. - The cycles $T_{m_{1}, \mathbb{C}}$ and $T_{m_{2}, \mathbb{C}}$ intersect properly on $S_{\mathbb{C}}$ if and only if the integer $m=m_{1} m_{2}$ is not a perfect square. In this case, the intersection $T_{m_{1}, \mathbb{C}} \times{ }_{S_{\mathbb{C}}} T_{m_{2}, \mathbb{C}}$ lies over the locus in $S_{\mathbb{C}}$ corresponding to pairs ( $E, E^{\prime}$ ) of elliptic curves with complex multiplication by orders in the imaginary-quadratic field $\mathbb{Q}(\sqrt{-m})$ of discriminant $\geq-4 m$. The intersection number is equal to

$$
\left(T_{m_{1}, \mathbb{C}} \cdot T_{m_{2}, \mathbb{C}}\right)=\sum_{\substack{t \in \mathbb{Z} \\ t^{2}<4 m}} \sum_{d \mid g c d\left(m_{1}, m_{2}, t\right)} d \cdot H\left(\frac{4 m-t^{2}}{d^{2}}\right)
$$

Here $H(n)$ denotes the Hurwitz class number (the number of $\mathrm{SL}_{2}(\mathbb{Z})$-equivalence classes of positive definite binary quadratic forms over $\mathbb{Z}$ with determinant $n$ ).

The second result is the theorem of Gross and Keating.
Theorem 1.2. - The cycles $T_{m_{1}}, T_{m_{2}}$ and $T_{m_{2}}$ intersect properly on $S$ if and only if there is no positive definite binary quadratic form over $\mathbb{Z}$ which represents the three integers $m_{1}, m_{2}, m_{3}$. In this case the intersection $T_{m_{1}} \times_{S} T_{m_{2}} \times_{S} T_{m_{3}}$ lies over the locus in $S$ corresponding to pairs $\left(E, E^{\prime}\right)$ of elliptic curves which are supersingular in some characteristic $p$ with $p<4 m_{1} m_{2} m_{3}$. The arithmetic intersection number is equal to

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)=\sum_{p<4 m_{1} m_{2} m_{3}} n(p) \log p
$$

where

$$
n(p)=\frac{1}{2} \cdot \sum_{Q}\left(\prod_{\substack{\ell \mid \Delta, \ell \neq p}} \beta_{\ell}(Q)\right) \cdot \alpha_{p}(Q)
$$

Here the sum is the taken over all positive definite integral ternary quadratic forms $Q$ with diagonal $\left(m_{1}, m_{2}, m_{3}\right)$ which are isotropic over $\mathbb{Q}_{\ell}$ for all $\ell \neq p$. Furthermore $\Delta=\frac{1}{2} \operatorname{det} Q$ and $\beta_{\ell}(Q)$ is a normalized representation density of $Q$ by the $\mathbb{Z}_{\ell}$-lattice $\mathrm{M}_{2}\left(\mathbb{Z}_{\ell}\right)$ with its norm form. Finally, $\alpha_{p}(Q)$ is the length of a certain local deformation space. Namely, one considers the universal deformation space of a triple of isogenies of formal groups of dimension 1 and height 2 over $\overline{\mathbb{F}}_{p}$. Here the passage from a global problem involving elliptic curves to a local problem on formal groups is provided by the Serre-Tate theorem. In [GK], Gross and Keating give completely explicit expressions for the factors $\beta_{\ell}(Q)$ and $\alpha_{p}(Q)$, comp. Chapters 5 and 13. They express these quantities in terms of new invariants of ternary quadratic forms over $\mathbb{Z}_{p}$ which are defined by them for this purpose (the Gross-Keating invariant in $\left(\mathbb{Z}_{\geq 0}\right)^{3}$ and the Gross-Keating epsilon factor in $\{ \pm 1\})$. This is especially striking in the cases when $\ell=2$ resp. $p=2$, in the other cases these invariants can be expressed in terms of classical quantities.

The invariant $\alpha_{p}(Q)$ is probably the most interesting ingredient in the formula above, and we now give a precise definition.

Let $G$ be a formal group of dimension 1 and height 2 over $\overline{\mathbb{F}}_{p}$. Let $W=W\left(\overline{\mathbb{F}}_{p}\right)$ be the ring of Witt vectors. The universal deformation of the pairs $(G, G)$ is then $\left(\Gamma, \Gamma^{\prime}\right)$ over the formal scheme $\hat{\mathcal{S}}=\operatorname{Spf} W\left[\left[t, t^{\prime}\right]\right]$. If now $f_{1}, f_{2}, f_{3}: G \rightarrow G$ are three endomorphisms $\neq 0$, we let $I_{i} \subset W\left[\left[t, t^{\prime}\right]\right]$ for $i=1,2,3$ be the minimum ideal such that $f_{i}$ lifts to a homomorphism $\tilde{f}_{i}: \Gamma \rightarrow \Gamma^{\prime}\left(\bmod I_{i}\right)$. Then $I_{i}$ defines a divisor $\hat{\mathcal{T}}_{i}$ on $\hat{\mathcal{S}}$. Consider

$$
\begin{equation*}
\left(\hat{\mathcal{T}}_{1} \cdot \hat{\mathcal{T}}_{2} \cdot \hat{\mathcal{T}}_{3}\right)=\operatorname{length}_{W} W\left[\left[t, t^{\prime}\right]\right] /\left(I_{1}+I_{2}+I_{3}\right) \tag{1.6}
\end{equation*}
$$

On $\operatorname{End}(G)$ we have the usual quadratic form Nm with values in $\mathbb{Z}_{p}$ (the norm form, after identifying $\operatorname{End}(G)$ with the maximal order in the quaternion division algebra over $\left.\mathbb{Q}_{p}\right)$. It turns out that (1.6) only depends on the quadratic form $Q\left(f_{1}, f_{2}, f_{3}\right)$ :
$(x, y, z) \mapsto \operatorname{Nm}\left(x f_{1}+y f_{2}+z f_{3}\right)$, and even only on its $\mathrm{GL}_{3}\left(\mathbb{Z}_{p}\right)$-equivalence class. We then set

$$
\begin{equation*}
\alpha_{p}(Q)=\left(\hat{\mathcal{T}}_{1} \cdot \hat{\mathcal{T}}_{2} \cdot \hat{\mathcal{T}}_{3}\right) \tag{1.7}
\end{equation*}
$$

for any triple $f_{1}, f_{2}, f_{3}$ with $Q\left(f_{1}, f_{2}, f_{3}\right)=Q$. The formula for $\alpha_{p}(Q)$ in terms of the Gross-Keating invariant $\left(a_{1}, a_{2}, a_{3}\right) \in\left(\mathbb{Z}_{\geq 0}\right)^{3}$ with $a_{1} \leq a_{2} \leq a_{3}$ is as follows. We note that for $p \neq 2$, in which case $Q$ can be diagonalized, the integers $a_{1}, a_{2}, a_{3}$ are simply the $p$-adic valuations of the diagonal entries.

$$
\begin{aligned}
\alpha_{p}(Q)= & \sum_{i=0}^{a_{1}-1}(i+1)\left(a_{1}+a_{2}+a_{3}-3 i\right) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-2\right) / 2}\left(a_{1}+1\right)\left(2 a_{1}+a_{2}+a_{3}-4 i\right) p^{i} \\
& +\frac{a_{1}+1}{2}\left(a_{3}-a_{2}+1\right) p^{\left(a_{1}+a_{2}\right) / 2}, \text { if } a_{1} \equiv a_{2}(\bmod 2) \\
\alpha_{p}(Q)= & \sum_{i=0}^{a_{1}-1}(i+1)\left(a_{1}+a_{2}+a_{3}-3 i\right) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-1\right) / 2}\left(a_{1}+1\right)\left(2 a_{1}+a_{2}+a_{3}-4 i\right) p^{i}, \\
& \text { if } a_{1} \not \equiv a_{2}(\bmod 2)
\end{aligned}
$$

The above results are the main focus of these notes. In the last chapter we reformulate Theorem 1.2 as a relation between the arithmetic intersection numbers and the Fourier coefficients of special values of derivatives of Siegel Eisenstein series, along the lines sketched in the introduction to $[\mathbf{G K}]$. The idea that this can be done is attributed there to S. Kudla and D. Zagier; in the intervening years Kudla and others have gone a long way towards proving such relations in much greater generality [Ku2, Ku3]. Let

$$
\begin{equation*}
E(\tau, s)=\sum \operatorname{det}(\mathbf{c} \tau+\mathbf{d})^{-2} \cdot \frac{\operatorname{det}(y)^{\frac{s}{2}}}{|\operatorname{det}(\mathbf{c} \tau+\mathbf{d})|^{s}} \tag{1.8}
\end{equation*}
$$

be the classical Siegel Eisenstein series of genus 3 and weight 2 for the full modular group. Here $\tau=x+i y \in \mathfrak{H}_{3}$ and $s \in \mathbb{C}$ is a complex parameter with large real part, and the sum is over representatives $\left(\gamma=_{\mathbf{c}}^{*} \begin{array}{c}* \\ \mathbf{d}\end{array}\right)$ of the left cosets of the Siegel parabolic in $\mathrm{Sp}_{3}(\mathbb{Z})$. Then $E(\tau, s)$ has a meromorphic continuation to the entire $s$-plane and vanishes at $s=0$. The derivative $E^{\prime}(\tau, 0)=\frac{\partial E}{\partial s}(\tau, 0)$ is a non-holomorphic modular form of weight 2 for $\mathrm{Sp}_{3}(\mathbb{Z})$ and has a Fourier expansion

$$
\begin{equation*}
E^{\prime}(\tau, 0)=\sum_{T \in \operatorname{Sym}_{3}(\mathbb{Z})^{\vee}} c^{\prime}(T, y) \cdot q^{T} \tag{1.9}
\end{equation*}
$$

where $q^{T}=\exp (2 \pi i \operatorname{tr}(T \tau))$, for any half-integral matrix $T$. It turns out that for positive-definite $T$ the coefficient $c^{\prime}(T, y) \equiv c^{\prime}(T)$ is independent of $y=\operatorname{Im}(\tau)$.

Theorem 1.3. - Let $m_{1}, m_{2}, m_{3}$ be a triple of positive integers such that there is no positive definite binary quadratic form over $\mathbb{Z}$ which represents $m_{1}, m_{2}$ and $m_{3}$. There
exists a constant $\kappa$ independent of $m_{1}, m_{2}, m_{3}$ such that

$$
\left(\mathcal{T}_{m_{1}} \cdot \mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}\right)=\kappa \cdot \sum_{\substack{T \in \operatorname{Sym}_{3}(\mathbb{Z})_{>0}^{\vee} \\ \operatorname{diag}(T)=\left(m_{1}, m_{2}, m_{3}\right)}} c^{\prime}(T)
$$

It should be pointed out that the left hand side in Theorem 1.3 also breaks into geometrically defined terms indexed by quadratic forms $T$ and that the identity above holds termwise.

## 2. Content of this book

We now explain what is to be found in this book and what is not. In Chapter 2, G. Vogel gives a review of classical results on modular polynomials. It is followed by the brief Chapter 3 by U . Görtz, where a certain sum of representation numbers is computed. In Chapter 4, U. Görtz proves the first part of Theorem 1.2, which states under which conditions the three divisors intersect in dimension 0 , and it is explained how the formula for the intersection number follows from the results of the later chapters. T. Wedhorn then investigates in Chapter 5 the quadratic space obtained as the space of homomorphisms between two supersingular elliptic curves with the degree form.

Next we come to the local theory. We have made an effort to give all the necessary background on formal groups, their deformations and the deformations of their endomorphisms. Chapter 6 by V. Meusers gives a summary of Lubin-Tate theory for formal groups, cf. [LT1]. In Chapter 7, E. Viehmann and K. Ziegler give the construction of the formal moduli space of formal groups, and more generally, formal $A$-modules, following Drinfeld [D]. Chapter 8 by S. Wewers is devoted to Gross' theory of canonical and quasi-canonical lifts, cf. [G], and in Chapter 9 V. Meuser explains an analogous theory in the split case, expanding on a remark in [G].

Since we will be interested in lifting isogenies of elliptic curves rather than just the curves themselves, we need to understand how endomorphisms of formal groups can be lifted. This was analyzed in much detail by Keating, see [K2] (which is based on Keating's unpublished Harvard thesis $[\mathbf{K 1}]$ ). This theory is presented here in Chapter 10 by E. Viehmann and Chapter 11 by I. Vollaard. Another ingredient we need is the theory of quadratic forms over $\mathbb{Z}_{\ell}$, including the delicate case $\ell=2$ treated in section 4 of $[\mathbf{G K}]$, comp. also the account of Yang [Y1]. This is dealt with in Chapter 12 by I. Bouw. We come back to the theory of quadratic forms, namely to the computation of certain representation densities, in T. Wedhorn's Chapter 15, where, however, we have merely quoted from Kitaoka $[\mathrm{Ki}]$ and Katsurada $[\mathbf{K a}]$ the facts that we use.

We then come to the investigation of the invariants $\alpha_{p}(Q)$, and to the proof of the explicit formula for them. Here the case $p=2$ causes additional complications. We provide two different proofs of the formula in that case - one which relies on laborious explicit computations (Chapter 13, by M. Rapoport), and another one which is more conceptual (Chapter 14, by S. Wewers). We feel that both are enlightening in their own way.

Whereas we have provided a lot of background information on the ingredients of the results in [GK], with Theorem 1.3 we were less ambitious, contenting ourselves with references to the appropriate papers (mostly of Kudla) to calculate the Fourier coefficients of $E^{\prime}(\tau, 0)$, see Chapter 16 by M. Rapoport et T. Wedhorn.

## 3. Perspective

We believe that Gross' theory of canonical and quasi-canonical liftings is going to have even more applications than have been found so far. We hope that our book can serve as a basis of future research. At the end of this introduction there is a list of references of which we are aware, where this theory is used. The theory was invented originally by Gross in connection with the proof of the Gross-Zagier formula. We note that in B. Conrad's recent account of the geometric ingredients of this proof $[\mathbf{C o}]$ the theory of quasi-canonical liftings is explicitly excluded; therefore our notes may also be viewed as a complement to Conrad's exposition. Also, Chapter III of [KRY] is based on the present notes (in loc. cit., only the intersection numbers ( $\mathcal{T}_{1} \cdot \mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}$ ) are needed).

The main ingredient of the proof of Theorem 1.2 is the determination of the quantity $\alpha_{p}(Q)$. This may be viewed as a special case of the following general problem. Let $G$ and $G^{\prime}$ be two $p$-divisible groups over $\overline{\mathbb{F}}_{p}$. The universal deformation of $\left(G, G^{\prime}\right)$ is then $\left(\Gamma, \Gamma^{\prime}\right)$ over the formal spectrum of a power series ring $R$ over $W$. Let $f: G \rightarrow G^{\prime}$ be an isogeny. The problem is to determine the minimal ideal $I$ in $R$ such that $f$ lifts to an isogeny $\tilde{f}: \Gamma \rightarrow \Gamma^{\prime}(\bmod I)$. Related to this question is the following problem: Let $I$ be the minimal ideal such that a given set of isogenies $f_{1}, \ldots, f_{r}: G \rightarrow G^{\prime}$ lifts to a set of isogenies $\tilde{f}_{1}, \ldots, \tilde{f}_{r}: \Gamma \rightarrow \Gamma^{\prime}(\bmod I)$. The problem is to determine when $I$ is of finite colength in $R$, and if so, to determine this colength explicitly.

To the sophisticated reader, it may seem curious that old-fashioned power series methods are used here to solve these problems in the case of $p$-divisible formal groups of dimension 1 and height 2 . It is natural to ask whether more recent methods, like Grothendieck-Messing lifting theory, Cartier theory, or the theory of displays can be used to solve this problem. Indeed, as Zink [Z] has shown, the theory of displays can be used in some instances to prove results in this direction. However, so far these methods have not succeeded in obtaining the full statement. In view of the fact that the cases of finite colength in [KR1, KR2] all reduce to the Gross-Keating problem, it is conceivable that the special case of the general problem studied here is the only one where a reasonable uniform answer exists. This might also explain the relative failure of the more generally applicable methods.

## 4. Acknowledgements

We thank all participants of the seminar for their interest and their pertinent questions, especially Ch. Kaiser and E. Lau. We also thank S. Kudla for his encouragement, for his many explanations and for his thorough reading of Chapter 16.

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## NOTATION

We assemble some of the notation which is used more or less systematically throughout the book.

For an integer $m>0$, we denote by $\varphi_{m} \in \mathbb{Z}[X, Y]$ the modular polynomial (see $[\mathbf{V g}]$ ). We write $S=\operatorname{Spec} \mathbb{Z}[X, Y]$, and denote by $T_{m} \subseteq S$ the divisor associated to $\varphi_{m}$. By $S_{\mathbb{C}}, T_{m, \mathbb{C}}$ etc. we denote the base change to $\mathbb{C}$, and by $\mathcal{S}, \mathcal{T}_{m}$ we denote the corresponding Deligne-Mumford stacks (see [Go2]).

In $[\mathbf{M e} 1],[\mathbf{V Z}],[\mathbf{W w} \mathbf{1}],[\mathbf{M e 2}],[\mathbf{V i}],[\mathbf{V l}]$ dealing with the local theory, the following situation is considered: $K$ is a field, complete with respect to a discrete valuation $v_{K}$, with ring of integers $\mathcal{O}_{K}$. We denote by $\mathfrak{p}$ the maximal ideal of $\mathcal{O}_{K}$, and by $\pi$ a uniformizer. The residue class field $\mathcal{O}_{K} / \mathfrak{p}$ of $K$ is assumed to be finite, of cardinality $q$, and $k$ is an extension field of $\mathcal{O}_{K} / \mathfrak{p}$, in fact in most cases it is an algebraic closure of the residue class field. Furthermore, $L$ is a quadratic extension of $K$, with ring of integers $\mathcal{O}_{L}$, and $M$ is the completion of the maximal unramified extension of $K$ (or in some places of $L$ ). By $D$ we denote 'the' quaternion division algebra over $K$. The maximal order of $D$ is denoted by $\mathcal{O}_{D}$, and $\Pi=\pi_{D}$ is a uniformizing element. In $[\mathbf{R}]$, [ $\mathbf{W} \mathbf{w} \mathbf{2}$ ] the special case where the base field is $\mathbb{Q}_{p}$ is considered, and the notation is slightly different: there $K / \mathbb{Q}_{p}$ is a quadratic extension (and $L$ denotes a quadratic space). Also, in $[\mathbf{W d 1}] D$ denotes a quaternion algebra over $\mathbb{Q}$.

The letters $F, G, H, \Gamma$ usually denote formal groups (or formal $\mathcal{O}_{K}$-modules etc.). Often, $G$ denotes the special fiber of a deformation $F$. Whereas mostly $F_{r}$ denotes a quasi-canonical lifting of level $r$ (and in particular $F_{0}$ denotes the canonical lift), see $[\mathbf{W w} \mathbf{w}]$, in $[\mathbf{V i}]$ and $[\mathbf{V l}] F_{n}$ denotes the base change $F \otimes_{k \llbracket t \rrbracket} k \llbracket t \rrbracket / t^{n}$ or $F \otimes_{A} A / \pi^{n}$.

If $R$ is a ring, and $(L, Q)$ is a quadratic space over $R$, i.e., a free $R$-module $L$ with a quadratic form $Q$, we associate to it the bilinear form $(x, y)=Q(x+y)-Q(x)-Q(y)$ and-after fixing a basis $\psi_{1}, \ldots, \psi_{n}$-the matrix $B=B(\boldsymbol{\psi})=\left(\left(\psi_{i}, \psi_{j}\right)\right)_{i, j}$ (in [B]) or the matrix $T=\left(\frac{1}{2}\left(\psi_{i}, \psi_{j}\right)\right)_{i, j}$ (in [R], [Wd2]). If $Q^{\prime}$ is another quadratic form, on $R^{m}$, say, then we denote by $R_{L}\left(Q^{\prime}\right)$ the representation number of $Q^{\prime}$ in $L$, i.e., the number of isometries from $\left(R^{m}, Q^{\prime}\right)$ to $(L, Q)$; see $[\mathbf{V g}],[\mathbf{G o 2}]$ and in particular [Wd2]. To a ternary quadratic form over $\mathbb{Z}_{\ell}$ we attach its Gross-Keating invariants $a_{1}, a_{2}, a_{3}$ and $\epsilon$, see $[\mathbf{B}]$. Finally, the numbers $\alpha_{p}(Q) \in \mathbb{Z}, \beta_{\ell}(Q) \in \mathbb{Z}$ which appear in the statement of the main theorem are defined in $[\mathbf{R}]$ and $[\mathbf{W d 2}]$, respectively.

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# 2. MODULAR POLYNOMIALS 

by

Gunther Vogel


#### Abstract

We introduce the classical modular polynomials and calculate (modulo the determination of a certain sum of representation numbers) the intersection number of two divisors defined by modular polynomials (Hurwitz's theorem).


Résumé (Polynômes modulaires). - On introduit les polynômes modulaires classiques et détermine (modulo le calcul d'une certaine somme de nombres de représentations) le nombre d'intersection de deux diviseurs définis par des polynômes modulaires (théorème de Hurwitz).

We introduce modular polynomials and prove some elementary properties. This is classical and well-known, see e.g. [L, §5]. In the second part, we compute the intersection numbers of the divisors defined by two modular polynomials in the 2dimensional complex plane. This computation, due to Gross and Keating ([GK]), re-proves the class number relations of Kronecker (Corollary 2.2).

We only consider elliptic curves over $\mathbb{C}$.

## 1. Modular Polynomials

Let $m \in \mathbb{N}$. Consider the elliptic curve $E=\mathbb{C} / \Gamma$ with $\Gamma=\mathbb{Z}+\mathbb{Z} \tau$ for some $\tau \in \mathbb{H}$.
Theorem 1.1 ([L, §5.3,5.1]). - There are canonical bijections between the following sets:
(i) isomorphism classes of isogenies $f: E_{1} \rightarrow E$ of degree $m$ (as yroup schemes over E),
(ii) subgroups $\Gamma_{1} \subseteq \Gamma$ of index $m$,

2000 Mathematics Subject Classification. - 11F32, 11F03, 11G15.
Key words and phrases. - Modular polynomials, representation number of a quadratic form, class number relations.
(iii) $\mathrm{SL}_{2}(\mathbb{Z}) \backslash\left\{A \in M_{2}(\mathbb{Z}) \mid \operatorname{det} A=m\right\}$, and
(iv) $\left\{\left.\left(\begin{array}{ll}a & b \\ 0 & d\end{array}\right) \in M_{2}(\mathbb{Z}) \right\rvert\, a d=m, a \geq 1\right.$ and $\left.0 \leq b<d\right\}$.

All of these sets have $\sigma_{1}(m)=\sum_{d \mid m} d$ elements.

Proof
(i) $\rightarrow$ (ii): Set $\Gamma_{1}:=f_{*} \pi_{1}\left(E_{1}\right) .($ ii $) \rightarrow(i):$ Set $E_{1}:=\mathbb{C} / \Gamma_{1}$.
(ii) $\leftrightarrow$ (iii): Choose a basis $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)\binom{1}{\tau}$ of $\Gamma_{1}$ with $A=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in M_{2}(\mathbb{Z})$.
(iii) $\leftrightarrow$ (iv): Left multiplication by matrices from $\mathrm{SL}_{2}(\mathbb{Z})$ corresponds to row operations. The matrices in (iv) are obviously inequivalent (the columns must be stabilized).

Now consider pairs $\left(j, j^{\prime}\right)$ of $j$-invariants of elliptic curves $E, E^{\prime}$ such that there is an isogeny $E \rightarrow E^{\prime}$ of degree $m$. These pairs are described by the divisor of a certain polynomial $\varphi_{m}$ :

For $j, j^{\prime} \in \mathbb{C}$ choose elliptic curves $E, E^{\prime}$ having $j$-invariants $j, j^{\prime}$, respectively. Set

$$
\varphi_{m}\left(j, j^{\prime}\right)=\varphi_{m}\left(j(E), j\left(E^{\prime}\right)\right):=\prod_{E_{1}^{\prime} \rightarrow E^{\prime}}\left(j(E)-j\left(E_{1}^{\prime}\right)\right) ;
$$

the product is over isomorphism classes of isogenies $E_{1}^{\prime} \rightarrow E^{\prime}$ of degree $m . \varphi_{m}$ does not depend on the choices made and is a polynomial of degree $\sigma_{1}(m)$ in $j$. For elliptic curves $E, E^{\prime}$, the condition $\varphi_{m}\left(j(E), j\left(E^{\prime}\right)\right)=0$ is equivalent to the existence of an isogeny $E \rightarrow E^{\prime}$ of degree $m$.

Define $\psi_{m}\left(j, j^{\prime}\right)$ by the same formula, but restrict the product to the isogenies which do not factor over some multiplication-by- $n$ map, $n>1$. In the above correspondence, these isogenies correspond to primitive matrices, i.e., matrices whose entries have no common divisor. We have

$$
\varphi_{m}=\prod_{n^{2} \mid m} \psi_{m / n^{2}}
$$

Obviously, $\varphi_{1}(X, Y)=\psi_{1}(X, Y)=X-Y$. As we will see below, $\varphi_{m}$ and $\psi_{m}$ are polynomials; they are called modular polynomials.

Theorem 1.2 ([L, §5.2])
(i) $\varphi_{m}, \psi_{m} \in \mathbb{Z}[X, Y]$.
(ii) $\psi_{m}(X, t)$ is irreducible over $\mathbb{C}(t)$.
(iii) For $m>1$, we have $\psi_{m}(X, Y)=\psi_{m}(Y, X)$. Consequently, $\varphi_{m}(X, Y)=$ $\pm \varphi_{m}(Y, X)$ ("-" precisely if $m$ is a square).

Proof
(i) First notice that the coefficients $k_{i}$ of

$$
\psi_{m}\left(X, j\left(\tau^{\prime}\right)\right)=\prod_{\operatorname{SL}_{2}(\mathbb{Z}) \backslash\left\{A \in M_{2}(\mathbb{Z}) \mid \operatorname{det} A=m, A \text { primitive }\right\}}\left(X-j\left(A \tau^{\prime}\right)\right) \in \mathcal{O}_{\mathbb{C}}[X]
$$

are holomorphic in $\tau^{\prime}$ and invariant under $\mathrm{SL}_{2}(\mathbb{Z})$. From the formula

$$
\begin{equation*}
\psi_{m}\left(X, j\left(\tau^{\prime}\right)\right)=\prod_{a, b, d}\left(X-j\left(\frac{a \tau^{\prime}+b}{d}\right)\right)=\prod_{a, b, d}\left(X-\frac{1}{\left(q^{\prime}\right)^{a / d} \zeta_{m}^{a b}}-744-\ldots\right) \tag{*}
\end{equation*}
$$

( $a, b, d$ as in 1.1 (iv) and $\left(\begin{array}{ll}a & b \\ 0 & d\end{array}\right)$ primitive, $\zeta_{m}:=\mathrm{e}^{2 \pi \mathrm{i} / m}$ ) we see that the $k_{i}$ are meromorphic at infinity. Since the $q$-expansion of the $j$-function has integral coefficients, we have

$$
k_{i} \in \mathbb{Z}\left[\zeta_{m}\right]\left[\left[q^{\prime}\right]\right]\left[\frac{1}{q^{\prime}}\right]
$$

Now there are polynomials $p_{i} \in \mathbb{Z}\left[\zeta_{m}, T\right]$ such that $k_{i}-p_{i}\left(j\left(q^{\prime}\right)\right)$ lies in $q^{\prime} \mathbb{Z}\left[\zeta_{m}\right]\left[\left[q^{\prime}\right]\right]$ and therefore, being a modular function, must vanish identically. Hence, $\psi_{m} \in$ $\mathbb{Z}\left[\zeta_{m}\right][X, Y]$.

There are two operations of $(\mathbb{Z} / m \mathbb{Z})^{\times}$: first, on matrices $\left(\begin{array}{ll}a & b \\ 0 & d\end{array}\right)$ as in 1.1 (iv) by

$$
\sigma\left(\begin{array}{ll}
a & b \\
0 & d
\end{array}\right):=\left(\begin{array}{cc}
a & \sigma b \\
0 & d
\end{array}\right) \quad\left(\text { via }(\mathbb{Z} / d \mathbb{Z})^{\times} \text {on }\{0, \ldots, d-1\} \cong \mathbb{Z} / d \mathbb{Z}\right)
$$

and the first product in $(*)$ is invariant under this operation. Second, $(\mathbb{Z} / m \mathbb{Z})^{\times}$ operates in a compatible way on $\mathbb{Z}\left[\zeta_{m}\right]$ by $\sigma \zeta_{m}=\zeta_{m}^{\sigma}$, and since the coefficients of $\varphi_{m}$ are invariant under this operation, we find that $\psi_{m} \in \mathbb{Z}[X, Y]$.
(ii) By mapping $t \mapsto j$, the field of meromorphic functions on $\mathbb{H}$ becomes an extension field of $\mathbb{C}(t)$ carrying an operation of the group $\mathrm{SL}_{2}(\mathbb{Z})$. By the elementary divisors theorem, it permutes the zeroes of $\psi_{m}(X, t)$ transitively, hence $\psi_{m}(X, t)$ is irreducible over $\mathbb{C}(t)$.
(iii) The condition $\psi_{m}\left(j(E), j^{\prime}(E)\right)=0$ is equivalent to the existence of an isogeny $E \rightarrow E^{\prime}$ of degree $m$ which does not factor over a multiplication-by- $n$ map for some $n>1$. This last property is also true for its dual isogeny, hence $\psi_{m}\left(j\left(E^{\prime}\right), j(E)\right)=0$. For a fixed $j_{0}^{\prime}$, the irreducible polynomial $\psi_{m}\left(X, j_{0}^{\prime}\right)$ is therefore a divisor of $\psi_{m}\left(j_{0}^{\prime}, X\right)$, and conversely. It follows that $\psi_{m}\left(j, j^{\prime}\right)= \pm \psi_{m}\left(j^{\prime}, j\right)$. If the "-" sign is correct, $\psi_{m}(t, t)$ vanishes identically, so $\psi_{m}(X, t)$ has a zero in $\mathbb{C}(t)$, hence the degree of $\psi_{m}(X, t)$ must be 1 . This is true precisely for $m=1$.

From the proof of (iii) we also see that $f_{m}(X):=\varphi_{m}(X, X)$ vanishes if and only if $m$ is a square. If $m$ is not a square, the degree of $f_{m}$ can be read off the $q$-expansion in $(*)$ : set $X=j\left(q^{\prime}\right)$, then because of $a \neq d$, the pole order of one factor is equal to $\max \{1, a / d\}$, hence the pole order of the entire product is

$$
\operatorname{deg} f_{m}=\sum_{a d=m} d \max \{1, a / d\}=\sum_{a d=m} \max \{a, d\}
$$

One also sees that the leading coefficient of $f_{m}$ is $\pm 1$.

## 2. Intersections

We first need to fix some notation. A quadratic space $(L, Q)$ over a ring $R$ consists of a free $R$-module $L$ of finite rank and a quadratic form $Q$ on $L$. The associated bilinear form on $L$ is defined by

$$
(x, y)=Q(x+y)-Q(x)-Q(y) .
$$

The determinant of $Q$ is the element of $R /\left(R^{\times}\right)^{2}$ given by the determinant of the matrix $\left(\left(b_{i}, b_{j}\right)\right)_{i j}$ for some basis $\left\{b_{i}\right\}$ of $L$. The diagonal of $Q$ with respect to some fixed basis $\left\{b_{i}\right\}$ is defined to be the $n$-tuple $\left(Q\left(b_{i}\right)\right)_{i}$ where $n$ is the rank of $L$.

For a quadratic form $F$ on $R^{m}$, we define the representation number $R_{L}(F)$ as the cardinality of the set

$$
\begin{aligned}
& \left\{\left(f_{i}\right) \in L^{m} \mid Q\left(x_{1} f_{1}+\cdots+x_{m} f_{m}\right)=F\left(x_{1}, \ldots, x_{m}\right) \text { for all } x \in R^{m}\right\} \\
= & \left\{\text { isometries }\left(R^{m}, F\right) \rightarrow(L, Q)\right\} .
\end{aligned}
$$

For $R=\mathbb{Z}$ and positive definite $Q$, this set is finite. (For each $x=e_{i}, i=1, \ldots, m$, there are only finitely many possible values of $x_{1} f_{1}+\cdots+x_{m} f_{m}=f_{i}$.)

For a positive integer $D$, let $H(D)$ be the number of $\mathrm{SL}_{2}(\mathbb{Z})$-equivalence classes of positive definite binary quadratic forms over $\mathbb{Z}$ with determinant $D$ (which is welldefined as an element of $\mathbb{Z}$ ), counting the forms equivalent to $e x_{1}^{2}+e x_{2}^{2}$ and $e x_{1}^{2}+$ $e x_{1} x_{2}+e x_{2}^{2}$ for some natural number $e$ with multiplicities $1 / 2$ and $1 / 3$, respectively. If the positive integer $m$ is not a square, we define

$$
G(m):=\sum_{\substack{t \in \mathbb{Z} \\ t^{2} \leq 4 m}} H\left(4 m-t^{2}\right) .
$$

Define $T_{m}:=V\left(\varphi_{m}\right) \subset \mathbb{A}_{\mathbb{C}}^{2}$.
Theorem 2.1 ([GK, 2.4]). - The curves $T_{m_{1}}$ and $T_{m_{2}}$ intersect properly if and only if $m=m_{1} m_{2}$ is not a square. In this case, their intersection is supported on pairs $\left(E, E^{\prime}\right)$ of elliptic curves with complex multiplication by orders whose discriminants satisfy $d(E), d\left(E^{\prime}\right) \geq-4 m$. The intersection number is

$$
T_{m_{1}} \cdot T_{m_{2}}=\sum_{\substack{t \in \mathbb{Z} \\ t^{2}<4 m}} \sum_{d \mid \operatorname{gcd}\left(m_{1}, m_{2}, t\right)} d \cdot H\left(\frac{4 m-t^{2}}{d^{2}}\right)=\sum_{n \mid \operatorname{gcd}\left(m_{1}, m_{2}\right)} n \cdot G\left(m / n^{2}\right)
$$

Proof. - If $m=m_{1} m_{2}$ is a square, $T_{m_{1}}$ and $T_{m_{2}}$ contain $V\left(\psi_{g}\right), g=\operatorname{gcd}\left(m_{1}, m_{2}\right)$, as a common component (note that $m_{1} / g$ and $m_{2} / g$ are coprime, hence squares themselves). Conversely, if $T_{m_{1}}$ and $T_{m_{2}}$ do not intersect properly, they must contain some $V\left(\psi_{g}\right)$ as a common component, but then $g=m_{1} / n_{1}^{2}=m_{2} / n_{2}^{2}$, so $m=g^{2} n_{1}^{2} n_{2}^{2}$ is a square.

For a pair of elliptic curves $\left(E, E^{\prime}\right)$ corresponding to an intersection point of $T_{m_{1}}$ and $T_{m_{2}}$, there are isogenies $f_{1}, f_{2}: E \rightarrow E^{\prime}$ of degrees $m_{1}$ and $m_{2}$, respectively. Then, $\alpha:={ }^{t} f_{2} f_{1}$ is an endomorphism of $E$ of degree $m$. Since $m$ is not a
square, $E$ has complex multiplication, and $\mathbb{Z}+\mathbb{Z} \alpha$ is a sublattice of End $E$. Hence, its discriminant $(\operatorname{Tr} \alpha)^{2}-4 m<0$ is divisible by $d(E)$, so

$$
d(E) \geq(\operatorname{Tr} \alpha)^{2}-4 m \geq-4 m
$$

Similarly, considering $\beta:=f_{2}{ }^{t} f_{1}$, it follows that $d\left(E^{\prime}\right) \geq-4 m$.
Next, we compute the local intersection number at some point $\left(j_{0}, j_{0}^{\prime}\right) \in \mathbb{C}^{2}$ corresponding to a pair of elliptic curves $\left(E, E^{\prime}\right)$. Set $u_{E}:=\frac{1}{2} \#$ Aut $E$, similarly for $E^{\prime}$. Choose $\tau_{0}^{\prime} \in \mathbb{H}$ such that $j\left(\tau_{0}^{\prime}\right)=j_{0}^{\prime}$. Locally at $\tau_{0}^{\prime}$, the map $j: \mathbb{H} \rightarrow \mathbb{C}$ is a branched covering of degree $u_{E^{\prime}}$, so the local intersection number in the $\left(j, j^{\prime}\right)$-plane is the intersection number in the $\left(j, \tau^{\prime}\right)$-plane divided by $u_{E^{\prime}}$.

In the $\left(j, \tau^{\prime}\right)$-plane, the $\varphi_{m_{i}}$ decompose into factors of the form

$$
j-j\left(A_{i} \tau^{\prime}\right) \quad \text { where } A_{i} \in M_{2}(\mathbb{Z}), \operatorname{det} A_{i}=m_{i}
$$

Therefore, it suffices to compute the local intersection number of two such factors, both vanishing at $\left(j_{0}, \tau_{0}^{\prime}\right)$. This number is the zero order of

$$
\begin{equation*}
j\left(A_{1} \tau^{\prime}\right)-j\left(A_{2} \tau^{\prime}\right) \tag{**}
\end{equation*}
$$

at $\tau^{\prime}=\tau_{0}^{\prime}$. Since $A_{1} \tau_{0}^{\prime}$ and $A_{2} \tau_{0}^{\prime}$ are $\mathrm{SL}_{2}(\mathbb{Z})$-equivalent, we may assume that $A_{1} \tau_{0}^{\prime}=$ $A_{2} \tau_{0}^{\prime}=: \tau_{0}$ and $c_{2}=0$. Locally at $\tau_{0}$,

$$
j(\tau)=j\left(\tau_{0}\right)+s \cdot\left(\tau-\tau_{0}\right)^{u_{E}}+\text { higher order terms }
$$

for some $s \neq 0$, hence $(* *)$ is of the form

$$
\begin{aligned}
& s\left(\frac{a_{1} \tau^{\prime}+b_{1}}{c_{1} \tau^{\prime}+d_{1}}-\frac{a_{1} \tau_{0}^{\prime}+b_{1}}{c_{1} \tau_{0}^{\prime}+d_{1}}\right)^{u_{E}}-s\left(\frac{a_{2} \tau^{\prime}+b_{2}}{d_{2}}-\frac{a_{2} \tau_{0}^{\prime}+b_{2}}{d_{2}}\right)^{u_{E}}+\text { h.o.t. } \\
= & s\left(\frac{\operatorname{det} A_{1}}{\left(c_{1} \tau_{0}^{\prime}+d_{1}\right)^{2}} \cdot\left(\tau^{\prime}-\tau_{0}^{\prime}\right)\right)^{u_{E}}-s\left(\frac{\operatorname{det} A_{2}}{d_{2}^{2}} \cdot\left(\tau^{\prime}-\tau_{0}^{\prime}\right)\right)^{u_{E}}+\text { h.o.t. }
\end{aligned}
$$

locally at $\tau_{0}^{\prime}$. We now claim that the two leading coefficients are different. (However, they have the same absolute value.) Otherwise, from

$$
s\left(\frac{\operatorname{det} A_{1}}{\left(c_{1} \tau_{0}^{\prime}+d_{1}\right)^{2}}\right)^{u_{E}}=s\left(\frac{\operatorname{det} A_{2}}{d_{2}^{2}}\right)^{u_{E}}
$$

we get

$$
\frac{c_{1} \tau_{0}^{\prime}+d_{1}}{\sqrt{m_{1}}}=\omega \cdot \frac{d_{2}}{\sqrt{m_{2}}}
$$

for some $2 u_{E}$-th root of unity $\omega$, implying

$$
\begin{equation*}
c_{1} \tau_{0}^{\prime}+d_{1}=\omega \cdot \frac{\sqrt{m_{1} m_{2}}}{a_{2}} \tag{***}
\end{equation*}
$$

The left-hand side is imaginary-quadratic, so by our assumption that $m=m_{1} m_{2}$ is not a square it follows that $\omega= \pm \mathrm{i}$ and $u_{E}=2$. But in this case, $\tau_{0}^{\prime}$ corresponds to an elliptic curve isogenous to $E \cong \mathbb{C} /\langle 1, \mathrm{i}\rangle$, hence $\tau_{0}^{\prime} \in \mathbb{Q}(\mathrm{i})$, contradicting $(* * *)$. Hence, the zero order of $(* *)$ at $\tau_{0}^{\prime}$ equals $u_{E}$.

Since the product decomposition of $\varphi_{m_{i}}$ contains

$$
\frac{1}{2 u_{E}} \cdot \#\left\{f_{i} \in \operatorname{Hom}\left(E, E^{\prime}\right) \mid \operatorname{deg} f_{i}=m_{i}\right\}
$$

factors vanishing at $\left(j_{0}, \tau_{0}^{\prime}\right)$, the local intersection number in the $\left(j, \tau^{\prime}\right)$-plane is

$$
\left(T_{m_{1}} \cdot T_{m_{2}}\right)_{\left(j_{0}, \tau_{0}^{\prime}\right)}=\frac{1}{4 u_{E}^{2}} \cdot \#\left\{\left(f_{1}, f_{2}\right) \in \operatorname{Hom}\left(E, E^{\prime}\right) \mid \operatorname{deg} f_{i}=m_{i}\right\} \cdot u_{E}
$$

Hence, the local intersection number in the $\left(j, j^{\prime}\right)$-plane is

$$
\left(T_{m_{1}} \cdot T_{m_{2}}\right)_{\left(j_{0}, j_{0}^{\prime}\right)}=\frac{1}{4 u_{E} u_{E^{\prime}}} \cdot \#\left\{\left(f_{1}, f_{2}\right) \in \operatorname{Hom}\left(E, E^{\prime}\right) \mid \operatorname{deg} f_{i}=m_{i}\right\}
$$

Such pairs ( $f_{1}, f_{2}$ ) correspond to representations of positive definite quadratic forms $Q\left(x_{1}, x_{2}\right)=\operatorname{deg}\left(x_{1} f_{1}+x_{2} f_{2}\right)$, hence

$$
\#\left\{\left(f_{1}, f_{2}\right) \in \operatorname{Hom}\left(E, E^{\prime}\right) \mid \operatorname{deg} f_{i}=m_{i}\right\}=\sum_{\substack{Q>0 \\ \operatorname{diag} Q=\left(m_{1}, m_{2}\right)}} R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q) .
$$

Therefore, the global intersection number is

$$
\begin{aligned}
T_{m_{1}} \cdot T_{m_{2}}= & \sum_{\substack{E, E^{\prime} \\
\text { ell.curves } / \mathbb{C} \operatorname{diag}}} \sum_{\substack{Q>0 \\
Q=\left(m_{1}, m_{2}\right)}} \frac{R_{\mathrm{Hom}\left(E, E^{\prime}\right)}(Q)}{4 u_{E} u_{E^{\prime}}} \\
= & \sum_{\substack{Q>0}} \sum_{E, E^{\prime}} \frac{R_{\mathrm{Hom}\left(E, E^{\prime}\right)}(Q)}{4 u_{E} u_{E^{\prime}}}
\end{aligned}
$$

By Proposition 1.1 in [Go1], the inner sum equals

$$
\sum_{d \mid e(Q)} d \cdot H\left(\frac{\operatorname{det} Q}{d^{2}}\right)
$$

In our case, $Q\left(x_{1}, x_{2}\right)=m_{1} x_{1}^{2}+t x_{1} x_{2}+m_{2} x_{2}^{2}$ for some $t \in \mathbb{Z}$ satisfying $t^{2}-4 m<0$ (as $Q$ is positive definite), so the above sum is equal to

$$
\sum_{d \mid \operatorname{gcd}\left(m_{1}, m_{2}, t\right)} d \cdot H\left(\frac{4 m-t^{2}}{d^{2}}\right)
$$

Putting everything together yields

$$
T_{m_{1}} \cdot T_{m_{2}}=\sum_{\substack{t \in \mathbb{Z} \\ t^{2}<4 m}} \sum_{d \mid \operatorname{gcd}\left(m_{1}, m_{2}, t\right)} d \cdot H\left(\frac{4 m-t^{2}}{d^{2}}\right)
$$

As a corollary, we get the class number relations of Kronecker and Hurwitz:
Corollary 2.2. - If $m$ is not a square,

$$
G(m)=\sum_{a d=m} \max \{a, d\}
$$

Proof. - By the remarks at the end of the preceding section,

$$
T_{1} \cdot T_{m}=\operatorname{deg} f_{m}=\sum_{a d=m} \max \{a, d\}
$$

Actually, using the convention $H(0):=\zeta(-1)=-1 / 12$, the above corollary is valid for all $m([\mathbf{W}, \S 116])$.

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# 3. A SUM OF REPRESENTATION NUMBERS 

by
Ulrich Görtz


#### Abstract

This article contains the proof of a formula stated in the paper by Gross and Keating on intersections of modular correspondences, for a certain sum of representation numbers.

Résumé (Une somme de nombres de représentations). - Cet article contient la preuve d'une formule donnée dans l'article de Gross et Keating sur les intersections de correspondances modulaires, pour une certaine somme de nombres de représentations.


## 1. Introduction

We prove a formula for a certain sum of representation numbers, stated in the paper of Gross and Keating $[\mathbf{G K}]$ without proof, which is used in $[\mathbf{V g}]$ in order to compute the intersection product of two modular divisors in $S_{\mathbb{C}}$. Let $Q$ be a positive definite binary quadratic form over $\mathbb{Z}$, say

$$
Q\left(x_{1}, x_{2}\right)=m_{1} x_{1}^{2}+t x_{1} x_{2}+m_{2} x_{2}^{2}
$$

The determinant of $Q$ is

$$
\operatorname{det}(Q)=4 m_{1} m_{2}-t^{2}(>0)
$$

and its content is

$$
e(Q)=\operatorname{gcd}\left(m_{1}, m_{2}, t\right)
$$

## Proposition 1.1

$$
\sum_{\substack{E, E^{\prime} \\ \text { ell. curves } / \mathbb{C}}} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{\# \operatorname{Aut}(E) \cdot \# \operatorname{Aut}\left(E^{\prime}\right)}=\sum_{d \mid e(Q)} d \cdot H\left(\operatorname{det}(Q) / d^{2}\right)
$$

2000 Mathematics Subject Classification. - 14K22.
Key words and phrases. - Elliptic curves, complex multiplication, representation densities.

Our argument is inspired by Hirzebruch's article $[\mathbf{H}]$, where the case $m_{1}=1$ is treated.
Acknowledgments. - I am grateful to Gunther Vogel for a discussion of this problem, and to Torsten Wedhorn for proof-reading.

## 2. Proof of the proposition

The sum on the left hand side extends over isomorphism classes of elliptic curves, and clearly the representation number $R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)$ is 0 unless $E$ and $E^{\prime}$ have complex multiplication and $\operatorname{End}(E) \otimes \mathbb{Q} \cong \operatorname{End}\left(E^{\prime}\right) \otimes \mathbb{Q}$. In particular, the sum is finite.

As in $[\mathbf{G K}]$, we denote by $H(D), D$ a positive integer, the number of $S L_{2}(\mathbb{Z})$ equivalence classes of positive definite binary quadratic forms over $\mathbb{Z}$ with determinant $D$, where the forms equivalent to $e x_{1}^{2}+e x_{2}^{2}$ and $e x_{1}^{2}+e x_{1} x_{2}+e x_{2}^{2}$ for some $e \in \mathbb{Z}$ are counted with multiplicity $1 / 2$ and $1 / 3$, respectively. A quadratic form is called primitive, if its content is 1 . We denote by $h(D)$ the number of primitive positive definite binary quadratic forms of discriminant $D$ if $D>4$, and we set $h(3)=\frac{1}{3}$, $h(4)=\frac{1}{2}$. We can also interpret $h(D)$ as the number of elliptic curves $E$ with complex multiplication, such that the endomorphism ring $\operatorname{End}(E)$ (which is an order in some imaginary quadratic number field) has discriminant $-D$, where each such $E$ is counted with multiplicity $2 / \# \operatorname{Aut}(E)$.

For a positive integer $N$ we denote by $\sigma_{1}(N)$ the sum of all divisors of $N$. Since clearly $H(D)=\sum_{d, d^{2} \mid D} h\left(D / d^{2}\right)$, we can then rewrite the right hand side of the formula as

$$
\sum_{d, d^{2} \mid \operatorname{det}(Q)} \sigma_{1}\left(\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)\right) h\left(\operatorname{det}(Q) / d^{2}\right) .
$$

Fix an elliptic curve $E$ with complex multiplication. We use the following notation:
Write $E=\mathbb{C} / \mathbb{Z} \oplus \mathbb{Z} \tau$ with $\tau \in \mathbb{H}$, and let $\alpha, \beta, \gamma \in \mathbb{Z}$, such that $\alpha \tau^{2}+\beta \tau+\gamma=0$, $\operatorname{gcd}(\alpha, \beta, \gamma)=1, \alpha>0$ (once $\tau$ is fixed, $\alpha, \beta$ and $\gamma$ are uniquely determined by these conditions).

If there exists an $E^{\prime}$, such that $R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q) \neq 0$, then there exists a natural number $d$ with

$$
\begin{equation*}
4 m_{1} m_{2}-t^{2}=\operatorname{det}(Q)=d^{2}\left(4 \alpha \gamma-\beta^{2}\right) \tag{2.1}
\end{equation*}
$$

Indeed, by assumption there exist $f_{i} \in \operatorname{Hom}\left(E, E^{\prime}\right), i=1,2$, such that $\operatorname{deg}\left(f_{i}\right)=m_{i}$ and $\operatorname{deg}\left(f_{1}+f_{2}\right)-\operatorname{deg}\left(f_{1}\right)-\operatorname{deg}\left(f_{2}\right)=t$. Let $g=f_{1}^{\vee} \circ f_{2}$. If we choose lattices $\Lambda, \Lambda^{\prime}$ such that $E \cong \mathbb{C} / \Lambda, E^{\prime} \cong \mathbb{C} / \Lambda^{\prime}$, then we get inclusions $\operatorname{Hom}\left(E, E^{\prime}\right) \subset \mathbb{C}$, $\operatorname{End}(E) \subset \mathbb{C}$, and have $g=m_{1} f_{1}^{-1} f_{2}$ (although $f_{1}$ and $f_{2}$ as complex numbers depend on the choice of $\Lambda$ and $\Lambda^{\prime}, g$ is independent of these choices). Since $g$ has norm $m_{1} m_{2}$ and trace $t$, the quadratic space generated by 1 and $g$ inside $\operatorname{End}(E)$ has determinant $4 m_{1} m_{2}-t=\operatorname{det}(Q)$. Since the determinant of the quadratic space
$\operatorname{End}(E)$ is $4 \alpha \gamma-\beta^{2}$, this implies the existence of $d$ as above. In particular, (2.1) implies that $\frac{t-d \beta}{2}, \frac{t+d \beta}{2} \in \mathbb{Z}$.

From now on, in addition to fixing $E$ as above, we let $g \in \mathbb{H}$ be the (unique) algebraic integer in $\mathbb{H}$ with norm $\mathrm{Nm}_{\mathbb{C} / \mathbb{R}} g=m_{1} m_{2}$ and trace $\operatorname{Tr}_{\mathbb{C} / \mathbb{R}} g=t$. We define
$\mathcal{D}_{i}=\left\{\left(E^{\prime}, f\right) ; E^{\prime}\right.$ an elliptic curve, $\left.f \in \operatorname{Hom}\left(E, E^{\prime}\right), \operatorname{deg}(f)=m_{i}, m_{i} \mid g f\right\} / \cong$
Here (and similarly below) two pairs $\left(E_{1}^{\prime}, f_{1}\right),\left(E_{2}^{\prime}, f_{2}\right)$ are called isomorphic if there exists an isomorphism $\varphi: E_{1}^{\prime} \rightarrow E_{2}^{\prime}$ such that $f_{2} \circ \varphi=f_{1}$. By definition of the sets $\mathcal{D}_{i}$, the set

$$
\begin{aligned}
& \left\{\left(E^{\prime}, f_{1}, f_{2}\right) ; E^{\prime} \text { ell. curve, } f_{i} \in \operatorname{Hom}\left(E, E^{\prime}\right)\right. \\
& \left.\operatorname{deg}\left(f_{i}\right)=m_{i}, \operatorname{deg}\left(f_{1}+f_{2}\right)=t+m_{1}+m_{2}\right\} / \cong
\end{aligned}
$$

maps bijectively to the disjoint union $\mathcal{D}_{1} \cup \mathcal{D}_{2}$, by sending a triple $\left(E^{\prime}, f_{1}, f_{2}\right)$ to $f_{1}$ or $f_{2}$, respectively, depending on whether $m_{1} f_{1}^{-1} f_{2} \in \mathbb{H}$ or $m_{2} f_{2}^{-1} f_{1} \in \mathbb{H}$, i. e. whether $m_{1} f_{1}^{-1} f_{2}=g$ or $m_{2} f_{2}^{-1} f_{1}=g$.

The key point in the proof of the proposition is the following lemma.
Lemma 2.1. - The set $\mathcal{D}_{i}$ can be identified with the set of matrices $\left(\begin{array}{cc}A & B \\ 0 & D\end{array}\right) \in M_{2}(\mathbb{Z})$, such that:
i) There exists $Z \mid \operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)$ such that $D=\frac{Z m_{i}}{\operatorname{gcd}\left(d \alpha, \frac{m_{i \beta}}{2}, m_{i}\right)}, A=\frac{m_{i}}{D}$.
ii) $0 \leq B<D$, such that $B$ satisfies a congruence of the form:

$$
B \equiv b \bmod \frac{D}{Z}
$$

where $b \in \mathbb{Z} / \frac{D}{Z} \mathbb{Z}$ is an element depending on $Z$.
Proof. - To ease the notation a little bit, we assume that $i=1$. Every matrix $M=\left(\begin{array}{cc}A & B \\ 0 & D\end{array}\right)$ with $A, B, D \in \mathbb{Z}_{\geq 0}, A D=m_{1}$ and $0 \leq B<D$ defines an isogeny

$$
E=\mathbb{C} / \mathbb{Z} \oplus \mathbb{Z} \tau \longrightarrow E^{\prime}:=\mathbb{C} / \mathbb{Z} \oplus \mathbb{Z}(M \tau), \quad x \longmapsto A x
$$

and —up to isomorphism- all isogenies of degree $m_{1}$ with source $E$ arise in this way (see $[\mathbf{V g}]$ ).

We need to find out under which conditions the isogeny $f$ corresponding to $A, B, D$ has the property that $m_{1} \mid g f$. This is equivalent to

$$
\frac{A g}{m_{1}} \mathbb{Z} \oplus \mathbb{Z} \tau \subseteq \mathbb{C} / \mathbb{Z} \oplus \mathbb{Z}(M \tau)
$$

hence to

$$
\begin{aligned}
& g \in D \mathbb{Z} \oplus \mathbb{Z}(A \tau+B), \\
& g \tau \in D \mathbb{Z} \oplus \mathbb{Z}(A \tau+B)
\end{aligned}
$$

It is not hard to check that $g=\frac{t+d \beta}{2}+d \alpha \tau$ and that $g \tau=-d \gamma+\frac{t-d \beta}{2} \tau$, and we find that the conditions above are equivalent to the following:

$$
\begin{align*}
& A|d \alpha, \quad A| \frac{t-d \beta}{2}  \tag{2.2}\\
& \frac{d \alpha}{A} B \equiv \frac{t+d \beta}{2} \bmod D  \tag{2.3}\\
& \frac{t-d \beta}{2 A} B \equiv-d \gamma \bmod D . \tag{2.4}
\end{align*}
$$

These congruences for $B$ are solvable if and only if

$$
\begin{equation*}
\left.\operatorname{gcd}\left(\frac{d \alpha}{A}, D\right) \right\rvert\, \frac{t+d \beta}{2}, \quad \text { and } \left.\operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right) \right\rvert\, d \gamma \tag{2.5}
\end{equation*}
$$

respectively, and they are solvable simultaneously if and only if in addition

$$
\frac{d \gamma}{\operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right)} \cdot \frac{d \alpha}{A \operatorname{gcd}\left(\frac{d \alpha}{A}, D\right)} \equiv \frac{t+d \beta}{2 \operatorname{gcd}\left(\frac{d \alpha}{A}, D\right)} \cdot \frac{t-d \beta}{2 A \operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right)} \bmod \frac{D}{l}
$$

where

$$
l=\operatorname{lcm}\left(\operatorname{gcd}\left(\frac{d \alpha}{A}, D\right), \operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right)\right)=\frac{\operatorname{gcd}\left(\frac{d \alpha}{A}, D\right) \operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right)}{\operatorname{gcd}\left(\frac{d \alpha}{A}, \frac{t-d \beta}{2 A}, D\right)}
$$

and this condition is equivalent to

$$
D \left\lvert\, \frac{d^{2} \alpha \gamma-\frac{(t+d \beta)(t-d \beta)}{4}}{A \operatorname{gcd}\left(\frac{d \alpha}{A}, \frac{t-d \beta}{2 A}, D\right)}=\frac{m_{1} m_{2}}{\operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)} .\right.
$$

From this we see that the above congruences for $B$ are simultaneously solvable if and only if

$$
\begin{equation*}
Z: \left.=\frac{D \operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)}{m_{1}} \right\rvert\, m_{2} \tag{2.6}
\end{equation*}
$$

(note that $Z \in \mathbb{Z}$ because $A \left\lvert\, \operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)\right.$ ) and that in this case the set of solutions is a residue class modulo $\frac{D}{Z} \mathbb{Z}$, as condition ii) asserts.

So for $A, D>0$ with $A D=m_{1}$, there exists a $B$ such that the triple $(A, B, D)$ gives rise to an element of $\mathcal{D}_{1}$ if and only if $A, D$ satisfy (2.2), (2.5) and (2.6), and what remains to show is that these conditions are equivalent to condition i) in the lemma.

However, given $(A, B, D)$, we have already defined the $Z$ in the lemma, such that $D$ and $A$ have got the desired form, so we only have to show that

1) if $(A, B, D)$ defines an element of $\mathcal{D}_{1}$, and $Z$ is defined as in (2.6), then $Z \mid m_{1}$, $Z \mid t$ and $Z \mid d$ (since we know already that $Z \mid m_{2}$ ),
2) if we have $Z \mid \operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)$ and define $A$ and $D$ as in i), then $A, D \in \mathbb{Z}$, and (2.2) and (2.5) automatically hold.
ad 1) Since $Z$ is a divisor of $D$, it is clear that $Z \mid m_{1}$. Note that $Z=$ $\operatorname{gcd}\left(\frac{d \alpha}{A}, \frac{t-d \beta}{2 A}, D\right)$, so obviously $Z \mid d \alpha$ and $Z \left\lvert\, \frac{t-d \beta}{2}\right.$. Furthermore, (2.2) implies
that $Z\left|\frac{t+d \beta}{2}, Z\right| d \gamma$. So for one thing, $Z \left\lvert\, \frac{t+d \beta}{2}\right.$ and $Z \left\lvert\, \frac{t-d \beta}{2}\right.$, hence $Z \mid t$ and $Z \mid d \beta$. In addition, we have seen that $Z|d \alpha, Z| d \beta$ and $Z \mid d \gamma$, and since $\operatorname{gcd}(\alpha, \beta, \gamma)=1$, we conclude that $Z \mid d$.
ad 2) Given a divisor $Z$ of $\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)$, we define $D=\frac{Z m_{1}}{\operatorname{gcd}\left(d \alpha, \frac{m_{1 \beta}}{2}, m_{1}\right)}, A=$ $\frac{m_{1}}{D}=\frac{g c d\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)}{Z}$. It is obvious that $D \in \mathbb{Z}$, and in order to prove that $A \in \mathbb{Z}$, all we need to show is that $Z \left\lvert\, \frac{t-d \beta}{2}\right.$. However, it is clear that $Z|t-d \beta, Z| t+d \beta$, and from (2.1) we get that $Z^{2} \left\lvert\, \frac{(t-d \beta)(t+d \beta)}{4}\right.$. Since $t-d \beta \equiv t+d \beta \bmod 2$, this implies $Z \left\lvert\, \frac{t-d \beta}{2}\right.$.

It remains to show that the conditions in (2.2) and (2.5) hold: It is clear that $A \mid d \alpha$ and $A \left\lvert\, \frac{t-d \beta}{2}\right.$. Next, let us show that $\left.\operatorname{gcd}\left(\frac{d \alpha}{A}, D\right) \right\rvert\, \frac{t+d \beta}{2}$. Since we have

$$
\operatorname{gcd}\left(\frac{d \alpha}{A}, D\right)=\frac{Z \operatorname{gcd}\left(d \alpha, m_{1}\right)}{\operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)}
$$

it suffices to show

$$
\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right) \operatorname{gcd}\left(d \alpha, m_{1}\right) \left\lvert\, \frac{t+d \beta}{2} \operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)\right.
$$

We use the following notation: for $x \in \mathbb{Z}$ such that $\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right) \mid x$, let $\tilde{x}=$ $\frac{x}{\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)}$. From (2.1) we get

$$
\frac{\tilde{t}-\tilde{d} \beta}{2} \frac{\tilde{t}+\tilde{d} \beta}{2}=\tilde{m_{1}} \tilde{m_{2}}-(\tilde{d})^{2} \alpha \gamma
$$

which implies

$$
\operatorname{gcd}\left(\tilde{d} \alpha, \tilde{m}_{1}\right) \left\lvert\, \frac{\tilde{t}+\tilde{d} \beta}{2} \operatorname{gcd}\left(\tilde{d} \alpha, \frac{\tilde{t}-\tilde{d} \beta}{2}, \tilde{m}_{1}\right) .\right.
$$

Multiplying both sides by $\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)^{2}$, we get the desired result.
Finally, in a similar way we can show that $\left.\operatorname{gcd}\left(\frac{t-d \beta}{2 A}, D\right) \right\rvert\, d \gamma$. Namely, it is enough to show

$$
\left.\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right) \operatorname{gcd}\left(\frac{t-d \beta}{2}, m_{1}\right) \right\rvert\, d \gamma \operatorname{gcd}\left(d \alpha, \frac{t-d \beta}{2}, m_{1}\right)
$$

and this follows from

$$
\tilde{m_{1}} \tilde{m_{2}}-\frac{\tilde{t}-\tilde{d} \beta}{2} \frac{\tilde{t}+\tilde{d} \beta}{2}=(\tilde{d})^{2} \alpha \gamma
$$

This concludes the proof of 2 ), and hence the proof of the lemma.
Corollary 2.2. - We fix $E$ as above, and use the same notation. Then

$$
\begin{aligned}
& \sum_{E^{\prime}} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{\# \operatorname{Aut}\left(E^{\prime}\right)} \\
= & \sum_{E^{\prime}} \frac{\#\left\{\left(f_{1}, f_{2}\right) \in \operatorname{Hom}\left(E, E^{\prime}\right)^{2} ; \operatorname{deg}\left(f_{i}\right)=m_{i}, \operatorname{deg}\left(f_{1}+f_{2}\right)=t+m_{1}+m_{2}\right\}}{\# \operatorname{Aut}\left(E^{\prime}\right)} \\
= & 2 \sigma_{1}\left(\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)\right) .
\end{aligned}
$$

Proof. - This follows from the lemma and the remark preceding it.
Proof of the proposition. - Using the corollary, we can now easily prove the proposition:

$$
\begin{aligned}
& \sum_{E, E^{\prime}} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{\# \operatorname{Aut}(E) \cdot \# \operatorname{Aut}\left(E^{\prime}\right)} \\
= & \sum_{\substack{d \\
d^{2} \mid \operatorname{det}(Q) \operatorname{disc}(\operatorname{End}(E))=-\operatorname{det}(Q) / d^{2}}}^{\# \operatorname{Aut}(E)} \sum_{E^{\prime}} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{\# \operatorname{Aut}\left(E^{\prime}\right)} \\
= & \sum_{\substack{d \\
d^{2} \mid \operatorname{det}(Q) \operatorname{disc}(\operatorname{End}(E))=-\operatorname{det}(Q) / d^{2}}} \frac{2 \sigma_{1}\left(\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)\right)}{\# \operatorname{Aut}(E)} \\
= & \sum_{\substack{d \\
d^{2} \mid \operatorname{det}(Q)}} \sigma_{1}\left(\operatorname{gcd}\left(m_{1}, m_{2}, t, d\right)\right) h\left(\operatorname{det}(Q) / d^{2}\right) . \quad \square
\end{aligned}
$$

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[^1]
# 4. ARITHMETIC INTERSECTION NUMBERS 

by

Ulrich Görtz


#### Abstract

We define the arithmetic intersection number of three modular divisors and interpret it from the point of view of algebraic stacks. A criterion is given when the intersection of three modular divisors is finite. Furthermore, the final result about the arithmetic intersection numbers, as given by Gross and Keating, is stated and the strategy of its proof, carried out in the subsequent chapters, is explained. Résumé (Nombres d'intersection arithmétiques). - On définit les nombres d'intersection arithmétiques de trois diviseurs modulaires, et on donne une interprétation du point de vue des champs algébriques. On en donne un critère pour que cette intersection soit finie. En plus, on indique le résultat final sur les nombres d'intersection arithmétiques, comme donné par Gross et Keating, et la stratégie de sa preuve, effectuée dans les chapitres suivants.


## 1. Introduction

Let us recall some notation: Let $m \geq 1$ be an integer. In $[\mathbf{V g}]$ we have defined the modular polynomial $\varphi_{m} \in \mathbb{Z}\left[j, j^{\prime}\right]$ (we regard $j, j^{\prime}$ as indeterminates). We denote by $T_{m} \subseteq \operatorname{Spec} \mathbb{Z}\left[j, j^{\prime}\right]$ the associated divisor. Write $S=\operatorname{Spec} \mathbb{Z}\left[j, j^{\prime}\right]$, and $S_{\mathbb{C}}=$ $\operatorname{Spec} \mathbb{C}\left[j, j^{\prime}\right]$.

In this chapter, we will first prove a criterion for the intersection of three modular divisors over Spec $\mathbb{Z}$ to be finite, which is analogous to the criterion of Hurwitz in the complex situation (see $[\mathbf{V g}]$ ).

In the second part we will prove, following $[\mathbf{G K}]$ and using results of later chapters, Gross' and Keating's explicit formula for the arithmetic intersection number: Fix positive integers $m_{1}, m_{2}$ and $m_{3}$. The arithmetic intersection number is, by definition,

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{S}:=\log \# \mathbb{Z}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right)
$$

[^2]Key words and phrases. - Modular divisors.

This number has a natural interpretation in the Arakelov theory for stacks (see below). In the proof, we use the properties of the invariants $\alpha_{p}(Q)$ and $\beta_{\ell}(Q)$ which will be established in later chapters. Altogether, this yields the proof of Theorem 1.2 in the introduction.
Acknowledgments. - I am grateful to all the participants of the ARGOS seminar for discussions and for feedback on these notes. In particular, I want to thank I. Bouw for her comments. I also profited from discussions with S. Kudla. Finally, I thank the anonymous referee for a number of helpful remarks.

## 2. Preliminaries, Notation

2.1. Quadratic forms and lattices in quadratic number fields. - There is a dictionary between binary quadratic forms (over $\mathbb{Z}$ ) and lattices in quadratic number fields (see [BS] II §7.5, in particular Satz 4). The exact statement we will use is the following.

Let $d<0$ be a square-free integer. Denote by $\mathcal{L}$ the set of $\mathbb{Z}$-lattices in $\mathbb{Q}(\sqrt{d})$ up to homothety, and denote by $\mathcal{F}$ the set of positive definite primitive binary quadratic forms over $\mathbb{Z}$ which split in $\mathbb{Q}(\sqrt{d})$, up to proper equivalence. Then there is a bijection

$$
\mathcal{L} \longrightarrow \mathcal{F}, \quad L \longmapsto \frac{N(\alpha x+\beta y)}{N(L)}
$$

where $N: \mathbb{Q}(\sqrt{d}) \rightarrow \mathbb{Q}$ denotes the norm, $N(L)=\operatorname{gcd}(N(l) ; l \in L \backslash\{0\})$, and $\alpha, \beta$ is a basis of $L$ such that $\frac{1}{i}(\alpha \bar{\beta}-\bar{\alpha} \beta)>0$ (here $\cdot$ denotes conjugation).
2.2. Stacks. - We mostly work with the coarse moduli space of (pairs of) elliptic curves, but in a few places it is more convenient to use the language of stacks. For the convenience of the reader, in this section we give a few references to the literature about the results that we need. A general reference is the book [LM] by Laumon and Moret-Bailly. See also Deligne's and Mumford's article $[\mathbf{D M}]$. For the stacks that we are concerned with the main reference is the book $[\mathbf{K M}]$ of Katz and Mazur: although superficially the language of stacks is not used there, it is obvious that their results can be understood as results about stacks.

We denote by $\mathcal{M}$ the moduli stack (over $\mathbb{Z}$ ) of elliptic curves; this is a DeligneMumford stack.

We denote by $\mathcal{T}_{m}$ the moduli space of isogenies of elliptic curves of degree $m$. (In $[\mathbf{K M}]$, the notation [ $m$-Isog] is used.) This is a Deligne-Mumford stack, too, and furthermore, we have:

Proposition 2.1. - The morphism $\mathcal{T}_{m} \rightarrow \mathcal{M}$ is finite and flat, and is étale over $\mathbb{Z}\left[\frac{1}{m}\right]$. The morphism $\mathcal{T}_{m} \rightarrow \mathcal{M} \times \mathcal{M}$ is finite and unramified.

Proof. - The first assertion is just [KM, 6.8.1], and the second one follows immediately from the rigidity theorem, see [KM, 2.4.2].

By relating the divisor $T_{m}$ (inside the coarse moduli space) defined by the modular polynomials $\varphi_{m}$ to the space $\mathcal{T}_{m}$, we get a description of the geometric points of $T_{m}$.

Lemma 2.2. - Let $m \geq 1$. A geometric point of $T_{m}$ corresponds to a pair $\left(E, E^{\prime}\right)$ of elliptic curves such that there exists an isogeny $E \longrightarrow E^{\prime}$ of degree $m$.

Proof. - In characteristic 0 this is basically the definition of $T_{m}$ and $\varphi_{m}$. In positive characteristic, we can prove this as follows: By mapping an isogeny to its source, we get a finite flat map from $\mathcal{T}_{m}$ to the moduli stack $\mathcal{M}$ of elliptic curves (see [KM, 6.8.1]). In particular, $\mathcal{T}_{m}$ is flat over $\mathbb{Z}$.

Now we have a map to the coarse moduli space $S$ of pairs of elliptic curves:

$$
F: \mathcal{T}_{m} \longrightarrow S, \quad\left(E \rightarrow E^{\prime}\right) \longmapsto\left(j(E), j\left(E^{\prime}\right)\right),
$$

and we get a diagram


Since $p \nmid \varphi_{m}(X, Y), \operatorname{div}\left(\varphi_{m}\right)$ is flat over $\mathbb{Z}_{p}$, and because $\operatorname{im} F_{\mathbb{Z}_{p}}$ is flat over $\mathbb{Z}_{p}$, too, we get $\operatorname{im} F_{\mathbb{Z}_{p}}=\operatorname{div}\left(\varphi_{m}\right)$. Obviously the geometric points of im $F_{\mathbb{Z}_{p}}$ correspond to pairs ( $E, E^{\prime}$ ) of elliptic curves such that there exists an isogeny $E \rightarrow E^{\prime}$ of degree $m$, so the lemma is proved.

We can express the arithmetic intersection number of three 'divisors' $\mathcal{T}_{m_{i}}$ in $\mathcal{S}:=$ $\mathcal{M} \times \mathcal{M}$ in terms of the complete local rings of their 'intersection' $\mathcal{X}:=\mathcal{T}_{m_{1}} \times \mathcal{S} \mathcal{T}_{m_{2}} \times \mathcal{S}$ $\mathcal{T}_{m_{3}}$. (Note however that $T_{m_{1}} \times{ }_{S} T_{m_{2}} \times{ }_{S} T_{m_{3}}$ is not the coarse moduli space of $\mathcal{X}$.)

Proposition 2.3. - Let $\mathcal{X}:=\mathcal{T}_{m_{1}} \times \mathcal{S} \mathcal{T}_{m_{2}} \times \mathcal{S} \mathcal{T}_{m_{3}}$. Then

$$
\begin{aligned}
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right) & :=\log \# \mathbb{Z}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right) \\
& =\frac{1}{2} \sum_{p} \log (p) \cdot \sum_{x \in \mathcal{X}\left(\overline{\mathbb{F}}_{p}\right)} \frac{1}{\# \operatorname{Aut}(x)} \lg \widehat{\mathcal{O}}_{\mathcal{X}, x}
\end{aligned}
$$

Proof. - We may assume that the intersection $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$ is finite, since otherwise both sides are infinite. (See the next section for a precise criterion, when this is the case.) The complete local ring of a geometric point in $\mathcal{M} \times \mathcal{M}$ is the universal
deformation ring of the corresponding pair of elliptic curves, and this ring is free of rank $\frac{\# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right)}{4}$ over the complete local ring in the corresponding point in the coarse moduli space. This gives us (see the remarks at the beginning of section 4 for details) that the local contribution to the intersection number at a point ( $E, E^{\prime}$ ) is

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{\left(E, E^{\prime}\right)}=\sum_{f_{i}, i=1,2,3} \frac{1}{2 \# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right)} \lg _{W} \widehat{\mathcal{O}}_{\mathcal{M} \times \mathcal{M},\left(E, E^{\prime}\right)} / I
$$

where the sum extends over triples of isogenies $f_{i}: E \rightarrow E^{\prime}, \operatorname{deg} f_{i}=m_{i}$, and where $I$ is the smallest ideal in $\widehat{O}_{\mathcal{M} \times \mathcal{M},\left(E, E^{\prime}\right)}$, such that $f_{1}, f_{2}$, and $f_{3}$ lift to isogenies between the universal deformations of $E, E^{\prime}$ modulo $I$.

Now if a triple $f_{1}, f_{2}, f_{3}$ corresponds to the point $x \in \mathcal{X}\left(\overline{\mathbb{F}}_{p}\right)$, then $\widehat{\mathcal{O}}_{\mathcal{M} \times \mathcal{M},\left(E, E^{\prime}\right)} / I=$ $\mathcal{O}_{\mathcal{X}, x}$. Another triple $\left(f_{1}^{\prime}, f_{2}^{\prime}, f_{3}^{\prime}\right)$ yields the same point in $\mathcal{X}$ if and only if there are automorphisms $\varphi$ of $E$ and $\varphi^{\prime}$ of $E^{\prime}$ such that $f_{i}^{\prime}=\varphi^{\prime} \circ f_{i} \circ \varphi^{-1}$ for $i=1,2,3$. Furthermore $\operatorname{Aut}_{\mathcal{X}}(x)$ is isomorphic to the group of $\left(\varphi, \varphi^{\prime}\right) \in \operatorname{Aut}(E) \times \operatorname{Aut}\left(E^{\prime}\right)$ such that $f_{i}=\varphi^{\prime} \circ f_{i} \circ \varphi^{-1}$ for $i=1,2,3$. Hence by splitting up the sum above according to classes of triples which map to the same point in $\mathcal{X}$, we get the claimed equality.
2.3. Notation. - We recall the following notation from $[\mathbf{V g}]$. For an elliptic curve $E$, we let $u_{E}:=\frac{1}{2} \# \operatorname{Aut}(E)$.

Furthermore, given a ring $R$, and a quadratic space $(L, D)$, for a quadratic form $Q$ on $R^{m}$ we define the representation number $R_{L}(Q)$ as the number of isogenies $\left(R^{m}, Q\right) \rightarrow(L, D)$.

## 3. When is $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$ finite?

We start with a lemma which guarantees the existence of elliptic curves such that the homomorphism module represents a given binary quadratic form.

Lemma 3.1. - Let $Q$ be a positive definite binary quadratic form over $\mathbb{Z}$. Then there exist elliptic curves $E, E^{\prime}$ (with complex multiplication) over $\mathbb{C}$ such that $Q \cong$ $\left(\operatorname{Hom}\left(E, E^{\prime}\right), \operatorname{deg}\right)$.

Proof. - By the dictionary between quadratic forms and lattices in imaginary quadratic number fields (see section 2), if $Q$ is a positive definite binary quadratic form over $\mathbb{Z}$ and $Q^{\prime}=\frac{1}{r} Q$ is the associated primitive form, then there exists $d<0$, an order $R_{f}=\mathbb{Z}+f \mathcal{O}_{\mathbb{Q}(\sqrt{d})} \subseteq \mathbb{Q}(\sqrt{d})$ and an ideal $\mathfrak{a} \subseteq R_{f}$ with $\mathbb{Z}$-basis $\alpha$, $\beta$, such that

$$
Q^{\prime}(x, y) \cong \frac{N(\alpha x+\beta y)}{N(\mathfrak{a})}
$$

For the elliptic curves $\mathbb{C} / R_{f r}$ and $\mathbb{C} / \mathfrak{a}$ we then have

$$
\operatorname{Hom}\left(\mathbb{C} / R_{f r}, \mathbb{C} / \mathfrak{a}\right)=\left\{\gamma \in \mathbb{C} ; \gamma R_{f r} \subseteq \mathfrak{a}\right\}=\mathfrak{a}
$$

and for $\gamma \in \operatorname{Hom}\left(\mathbb{C} / R_{f r}, \mathbb{C} / \mathfrak{a}\right)$,

$$
\operatorname{deg} \gamma=\left[\mathfrak{a}: \gamma R_{f r}\right]=r \cdot\left[\mathfrak{a}: \gamma R_{f}\right]=r \cdot \frac{N(\gamma)}{N(\mathfrak{a})}=Q(\gamma) .
$$

It has been shown already by Hurwitz that on $S_{\mathbb{C}}$, two divisors $T_{m_{1}}$ and $T_{m_{2}}$ intersect in dimension 0 if and only if $m_{1} m_{2}$ is not a square; see $[\mathbf{V g}]$. In other words, they intersect in dimension 0 if and only if there is no unary quadratic form $Q$ which represents both $m_{1}$ and $m_{2}$. The following proposition gives us a completely analogous criterion for the intersection of three $T_{m}$ 's on $S$.

Proposition 3.2.- The divisors $T_{m_{1}}, T_{m_{2}}$ and $T_{m_{3}}$ intersect in dimension 0 if and only if there is no positive definite binary quadratic form over $\mathbb{Z}$ which represents $m_{1}$, $m_{2}$ and $m_{3}$.

In this case the support of $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$ is contained in the zero cycle of pairs of supersingular elliptic curves in characteristic $p<4 m_{1} m_{2} m_{3}$.

Proof. - First suppose that $m_{1}, m_{2}, m_{3}$ are represented by the positive definite binary quadratic form $F$. Let $E, E^{\prime}$ be elliptic curves in characteristic 0 (with complex multiplication) such that $\operatorname{Hom}\left(E, E^{\prime}\right) \cong F$. Then $\left(E, E^{\prime}\right)$ corresponds to a point of $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$, so this intersection must have dimension $\geq 1$.

If, on the other hand, there is no positive definite binary quadratic form which simultaneously represents $m_{1}, m_{2}$ and $m_{3}$, then for all points ( $E, E^{\prime}$ ) of $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$ we must have $\operatorname{rk} \operatorname{Hom}\left(E, E^{\prime}\right)>2$, thus $E$ and $E^{\prime}$ are supersingular, and in particular live in positive characteristic.

Now fix a point $\left(E, E^{\prime}\right) \in S_{\mathbb{F}_{p}}$ which lies in the intersection $T_{m_{1}} \cap T_{m_{2}} \cap T_{m_{3}}$. To complete the proof of the proposition, we have to show that $p \leq 4 m_{1} m_{2} m_{3}$. There exist isogenies $f_{i} \in \operatorname{Hom}\left(E, E^{\prime}\right)$ of degree $m_{i}, i=1,2,3$.

Now consider the ternary quadratic form

$$
Q\left(x_{1}, x_{2}, x_{3}\right)=\operatorname{deg}\left(x_{1} f_{1}+x_{2} f_{2}+x_{3} f_{3}\right) .
$$

Since the matrix associated to $Q$ is symmetric and positive definite, its determinant is smaller or equal than the product of the diagonal entries (see [Be, ch. 8, Thm. 5]), i.e.,

$$
\Delta:=\frac{1}{2} \operatorname{det} Q \leq 4 m_{1} m_{2} m_{3} .
$$

Note that $\Delta \in \mathbb{Z}$ (see [B] Lemma 1.1).
Now the proposition follows from the following lemma.
Lemma 3.3. - With notation as above, we have

$$
p \mid \Delta
$$

Proof. - Let us first assume that $p>2$.
We recall the following theorem on quadratic forms over $\mathbb{Q}_{p}$, see $[\mathbf{S e}$, III Thm. 1, IV 2.1 and IV Thm. 6], for instance:

Theorem 3.4. - If $F$ is an anisotropic quadratic form of rank 4 over $\mathbb{Q}_{p}$, then its discriminant is a square, and its Hasse-Witt invariant $\varepsilon_{p}$ is -1 .

Here, if we write $F=\sum_{i=1}^{4} a_{i} x_{i}^{2}, a_{i} \in \mathbb{Q}_{p}$, then

$$
\begin{aligned}
\varepsilon_{p} & =\prod_{i<j}\left(a_{i}, a_{j}\right) \in\{1,-1\}, \quad \text { where }(x, y) \text { is the Hilbert symbol, } \\
(x, y) & =(-1)^{\alpha \beta \frac{p-1}{2}}\left(\frac{u}{p}\right)^{\beta}\left(\frac{v}{p}\right)^{\alpha}, \quad \text { if } x=p^{\alpha} u, y=p^{\beta} v, u, v \in \mathbb{Z}_{p}^{\times}, p \neq 2 .
\end{aligned}
$$

Now $\operatorname{Hom}\left(E, E^{\prime}\right) \otimes \mathbb{Q}$ is isomorphic, up to scaling the form, to $\operatorname{End}(E) \otimes \mathbb{Q}$ with the quadratic form deg. But $\operatorname{End}(E) \otimes \mathbb{Q}$ is the quaternion algebra over $\mathbb{Q}$ ramified exactly at $p$ and $\infty$, and the degree form corresponds to the reduced norm (see $[\mathbf{W d} \mathbf{1}$, 2.2]). Hence $\operatorname{det}\left(\left.\operatorname{deg}\right|_{\operatorname{Hom}\left(E, E^{\prime}\right)}\right)$ is a square. We also see that the quadratic form deg on $\operatorname{Hom}\left(E, E^{\prime}\right)$ is anisotropic over $\mathbb{Q}_{p}$, so its Hasse-Witt invariant $\varepsilon_{p}$ is -1 .

Since the $m_{i}$ are not simultaneously represented by a binary quadratic form, the $f_{i}$ are linearly independent over $\mathbb{Z}$. Now $\operatorname{Hom}\left(E, E^{\prime}\right)$ has square determinant and represents $Q$, so we have

$$
\operatorname{Hom}\left(E, E^{\prime}\right) \otimes \mathbb{Q} \cong Q \perp\langle\Delta\rangle
$$

where $\langle\Delta\rangle$ denotes the unary quadratic form $x \mapsto \Delta x^{2}$. Over $\mathbb{Z}_{p}$ we can diagonalize $Q$ :

$$
Q\left(x_{1}, x_{2}, x_{3}\right)=a x_{1}^{2}+b x_{2}^{2}+c x_{3}^{2}, \quad a, b, c \in \mathbb{Z}_{p} .
$$

Then $\Delta=4 a b c$ and $\varepsilon_{p}=-1$ implies $p \mid a b c$, by the formulas above.
For $p=2$ the bound $p \leq 4 m_{1} m_{2} m_{3}$ holds trivially, but the stronger assertion $p \mid \Delta$ is true in this case too. Namely, by [B] Prop. 4.7, the 2 -adic valuation of $\Delta$ is equal to the sum $a_{1}+a_{2}+a_{3}$ of the Gross-Keating invariants of $Q$ (see loc. cit.). Furthermore, since $Q$ is anisotropic, the $a_{i}$ cannot all be 0 (loc. cit. Lemma 5.3).

This concludes the proof of the lemma, and thus the proof of the proposition, as well.

We conclude this section by the following proposition which reformulates the criterion we obtained above in terms of ternary quadratic forms.

Proposition 3.5. - Let $m_{1}, m_{2}, m_{3}$ be positive integers. The following are equivalent:
(1) There exists no positive definite integral binary quadratic form $Q$ which represents $m_{1}, m_{2}$, and $m_{3}$.
(2) Every positive semi-definite half-integral symmetric matrix $T$ with diagonal entries $m_{1}, m_{2}, m_{3}$ is non-degenerate, i.e., $\operatorname{det} T \neq 0$.
(As usual, by half-integral we mean that the entries outside the diagonal lie in $\frac{1}{2} \mathbb{Z}$, and the diagonal entries are integers. We denote the set of half-integral symmetric $n \times n$ matrices by $\operatorname{Sym}(\mathbb{Z})^{\vee}$.)

Proof. - Given a positive semi-definite $T \in \operatorname{Sym}(\mathbb{Z})^{\vee}$ with $\operatorname{det} T=0$, we get a $Q$ as in (1) as follows: There exists an $x \in \mathbb{Z}^{3}$ such that ${ }^{t} x T x=0$, and we may assume that $x$ is not divisible, i.e., that it generates a direct summand in $\mathbb{Z}^{3}$. Choosing a complement, we get a positive-semidefinite binary quadratic form which represents the $m_{i}$. It could happen that this form is degenerate, but then we can clearly find a positive definite form which still represents all the three $m_{i}$.

On the other hand, given a $Q$ as in (1), choose $x_{i}, y_{i}$, such that $Q\left(x_{i}, y_{i}\right)=m_{i}$, $i=1,2,3$. The matrix $\left(\begin{array}{lll}x_{1} & x_{2} & x_{3} \\ y_{1} & y_{2} & y_{3}\end{array}\right)$ defines a map $\mathbb{Z}^{3} \rightarrow \mathbb{Z}^{2}$, and expressing the ternary quadratic form which we get as the composition of this map with $Q$, we obtain a positive semi-definite half-integral symmetric matrix $T$ with diagonal ( $m_{1}, m_{2}, m_{3}$ ) which is obviously degenerate.

## 4. A formula for the intersection number

From now on, we assume that $T_{m_{1}}, T_{m_{2}}$ and $T_{m_{3}}$ intersect in dimension 0 . We want to explain the final formula which we get for the intersection number, see Theorem 4.3 below. The proofs of the main steps will follow in later chapters.

We write

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{S}=\sum_{p} n(p) \log p
$$

with

$$
n(p)=\lg _{\mathbb{Z}_{p}} \mathbb{Z}_{p}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right)
$$

(and $n(p)=0$ for $\left.p>4 m_{1} m_{2} m_{3}\right)$.
Furthermore, $n(p)$ is the sum of the intersection multiplicities in points $\left(E, E^{\prime}\right)$ given by pairs of supersingular elliptic curves in characteristic $p$. Denote by $j^{(E)}, j^{\left(E^{\prime}\right)}$ their $j$-invariants.

Let $W=W\left(\overline{\mathbb{F}}_{p}\right)$ be the ring of Witt vectors of $\overline{\mathbb{F}}_{p}$, let $\tilde{j}^{(E)}, \tilde{j}^{\left(E^{\prime}\right)} \in W$ be lifts of $j^{(E)}, j^{\left(E^{\prime}\right)}$, respectively, and let $R_{0}$ be the completion of $W\left[j, j^{\prime}\right]$ in the ideal $\mathfrak{m}=$ $\left(p, j-\tilde{j}^{(E)}, j^{\prime}-\tilde{j}^{\left(E^{\prime}\right)}\right)$. Then

$$
R_{0} \cong W\left[\left[j-\tilde{j}^{(E)}, j^{\prime}-\tilde{j}^{\left(E^{\prime}\right)}\right]\right] .
$$

On the other hand, if $R$ denotes the universal deformation ring of the pair ( $E, E^{\prime}$ ), then $R \cong W\left[\left[t, t^{\prime}\right]\right]$, and $R_{0}$ is isomorphic to the $\operatorname{ring} R^{\operatorname{Aut}(E) \times \operatorname{Aut}\left(E^{\prime}\right)}$ of invariants under the finite group $\operatorname{Aut}(E) \times \operatorname{Aut}\left(E^{\prime}\right)$ (cf. $[\mathbf{K M}, 8.2 .3]$ ). Since $R_{0}$ is regular, $R$ is free over $R_{0}$ (see [Ma, Theorem 23.1]) and since $\pm$ id are the only automorphisms of the whole universal deformation, we have $\mathrm{rk}_{R_{0}} R=u_{E} u_{E^{\prime}}$.

We denote by $\left(\mathbb{E}, \mathbb{E}^{\prime}\right)$ the universal pair of elliptic curves over $\operatorname{Spf} R$.

Lemma 4.1. - In $R$, the modular polynomial $\varphi_{m}$ factors as follows:

$$
\varphi_{m}=\prod_{\substack{f: E \rightarrow E^{\prime} \text { isog. of } \\ \text { degree } m, \text { mod } \pm 1}} \varphi_{m, f}
$$

such that for each $f,\left(\varphi_{m, f}\right) \subseteq R$ is the smallest ideal $I \subseteq R$ such that $f$ lifts to an isogeny $\tilde{f}: \mathbb{E} \longrightarrow \mathbb{E}^{\prime}$ modulo $I$.

Proof. - Let $f: E \rightarrow E^{\prime}$ be an isogeny of degree $m$. Then its deformation functor $\operatorname{Def}_{f}$ is pro-represented by a closed subscheme of $\operatorname{Spf} R$ (by the rigidity theorem), and this closed subscheme is a divisor, say $\operatorname{div}\left(\varphi_{m, f}\right), \varphi_{m, f} \in R$. (This is proved in $[\mathbf{K M}$, (6.8)] if $m$ is a power of $p$, but the proof given there works in general. If $p$ does not divide $m$, then $\operatorname{Def}_{f}$ is actually smooth.)

Claim. - If $f$ and $g$ are isogenies $E \rightarrow E^{\prime}$ of degree $m$, then the elements $\varphi_{m, f}$ and $\varphi_{m, g}$ are coprime unless $f= \pm g$.

To prove the claim, suppose that $f$ and $g$ are given such that $\varphi_{m, f}$ and $\varphi_{m, g}$ are not coprime. Then $\operatorname{div}\left(\varphi_{m, f}\right)$ and $\operatorname{div}\left(\varphi_{m, g}\right)$ have a common component $C$. Now $C \otimes \mathbb{Q}$ must have dimension 1 , so $\operatorname{End}\left(\mathbb{E} \otimes_{\operatorname{Spf} R} C\right)=\operatorname{End}\left(\mathbb{E}^{\prime} \otimes_{\operatorname{Spf} R} C\right)=\mathbb{Z}$

By definition of $C$, we have isogenies $f, g: \mathbb{E} \otimes_{\operatorname{Spf} R} C \rightarrow \mathbb{E}^{\prime} \otimes_{\operatorname{Spf} R} C$ of degree $m$. Since ${ }^{t} f \circ f$ and ${ }^{t} f \circ g$ are elements in $\operatorname{End}(\mathbb{E} \otimes \operatorname{Spf} R C)=\mathbb{Z}$ of the same degree, we see that $f= \pm g$. This proves the claim.

Thus we get for the scheme-theoretic union

$$
\bigcup_{f \bmod \pm 1} \operatorname{Def}_{f}=\operatorname{div}\left(\prod_{f \bmod \pm 1} \varphi_{m, f}\right)
$$

Since

$$
\bigcup_{f \text { mod } \pm 1} \operatorname{Def}_{f}(S)=\operatorname{div}\left(\varphi_{m}\right)(S)
$$

for all $S \rightarrow \operatorname{Spf} R$, we obtain that (after possibly changing one of the $\varphi_{m, f}$ 's by a unit)

$$
\varphi_{m}=\prod_{f \bmod \pm 1} \varphi_{m, f}
$$

Lemma 4.2. - Let $A$ be a ring, $B$ an $A$-algebra, and let $x_{1}, \ldots, x_{n} \in B$. If none of the $x_{i}$ is a zero-divisor, then

$$
\lg _{\mathrm{A}} B /\left(x_{1} \cdots x_{n}\right)=\sum_{i=1}^{n} \lg _{A} B /\left(x_{i}\right)
$$

We can write

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)=\sum_{p} \log (p) \sum_{\left(E, E^{\prime}\right) \text { s.s. in char } p}\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{\left(E, E^{\prime}\right)},
$$

and by applying Lemma 4.1 to $\varphi_{m_{i}}$ for $i=1,2,3$, and applying lemma 4.2 successively, we get that the local contribution in a point $\left(E, E^{\prime}\right)$ is

$$
\begin{aligned}
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{\left(E, E^{\prime}\right)} & =\lg _{W} R_{0} /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right) \\
& =\sum_{f_{1}} \sum_{f_{2}} \sum_{f_{3}} \frac{1}{u_{E} u_{E^{\prime}}} \lg _{W} R /\left(\varphi_{\left.m_{1}, f_{1}, \varphi_{m_{2}, f_{2}}, \varphi_{m_{3}, f_{3}}\right)}\right. \\
& =\sum_{f_{i}, i=1,2,3} \frac{1}{u_{E} u_{E^{\prime}}} \lg _{W} R / I,
\end{aligned}
$$

where the sums are over isogenies $f_{i}: E \rightarrow E^{\prime}$ of degree $m_{i}$, up to $\pm 1$, and where $I$ is the smallest ideal in $R$ such that $f_{1}, f_{2}$ and $f_{3}$ lift to isogenies $\tilde{f}_{i}: \mathbb{E} \rightarrow \mathbb{E}^{\prime} \bmod I$.

We write, using the notation of $[\mathbf{R}]$,

$$
\alpha\left(f_{1}, f_{2}, f_{3}\right)=\lg _{W} R / I
$$

By the theorem of Serre-Tate, this global question about elliptic curves can be reduced to a local question about formal groups. This is the reason why we study deformations of isogenies between formal groups in detail in the following chapters.

From $\left[\mathbf{R}\right.$, Theorem 1.1] we get that $\alpha\left(f_{1}, f_{2}, f_{3}\right)$ depends only on the $\mathbb{Z}_{p^{-}}$ isomorphism class of the ternary quadratic form $Q:\left(x_{1}, x_{2}, x_{3}\right) \mapsto \operatorname{deg}\left(\sum x_{i} f_{i}\right)$. We thus write $\alpha_{p}(Q)$ instead of $\alpha\left(f_{1}, f_{2}, f_{3}\right)$. Loc. cit. gives an explicit expression for $\alpha_{p}(Q)$ in terms of the coefficients of $Q$. The number of occurrences of $Q$ in (4.1) is $\frac{1}{8} R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)$ (because we count the isogenies up to $\pm 1$, but the representation number counts each triple $\left(f_{1}, f_{2}, f_{3}\right)$ ). Furthermore, for a positive definite ternary form $Q, R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)=0$ unless $Q$ is isotropic over $\mathbb{Q}_{\ell}$ for all $\ell \neq p$, and anisotropic over $\mathbb{Q}_{p}$. The reason is that $\operatorname{Hom}\left(E, E^{\prime}\right) \otimes \mathbb{Q} \cong \operatorname{End}(E) \otimes \mathbb{Q}$, and $\operatorname{End}(E) \otimes \mathbb{Q}_{\ell} \cong M_{2}\left(\mathbb{Q}_{\ell}\right)$ for $\ell \neq p$, and $\operatorname{End}(E) \otimes \mathbb{Q}_{p}$ is a division algebra (see $[\mathbf{W d 1}$, 2.2]). On the other hand, in the latter case there exists a pair of supersingular elliptic curves $E, E^{\prime}$ in characteristic $p$, such that $Q$ is represented by $\operatorname{Hom}\left(E, E^{\prime}\right)$ (see [Wd1, Proposition 3.2]).

We have now

$$
n(p)=\frac{1}{8} \sum_{\left(E, E^{\prime}\right) \text { supersingular }}\left(\sum_{Q} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{u_{E} u_{E^{\prime}}} \alpha_{p}(Q)\right) .
$$

Further Corollary 4.4 in $[\mathbf{W d 1}]$ states that there are invariants $\beta_{\ell}(Q) \in \mathbb{Z}_{\geq 1}$ which depend only on the isomorphism class of the ternary form $Q$ over $\mathbb{Z}_{\ell}$, such that

$$
\begin{equation*}
\sum_{\left(E, E^{\prime}\right) \text { s.s. }} \frac{R_{\operatorname{Hom}\left(E, E^{\prime}\right)}(Q)}{u_{E} u_{E^{\prime}}}=4 \prod_{\substack{\ell \mid \Delta \\ \ell \neq p}} \beta_{\ell}(Q) . \tag{4.2}
\end{equation*}
$$

The invariants $\beta_{\ell}$ are computed explicitly in [Wd2, Proposition 2.1]. Altogether, we get the following theorem.

Theorem 4.3. - If $T_{m_{1}}, T_{m_{2}}$ and $T_{m_{3}}$ intersect in dimension 0, then

$$
\left(T_{m_{1}} \cdot T_{m_{2}} \cdot T_{m_{3}}\right)_{S}=\log \# \mathbb{Z}\left[j, j^{\prime}\right] /\left(\varphi_{m_{1}}, \varphi_{m_{2}}, \varphi_{m_{3}}\right)=\sum_{p} n(p) \log p
$$

with

$$
n(p)=\frac{1}{2} \sum_{Q}\left(\prod_{\substack{\ell \mid \Delta \\ \ell \neq p}} \beta_{\ell}(Q)\right) \alpha_{p}(Q)
$$

where the sum runs over all positive definite ternary quadratic forms $Q$ over $\mathbb{Z}$ with diagonal $\left(m_{1}, m_{2}, m_{3}\right)$ which are isotropic over $\mathbb{Q}_{\ell}$ for all $\ell \neq p$.

In this way we get a very explicit formula for the intersection numbers.

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[^3]
# 5. THE GENUS OF THE ENDOMORPHISMS OF A SUPERSINGULAR ELLIPTIC CURVE 

## by

Torsten Wedhorn


#### Abstract

We describe the genus of the quadratic space $\operatorname{Hom}\left(E^{\prime}, E\right)$ of homomorphisms of two supersingular elliptic curves $E$ and $E^{\prime}$ and study the map $\left(E^{\prime}, E\right) \mapsto$ $\operatorname{Hom}\left(E^{\prime}, E\right)$ from the set of pairs of supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$ to the set of proper classes in this genus. We show that this map is surjective and determine its fibres. In the last section we use the Minkowski-Siegel formula to express the mean value of the representation of a ternary quadratic form in this genus by local representation densities.


## Résumé (Le genre des endomorphismes d'une courbe elliptique supersingulière)

Nous décrivons le genre de l'espace quadratique $\operatorname{Hom}\left(E^{\prime}, E\right)$ des homomorphismes de deux courbes elliptiques supersingulières $E$ et $E^{\prime}$ et nous étudions l'application $\left(E^{\prime}, E\right) \mapsto \operatorname{Hom}\left(E^{\prime}, E\right)$ de l'ensemble des paires de courbes elliptiques supersingulières sur $\overline{\mathbb{F}}_{p}$ vers l'ensemble des classes propres dans ce genre. Dans le dernier paragraphe, on utilise la formule de Minkowski-Siegel pour exprimer la moyenne de la représentation d'une forme quadratique ternaire dans ce genre en termes de densités de représentation locales.

## Introduction

Let $p>0$ be a prime and let $D$ be the unique quaternion division algebra with center $\mathbb{Q}$ which is ramified precisely at $p$ and at infinity. The reduced norm Nrd is a quadratic form on $D$. We will study lattices and maximal orders in $D$. Recall that two lattices $\Lambda$ and $\Lambda^{\prime}$ are said to be in the same proper class if there exists a $g \in \mathrm{SO}(D, \mathrm{Nrd})$ such that $g \Lambda=\Lambda^{\prime}$.

We will relate the lattices and the maximal orders in $D$ to supersingular elliptic curves. Many of these results, although formulated somewhat differently, can already be found in $[\mathbf{D o}]$ (see also $[\mathbf{G Z}]$ ).

2000 Mathematics Subject Classification. - 11E08, 14K07, 11E12.
Key words and phrases. - Supersingular elliptic curve, quaternion algebra, genus, Minkowski-Siegel formula.

Fix a supersingular elliptic curve $E_{0}$ over $\overline{\mathbb{F}}_{p}$ set $O=\operatorname{End}\left(E_{0}\right)$. Then $O$ is a maximal order in the quadratic space $O \otimes_{\mathbb{Z}} \mathbb{Q}$, where the quadratic form is given by the degree, and we can and will identify the rational quadratic spaces $O \otimes_{\mathbb{Z}} \mathbb{Q}$ with $D$.

The first result is the following (proved in sections 2.9 and 2.15):
Theorem. - Consider isomorphism classes of pairs $(E, \varphi)$ where $E$ is a supersingular elliptic curve over $\overline{\mathbb{F}}_{p}$ and $\varphi: E \rightarrow E_{0}$ is a quasi-isogeny.
(1) The map $(E, \varphi) \mapsto \varphi \operatorname{Hom}\left(E_{0}, E\right)$ induces a bijection of the set of isomorphism classes of supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$ and the set of right ideal classes of $O$.
(2) The map $(E, \varphi) \mapsto \varphi \operatorname{End}(E) \varphi^{-1}$ induces a surjection from the set of isomorphism classes of supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$ to the set of conjugacy classes of maximal orders in $D$. Two supersingular elliptic curves $E$ and $E^{\prime}$ are sent to the same conjugacy class if and only if there exists a $\sigma \in \operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p}\right)$ such that $E^{\prime} \cong E^{(\sigma)}$.

For all pairs $\left(E^{\prime}, E\right)$ of supersingular elliptic curves it is possible to choose quasi-isogenies $\varphi: E \rightarrow E_{0}$ and $\varphi^{\prime}: E^{\prime} \rightarrow E_{0}$ with $\operatorname{deg}(\varphi)=\operatorname{deg}\left(\varphi^{\prime}\right)$. Then $\varphi \operatorname{Hom}\left(E^{\prime}, E\right) \varphi^{\prime-1}$ is a lattice in $D$ whose proper class is independent of the choice of $\varphi$ and $\varphi^{\prime}$. In this way we can consider $\operatorname{Hom}\left(E^{\prime}, E\right)$ as a proper class of lattices in $D$.

The second theorem describes these proper classes (see sections 3.1 and Proposition 3.2).

Theorem. - Let $\Lambda$ be a lattice in $D$. Then the proper class $[\Lambda]$ of $\Lambda$ is the proper class associated to $\operatorname{Hom}\left(E^{\prime}, E\right)$ if and only if $\Lambda$ is in the same genus as $O$.

It follows that the map $\left((E, \varphi),\left(E^{\prime}, \varphi^{\prime}\right)\right) \mapsto \varphi \operatorname{Hom}\left(E^{\prime}, E\right) \varphi^{\prime-1}$ induces a surjection $\left(E, E^{\prime}\right) \mapsto\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]$ from the set of pairs of isomorphism classes of supersingular elliptic curves onto the set of proper classes of lattices in $D$ which are locally isomorphic to $O$. The next theorem describes the fibres of this map and number of automorphisms of the quadratic space $\operatorname{Hom}\left(E, E^{\prime}\right)$ (see Proposition 3.3 and Corollary 3.5 ).

## Theorem

(1) Two pairs $\left(E, E^{\prime}\right)$ and $\left(F, F^{\prime}\right)$ are sent to the same proper class if and only if there exists a $\sigma \in \operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p}\right)$ such that $F=E^{(\sigma)}$ and $F^{\prime}=E^{\prime(\sigma)}$.
(2) For all $\left(E, E^{\prime}\right)$
$\# \operatorname{SO}\left(\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]\right)= \begin{cases}\# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right), & E, E^{\prime} \text { both defined over } \mathbb{F}_{p} ; \\ \frac{1}{2} \# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right), & \text { otherwise } .\end{cases}$

Now fix a positive definite ternary quadratic form $Q$ over $\mathbb{Z}$. By the theorems above we can consider the expression

$$
2\left(\sum_{E} \frac{1}{\# \operatorname{Aut}(E)}\right)^{-2} \sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(Q)}{\# \operatorname{Aut}\left(E^{\prime}\right) \# \operatorname{Aut}(E)}
$$

as the mean value of the representation of $Q$ by the genus of $\operatorname{End}\left(E_{0}\right)$ (here $E$ and $E^{\prime}$ run through all isomorphism classes of supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$, and $R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(Q)$ denotes the number of isometries $\left.Q \rightarrow \operatorname{Hom}\left(E^{\prime}, E\right)\right)$. Hence it can be expressed as a product of local representation densities $\alpha_{l}\left(Q, \operatorname{End}\left(E_{0}\right)\right)$ (see 4.3) by the Minkowski-Siegel formula. We obtain (theorem 4.3):

Theorem. - The mean value is given by

$$
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(Q)}{\# \operatorname{Aut}\left(E^{\prime}\right) \# \operatorname{Aut}(E)}=2\left(\frac{p-1}{12}\right)^{2} \frac{\pi^{4}}{p^{3}} \prod_{l} \alpha_{l}\left(Q, \operatorname{End}\left(E_{0}\right)\right)
$$

where $l$ runs through all prime numbers $l$.

This article is organized as follows. In the first section some definitions and results on quadratic spaces and quaternion algebras are recalled. The second section addresses the correspondence between supersingular elliptic curves, right ideal classes, and conjugacy classes of maximal orders. In the third section the above results on the quadratic spaces $\operatorname{Hom}\left(E^{\prime}, E\right)$ are proved. The Minkowski-Siegel formula is applied in the last section.

Acknowledgements. - I am very grateful to S. Kudla for his helpful remarks and to M. Rapoport, T. Yang and the referee for their comments.

## 1. Preliminaries on quadratic spaces and quaternion algebras

1.1. In this section we recall some definitions and results on quadratic spaces.

If $R$ is a commutative ring, a quadratic space over $R$ is a free $R$-module $M$ together with a map $Q: M \rightarrow R$, such that
(a) $Q(r m)=r^{2} Q(m)$ for all $r \in R$ and $m \in M$.
(b) The form $b_{Q}(x, y)=Q(x+y)-Q(x)-Q(y)$ is $R$-bilinear and nondegenerate (i.e., the $R$-linear map $M \rightarrow M^{*}$ corresponding to $b_{Q}$ is injective).

The map $Q$ is called the quadratic form of the quadratic space $(M, Q)$.
Two quadratic spaces $(M, Q)$ and $\left(M^{\prime}, Q^{\prime}\right)$ over $R$ are said to be isomorphic if there exists an $R$-linear isomorphism $f: M \rightarrow M^{\prime}$ such that $Q^{\prime}(f(m))=Q(m)$ for all $m \in M$. We then write $(M, Q) \cong\left(M^{\prime}, Q^{\prime}\right)$.

The group of automorphisms of a quadratic space will be denoted by $\mathrm{O}(M, Q)$, the subgroup of automorphisms $g \in \mathrm{O}(M, Q)$ with $\operatorname{det}(g)=1$ is denoted by $\mathrm{SO}(M, Q)$.
1.2. In the sequel we will only consider quadratic spaces $(M, Q)$ over integral domains $R$ whose field of fractions has characteristic not equal to 2 . Then we write $\operatorname{Sym}_{n}(R)^{\vee}$ for the set of symmetric matrices $n$ by $n$ matrices $A=\left(a_{i j}\right)$ with coefficients in Quot $(R)$ such that $a_{i i} \in R$ for all $i$ and such that $2 a_{i j} \in R$ for all $i, j$. Moreover, we denote by $B_{Q}$ the $\operatorname{Quot}(R)$-valued bilinear form

$$
B_{Q}: M \times M \longrightarrow \operatorname{Quot}(R), \quad(x, y) \longmapsto \frac{1}{2}(Q(x+y)-Q(x)-Q(y))
$$

Let $\mathcal{B}=\left(e_{1}, \ldots, e_{n}\right)$ be an $R$-basis of $M$. The matrix

$$
S_{Q}=\left(B_{Q}\left(e_{i}, e_{j}\right)\right) \in \operatorname{Sym}_{n}(R)^{\vee}
$$

is called the matrix associated to $(M, Q, \mathcal{B})$.
We denote by $\operatorname{det}(M)=\operatorname{det}((M, Q))$ the class of $\operatorname{det}\left(S_{Q}\right)$ modulo $\left(R^{\times}\right)^{2}$. This is independent of the choice of $\mathcal{B}$.
1.3. Very often we will consider quadratic spaces which arise as follows: Let $(V, Q)$ be a quadratic space over $\mathbb{Q}$ and let $\Lambda$ be a $\mathbb{Z}$-lattice of $V$ (i.e., a finitely generated $\mathbb{Z}$-submodule $\Lambda$ such that $\Lambda \mathbb{Q}=V)$. If $Q(\Lambda) \subset \mathbb{Z}$, the restriction of $Q$ to $\Lambda$ defines a quadratic form on $\Lambda$ over $\mathbb{Z}$.

If $l$ is a finite place of $\mathbb{Q}, \Lambda_{l}=\Lambda \otimes_{\mathbb{Z}} \mathbb{Z}_{l}$ is a lattice in the $\mathbb{Q}_{l}$-vector space $V_{l}=V \otimes \mathbb{Q}_{l}$. Recall that to give a $\mathbb{Z}$-lattice $\Lambda$ in $V$ is the same as to give a $\mathbb{Z}_{l}$-lattice $\Lambda_{l}$ for all $l$ such that there exists a $\mathbb{Z}$-lattice $\Gamma$ of $V$ with $\Lambda_{l}=\Gamma_{l}$ for almost all $l$.

Denote by $\mathbb{A}_{f}$ the ring of finite adeles of $\mathbb{Q}$. An element $g \in G L\left(V \otimes \mathbb{A}_{f}\right)$ is an element $\left(g_{l}\right) \in \prod_{l} G L\left(V_{l}\right)$ where $l$ runs over all finite places of $\mathbb{Q}$ such that $g_{l}\left(\Lambda_{l}\right)=\Lambda_{l}$ for almost all $l$ (this condition is independent of $\Lambda$ ). Hence $g=\left(g_{l}\right)$ acts on the set of lattices by setting

$$
g(\Lambda)=\bigcap_{l}\left(V \cap g_{l}\left(\Lambda_{l}\right)\right)
$$

We obtain an action of $G L\left(V \otimes \mathbb{A}_{f}\right)$ on the set of lattices in $V$ and in particular an action of the subgroups $\mathrm{O}\left(V \otimes \mathbb{A}_{f}\right)$ and $\mathrm{SO}\left(V \otimes \mathbb{A}_{f}\right)$.

Definition 1.1. - We say that two quadratic spaces $M$ and $M^{\prime}$ over $\mathbb{Z}$ are related if $M$ and $M^{\prime}$ are isomorphic over $\mathbb{Z}_{l}$ for all places $l$ of $\mathbb{Q}$ (with the convention $\mathbb{Z}_{\infty}=\mathbb{R}$ ).
1.4. If $M$ and $M^{\prime}$ are related, they are of course also isomorphic over $\mathbb{Q}_{l}$ for all places $l$ and hence they are isomorphic over $\mathbb{Q}$ by the weak approximation theorem for quadratic spaces. If we choose an isomorphism of rational quadratic spaces $M \otimes \mathbb{Q} \cong$ $M^{\prime} \otimes \mathbb{Q}$, we can consider $M$ and $M^{\prime}$ both as lattices in the same quadratic space $V$ over $\mathbb{Q}$. Moreover, the fact that $M$ and $M^{\prime}$ are related just means that there exists a $g \in \mathrm{O}(V)\left(\mathbb{A}_{f}\right)$ with $g(M)=M^{\prime}$. This leads us to the following definition:

Definition 1.2. - Let $V$ be a quadratic space over $\mathbb{Q}$. We say that two lattices $\Lambda$ and $\Lambda^{\prime}$ in $V$ are related if there exists a $g \in \mathrm{O}(V)\left(\mathbb{A}_{f}\right)$ such that $g(\Lambda)=\Lambda^{\prime}$.

An $\mathrm{O}(V)\left(\mathbb{A}_{f}\right)$-orbit of lattices in $V$ is called a genus.

Lemma 1.3. - Let $l$ be a prime number and let $(M, Q)$ be a quadratic space over $\mathbb{Z}_{l}$. Then there exists a reflection in $\mathrm{O}(M, Q)$.

Proof. - Let $x \in M$ be an element such that the $l$-adic valuation of $Q(x)$ is minimal among the elements in $M$. Then an easy calculation shows that the reflection associated to $x$ preserves $M$.

Corollary 1.4. - Let $V$ be a quadratic space over $\mathbb{Q}$. Two lattices $\Lambda$ and $\Lambda^{\prime}$ in $V$ are in the same genus if and only if there exists a $g \in \operatorname{SO}\left(V \otimes \mathbb{A}_{f}\right)$ such that $g(\Lambda)=\Lambda^{\prime}$.

Definition 1.5. - Let $V$ be a quadratic space over $\mathbb{Q}$. Two lattices $\Lambda$ and $\Lambda^{\prime}$ in $V$ are said to be in the same proper class or to be properly equivalent if there exists a $g \in \operatorname{SO}(V)$ such that $g(\Lambda)=\Lambda^{\prime}$.

They are in the same class or equivalent if there exists a $g \in \mathrm{O}(V)$ such that $g(\Lambda)=\Lambda^{\prime}$.

Obviously, every genus of a lattice is the disjoint union of classes and every class is the disjoint union of one or two proper classes. Moreover, it is well known (e.g., [Ki, 6.1.2]) that in each genus there are only finitely many proper classes.

The class of a lattice $\Lambda$ is equal to the proper class of $\Lambda$ if and only if there exists a $g \in \mathrm{O}(V)$ with $\operatorname{det}(g)=-1$ such that $g(\Lambda)=\Lambda$, i.e., if and only if $\mathrm{SO}(\Lambda) \neq \mathrm{O}(\Lambda)$.
1.5. We will be mostly interested in quadratic spaces which arise from quaternion algebras: By a quaternion algebra over a field $F$ we mean a central simple algebra $D$ over $F$ of dimension 4. We write Trd and Nrd for the reduced trace and the reduced norm on D , respectively, and we denote by $x \mapsto \bar{x}:=\operatorname{Trd}(x)-x$ the canonical involution on $D$.

Assume that $F$ is the field of fractions of Dedekind domain $A$ (e.g., $A=\mathbb{Z}$ or $A=\mathbb{Z}_{l}$ ). Let $\Lambda$ be some $A$-lattice of $D$. Then we set

$$
\begin{align*}
& O_{l}(\Lambda)=\{d \in D \mid d \Lambda \subset \Lambda\}  \tag{1.1}\\
& O_{r}(\Lambda)=\{d \in D \mid \Lambda d \subset \Lambda\} \tag{1.2}
\end{align*}
$$

These are orders in $D$. We call them the left order (resp. right order) of $\Lambda$. We say that $\Lambda$ is normal if $O_{l}(\Lambda)$ and $O_{r}(\Lambda)$ are maximal orders.

Lemma 1.6. - Let $F$ be a field with $\operatorname{char}(F) \neq 2$ and let $D$ be a quaternion algebra over $F$. We set

$$
S(D):=\left\{\left(d, d^{\prime}\right) \in D^{\times} \times D^{\times} \mid \operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)\right\}
$$

Consider the group homomorphism

$$
\begin{aligned}
\alpha: S(D) & \longrightarrow \mathrm{O}(D, \mathrm{Nrd}), \\
\left(d, d^{\prime}\right) & \longmapsto\left(\delta \longmapsto d \delta d^{\prime-1}\right) .
\end{aligned}
$$

Then we have

$$
\operatorname{Im}(\alpha)=\operatorname{SO}(D, \operatorname{Nrd}), \quad \operatorname{Ker}(\alpha)=F^{\times}
$$

where $F^{\times}$is embedded diagonally in $D^{\times} \times D^{\times}$.
Proof. - We give two proofs for this. The first is elementary: Clearly, we have $\operatorname{Ker}(\alpha)=F^{\times}$. Let $d \in D$ be an element with $\operatorname{Nrd}(d) \neq 0$ and denote by $\tau_{d}: D \rightarrow D$ the reflection with respect $d$. Then we have for $\delta \in D$ :

$$
\begin{equation*}
\tau_{d} \delta=\delta-\frac{\operatorname{Trd}(\delta \bar{d})}{\operatorname{Nrd}(d)} d=-d \bar{\delta} \bar{\delta}^{-1} \tag{1.3}
\end{equation*}
$$

Every element in $\mathrm{SO}(D)$ is a product of elements of the form $\tau_{d} \tau_{d^{\prime}}$ as $\operatorname{char}(F) \neq 2$. It follows from (1.3) that for all $\delta \in D$ we have

$$
\tau_{d} \tau_{d^{\prime}}(\delta)=d d^{\prime-1} \delta \bar{d}^{\prime} \bar{d}^{-1}
$$

and this proves that $\mathrm{SO}(D)$ is contained in the image of $\alpha$.
Conversely, let $\sigma: \delta \mapsto d \delta d^{\prime-1}$ with $\operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)$ be in the image of $\alpha$. The determinant of left or right multiplication with any element $\tilde{d} \in D$ is given by $\operatorname{Nrd}(\tilde{d})^{2}$ (this can be checked over an algebraic closure and for a matrix algebra this is elementary). Hence we see

$$
\operatorname{det}(\sigma)=\operatorname{Nrd}(d)^{2} \operatorname{Nrd}\left(d^{\prime}\right)^{-2}=1
$$

The second proof is as follows. By Hilbert 90 we have $H^{1}\left(F, \mathbb{G}_{m}\right)=0$ and therefore it suffices to show that $\alpha$ induces an exact sequence of algebraic groups over $F$

$$
1 \longrightarrow \mathbb{G}_{m, F} \longrightarrow S(D) \xrightarrow{\alpha} \mathrm{SO}(D, \mathrm{Nrd}) \longrightarrow 0 .
$$

We can replace $F$ by its algebraic closure. Then it is clear that $S(D)$ is a connected algebraic group of dimension 7 . This implies that $\operatorname{Im}(\alpha)$ must be contained in the conected component of 1 of $\mathrm{O}(D, \mathrm{Nrd})$ which is $\mathrm{SO}(D, \mathrm{Nrd})$. Again it is obvious that $\operatorname{Ker}(\alpha)=\mathbb{G}_{m}$. It follows that $\operatorname{dim}(\operatorname{Im}(\alpha))=6=\operatorname{dim}(\mathrm{SO}(D, \operatorname{Nrd}))$ which shows $\operatorname{Im}(\alpha)=\operatorname{SO}(D, \operatorname{Nrd})$.

Corollary 1.7. - Let $D$ be a quaternion algebra over $\mathbb{Q}$. And let $\Lambda$ and $\Lambda^{\prime}$ be $\mathbb{Z}$-lattices of $D$.
(1) They are in the same genus if and only if there exist $d=\left(d_{l}\right), d^{\prime}=\left(d_{l}^{\prime}\right) \in$ $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$with $\operatorname{Nrd}\left(d_{l}\right)=\operatorname{Nrd}\left(d_{l}^{\prime}\right)$ for all $l$ such that $d \Lambda=\Lambda^{\prime} d^{\prime}$.
(2) They are in the same proper class if and only if there exist $d, d^{\prime} \in D^{\times}$with $\operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)$ such that $d \Lambda=\Lambda^{\prime} d^{\prime}$.
1.6. From now on, $D$ will denote a quaternion algebra over $\mathbb{Q}$. For every place $v$ of $\mathbb{Q}$ we set $D_{v}=D \otimes_{\mathbb{Q}} \mathbb{Q}_{v}$. Then $D_{v}$ is a quaternion algebra over $\mathbb{Q}_{v}$. For all $v$ there are up to isomorphism two quaternion algebras over $\mathbb{Q}_{v}$. One is isomorphic to the ring of matrices $M_{2}\left(\mathbb{Q}_{v}\right)$ and the other one is a quaternion division algebra. If $D_{v} \cong M_{2}\left(\mathbb{Q}_{v}\right)$ we say that $D$ is split at $v$ otherwise $D$ is said to be ramified at $v$.

We know that $D$ is split at almost all places $v$ and that the number of ramified places is even. Conversely, for every set of places $S$ of $\mathbb{Q}$ with an even number of elements there exists a quaternion algebra $D$ over $\mathbb{Q}$ such that $D$ is ramified at $v$ if and only if $v \in S$.
1.7. Let $l$ be a prime number and let $D_{l}$ be a quaternion algebra over $\mathbb{Q}_{l}$. We recall some well-known facts on maximal orders and ideals in $D_{l}$ :

Assume first that $D_{l}=\operatorname{End}(V)$ where $V$ is a two-dimensional $\mathbb{Q}_{l}$-vector space. For every $\mathbb{Z}_{l}$-lattice $L$ in $V$ the ring

$$
\operatorname{End}(L)=\{d \in \operatorname{End}(V) \mid d(L) \subset L\}
$$

is a maximal order of $\operatorname{End}(V)$. Conversely, every maximal order $O$ is of this form. As $G L(V)=D^{\times}$acts transitively on the set of all lattices in $V$, we see that all maximal orders of $D$ are conjugate.

If we choose a basis for $L$, the Cartan decomposition can be written as

$$
G L_{2}\left(\mathbb{Q}_{l}\right)=G L_{2}\left(\mathbb{Z}_{l}\right) \cdot T \cdot G L_{2}\left(\mathbb{Z}_{l}\right)
$$

where $T$ consists of the diagonal matrices of the form $\operatorname{diag}\left(l^{a}, l^{b}\right)$ for $a, b \in \mathbb{Z}$. Using this decomposition, an easy calculation shows that the normalizer of a maximal order $O=\operatorname{End}(L)$ in $D_{l}^{\times}$is given by

$$
\begin{equation*}
N_{D_{l}^{\times}}(O)=l^{\mathbb{Z}} O^{\times} . \tag{1.4}
\end{equation*}
$$

A $\mathbb{Z}_{l}$-lattice $\Lambda$ is normal (see 1.5) if and only if there exist lattices $L$ and $L^{\prime}$ in $V$ such that

$$
\Lambda=\operatorname{Hom}\left(L^{\prime}, L\right)=\left\{d \in D_{l} \mid d\left(L^{\prime}\right) \subset L\right\}
$$

Conversely, $\operatorname{Hom}\left(L^{\prime}, L\right)$ is clearly a normal lattice of $\operatorname{End}(V)$ for all lattices $L, L^{\prime}$ of $V$. We have

$$
O_{l}\left(\operatorname{Hom}\left(L^{\prime}, L\right)\right)=\operatorname{End}(L), \quad O_{r}\left(\operatorname{Hom}\left(L^{\prime}, L\right)\right)=\operatorname{End}\left(L^{\prime}\right)
$$

and as a left $\operatorname{End}(L)$-module (resp. as a right $\operatorname{End}\left(L^{\prime}\right)$-module) $\operatorname{Hom}\left(L^{\prime}, L\right)$ is generated by any one element $d$ such that $d\left(L^{\prime}\right)=L$.
1.8. Now assume that $D_{l}$ is a quaternion division algebra over $\mathbb{Q}_{l}$. Then there exists a unique maximal order $\mathcal{O}_{D_{l}}$ of $D_{l}$, namely

$$
\mathcal{O}_{D_{l}}=\left\{d \in D_{l} \mid \operatorname{Nrd}(d) \in \mathbb{Z}_{l}\right\}
$$

Moreover, $\mathcal{O}_{D_{l}}$ has a unique maximal ideal $\mathfrak{m}$ which is a principal ideal. Every nonzero one-sided ideal of $\mathcal{O}_{D_{l}}$ is a power of $\mathfrak{m}$, in particular it is a two-sided ideal.

As $d \mathcal{O}_{D_{l}} d^{-1}$ is again a maximal order of $D_{l}$ for all $d \in D_{l}^{\times}$, we see that

$$
\begin{equation*}
N_{D_{l}^{\times}}\left(\mathcal{O}_{D_{l}}\right)=D_{l}^{\times} . \tag{1.5}
\end{equation*}
$$

## 2. Supersingular elliptic curves

2.1. From now on we fix a prime number $p$. We consider supersingular elliptic curves $E$ over $\overline{\mathbb{F}}_{p}$. Recall that any supersingular elliptic curve is already defined over $\mathbb{F}_{p^{2}}$. For two supersingular elliptic curves we denote by $\operatorname{Hom}\left(E^{\prime}, E\right)$ the set of homomorphism $E^{\prime} \rightarrow E$ which are defined over $\overline{\mathbb{F}}_{p}$. We set $\operatorname{End}(E)=\operatorname{Hom}(E, E)$.

We denote by $W\left(\overline{\mathbb{F}}_{p}\right)$ the ring of Witt vectors of $\overline{\mathbb{F}}_{p}$ and write $\sigma$ for the Frobenius on $W\left(\overline{\mathbb{F}}_{p}\right)$.
2.2. For any prime $l \neq p$ let $T_{l}(E)$ be the Tate module. It is a free $\mathbb{Z}_{l}$-module of rank 2 . For $l=p$ we denote by $T_{p}(E)$ the (covariant) Dieudonné module of $E$. It as a free $W\left(\overline{\mathbb{F}}_{p}\right)$-module of rank 2 with $\sigma$-linear operator $\Phi$ such that

$$
p T_{p}(E) \subsetneq \Phi\left(T_{p}(E)\right) \subsetneq T_{p}(E)
$$

where $\sigma$ is the Frobenius in $W\left(\overline{\mathbb{F}}_{p}\right)$. In fact, there exists a $W\left(\overline{\mathbb{F}}_{p}\right)$-basis $(e, f)$ of $T_{p}(E)$ such that $\Phi(e)=f$ and $\Phi(f)=p e$.

We denote by $\operatorname{Hom}\left(T_{p}\left(E^{\prime}\right), T_{p}(E)\right)$ the $\mathbb{Z}_{p}$-module of $W\left(\overline{\mathbb{F}}_{p}\right)$-linear homomorphisms $T_{p}\left(E^{\prime}\right) \rightarrow T_{p}(E)$ which commute with $\Phi$. It is easily checked that this is a free $\mathbb{Z}_{p^{-}}$ module of rank 4. Moreover, $\operatorname{End}\left(T_{p}(E)\right) \otimes_{\mathbb{Z}_{p}} \mathbb{Q}_{p}$ is "the" quaternion division algebra over $\mathbb{Q}_{p}$ and $\operatorname{End}\left(T_{p}(E)\right)$ is its maximal order.

We set $T^{p}(E)=\prod_{l \neq p} T_{l}(E), V^{p}(E)=T^{p}(E) \otimes_{\mathbb{Z}} \mathbb{Q}, V_{p}(E)=T_{p}(E) \otimes_{\mathbb{Z}} \mathbb{Q}$ and $V(E)=V^{p}(E) \times V_{p}(E)$.

These are free modules of rank 2 over the rings $\mathbb{Z}^{p}=\prod_{l \neq p} \mathbb{Z}_{l}, \mathbb{A}_{f}^{p}=\mathbb{Z}^{p} \otimes \mathbb{Q}, \mathbb{Q}_{p}$, and $\mathbb{A}_{f}$, respectively.
2.3. We fix a supersingular elliptic curve $E_{0}$, set $\mathcal{O}_{D}=\operatorname{End}\left(E_{0}\right)$ and $D=\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Q}$. It follows from 2.2 (see also 3.1 below) that $D$ is a quaternion division algebra over $\mathbb{Q}$ which is ramified precisely at $p$ and $\infty$. As $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Z}_{l}$ is a maximal order in $D_{l}=D \otimes_{\mathbb{Q}} \mathbb{Q}_{l}$ for all primes $l, \mathcal{O}_{D}$ is a maximal order in $D$.

The reduced norm Nrd is a positive definite quadratic form on $D$ and the induced homomorphism Nrd: $D^{\times} \rightarrow \mathbb{Q}^{>0}$ is surjective.
2.4. Denote by $Y^{p}$ be the set of $\widehat{\mathbb{Z}}^{p}$-lattices in $V^{p}\left(E_{0}\right)$ and by $Y_{p}$ the set of $W\left(\overline{\mathbb{F}}_{p}\right)$ lattices $L$ in $V_{p}\left(E_{0}\right)$ such that $p L \subsetneq \Phi(L) \subsetneq L$. Then $Y^{p} \times Y_{p}$ describes quasi-isogenies with target $E_{0}$ as follows: Consider pairs $(E, \varphi)$ where $E$ is a supersingular elliptic curve and $\varphi: E \rightarrow E_{0}$ a quasi-isogeny. We call two such pairs $(E, \varphi)$ and $\left(E^{\prime}, \varphi^{\prime}\right)$ equivalent if there exists a commutative diagram

where the vertical arrows are isomorphisms. Then $Y^{p} \times Y_{p}$ corresponds to the set $\mathcal{L}$ of equivalence classes of such pairs $(E, \varphi)$ as above (see $[\mathbf{M i}, 6]$ for the general description of an isogeny class of an abelian variety with endomorphisms).

The group $G L\left(V^{p}\left(E_{0}\right)\right)$ acts transitively on the elements in $Y^{p}$ and therefore we have a bijection

$$
Y^{p} \longleftrightarrow G L\left(V^{p}\left(E_{0}\right)\right) / G L\left(T^{p}\left(E_{0}\right)\right) .
$$

Moreover, if we denote by $\operatorname{Aut}\left(V_{p}\left(E_{0}\right)\right)$ the automorphisms of $\left.V_{p}\left(E_{0}\right)\right)$ which commute with $\Phi$, it follows from the existence of a normal form for lattices in $Y_{p}$ (see 2.2) that $\operatorname{Aut}\left(V_{p}\left(E_{0}\right)\right)$ acts transitively on $Y_{p}$. Therefore we have a bijection

$$
Y_{p} \longleftrightarrow \operatorname{Aut}\left(V_{p}\left(E_{0}\right)\right) / \operatorname{Aut}\left(T_{p}\left(E_{0}\right)\right) .
$$

If we choose isomorphisms $\alpha_{l}: \operatorname{End}\left(T_{l}\left(E_{0}\right)\right) \xrightarrow{\sim} \mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Z}_{l}$ for all primes $l$, we obtain a bijection

$$
\begin{equation*}
\mathcal{L} \longleftrightarrow\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times} \tag{2.1}
\end{equation*}
$$

which is independent of the choice of the $\alpha_{l}$ by the theorem of Skolem-Noether. Explicitly this bijection is given by the associating to $d=\left(d_{l}\right) \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$the equivalence class of the pair $(E, \varphi)$ such that $T_{l}(\varphi)\left(T_{l}(E)\right)=d_{l} T_{l}\left(E_{0}\right)$ for all primes $l$.
2.5. Let $d \in\left(D \otimes \mathbb{Q} \mathbb{A}_{f}\right)^{\times}$and let $[(E, \varphi)] \in \mathcal{L}$ be the associated quasi-isogeny. Then we have

$$
\operatorname{deg}(\varphi)=\prod_{l} l^{v_{l}\left(\operatorname{Nrd}\left(d_{l}\right)\right)}
$$

where $l$ runs over all prime numbers.
2.6. For example, the relative Frobenius $E_{0}^{(p)} \rightarrow E_{0}$ corresponds via the bijection (2.1) to the class of an element $\Pi$ in $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$which has a uniformizing element of $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Z}_{p}$ as $p$-th component and a unit of $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Z}_{l}$ as $l$-th component for all primes $l \neq p$.

More generally, if $O$ is any maximal order of $D$, we call an element $\Pi=\left(\Pi_{l}\right)_{l} \in$ $D \otimes_{\mathbb{Q}} \mathbb{A}_{f}$ a Frobenius element in a quaternion algebra with respect to $O$ if $\Pi_{l} \in\left(O \otimes_{\mathbb{Z}}\right.$ $\left.\mathbb{Z}_{l}\right)^{\times}$for all $l \neq p$ and $\Pi_{p}$ is a uniformizing element of $\mathcal{O}_{D_{p}}$.
2.7. Consider the natural map

$$
\begin{aligned}
\mathcal{L} & \longrightarrow \mathcal{I}:=\left\{\begin{array}{c}
\text { isomorphism classes of } \\
\text { supersingular elliptic curves over } \overline{\mathbb{F}}_{p}
\end{array}\right\}, \\
{[(E, \varphi)] } & \longmapsto E .
\end{aligned}
$$

Using the identification (2.1), two elements $d, d^{\prime} \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}$have the same image in $\mathcal{I}$ if and only if there exists a $\delta \in D^{\times}=\left(\operatorname{End}\left(E_{0}\right) \otimes_{\mathbb{Z}} \mathbb{Q}\right)^{\times}$such that $\delta d=d^{\prime}$. Hence we get:

Theorem 2.1. - There is a natural identification

$$
D^{\times} \backslash\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times} \longleftrightarrow\left\{\begin{array}{c}
\text { isomorphism classes of }  \tag{2.2}\\
\text { supersingular elliptic } \\
\text { curves over } \overline{\mathbb{F}}_{p}
\end{array}\right\}
$$

2.8. Let $O$ be any order in $D$. A $\mathbb{Z}$-lattice $\Lambda$ of $D$ is called a right ideal of $O$ if $O \subset O_{r}(\Lambda)$ (cf. 1.5). If $O$ is a maximal order, this is of course equivalent to $O=O_{r}(\Lambda)$. Two right ideals $\Lambda$ and $\Lambda^{\prime}$ of $O$ are said to be in the same right ideal class if there exists a $d \in D^{\times}$with $d \Lambda=\Lambda^{\prime}$.

Let $\Lambda$ be a lattice in $D$. It can be easily checked locally that $O_{l}(\Lambda)$ is a maximal order of $D$ if and only if $O_{r}(\Lambda)$ is a maximal order. Hence all right ideals of our fixed maximal order $\mathcal{O}_{D}$ are normal lattices in the sense of 1.5 .

By 1.7 and 1.8 we know that locally all right ideals of $\mathcal{O}_{D}$ are principal ideals. Hence it follows that for every right ideal $\Lambda$ of $\mathcal{O}_{D}$ there exists a $d \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$such that $\Lambda=d \mathcal{O}_{D}$. Therefore $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$acts transitively on the set of all right ideals of $\mathcal{O}_{D}$. Moreover, an easy local calculation shows that the stabilizer of the right ideal $\mathcal{O}_{D}$ in $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$is equal to $\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}$. Thus we have a natural bijection

$$
\begin{equation*}
\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times} \longleftrightarrow\left\{\text { right ideals of } \mathcal{O}_{D}\right\} . \tag{2.3}
\end{equation*}
$$

Composing this bijection with (2.1) we get a bijection of $\mathcal{L}$ with the set of right ideals of $\mathcal{O}_{D}$. Explicitly, this associates to each equivalence class $[(E, \varphi)]$ the right ideal $\varphi \operatorname{Hom}\left(E_{0}, E\right)$ of $\mathcal{O}_{D}=\operatorname{End}\left(E_{0}\right)$.
2.9. The bijection (2.3) induces a bijection of the set of right ideal classes of $\mathcal{O}_{D}$ with $D^{\times} \backslash\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}$. Composing this with (2.2) we obtain:

Proposition 2.2. - There exists a natural bijection

$$
\left\{\begin{array}{c}
\text { right ideal classes }  \tag{2.4}\\
\text { of } \mathcal{O}_{D}
\end{array}\right\} \longleftrightarrow\left\{\begin{array}{c}
\text { isomorphism classes of } \\
\text { supersingular elliptic } \\
\text { curves over } \\
\overline{\mathbb{F}}_{p}
\end{array}\right\} .
$$

2.10. By 1.7 and 1.8 we know that locally all maximal orders in $D$ are conjugate to each other. Therefore $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$acts transitively by conjugation on the set of maximal orders in $D$. The stabilizer of the maximal order $\mathcal{O}_{D}$ is the normalizer of $\mathcal{O}_{D}$ in $\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$. It can be computed locally and by (1.4) and (1.5) we have

$$
\begin{equation*}
N_{\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}}\left(\mathcal{O}_{D}\right)=\mathbb{Q}^{\times}\left(D_{p}^{\times} \times\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}^{p}\right)^{\times}\right) . \tag{2.5}
\end{equation*}
$$

2.11. Let $O$ be a maximal order of $D$ and $d \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$such that $O=d \mathcal{O}_{D} d^{-1}$. Then it follows at once from the definition of a Frobenius element in 2.6 that if $\Pi \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$is a Frobenius element with respect to $\mathcal{O}_{D}, d \Pi d^{-1}$ is a Frobenius element with respect to $O$.
2.12. For $d \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$let $[(E, \varphi)]$ be the associated element in $\mathcal{L}$ via the bijection (2.1). Let $\Pi$ be a Frobenius element with respect to $\mathcal{O}_{D}$. Then the pair $\left[\left(E^{\prime}, \varphi^{\prime}\right)\right]$ associated to $d \Pi$ is given by $E^{\prime}=E^{\left(p^{-1}\right)}\left(=E^{(p)}\right)$ and $\varphi^{\prime}=F_{E_{0}} \circ \varphi^{\left(p^{-1}\right)}$ where $F_{E_{0}}: E_{0}^{\left(p^{-1}\right)} \rightarrow E_{0}$ is the relative Frobenius.

Moreover, if $\Lambda=d \mathcal{O}_{D}$ is the right ideal of $\mathcal{O}_{D}$ corresponding to $[(E, \varphi)]$ via the bijection (2.3), the right ideal corresponding to $\left[\left(E^{\prime}, \varphi^{\prime}\right)\right]$ is given by $d \Pi \mathcal{O}_{D}=\tilde{\Pi} d \mathcal{O}_{D}=$ $\tilde{\Pi} \Lambda$ where $\tilde{\Pi}=d \Pi d^{-1}$ is a Frobenius element with respect to $O=O_{l}(\Lambda)$.
2.13. Let $\mathfrak{m}_{p} \subset \mathcal{O}_{D_{p}}$ be the maximal ideal. For a maximal order $O$ of $D$ let $\mathfrak{p}=O \cap \mathfrak{m}_{p}$ be the unique prime ideal of $O$ which lies over $p$. Let $\Lambda$ be any left ideal of $O$ and let $\Pi$ be a Frobenius element with respect to $O$. Then arguing locally one sees that

$$
\Pi \Lambda=\mathfrak{p} \Lambda
$$

2.14. It is easy to check that the canonical projection

$$
\begin{align*}
\left\{\text { right ideals of } \mathcal{O}_{D}\right\} & =\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times} \\
\longrightarrow\{\text { maximal orders of } D\} & =\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} / N_{\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)} \times\left(\mathcal{O}_{D}\right) \tag{2.6}
\end{align*}
$$

is given by $\Lambda \mapsto O_{l}(\Lambda)$.
The projection (2.6) induces a map from the set of right ideal classes of $\mathcal{O}_{D}$ to the set of $D^{\times}$-conjugacy classes of maximal orders in $D$ whose adelic version is the projection

$$
\begin{align*}
D^{\times} \backslash\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times} \longrightarrow & D^{\times} \backslash\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} / N_{\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}}\left(\mathcal{O}_{D}\right)  \tag{2.7}\\
& =D^{\times} \backslash\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(D_{p}^{\times} \times\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}^{p}\right)^{\times}\right) .
\end{align*}
$$

Lemma 2.3. - Let $\{O\}$ be a conjugacy class of maximal orders in D. Let $\mathcal{R}(\{O\})$ be the set of classes of right ideals of $\mathcal{O}_{D}$ which are sent to $\{O\}$ under the map (2.7). Then $\mathcal{R}(\{O\})$ consists either of one or two elements. It consists of one element if and only if $O$ contains an element $d$ with $\operatorname{Nrd}(d)=p$ for one (or equivalently for all) $O \in\{O\}$.

Proof. - We consider the map (2.7). Hence we are in the following situation: Let $G$ be a group, $H$ and $G_{2}$ be subgroups, and let $G_{1}$ be a normal subgroup of $G_{2}$. Consider the canonical projection

$$
\varpi: H \backslash G / G_{1} \longrightarrow H \backslash G / G_{2}
$$

Let $g_{0} \in G$ and set $F_{g_{0}}=\varpi^{-1}\left(\varpi\left(H g_{0} G_{1}\right)\right)$. Then $G_{1} \backslash G_{2}$ acts transitively from the right on $F_{g_{0}}$ and the stabilizer of $H g_{0} G_{1}$ is $\left(g_{0}^{-1} H g_{0} \cap G_{2}\right) G_{1} / G_{1}$.

In the special case of $(2.7)$ we have $H=D^{\times}, G=\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}, G_{1}=\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}$ and $G_{2}=D_{p}^{\times} \times\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}^{p}\right)^{\times}$. Therefore $G_{2} / G_{1}$ is the free cyclic group generated by a Frobenius element $\Pi$ (cf. 2.6). Moreover, $\Pi^{2} G_{1} \subset \mathbb{Q}^{\times} G_{1} \subset\left(g_{0}^{-1} H g_{0} \cap G_{2}\right) G_{1} / G_{1}$ for all $g_{0} \in G$. Therefore the fibres of (2.7) consist of at most two elements.

Now fix $g_{0}$ and let $O=g_{0} \mathcal{O}_{D} g_{0}^{-1}$ be the associated maximal order of $D$. A fibre $F_{g_{0}}$ consists of one element if and only if there exists in $g_{0}^{-1} D^{\times} g_{0}$ a Frobenius element. If there exists a $d \in D^{\times}$such that $g_{0}^{-1} d g_{0}$ is a Frobenius element (and hence an element of $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}$ ), we necessarily have $d \in g_{0} \mathcal{O}_{D} g_{0}^{-1}=O$. Moreover, $g_{0}^{-1} d g_{0}$ is a Frobenius element if and only if

$$
\operatorname{Nrd}\left(g_{0}^{-1} d g_{0}\right)_{l} \in \begin{cases}\mathbb{Z}_{l}^{\times}, & \text {if } l \neq p \\ p \mathbb{Z}_{p}, & \text { if } l=p\end{cases}
$$

Hence we see that $F_{g_{0}}$ consists of one element if and only if there exists an element $d \in O$ such that $\operatorname{Nrd}(d)=p$ (the case $\operatorname{Nrd}(d)=-p$ can of course not occur).
2.15. In 2.8 we have seen that every right ideal of $\mathcal{O}_{D}$ is of the form $\varphi \operatorname{Hom}\left(E_{0}, E\right)$ for some quasi-isogeny $\varphi: E \rightarrow E_{0}$. As we have

$$
O_{l}\left(\varphi \operatorname{Hom}\left(E_{0}, E\right)\right)=\varphi \operatorname{End}(E) \varphi^{-1}
$$

we see that every maximal order of $D$ is of the form $\varphi \operatorname{End}(E) \varphi^{-1}$. Moreover, the $D^{\times}$-conjugacy class of $\varphi \operatorname{End}(E) \varphi^{-1}$ depends only on $E$. We denote it by $\{\operatorname{End}(E)\}$.

Proposition 2.4. - Let $O$ be a maximal order and let $E$ be a supersingular elliptic curve such that $O$ is in the conjugacy class $\{\operatorname{End}(E)\}$. Let $\mathfrak{p}$ be the unique (twosided) prime ideal of $O$ which lies over $p$.
(1) The following assertions are equivalent:
(a) The elliptic curve $E$ is defined over $\mathbb{F}_{p}$.
(b) Up to isomorphism there exists a unique supersingular elliptic curve $E$ such that the conjugacy class of $O$ is equal to $\{\operatorname{End}(E)\}$.
(c) The prime ideal $\mathfrak{p}$ is a principal ideal.
(d) The subgroup $\mathbb{Q}^{\times} O^{\times}$of the normalizer $N_{D \times}(O)$ is of index 2 .
(2) If these equivalent conditions do not hold, we have:
(a) Up to isomorphism there are precisely two elliptic curves $E$ and $E^{\prime}$ such that $\{O\}=\{\operatorname{End}(E)\}=\left\{\operatorname{End}\left(E^{\prime}\right)\right\}$ and for them $E^{\prime} \cong E^{(p)}$.
(b) $N_{D^{\times}}(O)=\mathbb{Q}^{\times} O^{\times}$.

Proof. - It follows from Proposition 2.2 and Lemma 2.3 that (1)(b) is equivalent to the existence of a $d \in O$ such that $\operatorname{Nrd}(d)=p$. As $\mathfrak{p}=\mathfrak{m}_{p} \cap O$ such an element generates $\mathfrak{p}$. Conversely, for every generator $d$ of $\mathfrak{p}$ we have $\operatorname{Nrd}(d)=p$. This proves the equivalence of (1)(b) and (1)(c).

An easy calculation shows that $d \in N_{D \times}(O)$ implies $d \in l^{\mathbb{Z}}\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right)^{\times}$for all primes $l \neq p$. Therefore we have an injective homomorphism

$$
\nu: N_{D \times}(O) /\left(\mathbb{Q}^{\times} O^{\times}\right) \longleftrightarrow\left(D \otimes_{\mathbb{Q}} \mathbb{Q}_{p}\right)^{\times} / \mathbb{Q}^{\times}\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{p}\right)^{\times} \cong \mathbb{Z} / 2 \mathbb{Z}
$$

where the isomorphism is given by $v_{p} \circ$ Nrd. In particular, we see that if (1)(d) does not hold, (2)(b) holds.

The homomorphism $\nu$ is surjective if and only if there exists an element $d \in$ $N_{D^{\times}}(O)$ such that $v_{p}(\operatorname{Nrd}(d))=1$. For all $d \in N_{D^{\times}}(O)$ we have $v_{l}(\operatorname{Nrd}(d)) \in 2 \mathbb{Z}$ for all $l \neq p$. Therefore, if $d \in N_{D^{\times}}(O)$ satisfies $v_{p}(\operatorname{Nrd}(d))=1$, we can find a $\lambda \in \mathbb{Q}^{\times}$ such that $\operatorname{Nrd}(d \lambda)=p$. As $d \lambda \in N_{D^{\times}}(O)$, we can write $d \lambda=l^{m} a$ for all $l \neq p$ where $m \in \mathbb{Z}$ and $a \in\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right)^{\times}$. As $v_{l}(\operatorname{Nrd}(d \lambda))=0$, we have $d \lambda \in\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right)^{\times}$. Moreover, $\operatorname{Nrd}(d \lambda)=p$ also implies that $d \lambda \in O \otimes_{\mathbb{Z}} \mathbb{Z}_{p}=\mathcal{O}_{D_{p}}$. Hence $d \lambda \in O$. Altogether we have seen that $\nu$ is surjective if and only if there exists a $d \in O \cap N_{D^{\times}}(O)$ with $\operatorname{Nrd}(d)=p$. Such an element generates $\mathfrak{p}$ and thus (1)(d) implies (1)(c).

Conversely, for every $d \in O$ with $\operatorname{Nrd}(d) \in p^{\mathbb{Z}}$ we have $d\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right) d^{-1}=\left(O \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right)$ for all primes $l$ and hence $d \in N_{D \times}(O)$. Therefore the converse implication does also hold.

Next we show that for any supersingular elliptic curve $E$ we have

$$
\{\operatorname{End}(E)\}=\left\{\operatorname{End}\left(E^{\left(p^{-1}\right)}\right)\right\}
$$

Choose a quasi-isogeny $\varphi: E \rightarrow E_{0}$ and let $I=\varphi \operatorname{Hom}\left(E_{0}, E\right)$ be the corresponding right ideal of $\mathcal{O}_{D}$. The right ideal corresponding to the quasi-isogeny $E_{0}^{\left(p^{-1}\right)} \rightarrow E_{0}$ $\rightarrow^{\varphi} E$ (where the first arrow is the relative Frobenius) is the right ideal $\Pi I$ where $\Pi \in$ $\left(D \otimes \mathbb{Q} \mathbb{A}_{f}\right)^{\times}$is a Frobenius element with respect to the maximal order $O_{l}(I)$ (cf. 2.6). Using a local calculation it follows at once that $O_{l}(\Pi I)=\Pi O_{l}(I) \Pi^{-1}=O_{l}(I)$.

We have already seen in the proof of Lemma 2.3 that if $I=d \mathcal{O}_{D}$ and $I^{\prime}=d^{\prime} \mathcal{O}_{D}$ (with $d$, $d^{\prime} \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$) are two right ideals of $\mathcal{O}_{D}$ with $\left\{O_{l}(I)\right\}=\left\{O_{l}\left(I^{\prime}\right)\right\}$ then there exists a Frobenius element $\Pi$ with respect to $\mathcal{O}_{D}$ and an integer $n$ such that the right ideal class of $d^{\prime} \mathcal{O}_{D}$ is equal to the right ideal class of $d \Pi^{n} \mathcal{O}_{D}$. Writing $d \Pi=\tilde{\Pi} d$ for a Frobenius element $\tilde{\Pi}$ with respect to $O_{l}(I)$ (see 2.11), we see that the right ideal classes of $I^{\prime}$ and $\tilde{\Pi}^{n} I$ are equal. Let $E$ be the supersingular elliptic curve corresponding to the class of $I$ and let $E^{\prime}$ be the supersingular elliptic curve corresponding to the class of $I^{\prime}$. Then this implies that there exists an integer $n$ such that $E^{\prime} \cong E^{\left(p^{n}\right)}$ and this completes the proof.
2.16. Note that the proof of Proposition 2.4 also shows that every supersingular elliptic curve is already defined over $\mathbb{F}_{p^{2}}$.

Moreover we have seen:
Corollary 2.5. - Let $O$ be a maximal order of $D$ and set $N:=N_{D \times}(O)$. Consider the subgroups $O^{\times} \subset \mathbb{Q}^{\times} O^{\times} \subset N$. Then

$$
\begin{aligned}
O^{\times} & =\{d \in N \mid \operatorname{Nrd}(d)=1\} \\
\mathbb{Q}^{\times} O^{\times} & =\left\{d \in N \mid \operatorname{Nrd}(d) \in\left(\mathbb{Q}^{\times}\right)^{2}\right\}
\end{aligned}
$$

## 3. The genus of the quadratic space $\operatorname{Hom}\left(E^{\prime}, E\right)$

3.1. For any two supersingular elliptic curves $E$ and $E^{\prime}$ we will consider the free $\mathbb{Z}$-module $\operatorname{Hom}\left(E^{\prime}, E\right)$ together with the quadratic form given by the degree.

As $E$ and $E^{\prime}$ are supersingular, $\operatorname{Hom}\left(E^{\prime}, E\right)$ has rank 4 as a $\mathbb{Z}$-module. We have a canonical map of $\mathbb{Z}_{l}$-modules

$$
\alpha: \operatorname{Hom}\left(E^{\prime}, E\right) \otimes_{\mathbb{Z}} \mathbb{Z}_{l} \longrightarrow \operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right)
$$

As the union of the $l^{n}$-torsion $E\left[l^{n}\right]$ for $n \geq 1$ is scheme-theoretically dense in $E, \alpha$ is injective. Moreover, the cokernel of $\alpha$ is torsionfree: Indeed, let $\tau \in \operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right)$ be such that $l \tau=\alpha(\varphi)$ for some $\varphi \in \operatorname{Hom}\left(E^{\prime}, E\right) \otimes \mathbb{Z}_{l}$. We write $\varphi$ as the limit of sequence of $\varphi_{n} \in \operatorname{Hom}\left(E^{\prime}, E\right)$ converging to $\varphi$. For large $n$ we have $T_{l}\left(\varphi_{n}\right)=l \tau_{n}$ for some $\tau_{n} \in \operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right.$, and therefore the restriction of $\varphi_{n}$ to the $l$-torsion $E[l]$ is zero. But this implies that $\varphi$ is divisible by $l$.

As both sides have rank 4 , it follows that $\alpha$ is an isomorphism. Choosing an identification $T_{l}\left(E^{\prime}\right) \cong T_{l}(E)$, we can consider the right hand side as a lattice in the quaternion algebra $\operatorname{End}\left(T_{l}(E)\right) \otimes_{\mathbb{Z}_{l}} \mathbb{Q}_{l}$ and the quadratic form given by the degree on the left hand side corresponds via $\alpha$ to the reduced norm on the right hand side.

Therefore the isomorphism class of the quadratic space $\operatorname{Hom}\left(E^{\prime}, E\right) \otimes_{\mathbb{Z}} \mathbb{Z}_{l}$ is independent of $E$ and $E^{\prime}$ for all $l$. In other words (Definition 1.1), $\operatorname{Hom}\left(E^{\prime}, E\right)$ and $\operatorname{Hom}\left(F^{\prime}, F\right)$ are related for all supersingular elliptic curves $E, E^{\prime}, F$, and $F^{\prime}$.

Lemma 3.1. - Let $E$ and $E^{\prime}$ be two supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. Then there exists a quasi-isogeny $\varphi: E^{\prime} \rightarrow E$ of degree 1.

Proof. - As $\operatorname{Hom}\left(E^{\prime}, E\right)$ and $\operatorname{End}(E)$ are related, we can choose an isomorphism of quadratic spaces $\operatorname{Hom}\left(E^{\prime}, E\right) \otimes_{\mathbb{Z}} \mathbb{Q} \cong \operatorname{End}(E) \otimes_{\mathbb{Z}} \mathbb{Q}$. Via this isomorphism we identify $\operatorname{Hom}\left(E^{\prime}, E\right)$ with a sublattice of $\operatorname{End}(E) \otimes_{\mathbb{Z}} \mathbb{Q}$. Choose an integer $N \geq 1$ such that $N \operatorname{End}(E) \subset \operatorname{Hom}\left(E^{\prime}, E\right)$. Then $N \operatorname{id}_{E}$ corresponds to an isogeny $\varphi^{\prime}: E^{\prime} \rightarrow E$ of degree $N^{2}$ and $\varphi=(1 / N) \varphi^{\prime}$ is a quasi-isogeny of degree 1 .
3.2. If $E$ and $E^{\prime}$ are two supersingular elliptic curves, we can choose quasi-isogenies $\varphi: E \rightarrow E_{0}$ and $\varphi^{\prime}: E^{\prime} \rightarrow E_{0}$ such that $\operatorname{deg}(\varphi)=\operatorname{deg}\left(\varphi^{\prime}\right)$ by Lemma 3.1. Then

$$
\iota_{\varphi, \varphi^{\prime}}: \operatorname{Hom}\left(E^{\prime}, E\right) \longrightarrow D, \quad \alpha \longmapsto \varphi \circ \alpha \circ \varphi^{\prime-1}
$$

is an isometry. If we choose another pair $\left(\varphi_{1}, \varphi_{1}^{\prime}\right)$ as above, we have

$$
\iota_{\varphi, \varphi^{\prime}}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)=d \iota_{\varphi_{1}, \varphi_{1}^{\prime}}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right) d^{\prime-1}
$$

for $d, d^{\prime} \in D$ with $\operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)$. Hence it follows from Corollary 1.7 that the proper class of $\iota_{\varphi, \varphi^{\prime}}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)$ is independent of the choice of $\left(\varphi, \varphi^{\prime}\right)$. We denote this class by $\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]$.

Proposition 3.2. - Every proper class in the genus of $\operatorname{End}\left(E_{0}\right) \subset D$ is of the form $\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]$ for two supersingular elliptic curves $E$ and $E^{\prime}$.

Proof. - Let $\Lambda \subset D$ be a lattice in the genus of $\mathcal{O}_{D}$. By Corollary 1.7 there exist $d, d^{\prime} \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$with $\operatorname{Nrd}\left(d_{l}\right)=\operatorname{Nrd}\left(d_{l}^{\prime}\right)$ for all primes $l$ such that $\Lambda=d \mathcal{O}_{D} d^{\prime-1}$. Denote by $[(E, \varphi)],\left[\left(E^{\prime}, \varphi^{\prime}\right)\right] \in \mathcal{L}$ the quasi-isogenies associated to $d$, $d^{\prime}$, respectively, via bijection 2.1. Then $\varphi \operatorname{Hom}\left(E^{\prime}, E\right) \varphi^{\prime-1} \subset D$ depends only on the classes of $(E, \varphi)$ and $\left(E^{\prime}, \varphi^{\prime}\right)$ in $\mathcal{L}$. Recall that $\operatorname{Hom}\left(E^{\prime}, E\right) \otimes \mathbb{Z}_{l}=\operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right)$ for all primes $l$ (see 3.1). Therefore we have by 2.4 :

$$
\begin{aligned}
\varphi \operatorname{Hom}\left(E^{\prime}, E\right) \varphi^{\prime-1} & =\bigcap_{l}\left(\left(\varphi \operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right) \varphi^{\prime-1}\right) \otimes_{\mathbb{Z}} \mathbb{Z}_{l} \cap D\right) \\
& =\bigcap_{l}\left(T_{l}(\varphi) \operatorname{Hom}\left(T_{l}\left(E^{\prime}\right), T_{l}(E)\right) T_{l}\left(\varphi^{\prime}\right)^{-1} \cap D\right) \\
& =\bigcap_{l}\left(d_{l} \operatorname{Hom}\left(T_{l}\left(E_{0}\right), T_{l}\left(E_{0}\right)\right) d_{l}^{\prime-1} \cap D\right) \\
& =d \mathcal{O}_{D} d^{\prime-1}=\Lambda .
\end{aligned}
$$

Moreover, it follows from 2.5 that the condition $\operatorname{Nrd}\left(d_{l}\right)=\operatorname{Nrd}\left(d_{l}^{\prime}\right)$ implies $\operatorname{deg}(\varphi)=$ $\operatorname{deg}\left(\varphi^{\prime}\right)$. Thus $\Lambda$ lies in the proper class $\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]$.

Proposition 3.3. - Let $E, E^{\prime}, E_{1}$ and $E_{1}^{\prime}$ be supersingular elliptic curves. Then we have $\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]=\left[\operatorname{Hom}\left(E_{1}^{\prime}, E_{1}\right)\right]$ if and only if there exists an integer $n$ such that $E_{1}^{\prime} \cong E^{\prime\left(p^{n}\right)}$ and $E_{1} \cong E^{\left(p^{n}\right)}$.

Proof. - Choose quasi-isogenies $\varphi: E \rightarrow E_{0}, \varphi^{\prime}: E^{\prime} \rightarrow E_{0}, \varphi_{1}: E_{1} \rightarrow E_{0}$ and $\varphi_{1}^{\prime}: E_{1}^{\prime} \rightarrow E_{0}$ such that $\operatorname{deg}(\varphi)=\operatorname{deg}\left(\varphi^{\prime}\right)$ and $\operatorname{deg}\left(\varphi_{1}\right)=\operatorname{deg}\left(\varphi_{1}^{\prime}\right)$. We set $\Lambda=$ $\iota_{\varphi, \varphi^{\prime}}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)$ and $\Lambda_{1}=\iota_{\varphi_{1}, \varphi_{1}^{\prime}}\left(\operatorname{Hom}\left(E_{1}^{\prime}, E_{1}\right)\right)$. Let $d, d^{\prime}, d_{1}, d_{1}^{\prime} \in\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times}$ be elements such that the associated pairs in $\mathcal{L}$ via the bijection (2.1) are equal to $[(E, \varphi)],\left[\left(E^{\prime}, \varphi^{\prime}\right)\right],\left[\left(E_{1}, \varphi_{1}\right)\right],\left[\left(E_{1}^{\prime}, \varphi_{1}^{\prime}\right)\right]$ respectively. Then

$$
\Lambda=d \mathcal{O}_{D} d^{\prime-1}, \quad \Lambda_{1}=d_{1} \mathcal{O}_{D} d_{1}^{\prime-1}
$$

We set $\hat{\mathcal{O}}_{D}^{\times}:=\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}$. If $E_{1}^{\prime} \cong E^{\prime\left(p^{-n}\right)}$ and $E_{1} \cong E^{\left(p^{-n}\right)}$, we can choose $\varphi_{1}^{\prime}=F_{E_{0}}^{n} \circ \varphi^{\prime\left(p^{n}\right)}$ and $\varphi_{1}=F_{E_{0}}^{n} \circ \varphi^{\left(p^{n}\right)}$ where $F_{E_{0}}^{n}: E_{0}^{\left(p^{-n}\right)} \rightarrow E_{0}$ is the relative Frobenius. Then we have $d_{1}^{\prime} \hat{\mathcal{O}}_{D}^{\times}=d^{\prime} \Pi^{n} \hat{\mathcal{O}}_{D}^{\times}$and $d_{1} \hat{\mathcal{O}}_{D}^{\times}=d \Pi^{n} \hat{\mathcal{O}}_{D}^{\times}$for a Frobenius element $\Pi$ with respect to $\mathcal{O}_{D}$. Note that $\Pi^{n} \hat{\mathcal{O}}_{D}^{\times}=\mathcal{O}_{D}^{\times} \Pi^{n}$ and $\varepsilon \mathcal{O}_{D}=\mathcal{O}_{D}=\mathcal{O}_{D} \varepsilon$ for $\varepsilon \in \hat{\mathcal{O}}_{D}^{\times}$. Therefore

$$
\Lambda_{1}=d_{1} \mathcal{O}_{D} d_{1}^{\prime-1}=d \Pi^{n} \mathcal{O}_{D} \Pi^{-n} d^{\prime-1}=d \mathcal{O}_{D} d^{\prime-1}=\Lambda
$$

Conversely, assume that $\Lambda$ and $\Lambda_{1}$ are in the same proper class. By Corollary 1.7 there exist $\delta, \delta^{\prime} \in D^{\times}$with $\operatorname{Nrd}(\delta)=\operatorname{Nrd}\left(\delta^{\prime}\right)$ such that $\delta \Lambda=\Lambda_{1} \delta^{\prime}$. Then the maximal orders $O_{l}(\Lambda)$ and $O_{l}\left(\Lambda_{1}\right)$ are in the same conjugacy class. Hence we see that $\{\operatorname{End}(E)\}=\left\{\operatorname{End}\left(E_{1}\right)\right\}$ and then Proposition 2.4 implies that there exists an integer $n$ such that $E_{1} \cong E^{\left(p^{n}\right)}$. Considering $O_{r}(\Lambda)$ and $O_{r}\left(\Lambda_{1}\right)$ we see that there also exists an integer $n^{\prime}$ such that $E_{1}^{\prime} \cong E^{\prime\left(p^{n^{\prime}}\right)}$.

It remains to show that we can choose $n=n^{\prime}$. As all supersingular elliptic curves are defined over $\mathbb{F}_{p^{2}}$, we can assume that $n, n^{\prime} \in\{0,1\}$. If one elliptic curve $X$ is defined over $\mathbb{F}_{p}$ we have $X=X^{(p)}$ and therefore we can assume that all four elliptic curves are not defined over $\mathbb{F}_{p}$. Then we have to show that the following case cannot occur: $E_{1}^{\prime} \cong E^{\prime}$ and $E_{1} \not \approx E$. Note that we know already that $E_{1} \not \approx E$ implies $E_{1} \cong E^{(p)}$.

We can assume that $O_{r}(\Lambda)=O_{r}\left(\Lambda_{1}\right)=: O$. Then $\delta \Lambda=\Lambda_{1} \delta^{\prime}$ implies that $\delta^{\prime} \in N_{D \times}(O)$. As $E^{\prime}=E_{1}^{\prime}$ is not defined over $\mathbb{F}_{p}$ we have $N_{D \times}(O)=\mathbb{Q}^{\times} O^{\times}$by Proposition 2.4. Thus there exists a $\delta_{1} \in D^{\times}$with $\operatorname{Nrd}\left(\delta_{1}\right)=1$ such that $\delta_{1} \Lambda=\Lambda_{1}$. Now $E_{1}=E^{(p)}$ implies by 2.13 that there exists a $\delta_{1}^{\prime} \in D^{\times}$such that $\Lambda_{1}=\delta_{1}^{\prime} \mathfrak{p} \Lambda$ where $\mathfrak{p}$ is the prime ideal of $O_{l}(\Lambda)$ which lies over $(p)$. But this implies

$$
\delta_{1} O=\Lambda_{1} \Lambda^{-1}=\delta_{1}^{\prime} \mathfrak{p}
$$

and this is a contradiction as $\mathfrak{p}$ is not a principal ideal by Proposition 2.4.
3.3. Let $E$ and $E^{\prime}$ be two supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. Consider the natural map

$$
\begin{aligned}
\alpha: \operatorname{Aut}(E) \times \operatorname{Aut}\left(E^{\prime}\right) & \longrightarrow \mathrm{O}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right) \\
\left(\varphi, \varphi^{\prime}\right) & \longmapsto\left(x \longmapsto \varphi x \varphi^{\prime-1}\right)
\end{aligned}
$$

Proposition 3.4. The image of $\alpha$ lies in $\operatorname{SO}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)$ and the kernel of $\alpha$ consists of $\{ \pm 1\}$ (diagonally embedded in $\operatorname{Aut}(E) \times \operatorname{Aut}\left(E^{\prime}\right)$ ). The image of $\alpha$ is equal to $\mathrm{SO}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)$ if and only if $E$ or $E^{\prime}$ is not defined over $\mathbb{F}_{p}$. If both curves are defined over $\mathbb{F}_{p}$ the image of $\alpha$ has index 2 in $\operatorname{SO}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)$.

Proof. - We choose quasi-isogenies $\varphi: E \rightarrow E_{0}$ and $\varphi^{\prime}: E^{\prime} \rightarrow E_{0}$ of the same degree. Set $\Lambda=\varphi \operatorname{Hom}\left(E^{\prime}, E\right) \varphi^{\prime-1} \subset D, O_{l}:=O_{l}(\Lambda)$ and $O_{r}=O_{r}(\Lambda)$. By Lemma 1.6 we know

$$
\mathrm{SO}(\Lambda)=\left\{\left(d, d^{\prime}\right) \in D^{\times} \times D^{\times} \mid d \Lambda=\Lambda d^{\prime}, \operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)\right\} / \mathbb{Q}^{\times}
$$

where $\mathbb{Q}^{\times}$is embedded diagonally. Note that the condition $d \Lambda=\Lambda d^{\prime}$ already implies $\operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)$. Moreover, $d \Lambda=\Lambda d^{\prime}$ implies $O_{l}=O_{l}(d \Lambda)$ and hence $d \in N_{D \times}\left(O_{l}\right)$. Similarly $d^{\prime} \in N_{D \times}\left(O_{r}\right)$. Thus we see

$$
\operatorname{SO}(\Lambda)=\left\{\left(d, d^{\prime}\right) \in N_{D \times}\left(O_{l}\right) \times N_{D^{\times}}\left(O_{r}\right) \mid d \Lambda=\Lambda d^{\prime}\right\} / \mathbb{Q}^{\times}
$$

For $\left(d, d^{\prime}\right) \in \operatorname{SO}(\Lambda)$ we have by Corollary 2.5 :

$$
\begin{equation*}
d \in \mathbb{Q}^{\times} O_{l}^{\times} \Longleftrightarrow d^{\prime} \in \mathbb{Q}^{\times} O_{r}^{\times} . \tag{3.1}
\end{equation*}
$$

Further $\mathbb{Q}^{\times} \cap O_{l}^{\times}=\mathbb{Q}^{\times} \cap O_{r}^{\times}=\{ \pm 1\}$.

Now consider the case where $E$ is not defined over $\mathbb{F}_{p}$. By Proposition 2.4 this is equivalent to $N_{D^{\times}}\left(O_{l}\right)=\mathbb{Q}^{\times} O_{l}^{\times}$. Therefore

$$
\begin{aligned}
\mathrm{SO}(\Lambda) & =\left\{\left(d, d^{\prime}\right) \in \mathbb{Q}^{\times} O_{l}^{\times} \times \mathbb{Q}^{\times} O_{r}^{\times} \mid d \Lambda=\Lambda d^{\prime}\right\} / \mathbb{Q}^{\times} \\
& =\left\{\left(d, d^{\prime}\right) \in O_{l}^{\times} \times O_{r}^{\times}\right\} /\{ \pm 1\}
\end{aligned}
$$

which proves the proposition in this case. The case that $E^{\prime}$ is not defined over $\mathbb{F}_{p}$ is proved by the same argument.

It remains to consider the case that $E$ and $E^{\prime}$ are both defined over $\mathbb{F}_{p}$. Then $\mathbb{Q}^{\times} O_{l}^{\times} \subset N_{D^{\times}}\left(O_{l}\right)$ and $\mathbb{Q}^{\times} O_{r}^{\times} \subset N_{D^{\times}}\left(O_{r}\right)$ are subgroups of index 2 by Proposition 2.4. Moreover it follows from (3.1) that $\left\{\left(d, d^{\prime}\right) \in O_{l}^{\times} \times O_{r}^{\times}\right\} /\{ \pm 1\}$ is a subgroup of index 2 of $\operatorname{SO}(\Lambda)$. This finishes the proof.

Corollary 3.5. - Let $E$ and $E^{\prime}$ be two supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. Then we have

$$
\# \operatorname{SO}\left(\operatorname{Hom}\left(E^{\prime}, E\right)\right)= \begin{cases}\# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right), & \text { if } E, E^{\prime} \text { both defined over } \mathbb{F}_{p} \\ \frac{1}{2} \# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right), & \text { otherwise }\end{cases}
$$

3.4. We now sketch an alternative formulation in the language of groupoids of some of the above results suggested by the referee. We start with some general notations.

Let $X$ be a set and let $H$ be a group acting on $X$ from the left. Then we denote by $[H \backslash X]$ the category whose objects are the elements of $X$ and whose morphisms are for $x, x^{\prime} \in X$

$$
\operatorname{Hom}_{[H \backslash X]}\left(x, x^{\prime}\right)=\left\{h \in H \mid h x=x^{\prime}\right\}
$$

Composition of morphisms is given by the multiplication in the group H. Clearly, this category is a groupoid (i.e., every morphism is an isomorphism), and two objects $x, x^{\prime} \in X$ are isomorphic if and only if they are in the same $H$-orbit.

If $X$ is of the form $G / H^{\prime}$ for some group $G$ and some subgroup $H^{\prime}$ and if $H$ is a subgroup of $G$ acting in the natural way on $X$, we have by definition

$$
\begin{equation*}
\operatorname{Hom}_{\left[H \backslash\left(G / H^{\prime}\right)\right]}\left(g_{1} H^{\prime}, g_{2} H^{\prime}\right)=G \cap g_{2} H^{\prime} g_{1}^{-1} . \tag{3.2}
\end{equation*}
$$

We denote by $\mathcal{M}$ the following category. The objects are supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. For two such supersingular elliptic curves $E^{\prime}$ and $E$ the morphisms from $E^{\prime}$ to $E$ in $\mathcal{M}$ are by definition the isomorphisms of elliptic curves from $E^{\prime}$ to $E$ over $\overline{\mathbb{F}}_{p}$. Then Theorem 2.1 can be made more precise by saying that (2.2) induces an equivalence of categories

$$
\left[D^{\times} \backslash\left(\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}\right)\right] \approx \mathcal{M}
$$

This follows from (3.2) and the arguments in 3.2.
Let $\operatorname{SL}\left(\mathcal{O}_{D}\right)$ be the algebraic group over $\mathbb{Z}$ of elements with reduced Norm equal to 1 . Then Lemma 3.1 is by 2.5 equivalent to the fact that the canonical functor

$$
\left[\operatorname{SL}\left(\mathcal{O}_{D}\right)(\mathbb{Q}) \backslash\left(\operatorname{SL}\left(\mathcal{O}_{D}\right)\left(\mathbb{A}_{f}\right) / \operatorname{SL}\left(\mathcal{O}_{D}\right)(\widehat{\mathbb{Z}})\right)\right] \longrightarrow\left[D^{\times} \backslash\left(\left(D \otimes_{\mathbb{Q}} \mathbb{A}_{f}\right)^{\times} /\left(\mathcal{O}_{D} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)^{\times}\right)\right]
$$

is essentially surjective (which also can easily be proved directly). Moreover it follows from the definition of a morphism in the above categories that this functor is indeed an equivalence of categories.

Arguing as in the second proof of Lemma 1.6 one sees that the morphism

$$
\alpha: \mathrm{SL}(D) \times \mathrm{SL}(D) \rightarrow \mathrm{SO}(D), \quad\left(d, d^{\prime}\right) \longmapsto\left(\delta \longmapsto d \delta d^{\prime-1}\right)
$$

is a surjection of algebraic groups with $\operatorname{Ker}(\alpha)=\mu_{2}$ (embedded diagonally in $\operatorname{SL}(D) \times$ $\mathrm{SL}(D))$. Then $\alpha$ induces a functor

$$
\begin{align*}
& A:\left[\mathrm{SL}\left(\mathcal{O}_{D}\right)(\mathbb{Q}) \backslash\left(\mathrm{SL}\left(\mathcal{O}_{D}\right)\left(\mathbb{A}_{f}\right) / \mathrm{SL}\left(\mathcal{O}_{D}\right)(\widehat{\mathbb{Z}})\right)\right]^{2} \\
& \longrightarrow\left[\mathrm{SO}(D)(\mathbb{Q}) \backslash\left(\mathrm{SO}(D)\left(\mathbb{A}_{f}\right) / \mathrm{SO}\left(\mathcal{O}_{D}\right)(\widehat{\mathbb{Z}})\right)\right] . \tag{3.3}
\end{align*}
$$

As we have seen, we can identify the left hand side with $\mathcal{M} \times \mathcal{M}$. The set of isomorphism classes of the right side is the set of proper classes in the genus of $\mathcal{O}_{D}=\operatorname{End}\left(E_{0}\right)$ (Corollary 1.4 and 1.5). Then Proposition 3.2 asserts that $A$ is essentially surjective. Moreover, Proposition 3.3 describes the fibres of the induced surjective map on isomorphism classes and Proposition 3.4 is the analysis in which way $A$ fails to be fully faithful.

Finally, Proposition 4.1 below can also be expressed in the language of groupoids using the notion of direct and inverse image for functions (or more generally for sheaves) with respect to the functor $A$. We omit the details.

## 4. Local densities

4.1. Recall that for two quadratic spaces $Q$ and $L$ over $\mathbb{Z}$ we write $R_{L}(Q)$ for the number of isometries $Q \rightarrow L$. Note that $R_{L}(Q)$ depends only on the classes of $L$ and $Q$. In this section we are going to express

$$
\begin{equation*}
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(Q)}{u_{E} u_{E^{\prime}}} \tag{4.1}
\end{equation*}
$$

in terms of local densities. Here the $\left(E^{\prime}, E\right)$ runs through all pairs of isomorphism classes of supersingular elliptic curves over $\overline{\mathbb{F}}_{p}, Q$ is a fixed ternary positive definite quadratic form over $\mathbb{Z}$, and $u_{E}=\frac{1}{2} \# \operatorname{Aut}(E)$.

Proposition 4.1. - Fix a supersingular elliptic curve $E_{0}$. Then we have

$$
\begin{equation*}
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(Q)}{u_{E} u_{E^{\prime}}}=4 \sum_{L} \frac{R_{L}(Q)}{\# \mathrm{SO}(L)} . \tag{4.2}
\end{equation*}
$$

Here on the right hand side, $L$ runs through the proper classes of the genus of $\operatorname{End}\left(E_{0}\right)$.
Proof. - By Proposition 3.2 we know that the proper classes $\left[\operatorname{Hom}\left(E^{\prime}, E\right)\right]$ exhaust the genus of $\operatorname{End}\left(E_{0}\right)$ if $\left(E^{\prime}, E\right)$ runs through all pairs of supersingular elliptic curves. If $E$ and $E^{\prime}$ are both defined over $\mathbb{F}_{p}$, the class $L$ of $\operatorname{Hom}\left(E^{\prime}, E\right)$ occurs once in the sum on the left hand side of (4.2) by Proposition 3.3, and we have $\# \mathrm{SO}(L)=4 u_{E} u_{E^{\prime}}$ by

Corollary 3.5. Otherwise, the class $L$ of $\operatorname{Hom}\left(E^{\prime}, E\right)$ occurs twice in the sum on the left hand side of (4.2) and we have $\# \mathrm{SO}(L)=2 u_{E} u_{E^{\prime}}$. This proves the proposition.
4.2. Let $M$ be a quadratic space over $\mathbb{Z}_{p}$. We denote by $Q_{M}: M \rightarrow \mathbb{Z}_{p}$ its quadratic form and let $B_{M}$ the bilinear form given by

$$
B_{M}(x, y)=\frac{1}{2}\left(Q_{M}(x+y)-Q_{M}(x)-Q_{M}(y)\right)
$$

4.3. Let $M$ and $N$ be two quadratic spaces over $\mathbb{Z}_{p}$ of ranks $m$ and $n$, respectively. We choose bases $\left(\mu_{i}\right)$ and $\left(\nu_{i}\right)$ of $M, N$ respectively and let $T=\left(B_{M}\left(\mu_{i}, \mu_{j}\right)\right) \in$ $\operatorname{Sym}_{m}\left(\mathbb{Z}_{p}\right)^{\vee}$ and $S=\left(B_{N}\left(\nu_{i}, \nu_{j}\right)\right) \in \operatorname{Sym}_{n}\left(\mathbb{Z}_{p}\right)^{\vee}$ be the corresponding matrices. For $r \geq 0$ we define $A_{p^{r}}(M, N)=A_{p^{r}}(T, S)$ as

$$
\begin{aligned}
& \#\left\{\left.X \in M_{n, m}\left(\mathbb{Z}_{p} / p^{r} \mathbb{Z}_{p}\right)\right|^{t} X S X-T \in p^{r} \operatorname{Sym}_{m}\left(\mathbb{Z}_{p}\right)^{\vee}\right\} \\
= & \#\left\{\sigma: M / p^{r} M \rightarrow N / p^{r} N \mid Q_{N}(\sigma(x)) \equiv Q_{M}(x) \bmod p^{r}\right\} .
\end{aligned}
$$

For $r \gg 0$ we set

$$
\alpha_{p}(M, N)=\alpha_{p}(T, S)=2^{-\delta_{m n}}\left(p^{r}\right)^{m(m+1) / 2-m n} A_{p^{r}}(M, N)
$$

It is shown in $[\mathbf{K i}, 5.6]$ that this is independent of $r$ if $r$ is sufficiently big. We call $\alpha_{p}(T, S)$ the local representation density. Note that for $p=2$ our representation density is $2^{m(m-1) / 2}$-times the representation density $\alpha_{p}$ defined in [ $\left.\mathbf{K i}\right]$.
4.4. For any class $N$ of lattices in a positive definite quadratic space over $\mathbb{Q}$ we set

$$
o(N)=\# \mathrm{O}(N), \quad \text { so }(N)=\# \mathrm{SO}(N)
$$

Moreover, we set

$$
w(N):=\sum_{\substack{N^{\prime} \text { class } \\ \text { in } \operatorname{gen}(N)}} \frac{1}{o\left(N^{\prime}\right)}
$$

where $N^{\prime}$ runs through all classes within the genus of $N$. We call $w(N)$ the weight of $N$.

Then we have $o(N)=s o(N)$ if and only if there exist two proper classes in $N$. Therefore

$$
\begin{equation*}
w(N)=\frac{1}{2} \sum_{\substack{N^{\prime} \text { proper class } \\ \text { in gen }(N)}} \frac{1}{s o\left(N^{\prime}\right)} \tag{4.3}
\end{equation*}
$$

Let $M$ and $N$ be lattices in positive definite quadratic spaces over $\mathbb{Q}$. By the mean value of the representation of $M$ by $N$ we mean

$$
\begin{aligned}
m(M, N) & =w(N)^{-1} \sum_{\substack{N^{\prime} \text { class } \\
\text { in gen }(N)}} \frac{R_{N^{\prime}}(M)}{o\left(N^{\prime}\right)} \\
& =\left(\sum_{\substack{N^{\prime} \text { proper class } \\
\text { in gen }(N)}} \frac{1}{\operatorname{so(}\left(N^{\prime}\right)}\right)^{-1} \sum_{\substack{N^{\prime} \text { proper class } \\
\text { in gen }(N)}} \frac{R_{N^{\prime}}(M)}{s o\left(N^{\prime}\right)} .
\end{aligned}
$$

Clearly, $m(M, N)$ depends only on the genus of $N$.
4.5. We recall the Minkowski-Siegel formula (cf. $[\mathbf{K i}, 6.8]$ ): Let $M$ and $N$ be lattices in positive definite quadratic spaces over $\mathbb{Q}$. For any prime $l$ we define $\alpha_{l}(M, N):=$ $\alpha_{l}\left(M \otimes_{\mathbb{Z}} \mathbb{Z}_{l}, N \otimes_{\mathbb{Z}} \mathbb{Z}_{l}\right)$. Put $m:=\operatorname{rank}(M)$ and $n:=\operatorname{rank}(N)$. Set

$$
\epsilon_{m, n}:= \begin{cases}\frac{1}{2}, & \text { if either } n=m+1 \text { or } n=m>1 \\ 1, & \text { otherwise }\end{cases}
$$

We also define

$$
\begin{aligned}
& \alpha_{\infty}(M, N):=\pi^{m(2 n-m+1) / 4}\left(\prod_{i=0}^{m-1} \Gamma((n-i) / 2)^{-1}\right) \\
& \times(\operatorname{det}(N))^{-m / 2}(\operatorname{det}(M))^{(n-m-1) / 2}
\end{aligned}
$$

Here $\Gamma$ denotes the gamma function.

## Theorem 4.2 (Minkowski, Siegel)

$$
\begin{equation*}
m(M, N)=\epsilon_{m, n} 2^{-m(m-1) / 2} \alpha_{\infty}(M, N) \prod_{l} \alpha_{l}(M, N) \tag{4.4}
\end{equation*}
$$

where $l$ runs through all prime numbers.
Theorem 4.3. - Let $M$ be a positive definite ternary quadratic space over $\mathbb{Z}$ and let $N$ be the genus of $\operatorname{End}\left(E_{0}\right)$ for a supersingular elliptic curve over $\overline{\mathbb{F}}_{p}$. Then

$$
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(M)}{u_{E} u_{E^{\prime}}}=8\left(\frac{p-1}{12}\right)^{2} \frac{\pi^{4}}{p^{3}} \prod_{l} \alpha_{l}(M, N)
$$

where $l$ runs through all prime numbers $l$.
Proof. - We apply the Minkowski-Siegel formula: We have $m=3$ and $n=4$. We first compute $\operatorname{det}(N)$. As $N$ is positive definite, it suffices to compute $\operatorname{ord}_{l}\left(N \otimes \mathbb{Z}_{l}\right)$ for every prime $l$. For $l \neq p$ the quadratic $\mathbb{Z}_{l}$-space $N \otimes \mathbb{Z}_{l}$ is isomorphic to ( $M_{2}\left(\mathbb{Z}_{l}\right)$, det)
and with respect to the basis $\left(\begin{array}{ll}1 & 0 \\ 0 & 0\end{array}\right),\left(\begin{array}{ll}0 & 0 \\ 0 & 1\end{array}\right),\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right),\left(\begin{array}{cc}0 & 0 \\ -1 & 0\end{array}\right)$ the associated matrix is equal to

$$
S_{l}=\frac{1}{2}\left(\begin{array}{llll}
0 & 1 & &  \tag{4.5}\\
1 & 0 & & \\
& & 0 & 1 \\
& & 1 & 0
\end{array}\right)
$$

For $l=p$ we have $N \otimes \mathbb{Z}_{l} \cong\left(\mathcal{O}_{D_{p}}, \mathrm{Nrd}\right)$ and hence there exists a basis such that the associated matrix is equal to

$$
S_{p}=\left\{\begin{array}{lll}
\operatorname{diag}(1,-\delta, p,-\delta p), & \text { if } p \neq 2,  \tag{4.6}\\
\left(\begin{array}{cccc}
1 & 1 / 2 & & \\
1 / 2 & 1 & & \\
& & 2 & 1 \\
& & 1 & 2
\end{array}\right), & \text { if } p=2,
\end{array}\right.
$$

where $\delta$ is some element in $\mathbb{Z}_{p}^{\times} \backslash\left(\mathbb{Z}_{p}^{\times}\right)^{2}$ (see $[\mathbf{K i}, 5.2]$ ).
If follows that $\operatorname{det}(N)$ is equal to $2^{-4} p^{2}$ and hence we get

$$
\alpha_{\infty}(M, N)=\frac{\pi^{9 / 2}}{\Gamma(1) \Gamma(3 / 2) \Gamma(2)}\left(2^{-4} p^{2}\right)^{-3 / 2}=2^{7} \frac{\pi^{4}}{p^{3}}
$$

By (4.3) we can calculate the weight of $N$ as

$$
w(N)=\frac{1}{2} \sum_{\substack{N^{\prime} \text { proper class } \\ \text { in gen }(N)}} \frac{1}{\operatorname{son}\left(N^{\prime}\right)} .
$$

Using Proposition 3.3 and Corollary 3.5, it follows that

$$
\begin{aligned}
w(N) & =\frac{1}{2} \sum_{\left(E^{\prime}, E\right)} \frac{1}{\# \operatorname{Aut}(E) \# \operatorname{Aut}\left(E^{\prime}\right)} \\
& =\frac{1}{2}\left(\sum_{E} \frac{1}{\# \operatorname{Aut}(E)}\right)^{2} \\
& =\frac{1}{2}\left(\frac{p-1}{24}\right)^{2}
\end{aligned}
$$

where in the last equality we used Eichler's mass formula.

Now using Proposition 4.1, the Minkowski-Siegel formula tells us

$$
\begin{aligned}
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(M)}{u_{E} u_{E^{\prime}}} & =4 \sum_{\substack{N^{\prime} \text { in proper class } \\
\text { of gen }(N)}} \frac{R_{N^{\prime}}(M)}{\# \mathrm{SO}\left(N^{\prime}\right)} \\
& =8 w(N) m(M, N) \\
& =4\left(\frac{p-1}{24}\right)^{2} \epsilon_{3,4} 2^{-3} \alpha_{\infty}(M, N) \prod_{l} \alpha_{l}(M, N) \\
& =8\left(\frac{p-1}{12}\right)^{2} \frac{\pi^{4}}{p^{3}} \prod_{l} \alpha_{l}(M, N)
\end{aligned}
$$

4.6. In [Wd2] we will give explicit expressions for the local representation densities $\alpha_{l}(M, N)$ for arbitrary positive definite ternary quadratic spaces $M$ and all primes $l$. We deduce (cf. [GK, 6.23]):

Corollary 4.4. - Define $\beta_{l}(M)=\left(1-l^{-2}\right)^{-2} \alpha_{l}(M, N)$ and $\Delta(M)=4 \operatorname{det}(M)$. Assume that $M$ is a positive definite ternary quadratic space over $\mathbb{Z}$ which is isotropic over $\mathbb{Q}_{l}$ for all $l \neq p$. Then $M$ is anisotropic over $\mathbb{Q}_{p}$ and

$$
\sum_{\left(E^{\prime}, E\right)} \frac{R_{\operatorname{Hom}\left(E^{\prime}, E\right)}(M)}{u_{E} u_{E^{\prime}}}=4 \prod_{l \Delta \Delta(M), l \neq p} \beta_{l}(M) .
$$

Proof. - We have $\alpha_{p}(M, N)=2(p+1)^{2} p^{-1}$ by $[\mathbf{W d 2}$, Theorem 1.1], and hence $\beta_{p}(M)=2 p^{3} /(p-1)^{2}$. Moreover we know by [Wd2, Corollary 2.2], that $\beta_{l}(M)=1$ for all $l \neq p$ which do not divide $\Delta(M)$. Therefore we have

$$
\begin{aligned}
\prod_{l} \alpha_{l}(M, N) & =\frac{2 p^{3}}{(p-1)^{2}} \prod_{l}\left(1-l^{-2}\right)^{2} \prod_{l \mid \Delta(M), l \neq p} \beta_{l}(M) \\
& =\frac{2 p^{3}}{(p-1)^{2}} \zeta(2)^{-2} \prod_{l \Delta(M), l \neq p} \beta_{l}(M)
\end{aligned}
$$

As $\zeta(2)=\pi^{2} / 6$, the corollary follows from Theorem 4.3.

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# 6. LUBIN-TATE FORMAL GROUPS 

by

Volker Meusers


#### Abstract

We give an exposition of the theory of formal complex multiplication in local fields after Lubin and Tate. We recall the construction of Lubin-Tate modules, the structure of torsion points of their generic fibre and explicit local class field theory. We follow the original exposition of Lubin and Tate, and the exposition in Neukirch's book.


Résumé (Groupes formels de Lubin-Tate). - Nous donnons une exposition de la théorie de la multiplication complexe formelle dans les corps locaux d'après Lubin et Tate. On rappelle la construction des modules de Lubin-Tate, la structure de leurs modules de torsion de leur fibre générique et la théorie du corps de classes locale explicite. On suit l'article original de Lubin et Tate, et le livre de Neukirch.

## 1. Construction of Lubin-Tate Modules

Let $K$ be a field complete with respect to some discrete valuation. Let $\mathcal{O}_{K}$ be its ring of integers, $\mathfrak{p}$ its maximal ideal. Assume the residue field $\mathcal{O}_{K} / \mathfrak{p}$ to be finite and let $q$ be the number of its elements. Prime elements of $\mathcal{O}_{K}$ are denoted by $\pi$ or $\bar{\pi}$. Let $k$ be an algebraic closure of $\mathcal{O}_{K} / \mathfrak{p}$. Let $K^{\text {sep }}$ be a fixed separable closure of $K$ and $K^{\mathrm{nr}} \subseteq K^{\text {sep }}$ the maximal unramified extension of $K$. Let $M$ and $C$ denote the completions of $K^{\mathrm{nr}}$ and $K^{\text {sep }}$. Denote by $\mathcal{O}_{M}$ (resp. $\mathcal{O}_{C}$ ) the ring of integers of $M$ (resp. $C$ ). Let $\widehat{\mathcal{C}}$ be the category of complete local noetherian $\mathcal{O}_{K}$-algebras with residue field $k$.

Definition 1.1. - Let $i: \mathcal{O}_{K} \rightarrow R$ be an $\mathcal{O}_{K}$-algebra, e.g. $\mathcal{O}_{K}, \mathcal{O}_{M}$ or $k$. A formal $\mathcal{O}_{K}$-module over $R$ is a pair $\left(H, \gamma_{H}\right)$ consisting of a (one-dimensional commutative) formal group law $H(X, Y) \in R[[X, Y]]$ together with a ring homomorphism

2000 Mathematics Subject Classification. - 11S31, 14K22, 14L05.
Key words and phrases. - Lubin-Tate formal groups, local class field theory, complex multiplication.
$\gamma_{H}: \mathcal{O}_{K} \rightarrow \operatorname{End}_{R}(H) \subset R[[T]]$ given by sending an element $a \in \mathcal{O}_{K}$ to the endomorphism $\gamma_{H}(a)(T) \in R[[T]]$ of $H(X, Y)$. As a normalization condition we require that the $\mathcal{O}_{K}$-algebra structure on $R$ induced by the isomorphism

$$
\mathcal{O}_{K} \cong \xlongequal{\cong} \operatorname{Lie}(H),\left.a \longmapsto \frac{\partial \gamma_{H}(a)(T)}{\partial T}\right|_{T=0}
$$

agrees with the structure given by $i: \mathcal{O}_{K} \rightarrow R$, in other words we require $\gamma_{H}(a)(T)$ to be of the form

$$
\gamma_{H}(a)(T)=i(a) T+\cdots \in R[[T]] .
$$

We write $[a](T)$ for $\gamma_{H}(a)(T)$ and $a=i(a) \in R$ if no confusion is possible.
For $R \in \widehat{\mathcal{C}}$ write $H(R)$ for the abelian group $\left(\mathfrak{m}_{R},+_{H}\right)$ where we have set $x+_{H} y=H(x, y)$ for $x, y \in \mathfrak{m}_{R}$. This converges since $R$ is assumed to be complete. This group is also an (ordinary) $\mathcal{O}_{K}$-module by setting $a x=a \cdot{ }_{H} x=[a](x)$. Note that unless $\left(H, \gamma_{H}\right)$ is the formal additive group, i.e., $\left(\widehat{\mathbb{G}}_{a}(X, Y)=X+Y, \gamma_{\widehat{\mathbb{G}}_{a}}(a)(T)=a T\right)$, this $\mathcal{O}_{K^{-}}$-module structure is not the standard structure on $\mathfrak{m}_{R}$ as an ideal of $R$. For a finite extension $L \mid K$ with ring of integers $\mathcal{O}_{L} \in \widehat{\mathfrak{C}}$ and maximal ideal $\mathfrak{m}_{L} \subset \mathcal{O}_{L}$ we set $H(L)=H\left(\mathfrak{m}_{L}\right)$. Similarly for infinite extensions after completion.

The goal of this section is to construct, as for ordinary complex multiplication (see Remark 3.5 below), a formal $\mathcal{O}_{K}$-module $\left(G, \gamma_{G}\right)$ over $\mathcal{O}_{M}$ such that

$$
G[\mathfrak{p}]=\bigcap_{a \in \mathfrak{p}} \operatorname{Ker}(a)=G[\pi]
$$

is isomorphic to the kernel of the Frobenius $G \otimes k \rightarrow(G \otimes k)^{(q)}$ when reduced modulo the maximal ideal of $\mathcal{O}_{M}$. Lubin and Tate construct $G$ as a base change $G=H_{\pi} \otimes \mathcal{O}_{K}$ $\mathcal{O}_{M}$ of a formal $\mathcal{O}_{K}$-module $H_{\pi}$ over $\mathcal{O}_{K}$, the so called Lubin-Tate module associated to the prime element $\pi \in \mathcal{O}_{K}$. As we will see $H_{\pi}$ depends on the chosen $\pi$ while $G$ will be independent of it.
By our normalization condition $\gamma_{G}(\pi)(T)$ is of the form

$$
\gamma_{G}(\pi)(T)=\pi T+\cdots \in \mathcal{O}_{K}[[T]] .
$$

The condition on the Frobenius requires that

$$
\gamma_{G}(\pi)(T) \equiv T^{q} \bmod \pi
$$

This justifies the following definition:
Definition 1.2. - A power series $f(T)=\pi T+\cdots \in \mathcal{O}_{K}[[T]]$ such that

$$
f(T) \equiv T^{q} \bmod \pi
$$

is called a Lubin-Tate series associated to $\pi$. The set of Lubin-Tate series for $\pi$ is denoted by $\mathcal{F}_{\pi}$. A formal $\mathcal{O}_{K}$-module $\left(H, \gamma_{H}\right)$ over $\mathcal{O}_{K}$ with $\gamma_{H}(\pi)(T) \in \mathcal{F}_{\pi}$ is called Lubin-Tate module.

## Examples 1.3

(1) The simplest example of a Lubin-Tate-series is

$$
f(T)=\pi T+T^{q} \in \mathcal{F}_{\pi}
$$

(2) In the cyclotomic case, i.e., for $K=\mathbb{Q}_{p}, \mathcal{O}_{K}=\mathbb{Z}_{p}$ and $\pi=p \in \mathbb{Z}_{p}$ the polynomial

$$
f(T)=(T+1)^{p}-1=p T+p(\ldots)+T^{p} \in \mathcal{F}_{\pi}
$$

is a Lubin-Tate-series associated to $\pi=p$. One easily checks that in this case the formal multiplicative group

$$
\widehat{\mathbb{G}}_{m}(X, Y)=(1+X)(1+Y)-1
$$

is a Lubin-Tate module associated to $f(T)$.
The construction of Lubin-Tate-modules is based on the following lemma.
Lemma 1.4. - Let $\pi, \bar{\pi}$ be two prime elements of $M$ and $f(T) \in \mathcal{F}_{\pi}$ resp. $g(T) \in \mathcal{F}_{\bar{\pi}}$. Let $L\left(X_{1}, \ldots, X_{n}\right)=\sum_{i=1}^{n} a_{i} X_{i}$ be a linear form with coefficients in $\mathcal{O}_{M}$ such that

$$
\pi L\left(X_{1}, \ldots, X_{n}\right)=\bar{\pi} L^{\sigma}\left(X_{1}, \ldots, X_{n}\right)
$$

where $\sigma$ is the continuous extension of the Frobenius in $\operatorname{Gal}\left(K^{n r} \mid K\right)$ to M. Then there exists a unique power series $F\left(X_{1}, \ldots, X_{n}\right) \in \mathcal{O}_{M}\left[\left[X_{1}, \ldots, X_{n}\right]\right]$ such that

$$
\begin{equation*}
F\left(X_{1}, \ldots, X_{n}\right) \equiv L\left(X_{1}, \ldots, X_{n}\right) \bmod \left(X_{1}, \ldots, X_{n}\right)^{2} \tag{1.1}
\end{equation*}
$$

and

$$
\begin{equation*}
f\left(F\left(X_{1}, \ldots, X_{n}\right)\right)=F^{\sigma}\left(g\left(X_{1}\right), \ldots, g\left(X_{n}\right)\right) \tag{1.2}
\end{equation*}
$$

where $\left(X_{1}, \ldots, X_{n}\right)$ denotes the ideal generated by $X_{1}, \ldots, X_{n}$. If the coefficients of $f, g$, $L$ lie in $\mathcal{O}_{K}$ then $F$ also has coefficients in $\mathcal{O}_{K}$.

The idea of the proof is to construct $F$ inductively modulo powers of the ideal generated by $X_{1}, \ldots, X_{n}$ and then use the completeness of the power series ring. The induction starts with (1.1). For the induction step one plugs in (1.2) and uses that $f$ and $g$ are Lubin-Tate series to see that the coefficients are in $\mathcal{O}_{M}$. See $[\mathbf{N}]$ for a detailed proof.

We use the lemma to construct Lubin-Tate modules as follows:
For $f(T) \in \mathcal{F}_{\pi}$ let $H_{f}(X, Y)$ be the unique solution of the equations

$$
H_{f}(X, Y) \equiv X+Y \bmod (X, Y)^{2}
$$

and

$$
f\left(H_{f}(X, Y)\right) \equiv H_{f}(f(X), f(Y))
$$

For each $a \in \mathcal{O}_{K}$ and $f(T), g(T) \in \mathcal{F}_{\pi}$ let $[a]_{f, g}(T)$ be the unique solution of

$$
[a]_{f, g}(T) \equiv a T \bmod T^{2}
$$

and

$$
f\left([a]_{f, g}(T)\right) \equiv[a]_{f, g}(g(T))
$$

To simplify notations we shall write $[a]_{f}$ instead of $[a]_{f, f}$. The following theorem shows that the series $H_{f}(X, Y)$ together with $\gamma_{H_{f}}(a)(T)=[a]_{f}(T)$ is in fact a Lubin-Tate module associated to $f(T)$.

Theorem 1.5. - For any $f(T) \in \mathcal{F}_{\pi}$ the series $H_{f}(X, Y)$ is a formal group law over $\mathcal{O}_{K}$, i.e., the following identities hold:

$$
\begin{aligned}
H_{f}(X, Y) & =H_{f}(Y, X) \\
H_{f}\left(H_{f}(X, Y), Z\right) & =H_{f}\left(X, H_{f}(Y, Z)\right) \\
H_{f}(X, 0) & =X \\
H_{f}(0, Y) & =Y \\
H_{f}\left(X,[-1]_{f}(X)\right) & =0 .
\end{aligned}
$$

For $g, h \in \mathcal{F}_{\pi}$ and $a, b \in \mathcal{O}_{K}$ we have

$$
\begin{aligned}
H_{f}\left([a]_{f, g}(X),[a]_{f, g}(Y)\right) & =[a]_{f, g}\left(H_{g}(X, Y)\right) \\
{[a]_{f, g}\left([b]_{g, h}(T)\right) } & =[a b]_{f, h}(T) \\
{[a+b]_{f, g}(T) } & =H_{f}\left([a]_{f, g}(T),[b]_{f, g}(T)\right) \\
{[\pi]_{f}(T) } & =f(T) \\
{[1]_{f}(T) } & =T .
\end{aligned}
$$

In particular $\left(H_{f}(X, Y), \gamma_{H_{f}}\right)$ with $\gamma_{H_{f}}(a)(T)=[a]_{f}(T)$ is a Lubin-Tate-module such that $\gamma_{H_{f}}(\pi)(T)=f(T)$. For two series $f(T), g(T) \in \mathcal{F}_{\pi}$ we have the canonical isomorphism

$$
[1]_{f, g}(T): H_{g} \xrightarrow{\cong} H_{f}
$$

of formal $\mathcal{O}_{K}$-modules over $\mathcal{O}_{K}$.
The equalities in the Theorem are all true modulo squares and follow from the uniqueness assertion of Lemma 1.4. For a detailed proof see [ $\mathbf{N}$, proof of Theorem V.4.6].

Remark 1.6. - Although $H_{f}$ does not depend on the particular choice $f \in \mathcal{F}_{\pi}$ it does depend on the particular choice of the uniformizing element $\pi \in \mathcal{O}_{K}$. They become isomorphic over $\mathcal{O}_{M}$ because of the following lemma.

Lemma 1.7. - Let $\pi$ and $\bar{\pi}$ be two prime elements of $\mathcal{O}_{K}$ with $\pi=u \bar{\pi}$ for some unit $u \in \mathcal{O}_{K}^{\times}$. Let $\sigma$ be the Frobenius of $M$ as above. There exists some $\epsilon \in \mathcal{O}_{M}^{\times}$such that $u=\epsilon^{\sigma-1}$. Let $f(T) \in \mathcal{F}_{\pi}$ and $g(T) \in \mathcal{F}_{\bar{\pi}}$ be Lubin-Tate series. Then there exists a unique power series $\theta(X) \in \mathcal{O}_{M}[[X]]$ such that $\theta(X)=\epsilon X$ modulo $(X)^{2}$ and $f \circ \theta=\theta^{\sigma} \circ g$. Furthermore $\theta(X)$ induces an isomorphism $H_{g} \xlongequal{\cong} H_{f}$ of Lubin-Tate modules (defined over $\mathcal{O}_{M}$ ).

This is proved using Lemma 1.4. For a detailed proof see [ $\mathbf{N}$, Corollary V.2.3], and also [LT, Lemma 2].

## 2. Torsion points of the Generic Fibre

Now fix some $f \in \mathcal{F}_{\pi}$. We want to describe the structure of torsion points of the generic fibre of $H_{f}(C)$ as a Galois module. Recall that for every separable algebraic extension $K \subset L \subset C$ we set $H_{f}(L)=H_{f}\left(\widehat{\mathcal{O}}_{L}\right)$. If $L_{1} \subset L$ then $H_{f}\left(L_{1}\right) \subset H_{f}(L)$. If $L \mid L_{1}$ is Galois then $\operatorname{Gal}\left(L \mid L_{1}\right)$ operates naturally on $H_{f}(L)$ in a manner compatible with the $\mathcal{O}_{K}$-module structure. This results from the fact that the Galois group operates continuously on $\widehat{\mathcal{O}}_{L}$ and that $H_{f}$ is defined over $\mathcal{O}_{K} \subseteq \mathcal{O}_{L_{1}}$. In this way $H_{f}(L)$ becomes a $\operatorname{Gal}\left(L \mid L_{1}\right) \times \mathcal{O}_{K}$-module. For another $g \in \mathcal{F}_{\pi}$ the canonical map induced by $[1]_{f, g}(T)$ is an isomorphism of $\operatorname{Gal}\left(L \mid L_{1}\right) \times \mathcal{O}_{K}$-modules. It commutes with the inclusions $H_{f}\left(L_{1}\right) \subset H_{f}(L)$.

Set

$$
\Lambda_{f}=\bigcup_{m \geq 0} H_{f}(C)\left[\mathfrak{p}^{m}\right] \subset H_{f}(C)
$$

Then $\Lambda_{f}$ is a torsion $\mathcal{O}_{K}$-module, i. e., the union over its sub-modules $\Lambda_{f, m}=\Lambda_{f}\left[\mathfrak{p}^{m}\right]$. It is clear that the Galois extension $K \subset L_{\pi, m}=K\left(\Lambda_{f}[m]\right)$ does not depend on $f \in \mathcal{F}_{\pi}$. Let us denote its Galois group by $G_{\pi, m}=\operatorname{Gal}\left(L_{\pi, m} \mid K\right)$.

Theorem 2.1. - Let $\pi$ be a prime element of $\mathcal{O}_{K}$ and $f \in \mathcal{F}_{\pi}$.
(1) The $\mathcal{O}_{K}$-module $\Lambda_{f}$ is divisible.
(2) For each $m$, the $\mathcal{O}_{K}$-module $\Lambda_{f, m}$ is isomorphic to $\mathcal{O}_{K} / \mathfrak{p}^{m}$.
(3) The $\mathcal{O}_{K}$-module $\Lambda_{f}$ is isomorphic to $K / \mathcal{O}_{K}$.
(4) For each $\tau \in G_{\pi}$ there exists a unique $u_{\tau} \in \mathcal{O}_{K}^{\times}$such that $\tau \lambda=\left[u_{\tau}\right]_{f}(\lambda)$ for every $\lambda$ in $\Lambda_{f}$.
(5) The map $\tau \mapsto u_{\tau}$ is an isomorphism of $G_{\pi}$ onto the group $\mathcal{O}_{K}^{\times}$, under which the quotients $G_{\pi, m}$ of $G_{\pi}$ correspond to the quotients $\mathcal{O}_{K}^{\times} /\left(1+\mathfrak{p}^{m}\right)$ of $\mathcal{O}_{K}^{\times}$.

See [LT] for a proof.
Example 2.2. - In the cyclotomic case we get $1+\Lambda_{f, m}=\mu_{p^{m}}, 1+\Lambda_{f}=\mu_{p^{\infty}}$. We have $\widehat{\mathbb{G}}_{m}\left(\mathbb{Q}_{p}\right)=p \mathbb{Z}_{p}$ with addition given by the identification with $1+p \mathbb{Z}_{p} \subset \mathbb{Z}_{p}^{\times}$ as a multiplicative subgroup. In this case the multiplicative structure is given by exponentiating, i.e.,

$$
[a]_{f}(T)=\sum_{n=1}^{\infty}\binom{a}{n} T^{n}=(1+T)^{a}-1
$$

for $a \in \mathbb{Z}_{p}$.

## 3. Local Class Field Theory

Let $\pi \in \mathcal{O}_{K}$ be a fixed prime element. Since $L_{\pi}$ is totally ramified over $K$, it is linearly disjoint from $K^{\mathrm{nr}}$ over $K$, and the Galois group $\operatorname{Gal}\left(L_{\pi} K^{\mathrm{nr}} \mid K\right)$ is the product of $G_{\pi}=\operatorname{Gal}\left(L_{\pi} \mid K\right)$ and $\operatorname{Gal}\left(K^{\mathrm{nr}} \mid K\right)$. For each prime $\pi$ in $\mathcal{O}_{K}$, we can therefore define a homomorphism

$$
\rho_{\pi}: K^{\times} \longrightarrow \operatorname{Gal}\left(L_{\pi} K^{\mathrm{nr}} \mid K\right)
$$

such that
(1) For each unit $u \in \mathcal{O}_{K}^{\times}$, the automorphism $\rho_{\pi}(u)$ is the identity on $K^{\mathrm{nr}}$, and on $L_{\pi}$ the reciprocal $\tau_{u}^{-1}$ of the element $\tau_{u} \in G_{\pi}$ corresponding to $u$ by the isomorphism of the theorem; and
(2) $\rho_{\pi}(\pi)$ is the identity on $L_{\pi}$ and is the Frobenius automorphism $\sigma$ on $K^{\mathrm{nr}}$.

Thus for an arbitrary element $a=u \pi^{m} \in K^{\times}$we have, by definition:

$$
\rho_{\pi}(a)=\sigma^{m} \text { on } K^{\mathrm{nr}}
$$

and

$$
\lambda^{\rho_{\pi}(a)}=\left[u^{-1}\right]_{f}(\lambda) \text { for } \lambda \in \Lambda_{f} .
$$

Theorem 3.1. - The field $L_{\pi} K^{n r}$ and the homomorphism $\rho_{\pi}$ are independent of $\pi$.
This follows easily from Lemma 1.7. See $[\mathbf{L T}]$ for a detailed proof.
Corollary 3.2. - The field $L_{\pi} K^{n r}$ is the maximal abelian extension of $K$, and $\rho_{\pi}$ is the reciprocity law homomorphism for it, i.e.,

$$
\rho_{\pi}(a)=\left(a, L_{\pi} K^{n r} \mid K\right)
$$

for every $a \in K^{\times}$.
See $[\mathbf{L T}]$ for a proof.
Remark 3.3. - Note that while both the field $L_{\pi}$ and the reciprocity map $\rho_{\pi}$ can be defined in terms of a Lubin-Tate series alone, the proofs depend heavily on the extra structure given by the associated Lubin-Tate module.

Example 3.4. - In the cyclotomic case we get for $a=u p^{v_{p}(a)} \in \mathbb{Q}_{p}^{\times}$that

$$
\left(a, \mathbb{Q}_{p}(\zeta) \mid \mathbb{Q}_{p}\right)(\zeta-1)=\left[u^{-1}\right]_{f}(\zeta-1)
$$

or

$$
\left(a, \mathbb{Q}_{p}(\zeta) \mid \mathbb{Q}_{p}\right) \zeta=\zeta^{u^{-1}}
$$

if $\zeta=1+\lambda$ is a primitive $p^{m}$-th root of unity or in other words $\lambda=\zeta-1 \in \Lambda_{f, m}$.
Remark 3.5. - There are strong analogies with the classical theory of complex multiplication and explicit reciprocity laws for imaginary quadratic fields. In fact for every single statement presented here, there is an analogous one if one replaces the LubinTate modules by elliptic curves with complex multiplication. See for example [L].

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# 7. FORMAL MODULI OF FORMAL $\mathcal{O}_{K}$-MODULES 

by

Eva Viehmann \& Konstantin Ziegler


#### Abstract

We define formal $\mathcal{O}_{K}$-modules and their heights, following Drinfeld. To describe their universal deformations we introduce a formal cohomology group.

Résumé (Espaces de modules formels de $\mathcal{O}_{K}$-modules formels). - On définit les $\mathcal{O}_{K^{-}}$ modules formels et leurs hauteurs, suivant Drinfeld. Pour décrire leurs déformations universelles, on introduit un groupe de cohomologie formelle.


Notation. - Except in the proof of Lemma 2.1, all constant coefficients of power series are assumed to be 0 .

Acknowledgements. - During the preparation of Section 3 we profited from the talk given by S. Wewers in the ARGOS seminar. We thank I. Vollaard and W. Kroworsch for helpful comments on a preliminary version.

## 1. Formal modules

Let $A, R$ be commutative rings with 1 and $i: A \rightarrow R$ a homomorphism. We also write $a$ instead of $i(a)$ for the image of $a$ under $i$.

## Definition 1.1

1. A formal $A$-module over $R$ is a commutative formal group law $F(X, Y)=X+$ $Y+\cdots \in R[[X, Y]]$ together with a ring homomorphism $\gamma: A \rightarrow \operatorname{End}_{R}(F)$ such that the induced map $A \rightarrow \operatorname{End}_{R}(\operatorname{Lie} F) \cong R$ is equal to the structure map $i$.
2. For $a \in A$ we write $\gamma(a)(X)=[a]_{F}(X)=a X+\cdots \in R[[X]]$ for the corresponding endomorphism of $F$. We will also use the notation $X+{ }_{F} Y$ instead of $F(X, Y)$.

2000 Mathematics Subject Classification. - 14L05, 14B12, 13D10, 14K15.
Key words and phrases. - Formal module, formal group, universal deformation.
3. A homomorphism of formal A-modules over $R$ is a homomorphism $\varphi(X)$ : $F(X, Y) \rightarrow G(X, Y)$ of formal group laws $F(X, Y), G(X, Y)$ over $R$ such that $\varphi \circ \gamma_{F}(a)=\gamma_{G}(a) \circ \varphi$ for all $a \in A$. Denote by $\operatorname{Hom}_{R}(F, G)$ the set of homomorphisms from $F$ to $G$.

Definition 1.2. - For $r \geq 2$ let $\nu_{r}=p$, if $r$ is a power of a prime $p$, and $\nu_{r}=1$ else. Denote by

$$
C_{r}(X, Y)=\frac{1}{\nu_{r}}\left((X+Y)^{r}-X^{r}-Y^{r}\right)
$$

the modified binomial form of degree $r$.
Consider the functor which assigns to every $A$-Algebra $R$ the set of formal $A$ modules over $R$. It is represented by an algebra $\Lambda_{A}$ which is generated by the indeterminate coefficients of the series $F$ and $\gamma(a)$ and whose relations are those which are required by the condition that $(F, \gamma)$ is a formal module. It has a natural grading: the degree of a coefficient is one less than the degree of the corresponding monomial in $X, Y$. It is induced by the action of $\mathbb{G}_{m}$ on $\operatorname{Spf}(A[[t]])$. From this description (or by an elementary calculation) one sees that the grading is compatible with concatenation of power series. The elements of the form $a b$ with $\operatorname{deg} a, \operatorname{deg} b \geq 1$ generate a homogeneous ideal. Let $\tilde{\Lambda}_{A}$ be the quotient with induced grading $\tilde{\Lambda}_{A}=\bigoplus \tilde{\Lambda}_{A}^{n}$.

Denote by $\mathbb{G}_{a, R}$ the additive formal group law over $R$. With the canonical $R$-action $\gamma(a)=a X$, it becomes an $R$-module over $R$.

Lemma 1.3. - If $A$ is an infinite field, then for each formal $A$-module over $A$ there exists a unique isomorphism with $\mathbb{G}_{a, A}$ whose derivative at zero equals 1 . In this case there is a canonical isomorphism $\Lambda_{A} \cong A\left[c_{1}, c_{2}, \ldots\right]$ as graded algebras where $\operatorname{deg} c_{i}=i$.

To prove this lemma, one explicitly computes the desired isomorphism, compare [D, Prop. 1.2]. The $c_{i}$ correspond to the coefficients of a homomorphism to the additive formal group law together with the standard $A$-module structure.

From now on let $K$ be a complete discretely valued field with finite residue field $\mathbb{F}_{q}$, where $q=p^{l}$ for some prime $p$. Denote by $\mathcal{O}_{K}$ the ring of integers of $K$. Let $\pi$ be a uniformizer.

Theorem 1.4. - $\Lambda_{\mathcal{O}_{K}}$ and $\mathcal{O}_{K}\left[g_{1}, g_{2}, \ldots\right]$ are non-canonically isomorphic as graded algebras where $\operatorname{deg} g_{i}=i$.

Proof. - First we show that $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1} \cong \mathcal{O}_{K}$ as $\mathcal{O}_{K}$-modules for all $n \geq 2$. For each $i$ let $F_{i}$ and $[a]_{i}$ denote the polynomials of degree $i$ obtained from the universal formal module by leaving out all summands of higher degree. We write

$$
F_{n}(X, Y)=F_{n-1}(X, Y)+\sum_{i=1}^{n-1} c_{i} X^{i} Y^{n-i}
$$

and

$$
[a]_{n}=[a]_{n-1}+h(a) X^{n} .
$$

Then the $c_{i}$ and $h(a)$ generate $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$. As $F$ is a formal group law, we obtain $\sum_{i=1}^{n-1} c_{i} X^{i} Y^{n-i}=\alpha C_{n}(X, Y)$ (compare [ $\mathbf{H}$, Lemma 1.6.6]). Note that we need here that we consider elements in $\tilde{\Lambda}_{\mathcal{O}_{K}}$ and not in $\Lambda_{\mathcal{O}_{K}}$ itself. In particular, $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$ is generated by $\alpha$ and $h(a)$. The condition that $\gamma: \mathcal{O}_{K} \rightarrow \operatorname{End}(F)$ is a homomorphism implies that modulo $(X, Y)^{n+1}$ we have

$$
\begin{aligned}
{[a b]_{n-1}(X)+h(a b) X^{n}=[a]_{n-1}\left([b]_{n-1}(X)\right.} & \left.+h(b) X^{n}\right)+h(a)(b X)^{n} \\
F_{n-1}\left([a]_{n-1}(X)+h(a) X^{n},[b]_{n-1}(X)+h(b) X^{n}\right) & +\alpha C_{n}(a X, b X) \\
& =[a+b]_{n-1}(X)+h(a+b) X^{n},
\end{aligned}
$$

and

$$
\begin{aligned}
& {[a]_{n-1}\left(F_{n-1}(X, Y)+\alpha C_{n}(X, Y)\right)+h(a)(X+Y)^{n}} \\
& \quad=F_{n-1}\left([a]_{n-1}(X)+h(a) X^{n},[a]_{n-1}(Y)+h(a) Y^{n}\right)+\alpha C_{n}(a X, a Y)
\end{aligned}
$$

In $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$ this leads to the relations

$$
\begin{align*}
a h(b)+b^{n} h(a) & =h(a b)  \tag{1.1}\\
h(a+b)-h(a)-h(b) & =\alpha C_{n}(a, b)  \tag{1.2}\\
\left(a^{n}-a\right) \alpha & = \begin{cases}h(a) & \text { if } n \text { is not a power of a prime } \\
h(a) p^{\prime} & \text { if } n=p^{\prime \prime},\end{cases} \tag{1.3}
\end{align*}
$$

and these are all relations between the generators $\alpha, h(a)$ of $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$. If $n$ is invertible in $\mathcal{O}_{K}$, then (1.3) shows that each $h(a)$ is a multiple of $\alpha$. If $n$ is a power of $p$ (where $q=p^{l}$ ) but not of $q$ itself, then there exists an $a \in \mathcal{O}_{K}$ with $a^{n}-a \notin(\pi)$. From (1.1) we obtain $\left(a^{n}-a\right) h(b)=\left(b^{n}-b\right) h(a)$, thus $h(b)$ is a multiple of $h(a)$. Finally (1.2) shows that $\alpha$ is also a multiple of $h(a)$. Now let $n$ be a power of $q$. By choosing $h(a) \mapsto\left(a^{n}-a\right) / \pi$ and $\alpha \mapsto p / \pi$ we define an epimorphism of $\mathcal{O}_{K}$-modules $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1} \rightarrow \mathcal{O}_{K}$. It is well defined as (1.1)-(1.3) are the only relations of $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$. It remains to prove that $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$ is generated by $h(\pi)$. Let $M=\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1} /(h(\pi))$, and denote by $\bar{x} \in M$ the image of $x \in \tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1}$. Then (1.1) shows that $\pi \overline{h(b)}=\overline{h(\pi b)}=\pi^{n} \overline{h(b)}$, thus $\overline{h(\pi b)}=0$ for all $b \in \mathcal{O}_{K}$. Besides, (1.3) shows $\left(\pi^{n}-\pi\right) \bar{\alpha}=\overline{h(\pi) p}=0$, hence $\pi \bar{\alpha}=0$, and $M$ is an $\mathbb{F}_{q}$-vector space. As $n$ is a power of $q$, (1.1) reduces to $a \overline{h(b)}+b \overline{h(a)}=\overline{h(a b)}$. This shows

$$
\overline{h(a)}=\overline{h\left(a^{n}\right)}=n \overline{a^{n-1} h(a)}=0
$$

for all $a$. Then (1.2) implies that $C_{n}(a, b) \bar{\alpha}=0$ for all $a, b \in \mathbb{F}_{q}$. By [H, Lemma 21.3.2], there is an $x \in \mathbb{F}_{p}$ with $C_{n}(x, 1) \neq 0$ in $\mathbb{F}_{p}$. Thus $\bar{\alpha}=0$ and $M=0$.

Hence in all cases $\tilde{\Lambda}_{\mathcal{O}_{K}}^{n-1} \cong \mathcal{O}_{K}$, and we have an epimorphism of graded algebras $\mathcal{O}_{K}\left[g_{1}, g_{2}, \ldots\right] \rightarrow \Lambda_{\mathcal{O}_{K}}$. Here $g_{i}$ is a lift of a generator of $\tilde{\Lambda}_{\mathcal{O}_{K}}^{i}$. The construction of the isomorphism $\Lambda_{K} \cong K\left[c_{1}, c_{2}, \ldots\right]$ in Lemma 1.3 implies that the canonical morphism $\Lambda_{\mathcal{O}_{K}} \otimes K \rightarrow K\left[c_{1}, c_{2}, \ldots\right]$ which is compatible with the grading is also surjective. Comparing dimensions one sees that the epimorphism $\mathcal{O}_{K}\left[g_{1}, g_{2}, \ldots\right] \rightarrow \Lambda_{\mathcal{O}_{K}}$ is an isomorphism.

## 2. Heights

Let $\mathcal{O}_{K}$ be as above and let $R$ be a local $\mathcal{O}_{K}$-algebra of characteristic $p$ with residue field $k$.

Lemma 2.1. - Let $F, G$ be formal $\mathcal{O}_{K}$-modules over $R$ and let $\alpha \in \operatorname{Hom}_{R}(F, G) \backslash\{0\}$. Then there is a unique integer $h=\operatorname{ht}(\alpha) \geq 0$ and $\beta \in R[[X]]$ with $\alpha(X)=\beta\left(X^{q^{h}}\right)$ and $\beta^{\prime}(0) \neq 0$. The integer $h$ is called the height $\operatorname{ht}(\alpha)$ of $\alpha$.

This lemma is analogous to the corresponding result over a field, compare $[\mathbf{H}$, 18.3.1]. For $\alpha=0$ we set $\operatorname{ht}(\alpha)=\infty$.

Proof. - We first show that $\alpha(X)=\beta\left(X^{p^{n}}\right)$ for some $\beta$ with $\beta^{\prime}(0) \neq 0$. To do this we assume $\alpha(X) \neq 0$ with $(\partial \alpha / \partial X)(0)=0$ and show that $\alpha(X)=\beta\left(X^{p}\right)$ for some homomorphism $\beta$ of (not necessarily the same) formal group laws. The claim then follows by induction.

Partial differentiation of $\alpha(F(X, Y))=G(\alpha(X), \alpha(Y))$ with respect to $Y$ gives

$$
\frac{\partial \alpha}{\partial X}(F(X, Y)) \frac{\partial F}{\partial Y}(X, Y)=\frac{\partial G}{\partial Y}(\alpha(X), \alpha(Y)) \frac{\partial \alpha}{\partial X}(Y)
$$

Substituting $Y=0$ and using $(\partial \alpha / \partial X)(0)=0$ we obtain

$$
\frac{\partial \alpha}{\partial X}(X) \frac{\partial F}{\partial Y}(X, 0)=0
$$

As $(\partial F / \partial Y)(X, 0)=1+a_{1} X+\cdots \in R[[X]]^{\times}$, we obtain $\frac{\partial \alpha}{\partial X}(X)=0$. Hence $\alpha(X)=$ $\beta\left(X^{p}\right)$ for some $\beta \in R[[X]]$. Let $\sigma_{*} F$ be the formal group law obtained from $F$ by raising each coefficient to the $p$ th power. Then an easy calculation shows that $\beta$ is a homomorphism from $\sigma_{*} F$ to $G$.

We now have to show that $p^{n}$ is a power of $q$. Let $a \in \mathcal{O}_{K}$. Then

$$
[a]_{G}(\alpha(X))=\alpha\left([a]_{F}(X)\right)=\beta^{\prime}(0) i(a)^{p^{n}} X^{p^{n}}+\cdots
$$

and on the other hand

$$
[a]_{G}(\alpha(X))=\beta^{\prime}(0) i(a) X^{p^{n}}+\cdots .
$$

This implies $\beta^{\prime}(0)\left(i(a)-i\left(a^{p^{n}}\right)\right)=0$ with $\beta^{\prime}(0) \neq 0$, hence $i(a)-i\left(a^{p^{n}}\right)=i\left(a-a^{p^{n}}\right)$ maps to 0 in $k$. Thus $a^{p^{n}}=a$ for all $a \in \mathbb{F}_{q}$ and $p^{n}$ is a power of $q$.

Definition 2.2. - The height of a formal $\mathcal{O}_{K}$-module $F$ over $R$ is

$$
\operatorname{ht}(F)= \begin{cases}h & \text { if }[\pi]_{F} \text { has height } h \\ \infty & \text { if }[\pi]_{F}=0\end{cases}
$$

Remark 2.3. - This definition is different from the definition of height of a formal module given in $[\mathbf{H}]$, where it is defined as the height of the reduction of the module over the residue field.

Lemma 2.4. - Let $R$ be as above and let $\left(F, \gamma_{F}\right)$ be the formal $\mathcal{O}_{K}$-module corresponding to a homomorphism $\varphi: \Lambda_{\mathcal{O}_{K}} \rightarrow R$. Then $\operatorname{ht}(F)=\min \left\{i \mid \varphi\left(g_{q^{i}-1}\right) \neq 0\right\}$.

Proof. - In the proof of Theorem 1.4 we identified the generator $g_{q^{i}-1}$ of $\tilde{\Lambda}_{\mathcal{O}_{K}}^{q^{i}-1}$ with the coefficient of $X^{q^{i}}$ of $[\pi](X)$.

The following lemma reduces the examination of formal modules over fields and of their deformations to formal modules of an especially simple form. For a proof see [D, Prop. 1.7].

Lemma 2.5. - Let $(F, \gamma)$ be a formal $\mathcal{O}_{K}$-module of height $h<\infty$ over a separably closed field $k$ of characteristic $p>0$. Then $F$ is isomorphic to a formal module ( $F^{\prime}, \gamma^{\prime}$ ) over $k$ with

$$
\begin{aligned}
F^{\prime}(X, Y) & \equiv X+Y \quad\left(\bmod \operatorname{deg} q^{h}\right) \\
{[a]_{F^{\prime}}(X) } & \equiv a X \quad\left(\bmod \operatorname{deg} q^{h}\right), \\
{[\pi]_{F^{\prime}}(X) } & =X^{q^{h}} .
\end{aligned}
$$

Such modules are called normal modules.
Fix an integer $h>1$ and let $F_{0}$ be a formal $\mathcal{O}_{K}$-module of height $h$ over $k$. Assume that $R$ is a local artinian $\mathcal{O}_{K}$-algebra with maximal ideal $\mathfrak{m}$ and residue field $k$. Let $I \triangleleft R$ be an ideal. We set $\bar{R}=R / I$. If $F$ is a lift of $F_{0}$ over $R$, we set $\bar{F}:=F \otimes_{R} \bar{R}$.

Lemma 2.6. - Let $F, G$ be lifts of $F_{0}$ over $R$. Then the reduction map

$$
\begin{equation*}
\operatorname{Hom}_{R}(F, G) \rightarrow \operatorname{Hom}_{\bar{R}}(\bar{F}, \bar{G}) \tag{2.1}
\end{equation*}
$$

is injective.
Proof. - The reduction map in (2.1) is the composition of finitely many maps

$$
\operatorname{Hom}_{R_{n+1}}\left(F \otimes R_{n+1}, G \otimes R_{n+1}\right) \rightarrow \operatorname{Hom}_{R_{n}}\left(F \otimes R_{n}, G \otimes R_{n}\right),
$$

where $R_{n}=R / I_{n}$ with $I_{n}=I \cap \mathfrak{m}^{n}$. We may therefore assume that $\mathfrak{m} \cdot I=0$. Then $I$ is a finite dimensional $k$-vector space, and we have $I^{2}=0$. Let $\alpha(X)=a_{1} X+a_{2} X^{2}+\ldots$ be a homomorphism from $F$ to $G$ such that $\alpha(X) \equiv 0(\bmod I)$. We get

$$
\alpha\left([\pi]_{F}(X)\right)=[\pi]_{G}(\alpha(X))=0 .
$$

Since $\operatorname{ht}\left(F_{0}\right)<\infty$, we have $[\pi]_{F}(X) \neq 0(\bmod \mathfrak{m})$, thus $\alpha=0$ which proves the lemma.

From now on we may consider $\operatorname{Hom}_{R}(F, G)$ as a subset of $\operatorname{Hom}_{\bar{R}}(\bar{F}, \bar{G})$.

## 3. Deformations of modules, formal cohomology

Let $F$ be a formal $\mathcal{O}_{K}$-module of height $h<\infty$ over $k$, and let $M$ be a finite dimensional $k$-vector space. A symmetric 2-cocycle of $F$ with coefficients in $M$ is a
collection of power series $\Delta(X, Y) \in M[[X, Y]]$ and $\left\{\delta_{a}(X) \in M[[X]]\right\}_{a \in \mathcal{O}_{K}}$ satisfying

$$
\begin{align*}
\Delta(X, Y) & =\Delta(Y, X)  \tag{3.1}\\
\Delta(X, Y)+\Delta(F(X, Y), Z) & =\Delta(Y, Z)+\Delta(X, F(Y, Z))  \tag{3.2}\\
\delta_{a}(X)+\delta_{a}(Y)+\Delta\left([a]_{F}(X),[a]_{F}(Y)\right) & =i(a) \Delta(X, Y)+\delta_{a}(F(X, Y))  \tag{3.3}\\
\delta_{a}(X)+\delta_{b}(X)+\Delta\left([a]_{F}(X),[b]_{F}(X)\right) & =\delta_{a+b}(X)  \tag{3.4}\\
i(a) \delta_{b}(X)+\delta_{a}\left([b]_{F}(X)\right) & =\delta_{a b}(X) . \tag{3.5}
\end{align*}
$$

For any $\Psi \in M[[X]]$, the coboundary of $\Psi$ is the symmetric 2-cocycle ( $\Delta^{\Psi},\left\{\delta_{a}^{\Psi}\right\}$ ) with

$$
\begin{align*}
\Delta^{\Psi}(X, Y) & =\Psi(F(X, Y))-\Psi(X)-\Psi(Y)  \tag{3.6}\\
\delta_{a}^{\Psi}(X) & =\Psi\left([a]_{F}(X)\right)-i(a) \Psi(X) . \tag{3.7}
\end{align*}
$$

The coboundaries form a subspace of the vector space $Z^{2}(F, M)$ of symmetric 2cocycles. The quotient of the symmetric 2 -cocycles by the coboundaries is a $k$-vector space denoted $H^{2}(F, M)$.

The following lemma is due to Keating, see [K2, Lemma 2.1].
Lemma 3.1. - $A$ cocycle $\left(\Delta ;\left\{\delta_{a}\right\}\right) \in Z^{2}(F, M)$ is zero if and only if $\delta_{\pi}(X)=0$.
Proof. - If the cocyle is zero, then clearly $\delta_{\pi}(X)=0$. Assume conversely that $\delta_{\pi}(X)=0$. Substituting $a=\pi$ in (3.3) gives

$$
\Delta\left([\pi]_{F}(X),[\pi]_{F}(Y)\right)=0
$$

since $\delta_{\pi}(X)=0$ and $i(\pi)=0$. As $[\pi]_{F}(X) \neq 0$, this implies $\Delta(X, Y)=0$. Condition (3.5) with $a=\pi$ together with $\delta_{\pi}(X)=0$ shows $\delta_{\pi b}(X)=0$. The same formula with $b=\pi$ and $a$ arbitrary gives $\delta_{a}\left([\pi]_{F}(X)\right)=0$. This implies that $\delta_{a}(X)=0$, so all components of the cocycle are zero.

In the following let $R$ denote a local artinian $\mathcal{O}_{K}$-algebra with maximal ideal $\mathfrak{m}$ and residue field $k$. Let $I \subseteq \mathfrak{m}$ be an ideal with $\mathfrak{m} I=0$. Then $I$ is a $k$-vector space. We set $\bar{R}=R / I$. If $F_{0}$ is a formal module over $k$ and $F$ is a lift of $F_{0}$ over $R$, denote by $\bar{F}=F \otimes_{R} \bar{R}$ the reduction modulo $I$. The reduction modulo $\mathfrak{m}$ of power series over $R$ is denoted by $\cdot{ }^{*}$.

Proposition 3.2. - In the setting above let $F_{0}$ be a formal $\mathcal{O}_{K}$-module over $k$ and let $F, G \in R[[X, Y]]$ be formal $\mathcal{O}_{K}$-modules with $F^{*}=G^{*}=F_{0}$. For $\varphi(X) \in R[[X]]$ let $\bar{\varphi} \in \bar{R}[[X]]$ be the image. Assume that $\bar{\varphi}$ is a homomorphism from $\bar{F}$ to $\bar{G}$. Then

1. There is an element of $Z^{2}\left(F_{0}, I\right)$ defined by

$$
\begin{aligned}
\Delta & =\varphi(F(X, Y))-{ }_{G} \varphi(X)-{ }_{G} \varphi(Y) \\
\delta_{a} & =\varphi\left([a]_{F}(X)\right)-{ }_{G}[a]_{G}(\varphi(X)) .
\end{aligned}
$$

2. $\left(\Delta ;\left\{\delta_{a}\right\}_{a}\right)=0$ if and only if $\varphi(X) \in \operatorname{Hom}_{R}(F, G)$.
3. The class of $\left(\Delta ;\left\{\delta_{a}\right\}_{a}\right)$ in $H^{2}\left(F_{0}, I\right)$ is independent of the choice of the lift $\varphi$ of $\bar{\varphi}$. It vanishes if and only if $\bar{\varphi} \in \operatorname{Hom}_{R}(F, G) \subseteq \operatorname{Hom}_{\bar{R}}(\bar{F}, \bar{G})$. If $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ is the coboundary of $\psi$, the lift of $\bar{\varphi}$ to a homomorphism over $R$ is given by $\varphi-{ }_{G} \psi$.

Proof. - Applying $\varphi$ to the left hand side of the associativity law for $F$

$$
\begin{equation*}
\left(X+{ }_{F} Y\right)+{ }_{F} Z=X+{ }_{F}\left(Y+_{F} Z\right) \tag{3.8}
\end{equation*}
$$

and using the definition of $\Delta$, we get

$$
\begin{equation*}
\varphi(X)+{ }_{G} \varphi(Y)+{ }_{G} \varphi(Z)+{ }_{G} \Delta(X, Y)+{ }_{G} \Delta\left(X+{ }_{F} Y, Z\right) . \tag{3.9}
\end{equation*}
$$

Applying $\varphi$ to the right hand side of (3.8), we get

$$
\begin{equation*}
\varphi(X)+{ }_{G} \varphi(Y)+{ }_{G} \varphi(Z)+{ }_{G} \Delta\left(X, Y+_{F} Z\right)+_{G} \Delta(Y, Z) \tag{3.10}
\end{equation*}
$$

From (3.10) and (3.9) we obtain

$$
\begin{equation*}
\Delta(X, Y)+_{G} \Delta\left(X+_{F} Y, Z\right)=\Delta\left(X, Y+_{F} Z\right)+_{G} \Delta(Y, Z) . \tag{3.11}
\end{equation*}
$$

Using the assumption $\mathfrak{m} \cdot I=0$, we see that (3.11) implies the second cocycle rule

$$
\begin{equation*}
\Delta(X, Y)+\Delta\left(X+F_{0} Y, Z\right)=\Delta\left(X, Y+_{F_{0}} Z\right)+\Delta(Y, Z) \tag{3.12}
\end{equation*}
$$

The other cocycle rules are proved in a similar manner, replacing (3.8) by the commutativity resp. the distributivity law of $F$. This proves 1 .

Part 2 of the proposition is a straightforward consequence of the definition of $\left(\Delta ;\left\{\delta_{a}\right\}\right)$. To prove 3., we continue with the notation used in the proof of 1 . Let $\varphi^{\prime}(X)$ be another lift of $\bar{\varphi}$, and let $\left(\Delta^{\prime} ;\left\{\delta_{a}^{\prime}\right\}\right)$ be the cocycle it defines. We can write $\varphi^{\prime}=\varphi+{ }_{G} \psi$, with $\psi \in I[[X]]$. Then

$$
\begin{aligned}
\varphi^{\prime}\left([\pi]_{F}(X)\right) & =[\pi]_{G}(\varphi(X))+{ }_{G} \delta_{\pi}(X)+{ }_{G} \psi\left([\pi]_{F}(X)\right) \\
& =[\pi]_{G}\left(\varphi^{\prime}(X)\right)+{ }_{G}\left(\delta_{\pi}(X)+{ }_{G} \psi\left([\pi]_{F}(X)\right)\right) .
\end{aligned}
$$

For the second equality we have used that $I \mathfrak{m}=0$. We conclude that $\delta_{\pi}^{\prime}(X)-\delta_{\pi}(X)=$ $\psi\left([\pi]_{F}(X)\right)$ is the $\pi$-component of the coboundary of $\psi$. Then Lemma 3.1 implies that the two cocycles differ by the coboundary of $\psi$. Hence $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ and $\left(\Delta^{\prime} ;\left\{\delta_{a}^{\prime}\right\}\right)$ lie in the same class in $H^{2}\left(F_{0}, I\right)$. It follows from 2. that this class vanishes if and only if $\bar{\varphi} \in \operatorname{Hom}_{R}(F, G)$. This completes the proof of 3 . and the proposition.
Lemma 3.3. - In the setting of Proposition 3.2 let $(F, \gamma)$ be a lift of $F_{0}$ to $R$ and let $\bar{F}$ be the reduction to $\bar{R}$.

1. Proposition 3.2 defines a bijection between deformations of $\bar{F}$ to $R$ and cocycles in $Z^{2}\left(F_{0}, I\right)$. Its inverse is given by assigning to $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ the deformation $F_{\Delta}(X, Y)=X+{ }_{F} Y+{ }_{F} \Delta(X, Y)$ and $\gamma_{\delta}(a)=\gamma(a)+_{F} \delta_{a}$.
2. Two cocycles are in the same cohomology class if and only if the corresponding deformations are isomorphic via an isomorphism which lifts the identity of $\bar{F}$.

Proof. - For the first assertion we have to check that $\left(F_{\Delta}, \gamma_{\delta}\right)$ is a formal module. From $I^{2}=0$ we obtain that the equations (3.1) to (3.5) also hold with $F$ replaced by $F_{\Delta}$. These equations immediately imply that $\left(F_{\Delta}, \gamma_{\delta}\right)$ is a formal module. For $F_{\Delta}, F$ and $\varphi=X$ we obtain the cocycle $\left(\Delta,\left\{\delta_{a}\right\}\right)$. Then the second assertion follows from Proposition 3.2, 3.

Corollary 3.4. - Let $F_{0}, R$, and $I$ be as above with $\operatorname{char}(R)=p, \operatorname{ht}\left(F_{0}\right)=h$, and $\left(\Delta ;\left\{\delta_{a}\right\}\right) \in Z^{2}\left(F_{0}, I\right)$.

1. Let $g \leq h$. Then $\delta_{\pi}(X) \equiv 0\left(\bmod X^{q^{g-1}+1}\right)$ if and only if $\delta_{\pi} \in I\left[\left[X^{q^{g}}\right]\right]$.
2. The following are equivalent:
(a) The cocycle $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ is the coboundary of some $\psi(X) \in I[[X]]$.
(b) $\delta_{\pi} \in I\left[\left[X^{q^{h}}\right]\right]$.
(c) Let $(F, \gamma)$ be a lift of $F_{0}$ to a formal $\mathcal{O}_{K}$-module over $R$. Then the identity of $\bar{F}$ lifts to an isomorphism between $(F, \gamma)$ and $\left(F_{\Delta}, \gamma_{\delta}\right)$.
If these conditions are satisfied, $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ is the coboundary of $\psi=d \circ \beta^{-1}$ where $d\left(X^{q^{h}}\right)=\delta_{\pi}(X)$ and $\beta\left(X^{q^{h}}\right)=[\pi]_{F_{0}}(X)$.

Proof. - If $\delta_{\pi}(X) \equiv 0\left(\bmod X^{q^{g-1}+1}\right)$ then

$$
[\pi]_{F_{\Delta}}(X)=\delta_{\pi}(X)+_{F}[\pi]_{F}(X) \equiv 0 \quad\left(\bmod \left(X^{q^{g-1}+1}\right)\right),
$$

thus $\operatorname{ht}\left(F_{\Delta}\right)>g-1$. This shows that $\delta_{\pi}(X)=[\pi]_{F_{\Delta}}(X)-{ }_{F}[\pi]_{F}(X)$ is a power series in $X^{q^{g}}$. The other assertion of 1 . is trivial. The equivalence of (a) and (c) of 2. follows from Lemma 3.3. From Lemma 3.1 we see that $\left(\Delta ;\left\{\delta_{a}\right\}\right)=\left(\Delta^{\psi} ;\left\{\delta_{a}^{\psi}\right\}\right)$ for some $\psi$ if and only if $\delta_{\pi}(X)=\delta_{\pi}^{\psi}(X)=\psi\left([\pi]_{F}(X)\right)=\psi\left([\pi]_{F_{0}}(X)\right)$. Here the last two equations follow from $I \mathfrak{m}=0$. As $\operatorname{ht}\left(F_{0}\right)=h$, this implies (b). On the other hand assume (b) and let $d\left(X^{q^{h}}\right)=\delta_{\pi}(X)$ and $\beta\left(X^{q^{h}}\right)=[\pi]_{F_{0}}(X)$. Then the $\pi$-component of the coboundary of $\psi=d \circ \beta^{-1}$ is $\delta_{\pi}$.

Let $\hat{\mathcal{O}}_{K}^{n r}$ be the completion of the maximal unramified extension of $\mathcal{O}_{K}$. Denote by $\hat{\mathcal{O}}_{K}^{n r}[[t]]=\hat{\mathcal{O}}_{K}^{n r}\left[\left[t_{1}, \ldots, t_{h-1}\right]\right]$ the power series ring over $\hat{\mathcal{O}}_{K}^{n r}$ in $h-1$ variables. Let $k=\hat{\mathcal{O}}_{K}^{n r} /(\pi)$.

Lemma 3.5. - Let $\left(F, \gamma_{F}\right)$ be a normal $\mathcal{O}_{K}$-module over $k$ of height $h<\infty$. Then there exists a formal $\mathcal{O}_{K}$-module $(\Gamma, \gamma)$ over $\hat{\mathcal{O}}_{K}^{n r}[[t]]$ which over $k$ reduces to $F$ with the following property: For $1 \leq i \leq h-1$ denote by $\left(\Gamma_{i}, \gamma_{i}\right)$ the reduction to $\hat{\mathcal{O}}_{K}^{n r}[[t]] /\left(t_{1}, \ldots, t_{i-1}\right)$. Then

$$
\begin{equation*}
\gamma_{i}(\pi)(X) \equiv \pi X+t_{i} X^{q^{i}} \quad\left(\bmod \operatorname{deg}\left(q^{i}+1\right)\right) . \tag{3.13}
\end{equation*}
$$

Proof. - The module $F$ corresponds to a map $\bar{\varphi}: \Lambda_{\mathcal{O}_{K}} \cong \mathcal{O}_{K}\left[g_{1}, g_{2}, \ldots\right] \rightarrow k$ with $g_{i} \mapsto 0$ for all $i<q^{h}-1$. Let $\varphi: \Lambda_{\mathcal{O}_{K}} \rightarrow \hat{\mathcal{O}}_{K}^{n r}$ be a lift with the same property. We choose

$$
f_{i}= \begin{cases}t_{j} & \text { if } i=q^{j}-1 \text { with } 1 \leq j<h-1 \\ \varphi\left(g_{i}\right) & \text { else. }\end{cases}
$$

Let $\Gamma$ be the formal $\mathcal{O}_{K}$-module corresponding to the map $\Lambda_{\mathcal{O}_{K}} \rightarrow \hat{\mathcal{O}}_{K}^{n r}[[t]]$ which maps $g_{i}$ to $f_{i}$. Then for $\left(\Gamma_{i}, \gamma_{i}\right)$ we see that $g_{q^{i}-1}$ is the first generator which is mapped to a nonzero element in $\hat{\mathcal{O}}_{K}^{n r}[[t]] /\left(t_{1}, \ldots, t_{i-1}\right)$. From the description of $\tilde{\Lambda}_{\mathcal{O}_{K}}^{q^{i}-1}$ in the proof of Theorem 1.4 we see that $\gamma_{i}(\pi)(X)$ has the desired form.

Note that a proof of this result can also be found in [GH, Section 12].
Let $\left(F, \gamma_{F}\right)$ be a normal formal $\mathcal{O}_{K}$-module of height $h<\infty$ over $k$. Let $(\Gamma, \gamma)$ be the deformation over $\hat{\mathcal{O}}_{K}^{n r}[[t]]$ defined in Lemma 3.5. Let $\left(\Gamma^{i}, \gamma^{i}\right)$ be the reduction of $(\Gamma, \gamma)$ to $k\left[\left[t_{i}\right]\right] /\left(t_{i}\right)^{2}=R_{i}$ and let $\left(F, \gamma_{F}\right)_{R_{i}}$ be the base change of $\left(F, \gamma_{F}\right)$ to $R_{i}$.

Proposition 3.6. - For $F$ as above we have $\operatorname{dim}_{k} H^{2}(F, k)=h-1$. The cocycles $\left(\Delta^{i} ;\left\{\delta_{a}^{i}\right\}\right)$ associated to the pairs of deformations $\left(F, \gamma_{F}\right)_{R_{i}}$ and $\left(\Gamma^{i}, \gamma^{i}\right)$ with values in $t_{i} R_{i} \cong k$ satisfy

$$
\begin{equation*}
\delta_{\pi}^{i} \equiv t_{i} X^{q^{i}} \quad\left(\bmod \operatorname{deg} q^{i}+1\right) \tag{3.14}
\end{equation*}
$$

Their classes form a basis for $H^{2}(F, k)$.
Proof. - Equation (3.14) immediately follows from (3.13). Corollary 3.4, 2. shows that the $\pi$-components of coboundaries are power series in $X^{q^{h}}$. Thus (3.14) implies that the classes of the cocycles $\left(\Delta^{i} ;\left\{\delta_{a}^{i}\right\}\right)$ are linearly independent in $H^{2}(F, k)$. Let $\left(\Delta ;\left\{\delta_{a}\right\}\right) \in H^{2}(F, k)$. Then by Corollary 3.4, 1., $\delta_{\pi}$ is of the form $\beta\left(X^{q^{g}}\right)$ with $\beta^{\prime}(0) \neq 0$. If $g<h$ we subtract a suitable multiple of $\left(\Delta^{g} ;\left\{\delta_{a}^{g}\right\}\right)$ to annihilate the coefficient of $X^{q^{g}}$. In this way we can inductively represent the cocycle ( $\left.\Delta ;\left\{\delta_{a}\right\}\right)$ as a linear combination of the $\left(\Delta^{i} ;\left\{\delta_{a}^{i}\right\}\right)$ plus a cocycle whose $\pi$-component is congruent to 0 modulo $X^{q^{h-1}+1}$. Hence by Corollary 3.4, the cohomology class is a linear combination of the classes of the $\left(\Delta^{i} ;\left\{\delta_{a}^{i}\right\}\right)$.

Definition 3.7. - Let $R$ be a local ring with maximal ideal $\mathfrak{m}$. For a power series $f$ with coefficients in $R$ let $f^{*}$ be the reduction modulo $\mathfrak{m}$. A $*$-isomorphism between $\mathcal{O}_{K}$-modules $F, G$ over $R$ is an isomorphism $\varphi \in \operatorname{Hom}_{R}(F, G)$ with $\varphi^{*}(X)=X$.

Let $F$ be a fixed $\mathcal{O}_{K}$-module of height $h<\infty$ over $k=\hat{\mathcal{O}}_{K}^{n r} /(\pi)$. We consider the functor $\mathcal{D}_{F}$ which assigns to each complete local noetherian $\hat{\mathcal{O}}_{K}^{n r}$-algebra $R$ with residue field $k$ and maximal ideal $\mathfrak{m}$ the set of $*$-isomorphism classes of formal $\mathcal{O}_{K^{-}}$ modules over $R$ that modulo $\mathfrak{m}$ reduce to $F$.

Theorem 3.8 (Universal deformation). - Let $\left(F, \gamma_{F}\right)$ be an $\mathcal{O}_{K}$-module over $k$ of height $h<\infty$. Then $\mathcal{D}_{F}$ is represented by $\hat{\mathcal{O}}_{K}^{n r}[[t]]$.

Proof. - As $k$ is separably closed, Lemma 2.5 shows that we may assume ( $F, \gamma_{F}$ ) to be normal. Let $(\Gamma, \gamma)$ be the deformation over $\hat{\mathcal{O}}_{K}^{n r}[[t]]$ of Lemma 3.5. Let $\left(\Phi, \gamma_{\Phi}\right) \in$ $\mathcal{D}_{F}(R)$ for some complete local noetherian $\hat{\mathcal{O}}_{K}^{n r}$-algebra $R$ with residue field $k$ and maximal ideal $\mathfrak{m}$. As $R$ is complete, it is enough to show that for each $r \in \mathbb{N}$ the following holds: If the projection $\Phi_{r}$ of $\Phi$ to $R / \mathfrak{m}^{r}$ corresponds to a homomorphism $\varphi_{r}: \hat{\mathcal{O}}_{K}^{n r}[[t]] \rightarrow R / \mathfrak{m}^{r}$, then there is a unique lift $\varphi_{r+1}: \hat{\mathcal{O}}_{K}^{n r}[[t]] \rightarrow R / \mathfrak{m}^{r+1}$ of $\varphi_{r}$ corresponding to $\Phi_{r+1}$.

Let $\psi$ be any lift of $\varphi_{r}$ to $R / \mathfrak{m}^{r+1}[[X]]$. Then the pair of deformations $\psi(\Gamma, \gamma)$, $\left(\Phi_{r+1}, \gamma_{\Phi_{r+1}}\right)$ corresponds to an element of $H^{2}\left(F, \mathfrak{m}^{r} / \mathfrak{m}^{r+1}\right)$, hence to a uniquely defined linear combination of the $\Delta^{i}$ with coefficients $a_{i}$ in $\mathfrak{m}^{r} / \mathfrak{m}^{r+1}$. Let $\varphi_{r+1}\left(t_{i}\right)=$ $\psi\left(t_{i}\right)+a_{i}$. Then by Corollary 3.4, the deformations $\Phi_{r+1}$ and $\varphi_{r+1}(\Gamma, \gamma)$ of $F$ over $R / \mathfrak{m}^{r+1}$ are isomorphic via an isomorphism which lifts the given isomorphism over $R / \mathfrak{m}^{r}$. As the classes of the $\Delta^{i}$ are linearly independent, $\varphi_{r+1}$ is unique.

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# 8. CANONICAL AND QUASI-CANONICAL LIFTINGS 

by<br>Stefan Wewers


#### Abstract

The present note gives a detailed account of the paper of Gross on canonical and quasi-canonical liftings. These are liftings of formal $\mathcal{O}$-modules with extra endomorphisms, and thus correspond to CM-points in the universal deformation space.


Résumé (Relèvements canoniques et quasi-canoniques). - Nous donnons un exposé détaillé des travaux de Gross sur les relèvements canoniques et quasi-canoniques des $\mathcal{O}$-modules formels, qui correspondent aux points CM dans l'espace de déformations universel.

The present note gives a detailed account of Gross' paper [G] on canonical and quasi-canonical liftings. We make heavy use of results of Lubin and Tate [LT2] and Drinfeld [D] which are reviewed in [VZ]. All the results presented here have been generalized to the case of arbitrary finite height by J. K. Yu $[\mathbf{Y u}]$.

I thank Eva Viehmann, Inken Vollaard and Michael Rapoport for careful proofreading and helpful discussions.

## 1. Canonical lifts

In this section we study canonical lifts of a formal $\mathcal{O}_{K}$-module of height two with respect to a quadratic extension $L / K$. In particular, we prove the first main result of $[\mathbf{G}]$ which computes the endomorphism ring of the reduction of a canonical lift modulo some power of the prime ideal of $\mathcal{O}_{K}$.

2000 Mathematics Subject Classification. - 14L05, 14K22.
Key words and phrases. - Formal $\mathcal{O}$-modules, canonical liftings, Lubin-Tate theory.
1.1. Throughout this note, $K$ denotes a field which is complete with respect to a discrete valuation $v$, and whose residue class field is finite, with $q=p^{f}$ elements. We denote by $\mathcal{O}_{K}$ the ring of integers of $K$. We fix a prime element $\pi$ of $K$, and we assume that $v(\pi)=1$.

Let $i: \mathcal{O}_{K} \rightarrow R$ be an $\mathcal{O}_{K}$-algebra. Recall that a formal $\mathcal{O}_{K}$-module over $R$ is given by a commutative formal group law $F(X, Y)=X+Y+\cdots \in R[[X, Y]$ together with a ring homomorphism $\gamma: \mathcal{O}_{K} \rightarrow \operatorname{End}_{R}(F)$ such that the induced map $\mathcal{O}_{K} \rightarrow \operatorname{End}_{R}(\operatorname{Lie} F) \cong R$ is equal to the structure map $i$. Whenever this is not likely to be confusing, we will omit the maps $i$ and $\gamma$ from the notation. Given an element $a \in \mathcal{O}_{K}$, we write $[a]_{F}(X)=i(a) X+\cdots \in R[[X]$ for the corresponding endomorphism of $F$.

If $F_{1}, F_{2}$ are two formal $\mathcal{O}_{K}$-modules over $R$, we write $\operatorname{Hom}_{R}\left(F_{1}, F_{2}\right)$ for the group of homomorphisms $\alpha: F_{1} \rightarrow F_{2}$ of formal $\mathcal{O}_{K}$-modules, i.e., $\mathcal{O}_{K}$-linear homomorphisms of formal groups. Similarly, $\operatorname{End}_{R}(F)$ denotes the (in general noncommutative) ring of $\mathcal{O}_{K^{-}}$linear endomorphisms of $F$. Note that $\operatorname{End}_{R}(F)$ is an $\mathcal{O}_{K^{-}}$ algebra.
1.2. Let $k$ be an algebraic closure of the residue class field of $\mathcal{O}_{K}$. We regard $k$ as an $\mathcal{O}_{K}$-algebra, and write $\bar{a} \in k$ for the image of an element $a \in \mathcal{O}_{K}$.

Let $G$ be a formal $\mathcal{O}_{K}$-module over $k$ and let $\alpha \in k[[X]]$ be an endomorphism of $G$, with $\alpha \neq 0$. By [VZ, Lemma 2.1], there exists an integer $h=\operatorname{ht}(\alpha) \geq 0$, called the height of $\alpha$, such that $\alpha(X)=\beta\left(X^{q^{h}}\right)$, with $\beta^{\prime}(0) \neq 0$. It is easy to check that the function ht: $\operatorname{End}_{k}(G) \rightarrow \mathbb{Z}_{\geq 0} \cup\{\infty\}$ (we set ht $(0):=\infty$ ) is a valuation on the $\mathcal{O}_{K}$-algebra $\operatorname{End}_{k}(G)$. We say that the formal $\mathcal{O}_{K}$-module $G$ has height $h$, if the endomorphism $[\pi]_{G}$ has height $h$. In other words, the restriction of the valuation ht $v i a$ the structure map $\mathcal{O}_{K} \rightarrow \operatorname{End}_{k}(G)$ is equal to $h^{-1} \cdot v$.

We recall the following fundamental result.
Theorem 1.1. - For each natural number $h$, there exists a formal $\mathcal{O}_{K}$-module $G$ over $k$ of height $h$. It is unique up to isomorphism. The ring $\operatorname{End}_{k}(G)$ is isomorphic to the maximal order $\mathcal{O}_{D}$ of a division algebra $D$ of dimension $h^{2}$ over $K$, with invariant $\operatorname{inv}(D)=1 / h$.

Proof. - (Compare with [D], Proposition 1.7.) The existence of $G$ follows from Lubin-Tate theory, as follows. Let $L / K$ be the unramified extension of degree $h$. Extend the algebra map $\mathcal{O}_{K} \rightarrow k$ to $\mathcal{O}_{L}$, which gives $k$ the structure of an $\mathcal{O}_{L^{-}}$ algebra. Let $F$ be the Lubin-Tate module of $\mathcal{O}_{L}$ with respect to the prime element $\pi$, i.e., the (unique) formal $\mathcal{O}_{L}$-module over $\mathcal{O}_{L}$ such that $[\pi]_{F}=\pi X+X^{q^{h}}$, see $[\mathbf{L T 1}]$. By restriction, we may regard $F$ as a formal $\mathcal{O}_{K}$-module. Then $G:=F \otimes k$ is a formal $\mathcal{O}_{K}$-module of height $h$ over $k$.

The uniqueness of $G$ is more difficult. See e.g. [H, Theorem 21.9.1].

Let us sketch a proof of the last statement of Theorem 1.1. Set $H:=\operatorname{End}_{k}(G)$. We may assume that $G$ is the reduction to $k$ of the Lubin-Tate module for $\mathcal{O}_{L}$, where $L / K$ is unramified of degree $h$. Since the natural map $\mathcal{O}_{L}=\operatorname{End}(F) \rightarrow H$ is injective (see [VZ, Lemma 2.6]), we have $\mathcal{O}_{L} \subset H$. By construction, the group law $G(X, Y)=X+Y+\ldots$ and the endomorphisms $[a]_{G}(X)=\bar{a} X+\ldots$, for $a \in \mathcal{O}_{K}$, are power series with coefficients in $\mathbb{F}_{q}$. Moreover, we have $[\pi]_{G}(X)=X^{q^{h}}$. Hence the polynomial $\Pi(X):=X^{q}$ defines an element $\Pi \in H$ with $\Pi^{h}=\pi$. One checks that

$$
\Pi\left([a]_{G}(X)\right)=\left[a^{\sigma}\right]_{G}(\Pi(X))
$$

where $\sigma \in \operatorname{Gal}(L / K)$ is the Frobenius. From there, it is easy to see that the subalgebra $\mathcal{O}_{D}:=\mathcal{O}_{L}[\Pi]$ of $H$ is the maximal order of a division algebra $D$ of dimension $h^{2}$ over $K$, with invariant $1 / h$. It remains to be shown that $\mathcal{O}_{D}=H$.

Let $\alpha(X)=\bar{a} X+\ldots$ be an element of $H$. Since $\alpha$ commutes with $[\pi]_{G}(X)=X^{q^{h}}$, the coefficients of $\alpha$ lie in $\mathbb{F}_{q^{h}}=\mathcal{O}_{L} / \pi \mathcal{O}_{L}$. Let $a \in \mathcal{O}_{L}$ be a lift of $\bar{a}$. Then $\alpha-[a]_{G}$ is an endomorphism of $G$ with positive height, and therefore lies in the left ideal $H \cdot \Pi \subset H$. We have shown that the natural map

$$
\mathcal{O}_{D} \longrightarrow H /(H \cdot \Pi)
$$

is surjective. Now the desired equality $\mathcal{O}_{D}=H$ follows from the fact (which is easy to prove) that $H$ is complete with respect to the $\Pi$-adic topology.
1.3. For the rest of this note, we fix a formal $\mathcal{O}_{K^{K}}$-module $G$ of height two over $k$. By Theorem 1.1, $G$ is uniquely determined, up to isomorphism, and $\mathcal{O}_{D}:=\operatorname{End}_{k}(G)$ is the maximal order in a quaternion division algebra $D$ over $K$ with invariant $1 / 2$.

Let $L / K$ be a quadratic extension. Let $\pi_{L}$ denote a prime element of $L$. By [ $\mathbf{S}$, §XIII.3, Corollaire 3], there exists a $K$-linear embedding $\kappa: L \hookrightarrow D$. It is unique up to conjugation by elements of $D^{\times}$. We choose one such embedding and consider $L$, from now on, as a subfield of $D$. Note that $\mathcal{O}_{L} \subset \mathcal{O}_{D}$. Via this last embedding, we may regard $G$ as a formal $\mathcal{O}_{L}$-module over $k$. In particular, we obtain a map $\mathcal{O}_{L} \rightarrow \operatorname{End}(\operatorname{Lie} G)=k$, which extends the canonical morphism $\mathcal{O}_{K} \rightarrow k$.

Let $A$ be the strict completion of $\mathcal{O}_{L}$ with respect to $k$. In other words, $A$ is the completion of the maximal unramified extension of $\mathcal{O}_{L}$, together with a morphism $A \rightarrow k$ extending the morphism $\mathcal{O}_{L} \rightarrow k$.

Definition 1.2. - A canonical lift of $G$ with respect to the embedding $\kappa: L \hookrightarrow D$ is a lift $F$ of $G$ over $A$ in the category of $\mathcal{O}_{L}$-modules.

In more detail, a canonical lift is a formal $\mathcal{O}_{K}$-module $F$ over $A$, together with an isomorphism of $\mathcal{O}_{K}$-modules $\lambda: F \otimes k \xrightarrow{\sim} G$ and an isomorphism of $\mathcal{O}_{K^{-}}$-algebras $\gamma: \mathcal{O}_{L} \xrightarrow{\sim} \operatorname{End}(F)$, such that the following holds. First, the composition of $\gamma$ with the regular representation $\operatorname{End}(F) \rightarrow \operatorname{End}(\operatorname{Lie} F)=A$ is the canonical inclusion $\mathcal{O}_{L} \subset A$. Second, the composition of $\gamma$ with the inclusion $\operatorname{End}(F) \hookrightarrow \operatorname{End}(G)=\mathcal{O}_{D}$ induced by $\lambda$ is equal to $\kappa$. Note that $\gamma$ is uniquely determined by the lift $F$ and the first
condition. We will omit it from our notation and simply write $[a]_{F}: F \rightarrow F$ for the endomorphism $\gamma(a)$. Also, the fixed embedding $\kappa$ will mostly be understood, and we write $[a]_{G}: G \rightarrow G$ for the endomorphism $\kappa(a)$.

Since $G$ has height one as an $\mathcal{O}_{L}$-module, it follows from [VZ, Theorem 3.8], that a canonical lift $F$ is uniquely determined, up to $*$-isomorphism, by the embedding $\kappa$. On the other hand, using Lubin-Tate theory and the uniqueness statement of Theorem 1.1, we also conclude that a canonical lift $F$ exists, for any choice of $\kappa$. So it is justified to speak about the canonical lift $F$ of $G$, with respect to $\kappa$. By choosing a suitable parameter $X$ for $F$, we may always assume that

$$
\left[\pi_{L}\right]_{F}(X)=\pi_{L} X+X^{q^{2 / e}},
$$

where $e$ is the ramification index of the extension $L / K$.
1.4. Let $F$ be the canonical lift of $G$ over $A$, with respect to a fixed embedding $\kappa: L \hookrightarrow D$. For any positive integer $n$, we set

$$
A_{n}:=A / \pi_{L}^{n+1} A, \quad F_{n}:=F \otimes_{A} A_{n}, \quad H_{n}:=\operatorname{End}_{A_{n}}\left(F_{n}\right)
$$

Since $\mathcal{O}_{L} \subset H_{n}$ for all $n$, we may consider the rings $H_{n}$ as left $\mathcal{O}_{L}$-modules. We have a sequence of $\mathcal{O}_{L}$-linear maps, which are injective by [VZ, Lemma 2.6]:

$$
H_{n} \longleftrightarrow H_{n-1} \longleftrightarrow \cdots \hookrightarrow H_{0}=\mathcal{O}_{D} .
$$

We shall consider $H_{n}$ as an $\mathcal{O}_{L}$-submodules of $\mathcal{O}_{D}$. Since $A$ is complete, we have

$$
\cap_{n \geq 0} H_{n}=\mathcal{O}_{L}
$$

By [VZ, Proposition 3.2], we have an injective map

$$
H_{n-1} / H_{n} \longleftrightarrow H^{2}\left(G, M_{n}\right),
$$

where $M_{n}:=\left(\pi_{L}^{n}\right) /\left(\pi_{L}^{n+1}\right)$.
Lemma 1.3. - Fix $n \geq 1$ and let $\alpha$ be an element of $H_{n-1}-H_{n}$. Then $\left[\pi_{L}\right]_{G} \circ \alpha \in$ $H_{n}-H_{n+1}$. In other words, multiplication with $\pi_{L}$ induces an injective homomorphism of $\mathcal{O}_{L}$-modules

$$
H_{n-1} / H_{n} \longleftrightarrow H_{n} / H_{n+1} .
$$

Proof. - We may represent $\alpha$ by a power series $\alpha(X) \in A \llbracket[X]$, without constant coefficient, whose reduction modulo $\pi_{L}^{n}$ is an endomorphism of $F_{n-1}$. We write $\alpha_{n}$ for the reduction of $\alpha$ modulo $\pi_{L}^{n+1}$. Set

$$
\epsilon:=\alpha \circ[\pi]_{F}-{ }_{F}[\pi]_{F} \circ \alpha .
$$

Since $\alpha_{n-1}$ is an endomorphism of $F_{n-1}$, we have $\epsilon \equiv 0\left(\bmod \pi_{L}^{n}\right)$. Moreover, if $\left(\Delta,\left\{\delta_{a}\right\}\right) \in Z^{2}\left(G, M_{n}\right)$ denotes the cocycle associated to $\alpha_{n}$ by [VZ, Proposition 3.2], then we have

$$
\epsilon \equiv \delta_{\pi} \quad\left(\bmod \pi_{L}^{n+1}\right)
$$

By assumption, the endomorphism $\alpha_{n-1}$ of $F_{n-1}$ cannot be lifted to an endomorphism of $F_{n}$. Therefore, Corollary 3.4 of $[\mathbf{V Z}]$ shows that $\epsilon(X)=c X^{q}+\ldots$, with $c \in$ $\left(\pi_{L}^{n}\right)-\left(\pi_{L}^{n+1}\right)$.

Set

$$
\epsilon^{\prime}:=\left[\pi_{L}\right]_{F} \circ \alpha \circ[\pi]_{F}-{ }_{F}[\pi]_{F} \circ\left[\pi_{L}\right]_{F} \circ \alpha .
$$

Since $\left[\pi_{L}\right]_{F}$ is an endomorphism of $F$, we actually have $\epsilon^{\prime}=\left[\pi_{L}\right]_{F} \circ \epsilon$. Using our assumption $\left[\pi_{L}\right]_{F}(X)=\pi_{L} X+X^{q^{2 / e}}$ and the congruence $\epsilon \equiv 0\left(\bmod \pi_{L}^{n}\right)$, we see that

$$
\epsilon^{\prime}=\pi_{L} c X^{q}+\cdots \equiv 0 \quad\left(\bmod \pi_{L}^{n+1}\right)
$$

By [VZ, Corollary 3.4], this implies that $\left[\pi_{L}\right]_{F} \circ \alpha_{n}$ is an endomorphism of $F_{n}$, i.e., $\left[\pi_{L}\right] \circ \alpha \in H_{n}$. Moreover, if $\left(\Delta^{\prime},\left\{\delta_{a}^{\prime}\right\}\right) \in Z^{2}\left(G, M_{n+1}\right)$ denotes the cocycle associated to $\left[\pi_{L}\right] \circ \alpha_{n+1}$, then we have

$$
\epsilon^{\prime} \equiv \delta_{\pi}^{\prime} \quad\left(\bmod \pi_{L}^{n+2}\right)
$$

Since $\pi_{L} c \in\left(\pi_{L}^{n+1}\right)-\left(\pi_{L}^{n+2}\right)$, Corollary 3.4 of $[\mathbf{V Z}]$ shows that $\left[\pi_{L}\right]_{F} \circ \alpha_{n}$ cannot be lifted to an endomorphism of $F_{n}$. This means that $\left[\pi_{L}\right] \circ \alpha \notin H_{n+1}$.

We can now prove the main result of this section (Proposition 3.3 in $[\mathbf{G}]$ ).
Theorem 1.4. - For $n \geq 1$ we have $H_{n}=\mathcal{O}_{L}+\pi_{L}^{n} \mathcal{O}_{D}$.
Proof. - Each group $H_{n}$ is a submodule of the free rank-two $\mathcal{O}_{L}$-module $\mathcal{O}_{D}$ and contains the direct factor $\mathcal{O}_{L} \subset \mathcal{O}_{D}$. Therefore, the quotients $H_{n-1} / H_{n}$ are cyclic $\mathcal{O}_{L}$-modules. By Lemma 1.3, these quotients are killed by $\pi_{L}$. Hence $H_{n-1} / H_{n}$ is either 0 or isomorphic to $\mathcal{O}_{L} / \pi_{L} \mathcal{O}_{L}$. We claim that only the second case occurs. The case $n=1$ is dealt with in the following lemma.

Lemma 1.5. - We have $H_{1} \neq H_{0}=\mathcal{O}_{D}$.
We will prove this lemma in the next subsection. Lemma 1.3 says that left multiplication with $\pi_{L}$ induces an injective map $H_{n-1} / H_{n} \hookrightarrow H_{n} / H_{n+1}$. So by induction on $n$, Lemma 1.5 and the arguments preceding it show that $H_{n} / H_{n+1} \cong \mathcal{O}_{L} / \pi_{L} \mathcal{O}_{L}$ for all $n$ and that $\mathcal{O}_{D} / H_{n}$ is an $\mathcal{O}_{L}$-module of length $n$, killed by $\pi_{L}^{n}$. The theorem follows immediately.
1.5. We are now going to prove Lemma 1.5. We distinguish two cases.

Case 1: $L / K$ is unramified. In this case, we may assume that $\pi_{L}=\pi$ and hence $[\pi]_{F}=\pi X+X^{q^{2}}$. Then

$$
\mathcal{O}_{D}=\mathcal{O}_{L} \oplus \mathcal{O}_{L} \cdot \Pi
$$

where $\Pi=X^{q}$, see the proof of Theorem 1.1. Let $\alpha=\sum_{i \geq q} a_{i} X^{i} \in A_{1}[[X]]$ be a lift of $\Pi$ with leading term $X^{q}$. Let $\left(\Delta,\left\{\delta_{a}\right\}\right) \in Z^{2}\left(G, M_{1}\right)$ be the cocycle associated
to $\alpha$. Using Taylor expansion, we see that

$$
\begin{aligned}
\delta_{\pi}(X) & =\alpha\left([\pi]_{F_{1}}(X)\right)-F_{1}[\pi]_{F_{1}}(\alpha(X)) \\
& =\left(\alpha\left(X^{q^{2}}\right)+\pi \cdot \alpha^{\prime}\left(X^{q^{2}}\right) X\right)-F_{1}\left(\pi \alpha(X)+\alpha(X)^{q^{2}}\right) \\
& =-\pi X^{q}+\cdots \neq 0
\end{aligned}
$$

(Here we use the notation $\alpha^{\prime}:=\partial \alpha / \partial X$.) Therefore, by [VZ, Corollary 3.4], we have II $\notin H_{1}$.

Case 2: $L / K$ is ramified. Then $\pi_{L}$ satisfies an Eisenstein equation over $\mathcal{O}_{K}$, which we may normalize to

$$
\pi_{L}^{2}+a \pi_{L}+\pi=0
$$

with $a \in \pi \mathcal{O}_{K}$. Assuming, as usual, that $\left[\pi_{L}\right]_{F}=\pi_{L} X+X^{q}$, a short computation yields the congruence

$$
\begin{equation*}
[\pi]_{F}(X) \equiv-\pi_{L} X^{q}-X^{q^{2}}+\ldots \quad(\bmod \pi) \tag{1.1}
\end{equation*}
$$

Let $j \in \mathcal{O}_{D}$ be an element which generates an unramified quadratic extension of $K$. We may assume that $j(X)=\bar{u} X+\ldots$, where $\bar{u} \in k$ generates the quadratic extension of the residue class field of $\mathcal{O}_{K}$. Lift $j$ to a power series $\alpha(X)=u X+\cdots \in A_{1}[[X]]$ modulo $\pi$, and let $\left(\Delta,\left\{\delta_{a}\right\}\right) \in Z^{2}\left(G, M_{1}\right)$ be the associated cocycle. Then $u^{q} \not \equiv u$ $\left(\bmod \pi_{L}\right)$. Using the congruence (1.1), we compute

$$
\begin{aligned}
\delta_{\pi}(X) & =\alpha\left([\pi]_{F_{1}}(X)\right)-F_{1}[\pi]_{F_{1}}(\alpha(X)) \\
& =\left(u\left(-\pi_{L} X^{q}-X^{q^{2}}\right)+\ldots\right)-F_{1}\left(-\pi_{L} \alpha(X)^{q}-\alpha(X)^{q^{2}}\right) \\
& =\pi_{L}\left(u^{q}-u\right) X^{q}+\ldots
\end{aligned}
$$

As in Case 1, we use [VZ, Corollary 3.4], to conclude that $j \notin H_{1}$.

## 2. Isogenies and Tate modules

In this section we review the connection between the endomorphism ring and the isogeny classes of a formal $\mathcal{O}_{K}$-module on the one hand, and lattices inside the Tate module on the other hand. These results will be used in the following section on quasi-canonical lifts.
2.1. As in the previous sections, $K$ denotes a field which is complete with respect to a discrete valuation and has a finite residue field of order $q=p^{f}$. We let $k$ denote an algebraic closure of the residue field of $K$. Furthermore, $A$ is a flat local $\mathcal{O}_{K}$-algebra which is a complete discrete valuation ring with residue field $k$, and $M$ is the fraction field of $A$. We fix an algebraic closure $\bar{M}$ of $M$.

Let $F$ be a formal $\mathcal{O}_{K^{-}}$-module of finite height $h$ over $A$ (not necessarily a canonical lift). We write

$$
\Lambda(F):=F(\bar{M})_{\mathrm{tor}}=\cup_{n} F\left[\pi^{n}\right]
$$

for the torsion subgroup of $F$ and

$$
T(F):={\underset{\check{n}}{n}}_{\lim _{n}} F\left[\pi^{n}\right]
$$

for the Tate module of $F$. These are $\mathcal{O}_{K}$-modules with a continuous, $\mathcal{O}_{K}$-linear action of $\operatorname{Gal}(\bar{M} / M)$. As $\mathcal{O}_{K}$-modules, we have non-canonical isomorphisms

$$
\Lambda(F) \cong\left(K / \mathcal{O}_{K}\right)^{h}, \quad T(F) \cong \mathcal{O}_{K}^{h}
$$

Set $V(F):=T(F) \otimes \mathcal{O}_{K} K$; then we have a canonical short exact sequence of $\operatorname{Gal}(\bar{M} / M)-\mathcal{O}_{K}$-modules

$$
\begin{equation*}
0 \rightarrow T(F) \longrightarrow V(F) \longrightarrow \Lambda(F) \rightarrow 0 \tag{2.1}
\end{equation*}
$$

Let $A^{\prime}$ be a finite extension of $A$, and let $F^{\prime}$ be a formal $\mathcal{O}_{K}$-module over $A^{\prime}$. An isogeny between $F$ and $F^{\prime}$ defined over $A^{\prime}$ is a nonzero homomorphism $\alpha: F \otimes_{A} A^{\prime} \rightarrow$ $F^{\prime}$ of formal $\mathcal{O}_{K^{-}}$-modules. If such an isogeny exists, then we say that $F^{\prime}$ is isogenous to $F$ (over $A^{\prime}$ ). For simplicity, we shall write $\alpha: F \rightarrow F^{\prime}$, and consider $\alpha$ as a power series in $\mathcal{O}_{\bar{M}}[[X]$ whose coefficients generate a finite extension of $A$. We say that $\alpha$ is defined over $A^{\prime}$ if $\alpha \in A^{\prime}[[X]$.

Given an isogeny $\alpha: F \rightarrow F^{\prime}$, we obtain a diagram

with exact rows and columns. Note that $N$ is equal to the kernel of $\alpha$; it is a finite $\mathcal{O}_{K}$-submodule. A trivial version of the snake lemma shows that we have a canonical isomorphism $N \cong \operatorname{Coker} T(\alpha)$.

The following theorem states that every finite $\mathcal{O}_{K}$-submodule of $\Lambda(F)$ arises as the kernel of an isogeny. More precisely:

Theorem 2.1. - Let $N \subset \Lambda(F)$ be a finite $\mathcal{O}_{K}$-submodule, $\Gamma^{\prime} \subset \Gamma$ the stabilizer of $N$, $M^{\prime} \subset \bar{M}$ the fixed field of $\Gamma^{\prime}$ and $A^{\prime}$ the valuation ring of $M^{\prime}$. Then the formula

$$
\alpha(X):=\prod_{z \in N}\left(X-{ }_{F} z\right) \in A^{\prime}[[X]]
$$

defines an isogeny $\alpha: F \rightarrow F^{\prime}$ over $A^{\prime}$. It has the following properties.

1. $\operatorname{Ker}(\alpha)=N$.
2. Let $\beta: F \rightarrow F^{\prime \prime}$ be an isogeny with $N \subset \operatorname{Ker}(\beta)$. Then there exists a unique isogeny $\gamma: F^{\prime} \rightarrow F^{\prime \prime}$ with $\beta=\gamma \circ \alpha$.

Proof. - See [H, §35.2].
2.2. It will be more convenient for us to reformulate Theorem 2.1 in terms of lattices $T^{\prime} \subset V(F)$ (instead of finite subgroups $N \subset \Lambda(F)$ ). Let $F$ be a formal $\mathcal{O}_{K}$-module of finite height over $A$. Set $T:=T(F)$ and $V:=V(F)$.

Corollary 2.2. - 1. Let $T^{\prime} \subset V$ be an $\mathcal{O}_{K}$-lattice containing the lattice $T$ (a superlattice). Then there exists an isogeny $\alpha: F \rightarrow F^{\prime}$ such that $T^{\prime}=V(\alpha)^{-1}\left(T\left(F^{\prime}\right)\right)$. If $T^{\prime \prime}$ is a superlattice of $T^{\prime}$ and $\beta: F \rightarrow F^{\prime \prime}$ an isogeny with $T^{\prime \prime}=V(\beta)^{-1}\left(T\left(F^{\prime \prime}\right)\right)$, then there exists a unique isogeny $\gamma: F^{\prime} \rightarrow F^{\prime \prime}$ such that $\beta=\gamma \circ \alpha$.
2. Let $T^{\prime} \subset T$ be an $\mathcal{O}_{K}$-sublattice. Then there exists an isogeny $\alpha: F^{\prime} \rightarrow F$ such that $T^{\prime}=\operatorname{Im}(T(\alpha))$. If $T^{\prime \prime} \subset T^{\prime}$ is another sublattice, and $\beta: F^{\prime \prime} \rightarrow F$ is an isogeny such that $T^{\prime \prime}=\operatorname{Im}(T(\beta))$, then there exists a unique isogeny $\gamma: F^{\prime \prime} \rightarrow F^{\prime}$ with $\beta=\alpha \circ \gamma$.

Proof. - Given $T^{\prime}$ as in Part 1, we set $N:=T^{\prime} / T$. Via the short exact sequence (2.1), we consider $N$ as a (finite) $\mathcal{O}_{K}$-submodule of $V$. Let $\alpha: F \rightarrow F^{\prime}$ be the isogeny with kernel $N$, which exists by Theorem 2.1.1. Then the diagram (2.2) shows that $T^{\prime}=V(\alpha)^{-1}\left(T\left(F^{\prime}\right)\right)$. This proves the first assertion in Part 1. The second assertion follows from Theorem 2.1.2.

We are now going to prove Part 2 of the corollary. Let $T^{\prime} \subset T$ be a sublattice. Choose an integer $n$ such that $\pi^{n} T \subset T^{\prime}$. By Part 1 of the corollary, there exists an isogeny $\beta: F \rightarrow F^{\prime}$ such that $V(\beta)^{-1}\left(T\left(F^{\prime}\right)\right)=\pi^{-n} T^{\prime}$. The kernel of $\beta$ is isomorphic to $\pi^{-n} T^{\prime} / T$, which is an $\mathcal{O}_{K^{-}}$module killed by $\pi^{n}$. Therefore, Theorem 2.1.2 shows that there exists an isogeny $\alpha: F^{\prime} \rightarrow F$ with $\alpha \circ \beta=\left[\pi^{n}\right]_{F}$. By construction, we have

$$
\operatorname{Im}(T(\alpha))=\pi^{n} \cdot V(\beta)^{-1}\left(T\left(F^{\prime}\right)\right)=T^{\prime}
$$

This proves the first assertion of Part 2. The proof of the second assertion is left to the reader.
2.3. Let $F, T$ and $V$ be as before. The faithful representation of $\operatorname{End}(F)$ on $V$ extends to a faithful representation

$$
\operatorname{End}^{0}(F):=\operatorname{End}(F) \otimes \mathcal{O}_{K} K \longleftrightarrow \operatorname{End}_{K}(V) .
$$

We will from now on consider elements of $\operatorname{End}^{0}(F)$ as elements of $\operatorname{End}_{K}(V)$.
Let $T^{\prime}, T^{\prime \prime}$ be $\mathcal{O}_{K^{-}}$-superlattices of $T$ inside $V$. Let $\alpha: F \rightarrow F^{\prime}$ and $\beta: F \rightarrow F^{\prime \prime}$ be the corresponding isogenies, as in Corollary 2.2.1. We identify $V\left(F^{\prime}\right)$ and $V\left(F^{\prime \prime}\right)$ with $V$, via the isomorphisms $V(\alpha)$ and $V(\beta)$. Then $T^{\prime}=T\left(F^{\prime}\right)$ and $T^{\prime \prime}=T\left(F^{\prime \prime}\right)$.

Corollary 2.3. - The map which sends a homomorphism $\psi: F^{\prime} \rightarrow F^{\prime \prime}$ to the induced endomorphism $\tilde{\psi}: V \cong V\left(F^{\prime}\right) \rightarrow V\left(F^{\prime \prime}\right) \cong V$ is a bijection

$$
\operatorname{Hom}\left(F^{\prime}, F^{\prime \prime}\right) \xrightarrow{\sim}\left\{\tilde{\psi} \in \operatorname{End}^{0}(F) \mid \tilde{\psi}\left(T^{\prime}\right) \subset T^{\prime \prime}\right\}
$$

Proof. - Let $\psi: F^{\prime} \rightarrow F^{\prime \prime}$ be a homomorphism and $\tilde{\psi} \in \operatorname{End}_{K}(V)$ the induced endomorphism of $V$. By definition, we have $\tilde{\psi}\left(T^{\prime}\right) \subset T^{\prime \prime}$. We have to show that $\tilde{\psi} \in$ $\operatorname{End}^{0}(F)$. Set $\gamma:=\psi \circ \alpha: F \rightarrow F^{\prime \prime}$. The isogeny $\gamma$ corresponds, via Corollary 2.2.2, to the sublattice $\tilde{\psi}(T) \subset T^{\prime \prime}$. From the same point of view, the isogeny $\beta: F \rightarrow F^{\prime \prime}$ corresponds to the sublattice $T \subset T^{\prime \prime}$. Choose an integer $n$ such that $\pi^{n} \tilde{\psi}(T) \subset T$. Then by Corollary 2.2 .2 , there exists an endomorphism $\phi: F \rightarrow F$ such that $\beta \circ \phi=$ $\underset{\sim}{\gamma} \circ\left[\pi^{n}\right]_{F}$. One checks that $\phi=\pi^{n} \tilde{\psi}$, as elements of $\operatorname{End}_{K}(V)$, which shows that $\tilde{\psi} \in \operatorname{End}^{0}(F)$.

Conversely, let $\tilde{\psi}$ be an element of $\operatorname{End}^{0}(F) \subset \operatorname{End}_{K}(V)$ with $\tilde{\psi}\left(T^{\prime}\right) \subset T^{\prime \prime}$. By definition, we can write $\tilde{\psi}=\pi^{-n} \phi$ for some endomorphism $\phi: F \rightarrow F$. The isogeny $\alpha \circ\left[\pi^{n}\right]_{F}: F \rightarrow F^{\prime}$ (resp. the isogeny $\beta \circ \phi: F \rightarrow F^{\prime \prime}$ ) corresponds, via Corollary 2.2.1, to the superlattice $\pi^{-n} T^{\prime} \supset T$ (resp. the superlattice $\phi^{-1}\left(T^{\prime \prime}\right) \supset T$ ). The assumption $\tilde{\psi}\left(T^{\prime}\right) \subset T^{\prime \prime}$ together with $\tilde{\psi}=\pi^{-n} \phi$ implies $\pi^{-n} T^{\prime} \subset \phi^{-1}\left(T^{\prime \prime}\right)$. Therefore, by Corollary 2.2.1, there exists an isogeny $\psi: F^{\prime} \rightarrow F^{\prime \prime}$ with $\psi \circ \alpha \circ\left[\pi^{n}\right]_{F}=\beta \circ \phi$. By construction, $\tilde{\psi}$ is the image of $\psi$ under the embedding $\operatorname{Hom}\left(F^{\prime}, F^{\prime \prime}\right) \hookrightarrow \operatorname{End}_{K}(V)$. This concludes the proof of the corollary.

## 3. Quasi-canonical lifts

A quasi-canonical lift is a lift whose endomorphism ring is an order in a quadratic extension $L / K$. In this section we show that every quasi-canonical lift is isogenous to a canonical lift, and we determine the set of isomorphism classes of all quasi-canonical lifts together with its natural Galois action.
3.1. We now come back to the situation of Section 1. In particular, $G$ is the (unique) formal $\mathcal{O}_{K^{-}}$-module of height two over $k$. We fix a quadratic extension $L / K$, an $\mathcal{O}_{K^{-}}$ linear embedding $\kappa: \mathcal{O}_{L} \hookrightarrow \mathcal{O}_{D}:=\operatorname{End}_{k}(G)$. We denote by $F$ the canonical lift of $G$ with respect to $\kappa$. Recall that $F$ is defined over $A$, the strict completion of $\mathcal{O}_{L}$ with respect to the map $\mathcal{O}_{K} \rightarrow k$ induced by the $\mathcal{O}_{L^{-}}$-action on $\operatorname{Lie}(G)$.

Let $M$ denote the fraction field of $A, \bar{M}$ an algebraic closure of $M$ and $\Gamma:=$ $\operatorname{Gal}(\bar{M} / M)$. We let $T:=T(F)$ denote the Tate-module of $F$ and $V:=T \otimes_{\mathcal{O}_{K}} K$. Note that $T$ has the structure of a free $\mathcal{O}_{L}$-module of rank one, and that the $\Gamma$-action on $T$ is continuous and $\mathcal{O}_{L}$-linear. By Lubin-Tate theory, the resulting homomorphism

$$
\begin{equation*}
\rho: \Gamma=\operatorname{Gal}(\bar{M} / M) \longrightarrow \mathcal{O}_{L}^{\times} \tag{3.1}
\end{equation*}
$$

yields an isomorphism $\Gamma^{\text {ab }} \xrightarrow{\sim} \mathcal{O}_{L}^{\times}$. Identifying $\Gamma$ with the inertia subgroup of $\operatorname{Gal}(\bar{L} / L)$, the homomorphism (3.1) is the inverse of the reciprocity map $L^{\times} \rightarrow \operatorname{Gal}(\bar{L} / L)^{\text {ab }}$ of local class field theory, restricted to $\mathcal{O}_{L}^{\times}$. See [LT1].

Fix an integer $s \geq 0$. Let

$$
\mathcal{O}_{s}:=\mathcal{O}_{K}+\mathcal{O}_{L} \cdot \pi^{s}
$$

denote the order of $\mathcal{O}_{L}$ generated by $\mathcal{O}_{K}$ and the ideal $\mathcal{O}_{L} \cdot \pi^{s}$. It is easy to see that every order of $\mathcal{O}_{L}$ containing $\mathcal{O}_{K}$ is equal to $\mathcal{O}_{s}$, for some $s$. Let $M_{s} / M$ be the ring class field of $\mathcal{O}_{s}^{\times}$, i.e., the fixed field of the subgroup $\Gamma_{s} \subset \Gamma$, where $\Gamma_{s}$ is the inverse image of $\mathcal{O}_{s}^{\times} \subset \mathcal{O}_{L}^{\times}$under the inverse reciprocity homomorphism (3.1). In other words, we have

$$
\operatorname{Gal}\left(M_{s} / M\right) \cong \mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times}
$$

An easy computation shows that, for $s \geq 1$,

$$
\left[M_{s}: M\right]=\left|\mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times}\right|=\left\{\begin{array}{cl}
q^{s-1}(q+1), & \text { if } L / K \text { is unramified } \\
q^{s}, & \text { if } L / K \text { is ramified }
\end{array}\right.
$$

Definition 3.1. - A quasi-canonical lift of $G$ of level $s$ (with respect to the embedding $\left.\kappa: \mathcal{O}_{L} \hookrightarrow \mathcal{O}_{D}\right)$ is a lift $F$ of $G$, defined over some finite extension $A^{\prime} / A$, together with an $\mathcal{O}_{K}$-algebra isomorphism $\gamma: \mathcal{O}_{s} \xrightarrow{\sim} \operatorname{End}\left(F^{\prime}\right)$, such that the following holds.

1. The composition of $\gamma$ with the representation $\operatorname{End}\left(F^{\prime}\right) \hookrightarrow \operatorname{End}\left(\operatorname{Lie} F^{\prime}\right)=A^{\prime}$ is the canonical embedding $\mathcal{O}_{s} \hookrightarrow A^{\prime}$.
2. The composition of $\gamma$ with the embedding $\operatorname{End}\left(F^{\prime}\right) \hookrightarrow \mathcal{O}_{D}$ is equal to the restriction of $\kappa$ to $\mathcal{O}_{s} \subset \mathcal{O}_{L}$.

To ease the notation, we will usually omit the isomorphism $\gamma$ and the embedding $\kappa$ from our notation. Note that a quasi-canonical lift of level 0 is the same thing as a canonical lift (which exists and is unique). For general $s$, we have the following result.

Theorem 3.2. - Let $\mathcal{O}_{M_{*}}$ denote the ring of integers of $M_{s}$.

1. Let $F^{\prime}$ be a quasi-canonical lift of level s. Then there exists an isogeny

$$
\alpha: F \longrightarrow F^{\prime}
$$

of degree $q^{s}$, defined over $\mathcal{O}_{M_{s}}$. It is unique up to composing $\alpha$ with an element of $\operatorname{Aut}(F)=\mathcal{O}_{L}^{\times}$. In particular, $F^{\prime}$ can be defined over $\mathcal{O}_{M_{s}}$.
2. The set of *-isomorphism classes of all quasi-canonical lifts of level s is a principal homogeneous space under the action of $\operatorname{Gal}\left(M_{s} / M\right)$.

Remark 3.3. - The proof of this theorem will show that the action of $\operatorname{Gal}\left(M_{s} / M\right)$ on the set of $*$-isomorphism classes can be described as follows. Let $\left(F^{\prime}, \lambda\right)$ be a quasicanonical lift of level $s$ (with $\lambda: F^{\prime} \otimes k \xrightarrow{\sim} G$ ), and $\sigma \in \Gamma$. Then the lift $\left(F^{\prime}, \lambda\right)^{\sigma}$ is $*$-isomorphic to the lift $\left(F^{\prime},\left[\rho(\sigma)^{-1}\right]_{G} \circ \lambda\right)$. Therefore, by Theorem 3.2.2, two quasicanonical lifts of the same level are always isomorphic as formal $\mathcal{O}_{K}$-modules.
3.2. Let $\alpha: F \rightarrow F^{\prime}$ and $\beta: F \rightarrow F^{\prime \prime}$ be two isogenies with source $F$. We say that $\alpha$ and $\beta$ are isomorphic if there exists an isomorphism of formal $\mathcal{O}_{K}$-modules $\gamma: F^{\prime} \xrightarrow{\sim} F^{\prime \prime}$ with $\beta=\gamma \circ \alpha$.

Fix an isogeny $\alpha: F \rightarrow F^{\prime}$. To simplify the notation, we will identify $V\left(F^{\prime}\right)$ with $V$ via the isomorphism $V(\alpha)$. Then, by Corollary 2.2.1, $\alpha$ corresponds, up to isomorphism, to an $\mathcal{O}_{K}$-superlattice $T^{\prime} \supset T$ in $V$. Moreover, by Corollary 2.3, $\alpha$ induces an isomorphism of $\mathcal{O}_{K}$-algebras

$$
\begin{equation*}
\operatorname{End}\left(F^{\prime}\right) \xrightarrow{\sim}\left\{\phi \in L=\operatorname{End}^{0}(F) \mid \phi\left(T^{\prime}\right) \subset T^{\prime}\right\} \tag{3.2}
\end{equation*}
$$

This exhibits $\operatorname{End}\left(F^{\prime}\right)$ as an order of $\mathcal{O}_{L}$.
Lemma 3.4. - Let $T$ be a free $\mathcal{O}_{L}$-module of rank one, $V:=T \otimes \mathcal{O}_{K} K$. Let $T^{\prime} \supset T$ be an $\mathcal{O}_{K}$-superlattice in $V$. Then there exists a generator $t$ of $T$ (i.e., $T=\mathcal{O}_{L} \cdot t$ ) and integers $n, s \geq 0$ such that

$$
\pi_{L}^{n} \cdot T^{\prime}=\left(\mathcal{O}_{K} \cdot \pi^{-s}+\mathcal{O}_{L}\right) \cdot t
$$

Moreover, the multiplicator $\mathcal{O}_{T^{\prime}}$ of $T^{\prime}$ is equal to the order $\mathcal{O}_{s} \subset \mathcal{O}_{L}$.
Proof. -- For $T^{\prime} \supset T$ as in Part 1, define

$$
n:=\max \left\{n^{\prime} \mid \pi_{L}^{n^{\prime}} T^{\prime} \supset T\right\}, \quad s:=\min \left\{s^{\prime} \mid \pi^{s^{\prime}} \pi_{L}^{n} T^{\prime} \subset T\right\} .
$$

Then $\pi_{L}^{n} T^{\prime} / T$ is a cyclic $\mathcal{O}_{K}$-module, generated by an element of the form $\pi^{-s} t$. Moreover, any $t$ with this property is a generator of $T$. It follows that $\pi_{L}^{n} T^{\prime}=$ $\left(\mathcal{O}_{K} \cdot \pi^{-s}+\mathcal{O}_{L}\right) \cdot t$. The proof of the fact that $\mathcal{O}_{s}$ is the multiplicator of $T^{\prime}$ is standard and left to the reader.

A superlattice $T^{\prime} \supset T$ is called minimal of level $s$ if $T^{\prime}=\left(\mathcal{O}_{K} \cdot \pi^{-s}+\mathcal{O}_{L}\right) \cdot t$, for some generator $t$ of $T$. The corresponding isogenies $\alpha: F \rightarrow F^{\prime}$ are also called minimal of level $s$. We let $X_{s}$ denote the set of isomorphism classes of minimal isogenies of level $s$. The Galois group $\Gamma$ acts on $X_{s}$, in a natural way. There is also an action of $\mathcal{O}_{L}^{\times}$ on $X_{s}$, given by composing $\alpha: F \rightarrow F^{\prime}$ with the automorphism $[a]_{F}: F \xrightarrow{\sim} F$, for $a \in \mathcal{O}_{L}^{\times}$.

Proposition 3.5. - The actions of $\Gamma$ and $\mathcal{O}_{L}^{\times}$on $X_{s}$ are anti-compatible via the reciprocity homomorphism $\rho: \Gamma \rightarrow \mathcal{O}_{L}^{\times}$, i.e., for $\sigma \in \Gamma$ there exists an isomorphism $\gamma_{\sigma}:\left(F^{\prime}\right)^{\sigma} \xrightarrow{\sim} F^{\prime}$ such that the diagram

commutes. Furthermore, $X_{s}$ is a principal homogeneous space under the induced action of $\operatorname{Gal}\left(M_{s} / M\right) \cong \mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times}$.

Proof. - If the isogeny $\alpha: F \rightarrow F^{\prime}$ corresponds to the lattice $T^{\prime}$, then $\alpha \circ[a]_{F}$ : $F \rightarrow F^{\prime}$, for $a \in \mathcal{O}_{L}$, corresponds to the lattice $a^{-1} \cdot T^{\prime}$. Therefore, it follows immediately from Lemma 3.4 that the action of $\mathcal{O}_{L}^{\times}$on $X_{s}$ is transitive, and the stabilizer of each element is equal to $\mathcal{O}_{s}^{\times}$. To see that this action is compatible with the Galois action, fix an element $\sigma \in \Gamma$. Clearly, the kernel of $\alpha^{\sigma}$ can be identified with $\left(T^{\prime} / T\right)^{\sigma}=\rho(\sigma) \cdot T^{\prime} / T$. Since this is also the kernel of $\alpha \circ\left[\rho(\sigma)^{-1}\right]_{F}$, the existence of $\gamma_{\sigma}$ follows from Theorem 2.1. The proposition is proved.

Proof of Theorem 3.2. -. We first prove Part 1 of the theorem. Let $F^{\prime}$ be a quasicanonical lift of level $s$. Set $T^{\prime}:=T\left(F^{\prime}\right)$ and $V^{\prime}:=T^{\prime} \otimes \mathcal{O}_{K} K$. The isomorphism $\mathcal{O}_{s} \xrightarrow{\sim}$ $\operatorname{End}\left(F^{\prime}\right)$ extends to an isomorphism $L \xrightarrow{\sim} \operatorname{End}^{0}\left(F^{\prime}\right)$, which gives $V^{\prime}$ the structure of an $L$-vector space of dimension one and identifies $\mathcal{O}_{s}$ with the multiplicator of the lattice $T^{\prime} \subset V^{\prime}$.

Let $T^{\prime \prime} \subset T^{\prime}$ be a maximal $\mathcal{O}_{L}$-submodule of rank one. Then $T^{\prime}=\left(\mathcal{O}_{K} \cdot \pi^{-s}+\mathcal{O}_{L}\right) \cdot t$ for some generator $t$ of $T^{\prime \prime}$, by Lemma 3.4. Let $\alpha: F^{\prime \prime} \rightarrow F^{\prime}$ be an isogeny with $\operatorname{Im}(T(\alpha))=T^{\prime \prime}$, see Corollary 2.2.2. By Corollary $2.3, \alpha$ induces an isomorphism

$$
\operatorname{End}\left(F^{\prime \prime}\right) \cong\left\{\phi \in \operatorname{End}^{0}\left(F^{\prime}\right)=L \mid \phi\left(T^{\prime \prime}\right)=T^{\prime \prime}\right\} \cong \mathcal{O}_{L}
$$

Therefore, $F^{\prime \prime} \cong F$ as formal $\mathcal{O}_{K}$-modules. Choosing an arbitrary isomorphism $F^{\prime \prime} \cong$ $F$, we can regard $\alpha: F \cong F^{\prime \prime} \rightarrow F^{\prime}$ as an element of $X_{s}$. Since $\mathcal{O}_{L}^{\times}$acts transitively on $X_{s}$, by Proposition 3.5, we have proved Part 1 of Theorem 3.2.

Now we prove Part 2 of the theorem. In view of Part 1 and Proposition 3.5, we only need to show the following. For every minimal isogeny $\alpha: F \rightarrow F^{\prime}$ of level $s$, there exists an isomorphism $\lambda: F^{\prime} \otimes k \xrightarrow{\sim} G$ which makes $F^{\prime}$ a quasi-canonical lift. For this, we may assume that the isogeny $\alpha$ is given, as a power series with coefficients in $\mathcal{O}_{M_{*}}$, by the formula of Theorem 2.1:

$$
\alpha(X):=\prod_{\gamma \in \operatorname{Ker}(\alpha)}\left(X-{ }_{F} \gamma\right) .
$$

Here $\operatorname{Ker}(\alpha)$ is simply considered as a subset of the maximal ideal of the ring of integers of $\bar{M}$. Therefore, the reduction of $\alpha$ to $k$ is $\bar{\alpha}(X)=X^{q^{*}}$. By the proof of Theorem 1.1, we may assume that $\Pi(X):=X^{q}$ is an endomorphism of $G$ and lies in the normalizer of $\mathcal{O}_{L}=\operatorname{End}(F) \subset \mathcal{O}_{D}$. In particular, $\bar{\alpha}=\Pi^{s}$ is an endomorphism of $G$. Therefore, $F^{\prime} \otimes k$ is actually equal to $G$. We define the isomorphism $\lambda: F^{\prime} \otimes k \xrightarrow{\sim} G$ as the identity and claim that $\left(F^{\prime}, \lambda\right)$ is a quasi-canonical lift.

By construction, we have an isomorphism

$$
\begin{equation*}
\operatorname{End}\left(F^{\prime}\right) \cong\left\{\phi \in L=\operatorname{End}^{0}(F) \mid \phi\left(T^{\prime}\right)=T^{\prime}\right\} \cong \mathcal{O}_{s} \tag{3.3}
\end{equation*}
$$

Hence the image of the natural injection $\operatorname{End}\left(F^{\prime}\right) \hookrightarrow \operatorname{End}\left(\right.$ Lie $\left.F^{\prime}\right)=\mathcal{O}_{M_{s}}$ is an $\mathcal{O}_{K^{-}}$ algebra isomorphic to $\mathcal{O}_{s}$. It must therefore be equal to $\mathcal{O}_{s}$. Let $\gamma: \mathcal{O}_{s} \xrightarrow{\sim} \operatorname{End}\left(F^{\prime}\right)$ be the resulting isomorphism. Then Condition 1 of Definition 3.1 holds by construction.

Let $\kappa^{\prime}: \mathcal{O}_{s} \hookrightarrow \mathcal{O}_{D}$ be the composition of $\gamma$ with the embedding $\operatorname{End}\left(F^{\prime}\right) \hookrightarrow \mathcal{O}_{D}$ induced by the identification $F^{\prime} \otimes k=G$. We have to show that $\kappa^{\prime}$ is equal to the restriction of $\kappa$ to $\mathcal{O}_{s}$ (see Condition 2 of Definition 3.1). Tracing back the definitions, we see that $\kappa^{\prime}=\left(\left.\kappa\right|_{\mathcal{O}_{s}}\right)^{\bar{\alpha}}$ is the conjugate of $\left.\kappa\right|_{\mathcal{O}_{s}}$ by $\bar{\alpha}=\Pi^{s} \in \mathcal{O}_{D}$. Since we assumed $\Pi$ to lie in the normalizer of the image of $\kappa$, we have already proved that $\kappa^{\prime}$ and $\left.\kappa\right|_{\mathcal{O}_{s}}$ have the same image and are equal up to composition with an element of $\operatorname{Gal}(L / K) \cong \mathbb{Z} / 2$. However, if $L / K$ is ramified, then the assumption that $\Pi$ normalizes $\mathcal{O}_{L}$ already implies that $\Pi \in \mathcal{O}_{L}$, and we get $\kappa^{\prime}=\left.\kappa\right|_{\mathcal{O}_{s}}$ as desired. Now assume that $L / K$ is unramified. Then it suffices to show that $\kappa^{\prime}$ and $\left.\kappa\right|_{\mathcal{O}_{s}}$ agree modulo the maximal ideal $\mathcal{O}_{D} \cdot \Pi$. But this is a consequence of Condition 1 of Definition 3.1. This concludes the proof of Theorem 3.2.

## 4. Canonical subgroups

The main result of this section is Proposition 4.6 which computes the valuation of the formal modulus of a quasi-canonical lift. The heart of the proof of this proposition is the study of canonical subgroups and their behavior under isogenies. The relevance of canonical subgroups was first pointed out in $[\mathbf{L}]$.
4.1. We continue with the notation used in the last section. In particular, $A$ is the completion of the maximal unramified extension of $\mathcal{O}_{L}$ and $M$ the fraction field of $A$. We choose an algebraic closure $\bar{M}$ of $M$ and let $v: \bar{M} \rightarrow \mathbb{Q} \cup\{\infty\}$ denote the exponential rank-one valuation with $v(\pi)=1$.

Let $M^{\prime} / M$ be some finite extension, and let $A^{\prime}$ denote the valuation ring of $M^{\prime}$. Throughout this section, we will implicitly assume that the extension $M^{\prime} / M$ is 'sufficiently large'. In practice this will mean that sometimes we have to enlarge $M^{\prime}$ in order to make certain torsion points $M^{\prime}$-rational.

For the moment, we fix an arbitrary lift $F$ of the formal $\mathcal{O}_{K}$-module $G$, defined over $A^{\prime}$ (not necessarily the canonical lift). By [VZ, Theorem 3.8], $F$ is isomorphic to the pullback of the universal deformation $\tilde{F}$ of $G$ via a unique $\widehat{\mathcal{O}}_{K}^{\mathrm{nr}}$-algebra morphism $R^{\text {univ }} \rightarrow A$. Moreover, $R^{\text {univ }}$ can be written as a power series algebra $\widehat{\mathcal{O}}_{K}^{\text {nr }} \llbracket u \rrbracket$. (The proof of this result in [ $\mathbf{V Z}]$ does not provide us with a natural choice of the parameter $u$, but this is irrelevant for us by Remark 4.2 below. See $[\mathbf{H G}]$ for a more explicit choice of the parameter $u$.)

Definition 4.1. - The image of the parameter $u$ under the morphism $R^{\text {univ }} \rightarrow A^{\prime}$ corresponding to $F$ is denoted by $u(F)$ and is called the formal modulus of the lift $F$. The rational number $v(F):=\min \{v(u(F)), 1\}$ is called the valuation of $F$.

Remark 4.2. - It is clear that the valuation $v(F)$ is actually independent of the choice of the parameter $u$. Therefore, $v(F)$ depends only on the isomorphism class of $F$ as a formal $\mathcal{O}_{K}$-module, and not on the chosen isomorphism $\lambda: F \otimes k \xrightarrow{\sim} G$. Indeed,
a unit $\gamma \in \mathcal{O}_{D}^{\times}$induces an automorphism $\tilde{\gamma}$ of the universal deformation space of $G$ (which sends the pair $(F, \lambda)$ to the pair $(F, \gamma \circ \lambda)$ ). Applying the automorphism $\tilde{\gamma}$ amounts to replacing the parameter $u$ by $u^{\prime}:=\tilde{\gamma}^{*} u$.

Definition 4.3. - A sub- $\mathcal{O}_{K}$-module $H \subset F[\pi]$ of length one is called a canonical subgroup if

$$
v(x)>v(y)
$$

for all $x \in H$ and $y \in F[\pi]-H$.
Note that a canonical subgroup, if it exists, is unique. We may therefore speak about the canonical subgroup of $F$. The two last definitions are related to each other in the following manner.

## Proposition 4.4

1. Write $[\pi]_{F}=\sum_{i \geq 1} a_{i} X^{i}$, with $a_{i} \in A^{\prime}$. Then $v(F)=\min \left\{v\left(a_{q}\right), 1\right\}$.
2. The lift $F$ has a canonical subgroup if and only if

$$
v(F)<\frac{q}{q+1} .
$$

Proof. - It follows from the proof of [VZ, Theorem 3.8], that we can choose for the parameter $u$ defining the isomorphism $\left.R^{\text {univ }} \cong \widehat{\mathcal{O}}_{K}^{\mathrm{nr}}[u]\right]$ the $q$ th coefficient of $[\pi]_{\tilde{F}}$, where $\tilde{F}$ is the universal deformation of $G$. Therefore, Part 1 of the proposition is a direct consequence of the definition of $v(F)$. Now Part 2 is easily seen by looking at the Newton polygon of $[\pi]_{F}$. Indeed, the slope filtration on the set $F[\pi]-\{0\}$ is also a filtration of $\mathcal{O}_{K}$-modules. But as an $\mathcal{O}_{K}$-module, $F[\pi]$ has length two, so there can be at most two finite negative slopes. Also, breaks occur only at $i=1, q^{2}$ and possibly at $i=q$. Since $v\left(a_{1}\right)=1$ and $v\left(a_{q^{2}}\right)=0$, we have a break at $i=q$ if and only if $v(F)<q /(q+1)$.
4.2. Fix a lift $F$ of $G$ defined over $A^{\prime}$ and a sub- $\mathcal{O}_{K}$-module $H \subset F[\pi]$ of length one. Let $\alpha: F \rightarrow F^{\prime}$ be the isogeny with kernel $H$, defined by Theorem 2.1. Recall that $\alpha$ is given by the power series

$$
\begin{equation*}
\alpha(X):=\prod_{x \in H}\left(X-{ }_{F} x\right) . \tag{4.1}
\end{equation*}
$$

Let us choose an isomorphism $\lambda^{\prime}: F^{\prime} \otimes k \xrightarrow{\sim} G$. We will use $\lambda^{\prime}$ as an identification, i.e., we will regard $F^{\prime}$ as a lift of $G$. As in Section 3.2, one can choose $\lambda^{\prime}$ in such a way that $\alpha \otimes k$ gets identified with the isogeny $\Pi=X^{q}: G \rightarrow G$. However, this choice is not at all canonical. In what follows, we are mainly interested in relating the two valuations $v(F)$ and $v\left(F^{\prime}\right)$. By Remark 4.2, the choice that we have made is irrelevant for this problem.

Let $\beta: F \rightarrow F^{\prime}$ be the unique isogeny such that $[\pi]_{F}=\beta \circ \alpha$. Then $H^{\prime}:=\operatorname{ker}(\beta)$ is equal to the image of $F[\pi]$ under the isogeny $\alpha$. Clearly, $H^{\prime}$ is an $\mathcal{O}_{K}$-module of length one.

## Proposition 4.5

1. Suppose that $H$ is the canonical subgroup of $F$. There are two cases:
(a) If $v(F) \leq \frac{1}{q}$ then $v\left(F^{\prime}\right)=q \cdot v(F)$ and $H^{\prime}$ is not canonical.
(b) If $\frac{1}{q}<v(F)<\frac{q}{q+1}$ then $v\left(F^{\prime}\right)=1-v(F)$ and $H^{\prime}$ is the canonical subgroup of $F^{\prime}$.
2. Suppose that $H$ is not the canonical subgroup of $F$. Again we have two cases:
(a) If $v(F) \leq \frac{q}{q+1}$ then $v\left(F^{\prime}\right)=q^{-1} \cdot v(F)$.
(b) If $v(F) \geq \frac{q}{q+1}$ then $v\left(F^{\prime}\right)=1 /(q+1)$.

In both cases, $H^{\prime}$ is the canonical subgroup of $F^{\prime}$.
Proof. - Suppose that $H$ is canonical. By Proposition 4.4, we have $v(F)<q /(q+1)$. Moreover, the proof of this proposition shows that the Newton polygon of $[\pi]_{F}$ has exactly two finite negative slopes, namely

$$
s_{1}=-\frac{1-v(F)}{q-1}, \quad s_{2}=-\frac{v(F)}{q^{2}-q} .
$$

Here $s_{1}$ is the slope above the interval $[1, q]$ and corresponds to the canonical subgroup, whereas $s_{2}$ is the slope above $\left[q, q^{2}\right]$.

Pick an element $y \in F[\pi]-H$; then $v(y)=-s_{2}=v(F) /\left(q^{2}-q\right)$. It follows from (4.1) that the element $z:=\alpha(y) \in H^{\prime}$ has valuation

$$
v(z)=\sum_{x \in H} v\left(y-{ }_{F} x\right)=q \cdot v(y)=\frac{v(F)}{q-1} .
$$

Now if $v(F) \leq 1 / q$ then $v(z) \leq 1 /\left(q^{2}-q\right)$. This means that $-v(z)$ is equal to the slope of the Newton polygon of $[\pi]_{F^{\prime}}$ above the interval $\left[q, q^{2}\right]$. We conclude that

$$
v\left(F^{\prime}\right)=\left(q^{2}-q\right) \cdot v(z)=q \cdot v(F)
$$

and that $H^{\prime}$ is not the canonical subgroup of $F^{\prime}$. On the other hand, if $v(F)>1 / q$ then $v(z)>1 /\left(q^{2}-q\right)$. Therefore, $v(z)$ is equal to the slope above the interval $[1, q]$. We conclude that

$$
v\left(F^{\prime}\right)=1-(q-1) \cdot v(F)=1-v(F)
$$

and that $H^{\prime}$ is the canonical subgroup of $F^{\prime}$. This finishes the proof of Case 1. The proof of Case 2 is similar and left to the reader.
4.3. Let us now assume that the lift $F$ is the canonical lift of $G$ with respect to some fixed embedding $\kappa: L \hookrightarrow D$. Note that we have $v(F)=1$ if $L / K$ is unramified and $v(F)=1 / 2$ if $L / K$ is ramified. In the former case, $F$ has no canonical subgroup, whereas in the latter case the canonical subgroup of $F$ is the kernel of $\left[\pi_{L}\right]_{F}$.

For $s=1,2, \ldots$, we define isogenies $\alpha_{s}: F \rightarrow F_{s}$ inductively, as follows. First, choose a non-canonical $\mathcal{O}_{K}$-submodule $H \subset F[\pi]$ of height one. Set $F_{1}:=F / H$ and let $\alpha_{1}: F \rightarrow F_{1}$ be the natural projection. For $s \geq 1$, choose a non-canonical $\mathcal{O}_{K^{-}}$ submodule $H_{s} \subset F_{s}[\pi]$ of height one, set $F_{s+1}:=F_{s} / H_{s}$ and let $\alpha_{s+1}: F \rightarrow F_{s+1}$ be
the composition of $\alpha_{s}$ with the natural projection $F_{s} \rightarrow F_{s+1}$. As we have seen in the last section, we can see $F_{s}$ as a lift of $G$ in such a way that the isogeny $\alpha_{s}$ reduces to the endomorphism $\Pi^{s}: G \rightarrow G$ modulo the maximal ideal of $A^{\prime}$. This choice is by no means canonical; however, for the statement of the next proposition, the choice that we have made is irrelevant, see Remark 3.3 and Remark 4.2.

Proposition 4.6. - The lift $F_{s}$ is quasi-canonical of level s, and we have

$$
v\left(F_{s}\right)=\left\{\begin{array}{cl}
\frac{1}{q^{s-1}(q+1)}, & \text { if } L / K \text { is unramified and } s \geq 1 \\
\frac{1}{2 q^{s}}, & \text { if } L / K \text { is ramified }
\end{array}\right.
$$

Proof. - We proceed by induction over $s$. We start the induction at $s=1$ if $L / K$ is unramified and at $s=0$ in the ramified case (one has to be careful with the notation: plugging in $s=0$ into $F_{s}$ should be understood as $\left.F\right)$. If $L / K$ is unramified, then $v(F)=1>q /(q+1)$, and Proposition 4.5, Case 2(b), shows that $v\left(F_{1}\right)=1 /(q+1)$. This is indeed as in the statement of the proposition. The statement of the proposition is also true for $s=0$ if $L / K$ is ramified.

Suppose now that $s \geq 1$ or that $L / K$ is ramified. Then $v\left(F_{s}\right) \leq q /(q+1)$, so Proposition 4.5, Case 2(a), shows that

$$
v\left(F_{s+1}\right)=\frac{v\left(F_{s}\right)}{q} .
$$

We see that the formula for $v\left(F_{s}\right)$ follows by induction.
Since $F_{s}$ is isogenous to $F$, it is a quasi-canonical lift of some level. By construction, the isogeny $\alpha_{s}: F \rightarrow F_{s}$ has degree $q^{s}$. Let $n$ be the maximal integer such that $\alpha_{s}$ factors over $\left[\pi_{L}^{n}\right]: F \rightarrow F$. The proof of Theorem 3.2 shows that $F_{s}$ is quasi-canonical of level $s^{\prime}:=s-2 n / e$.

Suppose $n>0$. By the induction hypothesis, $F_{s^{\prime}}$ is quasi-isogenous of level $s^{\prime}$. Therefore, by Remark 3.3, $F_{s^{\prime}}$ and $F_{s}$ are isomorphic as formal $\mathcal{O}_{K^{\prime}}$-modules. But then we have $v\left(F_{s}\right)=v\left(F_{s^{\prime}}\right)$. This gives a contradiction with the formula for $v\left(F_{s}\right)$ which we have already proved. We conclude that $n=0$, i.e., that $F_{s}$ is quasi-canonical of level $s$.

Corollary 4.7. - Let $F_{s}$ be a quasi-canonical lift of level $s$ and $\mathcal{O}_{M_{s}} / A$ be the smallest extension over which it can be defined. Then the formal modulus $u\left(F_{s}\right) \in \mathcal{O}_{M_{s}}$ of $F_{s}$ is a uniformizer for the valuation ring $\mathcal{O}_{M_{s}}$.

Proof. - It follows from Theorem 3.2 that $\mathcal{O}_{M_{s}}$ is the ring of integers of the extension $M_{s} / M$, the ring class field of $\mathcal{O}_{s}^{\times}$. Moreover, we may assume that $F_{s}$ is the lift constructed before Proposition 4.5. Therefore, the formula for $v\left(F_{s}\right)$ in Proposition 4.5 shows that the valuation of $u\left(F_{s}\right)$ is equal to the reciprocal of the degree $\left[M_{s}: M\right]$. This concludes the proof.

Corollary 4.8. - Let $F_{s}$ and $F_{s+1}$ be quasi-canonical lifts of level $s$ and $s+1$, respectively. Let $\beta: F_{s} \rightarrow F_{s+1}$ be an isogeny of height one. Then $H:=\operatorname{ker}(\beta)$ is not the canonical subgroup, and

$$
v(\operatorname{Lie}(\beta))=v\left(F_{s+1}\right)
$$

Proof. - We note that $\beta$ identifies $F_{s+1}$ with the quotient $F_{s} / H$. It follows from the proof of Proposition 4.6 that $H$ is not the canonical subgroup of $F_{s}$ and that therefore the nonzero elements $x \in H$ have valuation

$$
v(x)=\frac{v\left(F_{s}\right)}{q^{2}-q} .
$$

Set $b:=\operatorname{Lie}(\beta)$. The formula for $\beta$ in terms of $H$ (see Theorem 2.1) shows that

$$
v(b)=\sum_{x \in H-\{0\}} v(x)=\frac{v\left(F_{s}\right)}{q} .
$$

By Corollary 4.7 , this is equal to $v\left(F_{s+1}\right)$.

## 5. Some complements

We prove some technical results which are needed in $[\mathbf{R}]$.
5.1. Let $K$ and $k$ be as before. Let $G$ be the formal $\mathcal{O}_{K}$-module of height two over $k$, with endomorphism ring $\mathcal{O}_{D}$. We have seen in $[\mathbf{V Z}]$ that the formal cohomology group $H^{2}(G, k)$ has dimension $h-1=1$. Therefore, the universal deformation ring of $G$ is $W[t]]$ (where $W=\hat{\mathcal{O}}_{K}^{\mathrm{nr}}$ is the completion of the maximal unramified extension of $\mathcal{O}_{K}$ ).

Let $A$ be a complete local $\mathcal{O}_{K}$-algebra with residue field $k$ and $I \triangleleft A$ an ideal with $\mathfrak{m}_{A} \cdot I=0$. Set $\bar{A}:=A / I$. Let $F, F^{\prime}$ be two deformations of $G$ over $A$ and $\bar{\alpha}: F \otimes \bar{A} \rightarrow F^{\prime} \otimes \bar{A}$ a homomorphism which is defined modulo $I$. Then the obstruction for lifting $\bar{\alpha}$ to a homomorphism $\alpha: F \rightarrow F^{\prime}$ is an element of the $k$-vector space

$$
H^{2}(G, I) \cong H^{2}(G, k) \otimes_{k} I
$$

Indeed, as in [VZ, Section 3], a lift $\alpha(X) \in A[[X]$ of $\bar{\alpha}$ as a power series defines a cocycle ( $\Delta ; \delta_{a}$ ),

$$
\begin{aligned}
\Delta(X, Y) & :=\alpha\left(X+{ }_{F} Y\right)-{ }_{F^{\prime}} \alpha(X)-F_{F^{\prime}} \alpha(Y), \\
\delta_{a}(X) & :=\alpha\left([a]_{F}(X)\right)-{ }_{F^{\prime}}[a]_{F^{\prime}}(\alpha(X)) .
\end{aligned}
$$

The cohomology class of this cocycle is independent of the chosen lift $\alpha$. It vanishes if and only if there exists some lift $\alpha$ which is a homomorphism $F \rightarrow F^{\prime}$. If this is the case, then the lift which is a homomorphism is unique.

Let $F$ be the universal deformation of $G$ over $W[[t]]$, and let $F^{\prime}$ be another universal deformation over $W\left[\left[t^{\prime}\right]\right]$. Hence the pair $\left(F, F^{\prime}\right)$ is defined over the formal scheme $\mathcal{S}=\operatorname{Spf} R$, where $R:=W\left[t t, t^{\prime}\right]$.

Proposition 5.1. - Let $\alpha: G \rightarrow G$ be an isogeny, i.e., $\alpha \neq 0$. Let $J$ be the minimal ideal of $R$ such that $\alpha$ lifts to an isogeny $F \rightarrow F^{\prime}$ modulo $J$. Then the closed formal subscheme $\mathcal{T}$ of $\mathcal{S}$ defined by $J$ is a relative divisor over $\operatorname{Spf} W$.

Proof. - We have to show that $J$ is generated by one element which is neither a unit nor divisible by $p$. Suppose, for the moment, that $\alpha \notin \mathcal{O}_{K}$ and set $L=K(\alpha) \subset D$. Let $M$ be the completion of the maximal unramified extension of $L$ and $F_{1}$ the canonical lift of $G$ with respect to $\mathcal{O}_{L} \subset \mathcal{O}_{D}$ (which is defined over $\mathcal{O}_{M}$ ). There is a unique homomorphism of $\mathcal{O}_{K}$-algebras $\left.\varphi: W\left[t t, t^{\prime}\right]\right] \rightarrow \mathcal{O}_{M}$ which induces the identity on $k$, such that the pair $\left(F_{1}, F_{1}\right)$ is $*$-isomorphic to the pullback of the pair $\left(F, F^{\prime}\right)$ via $\varphi$. By construction, $J$ is contained in the kernel of $\varphi$. This shows $J \neq R$, at least if $\alpha \notin K$. The case $\alpha \in K$ is handled in a similar way.

Suppose that $J \subset(\pi)$. This means that $\alpha$ lifts to an isogeny $F \rightarrow F^{\prime}$ over $k\left[\left[t, t^{\prime}\right]\right]$. Setting $t^{\prime}=0$, the isogeny $\alpha$ would then induce an isogeny between $F \otimes_{W \llbracket t]]} k((t))$ and $G \otimes_{k} k((t))$. But $F \otimes_{W[t t]} k((t))$ has height $h-1=1$ (see $\left.[\mathbf{V Z}]\right)$ and is therefore not isogenous to the height-two module $G \otimes_{k} k((t))$. This gives a contradiction and shows that $J \not \subset(\pi)$.

Let $\mathfrak{m}$ denote the maximal ideal of $R$. Set $A:=R / \mathfrak{m} J$ and $I:=J / \mathfrak{m} J$. Then $\mathfrak{m} \cdot I=0$, and $\bar{A}=A / I \cong R / J$. Clearly, $\alpha$ lifts to a homomorphism $F \otimes \bar{A} \rightarrow F^{\prime} \otimes \bar{A}$ but not to a homomorphism $F \otimes A \rightarrow F^{\prime} \otimes A$. The responsible obstruction is a nonzero element in

$$
H^{2}(G, I) \cong H^{2}(G, k) \otimes_{k} I \cong I
$$

Let $f$ be the image of this obstruction in $I$. The element $f$ depends on the choice of an isomorphism $H^{2}(G, k) \cong k$, but the ideal $(f) \triangleleft A$ does not. Clearly, $\alpha$ lifts to a homomorphism $F \otimes A^{\prime} \rightarrow F^{\prime} \otimes A^{\prime}$ over the ring $A^{\prime}=A /(f)$. This implies $I=(f)$. Now Nakayama's Lemma shows that $J$ is generated by one element. The proposition is proved.
5.2. Let $A$ be the ring of integers of a finite extension of the fraction field of $W$. Let $\lambda$ denote a uniformizer of $A$. For each positive integer $n$, we set $A_{n}:=A /\left(\lambda^{n+1}\right)$ and $M_{n}:=\left(\lambda^{n}\right) /\left(\lambda^{n+1}\right)$.

Let $F_{1}, F_{2}, F_{3}$ be three lifts of $G$ over $A$. We define

$$
H_{n}:=\operatorname{Hom}\left(F_{1} \otimes A_{n}, F_{2} \otimes A_{n}\right), \quad H_{n}^{\prime}:=\operatorname{Hom}\left(F_{1} \otimes A_{n}, F_{3} \otimes A_{n}\right) .
$$

As for endomorphisms, the natural reduction maps $H_{n}, H_{n}^{\prime} \rightarrow \operatorname{End}(G)=\mathcal{O}_{D}$ are injective. We will consider $H_{n}$ and $H_{n}^{\prime}$ as subsets of $\mathcal{O}_{D}$. Note that $H_{n}$ and $H_{n}^{\prime}$ are in fact sub- $\mathcal{O}_{K}$-modules of $\mathcal{O}_{D}$. The obstruction theory reviewed above gives injective maps

$$
\kappa_{n}: H_{n-1} / H_{n} \hookrightarrow H^{2}\left(G, M_{n}\right), \quad \kappa_{n}^{\prime}: H_{n-1}^{\prime} / H_{n}^{\prime} \hookrightarrow H^{2}\left(G, M_{n}\right)
$$

Proposition 5.2. - Let $\alpha: G \rightarrow G$ be an isogeny defined over $k$ which does not lift to a homomorphism $F_{1} \rightarrow F_{2}$. Let $n$ be the unique positive integer such that $\alpha \in$
$H_{n-1}-H_{n}$. Let $\beta: F_{2} \rightarrow F_{3}$ be an isogeny defined over $A$, and let $m$ denote the valuation of $b:=\operatorname{Lie}(\beta) \in A$. We make the following assumptions:

1. $\beta$ has height one.
2. $m \leq(q-1) n$.

Then $\beta \circ \alpha \in H_{n+m-1}^{\prime}-H_{n+m}^{\prime}$.
Proof. - (compare with the proof of Lemma 1.3) We may represent $\alpha$ as a power series with coefficients in $A$ without constant coefficient such that $\alpha_{n-1}$, the reduction of $\alpha$ modulo $\lambda^{n}$, is a homomorphism $F_{1} \otimes A_{n-1} \rightarrow F_{2} \otimes A_{n-1}$. We define

$$
\epsilon:=\alpha \circ[\pi]_{F_{1}}-F_{2}[\pi]_{F_{2}} \circ \alpha
$$

Then $\epsilon \equiv 0\left(\bmod \lambda^{n}\right)$. Moreover, we have $\epsilon \equiv \delta_{\pi}\left(\bmod \lambda^{n+1}\right)$, where $\left(\Delta,\left\{\delta_{a}\right\}\right)$ denotes the cocycle associated to $\alpha_{n}$. The assumption $\alpha \notin H_{n}$ implies $\epsilon(X)=c X^{q}+$ $\ldots$, with $\operatorname{ord}_{\lambda}(c)=n$. Similarly, define

$$
\epsilon^{\prime}:=\beta \circ \alpha \circ[\pi]_{F_{1}}-F_{3}[\pi]_{F_{3}} \circ \beta \circ \alpha .
$$

Then $\epsilon^{\prime}=\beta \circ \epsilon$. Write $\beta(X)=\sum_{i} b_{i} X^{i}$. It follows from Assumption 1 that the Newton polygon of $\beta$ has slope $-m /(q-1)$ over $[1, \ldots, q]$. This means that

$$
\operatorname{ord}_{\lambda}\left(b_{i}\right) \geq \frac{q-i}{q-1} \cdot m, \quad i=1, \ldots, q
$$

(with equality for $i=1, q$ ). Now Assumption 2, together with an easy calculation, shows that

$$
\epsilon^{\prime}=\beta \circ \epsilon=b_{1} c X^{q}+\cdots \equiv 0 \quad\left(\lambda^{n+m}\right)
$$

Since $\operatorname{ord}_{\lambda}\left(b_{1} c\right)=n+m$, we conclude as in the proof of Lemma 1.3 that $\beta \circ \alpha \in$ $H_{n+m-1}^{\prime}-H_{n+m}^{\prime}$.

Corollary 5.3. - Suppose that $F_{1}, F_{2}, F_{3}$ are quasi-canonical liftings of $G$ of level $r, s, s+1$ (with respect to some embedding $\kappa: L \hookrightarrow D$ ). Suppose that $r \leq s$. Suppose, moreover, that $A$ is the minimal $\mathcal{O}_{K}$-algebra over which the lifts $F_{1}, F_{2}, F_{3}$ can be defined. (By Theorem 3.2 and Corollary 4.7, A is the ring of integers of the ring class extension of $\mathcal{O}_{s+1}$.)

Let $\alpha: G \rightarrow G$ be an element of $\mathcal{O}_{D}$ and $\beta: F_{2} \rightarrow F_{3}$ an isogeny of height one, defined over $A$. We assume that $\alpha$ does not lift to a homomorphism $F_{1} \rightarrow F_{2}$. Let $n$ be the maximal integer such that $\alpha$ can be lifted to a homomorphism $F_{1} \rightarrow F_{2}$ modulo $\lambda^{n}$. Then $\beta \circ \alpha$ can be lifted to a homomorphism $F_{1} \rightarrow F_{3}$ modulo $\lambda^{n+1}$, but not modulo $\lambda^{n+2}$.

Proof. - It follows from Corollary 4.8 that $\operatorname{ord}_{\lambda}(\operatorname{Lie}(\beta))=1$. Hence we can apply Proposition 5.2, which proves the corollary.

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[^4]
# 9. CANONICAL AND QUASI-CANONICAL LIFTINGS IN THE SPLIT CASE 

by

Volker Meusers


#### Abstract

Following Gross we sketch a theory of quasi-canonical liftings when the formal $\mathcal{O}_{K^{-}}$module of height two and dimension one is replaced by a divisible $\mathcal{O}_{K^{-}}$ module of height one and dimension one in the sense of Drinfel'd.


Résumé (Relèvements canoniques et quasi-canoniques dans le cas déployé). - Suivant Gross, on donne une théorie de relèvements quasi-canoniques dans le cas où le $\mathcal{O}_{K^{-}}$ module de hauteur deux et de dimension un est remplacé par un $\mathcal{O}_{K}$-module divisible de hauteur un et de dimension un au sens de Drinfel'd.

In this paper, we follow up on a remark by Gross $[\mathbf{G}]$ and discuss a theory of quasi-canonical liftings when the formal $\mathcal{O}_{K}$-module of height two and dimension one considered in $[\mathbf{W w} \mathbf{1}]$ is replaced by a divisible $\mathcal{O}_{K}$-module of height one and dimension one in the sense of Drinfel'd [D]. In this situation the statements analogous to those in $[\mathbf{G}],[\mathbf{W w 1}]$ are easy consequences of Lubin-Tate theory and of a slight modification of the Serre-Tate theorem for ordinary elliptic curves, as discussed in the appendix to [Mes].

## 1. Formal moduli of divisible $\mathcal{O}_{K}$-modules

Let $K$ be a field complete with respect to some discrete valuation. Let $\mathcal{O}_{K}$ be its ring of integers, $\mathfrak{p}=(\pi)$ its maximal ideal. We assume the residue field $\mathcal{O}_{K} / \mathfrak{p}$ to be finite and let $q$ denote the number of its elements. For any non-zero ideal $\mathfrak{a} \subset \mathcal{O}_{K}$ we set $N(\mathfrak{a}):=\left|\mathcal{O}_{K} / \mathfrak{a}\right|$, i.e., $N\left(\mathfrak{p}^{s}\right)=q^{s}$. Let $k$ be an algebraic closure of $\mathcal{O}_{K} / \mathfrak{p}$. Let $M$ be the completion of the maximal unramified extension of $K$ in some fixed separable closure $K^{\text {sep }}$. Denote the completion of $K^{\text {sep }}$ by $C$. Let $\mathcal{O}_{M}$ and $\mathcal{O}_{C}$ be the rings of integers in $M$ and $C$ respectively.

Following $[\mathbf{D}, \S 4]$ a formal group is a group object in the category of formal schemes. For example any group scheme or any discrete group is a formal group in this sense.

2000 Mathematics Subject Classification. - 11G15, 14K07, 14K22, 14L05.
Key words and phrases. - Quasi-canonical liftings, complex multiplication, Lubin-Tate formal groups, Serre-Tate theorem.

For a formal group $F$ let us denote by $F^{\circ}$ its connected component. Let $\widehat{\mathcal{C}}$ be the category of complete local noetherian $\mathcal{O}_{M}$-algebras with residue field $k$.

Definition 1.1. - Let $R \in \widehat{\mathcal{C}}$. A divisible $\mathcal{O}_{K}$-module over $R$ is a pair $F$, where $F$ is a formal group over $R$ and $\gamma_{F}: \mathcal{O}_{K} \rightarrow \operatorname{End}_{R}(F)$ is a homomorphism such that $F^{\circ}$ is a formal $\mathcal{O}_{K}$-module of height $h<\infty$ in the sense of $[\mathbf{V Z}]$, and such that

$$
F / F^{\circ} \cong\left(K / \mathcal{O}_{K}\right)_{\operatorname{Spf}(R)}^{j}
$$

for some $j<\infty$. The pair $(h, j)$ will be called type of $F$.
To ease the notation, we will suppress the structure map $\gamma_{F}$ of an $\mathcal{O}_{K}$-module $F$ and simply write $F$.

Drinfel'd shows that a divisible $\mathcal{O}_{K}$-module over $k$ is up to isomorphism given by its type $(h, j)$ (see $[\mathbf{D}, \S 4]$ ).

Example 1.2.- For $K=\mathbb{Q}_{p}, \mathcal{O}_{K}=\mathbb{Z}_{p}$ the product group $G=\widehat{\mathbb{G}}_{m, R} \times\left(\mathbb{Q}_{p} / \mathbb{Z}_{p}\right)_{R}$ is an example of a divisible module of type $(h, j)=(1,1)$ over $R$.

If $R \in \widehat{\mathcal{C}}$ is artinian then the category of fppf-abelian sheaves on $R$ with $\mathcal{O}_{K^{-}}$ structure is an abelian category, the category of $\mathcal{O}_{K}$-modules over $R$. It is useful to view the category of divisible $\mathcal{O}_{K}$-modules over $R$ as a full sub-category of this category.

Definition 1.3. - Fix a divisible $\mathcal{O}_{K}$-module $G$ over $k$. A deformation of $G$ to $R \in$ $\widehat{\mathcal{C}}$ is a pair $(F, \psi)$ consisting of a divisible $\mathcal{O}_{K}$-module $F$ over $R$ together with an isomorphism $\psi: F \otimes_{R} k \stackrel{\cong}{\cong} G$ of $\mathcal{O}_{K}$-modules.

The deformations of $G$ to $R \in \widehat{\mathcal{C}}$ form a category in a natural way. One checks that it is a groupoid and moreover that no object of this groupoid has non-trivial automorphisms. The last point is due to the fact that for a deformation $F$ the isomorphism $\psi$ is part of the data. Nevertheless we often omit $\psi$ from the notation.

Definition 1.4. - For any $R \in \widehat{\mathcal{C}}$ let us denote by $\mathcal{D}_{G}(R)$ the set of isomorphism classes of the groupoid of deformations of $G$ to $R$. Then $\mathcal{D}_{G}$ becomes a set-valued functor on $\widehat{\mathcal{C}}$.

Fix a formal $\mathcal{O}_{K}$-module $H_{0}$ of height $h=1$ over $k$. It has a trivial deformation space, i.e., $\mathcal{D}_{H_{0}}(R)=\{$ point $\}$ for any $R \in \widehat{\mathcal{C}}$. More precisely $\mathcal{D}_{H_{0}}$ is representable by $\mathcal{O}_{M}$. This follows easily from the uniqueness of Lubin-Tate modules (see [Me1]; see also Remark 1.11(ii) for a far more general result of Drinfel'd). Let us denote by $H$ the unique lift of $H_{0}$ to $\mathcal{O}_{M}$. We assume, as we may, that $H$ is given as the base change

$$
H=H_{f} \otimes_{\mathcal{O}_{K}} \mathcal{O}_{M}
$$

where $H_{f}$ is the Lubin-Tate module over $\mathcal{O}_{K}$ corresponding to some fixed prime element $\pi \in \mathcal{O}_{K}$ and some fixed Lubin-Tate series $f \in \mathcal{F}_{\pi}$. Recall from [Me1, Lemma 1.7] that the isomorphism class of $H$ does not depend on these choices. Recall further that for any $R \in \widehat{\mathcal{C}}$ we have $H(R)=\mathfrak{m}_{R}$ as a set. The $\mathcal{O}_{K}$-module structure is given as follows: For $q, q^{\prime} \in H(R)$ and $z \in \mathcal{O}_{K}$ we have $q+_{H} q^{\prime}=H\left(q, q^{\prime}\right)$ and $z \cdot{ }_{H} q=[z]_{f}(q)$. We often omit the subscript $H$ from the notation.

Now fix some divisible $\mathcal{O}_{K}$-module $G$ over $k$ of height $h=1$ such that there is an isomorphism $G / G^{\circ} \cong\left(K / \mathcal{O}_{K}\right)_{k}$. Fix an isomorphism of divisible $\mathcal{O}$-modules

$$
r: G \stackrel{\cong}{\cong} H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}
$$

where $H$ is the unique lift of $G^{\circ}$ to $\mathcal{O}_{M}$ as above. Two such isomorphisms differ by an element of the automorphism group of the right hand side. This group is described by the following easy but important lemma.

## Lemma 1.5

(1) We have

$$
\operatorname{Hom}_{\mathcal{O}_{K}, k}\left(\left(K / \mathcal{O}_{K}\right)_{k}, H_{0}\right)=\{0\}=\operatorname{Hom}_{\mathcal{O}_{K}, k}\left(H_{0},\left(K / \mathcal{O}_{K}\right)_{k}\right)
$$

and

$$
\operatorname{End}_{\mathcal{O}_{K}, k}\left(H_{0}\right)=\mathcal{O}_{K}=\operatorname{End}_{\mathcal{O}_{K}, k}\left(\left(K / \mathcal{O}_{K}\right)_{k}\right)
$$

(2) In particular there is a canonical isomorphism

$$
\mathcal{O}_{K} \times \mathcal{O}_{K} \longrightarrow \operatorname{End}_{\mathcal{O}_{K}, k}\left(H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}\right) .
$$

It induces an isomorphism

$$
\mathcal{O}_{K}^{\times} \times \mathcal{O}_{K}^{\times} \longrightarrow \operatorname{Aut}_{\mathcal{O}_{K}, k}\left(H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}\right) .
$$

Proof. - It clearly suffices to prove the first point. We have

$$
\operatorname{Hom}_{\mathcal{O}_{K}, k}\left(\left(K / \mathcal{O}_{K}\right)_{k}, H_{0}\right)=\operatorname{Hom}_{\mathcal{O}_{K}}\left(K / \mathcal{O}_{K}, H_{0}(k)\right)=\{0\}
$$

by adjunction and because $H_{0}(k)=\{0\}$. We have

$$
\operatorname{Hom}_{\mathcal{O}_{K}, k}\left(H_{0},\left(K / \mathcal{O}_{K}\right)_{k}\right)=\operatorname{Hom}_{\mathcal{O}_{K}, k}\left(H_{0},\left(K / \mathcal{O}_{K}\right)_{k}^{\circ}\right)=\{0\}
$$

because $H_{0}$ is connected and $\left(K / \mathcal{O}_{K}\right)^{\circ}=\{0\}$. We have

$$
\operatorname{End}_{\mathcal{O}_{K}, k}\left(H_{0}\right)=\mathcal{O}_{K}
$$

because by Lubin-Tate theory every endomorphism of $H_{0}$ is uniquely given by its differential at zero. We have

$$
\operatorname{End}_{\mathcal{O}_{K}, k}\left(\left(K / \mathcal{O}_{K}\right)_{k}\right)=\operatorname{End}_{\mathcal{O}_{K}}\left(K / \mathcal{O}_{K}\right)
$$

by adjunction. Since the natural map

$$
\mathcal{O}_{K} \longrightarrow \operatorname{End}_{\mathcal{O}_{K}}\left(K / \mathcal{O}_{K}\right)
$$

is well known to be an isomorphism we are done.
We want to sketch a proof of the following theorem (compare the analogous statement in [VZ, Theorem 3.8]):

Theorem 1.6 (Universal deformation). - For any $R \in \widehat{\mathcal{C}}$ and fixed isomorphism $r$ there is a natural isomorphism

$$
\eta_{R}: \mathcal{D}_{G}(R) \xrightarrow{\cong} H(R) .
$$

In particular $\mathcal{D}_{G}$ can be given the structure of an $\mathcal{O}_{K}$-module (depending on $r$ of course). Since we assume $H=H_{f} \otimes \mathcal{O}_{K} \mathcal{O}_{M}$, the $\mathcal{O}_{K}$-module structure is given by Lubin-Tate theory as recalled above.

The proof will take up the rest of this section. One proceeds as in [Mes, appendix]: In the course of the proof we will identify both, $\mathcal{D}_{G}(R)$ and $H(R)$ for $R \in \widehat{\mathcal{C}}$ artinian, with a certain Ext-group. So let us briefly recall the definition and some basic properties of these groups. A careful discussion can be found in [Mt, chapter VII].

For objects $M^{\prime \prime}$ and $M^{\prime}$ of an abelian category $\mathcal{A}$ let

$$
{\mathcal{E} x t_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)}
$$

denote the groupoid of extensions $(M, p, i): M^{\prime} \stackrel{i}{\hookrightarrow} M \xrightarrow{p} M^{\prime \prime}$. It is well known that the map

$$
\begin{aligned}
\operatorname{Hom}_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right) & \longrightarrow \operatorname{Aut}_{\mathcal{E x t} t_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)}((M, p, i)) \\
\varphi & \longmapsto \operatorname{id}_{M}+i \circ \varphi \circ p
\end{aligned}
$$

is an isomorphism of groups. In particular the automorphism group of ( $M, p, i$ ) is trivial if and only if $\operatorname{Hom}_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)$ is. Let

$$
\operatorname{Ext}_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)
$$

be the class of isomorphism classes of $\mathcal{E} x t_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)$. Assume it to be a set. Sometimes we will not distinguish an extension from its isomorphism class. Using Baeraddition $\operatorname{Ext}_{\mathcal{A}}\left(M^{\prime \prime}, M^{\prime}\right)$ becomes an abelian group in the usual way. For $N^{\prime} \in \mathcal{A}$ let

$$
\begin{equation*}
\delta_{(M, p, i), N^{\prime}}: \operatorname{Hom}_{\mathcal{A}}\left(M^{\prime}, N^{\prime}\right) \longrightarrow \operatorname{Ext}_{\mathcal{A}}\left(M^{\prime \prime}, N^{\prime}\right) \tag{1.1}
\end{equation*}
$$

be the boundary homomorphism.
Apply this in the case that $\mathcal{A}$ is the category of $\mathcal{O}_{K}$-modules on some fixed artinian $R \in \widehat{\mathcal{C}}$. In this case the Ext-groups are in fact $\mathcal{O}_{K}$-modules.

Definition 1.7. - Let $R \in \widehat{\mathcal{C}}$ be artinian. For any two $\mathcal{O}_{K}$-modules $M^{\prime}$ and $M^{\prime \prime}$ over $R$ let

$$
\operatorname{Ext}_{\mathcal{O}_{K}, R}\left(M^{\prime \prime}, M^{\prime}\right)
$$

denote the $\mathcal{O}_{K}$-module of extension classes of $M^{\prime \prime}$ by $M^{\prime}$ constructed above.
Recall that we view the category of divisible $\mathcal{O}_{K}$-modules on $\operatorname{artinian} R$ as a full sub-category of the category of all $\mathcal{O}_{K}$-modules.

Lemma 1.8 (compare [Mes, I.2.4.3]). - Let $R \in \widehat{\mathcal{C}}$ be artinian. Given an extension of the form

$$
H_{R} \stackrel{i}{\longleftrightarrow} F \stackrel{p}{\longleftrightarrow}\left(K / \mathcal{O}_{K}\right)_{R}
$$

of $\mathcal{O}_{K}$-modules over $R$, then $F$ is a divisible $\mathcal{O}_{K}$-module such that $F^{\circ} \cong H_{R}$ and $F / F^{\circ} \cong\left(K / \mathcal{O}_{K}\right)_{R}$. If one uses the isomorphism $r: G \xlongequal{\rightrightarrows} H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}$ then $F$ becomes a deformation of $G$ to $R$. This association yields a functor between the groupoid of extensions of $\left(K / \mathcal{O}_{K}\right)_{R}$ by $H_{R}$ and the groupoid of deformations of $G$ to $R$.

Proof. - Since $\left(K / \mathcal{O}_{K}\right)_{R}$ is totally disconnected and $H_{R}$ is connected it follows that $i: H_{R} \xlongequal{\cong} F^{\circ}$. The snake lemma implies that $p$ induces an isomorphism $p^{\prime}: F / F^{\circ} \xlongequal{\rightrightarrows}$ $\left(K / \mathcal{O}_{K}\right)_{R}$. It follows that $F$ is divisible. Since $H_{R}(k)=\{0\}$ the extension $H_{R} \hookrightarrow$ $F \rightarrow\left(K / \mathcal{O}_{K}\right)_{R}$ yields an injective map $F(k) \hookrightarrow\left(K / \mathcal{O}_{K}\right)_{R}(k)=K / \mathcal{O}_{K}$. Since $k$ is algebraically closed it is an isomorphism. This isomorphism gives us a canonical splitting map $\left(K / \mathcal{O}_{K}\right)_{k} \hookrightarrow F \otimes k$. Thus the extension is canonically split over $k$. Together with the identification $r: G \stackrel{\cong}{\rightrightarrows} H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}$ we get an isomorphism $\psi: F \otimes k \stackrel{\cong}{\rightrightarrows} G$ such that the pair $(F, \psi)$ is a deformation of $G$. One checks that it is functorial.

Proposition 1.9 (compare [Mes, appendix Prop.2.1]). - Assume $R \in \widehat{\mathcal{C}}$ to be artinian. Then the functor of the preceding lemma is an equivalence of groupoids and there is a natural isomorphism

$$
\epsilon_{R}: \mathcal{D}_{G}(R) \xrightarrow{\cong} \operatorname{Ext}_{\mathcal{O}_{K}, R}\left(\left(K / \mathcal{O}_{K}\right)_{R}, H_{R}\right) .
$$

Proof. - fully faithful: It is enough to see that every object in either groupoid has a trivial automorphism group. For deformations, this was noted above. For extensions, recall that the automorphism group is isomorphic to $\operatorname{Hom}_{\mathcal{O}_{K}, R}\left(\left(K / \mathcal{O}_{K}\right)_{R}, H_{R}\right)=\{0\}$.
essentially surjective: Let $F$ be a deformation of $G$ to $R$. We need to define homomorphisms $i: H_{R} \hookrightarrow F$ and $p: F \rightarrow\left(K / \mathcal{O}_{K}\right)_{R}$ such that $p \circ i=0$. For this we let $p$ on $R$-valued points be defined as follows:

$$
F(R) \longrightarrow F(k)=F \otimes k(k) \underset{r \circ \psi}{\cong} H_{0}(k) \times\left(K / \mathcal{O}_{K}\right)_{k}(k) \underset{\text { pr }_{2}}{\longrightarrow} K / \mathcal{O}_{K}=\left(K / \mathcal{O}_{K}\right)_{R}(R) .
$$

Since $K / \mathcal{O}_{K}$ is discrete the kernel of $p$ equals $F^{\circ}$. Because $R$ is artinian local it follows that $F^{\circ} \otimes k=(F \otimes k)^{\circ} \cong G^{\circ} \cong H_{0}$. Since $H_{R}$ is the unique lift of $H_{0}$ to $R$ it follows that $F^{\circ}$ is isomorphic to $H_{R}$ and we get the map $i: H_{R} \cong F^{\circ} \hookrightarrow F$. This proves the first assertion. The second follows by passage to isomorphism classes.

To calculate the Ext-group, we use
Proposition 1.10. - For any artinian $R \in \widehat{\mathcal{C}}$ the connecting homomorphism associated to the sequence $\mathcal{O}_{K} \hookrightarrow K \longrightarrow K / \mathcal{O}_{K}$ is an isomorphism

$$
\delta_{R}: H(R)=\operatorname{Hom}_{\mathcal{O}_{K}, R}\left(\mathcal{O}_{K}, H_{R}\right) \xrightarrow{\cong} \operatorname{Ext}_{\mathcal{O}_{K}, R}\left(\left(K / \mathcal{O}_{K}\right)_{R}, H_{R}\right)
$$

Proof. - Assume $\mathfrak{m}_{R}^{n+1}=0$ for some $n \gg 0$. Then $H$ is killed by $\mathfrak{p}^{n}$ (compare $[\mathbf{K}$, Lemma 1.1.2]). Associated to the short exact sequence

$$
\left(\mathcal{O}_{K}\right)_{R} \stackrel{i}{\longrightarrow} K_{R} \xrightarrow{p}\left(K / \mathcal{O}_{K}\right)_{R}
$$

and $H_{R}$ we have the boundary map (1.1)

$$
\delta_{\left(K_{R}, p, i\right), H_{R}}: \operatorname{Hom}_{\mathcal{O}_{K}, R}\left(\left(\mathcal{O}_{K}\right)_{R}, H_{R}\right) \longrightarrow \operatorname{Ext}_{\mathcal{O}_{K}, R}\left(\left(K / \mathcal{O}_{K}\right)_{R}, H_{R}\right)
$$

If we identify $H(R)$ with $\operatorname{Hom}_{\mathcal{O}_{K}, R}\left(\left(\mathcal{O}_{K}\right)_{R}, H_{R}\right)$ this gives us the desired map $\delta_{R}$. Because the prime element $\pi \in \mathcal{O}_{K}$ acts invertibly on $K$ and nilpotently on $H$ one sees easily that

$$
\operatorname{Hom}_{\mathcal{O}_{K}, R}(K, H)=\{0\}=\operatorname{Ext}_{\mathcal{O}_{K}, R}(K, H)
$$

By the exactness of the long Ext-sequence, it follows that $\delta_{R}$ is an isomorphism.
Proof of Theorem 1.6. - Combining Proposition 1.9 and Proposition 1.10 we get the desired isomorphism for artinian $R \in \widehat{\mathcal{C}}$ as

$$
\eta_{R}=\delta_{R}^{-1} \circ \epsilon_{R}
$$

For general $R$ we can pass to the limit over its artinian quotients.

## Remark 1.11

(i) How does one calculate the inverse of $\delta_{R}$ ? For $R=k$ both sides are trivial and so is $\delta_{k}$. In the general case $\delta_{R}^{-1}$ can be computed by an approximation process with respect to the "p-adic topology" on both $\operatorname{Ext}_{\mathcal{O}_{K}, R}\left(\mathcal{O}_{K}, H_{R}\right)$ and $H(R)$. For details we refer to [K, page 151f], [Mes, appendix].
(ii) In particular it follows from this theorem that the formal moduli space of the divisible module $G=H_{0} \times\left(K / \mathcal{O}_{K}\right)_{k}$ is representable by a formal power series ring in one variable over $\mathcal{O}_{M}$. More generally, Drinfel'd shows that the formal moduli space of a divisible module of type $(h, j)$ over $k$ is representable by a power series ring in $h+j-1$ variables (compare [D, Prop.4.5]).

Definition 1.12. - For $R \in \widehat{\mathcal{C}}$ and fixed $r$, let $F$ be a lift of $G$ to $R$. Let us set

$$
q(F, r)=\eta_{R}(\text { isom. class of } F) \in H(R) .
$$

We simply write $q(F)$ if $\gamma_{F}$ and $r$ are understood. As in [ $\mathbf{W w} \mathbf{1}$ ], Definition 4.1 we refer to the element $q(F) \in H(R)=\mathfrak{m}_{R}$ as the formal modulus or coordinate of the lift $F$.

Example 1.13. - If $K=\mathbb{Q}_{p}, \mathcal{O}_{K}=\mathbb{Z}_{p}$, and $H=\widehat{\mathbb{G}}_{m}$ we are in the situation of [Mes], Appendix. If we let $q_{\text {Tate }}(F) \in 1+\widehat{\mathbb{G}}_{m}(R)$ denote the coordinate introduced in [Mes], then the relations are simply

$$
q_{\text {Tate }}(F)=1+q(F) \in 1+\widehat{\mathbb{G}}_{m}(R)
$$

and

$$
q_{\text {Tate }}(F)^{z}=(1+q(F))^{z}=1+z \cdot \widehat{\mathbb{G}}_{m} q(F)
$$

## 2. Lifting endomorphisms

Let $F$ and $F^{\prime}$ be deformations of $G$ to $R$ with coordinates $q=q(F), q^{\prime}=q\left(F^{\prime}\right) \in$ $H(R)$. We want to describe in terms of our chosen coordinates which endomorphisms $\rho_{0} \in \operatorname{End}_{\mathcal{O}_{K}, R}(G)$ lift to homomorphisms $\rho: F \rightarrow F^{\prime}$.

Proposition 2.1 (compare [Mes, Appendix Prop.3.3]). - Let $\rho_{0}: F_{0} \rightarrow F_{0}^{\prime}$ be given by multiplication by $z_{1}$ on $\left(K / \mathcal{O}_{K}\right)_{R}$ and by multiplication by $z_{0}$ on $H(R)$. Then $\rho_{0}$ lifts to a (necessarily unique) homomorphism $\rho: F \rightarrow F^{\prime}$ if and only if we have the equality

$$
z_{1} q-z_{0} q^{\prime}=0 \in H(R)
$$

where the last expression is more precisely written as $\left[z_{1}\right]_{H}(q)-_{H}\left[z_{0}\right]_{H}\left(q^{\prime}\right)$.
Sketch of proof. - This follows from rigidity (see [VZ, Lemma 2.6], for formal $\mathcal{O}_{K^{-}}$ modules), the description of lifts in terms of extensions and the following well known and simple lemma applied to $M^{\prime}=N^{\prime}=H, M^{\prime \prime}=N^{\prime \prime}=K / \mathcal{O}_{K}$ and $\varphi=z_{1}$ and $\psi=z_{0}$.

Lemma 2.2 (compare [CE, chap.XIV, exercise 18]). - Let

be a commutative diagram in an arbitrary abelian category. Then it can be completed by a homomorphism $\rho: M \rightarrow N$ if and only if the extension obtained by pushing out the upper sequence along $\varphi$ is isomorphic to the extension obtained by pulling back the lower sequence along $\psi$.

Example 2.3. - For reasons explained above (see [Me1, Example 1.3]), the analogous formula of [Mes], Appendix reads:

$$
\left(q_{\text {Tate }}\right)^{z_{1}}\left(q_{\text {Tate }}^{\prime}\right)^{-z_{0}}=(1+q)^{z_{1}}\left(1+q^{\prime}\right)^{-z_{0}}=1+\left(z_{1} q-_{\widehat{\mathbb{G}}_{m}} z_{0} q^{\prime}\right)=1 .
$$

Specialize to $R=\mathcal{O}_{C}$. As a consequence of proposition (1.9) we can describe the ring of endomorphisms of a lift $F$ of $F_{0}$ to $\mathcal{O}_{C}$.

Corollary 2.4. - Let $F$ be a lift $F$ of $G$ to $\mathcal{O}_{C}$ with $q=q(F, r) \in H\left(\mathcal{O}_{C}\right)$. Then there are two cases:
(i) If the annihilator of $q$ in $\mathcal{O}_{K}$ is zero then the endomorphism ring of $F$ equals $\mathcal{O}_{K}$.
(ii) If the annihilator of $q$ in $\mathcal{O}_{K}$ is $\mathfrak{p}^{s}$ for some $0 \leq s<\infty$ then the endomorphism ring of $F$, as a subring of the ring of endomorphisms of $G$, is strictly bigger then $\mathcal{O}_{K}$ and is isomorphic to

$$
\operatorname{End}_{\mathcal{O}_{K}, \mathcal{O}_{C}}(F) \cong\left\{\left(z_{0}, z_{1}\right) \in \mathcal{O}_{K} \times \mathcal{O}_{K} \mid z_{0}-z_{1} \in \mathfrak{p}^{s}\right\} \subseteq \mathcal{O}_{K} \times \mathcal{O}_{K}
$$

Proof. - This follows directly from the proposition with $q=q^{\prime}$. Note that in this case

$$
\left(z_{1} \cdot{ }_{H} q\right)-_{H}\left(z_{0} \cdot \cdot_{H} q\right)=\left(z_{1}-z_{0}\right) \cdot{ }_{H} q=0 \in H(R) .
$$

## 3. Quasi-canonical lifts in the split case

We now show that the results on canonical and quasi-canonical liftings in [Ww1] and $[\mathbf{G}]$ have analogues in the present case. To bring out this analogy we introduce the following definitions:

## Definition 3.1

(i) Set $L=K \times K$ and $\mathcal{O}_{L}=\mathcal{O}_{K} \times \mathcal{O}_{K}$. Embed $K$ resp. $\mathcal{O}_{K}$ diagonally into $L$ resp. $\mathcal{O}_{L}$.
(ii) From Lemma 1.5 we get an $\mathcal{O}_{K}$-linear isomorphism

$$
\kappa: \mathcal{O}_{L} \xrightarrow{\cong} \operatorname{End}_{\mathcal{O}_{K}, k}(G)
$$

(iii) The "completion of the maximal unramified extension" of $L$ is given by $M_{L}=$ $M \times M$ whose "separable closure" is $M_{L}^{\text {sep }}=M^{\text {sep }} \times M^{\text {sep }}$.
(iv) Set

$$
\Gamma_{L}=\operatorname{Gal}\left(M_{L}^{\mathrm{sep}} \mid M_{L}\right)=\operatorname{Gal}\left(M^{\mathrm{sep}} \mid M\right) \times \operatorname{Gal}\left(M^{\mathrm{sep}} \mid M\right)
$$

By Lubin-Tate theory we have a reciprocity isomorphism

$$
\rho_{K}^{\mathrm{ab}}: \operatorname{Gal}\left(M^{\mathrm{sep}} \mid M\right)^{\mathrm{ab}} \xlongequal{\cong} \mathcal{O}_{K}^{\times} .
$$

It induces a reciprocity isomorphism

$$
\rho_{L}^{\mathrm{ab}}=\left(\rho_{K}^{\mathrm{ab}}, \rho_{K}^{\mathrm{ab}}\right): \Gamma_{L}^{\mathrm{ab}} \xrightarrow{\cong} \mathcal{O}_{L}^{\times} .
$$

(v) For any integer $s \geq 0$ let

$$
\mathcal{O}_{s}=\mathcal{O}_{K}+\mathfrak{p}^{s} \mathcal{O}_{L}=\left\{\left(z_{0}, z_{1}\right) \in \mathcal{O}_{L} \mid z_{0}-z_{1} \in \mathfrak{p}^{s}\right\}
$$

be the "order" containing $\mathcal{O}_{K}$ of conductor $\mathfrak{p}^{s}$ or level $s$ in $\mathcal{O}_{L}$.
(vi) For $s \geq 1$ let $M_{s} \mid M$ be the fixed field in $M^{\text {sep }}$ of the inverse image under the reciprocity isomorphism $\rho_{K}^{\text {ab }}$ of $\left(1+\mathfrak{p}^{s}\right) \subset \mathcal{O}_{K}^{\times}$in $\operatorname{Gal}\left(M^{\text {sep }} \mid M\right)$ ), i.e., such that reciprocity gives an isomorphism

$$
\rho_{K}^{\mathrm{ab}}: \operatorname{Gal}\left(M_{s} \mid M\right) \xrightarrow{\cong} \mathcal{O}_{K}^{\times} /\left(1+\mathfrak{p}^{s}\right)
$$

Remark 3.2. - One easily sees that the map $\mathcal{O}_{L}^{\times} \rightarrow \mathcal{O}_{K}^{\times}$given by sending $(x, y) \in \mathcal{O}_{L}^{\times}$ to the quotient $x y^{-1} \in \mathcal{O}_{K}^{\times}$induces an isomorphism

$$
\mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times} \xrightarrow{\cong} \mathcal{O}_{K}^{\times} /\left(1+\mathfrak{p}^{s}\right) .
$$

If we let $\Gamma_{s} \subset \Gamma_{L}^{\mathrm{ab}}$ be the inverse image of $\mathcal{O}_{s}^{\times}$in $\Gamma_{L}^{\mathrm{ab}}$ under $\rho_{L}^{\mathrm{ab}}$, then we have the following commutative diagram

where " $\cong$ " denotes isomorphisms. In this sense we may consider $M_{s} \mid M$ to be the "ring class field" of the "order" $\mathcal{O}_{s} \subseteq \mathcal{O}_{L}$.

Definition 3.3. - A quasi-canonical lift of $G$ of level $s \geq 0$ ( with respect to $\kappa$ ) is a lift $F$ of $G$ to $\mathcal{O}_{C}$ already defined over the ring of integers of some finite extension of $M$, together with an $\mathcal{O}_{K}$-algebra isomorphism $\mathcal{O}_{s} \xlongequal{\cong} \operatorname{End}_{\mathcal{O}_{K}, \mathcal{O}_{C}}\left(F^{\prime}\right)$. A quasi-canonical lift of level $s=0$ is also called canonical.

Proposition 3.4 (compare [Ww1, §1.3]). - Let $F$ be a lift of $G$. Then the following statements are equivalent:
(1) The lift $F$ is canonical, i.e., defined over some finite extension of $M$ and such that $\operatorname{End}_{\mathcal{O}_{K}, \mathcal{O}_{C}}(F)=\operatorname{End}_{\mathcal{O}_{K}, k}(G) \cong \mathcal{O}_{K} \times \mathcal{O}_{K}$.
(2) The lift $F$ is isomorphic to $H_{\mathcal{O}_{M}} \times\left(K / \mathcal{O}_{K}\right)_{\mathcal{O}_{M}}$.

In particular there exists a canonical lift and it is unique up to unique isomorphism.
The formal modulus of a canonical lift $F_{c a n}$ is $q\left(F_{c a n}\right)=0$ and thus independent of the chosen isomorphism $r$.

Proof. - Clearly, the lift $F=H_{\mathcal{O}_{M}} \times\left(K / \mathcal{O}_{K}\right)_{\mathcal{O}_{M}}$ is canonical. To show that any canonical lift is isomorphic to the product, note that the endomorphism ring of a canonical lift contains the images $e_{\text {inf }}$ and $e_{\text {et }}$ of $(1,0) \in \mathcal{O}_{L}$ and $(0,1) \in \mathcal{O}_{L}$. They satisfy $e_{\mathrm{inf}}^{2}=e_{\mathrm{et}}^{2}=1$ and $e_{\mathrm{inf}}+e_{\mathrm{et}}=1$ and hence define a splitting

$$
F \cong \operatorname{Im}\left(e_{\mathrm{inf}}\right) \times \operatorname{Im}\left(e_{\mathrm{et}}\right)
$$

as claimed. Given two canonical lifts, the element $(1,1) \in \mathcal{O}_{L}$ induces a canonical isomorphism. For the last claim simply observe that the split extension is the image of $0 \in H\left(\mathcal{O}_{C}\right)$ under $\delta_{\mathcal{O}_{C}}$ by construction.

Proposition 3.5 (compare [Ww1, §3] and [G, Prop.5.3])
(1) Quasi-canonical liftings $F_{s}$ exist for all levels $s \geq 0$.
(2) Liftings of level $s$ are rational over the ring of integers $\mathcal{O}_{M_{s}}$ of $M_{s}$. Their isomorphism classes are permuted simply transitively under the action of the Galois group

$$
\operatorname{Gal}\left(M_{s} \mid M\right) \cong \mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times} \cong\left(\mathcal{O}_{L} / \mathfrak{p}^{s} \mathcal{O}_{L}\right)^{\times} /\left(\mathcal{O}_{K} / \mathfrak{p}^{s}\right)^{\times}
$$

which has order

$$
\left|\operatorname{Gal}\left(M_{s} \mid M\right)\right|=\left\{\begin{array}{cl}
q^{s}\left(1-\frac{1}{q}\right) & : \quad s \geq 1 \\
1 & : \quad s=0
\end{array}\right.
$$

In particular $M_{s}$ is the smallest extension of $M$ over which a quasi-canonical lift can be defined.
(3) The formal modulus $q\left(F_{s}\right) \in H\left(\mathcal{O}_{M_{s}}\right)=H\left(\mathcal{O}_{C}\right)$ of a quasi-canonical lift of level $s$ is a uniformizing element of $\mathcal{O}_{M_{s}}$. In particular, for $s \geq 1$ the $\mathcal{O}_{K}$-modules $F_{s}$ and $F_{\text {can }}$ are not isomorphic over $\mathcal{O}_{M_{s}} / \mathfrak{m}_{M_{s}}^{2}$.
Proof. - For the first point recall that it follows from Lubin-Tate theory that $H\left(\mathcal{O}_{C}\right)_{\text {torsion }} \cong K / \mathcal{O}_{K}$ as $\mathcal{O}_{K}$-modules. Thus there are elements $q_{s} \in H\left(\mathcal{O}_{C}\right)$ with annihilator $\mathfrak{p}^{s}$ for any given $s \geq 0$. This implies the existence of a lift $F_{s} / \mathcal{O}_{C}$ with formal modulus $q_{s}$. By Corollary 2.4 the endomorphism ring of $F_{s}$ is isomorphic to $\mathcal{O}_{s}$. If $s=0$ then $F_{\text {can }}=H \times K / \mathcal{O}_{K}$ is a canonical lift and it is clearly defined over $M$. If $s \geq 1$ then the stabilizer of the formal modulus $q_{s}$, i.e., $1+\operatorname{Ann}\left(q_{s}\right)$, equals $1+\mathfrak{p}^{s} \subset \mathcal{O}_{K}^{\times}$. Thus again by Lubin-Tate theory its isomorphism class is stable under the Galois group $\operatorname{Gal}\left(M^{\text {sep }} \mid M_{s}\right)$ since the identification of $\mathcal{D}_{F_{0}}\left(\mathcal{O}_{C}\right)$ with $H\left(\mathcal{O}_{C}\right)$ is compatible with the action of $\operatorname{Gal}\left(M^{\text {sep }} \mid M\right)$. Since deformations have no non-trivial automorphisms, this induces a Galois action on the chosen lift $F_{s} / \mathcal{O}_{C}$ itself. It follows that $F_{s}$ descends to a formal $\mathcal{O}_{K}$-module over $\mathcal{O}_{M_{s}}=\mathcal{O}_{C} \cap M_{s}$.

For the second point note that the first isomorphism follows from Remark 3.2. One checks easily that the natural map

$$
\mathcal{O}_{L}^{\times} / \mathcal{O}_{s}^{\times} \longrightarrow\left(\mathcal{O}_{L} / \mathfrak{p}^{s} \mathcal{O}_{L}\right)^{\times} /\left(\mathcal{O}_{K} / \mathfrak{p}^{s}\right)^{\times}
$$

is an isomorphism. For $s \geq 1$ it follows from Lubin-Tate theory that

$$
\left|\mathcal{O}_{K}^{\times} / 1+\mathfrak{p}^{s}\right|=N(\mathfrak{p})^{s-1}(N(\mathfrak{p})-1)=\left|\operatorname{Gal}\left(M_{s} \mid M\right)\right|
$$

as claimed .
The last point also follows from Lubin-Tate theory (see [Me1]), for one knows that $N_{M_{s} \mid M}\left(-q_{s}\right)=\pi$ and hence

$$
v_{M_{s} \mid M}\left(q_{s}\right)=\frac{1}{\left[M_{s}: M\right]} v_{M}\left(N_{M_{s} \mid M}\left(q_{s}\right)\right)=\frac{1}{\left[M_{s}: M\right]}
$$

as claimed. Therefore $q_{s} \in \mathfrak{m}_{M_{s}} \backslash \mathfrak{m}_{M_{s}}^{2}$ for $s \geq 1$. But the canonical lift has formal modulus $q_{\text {can }}=0 \in \mathfrak{m}_{M_{s}}^{2}$. It follows that $q_{s} \not \equiv q_{\text {can }} \bmod \mathfrak{m}_{M_{s}}^{2}$.

## Remark 3.6

(i) The degree formula in the proposition can be written in a uniform way as

$$
\left|\operatorname{Gal}\left(M_{s} \mid M\right)\right|=N\left(\mathfrak{p}^{s}\right) \prod_{\mathfrak{l} \mid \mathfrak{p}^{s}}\left(1-\left(\frac{L}{\mathfrak{l}}\right) \frac{1}{N(\mathfrak{l})}\right)
$$

where one formally sets

$$
\left(\frac{L}{\mathfrak{l}}\right)=+1,-1,0
$$

according as $\mathfrak{l}=\mathfrak{p}$ is split (our case), inert or ramified (the cases treated in [Ww1]) in the extension $L \mid K$.
(ii) Let $E_{0}$ be an ordinary elliptic curve over $\overline{\mathbb{F}}_{p}$. Then one knows that its endomorphism ring is isomorphic to some order $\mathcal{O} \subset L$ in some imaginary quadratic field $L$. Let $c_{0} \in \mathbb{Z}$ be the conductor of $\mathcal{O}$. It is known that $p$ does not divide $c_{0}$. Set $c_{s}=p^{s} c_{0}$ and $\mathcal{O}_{s}=\mathbb{Z}+p^{s} \mathcal{O}$. Let $M_{s} \mid L$ be the ring class field of the order $\mathcal{O}_{s}$. For example if $c_{0}=1$ and $s=0$ then $M_{s}=M$ is the Hilbert class field of $L$, i.e., the maximal unramified abelian extension of $L$. In this situation one has Deuring's lifting theorem (compare [ $\mathbf{L}$, chap. $13, \S 4, \S 5]$ ). It guarantees the existence of an elliptic curve $E_{s}$ over $M_{s}$ with complex multiplication by $\mathcal{O}_{s}$ and such that the reduction of $E_{s}$ at some prime of degree one over $p$ is isomorphic to $E_{0}$ (same notational conflict as in the local case). The $j$-invariants of the different curves $E_{s}$ are permuted simply transitively by the Galois group $\operatorname{Gal}\left(M_{s} \mid M\right)$. By the well known formula for the class numbers of orders in imaginary quadratic fields (see [S, exercise 4.12]) the Galois group has order

$$
\left|\operatorname{Gal}\left(M_{s} \mid M\right)\right|=\frac{h\left(\mathcal{O}_{s}\right)}{h(\mathcal{O})}=\frac{\left|\mathcal{O}_{s}^{\times}\right|}{\left|\mathcal{O}^{\times}\right|} \cdot \frac{c_{s}}{c_{0}} \prod_{l \left\lvert\, \frac{c_{s}}{c_{0}}\right.}\left(1-\left(\frac{L}{l}\right) \frac{1}{l}\right) .
$$

where the symbol $\left(\frac{L}{l}\right)$ is defined as in (i). The extra factor $\frac{\left|\mathcal{O}_{s}^{\times}\right|}{\left|\mathcal{O}^{\times}\right|}$is due to the presence of nontrivial automorphisms in this situation. It is trivial for $L \neq \mathbb{Q}(i), \mathbb{Q}\left(e^{\frac{2 \pi i}{3}}\right)$. This statement of a global nature is thus completely analogous to the local statement of Proposition 3.5.

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# 10. LIFTING ENDOMORPHISMS OF FORMAL $\mathcal{O}_{K}$-MODULES 

by

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#### Abstract

We present Keating's results on lifts of endomorphisms of formal $\mathcal{O}_{K^{-}}$ modules over a power series ring. Let $k$ be a separably closed field of characteristic $p$. Let $K$ a complete discretely valued field of characteristic $p$ with finite residue field and $\mathcal{O}_{K}$ its ring of integers. Let $F$ be a formal $\mathcal{O}_{K}$-module over $k[[t]]$ with generic fiber of height $h-1$ and special fiber of height $h$. We compute the endomorphism ring of the reduction of $F$ to $k[[t]] /\left(t^{n+1}\right)$.


Résumé (Relèvements des endomorphismes de $\mathcal{O}_{K}$-modules formels). - On présente les résultats de Keating sur les relèvements des endomorphismes de $\mathcal{O}_{K}$-modules formels sur un anneau de séries formelles. Soit $k$ un corps de charactéristique $p$ séparablement clos. Soit $\mathcal{O}_{K}$ un anneau de valuation discrète complète de charactéristique $p$ à corps résiduel fini et soit $K$ son corps des fractions. Soit $F$ un $\mathcal{O}_{K}$-module formel sur $k[[t]]$ à fibre générique de hauteur $h-1$ et fibre speciale de hauteur $h$. On calcule l'anneau des endomorphismes de la réduction de $F$ à $k[[t]] /\left(t^{n+1}\right)$.

The following is an exposition of results in [K2].

## 1. Results

Let $K$ be a complete discretely valued field of arbitrary characteristic with finite residue field $\mathbb{F}_{q}$, where $q=p^{f}$ for some prime $p$. Denote by $\mathcal{O}_{K}$ the ring of integers in $K$ and let $\pi$ be a uniformizer.

Let $1<h<\infty$ and let $\left(F_{0}, \gamma_{0}\right)$ be a formal $\mathcal{O}_{K}$-module of height $h$ over a field $k$ of characteristic $p$ with $\mathcal{O}_{K}$-algebra structure $i: \mathcal{O}_{K} \rightarrow k$. The discrete valuation ring $R=k[[t]]$ has a canonical $\mathcal{O}_{K}$-algebra structure given by $\mathcal{O}_{K} \rightarrow k \hookrightarrow k[[t]]=R$. Let $(F, \gamma)$ be a deformation of $F_{0}$ of height $g=h-1$ over $R$, that is a formal $\mathcal{O}_{K}$-module $F$ over $R$ with $F \equiv F_{0}(\bmod (t))$. For $a \in \mathcal{O}_{K}$ let $[a]_{F}=\gamma(a) \in \operatorname{End}_{R}(F)$. Let
$R_{n}=R /\left(t^{n+1}\right), M_{n}=\left(t^{n} R_{n}\right)$, and $F_{n}=F \otimes_{R} R_{n}$. As $F$ is an $\mathcal{O}_{K}$-module of height $g<h$, Lemma 2.1 of $[\mathbf{V Z}]$ shows that

$$
\begin{equation*}
[\pi]_{F}(X)=a_{0} X^{q^{g}}+\cdots \in R\left[\left[X^{q^{g}}\right]\right] \tag{1.1}
\end{equation*}
$$

with $a_{0} \in R \backslash\{0\}$ and $v_{t}\left(a_{0}\right) \geq 1$.
The aim is to compute the endomorphism ring

$$
H_{n}=\operatorname{End}_{R_{n}}\left(F_{n}\right)=\operatorname{End}_{R}\left(F_{n}\right)
$$

for every $n$. Lemma 2.6 of $[\mathbf{V Z}]$ implies that the reduction maps $R_{n} \rightarrow R_{n-1}$ induce injections $H_{n} \hookrightarrow H_{n-1}$. By [Ww1, Theorem 1.1] we have $H_{0}=\operatorname{End}_{k}\left(F_{0}\right)=\mathcal{O}_{D}$, where $\mathcal{O}_{D}$ is the maximal order in a division algebra $D$ of degree $h^{2}$ and invariant $\frac{1}{h}$ over $K$. Hence the rings $H_{n}$ can be identified with $\mathcal{O}_{K^{-}}$-subalgebras of $\mathcal{O}_{D}$. Obviously $\mathcal{O}_{K} \subseteq \bigcap_{n} H_{n}$. Let $\pi_{D}$ be a uniformizer of $\mathcal{O}_{D}$.

For $m \geq 0$ we define

$$
a(g m)=\frac{\left(q^{h}-1\right)\left(q^{g m}-1\right)}{\left(q^{g}-1\right)(q-1)}
$$

Theorem 1.1. - Let $F_{0}, F$, and $a_{0}$ be as above, $v_{t}\left(a_{0}\right)=1, k$ separably closed, and

$$
f_{0} \in\left(\mathcal{O}_{K}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{K}+\pi_{D}^{l+1} \mathcal{O}_{D}\right) \subseteq \mathcal{O}_{D}
$$

for some $l \geq 0$. Write $l=h m+b$ with integers $m, b$ and $0 \leq b<h$. Then $f_{0} \in$ $H_{n-1} \backslash H_{n}$ for

$$
n=a(g m)+q^{g m} \frac{q^{b}-1}{q-1}+1
$$

Using this result, we can easily calculate $H_{n}$ :
Theorem 1.2. - In the setting of Theorem 1.1 we have

$$
H_{n}=\mathcal{O}_{K}+\pi_{D}^{j(n)} \mathcal{O}_{D}
$$

with $j(n)=h m+b$, where $m$ and $b$ are the uniquely defined integers with $0 \leq b<h$ and

$$
\begin{aligned}
a(g m)-q^{g m}+1 \leq n<a(g m)+1 & \text { if } \quad b=0 \\
a(g m)+q^{g m} \frac{q^{b-1}-1}{q-1}+1 \leq n<a(g m)+q^{g m} \frac{q^{b}-1}{q-1}+1 & \text { if } \quad 0<b<h .
\end{aligned}
$$

Keating ([K2, Thm. 3.4]) also calculates $H_{n}$ without making the assumptions $v_{t}\left(a_{0}\right)=1$ and $k$ separably closed.

To prove Theorem 1.1 we need the following two propositions.
Proposition 1.3. - Let $f_{0}$ be as in Theorem 1.1 with $0 \leq l \leq h$. Then $f_{0} \in H_{n-1} \backslash H_{n}$ where $n=\frac{q^{l}-1}{q-1}+1$.

Proposition 1.4. - Let $f_{0} \in \mathcal{O}_{K}+\pi_{D} \mathcal{O}_{D}$ and $n>1$ such that $f_{0} \in H_{n-1} \backslash H_{n}$. Then $[\pi]_{F_{0}} \circ f_{0} \in H_{n^{\prime}-1} \backslash H_{n^{\prime}}$ where $n^{\prime}=q^{g} n+\frac{q^{g}-1}{q-1}+1$.

The first proposition says that the theorem holds for $l \leq h$. The second calculates the maximal lifting of $[\pi]_{F_{0}} \circ f_{0}$ given the maximal lifting of $f_{0}$. As the elements of $\mathcal{O}_{K}$ lift to all levels, an inductive argument shows that the two propositions imply Theorem 1.1.

## 2. Proof of the Lifting Theorem

We use the notation of the preceding section. Without further mention we assume that the constant coefficients of all power series in this section are 0 .

Let $f_{n+1} \in R_{n+1}[[X]]$ be a lift of $f_{n} \in H_{n}$. We recall the definition of the associated symmetric 2-cocycle with coefficients in $\left(t^{n+1}\right) k[[t]] /\left(t^{n+2}\right)$ from [VZ, Prop. 3.2]:

$$
\begin{aligned}
\Delta(X, Y) & =f_{n+1}\left(X+_{F_{n+1}} Y\right)-F_{n+1} f_{n+1}(X)-F_{n+1} f_{n+1}(Y) \\
\delta_{a}(X) & =f_{n+1} \circ[a]_{F_{n+1}}(X)-F_{n+1}[a]_{F_{n+1}} \circ f_{n+1}(X)
\end{aligned} \quad(a \in A)
$$

The cocycle vanishes if and only if $f_{n+1} \in H_{n+1}=\operatorname{End}_{R_{n+1}}(F)$.
Corollary 2.1. - Let $f_{n}(X) \in R_{n}[[X]]$ be a lift of $f_{0}(X) \in H_{0}$. Then $f_{n} \in H_{n}$ if and only if $f_{n}$ commutes with $[\pi]_{F_{n}}$.

Proof. - This follows by induction on $i \in[0, n]$ from [VZ, Prop. 3.2, 2. and Lemma 3.1].

Lemma 2.2. - Let $f_{n} \in H_{n}$ with $f_{n}(X) \in R_{n}\left[\left[X^{q^{r}}\right]\right]$ of the form $f_{n}(X)=b_{0} X^{q^{r}}+\cdots$ for some $r>0$ and $b_{0} \in R_{n} \backslash\{0\}$.

1. There exists a lift $f_{n+1} \in H_{n+1}$ of $f_{n}$.
2. If the degree of the leading term of $\delta_{\pi}$ is greater than $q^{r+g}$, the leading term of $f_{n+1}$ has degree $q^{r}$.
3. Otherwise the degree of the leading term of $\delta_{\pi}$ is $q^{r+g}$ and the leading term of $f_{n+1}$ has degree $q^{r-1}$.

Proof. - Let $f_{n+1}^{\prime} \in R_{n+1}\left[\left[X^{q^{r}}\right]\right]$ be an arbitrary lift of $f_{n}$ and let $\left(\Delta ;\left\{\delta_{a}\right\}\right)$ be the corresponding cocycle. As $[\pi]_{F_{n+1}}(X) \in R_{n+1}\left[\left[X^{q^{g}}\right]\right]$, the degree of the leading term of $\delta_{\pi}$ is at least $q^{r+g}$. The assumption $r>0$ together with [VZ, Cor. 3.4 1.] implies that $\delta_{\pi}(X) \in M_{n+1}\left[\left[X^{q^{h}}\right]\right]$. Hence by [VZ, Cor. 3.4, 2. and Prop. 3.2, 3.] we get a lift $f_{n+1}=f_{n+1}^{\prime}-F_{n+1} \varepsilon$. If the degree of the leading term of $\delta_{\pi}$ is $q^{r+g}$, the leading term of $\varepsilon=d \circ p^{-1}$ has degree $q^{r+g-h}=q^{r-1}$. Otherwise the leading term of $\varepsilon$ has degree greater than $q^{r-1}$. Thus the degree of the leading term of $f_{n+1}$ is greater than $q^{r-1}$. As it is a lift of $f_{n}$, its degree is at most $q^{r}$. But the degree has to be a power of $q$, as $f_{n+1} \in H_{n+1}=\operatorname{End}_{R_{n+1}}\left(F_{n+1}\right)$ (compare [VZ, Lemma 2.1]). Hence it has to be $q^{r}$.

Let $F, F_{0}, a_{0}$ as in Theorem 1.1. Let $f_{0} \in H_{0}$ and let $f_{n-1} \in H_{n-1}$ be a lift, i.e., $f_{n-1} \equiv f_{0}(\bmod (t))$. We can write $f_{n-1}(X)=b_{0} X^{q^{r}}+\cdots$ for some $r \geq 0$ and $b_{0} \in R_{n-1} \backslash\{0\}$. Let $m=v_{t}\left(b_{0}\right)$. As $b_{0} \neq 0$ in $R_{n-1}$, we have $m<n$.

Lemma 2.3. - Suppose that in the above situation $m+q^{r}<q^{g} m+1$. Then

1. $m+q^{r} \geq n$.
2. If $m+q^{r}>n$, then $f_{n-1}$ lifts to $f_{n}^{\prime} \in H_{n}$ of the form $f_{n}^{\prime}(X)=b_{0}^{\prime} X^{q^{r}}+\cdots$ with $v_{t}\left(b_{0}^{\prime}\right)=m$.
3. If $m+q^{r}=n$ and $r>0$, then $f_{n-1}$ lifts to $f_{n}^{\prime} \in H_{n}$ of the form $f_{n}^{\prime}(X)=$ $b_{0}^{\prime} X^{q^{r-1}}+\cdots$ with $v_{t}\left(b_{0}^{\prime}\right)=n=m+q^{r}$.
4. A lift of $f_{n-1}$ to $H_{n}$ again satisfies the assumption of the lemma.
5. If $m+q^{r}=n$ and $r=0$, then $f_{n-1}$ does not lift to an element of $H_{n}$.
6. $f_{n-1}$ lifts to $H_{n^{\prime}-1}$ with $n^{\prime}=m+\frac{q^{r+1}-1}{q-1}$ but not to $H_{n^{\prime}}$.

Proof. - We assumed $v_{t}\left(a_{0}\right)=1$, where $a_{0}$ is as in (1.1). By the assumption of the lemma the valuation of the leading coefficient of $\delta_{\pi}$ is

$$
v_{t}\left(b_{0} a_{0}^{q^{r}}-a_{0} b_{0}^{q^{g}}\right)=m+q^{r} .
$$

As $f_{n-1} \in H_{n-1}$, this coefficient is in $\left(t^{n}\right)$ and 1. follows. Since $f_{n-1} \in H_{n-1}$ has leading term of degree $q^{r}$, it has to be a power series in $X^{q^{r}}$. Let $f_{n} \in R_{n}\left[\left[X^{q^{r}}\right]\right]$ be a lifting with

$$
f_{n}(X)=\tilde{b}_{0} X^{q^{r}}+\cdots
$$

For the corresponding 2-cocycle we have

$$
\delta_{\pi}(X)=\left(\tilde{b}_{0} a_{0}^{q^{r}}-a_{0} \tilde{b}_{0}^{q^{g}}\right) X^{q^{r+g}}+\cdots
$$

If $m+q^{r}>n$, the first term vanishes modulo $\left(t^{n+1}\right)$, so the degree of $\delta_{\pi}$ is greater than $q^{r+g}$. Besides, $r$ has to be positive, thus 2. follows from Lemma 2.2, 2. If $m+q^{r}=n$, the leading term of $\delta_{\pi}$ has degree $q^{r+g}$, hence 3. follows from Lemma 2.2,3. The values of $m$ and $r$ for a lift of $f_{n-1}$ to $H_{n}$ are either the same as for $f_{n-1}$ or they change to $n$ and $r-1$. In both cases the assumption of the lemma is satisfied. This shows 4 . for the case $r>0$. If $r=0$, the statement is trivial. To show 5 ., we assume that there exists a lift. By 4. it satisfies the assumption of the lemma with $r=0$ and $m=n-1$. Thus 1 . implies $n-1+1 \geq n+1$ which is a contradiction. The last assertion follows by applying $2 .-4$. until the assumption of 5 . holds.

Proof of Proposition 1.4. - We write $f_{0}=a+f_{0}^{\prime}$ with $a \in \mathcal{O}_{K}$ and $f_{0}^{\prime} \in \pi_{D} \mathcal{O}_{D}$. Since $\mathcal{O}_{K} \subseteq H_{n}$ for every $n$, we may assume $f_{0}=f_{0}^{\prime} \in \pi_{D} \mathcal{O}_{D}$. We assumed that $f_{0}$ lifts to $f_{n-1} \in H_{n-1}$ with $f_{n-1}(X)=b_{0} X^{q^{r}}+\cdots$ for some $b_{0} \in R_{n-1} \backslash\{0\}$ but not to $H_{n}$. Lemma 2.2 implies $r=0$. We have $f_{n-1} \equiv f_{0}(\bmod (t))$ and $f_{0} \in \pi_{D} \mathcal{O}_{D}$. Hence $m=v_{t}\left(b_{0}\right)>0$ and the assumption of Lemma 2.3 is satisfied. As $f_{n-1}$ does not lift to $H_{n}$, the lemma shows that $v_{t}\left(b_{0}\right)=n-1$.

We lift $f_{n-1}$ arbitrarily to $f(X) \in R[[X]]$. This lift is unique modulo $\left(t^{n}\right)$, so $[\pi]_{F} \circ f$ is unique modulo $\left(t^{q^{g} n+1}\right)$. Here we use that $[\pi]_{F}(X)=a_{0} X^{q^{g}}+\cdots$ with $v_{t}\left(a_{0}\right)=1$. We next show that the reduction $\phi_{q^{g} n}$ of $[\pi]_{F} \circ f$ modulo $\left(t^{q^{g} n+1}\right)$ is in $H_{q^{g}}$. By Lemma 2.1, 2. it suffices to verify that

$$
[\pi]_{F} \circ\left([\pi]_{F} \circ f\right) \equiv\left([\pi]_{F} \circ f\right) \circ[\pi]_{F} \quad\left(\bmod \left(t^{q^{g} n+1}\right)\right)
$$

which follows from $[\pi]_{F} \circ f \equiv f \circ[\pi]_{F}\left(\bmod \left(t^{n}\right)\right)$. Now we determine the maximal lifting of $\phi_{q^{g}}$. We have

$$
\begin{aligned}
\phi_{q^{g} n} & \equiv[\pi]_{F} \circ f(X) \\
& \equiv a_{0} b_{0}^{q^{g}} X^{q^{g}}+\cdots \quad\left(\bmod \left(t^{q^{g} n+1}\right)\right)
\end{aligned}
$$

and $v_{t}\left(a_{0} b_{0}^{q^{g}}\right)=1+q^{g}(n-1)$. The assumption of Lemma 2.3 applied to $\phi_{q^{g} n}$ reads

$$
\left(1+q^{g}(n-1)\right)+q^{g}<q^{g}\left(1+q^{g}(n-1)\right)+1 .
$$

It is satisfied because $n>1$ and $g \geq 1$. We get that $\phi_{q^{g} n}$ lifts to $H_{n^{\prime}-1}$ but not to $H_{n^{\prime}}$ where $n^{\prime}=q^{g} n+\frac{q^{g}-1}{q-1}+1$.

For the proof of Proposition 1.3 we need the following lemma.
Lemma 2.4. - Let $f_{n-1} \in H_{n-1}$ with $f_{n-1}(X)=b_{0} X^{q^{r}}+\cdots$ and $m=v_{t}\left(b_{0}\right)$ as before. Assume $m+q^{r}>q^{g} m+1=n$. Then $f_{n-1}$ lifts to $f_{n}^{\prime} \in H_{n}$ of the form $f_{n}^{\prime}(X)=b_{0}^{\prime} X^{q^{r-1}}+\cdots$ with $v_{t}\left(b_{0}^{\prime}\right)=q^{g} m+1$.
Proof. - We lift $f_{n-1}$ arbitrarily to $f_{n}(X) \in R_{n}\left[\left[X^{q^{r}}\right]\right]$ with $f_{n}(X)=\tilde{b}_{0} X^{q^{r}}+\cdots$. As before let $\left(\Delta,\left\{\delta_{a}\right\}\right)$ be the corresponding cocycle with coefficients in $M_{n}$. We have

$$
\begin{aligned}
\delta_{\pi}(X) & =f_{n} \circ[\pi]_{F_{n}}(X)-F_{n}[\pi]_{F_{n}} \circ f_{n}(X) \\
& =\left(\tilde{b}_{0} a_{0}^{q^{r}}-a_{0} \tilde{b}_{0}^{q^{g}}\right) X^{q^{r+g}}+\cdots .
\end{aligned}
$$

The assumptions imply $v_{t}\left(\tilde{b}_{0} a_{0}^{q^{r}}-a_{0} \tilde{b}_{0}^{q^{g}}\right)=q^{g} m+1=n$. Therefore the first nonvanishing term of $\delta_{\pi}$ has degree $q^{r+g}$. As $m<n$, the assumption $m+q^{r}>n$ implies $r>0$, and by Lemma 2.2, 3., $f_{n-1}$ lifts to $f_{n}^{\prime} \in H_{n}$ with leading term of degree $q^{r-1}$.
Proof of Proposition 1.3
Case 1: $l=0$. In this case $f_{0} \in \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\pi_{D} \mathcal{O}_{D}\right)$ has the form $f_{0}(X)=b_{0} X+\cdots$ with $b_{0} \in \mathbb{F}_{q^{h}} \backslash \mathbb{F}_{q}$. Let $f_{1}(X) \in R_{1}[[X]]$ with $f_{1}(X)=b_{0} X+\cdots$ be an arbitrary lift of $f_{0}$. We have to show that $f_{1} \notin H_{1}$. For the corresponding cocycle with coefficients in $M_{1}$ we have $\delta_{\pi}(X)=\left(b_{0} a_{0}-a_{0} b_{0}^{q^{g}}\right) X^{q^{g}}+\cdots$. Since $b_{0} \in \mathbb{F}_{q^{h}} \backslash \mathbb{F}_{q}$, we have $v_{t}\left(b_{0}-b_{0}^{q^{g}}\right)=0$. Thus $b_{0} a_{0}-a_{0} b_{0}^{q^{g}}$ is nonzero in $R_{1}$ and $\delta_{\pi}(X) \neq 0$, which shows $f_{1} \notin H_{1}$.

Case 2: $0<l<h$. Here

$$
\left(\mathcal{O}_{K}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{K}+\pi_{D}^{l+1} \mathcal{O}_{D}\right) \subseteq \mathcal{O}_{K}+\left(\pi_{D}^{l} \mathcal{O}_{D} \backslash \pi_{D}^{l+1} \mathcal{O}_{D}\right)
$$

As elements of $\mathcal{O}_{K}$ lift to all levels, it is enough to consider $f_{0} \in \pi_{D}^{l} \mathcal{O}_{D} \backslash \pi_{D}^{l+1} \mathcal{O}_{D}$. Then $f_{0}$ is of the form $f_{0}(X)=b_{0} X^{q^{l}}+\cdots$ with $b_{0} \in k^{\times}$. As $m=0$, Lemma 2.4 shows that $f_{0}$ lifts to $f_{1}^{\prime} \in H_{1}$ of the form $f_{1}^{\prime}(X)=b_{0}^{\prime} X^{q^{l-1}}+\cdots$ with $v_{t}\left(b_{0}^{\prime}\right)=1$. For $f_{1}^{\prime}$ the assumption of Lemma 2.3 is satisfied, so $f_{0}$ lifts to $H_{n-1}$ but not to $H_{n}$ with $n=\frac{q^{l}-1}{q-1}+1$.

Case 3: $l=h$. Here

$$
\left(\mathcal{O}_{K}+\pi_{D}^{h} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{K}+\pi_{D}^{h+1} \mathcal{O}_{D}\right) \subseteq \mathcal{O}_{K}+\pi_{D}^{h}\left(\mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\pi_{D} \mathcal{O}_{D}\right)\right)
$$

Similarly to the second case it suffices to consider $f_{0} \in \pi_{D}^{h}\left(\mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\pi_{D} \mathcal{O}_{D}\right)\right)$, that is $f_{0}=\pi g_{0}$ for some $g_{0} \in \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\pi_{D} \mathcal{O}_{D}\right)$. Then $g_{0}(X)=b_{0} X+\cdots$ with $b_{0} \in \mathbb{F}_{q^{h}} \backslash \mathbb{F}_{q}$. Let $g(X) \in R[[X]]$ be an arbitrary lift of $g_{0}$. From

$$
[\pi]_{F} \circ g \equiv g \circ[\pi]_{F} \quad(\bmod (t))
$$

we get

$$
[\pi]_{F} \circ\left([\pi]_{F} \circ g\right) \equiv\left([\pi]_{F} \circ g\right) \circ[\pi]_{F} \quad\left(\bmod \left(t^{q^{g}+1}\right)\right) .
$$

Lemma 2.1 shows that

$$
\begin{aligned}
f_{q^{g}}(X) & \equiv[\pi]_{F} \circ g(X) \\
& \equiv a_{0} b_{0}^{q^{g}} X^{q^{g}}+\cdots \quad\left(\bmod \left(t^{q^{g}+1}\right)\right)
\end{aligned}
$$

is in $H_{q^{g}}$. Let $f_{q^{g}+1} \in R_{q^{g}+1}\left[\left[X^{q^{g}}\right]\right]$ with $f_{q^{g}+1}(X)=c_{0} X^{q^{g}}+\cdots$ and $c_{0} \equiv a_{0} b_{0}^{q^{g}}$ $\left(\bmod \left(t^{q^{g}+1}\right)\right)$ be a lift of $f_{q^{g}}$. The corresponding cocycle satisfies

$$
\begin{aligned}
\delta_{\pi}(X) & =f_{q^{g}+1} \circ[\pi]_{F_{q^{g}+1}}(X)-F_{q^{g}+1}[\pi]_{F_{q^{g}+1}} \circ f_{q^{g}+1}(X) \\
& =\left(c_{0} a_{0}^{q^{g}}-a_{0} c_{0}^{q^{g}}\right) X^{q^{2 g}}+\cdots \\
& \equiv a_{0}^{q^{g}+1}\left(b_{0}^{q^{g}}-b_{0}^{q^{2 g}}\right) X^{q^{g}}+\cdots \quad\left(\bmod \left(t^{q^{g}+2}\right)\right) .
\end{aligned}
$$

Since $b_{0} \in \mathbb{F}_{q^{h}} \backslash \mathbb{F}_{q}$, we have $b_{0}^{q^{g}} \neq b_{0}^{q^{2 g}}$ in $\mathbb{F}_{q^{h}}$. Hence $a_{0}^{q^{g}+1}\left(b_{0}^{q^{g}}-b_{0}^{q^{2 g}}\right)$ is nonzero in $M_{q^{g}+1}$ and $\delta_{\pi}$ has leading term of degree $q^{2 g}=q^{r+g}$. So Lemma 2.2,3. shows that $f_{q^{g}}$ lifts to $f_{q^{g}+1}^{\prime} \in H_{q^{g}+1}$ with leading term of degree $q^{g-1}$. As $f_{q^{g}+1}^{\prime}$ satisfies the assumption of Lemma 2.3, it lifts to $H_{n-1}$ but not to $H_{n}$ where $n=\frac{q^{h}-1}{q-1}+1$.

## References

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[^5]
# 11. ENDOMORPHISMS OF QUASI-CANONICAL LIFTS 

by

Inken Vollaard


#### Abstract

We present Keating's result on the locus of deformation of an endomorphism of a quasi-canonical lifting. At the same time, this determines the endomorphism ring of the reduction of quasi-canonical liftings to Artin rings. Résumé (Endomorphismes de relèvements quasi-canoniques). - On donne le résultat de Keating concernant le lieu de déformation d'un endomorphisme d'un relèvement quasi-canonique. En même temps, ceci détermine l'anneau des endomorphismes de la réduction d'un relèvement quasi-canonique à des anneaux artiniens.


In this paper we prove a lifting theorem for endomorphisms of a formal $\mathcal{O}_{K}$-module to a quasi-canonical lift. For the canonical lift, a similar lifting theorem is proved in [Ww1]. This work is due to K. Keating ([K1]).

I thank S. Wewers for helpful comments on this manuscript.

## 1. Notation

Let $K$ be a complete discretely valued field, let $\mathcal{O}_{K}$ be its ring of integers and let $\pi$ be a uniformizing element of $\mathcal{O}_{K}$. We will assume that the residue field of $\mathcal{O}_{K}$ is equal to the field $\mathbb{F}_{q}$ of characteristic $p$. Denote by $k$ an algebraic closure of $\mathbb{F}_{q}$. Let $L$ be a quadratic extension of $K$ and let $A=\widehat{\mathcal{O}}_{L}^{u r}$ be the completion of the maximal unramified extension of $\mathcal{O}_{L}$. Denote by $M$ the quotient field of $A$.

Let $F_{0}$ be a formal $\mathcal{O}_{K}$-module of height 2 over $k$. By [ $\mathbf{W w} \mathbf{1}$ ] Theorem 1.1, the ring of $\mathcal{O}_{K}$-linear endomorphisms $\operatorname{End}_{k} F_{0}$ is isomorphic to the maximal order $\mathcal{O}_{D}$ in a division algebra $D$ of dimension 4 over $K$ and invariant $1 / 2$. We identify $\operatorname{End}_{k} F_{0}$ with $\mathcal{O}_{D}$. Let $F$ be the canonical lift of $F_{0}$ over $A$ with respect to an embedding

$$
\mathcal{O}_{L} \longleftrightarrow \mathcal{O}_{D}
$$

2000 Mathematics Subject Classification. - 14L05, 11G07, 11S31, 14K07.
Key words and phrases. - Quasi-canonical lift, endomorphism ring, formal module.

We consider a quasi-canonical lift $F^{\prime}$ of $F_{0}$ of level $s$ ([Ww1, Def. 3.1]). By definition, $\operatorname{End}_{A^{\prime}} F^{\prime}$ is an order $\mathcal{O}_{s}:=\mathcal{O}_{K}+\pi^{s} \mathcal{O}_{L}$ in $\mathcal{O}_{L}$. Note that a quasi-canonical lift of level 0 is a canonical lift and therefore can be defined over $A$. A quasi-canonical lift of level $s \geq 1$ can be defined over a totally ramified Galois extension $M^{\prime} / M$ of degree

$$
\left[M^{\prime}: M\right]= \begin{cases}q^{s}+q^{s-1} & \text { if } L / K \text { is unramified } \\ q^{s} & \text { if } L / K \text { is ramified }\end{cases}
$$

([Ww1, Thm. 3.2]). Denote by $A^{\prime}$ the ring of integers of $M^{\prime}$ and denote by $\pi^{\prime}$ a uniformizing element of $A^{\prime}$. If $s$ is equal to 0 , the ring $A^{\prime}$ is equal to $A$. Let $e_{s}=e\left(A^{\prime} / \mathcal{O}_{K}\right)$ be the ramification index of $A^{\prime}$ over $\mathcal{O}_{K}$, i.e.,

$$
e_{s}= \begin{cases}2 q^{s} & \text { if } L / K \text { is ramified. } \\ q^{s}+q^{s-1} & \text { if } L / K \text { is unramified and } s \neq 0 \\ 1 & \text { if } L / K \text { is unramified and } s=0\end{cases}
$$

By [Ww1, Proposition 4.4 and Proposition 4.6], the endomorphism $[\pi]_{F^{\prime}}$ is given by a power series

$$
\begin{equation*}
[\pi]_{F^{\prime}}=\pi X+\cdots+u X^{q}+\cdots+v X^{q^{2}}+\cdots \in A^{\prime}[[X]] \tag{1.1}
\end{equation*}
$$

with $v_{\pi^{\prime}}(u)=1$ and $v_{\pi^{\prime}}(v)=0$.
Denote by $A_{n}^{\prime}=A^{\prime} /\left(\pi^{\prime}\right)^{n+1}$ the reduction of $A^{\prime}$ modulo $\left(\pi^{\prime}\right)^{n+1}$ and by $F_{n}^{\prime}=$ $F^{\prime} \otimes_{A^{\prime}} A_{n}^{\prime}$ the reduction of $F^{\prime}$ to $A_{n}^{\prime}$. We obtain

$$
\mathcal{O}_{s}=\operatorname{End}_{A^{\prime}} F^{\prime} \subset \cdots \subset \operatorname{End}_{A_{n}^{\prime}} F_{n}^{\prime} \subset \cdots \subset \operatorname{End}_{k} F_{0}=\mathcal{O}_{D}
$$

([VZ, Lem. 2.6]), hence we will consider $\operatorname{End}_{A_{n}^{\prime}} F_{n}^{\prime}$ as a subring of $\operatorname{End}_{k} F_{0}=\mathcal{O}_{D}$. We write End $F_{n}^{\prime}$ instead of $\operatorname{End}_{A_{n}^{\prime}} F_{n}^{\prime}$.

For $n \leq e_{s}$ the ring $A^{\prime} /\left(\pi^{\prime}\right)^{n}$ is of characteristic $p$ and one can define the height of the module $F_{n}^{\prime}\left(\left[\mathbf{V Z}\right.\right.$, Def. 2.2]). By construction, $F_{n}^{\prime}$ is of height 1 if $0<n \leq e_{s}$ and $F_{0}$ is of height 2 . Denote by $a_{i}$ the coefficients of $[\pi]_{F^{\prime}}$. Then $v_{\pi^{\prime}}\left(a_{i}\right) \geq e_{s}$ if $q \nmid i$, and $v_{\pi^{\prime}}\left(a_{i}\right) \geq e_{s}$ if $q \mid i$ and $q^{2} \nmid i$.

## 2. Results

The goal of this paper is to compute the endomorphism rings End $F_{n}^{\prime}$ as subrings of $\mathcal{O}_{D}$. In the case of the canonical lift, these rings are calculated in [Ww1]. Denote by $a(k)$ the rational number

$$
a(k)=\frac{\left(q^{k}-1\right)(q+1)}{q-1}
$$

for every integer $k$. We have $a(0)=0$ and $a(k)=(q+1)\left(\sum_{i=0}^{k-1} q^{i}\right)$ for $k \geq 1$.

Theorem 2.1. - Let $F^{\prime}$ be a quasi-canonical lift of $F_{0}$ of level $s$. Let $l \geq 0$ be an integer and let

$$
f_{0} \in\left(\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}\right)
$$

Then $f_{0}$ lifts to End $F_{n_{l}-1}^{\prime}$ and not to End $F_{n_{l}}^{\prime}$ with

$$
n_{l}=n_{l}(s)=n_{l}(L / K, s)= \begin{cases}a\left(\frac{l}{2}\right)+1 & \text { if } l \leq 2 s \text { and } l \text { even } . \\ a\left(\frac{l-1}{2}\right)+q^{\frac{l-1}{2}}+1 & \text { if } l \leq 2 s \text { and } l \text { odd. } \\ a(s-1)+q^{s-1}+\left(\frac{l+1}{2}-s\right) e_{s}+1 & \text { if } l \geq 2 s-1 .\end{cases}
$$

Remark 2.2. - The rational number $n_{l}$ of the theorem is an integer. Indeed, if $L / K$ is ramified, the ramification index $e_{s}$ is even. If $L / K$ is unramified and $l \geq 2 s$ is even, then

$$
\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}=\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}
$$

Theorem 2.3. - Consider the same situation as in Theorem 2.1. Then End $F_{n}^{\prime}=$ $\mathcal{O}_{s}+\pi_{D}^{j(n)} \mathcal{O}_{D}$ where

$$
j(n)= \begin{cases}2 k & \text { if } \left.n \in] a(k-1)+q^{k-1} ; a(k)\right] \text { for } k<s . \\ 2 k+1 & \text { if } \left.n \in] a(k) ; a(k)+q^{k}\right] \text { for } k<s . \\ k & \text { if } \left.n \in] a(s-1)+q^{s-1}+\left(\frac{k}{2}-s\right) e_{s} ; a(s-1)+q^{s-1}+\left(\frac{k+1}{2}-s\right) e_{s}\right] \\ & \text { for } k \geq 2 s .\end{cases}
$$

Note that the above intervals form a disjoint cover of the set of positive integers. The integer $j(n)$ is uniquely determined unless $L / K$ is unramified and $j(n) \geq 2 s$. In this case we have $\mathcal{O}_{s}+\pi_{D}^{j(n)} \mathcal{O}_{D}=\mathcal{O}_{s}+\pi_{D}^{j(n)+1} \mathcal{O}_{D}$ for every even $j(n)$.

Proof. - This theorem follows from Theorem 2.1.
Remark 2.4. - If $F^{\prime}$ is the canonical lift of $F_{0}$, i.e., if $s=0$, Theorem 2.1 and Theorem 2.3 have already been proved in [Ww1] Theorem 1.4. We obtain in this case

$$
\text { End } F_{n}=\mathcal{O}_{L}+\pi_{L}^{n} \mathcal{O}_{D}
$$

and

$$
n_{l}(0)= \begin{cases}l+1 & \text { if } L / K \text { is ramified } \\ \frac{l+1}{2} & \text { if } L / K \text { is unramified }\end{cases}
$$

## 3. Proofs

We will assume in the following that $s$ is greater or equal than 1 . We will split the proof of Theorem 2.1 into two propositions similar to the proof of Theorem 1.1 in $[\mathbf{V i}]$.

Proposition 3.1. - Let $l \leq 2 s+1$ and let $s \geq 1$. Let

$$
f_{0} \in\left(\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}\right)
$$

Then $f_{0}$ lifts to End $F_{n_{l}-1}^{\prime} \backslash$ End $F_{n_{l}}^{\prime}$ with

$$
n_{l}= \begin{cases}a\left(\frac{l}{2}\right)+1 & \text { if } l \leq 2 s \text { and } l \text { even } . \\ a\left(\frac{l-1}{2}\right)+q^{\frac{l-1}{2}}+1 & \text { if } l \leq 2 s \text { and } l \text { odd. } \\ a(s-1)+q^{s-1}+e_{s}+1 & \text { if } l=2 s+1\end{cases}
$$

Proposition 3.2. - Let $s \geq 1$ and let $f_{0} \in \operatorname{End} F_{n-1}^{\prime} \backslash$ End $F_{n}^{\prime}$ with $n \geq \frac{e_{s}-1}{q-1}$. Then $[\pi] \circ f_{0}$ lifts to End $F_{n^{\prime}-1}^{\prime} \backslash$ End $F_{n^{\prime}}^{\prime}$ with $n^{\prime}=n+e_{s}$.

Proof of Theorem 2.1. - Theorem 2.1 follows by induction from Proposition 3.1 and Proposition 3.2. Let $l>2 s+1$ and let $f_{0} \in\left(\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}\right)$. Write $f_{0}=c+[\pi]_{F_{0}} \circ g_{0}$ with $c \in \mathcal{O}_{s}$ and $g_{0} \in \pi_{D}^{l-2} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l-1} \mathcal{O}_{D}\right)$. By induction $g_{0}$ lifts to End $F_{n_{l-2}-1}^{\prime} \backslash$ End $F_{n_{l-2}}^{\prime}$ with

$$
\begin{aligned}
n_{l-2} & =a(s-1)+q^{s-1}+\left(\frac{l-1}{2}-s\right) e_{s}+1 \\
& =\frac{2 q^{s}-2}{q-1}+\left(\frac{l-1}{2}-s\right) e_{s} \geq \frac{e_{s}-1}{q-1}
\end{aligned}
$$

By Proposition 3.2 the endomorphism $[\pi]_{F_{0}} \circ g_{0}$, hence $f_{0}$, lifts to End $F_{n^{\prime}-1}^{\prime} \backslash$ End $F_{n^{\prime}}^{\prime}$ with $n^{\prime}=n_{l-2}+e_{s}=n_{l}$.

Remark 3.3. - We now split the proof of Proposition 3.1 into two cases. As we will see below, we can use the results of $[\mathbf{V i}]$ in the case $n_{l}+1 \leq e_{s}$. Note that $n_{l}$ is a strictly increasing sequence.

An easy computation shows that there exists an integer $l_{0}$ such that $n_{l_{0}}+1 \leq e_{s}<$ $n_{l_{0}+1}$. We obtain

- $l_{0}=2 s$ if $L / K$ is ramified and $q \geq 3$.
$-l_{0}=2 s-1$ if $\left\{\begin{array}{l}L / K \text { is unramified and } q \geq 3 . \\ L / K \text { is unramified, } q=2 \text { and } s=1 . \\ L / K \text { is ramified and } q=2 .\end{array}\right.$
$-l_{0}=2 s-2$ if $L / K$ is unramified, $q=2$ and $s \neq 1$.
Proof of Proposition 3.1 in the case of $n_{l}+1 \leq e_{s}$. -- Since $A^{\prime} / \widehat{\mathcal{O}}_{K}$ is a totally ramified extension of ramification index $e_{s}$, we obtain for $n \leq e_{s}$ an isomorphism of $\mathcal{O}_{K^{-}}$-algebras

$$
\begin{aligned}
A^{\prime} /\left(\pi^{\prime}\right)^{n} & \cong\left(\widehat{\mathcal{O}}_{K} /(\pi)\right)\left[\pi^{\prime}\right] /\left(\pi^{\prime}\right)^{n} \\
& \cong k[t] /(t)^{n} .
\end{aligned}
$$

Let $f_{0} \in\left(\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}\right)$ with $n_{l}+1 \leq e_{s}$, i.e., with $l \leq l_{0}$ (Rem. 3.3). Then $F_{n_{l}}^{\prime}$ is a lift of $F_{0}$ of height 1 over $k[t] /(t)^{n_{l}+1}$ and we will prove the proposition by using the results of $[\mathbf{V i}]$.

We have $\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}=\mathcal{O}_{K}+\pi_{D}^{l} \mathcal{O}_{D}$ for $l \leq 2 s$ and by an easy computation $\mathcal{O}_{s}+\pi_{D}^{2 s+1} \mathcal{O}_{D}=\mathcal{O}_{K}+\pi_{D}^{2 s+1} \mathcal{O}_{D}$ if $L / K$ is ramified. Hence [ $\mathbf{V i}$ ] Theorem 1.1, shows that $f_{0}$ lifts to End $F_{n_{l}-1}^{\prime} \backslash$ End $F_{n_{l}}^{\prime}$. This proves the proposition in this case.
3.1. Let $f_{0}$ be an element of End $F_{n-1}^{\prime}$. By $f_{n-1} \in A_{n-1}^{\prime}[[X]]$ we always denote the unique lift of $f_{0}$ as an endomorphism of $F_{n-1}^{\prime}$. Let $f \in A^{\prime}[[X]]$ be a lift of $f_{n-1}$ as a power series without constant coefficient. As we are interested in endomorphisms of formal groups, we make the general assumption that all power series in this article have no constant coefficient. We write $f_{k}$ for the residue class of the power series $f$ in $A_{k}^{\prime}[[X]]$. Denote by $\epsilon$ the commutator

$$
\epsilon=f \circ[\pi]_{F^{\prime}}-_{F^{\prime}}[\pi]_{F^{\prime}} \circ f \in A^{\prime}[[X]]
$$

using the additive operation on $A^{\prime}[[X]]$ induced by $F^{\prime}$. Then $\epsilon$ has coefficients in $\left(\pi^{\prime}\right)^{n}$ because $f_{n-1}$ is an endomorphism of $F_{n-1}^{\prime}$.

The main technique to prove the lifting theorem is the cohomology theory as in [VZ]. Denote by $I_{n}$ the 1 -dimensional $k$-vector space $\left(\pi^{\prime}\right)^{n} /\left(\pi^{\prime}\right)^{n+1}$. Consider the cohomology group $H^{2}\left(F_{0}, I_{n}\right)$ as in $[\mathbf{V Z}]$ Chapter 3. For $f_{n-1} \in \operatorname{End} F_{n-1}^{\prime}$ one can define a cocycle $\left(\Delta,\left\{\delta_{a}\right\}\right) \in H^{2}\left(F_{0}, I_{n}\right)$. Then $f_{n-1}$ lifts to End $F_{n}^{\prime}$ if and only if $\left(\Delta,\left\{\delta_{a}\right\}\right) \equiv 0$, i.e., if and only if $\delta_{\pi}$ is a power series in $X^{q^{2}}([\mathbf{V Z}]$ Prop. 3.2, Cor. 3.4). We have

$$
\delta_{\pi}=\epsilon \bmod \left(\pi^{\prime}\right)^{n+1} .
$$

Lemma 3.4. - The cohomology group $H^{2}\left(F_{0}, I_{n}\right)$ is a $k$-vector space of dimension 1. For a cocycle $\left(\Delta,\left\{\delta_{a}\right\}\right) \in H^{2}\left(F_{0}, I_{n}\right)$, the element $\delta_{\pi}=\beta\left(X^{q}\right)$ is a power series in $X^{q}$ and $\left(\Delta,\left\{\delta_{a}\right\}\right) \not \equiv 0$ if and only if $\beta^{\prime}(0) \neq 0$.

Proof. - By [VZ, Lemma 2.5], every formal module over $k$ is isomorphic to a normal module. Then [VZ, Proposition 3.6], shows that $H^{2}\left(F_{0}, k\right)$ is a $k$-vector space of dimension 1. A basis is given by a cocycle $\left(\Delta,\left\{\delta_{a}\right\}\right)$ such that $\delta_{\pi}=\beta\left(X^{q}\right)$ is a power series in $X^{q}$ with $\beta^{\prime}(0) \neq 0$. This proves the lemma.

Remark 3.5. - Let $f_{0} \in$ End $F_{n-1}^{\prime}$. By Lemma 3.4 the power series $\epsilon$ is a power series in $X^{q}$ modulo $\left(\pi^{\prime}\right)^{n+1}$,

$$
\begin{equation*}
\epsilon \equiv a X^{q}+\ldots \bmod \left(\pi^{\prime}\right)^{n+1} \tag{3.1}
\end{equation*}
$$

Furthermore, $v_{\pi^{\prime}}(a) \geq n$ and $v_{\pi^{\prime}}(a)=n$ if and only if $f_{0} \notin \operatorname{End} F_{n}^{\prime}$.
Lemma 3.6. - Let $f_{0} \in \operatorname{End} F_{n-1}^{\prime}$ and let $k=\min \left\{n+e_{s}, 1+q n\right\}$. Then $[\pi]_{F_{0}} \circ f_{0}$ lifts to End $F_{k-1}^{\prime}$.
(i) If $1+q n<n+e_{s}$, the endomorphism $[\pi]_{F_{0}} \circ f_{0}$ lifts to End $F_{k}^{\prime}$.
(ii) If $k=n+e_{s}$ and $f_{0} \notin \operatorname{End} F_{n}^{\prime}$, the endomorphism $[\pi]_{F_{0}} \circ f_{0}$ does not lift to End $F_{k}^{\prime}$.

Proof. - We use the notations of 3.1. By equation (1.1) we obtain

$$
\begin{align*}
{[\pi]_{F^{\prime}} \circ f \circ[\pi]_{F^{\prime}}-F^{\prime}[\pi]_{F^{\prime}}^{2} \circ f } & =[\pi]_{F^{\prime}} \circ \epsilon \\
& =\pi \epsilon+\cdots+u \epsilon^{q}+\cdots+v \epsilon^{q^{2}}+\ldots \tag{3.2}
\end{align*}
$$

Since $\epsilon$ has coefficients in $\left(\pi^{\prime}\right)^{n}$, we have

$$
[\pi]_{F^{\prime}} \circ \epsilon \equiv 0 \bmod \left(\pi^{\prime}\right)^{k} .
$$

Thus $[\pi]_{F_{k-1}^{\prime}} \circ f_{k-1}$ commutes with $[\pi]_{F_{k-1}^{\prime}}$, hence it is an element of End $F_{k-1}^{\prime}([\mathbf{V Z}$, Cor. 3.1]). We obtain by (3.2)

$$
\begin{equation*}
\delta_{\pi}\left([\pi]_{F_{k-1}^{\prime}} \circ f_{k-1}\right) \equiv \pi \epsilon+\cdots+u \epsilon^{q}+\cdots+v \epsilon^{q^{2}}+\ldots \bmod \left(\pi^{\prime}\right)^{k+1} \tag{3.3}
\end{equation*}
$$

If $1+q n<n+e_{s}$, the power series (3.3) is a power series in $X^{q^{2}}$ as $\epsilon$ is a power series in $X^{q}$ modulo $\left(\pi^{\prime}\right)^{n+1}$. Hence $[\pi]_{F_{0}} \circ f_{0}$ lifts to End $F_{k}^{\prime}$.

If $k=n+e_{s}$, we obtain

$$
\delta_{\pi}\left([\pi]_{F_{k-1}^{\prime}} \circ f_{k-1}\right) \equiv \pi a X^{q}+\ldots \bmod \left(\pi^{\prime}\right)^{k+1}
$$

with $v_{\pi^{\prime}}(\pi a)=n+e_{s}$. Hence $[\pi]_{F_{n^{\prime}-1}^{\prime}} \circ f_{n^{\prime}-1}$ does not lift to End $F_{k}^{\prime}$.
Proof of Proposition 3.2. - Since $n \geq \frac{e_{s}-1}{q-1}$, we obtain $\min \left\{n+e_{s}, 1+q n\right\}=n+e_{s}=$ $n^{\prime}$. The proposition follows from Lemma 3.6.

Proof of Proposition 3.1 in the case of $n \geq e_{s}$. - By Remark 3.3 we have to prove the following cases.

1. $L / K$ unramified and $l=2 s+1$.
2. $L / K$ ramified and $l=2 s+1$.
3. $L / K$ ramified, $q=2$ and $l=2 s$.
4. $L / K$ unramified, $q=2, l=2 s-1$ and $s \neq 1$.

Note that $l \geq 2$. Let $f_{0}$ be an element of $\left(\mathcal{O}_{s}+\pi_{D}^{l} \mathcal{O}_{D}\right) \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l+1} \mathcal{O}_{D}\right)$. Write $f_{0}=c+[\pi]_{F_{0}} \circ g_{0}$ with $c \in \mathcal{O}_{s}$ and $g_{0} \in \pi_{D}^{l-2} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{s}+\pi_{D}^{l-1} \mathcal{O}_{D}\right)$. Since elements of $\mathcal{O}_{s}$ lift to End $F^{\prime}$, it is enough to show that $[\pi]_{F_{0}} \circ g_{0}$ satisfies the claim. As $g_{0}$ is an element of $\pi_{D}^{2 s-1} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{s}+\pi_{D}^{2 s} \mathcal{O}_{D}\right)$, it lifts to End $F_{n-1}^{\prime} \backslash$ End $F_{n}^{\prime}$ with $n=n_{l-2}$. We have

$$
\begin{equation*}
n=a(s-1)+q^{s-1}+1=\frac{2 q^{s}-2}{q-1} \tag{3.4}
\end{equation*}
$$

In the first case, we obtain $n_{l-2} \geq \frac{e_{s}-1}{q-1}$ and the claim follows from the case $l=2 s-1$ from Proposition 3.2.

Now consider the other cases. Note that in these cases $n+1 \leq e_{s}$ (Rem. 3.3). Let $n^{\prime}=n_{l}$. We have to show that $[\pi]_{F_{0}} \circ g_{0}$ lifts to End $F_{n^{\prime}-1}^{\prime} \backslash$ End $F_{n^{\prime}}^{\prime}$. An easy
calculation shows that in each case $n^{\prime}=q n+2$. By equation (3.4) we see that $e_{s}+n=q n$ in the second case, and $e_{s}+n>q n+2$ in the other cases. Now we can use Lemma 3.6 (ii) to see that $\pi \circ g_{0}$ lifts to End $F_{n^{\prime}-1}^{\prime}$. Let $h_{n^{\prime}-1} \in \operatorname{End} F_{n^{\prime}-1}^{\prime}$ be a lift of $[\pi]_{F_{0}^{\prime}} \circ g_{0}$. It remains to show that $h_{n^{\prime}-1}$ does not lift to End $F_{n^{\prime}}^{\prime}$, i.e., $\delta_{\pi}\left(h_{n^{\prime}-1}\right) \in A_{n^{\prime}}^{\prime}[[X]]$ is not equal to zero modulo $(X)^{q^{2}}$.

Let $h_{n^{\prime}} \in A_{n^{\prime}}^{\prime}[[X]]$ be a lift of $h_{n^{\prime}-1}$ as a power series. Then $h_{n^{\prime}}=[\pi]_{F_{n^{\prime}}^{\prime}} \circ g_{n^{\prime}}+F_{n^{\prime}}^{\prime} \psi$ with a power series $\psi=b X+\cdots \in\left(\pi^{\prime}\right)^{n^{\prime}-1}[[X]]$. Using the notation of 3.1, we obtain

$$
\begin{align*}
\delta_{\pi}\left(h_{n^{\prime}-1}\right) & =\delta_{\pi}\left(\pi_{F_{n^{\prime}-1}^{\prime}} \circ g_{n^{\prime}-1}\right)+{ }_{F_{n^{\prime}}^{\prime}} \psi \circ[\pi]_{F_{n^{\prime}}^{\prime}}-F_{F^{\prime}}^{\prime}[\pi]_{F_{n^{\prime}}^{\prime}} \circ \psi \\
& \equiv[\pi]_{F_{n^{\prime}}^{\prime}} \circ \epsilon+_{F_{n^{\prime}}^{\prime}} \psi \circ[\pi]_{F_{n^{\prime}}^{\prime}} \bmod \left(\pi^{\prime}\right)^{n^{\prime}+1} . \tag{3.5}
\end{align*}
$$

By (3.1) we obtain from equation (3.5)

$$
\delta_{\pi}\left(h_{n^{\prime}-1}\right) \equiv(\pi a+b u) X^{q}+\ldots \bmod \left(\pi^{\prime}\right)^{n^{\prime}+1}
$$

It is sufficient to prove the following claim.
Claim. - We have

$$
\pi a+b u \not \equiv 0 \bmod \left(\pi^{\prime}\right)^{n^{\prime}+1}
$$

Indeed, we have $\delta_{\pi}\left(h_{n^{\prime}-1}\right) \equiv 0 \bmod \left(\pi^{\prime}\right)^{n^{\prime}}$ since $h_{n^{\prime}-1}$ is an endomorphism. We obtain from equation (3.5) that

$$
\begin{aligned}
\delta_{\pi}\left(h_{n^{\prime}-1}\right) & \equiv\left(u \epsilon^{q}+\cdots+v \epsilon^{q^{2}}+\ldots\right)+F_{n^{\prime}}^{\prime}\left(b v X^{q^{2}}+\ldots\right) \bmod \left(\pi^{\prime}\right)^{n^{\prime}} \\
& \equiv\left(u a^{q}+b v\right) X^{q^{2}}+\ldots \bmod \left(\pi^{\prime}\right)^{n^{\prime}}
\end{aligned}
$$

hence we have

$$
\begin{equation*}
u a^{q}+b v \equiv 0 \bmod \left(\pi^{\prime}\right)^{n^{\prime}} \tag{3.6}
\end{equation*}
$$

Since $v_{\pi^{\prime}}(a)=n$ (Rem. 3.5) and $n^{\prime}=q n+2$, we obtain that $v_{\pi^{\prime}}(b)=n^{\prime}-1$.

We first consider the last two cases. In these cases, we have $e_{s}+n>n^{\prime}$. Therefore, $\pi a \equiv 0 \bmod \left(\pi^{\prime}\right)^{n^{\prime}+1}$ and the claim is satisfied. Thus the proposition is proved in these cases.

Now consider the second case. Let

$$
g=\alpha X+\cdots \in A^{\prime}[[X]] .
$$

Since $n+1 \leq e_{s}$, we obtain from the definition of $\epsilon$

$$
\epsilon \equiv u\left(\alpha-\alpha^{q}\right) X^{q}+\ldots \bmod \left(\pi^{\prime}\right)^{n+1}
$$

hence

$$
a=u\left(\alpha-\alpha^{q}\right)
$$

As $v_{\pi^{\prime}}(a)=n$, we have $v_{\pi^{\prime}}(\alpha)=n-1$.

Using equation (3.6), we obtain

$$
\pi a+b u \equiv \pi a-v^{-1} u^{2} a \equiv \pi u \alpha-v^{-1} u^{q+2} \alpha^{q} \bmod \left(\pi^{\prime}\right)^{n^{\prime}+1} .
$$

The idea is to analyze the solutions of the equation

$$
\begin{equation*}
\pi \alpha-v^{-1} u^{q+1} \alpha^{q} \equiv 0 \bmod \left(\pi^{\prime}\right)^{n^{\prime}} \tag{3.7}
\end{equation*}
$$

There are $q$ different solutions of this equation for $\alpha \in\left(\pi^{\prime}\right)^{n-1} /\left(\pi^{\prime}\right)^{n}$. We will identify these solutions as first coefficients of endomorphisms corresponding to elements of $\mathcal{O}_{s}$.

Consider the following general situation. Let $f_{0}$ and $f_{0}^{\prime}$ be two elements of $\pi_{D}^{2 s+1} \mathcal{O}_{D}$ which are not equivalent modulo $\pi_{D}^{2 s+2} \mathcal{O}_{D}$. As before, we write $f_{0}=[\pi]_{F_{0}^{\prime}} \circ g_{0}$ and $f_{0}^{\prime}=[\pi]_{F_{0}^{\prime}} \circ g_{0}^{\prime}$. We obtain

$$
g_{0}-g_{0}^{\prime} \in \pi_{D}^{2 s-1} \mathcal{O}_{D} \backslash \pi_{D}^{2 s} \mathcal{O}_{D}=\pi_{D}^{2 s-1} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{s}+\pi_{D}^{2 s} \mathcal{O}_{D}\right)
$$

Hence the endomorphism $g_{0}-g_{0}^{\prime}$ lifts to End $F_{n-1}^{\prime} \backslash$ End $F_{n}^{\prime}$. Write $g=\alpha X+\ldots$ and $g^{\prime}=\alpha^{\prime} X+\ldots$ as before. We obtain $v_{\pi^{\prime}}\left(\alpha-\alpha^{\prime}\right)=n-1$, hence $\alpha$ and $\alpha^{\prime}$ are not equivalent modulo $\left(\pi^{\prime}\right)^{n}$. Thus different equivalence classes of endomorphisms belong to different equivalence classes of coefficients. As $L / K$ is a ramified extension in the division algebra $D$, we have

$$
\begin{aligned}
\left(\mathcal{O}_{s} \cap \pi_{D}^{2 s+1} \mathcal{O}_{D}\right) / \pi_{D}^{2 s+2} \mathcal{O}_{D} & =\left(\left(\mathcal{O}_{K}+\pi^{s} \mathcal{O}_{L}\right) \cap \pi_{D}^{2 s+1} \mathcal{O}_{D}\right) / \pi_{D}^{2 s+2} \mathcal{O}_{D} \\
& =\pi^{s+1} \mathcal{O}_{L} /\left(\pi_{D}^{2 s+2} \mathcal{O}_{D} \cap \pi^{s+1} \mathcal{O}_{L}\right) \\
& =\mathcal{O}_{L} / \pi_{L} \mathcal{O}_{L} \cong \mathbb{F}_{q}
\end{aligned}
$$

Thus the $q$ different solutions of (3.7) correspond to the equivalence classes of endomorphisms of $\mathcal{O}_{s}$ in End $F_{n^{\prime}}^{\prime}$. By our assumption $[\pi]_{F_{0}^{\prime}} \circ g_{0} \notin \mathcal{O}_{s}+\pi_{D}^{2 s+2} \mathcal{O}_{D}$, hence equation (3.7) is not satisfied which proves the claim.

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# 12. INVARIANTS OF TERNARY QUADRATIC FORMS 

by

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#### Abstract

This paper deals with Gross-Keating invariants of ternary quadratic forms over $\mathbb{Z}_{\ell}$. The main technical difficulties arise in residue characteristic $\ell=2$. In this case, we define the Gross-Keating invariants in terms of a normal form. We give an alternative, less computational approach for anisotropic quadratic forms.


## Résumé (Invariants de Gross-Keating pour les formes quadratiques ternaires)

Cet article concerne les invariants de Gross-Keating pour les formes quadratiques ternaires sur $\mathbb{Z}_{\ell}$. Les difficultés principales n'apparaissent qu'en caractéristique résiduelle $\ell=2$. Dans ce cas, nous déterminons les invariants de Gross-Keating en termes d'une forme normale. Pour les formes anisotropes nous donnons une approche plus directe.

This note provides details on [GK, Section 4]. The main goal is to define and compute the Gross-Keating invariants $a_{1}, a_{2}, a_{3}$ of ternary quadratic forms over $\mathbb{Z}_{\ell}$ (Definition 1.2). If $a_{1} \equiv a_{2} \bmod 2$ and $a_{3}>a_{2}$ we define an additional invariant $\epsilon \in\{ \pm 1\}$ (Definition 2.7, Definition 4.8). If $\ell \neq 2$ every quadratic form over $\mathbb{Z}_{\ell}$ is diagonalizable, and it is easy to determine these invariants from the diagonal form (Section 2). If $\ell=2$ not every quadratic form is diagonalizable. Moreover, even for diagonal quadratic forms it is not straightforward to determine the Gross-Keating invariants. We determine a normal form in Section 3 and compute the invariants in terms of this normal form (Section 4). In Section 5 we determine explicitly when a ternary quadratic form is anisotropic. A complete table can be found in Proposition 5.2 (non diagonalizable case) and Theorem 5.7 (diagonalizable case). In Section 6, we give an alternative definition of the Gross-Keating invariants for anisotropic quadratic forms. The results of Section 6 are due to Stefan Wewers, following a hint in [GK, Section 4].

Our main reference on quadratic forms over $\mathbb{Z}_{\ell}$ is $[\mathbf{C}$, Chapter 8]. Most of the results of this paper can also be found in the work of Yang, in a somewhat different

2000 Mathematics Subject Classification. - 11E08.
Key words and phrases. - Ternary quadratic forms, Gross-Keating invariants, anisotropic quadratic forms.
form. The Gross-Keating invariants are computed in [Y1, Appendix B]. The question whether a given form over $\mathbb{Z}_{2}$ is isotropic or not (Section 5) is discussed in [Y2].

I would like to thank M. Rapoport for comments on an earlier version.

## 1. Definition of the invariants $a_{i}$

In this section we give the general definition of the Gross-Keating invariants $a_{i}$ of quadratic forms over $\mathbb{Z}_{\ell}$ which are used in $[\mathbf{G K}]$.

Let $L$ be a free $\mathbb{Z}_{\ell}$-module of rank $n$ and choose a (for the moment) arbitrary basis $\psi=\left\{\psi_{1}, \psi_{2}, \ldots, \psi_{n}\right\}$. For the application to $[\mathbf{G K}]$ we are only interested in the case $n=3$ of ternary quadratic forms. Let $(L, Q)$ be an integral quadratic form over $\mathbb{Z}_{\ell}$, that is,

$$
Q(x)=Q\left(\sum x_{i} \psi_{i}\right)=\sum_{i \leq j} b_{i j} x_{i} x_{j}, \quad \text { with } b_{i j} \in \mathbb{Z}_{\ell} .
$$

Put $b_{j i}=b_{i j}$ for $j>i$. If we want to stress the dependence of the $b_{i j}$ on the basis, we write $b_{i j}(\boldsymbol{\psi})$ for $b_{i j}$. We write $(x, y)=Q(x+y)-Q(x)-Q(y)$ for the corresponding symmetric bilinear form and $B=\left(\left(\psi_{i}, \psi_{j}\right)\right)$ for the corresponding matrix. Note that

$$
B=\left(B_{i j}\right), \quad \text { where } \quad B_{i j}= \begin{cases}b_{i j}, & \text { if } i<j \\ 2 b_{i j}, & \text { if } i=j .\end{cases}
$$

In the rest of the paper we only use the $b_{i j}$ and not the $B_{i j}$, for simplicity. We denote by ord the $\ell$-adic valuation on $\mathbb{Z}_{\ell}$. We always suppose that $Q$ is regular, that is, $\operatorname{det}(B) \neq 0$.

Changing the basis multiplies the determinant of $B$ by an element of $\left(\mathbb{Z}_{\ell}^{\times}\right)^{2}$. Therefore the determinant is a well defined element of $\mathbb{Z}_{\ell} /\left(\mathbb{Z}_{\ell}^{\times}\right)^{2}$.

Lemma 1.1. - Suppose that either $\ell \neq 2$ or $n$ is odd. Define

$$
\Delta=\Delta(Q)=\frac{1}{2} \operatorname{det}(B)
$$

Then $\Delta \in \mathbb{Z}_{\ell}$.
Proof. - The lemma is obvious if $\ell \neq 2$. Suppose that $\ell=2$ and $n$ odd. Write $\Delta=\sum_{\sigma \in S_{n}} 2^{\delta(\sigma)} d(\sigma)$, where $d(\sigma)=(-1)^{\operatorname{sgn}(\sigma)} \prod_{i=1}^{n} b_{i \sigma(i)}$ and $\delta(\sigma)+1$ is the number of $i \in\{1,2, \ldots, n\}$ which are fixed by $\sigma$. The only problematic terms are those with $\delta(\sigma)=-1$. Suppose that $\sigma$ acts without fixed points on $\{1,2, \ldots, n\}$. Then $\sigma^{-1} \neq \sigma$, since $n$ is odd. The matrix $\left(\left(\psi_{i}, \psi_{j}\right)\right)$ is symmetric. It follows that $d(\sigma)=d\left(\sigma^{-1}\right)$, hence $2^{\delta(\sigma)} d(\sigma)+2^{\delta\left(\sigma^{-1}\right)} d\left(\sigma^{-1}\right) \in \mathbb{Z}_{\ell}$.

We now come to the definition of the Gross-Keating invariants of a quadratic form. Let $\boldsymbol{\psi}=\left(\psi_{1}, \psi_{2}, \ldots, \psi_{n}\right)$ be a basis of $L$. We write $S(\boldsymbol{\psi})$ for the set of tuples $\mathbf{y}=\left(y_{1}, y_{2}, \ldots, y_{n}\right) \in \mathbb{Z}^{n}$ such that

$$
\begin{equation*}
y_{1} \leq y_{2} \leq \cdots \leq y_{n}, \quad \frac{y_{i}+y_{j}}{2} \leq \operatorname{ord}\left(b_{i j}(\boldsymbol{\psi})\right) \quad \text { for } 1 \leq i \leq j \leq n \tag{1.1}
\end{equation*}
$$

Let $S=\cup S(\boldsymbol{\psi})$. We order tuples $\left(y_{1}, \ldots, y_{n}\right) \in S$ lexicographically, as follows. For given $\left(y_{1}, \ldots, y_{n}\right),\left(z_{1}, \ldots, z_{n}\right) \in S$, let $j$ be the largest integer such that $y_{i}=z_{i}$ for all $i<j$. Then $\left(y_{1}, \ldots, y_{n}\right)>\left(z_{1}, \ldots, z_{n}\right)$ if $y_{j}>z_{j}$.

Definition 1.2. - The Gross-Keating invariants $a_{1}, \ldots, a_{n}$ are the maximum of $\left(y_{1}, \ldots, y_{n}\right) \in S$. A basis $\boldsymbol{\psi}$ is called optimal if $\left(a_{1}, \ldots, a_{n}\right) \in S(\boldsymbol{\psi})$.

If $\psi$ is optimal, then

$$
\begin{equation*}
a_{i}+a_{j} \leq 2 \operatorname{ord}\left(b_{i j}(\boldsymbol{\psi})\right) \quad \text { for } \quad 1 \leq i \leq j \leq n, \quad \text { and } \quad a_{1} \leq a_{2} \leq \cdots \leq a_{n} \tag{1.2}
\end{equation*}
$$

Since $\Delta$ is well defined up to $\left(\mathbb{Z}_{\ell}^{\times}\right)^{2}$, the integer ord $(\Delta)$ is well defined. The following lemma will be useful in computing the Gross-Keating invariants.

## Lemma 1.3

(a) Suppose that $n$ is odd, then

$$
\operatorname{ord}(\Delta) \geq a_{1}+a_{2}+\cdots+a_{n}
$$

(b) We have

$$
a_{1}=\min _{x, y \in L} \operatorname{ord}(x, y) .
$$

(c) Define $\rho:=\min _{A} \operatorname{ord}(\operatorname{det}(A))$, where $A$ runs through the 2 by 2 minors of $B$. Then

$$
a_{1}+a_{2} \leq \rho .
$$

Proof. - This lemma is proved in [Y1, Lemma B.1, Lemma B.2]. Note that the matrix $T$ in $[\mathbf{Y 1}]$ differs by a factor 2 from our matrix $B$. Let $\varphi$ be an optimal basis. We use the notation of the proof of Lemma 1.1.

First suppose that $\ell=2$. Write $\mathbb{S}$ for the set of equivalence classes in $S_{n}$ under the equivalence relation $\sigma \sim \sigma^{-1}$. The proof of Lemma 1.1 shows that $\Delta=$ $\sum_{\sigma \in \mathbb{S}}(-1)^{\operatorname{sgn}(\sigma)} 2^{\delta^{\prime}(\sigma)} d(\sigma)$, where $\delta^{\prime}(\sigma) \geq 0$. The choice of $\varphi$ implies that

$$
\operatorname{ord}\left(2^{\delta^{\prime}(\sigma)} d(\sigma)\right)=\delta^{\prime}(\sigma)+\operatorname{ord}\left(\prod_{i} b_{i \sigma(i)}\right) \geq \sum_{i=1}^{n} \frac{a_{i}+a_{\sigma(i)}}{2}=\sum_{i=1}^{n} a_{i} .
$$

This proves (a) in this case.
If $\ell \neq 2$, define $\delta^{\prime}(\sigma)=0$ for all $\sigma \in S_{n}$. Then the proof works also in this case.
Since $a_{1} \leq a_{2} \leq \cdots \leq a_{n}$, it follows from (1.2) that ord $\left(b_{i j}(\varphi)\right) \geq a_{1}$ for all $i \leq j$. On the other hand, it is obvious that $a_{1} \geq \min _{x, y \in L}$ ord $(x, y)$. This implies (b).

Part (c) is similar to (a), compare to Lemma B1.ii in [Y1]. Let $i_{1}, i_{2}, j_{1}, j_{2} \in$ $\{1,2, \ldots, n\}$ be integers such that $i_{1} \neq i_{2}$ and $j_{1} \neq j_{2}$. Write $B\left(i_{1}, i_{2} ; j_{1}, j_{2}\right)$ for the corresponding minor of $B$. After renumbering, we may suppose that $i_{1} \neq j_{2}$ and $i_{2} \neq j_{1}$. Then $\operatorname{det}\left(B\left(i_{1}, i_{2} ; j_{1}, j_{2}\right)\right)= \pm\left(2^{\alpha} b_{i_{1}, j_{1}} b_{i_{2}, j_{2}}-b_{i_{1}, j_{2}} b_{i_{2} j_{1}}\right)$, where $\alpha \in$ $\{0,1,2\}$ is the number of equalities $i_{1}=j_{1}, i_{2}=j_{2}$ that hold. We conclude that $\operatorname{ord}\left(\operatorname{det}\left(B\left(i_{1}, i_{2} ; j_{1}, j_{2}\right)\right) \geq\left(a_{i_{1}}+a_{i_{2}}+a_{j_{1}}+a_{j_{2}}\right) / 2 \geq a_{1}+a_{2}\right.$. (Here we use that $a_{1} \leq a_{2} \leq \cdots \leq a_{n}$ and $i_{1} \neq i_{2}$ and $j_{1} \neq j_{2}$.) This proves (c).

## 2. Definition of the Gross-Keating invariants for $\ell \neq 2$

We start this section with an elementary lemma which holds without assumption on $\ell$.

Lemma 2.1. - Choose a basis $\boldsymbol{\psi}=\left(\psi_{1}, \ldots, \psi_{n}\right)$ of L. Let $\gamma_{1}, \ldots, \gamma_{m} \in L$ be linearly independent. The following are equivalent.
(a) There exists $\gamma_{m+1}, \ldots, \gamma_{n} \in L$ such that the $\left(\gamma_{i}\right)$ form a basis.
(b) The matrix $\left(\gamma_{1}, \ldots, \gamma_{m}\right)$, expressing the $\gamma_{i}$ in terms of the basis $\boldsymbol{\psi}$, contains a $m \times m$ minor whose determinant is a $p$-adic unit.
(c) If $\sum_{i=1}^{n} v_{i} \gamma_{i} \in L$ for some $v_{i} \in \mathbb{Q}_{\ell}$, then $v_{i} \in \mathbb{Z}_{\ell}$.

Proof. - This is straightforward. See also [C, Chapter 8, Lemma 2.1].
In particular, a vector $\alpha=\sum_{i} \alpha_{i} \psi_{i} \in L$ is part of a basis of $L$ if and only if $\min _{j} \operatorname{ord}\left(\alpha_{j}\right)=0$. We call such vectors primitive.

We have that

$$
\begin{equation*}
2(x, y)=2[Q(x+y)-Q(x)-Q(y)]=(x+y, x+y)-(x, x)-(y, y) \tag{2.1}
\end{equation*}
$$

If $\ell \neq 2$, this implies that

$$
\begin{equation*}
\min _{x, y \in L} \operatorname{ord}(x, y)=\min _{x \in L} \operatorname{ord}(x, x) . \tag{2.2}
\end{equation*}
$$

In the rest of this section, we suppose that $\ell \neq 2$. There is a $x \in L$ for which the minimum in (2.2) is attained. This vector $x$ is primitive. Lemma 2.1 implies that $x$ can be extended to a basis of $L$. We will see in Section 4 that (2.2) does not hold for $\ell=2$; this is the main reason why things are more difficult for $\ell=2$.

Proposition 2.2. - Suppose that $\ell \neq 2$. Then there exists a basis $\boldsymbol{\psi}$ of $L$ such that $Q(x)=Q\left(\sum x_{i} \psi_{i}\right)=\sum_{i} b_{i i} x_{i}^{2}, \quad$ where $\quad \operatorname{ord}\left(b_{11}\right) \leq \operatorname{ord}\left(b_{22}\right) \leq \cdots \leq \operatorname{ord}\left(b_{n n}\right)$.

Proof. - Our proof follows [C, Chapter 8, Theorem 3.1].
The discussion before the statement of the theorem shows that we may choose $\varphi_{1}$ such that

$$
\operatorname{ord}\left(Q\left(\varphi_{1}\right)\right)=\operatorname{ord}\left(\varphi_{1}, \varphi_{1}\right)=\min _{x, y \in L} \operatorname{ord}(x, y)
$$

Here we use the equality (2.2).
Choose $\varphi_{2}, \ldots, \varphi_{n} \in L$ such that $\varphi=\left\{\varphi_{1}, \varphi_{2} \ldots, \varphi_{n}\right\}$ is a basis of $L$. As before we write $Q\left(\sum_{i} x_{i} \varphi_{i}\right)=\sum_{1 \leq i \leq j \leq n} b_{i j}(\boldsymbol{\varphi}) x_{i} x_{j}$. Then

$$
Q(x)=b_{11}\left(x_{1}+\frac{b_{12}}{2 b_{11}} x_{2}+\cdots \frac{b_{1 n}}{2 b_{11}} x_{n}\right)^{2}+\tilde{Q}\left(x_{2}, \ldots, x_{n}\right)
$$

for some integral quadratic form $\tilde{Q}$ in $n-1$ variables.

We define a new basis by $\psi_{1}=\varphi_{1}$, and $\psi_{i}=\varphi_{i}-\left(b_{1 i} / 2 b_{11}\right) \varphi_{1}$ for $i \neq 1$. The choice of $\psi_{1}$ ensures that $\psi_{i} \in L$, since $e=\operatorname{ord}\left(2 b_{11}\right) \leq \operatorname{ord}\left(b_{1 i}\right)$. With respect to this new basis, the quadratic form is

$$
Q(x)=b_{11}(\boldsymbol{\psi}) x_{1}^{2}+\tilde{Q}\left(\sum_{i \geq 2} x_{i} \psi_{i}\right)
$$

The proposition follows by induction.
Remark 2.3. - Cassels ([C, Chapter 8, Theorem 3.1]) proves a stronger statement than Proposition 2.2. Namely, he gives a list of pairwise nonisomorphic quadratic forms such that every integral quadratic form is isomorphic to one of these. This stronger statement implies that the definition of the invariants $a_{i}$ of Proposition 2.6 does not depend of the choice of the orthogonal basis.

We can give a simpler definition of the invariants $a_{i}$ in terms of a basis $\psi$ as in Proposition 2.2. If $\gamma \in L$ is an element such that $Q(\gamma) \neq 0$, we may define a reflection $\tau_{\gamma}$ by

$$
\tau_{\gamma}(x)=x-\frac{2(x, \gamma)}{(\gamma, \gamma)} \gamma
$$

This is the reflection in the orthogonal complement of $\gamma$. Clearly, $\tau_{\gamma}$ is defined over $\mathbb{Z}_{\ell}$ if and only if $\operatorname{ord}(\gamma, \gamma)=\min _{x \in L}$ ord $(x, x)$. (In fact, this also holds for $\ell=2$.) Since $\tau_{\gamma}$ is a reflection, it is clearly invertible. The following lemma is a partial analog of Witt's Lemma ([C, Corollary to Theorem 2.4.1]) which holds for quadratic forms over fields.

Lemma 2.4. - Suppose that $\psi, \varphi \in L$ satisfy

$$
Q(\psi)=Q(\varphi), \quad \operatorname{ord}(Q(\psi))=\operatorname{ord}(Q(\varphi))=\min _{x \in L} \operatorname{ord}(Q(x))
$$

Then there exists an integral isometry $\sigma$ of $(L, Q)$ such that $\sigma(\psi)=\varphi$. Moreover, $\sigma$ may be taken as a product of reflections $\tau_{\gamma}$.

Proof. - This is [C, Lemma 8.3.3]. Our assumptions on $\psi$ and $\varphi$ imply that $Q(\psi+\varphi)+Q(\psi-\varphi)=2 Q(\psi)+2 Q(\varphi)=4 Q(\psi)$. Since $\operatorname{ord}(Q(\psi))=\operatorname{ord}(\psi, \psi)=$ $\min _{x \in L}$ ord $(x, x)=: e$, it follows that one of the following holds:
(a) $\operatorname{ord} Q(\psi+\varphi)=e$,
(b) $\operatorname{ord} Q(\psi-\varphi)=e$.

Since $\ell \neq 2$, it is also possible that both hold. If (a) holds, then $\tau_{\psi+\varphi}$ is integral and sends $\psi$ to $\varphi$. If (b) holds, define $\sigma=\tau_{\psi-\varphi} \circ \tau_{\psi}$.

Lemma 2.5. - Suppose $u, v \in \mathbb{Z}_{\ell}^{\times}$. Then $u x_{1}^{2}+v x_{2}^{2} \sim_{\mathbb{Z}_{\ell}} x_{1}^{2}+u v x_{2}^{2}$.

Proof. - This is proved in the second corollary to [C, Lemma 8.3.3]. We give the idea. Since $\ell \neq 2$, there exists $a, c \in \mathbb{Z}_{\ell}$ such that $a^{2} u+c^{2} v=1$. We may assume that $a$ is a unit. Then

$$
C=\left(\begin{array}{cc}
a & -c v \\
c & a u
\end{array}\right)
$$

defines the equivalence of the lemma.

## Proposition 2.6

(a) Let $\boldsymbol{\psi}=\left(\psi_{1}, \psi_{2}, \ldots, \psi_{n}\right)$ be an orthogonal basis of $L$ as in Proposition 2.2 Write $Q(x)=\sum_{i} b_{i} x_{i}^{2}$. Then the invariants $a_{i}$ (Definition 1.2) satisfy

$$
a_{i}=\operatorname{ord}\left(b_{i}\right) .
$$

In particular, $\boldsymbol{\psi}$ is optimal.
(b) Suppose that $n$ is odd. Then

$$
\operatorname{ord}(\Delta)=a_{1}+\cdots+a_{n}
$$

Proof. - Let $\varphi$ be a basis such that the inequalities (1.2) hold. We claim that ord $\left(\varphi_{1}, \varphi_{1}\right)=a_{1}$. Part (b) of Lemma 1.3 implies that $a_{1}=\min _{x \in L}$ ord $(x, x)$. The choice of $\boldsymbol{\varphi}$ implies moreover that ord $\left(\varphi_{1}, \varphi_{1}\right)=\min _{x \in L}$ ord $(x, x)$. The definition of $a_{1}$ implies therefore that $a_{1}=\operatorname{ord}\left(\varphi_{1}, \varphi_{1}\right)$.

We apply the diagonalization process of the proof of Proposition 2.2 to the basis $\varphi$. Define $\psi_{1}=\varphi_{1}$ and $\psi_{i}=\varphi_{i}-\left(b_{1 i} / 2 b_{11}\right) \varphi_{1}$ for $i \neq 1$. One computes that

$$
\left(\psi_{j}, \psi_{1}\right)=0, \quad\left(\psi_{j}, \psi_{j}\right)=\frac{b_{1 j}^{2}}{2 b_{11}}+2 b_{j j}, \quad\left(\psi_{i}, \psi_{j}\right)=-\frac{b_{1 i} b_{1 j}}{2 b_{11}}+b_{i j}
$$

for $j \neq 1$ and $i \neq 1, j$. The inequalities (1.2) imply that $\operatorname{ord}\left(\psi_{j}, \psi_{j}\right) \geq a_{j}$ and 2 ord $\left(\psi_{i}, \psi_{j}\right) \geq a_{i}+a_{j}$. Therefore the new basis also satisfies the inequalities (1.2). This implies that there exists an orthogonal basis $\boldsymbol{\psi}$ which satisfies (1.2). It follows that the Gross-Keating invariants $\left(a_{1}, \ldots, a_{n}\right)$ are the maximum of $\cup S(\boldsymbol{\psi})$, where the union is taken over the orthogonal bases and $\cup S(\boldsymbol{\psi})$ is as in (1.1).

Let $\boldsymbol{\varphi}$ and $\boldsymbol{\psi}$ be two orthogonal bases. Write $Q(x)=b_{1} x_{2}^{2}+b_{2} x_{2}^{2}+\cdots+b_{n} x_{n}^{2}$ with respect to the basis $\psi$ and $Q(x)=d_{1} x_{1}^{2}+d_{2} x_{2}^{2}+\cdots+d_{n} x_{n}^{2}$ with respect to the basis $\varphi$. We suppose that $\operatorname{ord}\left(b_{1}\right) \leq \operatorname{ord}\left(b_{2}\right) \leq \cdots \leq \operatorname{ord}\left(b_{n}\right)$ and $\operatorname{ord}\left(d_{1}\right) \leq \operatorname{ord}\left(d_{2}\right) \leq \cdots \leq$ $\operatorname{ord}\left(d_{n}\right)$. We suppose moreover that $\varphi$ satisfies (1.2). (Such $\varphi$ exists by the above argument.) We have to show that $\boldsymbol{\psi}$ satisfies (1.2), also. Write $C=\left(c_{i j}\right)$ for the change of basis matrix expressing $\varphi$ in terms $\boldsymbol{\psi}$. As before, Lemma 1.3.(b) implies that $\operatorname{ord}\left(b_{1}\right)=\operatorname{ord}\left(d_{1}\right)=a_{1}$. Write $b_{1}=u d_{1}$, for some unit $u$.

Suppose that ord $\left(b_{2}\right)>\operatorname{ord}\left(b_{1}\right)$. Then

$$
d_{1}=\sum_{j=1}^{n} c_{j 1}^{2} b_{j} \equiv c_{11}^{2} b_{1} \bmod \ell^{a_{1}+1}
$$

This implies that $u$ is a quadratic residue. To prove the claim, we may therefore assume that $Q\left(\psi_{1}\right)=Q\left(\varphi_{1}\right)$ in this case.

Suppose that ord $\left(b_{1}\right)=\operatorname{ord}\left(b_{2}\right)$. Then Lemma 2.5 implies that $Q$ is $\mathbb{Z}_{\ell}$-equivalent to $d_{1} x_{1}^{2}+u b_{2} x_{2}^{2}+b_{3} x_{3}^{2}+\cdots$. Hence also in this case we may assume that $Q\left(\psi_{1}\right)=Q\left(\varphi_{1}\right)$.

Lemma 2.4 implies that there exists an isometry $\sigma$ of $Q$ which sends $\psi_{1}$ to $\varphi_{1}$. Then $D:=\sigma^{-1} C$ fixes $\psi_{1}$. Write

$$
D=\left(\begin{array}{cc}
1 & D_{1} \\
0 & D_{2}
\end{array}\right), \quad B:=\left(\begin{array}{ccc}
2 b_{1} & & 0 \\
& \ddots & \\
0 & & 2 b_{n}
\end{array}\right)
$$

where $D_{2}$ is an $(n-1) \times(n-1)$ matrix. One computes that

$$
D^{t} B D=\left(\begin{array}{cc}
2 \gamma^{2} b_{1} & 2 \gamma D_{1} \\
2 \gamma D_{1}^{t} & *
\end{array}\right)
$$

Our assumption implies that $D^{t} B D$ is a diagonal matrix, with diagonal entries $2 d_{i}$. This implies that $D_{1}=(0, \ldots, 0)$. We conclude that $D$ restricts to an integral and invertible map from the sublattice of $L$ spanned by $\psi_{2}, \ldots, \psi_{n}$ to the sublattice spanned by $\varphi_{2}, \ldots, \varphi_{n}$. This implies (a).

Part (b) follows immediately from (a).
Definition 2.7. - Suppose that $n=3$ and $\ell \neq 2$. Assume $a_{1} \equiv a_{2} \bmod 2$, and $a_{3}>a_{2}$. Choose a basis $\boldsymbol{\psi}=\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ of $L$ as in Proposition 2.2. Write $b_{i i}=\ell^{a_{i}} u_{i}$. We define an invariant $\epsilon=\epsilon(\boldsymbol{\psi})$ by the Legendre symbol

$$
\begin{equation*}
\epsilon=\left(\frac{-u_{1} u_{2}}{\ell}\right) \tag{2.3}
\end{equation*}
$$

Lemma 2.8. - Assumptions and notations are as in Definition 2.7.
(a) The invariant $\epsilon(\boldsymbol{\psi})$ does not depend on the choice of the orthogonal basis $\boldsymbol{\psi}$.
(b) We have that $\epsilon=1$ if and only if the subspace of $L \otimes_{\mathbb{Z}_{\ell}} \mathbb{Q}_{\ell}$ spanned by $\psi_{1}$ and $\psi_{2}$ is isotropic.

Proof. - Let $\boldsymbol{\psi}=\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ be a basis of $L$ as in Proposition 2.6, in particular $\boldsymbol{\psi}$ is orthogonal and the valuation of $b_{i}=\left(\psi_{i}, \psi_{i}\right) / 2$ is equal to $a_{i}$, for $i=1,2,3$.

Suppose that $a_{2} \equiv a_{1} \bmod 2$ and $a_{3}>a_{2}$. Write $a_{2}=a_{1}+2 \gamma$. Write $Q^{\prime}$ for the restriction of $Q$ to the sublattice of $L$ spanned by $\psi_{1}$ and $\psi_{2}$. Then $Q^{\prime}(x)=b_{1} x_{1}^{2}+b_{2} x_{2}^{2}$ is equivalent to $\ell^{a_{1}}\left(x_{1}^{2}+u_{1} u_{2} \ell^{2 \gamma} x_{2}^{2}\right)$ (Lemma 2.5). It follows that $Q^{\prime}$ is isotropic if $\epsilon=1$ and anisotropic if $\epsilon=-1$. This proves (b).

Let $\boldsymbol{\varphi}$ be another orthogonal basis and write $Q\left(\sum_{i} x_{i} \varphi_{i}\right)=d_{1} x_{1}^{2}+d_{2} x_{2}^{2}+d_{3} x_{3}^{2}$. We assume that $\operatorname{ord}\left(d_{i}\right)=a_{i}$. Write $C$ for the matrix expressing $\varphi$ in terms of $\boldsymbol{\psi}$. The argument of the proof of Proposition 2.6 together with the assumption that $a_{2}<a_{3}$
implies that there exists an isometry $\sigma$ such that

$$
\sigma^{-1} C=\left(\begin{array}{ccc}
v_{1} & 0 & 0 \\
0 & v_{2} & 0 \\
0 & 0 & v_{3}
\end{array}\right)
$$

where the $v_{i}$ are units. This shows that $d_{i}=v_{i}^{2} b_{i}$. The lemma follows.

## 3. A normal form for quadratic forms over $\mathbb{Z}_{2}$

Not every quadratic form over $\mathbb{Z}_{2}$ is diagonalizable. In this section we give a normal form for ternary quadratic forms over $\mathbb{Z}_{2}$, following [ $\mathbf{C}$, Section 8.4]. Cassels uses a slightly stronger notion of integrality, namely he supposes that $b_{i j} / 2 \in \mathbb{Z}_{\ell}$, for all $i \neq j$. However, this does not make any difference.

Lemma 3.1. - Suppose $\ell=2$. Let $Q$ be a regular quadratic form over $\mathbb{Z}_{2}$. Then $Q$ is $\mathbb{Z}_{2}$-equivalent to a sum of quadratic forms of the form

$$
\begin{equation*}
2^{e} u x^{2} \tag{3.1}
\end{equation*}
$$

for $e \in \mathbb{Z}_{\geq 0}$ and $u \in \mathbb{Z}_{2}^{\times}$, and

$$
\begin{equation*}
2^{e}\left(b_{1} x_{1}^{2}+u x_{1} x_{2}+b_{2} x_{2}^{2}\right) \tag{3.2}
\end{equation*}
$$

with $e \in \mathbb{Z}_{\geq 0}$, and $u \in \mathbb{Z}_{2}^{\times}$.
The equality (2.1) holds for $\ell=2$, but (2.2) does not. However, (2.1) implies that

$$
\min _{x, y \in L} \text { ord }(x, y)+1 \geq \min _{x \in L} \operatorname{ord}(x, x) .
$$

Therefore $\min _{x, y \in L}$ ord $(x, y)$ equals either $\min _{x \in L}$ ord $(x, x)$ or $\min _{x \in L}$ ord $(x, x)-1$. Proof. - Let $e=\min _{x, y \in L}$ ord $(x, y)$. We distinguish two cases.
(a) There exists a $\gamma \in L$ such that $\operatorname{ord}(\gamma, \gamma)=e$.
(b) For all $\gamma \in L$ we have that ord $(\gamma, \gamma)>e$.

Suppose we are in case (a). Then ord $\left(\psi_{1}, \psi_{i}\right) \geq e$, by definition. We can now proceed as in the proof of Proposition 2.2. Namely, $2 b_{11}=2 Q\left(\psi_{1}\right)=\left(\psi_{1}, \psi_{1}\right)$. Therefore $b_{11}$ has valuation $e-1$. For $i \neq 1$, we have that $\operatorname{ord}\left(b_{1 i}\right)=\operatorname{ord}\left(\psi_{1}, \psi_{i}\right) \geq e$. Therefore

$$
\varphi_{i}=\psi_{i}-\left(\frac{b_{1 i}}{2 b_{11}}\right) \psi_{1}
$$

is an element of $L$ and $\psi_{1}, \varphi_{2}, \ldots, \varphi_{n}$ form a basis. With respect to this basis the quadratic form $Q$ becomes $Q(x)=b_{11} x_{1}^{2}+\tilde{Q}\left(x_{2}, \ldots, x_{n}\right)$, for some quadratic form $\tilde{Q}$ in $n-1$ variables.

Suppose we are in case (b). Then ord $(\gamma, \gamma)>e$ for all $\gamma \in L$. We may choose $\psi_{1}, \psi_{2} \in L$ such that ord $\left(\psi_{1}, \psi_{2}\right)=e$. The definition of $e$ implies that $\left(\psi_{1}+\psi_{2}\right) / 2 \notin L$. Lemma 2.1 implies therefore that $\psi_{1}, \psi_{2}$ can be extended to a basis $\psi_{1}, \ldots, \psi_{n}$ of $L$.

The choice of $\psi_{1}$ and $\psi_{2}$ implies that the determinant of the matrix

$$
\left(\begin{array}{cc}
2 b_{11} 2^{-e} & b_{12} 2^{-e} \\
b_{12} 2^{-e} & 2 b_{22} 2^{-e}
\end{array}\right)
$$

is a unit in $\mathbb{Z}_{\ell}$. Therefore we can find $\lambda_{1}^{J}, \lambda_{2}^{j}$ such that

$$
-2 \lambda_{1}^{j} b_{11}-\lambda_{2}^{j} b_{12}+b_{1 j}=0, \quad-2 \lambda_{2}^{j} b_{22}-\lambda_{1}^{j} b_{12}+b_{2 j}=0,
$$

for $j=3, \ldots, n$. Define $\varphi_{j}=\psi_{j}-\lambda_{1}^{j} \psi_{1}-\lambda_{2}^{j} \psi_{2}$. The choice of the $\lambda_{i}^{j}$ implies that $\left(\varphi_{j}, \psi_{1}\right)=\left(\varphi_{j}, \psi_{2}\right)=0$, for $j=3, \ldots, n$.

With respect to the basis $\left(\psi_{1}, \psi_{2}, \varphi_{3}, \ldots, \varphi_{n}\right)$ the quadratic form $Q$ becomes

$$
Q(x)=2^{e}\left(b_{11} x_{1}^{2}+b_{12} x_{1} x_{2}+b_{22} x_{2}^{2}\right)+\tilde{Q}\left(x_{3}, \ldots, x_{n}\right) .
$$

This proves the lemma.
Lemma 3.2. - Let $Q_{2}(x)=b_{11} x_{1}^{2}+b_{12} x_{1} x_{2}+b_{22} x_{2}^{2}$ be a binary quadratic form over $\mathbb{Z}_{2}$ and $L_{2}$ the corresponding free $\mathbb{Z}_{2}$-lattice of rank two.
(a) If $\min \left(\operatorname{ord}\left(b_{11}\right), \operatorname{ord}\left(b_{22}\right)\right)<\operatorname{ord}\left(b_{12}\right)$ then $Q_{2}$ is diagonalizable.
(b) Suppose that $Q_{2}$ is not diagonalizable. Then $Q_{2}$ is anisotropic if and only if $\operatorname{ord}\left(b_{12}\right)=\operatorname{ord}\left(b_{11}\right)=\operatorname{ord}\left(b_{22}\right)$.
(c) Suppose $Q_{2}$ is anisotropic and not diagonalizable. Then $Q_{2}$ is equivalent to

$$
2^{e}\left(x_{1}^{2}+x_{1} x_{2}+x_{2}^{2}\right),
$$

for some e.
(d) Suppose that $Q_{2}$ is isotropic and not diagonalizable. Then $Q_{2}$ is equivalent to

$$
2^{e} x_{1} x_{2}
$$

for some e.

Proof. - Part (a) follows from the proof of Lemma 3.1.
Suppose that $Q_{2}$ is not diagonalizable. Then $\operatorname{ord}\left(b_{12}\right) \leq \min \left(\operatorname{ord}\left(b_{11}\right), \operatorname{ord}\left(b_{22}\right)\right)$, by (a). Part (b) is an elementary Hilbert-symbol computation using [ $\mathbf{S}$, Theorem IV.6].

Suppose that $Q_{2}$ is anisotropic and not diagonalizable. Then (b) implies that $e:=$ $\operatorname{ord}\left(b_{12}\right)=\operatorname{ord}\left(b_{11}\right)=\operatorname{ord}\left(b_{22}\right)$. Part (c) now follows from an elementary computation.

Suppose that $Q_{2}$ is isotropic and not diagonalizable. There exists a primitive vector $\psi_{1}$ such that $Q\left(\psi_{1}\right)=0$. Lemma 2.1 together with the fact that the quadratic form is nondegenerate, implies that there exists a vector $\psi_{2} \in L_{2}$ such that $\psi_{1}, \psi_{2}$ form a basis of $L_{2}$ and $\left(\psi_{1}, \psi_{2}\right) \neq 0$. After multiplying $\psi_{2}$ with a unit, we may suppose that $\left(\psi_{1}, \psi_{2}\right)=2^{e}$, for some $e \geq 0$.

We claim that $\operatorname{ord}\left(\psi_{2}, \psi_{2}\right)>\operatorname{ord}\left(\psi_{1}, \psi_{2}\right)$. Namely, if $\operatorname{ord}\left(\psi_{2}, \psi_{2}\right) \leq \operatorname{ord}\left(\psi_{1}, \psi_{2}\right)$ then $Q_{2}$ is diagonalizable by (a), but this contradicts our assumptions. Therefore

$$
\psi_{2}^{\prime}:=\psi_{2}-\frac{\left(\psi_{2}, \psi_{2}\right)}{2\left(\psi_{1}, \psi_{2}\right)} \psi_{1} \in L_{2} .
$$

Now $\psi_{1}, \psi_{2}^{\prime}$ form a basis of $L$ and $\left(\psi_{2}^{\prime}, \psi_{2}^{\prime}\right)=0$. This proves (d).
Proposition 3.3. - Let $(L, Q)$ be a ternary quadratic form over $\mathbb{Z}_{2}$. One of the following two possibilities occurs.
(a) The form $Q$ is diagonalizable; there exists a basis such that

$$
Q(x)=b_{1} x_{1}^{2}+b_{2} x_{2}^{2}+b_{3} x_{3}^{3}, \quad \text { with } \quad 0 \leq \operatorname{ord}\left(b_{1}\right) \leq \operatorname{ord}\left(b_{2}\right) \leq \operatorname{ord}\left(b_{3}\right) .
$$

(b) The form $Q$ is not diagonalizable; there exists a basis such that
$Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+x_{2} x_{3}+v x_{3}^{2}\right), \quad$ with $\quad v \in\{0,1\}, \quad \mu_{i} \geq 0 \quad$ and $\quad u_{1} \in \mathbb{Z}_{2}^{\times}$.
Proof. - This follows immediately from Lemma 3.1 and Lemma 3.2.
This classification is the same as the classification used (but not explicitly stated) in [Y1, Appendix B]. Note that Yang's matrix $T$ differs by a factor 2 from the matrix $B$ we use. In particular, the invariant $\beta$ used in [Y1, Proposition B.4] satisfies $\beta \geq-1$ rather than $\beta \geq 0$.

## 4. The Gross-Keating invariants for $\ell=2$

In this section we compute the Gross-Keating invariants of ternary quadratic forms $(L, Q)$ over $\mathbb{Z}_{2}$ in terms of the normal form of Proposition 3.3. The computation of the $a_{i}$ can be found in Proposition 4.1 (non-diagonalizable case) and Proposition 4.2 (diagonalizable case). The computation of $\epsilon$ can be found in Proposition 4.9. This section is based on [ $\mathbf{Y} 1$, Appendix B].

We start by considering quadratic forms which are not diagonalizable. Recall from Proposition 3.3 that if $Q$ is not diagonalizable then there exists a basis $\boldsymbol{\psi}$ of $L$ with respect to which we have

$$
\begin{equation*}
Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+x_{2} x_{3}+v x_{3}^{2}\right), \quad \text { with } \quad v \in\{0,1\}, \quad u_{1} \in \mathbb{Z}_{2}^{\times} \tag{4.1}
\end{equation*}
$$

We do not suppose that $\mu_{1} \leq \mu_{2}$.
Proposition 4.1. - Suppose that $Q$ is given by (4.1). Then

$$
\left(a_{1}, a_{2}, a_{3}\right)= \begin{cases}\left(\mu_{1}, \mu_{2}, \mu_{2}\right), & \text { if } \mu_{1} \leq \mu_{2} \\ \left(\mu_{2}, \mu_{2}, \mu_{1}\right), & \text { if } \mu_{1}>\mu_{2} .\end{cases}
$$

Proof. - Lemma 1.3.(b) implies that $a_{1}=\min \left(\mu_{1}, \mu_{2}\right)$. We distinguish two cases.
Suppose that $\mu_{1} \leq \mu_{2}$. Then $a_{1}=\mu_{1}$ and $\operatorname{ord}(\Delta)=\mu_{1}+2 \mu_{2} \geq a_{1}+a_{2}+a_{3}$ (Lemma 1.3.(a)). Therefore $a_{2} \leq\left(a_{2}+a_{3}\right) / 2 \leq \mu_{2}$. The existence of a basis $\boldsymbol{\psi}$ as in (4.1) implies that $\left(\mu_{1}, \mu_{2}, \mu_{2}\right) \in S(\boldsymbol{\psi})$. We conclude that $a_{2}=a_{3}=\mu_{2}$.

Suppose that $\mu_{1}>\mu_{2}$. In this case we have that $a_{1}=\mu_{2}$. Recall that we defined $\rho$ as the minimum of the valuation of the determinant of the $2 \times 2$-minors of $B$. One computes that $\rho=\min \left(2 \mu_{2}, 1+\mu_{1}+\mu_{2}\right)=2 \mu_{2}$, since we assumed that $\mu_{1} \geq$ $\mu_{2}+1$. Lemma 1.3.(c) implies that $\rho \geq a_{1}+a_{2}$, hence $a_{2} \leq \mu_{2}$. The existence of a basis $\boldsymbol{\psi}$ as in (4.1) implies that $\left(\mu_{2}, \mu_{2}, \mu_{1}\right) \in S(\boldsymbol{\psi})$. We conclude that $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{2}, \mu_{2}, \mu_{1}\right)$.

We now consider diagonalizable quadratic forms $Q$. Contrary to the situation for $\ell \neq 2$, a basis $\boldsymbol{\psi}$ which diagonalizes $Q$ is not optimal (Definition 1.2).

Proposition 4.2. - Suppose that $Q$ is diagonalizable. Let $\boldsymbol{\psi}$ be a basis of $L$ such that

$$
\begin{equation*}
Q(x)=b_{1} x_{1}^{2}+b_{2} x_{2}^{2}+b_{3} x_{3}^{2}, \text { with } b_{i}=u_{i} 2^{\mu_{i}}, u_{i} \in \mathbb{Z}_{2}^{\times} \text {and } \mu_{1} \leq \mu_{2} \leq \mu_{3} . \tag{4.2}
\end{equation*}
$$

(a)Suppose that $\mu_{1} \not \equiv \mu_{2} \bmod 2$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right)$.
(b) Suppose that $\mu_{1} \equiv \mu_{2} \bmod 2$.
(i) If $u_{1}+u_{2} \equiv 2 \bmod 4$ or $\mu_{3} \leq \mu_{2}+1$, then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$.
(ii) Otherwise, $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+2, \mu_{3}\right)$.

The proof of this proposition is divided in several lemmas. We use the notation of Proposition 4.2. In particular, $\boldsymbol{\psi}$ is a basis of $L$ with respect to which $Q$ is as in (4.2). Let $\boldsymbol{\varphi}$ be an optimal basis, i.e., suppose that the inequalities (1.2) hold. We write $C=\left(c_{i j}\right)$ for the change of basis matrix expressing $\boldsymbol{\varphi}$ in terms of $\boldsymbol{\psi}$. We write the quadratic form $Q$ in terms of the basis $\varphi$ as $Q(x)=\sum_{i \leq j} d_{i j} x_{i} x_{j}$. In other words, the $d_{i j}$ are the coefficients of the matrix obtained by dividing the diagonal elements of $C^{t} B C$ by two. One computes that

$$
\begin{equation*}
d_{i i}=c_{1 i}^{2} b_{1}+c_{2 i}^{2} b_{2}+c_{3 i}^{2} b_{3} \tag{4.3}
\end{equation*}
$$

Lemma 4.3. - Suppose that $Q$ is diagonal and $\mu_{1} \not \equiv \mu_{2} \bmod 2$. Then $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right)$.

Proof. - We have already seen that $a_{1}=\mu_{1}$. Therefore it follows from the definition of the $a_{i}$ that $a_{2} \geq \mu_{2}$. We claim that $a_{2}=\mu_{2}$. Suppose that $a_{2}>\mu_{2}$.

Write $\mu_{2}=\mu_{1}+2 \gamma+1$. The inequalities (1.2) imply that ord $\left(d_{22}\right) \geq a_{2} \geq \mu_{2}+1$ and $\operatorname{ord}\left(d_{33}\right) \geq a_{3} \geq a_{2} \geq \mu_{2}+1$. Since $\mu_{1} \not \equiv \mu_{2} \bmod 2$, it follows from (4.3) that $\operatorname{ord}\left(c_{12}\right) \geq \gamma+1$ and $\operatorname{ord}\left(c_{13}\right) \geq \gamma+1$.

We first suppose that $\mu_{3}>\mu_{2}$. Then ord $\left(c_{22}\right) \geq 1$ and $\operatorname{ord}\left(c_{33}\right) \geq 1$. But this implies that $\operatorname{det}(C) \equiv 0 \bmod 2$. This gives a contradiction.

If $\mu_{2}=\mu_{3}$, we proceed similarly. In this case $c_{22} \equiv c_{32} \bmod 2$ and $c_{23} \equiv c_{33} \bmod 2$. This implies again that $\operatorname{det}(C) \equiv 0 \bmod 2$. We conclude that $a_{2}=\mu_{2}$.

Since $\operatorname{ord}(\Delta)=\operatorname{ord}(\operatorname{det}(B))+2=\mu_{1}+\mu_{2}+\mu_{3}+2$, it follows from Lemma 1.3.(a) that $a_{3} \leq \mu_{3}+2$. To show that $a_{3}=\mu_{3}+2$ it suffices to find a basis $\varphi$ such that $\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right) \in S(\varphi)$. We now construct such a basis.

Our assumptions imply that $\mu_{3}$ is congruent to $\mu_{1}$ or $\mu_{2}$ (modulo 2). We suppose that $\mu_{3} \equiv \mu_{1} \bmod 2$. (The case $\mu_{3} \equiv \mu_{2} \bmod 2$ is similar.) Write $\mu_{2}=\mu_{1}+2 \gamma+1$ and $\mu_{3}=\mu_{1}+2 \lambda$. We distinguish two cases:
$-u_{1}+u_{3} \equiv 0 \bmod 4$,
$-u_{1}+u_{3} \equiv 2 \bmod 4$.
In the first case define

$$
C=\left(\begin{array}{ccc}
1 & 0 & 2^{\lambda} \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+b_{2} x_{2}^{2}+2^{\lambda+1} b_{1} x_{1} x_{3}+\left(b_{3}+2^{2 \lambda} b_{1}\right) x_{3}^{2}$.
In the second case we define

$$
C=\left(\begin{array}{ccc}
1 & 0 & 2^{\lambda} \\
0 & 1 & 2^{\lambda-\gamma} \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+b_{2} x_{2}^{2}+2^{\lambda+1} b_{1} x_{1} x_{3}+\left(b_{3}+2^{2 \lambda} b_{1}+\right.$ $\left.2^{2(\lambda-\gamma)} b_{2}\right) x_{3}^{2}+2^{\lambda-\gamma+1} b_{2} x_{2} x_{3}$. It is easy to check that the basis $\varphi$ corresponding to $C$ $\operatorname{satisfies}(1.2)$ for $a_{1}=\mu_{1}, a_{2}=\mu_{2}$ and $a_{3}=\mu_{3}+2$. This proves the lemma.

The proof of Lemmas 4.4, 4.5 and 4.6 follows the same pattern as the proof of Lemma 4.3.

Lemma 4.4. - Suppose that $Q$ is diagonalizable, $\mu_{1} \equiv \mu_{2} \bmod 2$ and $\mu_{3} \leq \mu_{2}+1$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$.

Proof. - Since $a_{1}=\mu_{1}$ and $\operatorname{ord}(\Delta)=\mu_{1}+\mu_{2}+\mu_{3}+2$ it follows from Lemma 1.3 that $a_{1}+2 a_{2} \leq a_{1}+a_{2}+a_{3} \leq \mu_{1}+\mu_{2}+\mu_{3}+2 \leq \mu_{1}+2 \mu_{2}+3$. This implies that $a_{2} \leq \mu_{2}+1$.

We now construct a basis $\boldsymbol{\varphi}$ such that $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right) \in S(\boldsymbol{\varphi})$. The lemma follows from this. Let $C$ be the corresponding change of basis matrix. Write $\mu_{2}=\mu_{1}+2 \gamma$.

If $\mu_{2}=\mu_{3}$ define

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 2^{\gamma} \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+\left(2^{2 \gamma} b_{1}+b_{2}\right) x_{2}^{2}+2^{\gamma+1} b_{1}\left(x_{1} x_{2}+\right.$ $\left.x_{1} x_{3}\right)+\left(b_{3}+2^{2 \gamma} b_{1}\right) x_{3}^{2}+2^{1+2 \gamma} b_{1} x_{2} x_{3}$.

If $\mu_{3}=\mu_{2}+1$ and $u_{1}+u_{2} \equiv 2 \bmod 4$ define

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 2^{\gamma} \\
0 & 1 & 1 \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+\left(b_{2}+2^{2 \gamma} b_{1}\right) x_{2}^{2}+2^{\gamma+1} b_{1}\left(x_{1} x_{2}+\right.$ $\left.x_{1} x_{3}\right)+\left(b_{3}+2^{2 \gamma} b_{1}+b_{2}\right) x_{3}^{2}+\left(2^{2 \gamma+1} b_{1}+2 b_{2}\right) x_{2} x_{3}$.

If $\mu_{3}=\mu_{2}+1$ and $u_{1}+u_{2} \equiv 0 \bmod 4$ define

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 2^{\gamma} \\
0 & 1 & 1 \\
0 & 1 & 2
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+\left(2^{2 \gamma} b_{1}+b_{2}+b_{3}\right) x_{2}^{2}+2^{\gamma+1} b_{1}\left(x_{1} x_{2}+\right.$ $\left.x_{1} x_{3}\right)+\left(4 b_{3}+2^{2 \gamma} b_{1}+b_{2}\right) x_{3}^{2}+\left(2^{2 \gamma+1} b_{1}+2 b_{2}+4 b_{3}\right) x_{2} x_{3}$.

In each of these cases one checks that $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right) \in S(\boldsymbol{\varphi})$.
Lemma 4.5. - Suppose that $Q$ is diagonal, $\mu_{1} \equiv \mu_{2} \bmod 2$ and $u_{1}+u_{2} \equiv 2 \bmod 4$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$.

Proof. - By Lemma 4.4 we may assume that $\mu_{3} \geq \mu_{2}+2$. We claim that $a_{2} \leq \mu_{2}+1$. Suppose that $a_{2} \geq \mu_{2}+2$. As before, we suppose that $\varphi$ is an optimal basis. As before, we write $C=\left(c_{i j}\right)$ for the change of basis matrix and $D=C^{t} B C=\left(d_{i j}\right)$ for the matrix corresponding to the new basis. Write $\mu_{2}=\mu_{1}+2 \gamma$.

The assumption $a_{2} \geq \mu_{2}+2$ implies that ord $\left(d_{22}\right) \geq a_{2} \geq \mu_{2}+2$ and $\operatorname{ord}\left(d_{33}\right) \geq$ $a_{3} \geq a_{2} \geq \mu_{2}+2$. It follows from (4.3) that ord $\left(c_{12}\right) \geq \gamma$ and $\operatorname{ord}\left(c_{13}\right) \geq \gamma$. Suppose that $\operatorname{ord}\left(c_{12}\right)=\gamma$. Then $\operatorname{ord}\left(c_{22}\right)=1$ and $d_{22} \equiv 2^{\mu_{2}}\left(u_{1}+u_{2}\right) \not \equiv 0 \bmod 2^{\mu_{2}+2}$. This gives a contradiction. Similarly, we obtain a contradiction if ord $\left(c_{13}\right)=\gamma$. Therefore $\operatorname{ord}\left(c_{1 j}\right)>\gamma$ for $j=2,3$ and $d_{22} \equiv c_{22}^{2} b_{2} \bmod 2^{\mu_{2}+2}$. Since $\operatorname{ord}\left(d_{22}\right) \geq \mu_{2}+2$ and $\operatorname{ord}\left(b_{2}\right)=\mu_{2}$, we conclude that ord $\left(c_{22}\right)>0$. Similarly, $d_{33} \equiv c_{23}^{2} b_{2} \bmod 2^{\mu_{2}+2}$; this implies that $\operatorname{ord}\left(c_{23}\right)>0$. But then $\operatorname{det}(C) \equiv 0 \bmod 2$. This gives a contradiction. We conclude that $a_{2} \leq \mu_{2}+1$.

To prove the lemma, we construct a basis $\varphi$ such that $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right) \in S(\boldsymbol{\varphi})$. We distinguish two subcases:
$-\mu_{3} \equiv \mu_{1} \bmod 2$,
$-\mu_{3} \not \equiv \mu_{1} \bmod 2$.
Suppose that $\mu_{3} \equiv \mu_{1} \bmod 2$. Write $\mu_{2}=\mu_{1}+2 \gamma$ and $\mu_{3}=\mu_{1}+2 \lambda$. Let $\varphi$ be the basis of $L$ corresponding to the change of basis matrix

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 2^{\lambda} \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+\left(2^{2 \gamma} b_{1}+b_{2}\right) x_{2}^{2}+2^{\gamma+1} b_{1} x_{1} x_{2}+$ $2^{\lambda+1} b_{1} x_{1} x_{3}+\left(b_{3}+2^{2 \lambda} b_{1}\right) x_{3}^{2}+2^{\gamma+\lambda+1} b_{1} x_{2} x_{3}$.

Suppose that $\mu_{3} \not \equiv \mu_{1} \bmod 2$. Write $\mu_{2}=\mu_{1}+2 \gamma$ and $\mu_{3}=\mu_{1}+2 \lambda+1$. Let $\varphi$ be the basis of $L$ corresponding to the change of basis matrix

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 2^{\lambda} \\
0 & 1 & 2^{\lambda-\gamma} \\
0 & 0 & 1
\end{array}\right)
$$

With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+\left(2^{2 \gamma} b_{1}+b_{2}\right) x_{2}^{2}+2^{\gamma+1} b_{1} x_{1} x_{2}+$ $2^{\lambda+1} b_{1} x_{1} x_{3}+\left(b_{3}+2^{2 \lambda} b_{1}+2^{2(\lambda-\gamma)} b_{2}\right) x_{3}^{2}+\left(2^{\gamma+\lambda+1} b_{1}+2^{\lambda-\gamma+1} b_{2}\right) x_{2} x_{3}$.

In each of these cases one checks that $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right) \in S(\varphi)$.
Lemma 4.6. - Suppose that $Q$ is diagonal, $\mu_{1} \equiv \mu_{2} \bmod 2, \mu_{3} \geq \mu_{2}+2$ and $u_{1}+u_{2} \equiv$ $0 \bmod 4$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+2, \mu_{3}\right)$.

Proof. - Write $\mu_{2}=\mu_{1}+2 \gamma$. We already know that $a_{1}=\mu_{1}$. We claim that $a_{2} \leq \mu_{2}+2$. Suppose $a_{2} \geq \mu_{2}+3$. The same reasoning as in the beginning of the proof of Lemma 4.4 shows that we may assume that $\mu_{3} \geq \mu_{2}+4$. If $c_{22} \equiv c_{23} \equiv 0$ $\bmod 2$, we conclude as in the proof of Lemma 4.5 that $\operatorname{det}(C) \equiv 0 \bmod 2$. This gives a contradiction, hence either $c_{22}$ or $c_{23}$ is a unit.

Suppose that $c_{22}$ is a unit. (The argument in the case that $c_{23}$ is a unit is similar, and we omit it.) Then $\operatorname{ord}\left(c_{12}\right)=\gamma$. One computes that

$$
\begin{equation*}
d_{12} \equiv 2 c_{12} c_{11} b_{1}+2 c_{21} c_{22} b_{2} \bmod 2^{\mu_{2}+3} \tag{4.4}
\end{equation*}
$$

It follows from (1.2) that $2 \operatorname{ord}\left(d_{12}\right) \geq a_{1}+a_{2} \geq \mu_{1}+\mu_{2}+3=2 \mu_{1}+2 \gamma+3$. Hence

$$
\begin{equation*}
\operatorname{ord}\left(d_{12}\right) \geq \mu_{1}+\gamma+2 \tag{4.5}
\end{equation*}
$$

Recall that Lemma 1.3.(b) implies that ord $\left(d_{11}\right)=a_{1}$.
First suppose that $\mu_{1}<\mu_{2}$, that is $\gamma \neq 0$. Since $d_{11}$ has valuation $a_{1}, c_{11}$ is a unit. It follows from (4.4) that ord $\left(d_{12}\right)=\mu_{1}+\gamma+1$. This contradicts (4.5).

Now suppose that $\mu_{1}=\mu_{2}$. Since $d_{11} \equiv c_{12}^{2} b_{1}+c_{21}^{2} b_{2} \bmod 2^{\mu_{1}+1}$. Since $d_{11}$ has valuation $a_{1}=\mu_{1}$, it follows that either
(i) $c_{12} \equiv 1 \bmod 2$ and $c_{21} \equiv 0 \bmod 2$, or
(ii) $c_{12} \equiv 0 \bmod 2$ and $c_{21} \equiv 1 \bmod 2$.

Since $\operatorname{ord}\left(d_{12}\right) \geq \mu_{1}+2$, it follows from (4.4) that (i) holds and that $c_{11} \equiv 0 \bmod 2$. One computes that

$$
d_{23} \equiv 2 c_{12} c_{13} b_{1}+2 c_{22} c_{23} b_{2} \equiv 2 c_{13} b_{1}+2 c_{23} b_{2} \bmod 2^{\mu_{1}+2}
$$

since $c_{12}$ and $c_{22}$ are units. It follows that $c_{13} \equiv c_{23} \bmod 2$. But this implies that $\operatorname{det}(C) \equiv 0 \bmod 2$. (In case $u_{1}+u_{2} \equiv 4 \bmod 8$ one could alternatively argue as in the proof of Lemma 4.5.)

Let $\varphi$ be the basis of $L$ corresponding to the change of basis matrix

$$
C=\left(\begin{array}{ccc}
1 & 2^{\gamma} & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

Then $b_{22}(\boldsymbol{\varphi}) \equiv 0 \bmod 2^{\mu_{2}+2}$. With respect to the new basis we have $Q(x)=b_{1} x_{1}^{2}+$ $\left(2^{2 \gamma} b_{1}+b_{2}\right) x_{2}^{2}+2^{\gamma+1} b_{1} x_{1} x_{2}+b_{3} x_{3}^{2}$. Therefore $\left(\mu_{1}, \mu_{2}+1, \mu_{3}\right) \in S(\boldsymbol{\varphi})$. This proves the lemma.

The following proposition is an immediate consequence of the computation of the invariants $a_{i}$. It illustrates that the $a_{i}$ satisfy similar properties for $\ell=2$ and $\ell \neq 2$, which is not so clear from the definition.

Proposition 4.7. - Let $Q$ be a ternary quadratic form over $\mathbb{Z}_{\ell}$ for $\ell \geq 2$. Then

$$
\operatorname{ord}(\Delta)=a_{1}+a_{2}+a_{3} .
$$

Proof. - For $\ell \neq 2$ this is Proposition 2.6.(b). For $\ell=2$ the theorem follows from the Propositions 4.1 and 4.2.

In the rest of this section we define the Gross-Keating invariant $\epsilon$ for $\ell=2$ and show that it is well defined (compare to Lemma 2.8).

Definition 4.8. - Suppose that $a_{1} \equiv a_{2} \bmod 2$ and $a_{3}>a_{2}$. Let $\varphi$ be an optimal basis. We define $\epsilon=\epsilon(\boldsymbol{\varphi})$ by $\epsilon=1$ if the subspace of $L \otimes_{\mathbb{Z}_{2}} \mathbb{Q}_{2}$ spanned by $\varphi_{1}$ and $\varphi_{2}$ is isotropic, and $\epsilon=-1$, otherwise.

Proposition 4.9. - Suppose that $a_{1} \equiv a_{2} \bmod 2$ and $a_{3}>a_{2}$.
(a) The invariant $\epsilon$ does not depend on the choice of the basis.
(b) (i) If $Q$ is not diagonalizable we may write $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+\right.$ $\left.x_{2} x_{3}+v x_{3}^{2}\right)$ with $v \in\{0,1\}$ and $\mu_{1}>\mu_{2}$. In this case

$$
\epsilon=(-1)^{v}
$$

(ii) If $Q$ is diagonalizable we may write $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+u_{2} 2^{\mu_{2}} x_{2}^{2}+u_{3} 2^{\mu_{3}} x_{3}^{2}$ with $u_{1}+u_{2} \equiv 0 \bmod 4, \mu_{1} \equiv \mu_{2} \bmod 2$ and $\mu_{3} \geq \mu_{2}+2$. We have that

$$
\epsilon=(-1)^{\left(u_{1}+u_{2}\right) / 4}
$$

Proof. - The fact that one of the two cases of (b) holds follows immediately from Propositions 4.1 and 4.2 .

Suppose that $Q$ is not diagonalizable. Write $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+x_{2} x_{3}+\right.$ $\left.v x_{3}^{2}\right)$, as in the statement of the proposition, and let $\psi$ be the corresponding basis. Write $Q_{2}$ for the restriction of $Q$ to the sublattice spanned by the basis vectors $\psi_{2}, \psi_{3}$. Lemma 3.2 implies that $Q_{2}$ is isotropic if and only $v=0$. This implies that $\epsilon(\boldsymbol{\psi})=(-1)^{v}$.

We now show that $\epsilon$ is well defined in this case. It suffices to show that $\epsilon(\varphi)=$ $\epsilon(\boldsymbol{\psi})$ for optimal bases $\varphi$ and $\boldsymbol{\psi}$ with respect to which $Q$ is in a normal form as in Proposition 3.3. By assumption, $Q$ is not diagonalizable. (In fact, it follows from Proposition 4.2 that no quadratic form $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+x_{2} x_{3}+\right.$ $v x_{3}^{2}$ ) with $v \in\{0,1\}$ and $\mu_{1}>\mu_{2}$ is diagonalizable. Hence we could have dropped this assumption from the statement of the proposition.) Write $Q^{\prime}(x)=u_{1}^{\prime} 2^{\mu_{1}} x_{1}^{2}+$ $2^{\mu_{2}}\left(v^{\prime} x_{2}^{2}+x_{2} x_{3}+v^{\prime} x_{3}^{2}\right)$ for $Q$ expressed with respect to the basis $\varphi$. Since $\Delta(Q)=$ $\Delta\left(Q^{\prime}\right)$ we have that $u_{1}\left(4 v^{2}-1\right)=u_{1}^{\prime}\left(4\left(v^{\prime}\right)^{2}-1\right)$, therefore $v=v^{\prime}$ implies that $u_{1}=u_{1}^{\prime}$.

Hence, to show that $\epsilon(\boldsymbol{\varphi})=\epsilon(\boldsymbol{\psi})$, it suffices to show that $v=v^{\prime}$. We assume that $v=1$ and $v^{\prime}=0$, and derive a contradiction.

The basis vector $\varphi_{2}$ is isotropic. Write $\varphi_{2}=c_{1} \psi_{1}+c_{2} \psi_{2}+c_{3} \psi_{3}$. The fact that $Q\left(\varphi_{2}\right)=0$ implies that $\mu_{1} \equiv \mu_{2} \bmod 2$. Moreover, it follows that $\operatorname{ord}\left(c_{j}\right) \geq\left(\mu_{1}-\right.$ $\left.\mu_{2}\right) / 2>0$ for $j=2,3$. Since $\varphi_{2}$ is primitive, it follows that $c_{1} \equiv 1 \bmod 2$. An easy computation shows that ord $\left(\varphi_{2}, \psi_{i}\right)>\mu_{2}$ for $i=1,2,3$. In particular ord $\left(\varphi_{2}, \varphi_{3}\right)>$ $\mu_{2}$. But this contradicts the assumption that ord $\left(\varphi_{2}, \varphi_{3}\right)=\mu_{2}$.

Next we assume that $Q$ is diagonalizable, and let $Q(x)$ be as in the statement of (b.ii). Write $\boldsymbol{\psi}$ for the corresponding basis of $L$. Let $Q_{2}$ be the restriction of $Q$ to the subspace spanned by $\psi_{1}, \psi_{2}$. Then $Q_{2}$ is isotropic if and only if $-\operatorname{det}(Q)$ is a square ( $[\mathbf{S}$, Theorem IV.6]). It is easy to see that this happens if and only if $u_{1}+u_{2} \equiv 0 \bmod 8$.

We now show that $\epsilon$ is independent of the choice of the optimal basis in this case. Let $\varphi$ be an optimal basis. Let $C=\left(c_{i j}\right)$ be the corresponding change of basis matrix expressing $\varphi$ in terms of $\boldsymbol{\psi}$. Write $\mu_{1}=\mu_{2}+2 \gamma$.

We suppose that $\mu_{2}>\mu_{1}$, that is $\gamma>0$. (The case $\mu_{1}=\mu_{2}$ is analogous and left to the reader.) We use the notation of the proof of Lemma 4.6. In particular, we write $Q(x)=\sum_{i \leq j} d_{i j} x_{i} x_{j}$ for the representation of $Q$ in terms of the basis $\varphi$.

We showed in the proof of Lemma 4.6 that either $c_{22}$ or $c_{23}$ is a unit. Suppose that $c_{22} \equiv 0 \bmod 2$ and $c_{23} \equiv 1 \bmod 2$. It follows that $\operatorname{ord}\left(d_{33}\right) \geq a_{3}=\mu_{3} \geq \mu_{2}+3$. Therefore (4.3) implies that $\operatorname{ord}\left(c_{13}\right)=\gamma$. We showed in the proof of Lemma 4.6 that $c_{11}$ is a unit. Since $d_{13} \equiv 2 c_{11} c_{13} b_{1}+2 c_{21} c_{23} b_{2} \bmod 2^{\mu_{3}+1}$, we conclude that $2 \operatorname{ord}\left(d_{13}\right)=2+2 \gamma+2 \mu_{1}=\mu_{1}+\mu_{2}+2$. (Here we use that $\gamma>0$.) But this contradicts $2 \operatorname{ord}\left(d_{13}\right) \geq a_{1}+a_{3}=\mu_{1}+\mu_{3} \geq \mu_{1}+\mu_{2}+3$. We conclude that $c_{22}$ is a unit. Recall from the proof of Lemma 4.6 that this implies that $c_{12} \equiv 1 \bmod 2$ and $c_{21} \equiv 0 \bmod 2$. Therefore the determinant of the submatrix

$$
\tilde{C}=\left(\begin{array}{ll}
c_{11} & c_{12} \\
c_{21} & c_{22}
\end{array}\right)
$$

of $C$ is a unit. We may define

$$
D=\left(\begin{array}{cc}
\tilde{C}^{-1} & 0 \\
0 & 1
\end{array}\right)
$$

With respect to the basis corresponding to $C D$, the quadratic form $Q$ becomes $Q(x)=$ $\left(b_{1}+\delta_{1}^{2} b_{3}\right) x_{1}^{2}+\left(b_{2}+\delta_{2}^{2} b_{3}\right) x_{2}^{2}+2 \delta_{1} b_{3} x_{1} x_{2}+x_{3}$ (other terms), for certain $\delta_{1}, \delta_{2} \in \mathbb{Z}_{2}$. Since $\operatorname{ord}\left(b_{3}\right) \geq \operatorname{ord}\left(b_{2}\right)+3$ this implies that the subspace spanned by $\varphi_{1}$ and $\varphi_{2}$ is isotropic if and only if the space spanned by $\psi_{1}$ and $\psi_{2}$ is isotropic.

## 5. Anisotropic quadratic forms

The goal is to classify all anisotropic ternary quadratic forms over $\mathbb{Z}_{2}$, starting from the normal form of Proposition 3.3. We will see that for anisotropic forms we may
choose an optimal basis $\varphi$ so that $\operatorname{ord}\left(Q\left(\varphi_{i}\right)\right)=a_{i}$, similar to what we had for $\ell \neq 2$ (Corollary 5.8).

Proposition 5.1. - Let $Q$ be a ternary quadratic form over $\mathbb{Q} \ell$. Write $Q(x)=b_{1} x_{1}^{2}+$ $b_{2} x_{2}^{2}+b_{3} x_{3}^{2}$. We denote by $\operatorname{det}(Q)=b_{1} b_{2} b_{3}$ the determinant of $Q$. Then $Q$ is isotropic if and only if

$$
(-1,-\operatorname{det}(Q))=\prod_{i<j}\left(b_{i}, b_{j}\right)
$$

Here $(\cdot, \cdot)$ denotes the Hilbert symbol.
Proof. - This is [S, Theorem IV.6.ii].
Proposition 5.2. - Let $Q$ be a ternary quadratic form over $\mathbb{Z}_{2}$ which is not diagonalizable. Let $\boldsymbol{\psi}$ be an optimal basis such that $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(v x_{2}^{2}+x_{2} x_{3}+v x_{3}^{2}\right)$ with $v \in\{0,1\}$. Then $Q$ is isotropic if and only if $v=0$ or $\mu_{1} \equiv \mu_{2} \bmod 2$.

Proof. - If $v=0$ then $Q$ is obviously isotropic. Therefore suppose that $v=1$. To decide whether $Q$ is isotropic, we may consider $Q$ as quadratic form over $\mathbb{Q}_{2}$. We have $Q(x) \sim_{\mathbb{Q}_{2}} u_{1} 2^{\mu_{1}} x_{1}^{2}+2^{\mu_{2}}\left(x_{2}^{2}+3 x_{3}^{2}\right)$. The proposition follows from Proposition 5.1 by direct verification using the formula for the Hilbert symbol [S, Theorem III.1].

Lemma 5.3. - Let $Q$ be a ternary quadratic form over $\mathbb{Z}_{\ell}$. We do not assume that $\ell=2$. Suppose that $a_{1} \equiv a_{2} \equiv a_{3} \bmod 2$. Then $Q$ is isotropic.

Proof. - If $Q$ is not diagonalizable then the lemma follows from Proposition 5.2, since $\left(a_{1}, a_{2}, a_{3}\right) \in\left\{\left(\mu_{1}, \mu_{2}, \mu_{2}\right),\left(\mu_{2}, \mu_{2}, \mu_{1}\right)\right\}$.

Suppose that $Q$ is diagonalizable. Write $Q(x)=u_{1} \ell^{\mu_{1}} x_{1}^{2}+u_{2} \ell^{\mu_{2}} x_{2}^{2}+u_{3} \ell^{\mu_{3}} x_{3}^{2}$. If $\ell \neq 2$ we have that $\mu_{i}=a_{i}$ hence $\mu_{1} \equiv \mu_{2} \equiv \mu_{3} \bmod 2$. To show that $Q$ is isotropic, it suffices to consider $Q$ over $\mathbb{Q}_{\ell}$. After multiplying the basis vectors by a suitable constant, we may assume that $\mu_{1}=\mu_{2}=\mu_{3}=0$. The lemma now follows immediately from Proposition 5.1, since the Hilbert symbol is trivial on units for $\ell \neq 2$.

Suppose that $\ell=2$ and $Q$ is diagonalizable. Proposition 4.2 implies that $\mu_{1} \equiv$ $\mu_{2} \equiv \mu_{3} \bmod 2$ and $u_{1}+u_{2} \equiv 0 \bmod 4$. As for $\ell \neq 2$, it is no restriction to suppose that $Q(x)=u_{1} x_{2}^{2}+u_{2} x_{2}^{2}+u_{3} x_{3}^{2}$. One computes that this quadratic form is anisotropic if and only if $u_{1} \equiv u_{2} \equiv u_{3} \bmod 4$. Hence in our case $Q$ is isotropic.

For future reference we record from the proof of Lemma 5.3 when a diagonal ternary form over $\mathbb{Z}_{2}$ is anisotropic.

Lemma 5.4. - Let $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+u_{2} 2^{\mu_{2}} x_{2}^{2}+u_{3} 2^{\mu_{3}} x_{3}^{2}$ be a diagonal, ternary quadratic form over $\mathbb{Z}_{2}$. Suppose that $\mu_{1} \equiv \mu_{2} \equiv \mu_{3} \bmod 2$. Then $Q$ is anisotropic if and only if $u_{1} \equiv u_{2} \equiv u_{3} \bmod 4$.

Lemma 5.5. - Let $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+u_{2} 2^{\mu_{2}} x_{2}^{2}+u_{3} 2^{\mu_{3}} x_{3}^{2}$ be a diagonal, ternary quadratic form over $\mathbb{Z}_{2}$. Suppose that $\mu_{1} \equiv \mu_{2} \bmod 2$ and $\mu_{3} \not \equiv \mu_{1} \bmod 2$.
(a) Suppose that $u_{1} \equiv u_{2} \equiv u_{3} \bmod 4$. Then $Q$ is anisotropic if and only if $u_{2} \equiv \pm u_{1}$ $\bmod 8$.
(b) Suppose that the $u_{i}$ are not all equivalent modulo 4. Then $Q$ is anisotropic if and only if $u_{2} \equiv \pm 3 u_{1} \bmod 8$.

Proof. - The proof is similar to the proof of Lemma 5.3 and is left to the reader.
Notation 5.6. - Let $Q$ be a ternary quadratic form with Gross-Keating invariants $\left(a_{1}, a_{2}, a_{3}\right)$. For every $1 \leq i<j \leq 3$ we define

$$
\delta_{i j}=\left\lceil\frac{a_{i}+a_{j}}{2}\right\rceil,
$$

where $\lceil a\rceil$ is the smallest integer greater than or equal to $a$.
Theorem 5.7. - Let $Q(x)=u_{1} 2^{\mu_{1}} x_{1}^{2}+u_{2} 2^{\mu_{2}} x_{2}^{2}+u_{3} 2^{\mu_{3}} x_{3}^{2}$ be a diagonal anisotropic quadratic form over $\mathbb{Z}_{2}$ with $\mu_{1} \leq \mu_{2} \leq \mu_{3}$. Then one of the following cases occurs.
(a) Suppose $\mu_{1} \equiv \mu_{3} \not \equiv \mu_{2} \bmod 2$ and $u_{1} \equiv 3 u_{3} \bmod 8$. Then $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right)$ and $a_{1} \not \equiv a_{2} \bmod 2$. There exists an optimal basis with respect to which

$$
Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} u_{2} x_{2}^{2}+2^{\delta_{13}} u_{1} x_{1} x_{3}+2^{a_{3}} u_{1} x_{3}^{2}
$$

(b) Suppose $\mu_{1} \equiv \mu_{3} \not \equiv \mu_{2} \bmod 2$ and $u_{1} \equiv u_{3} \bmod 4$. Then $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right)$ and $a_{1} \not \equiv a_{2} \bmod 2$. Moreover, $u_{2} \equiv u_{1} \bmod 4$ if $u_{3} \equiv u_{1} \bmod 8$ and $u_{2} \equiv-u_{1} \bmod 4$ if $u_{3} \equiv 5 u_{1} \bmod 8$. There exists an optimal basis with respect to which

$$
Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} u_{2} x_{2}^{2}+2^{\delta_{13}} u_{1} x_{1} x_{3}+2^{\delta_{23}} u_{2} x_{2} x_{3}+2^{a_{3}} u_{1} v x_{3}^{2} .
$$

Here $v=\left(u_{1}+u_{2}\right) / 2$ if $u_{2} \equiv u_{1} \bmod 4$ and $v=\left(3 u_{1}+u_{2}\right) / 2$ if $u_{2} \equiv-u_{1} \bmod 4$.
(c) Suppose $\mu_{1} \not \equiv \mu_{2} \equiv \mu_{3} \bmod 2$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}, \mu_{3}+2\right)$ and $a_{2} \not \equiv a_{1}$ $\bmod 2$. The quadratic form with respect to an optimal basis is as in (a) and (b) with the role of $x_{1}$ and $x_{2}$ reversed.
(d) Suppose $\mu_{1} \equiv \mu_{2} \bmod 2$ and $\mu_{2}=\mu_{3}$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$ and $a_{1} \not \equiv a_{2} \bmod 2$. Moreover, $u_{1} \equiv u_{2} \equiv u_{3} \bmod 4$. There exists an optimal basis with respect to which
$Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{13}} u_{1}\left(x_{1} x_{2}+x_{1} x_{3}\right)+2^{\delta_{23}} u_{1} x_{2} x_{3}+2^{a_{3}} v_{3} x_{3}^{2}$.
Here $v_{i}=\left(u_{1}+u_{i}\right) / 2$ for $i=2,3$.
(e) Suppose $\mu_{1} \equiv \mu_{2} \bmod 2$, $\mu_{3}=\mu_{2}+1$ and $u_{1} \equiv u_{2} \bmod 4$. Then $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$ and $a_{2} \not \equiv a_{1} \bmod 2$. Moreover, $u_{2} \equiv u_{1} \bmod 8$ if $u_{3} \equiv u_{1}$ $\bmod 4$ and $u_{2} \equiv 5 u_{1} \bmod 8$ if $u_{3} \equiv-u_{1} \bmod 4$. There exists an optimal basis with respect to which
$Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{13}} u_{1}\left(x_{1} x_{2}+x_{1} x_{3}\right)+2^{\delta_{23}} v_{2} x_{2} x_{3}+2^{a_{3}} v_{3} x_{3}^{2}$.
Here $v_{2}=\left(u_{1}+u_{2}\right) / 2$ and $v_{3}=\left(u_{1}+u_{3}\right) / 2\left(\right.$ resp. $\left.\left(3 u_{1}+u_{3}\right) / 2\right)$ depending on whether $u_{3} \equiv u_{1} \bmod 4$ or not.
(f) Suppose $\mu_{1} \equiv \mu_{2} \bmod 2, \mu_{3}=\mu_{2}+1$ and $u_{1} \equiv-u_{2} \bmod 4$. Then $\left(a_{1}, a_{2}, a_{3}\right)=$ $\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$ and $a_{1} \equiv a_{2} \bmod 2$. Moreover, $u_{2} \equiv 3 u_{1} \bmod 8$. There exists an optimal basis with respect to which
$Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{13}} u_{1}\left(x_{1} x_{2}+x_{1} x_{3}\right)+2^{\delta_{23}} v_{23} x_{2} x_{3}+2^{a_{3}} v_{3} x_{3}^{2}$.
Here $v_{2}=\left(u_{1}+u_{2}+2 u_{3}\right) / 2, v_{23}=\left(u_{1}+u_{2}+4 u_{3}\right) / 2$ and $v_{3}=u_{1}+2 u_{3}$.
(g) Suppose $\mu_{1} \equiv \mu_{2} \equiv \mu_{3} \bmod 2$ and $u_{1} \equiv u_{2} \bmod 4$ and $\mu_{3} \geq \mu_{2}+2$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$ and $a_{2} \not \equiv a_{1} \bmod 2$. Moreover, $u_{3} \equiv u_{1} \bmod 4$. There exists an optimal basis with respect to which

$$
Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{12}} u_{1} x_{1} x_{2}+2^{\delta_{13}} u_{1} x_{1} x_{3}+2^{\delta_{23}} u_{1} x_{2} x_{3}+2^{a_{3}} v_{3} x_{3}^{2}
$$

Here $v_{i}=\left(u_{1}+u_{i}\right) / 2$ for $i=2,3$.
(h) Suppose $\mu_{1} \equiv \mu_{2} \not \equiv \mu_{3} \bmod 2$ and $u_{1} \equiv u_{2} \bmod 4$ and $\mu_{3} \geq \mu_{2}+2$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+1, \mu_{3}+1\right)$ and $a_{2} \not \equiv a_{1} \bmod 2$. One of the following two cases holds:

$$
\left\{\begin{array}{l}
u_{2} \equiv u_{1} \bmod 8 \text { and } u_{3} \equiv u_{1} \bmod 4 \\
u_{2} \equiv 5 u_{1} \bmod 8 \text { and } u_{3} \equiv-u_{1} \bmod 4 .
\end{array}\right.
$$

There exists an optimal basis with respect to which

$$
Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{12}} u_{1} x_{1} x_{2}+2^{\delta_{13}} u_{1} x_{1} x_{3}+2^{\delta_{23}} v_{2} x_{2} x_{3}+2^{a_{3}} v_{3} x_{3}^{2}
$$

Here $v_{2}=\left(u_{1}+u_{2}\right) / 2$ and $v_{3}=\left(u_{1}+u_{3}\right) / 2$ (resp. $\left.v_{3}=\left(3 u_{1}+u_{3}\right) / 2\right)$ depending on whether $u_{1} \equiv u_{3} \bmod 4$ or not.
(i) Suppose $\mu_{1} \equiv \mu_{2} \not \equiv \mu_{3} \bmod 2, \mu_{3} \geq \mu_{2}+2$ and $u_{2} \equiv 3 u_{1} \bmod 8$. Then $\left(a_{1}, a_{2}, a_{3}\right)=\left(\mu_{1}, \mu_{2}+2, \mu_{3}\right)$ and $a_{1} \equiv a_{2} \bmod 2$. There exists an optimal basis with respect to which

$$
Q(x)=2^{a_{1}} u_{1} x_{1}^{2}+2^{a_{2}} v_{2} x_{2}^{2}+2^{\delta_{12}} u_{1} x_{1} x_{2}+2^{a_{3}} u_{3} x_{3}^{2}
$$

Here $v_{2}=\left(u_{1}+u_{2}\right) / 2$.
Proof. - This follows from the results of Section 4 together with the Lemmas 5.4, 5.5.

Corollary 5.8. - Suppose that $Q$ is anisotropic. Then there exists an optimal basis $\varphi$ such that

$$
\operatorname{ord}\left(b_{i i}(\boldsymbol{\varphi})\right)=a_{i}
$$

for $i=1,2,3$.
Proof. - This follows immediately from Theorem 5.7 (diagonal case) and Proposition 5.2 (non-diagonal case).

In Section 6, we give a more conceptual proof of Corollary 5.8. In fact, we prove that any optimal basis has the property in Corollary 5.8. The following lemma gives a list of the small cases.

Lemma 5.9. - Let $Q$ be an anisotropic ternary quadratic form over $\mathbb{Z}_{2}$ and suppose that $a_{3} \leq 1$. Then one of the following possibilities occurs.
(a) We have $\left(a_{1}, a_{2}, a_{3}\right)=(0,0,1)$. In this case $Q$ is not diagonalizable; it is of the form

$$
Q(x)=x_{1}^{2}+x_{1} x_{2}+x_{2}^{2}+u_{3} 2 x_{3}^{2} .
$$

(b) We have $\left(a_{1}, a_{2}, a_{3}\right)=(0,1,1)$ and $Q$ is not diagonalizable. Then $Q$ is of the form

$$
Q(x)=u_{1} x_{1}^{2}+2\left(x_{2}^{2}+x_{2} x_{3}+x_{3}^{2}\right) .
$$

(c) We have $\left(a_{1}, a_{2}, a_{3}\right)=(0,1,1)$ and $Q$ is diagonalizable. Then $Q$ is as in Theorem 5.7.(d) with $a_{1}=\delta_{13}=0$ and $a_{2}=a_{3}=\delta_{23}=1$.

## 6. Alternative version of the Gross-Keating invariants for anisotropic forms

We fix an arbitrary prime number $\ell$ and a free quadratic module $(L, Q)$ over $\mathbb{Z}_{\ell}$ of rank $n$. We assume that $(L, Q)$ is anisotropic, i.e., that $Q(\psi)=0$ implies $\psi=$ 0 . Under this assumption, there is an alternative definition of the Gross-Keating invariants and a very useful characterization of optimal bases; see the remark at the end of section 4 in $[\mathbf{G K}]$. In this section we do not suppose that $n=3$ to streamline some arguments. Recall that $n \geq 5$ implies that $(L, Q)$ is isotropic ([S, Theorem IV.6]). Therefore the only additional case is anisotropic quadratic forms in four variables.

We define a function $v: L \rightarrow \mathbb{Z} \cup\{\infty\}$ by the rule

$$
v(\psi):=\operatorname{ord}_{\ell} Q(\psi) .
$$

For $\psi \in L$ and $x \in \mathbb{Z}_{p}$ we have

$$
\begin{equation*}
v(x \psi)=2 \operatorname{ord}_{\ell}(x)+v(\psi) . \tag{6.1}
\end{equation*}
$$

Lemma 6.1. - The function $v$ satisfies the triangle inequality

$$
\begin{equation*}
v\left(\psi+\psi^{\prime}\right) \geq \min \left(v(\psi), v\left(\psi^{\prime}\right)\right) \tag{6.2}
\end{equation*}
$$

Moreover, if the inequality in (6.2) is strict we have $v(\psi)=v\left(\psi^{\prime}\right)$.

Proof. - If $\psi$ and $\psi^{\prime}$ are linearly dependent the claim is obvious. We may hence assume that they are linearly independent. For $x, y \in \mathbb{Z}_{\ell}$ we write

$$
Q\left(x \psi+y \psi^{\prime}\right)=a x^{2}+y^{2} b+c x y .
$$

Suppose that $v\left(\psi+\psi^{\prime}\right)<v(\psi), v\left(\psi^{\prime}\right)$. Then $\operatorname{ord}_{\ell}(a+b+c)<\operatorname{ord}_{\ell}(a), \operatorname{ord}_{\ell}(b)$. The usual triangle inequality for ord ${ }_{\ell}$ implies

$$
\operatorname{ord}_{\ell}(c)=\operatorname{ord}_{\ell}(a+b+c)<\operatorname{ord}_{\ell}(a), \operatorname{ord}_{\ell}(b) .
$$

Lemma 3.2.(b) implies that $(L, Q)$ is isotropic. This and proves (6.2). The second assertion of the lemma follows from (6.2), applied to a suitable combination of the vectors $\pm \psi, \pm \psi^{\prime}$ and $\psi+\psi^{\prime}$.

Remark 6.2. - If $n \leq 3$, one gets an alternative proof of Lemma 6.1 by noting that $(L, Q)$ is represented by the quaternion division algebra $D$ over $\mathbb{Q}_{\ell}$, equipped with its norm form. The function $v$ is then the restriction of the standard valuation of $D$.

Let $\boldsymbol{\psi}=\left(\psi_{i}\right)$ be a basis of $L$. For $i=1, \ldots, n$, let $L_{i-1} \subset L$ be the subspace (of rank $i-1$ ) spanned by $\psi_{1}, \ldots, \psi_{i-1}$. We define a function $\tilde{v}_{i}: L / L_{i-1} \rightarrow \mathbb{Z}_{\geq 0} \cup\{\infty\}$ by the rule

$$
\tilde{v}_{i}\left(\psi+L_{i-1}\right):=\max \left(v\left(\psi^{\prime}\right) \mid \psi^{\prime} \in \psi+L_{i-1}\right) .
$$

Note that $\tilde{v}_{i}(\psi)=\infty$ if and only of $\psi \in L_{i-1}$.
Definition 6.3. - A basis $\boldsymbol{\psi}=\left(\psi_{i}\right)$ of $L$ is called ideal, if

$$
v\left(\psi_{i}\right)=\tilde{v}_{i}\left(\psi_{i}+L_{i-1}\right)=\min _{\psi \in L}\left(\tilde{v}_{i}\left(\psi+L_{i-1}\right)\right)
$$

holds for $i=1, \ldots, n$.
It is clear that there exists an ideal basis of $L$. The next lemma gives a useful characterization of an ideal basis.

Lemma 6.4. - $A$ basis $\boldsymbol{\psi}=\left(\psi_{i}\right)$ of $L$ is ideal if and only if

$$
\begin{equation*}
v\left(\psi_{i}\right) \leq v\left(\psi_{j}\right) \quad \text { for } i \leq j \tag{6.3}
\end{equation*}
$$

and for all $\left(x_{i}\right) \in \mathbb{Z}_{\ell}^{n}$ we have

$$
\begin{equation*}
v\left(\sum_{i} x_{i} \psi_{i}\right)=\min _{i} v\left(x_{i} \psi_{i}\right) . \tag{6.4}
\end{equation*}
$$

Proof. - Let $\boldsymbol{\psi}=\left(\psi_{i}\right)$ be a basis of $L$. If (6.3) and (6.4) hold, then one easily checks from Definition 6.3 that $\boldsymbol{\psi}$ is ideal.

Conversely, suppose that $\boldsymbol{\psi}$ is ideal. The inequality (6.3) follows directly from Definition 6.3. It remains to prove (6.4). Fix $\left(x_{i}\right) \in \mathbb{Z}_{\ell}^{n}$ and $k$ with $1 \leq k \leq n$. Set $\varphi_{k}:=\sum_{i<k} x_{i} \psi_{i}$. We claim that

$$
\begin{equation*}
v\left(\varphi_{k}+x_{k} \psi_{k}\right)=\min \left(v\left(\varphi_{k}\right), v\left(x_{k} \psi_{k}\right)\right) \tag{6.5}
\end{equation*}
$$

From this claim, (6.4) follows by induction.

For $k=1$, the claim is obvious. To prove it for $k>1$ we may assume that it holds for $k^{\prime}=k-1$. Also, by the triangle inequality (6.2), the left hand side of (6.5) is greater than or equal to the right hand side. Suppose that the left hand side is strictly greater than the right hand side. Then we have $v\left(\varphi_{k}\right)=v\left(x_{k} \psi_{k}\right)$. Using (6.1), (6.3) and the claim for $k^{\prime}=k-1$, we find that $\operatorname{ord}_{\ell}\left(x_{k}\right) \leq \operatorname{ord}_{\ell}\left(x_{i}\right)$ for all $i \leq k$. After dividing by $x_{k}$, we may therefore assume that $x_{k}=1$. However, by the definition of an ideal basis we have

$$
v\left(\varphi_{k}\right)=v\left(\psi_{k}\right) \geq v\left(\varphi_{k}+\psi_{k}\right)
$$

This contradicts our assumption and proves the claim.
Let us fix an ideal basis $\boldsymbol{\psi}=\left(\psi_{1}, \ldots, \psi_{n}\right)$ of $L$, and set

$$
a_{i}:=v\left(\psi_{i}\right), \quad i=1, \ldots, n
$$

We want to show that the $a_{i}$ are the Gross-Keating invariants of $(L, Q)$. We first check that $\left(a_{i}\right)$ lies in the set $S$ (Section 1). For this we write the quadratic form $Q$ as follows:

$$
Q\left(\sum_{i} x_{i} \psi_{i}\right)=\sum_{i \leq j} b_{i j} x_{i} x_{j}
$$

We set $a_{i j}:=\operatorname{ord}_{\ell}\left(b_{i j}\right)$. Note that $a_{i}=a_{i i}$.
Proposition 6.5. - For $1 \leq i \leq j \leq n$ we have

$$
a_{i j} \geq \frac{a_{i}+a_{j}}{2}
$$

Proof. - The case $i=j$ being trivial, we may assume that $i<j$. Our proof is by contradiction. First we assume that $2 a_{i j}+1<a_{i}+a_{j}$. We set $c:=\max \left(a_{i j}-a_{i}+1,0\right)$ and look at the right hand side of

$$
Q\left(\ell^{c} \psi_{i}+\psi_{j}\right)=b_{i i} \ell^{2 c}+b_{j j}+b_{i j} \ell^{c}
$$

The three terms of this sum have $\ell$-valuation $a_{i}+2 c, a_{j}$ and $a_{i j}+c$, respectively. By our choice of $c$ we have

$$
a_{i j}+c<\min \left(a_{i}+2 c, a_{j}\right)
$$

It follows that

$$
v\left(\ell^{c} \psi_{i}+\psi_{j}\right)=a_{i j}+c<\min \left(v\left(\ell^{c} \psi_{i}\right), v\left(\psi_{j}\right)\right)
$$

This contradicts the triangle inequality and excludes the case $2 a_{i j}+1<a_{i}+a_{j}$.
It remains to exclude the case $2 a_{i j}+1=a_{i}+a_{j}$. Since $a_{i} \leq a_{j}$ we have $c:=$ $a_{i j}-a_{i} \geq 0$. Let $x \in \mathbb{Z}_{\ell}^{\times}$be a $\ell$-adic unit. Then

$$
\begin{equation*}
Q\left(\ell^{c} x \psi_{i}+\psi_{j}\right)=b_{i i} \ell^{2 c} x^{2}+b_{j j}+b_{i j} \ell^{c} x \tag{6.6}
\end{equation*}
$$

By our choice of $c$ we have

$$
a_{i}+2 c=a_{j}-1=a_{i j}+c
$$

We see that on the right hand side of (6.6), the first and the last term have the minimal valuation $a_{j}-1$, while the middle term has valuation $a_{j}$. Therefore, for an appropriate choice of $x$, we get

$$
v\left(\ell^{c} x \psi_{i}+\psi_{j}\right) \geq a_{j}>\min \left(v\left(\ell^{c} x \psi_{i}\right), v\left(\psi_{j}\right)\right)
$$

But this contradicts Lemma 6.4, (6.4). The proposition follows.
Proposition 6.6. - An ideal basis is also optimal (Definition 1.2). Moreover, if $\boldsymbol{\psi}=$ $\left(\psi_{i}\right)$ is an ideal basis of L, then $\left(a_{i}:=v\left(\psi_{i}\right)\right)$ are the Gross-Keating invariants of $(L, Q)$.

Proof. - The previous proposition says that $\left(a_{i}\right)$ is an element of $S$. It remains to show that $\left(a_{i}\right)$ is a maximal element, with respect to the lexicographical ordering.

Let $\boldsymbol{\psi}^{\prime}=\left(\psi_{i}^{\prime}\right)$ be an arbitrary basis of $L$, and let $\left(a_{i}^{\prime}\right)$ be an element of $S\left(\boldsymbol{\psi}^{\prime}\right)$ (Section 1). We will show that $a_{k}^{\prime} \leq a_{k}$ for $k=1, \ldots, n$, which proves the proposition. Write

$$
\psi_{i}^{\prime}=\sum_{j} x_{i j} \psi_{j}, \quad \text { with } \quad\left(x_{i j}\right) \in \mathrm{GL}_{n}\left(\mathbb{Z}_{\ell}\right)
$$

The condition $\left(a_{i}^{\prime}\right) \in S\left(\boldsymbol{\psi}^{\prime}\right)$ together with Lemma 6.4 shows that

$$
\begin{equation*}
a_{i}^{\prime} \leq v\left(\psi_{i}^{\prime}\right)=\min _{j}\left(a_{j}+2 \operatorname{ord}_{\ell}\left(x_{i j}\right)\right) . \tag{6.7}
\end{equation*}
$$

Using that $\left(x_{i j}\right)$ is invertible, one shows that there exists at least one pair of indices (ij) with $k \leq i$ and $j \leq k$ such that $x_{i j}$ is a unit. Applying (6.7) and (6.3) we get

$$
a_{k}^{\prime} \leq a_{i}^{\prime} \leq a_{j} \leq a_{k}
$$

This is what we had to prove.
Corollary 6.7. - Let $\boldsymbol{\psi}=\left(\psi_{i}\right)$ be an ideal basis of $L$ and $\left(y_{i}\right) \in \mathbb{Q}_{\ell}^{n}$ with $y_{i} \neq 0$. Set $\psi^{\prime}:=\left(\psi_{i}^{\prime}\right)$, where $\psi_{i}^{\prime}:=y_{i} \psi_{i} \in L \otimes_{\mathbb{Z}_{\ell}} \mathbb{Q}_{\ell}$, and let $L^{\prime}$ denote the $\mathbb{Z}_{\ell}$-lattice spanned by $\boldsymbol{\psi}^{\prime}$. Let $\left(a_{i}\right)$ be the Gross-Keating invariants of $L$.
(a) The basis $\psi^{\prime}$ of $L^{\prime}$ is ideal.
(b) The Gross-Keating invariants of $L^{\prime}$ are the numbers

$$
a_{i}^{\prime}:=a_{i}+2 \operatorname{ord}_{\ell}\left(y_{i}\right),
$$

in some order.
Proof. - Choose an integer $r$ such that $\ell^{r} y_{i} \in \mathbb{Z}_{\ell}$, for all $i$. For $\left(x_{i}\right) \in \mathbb{Z}_{\ell}^{n}$, Lemma 6.4 shows that

$$
\begin{aligned}
v\left(\sum_{i} x_{i} \psi_{i}^{\prime}\right) & =v\left(\sum_{i} \ell^{r} x_{i} y_{i} \psi_{i}\right)-2 r \\
& =\min _{i}\left(v\left(\ell^{r} x_{i} y_{i} \psi_{i}\right)\right)-2 r \\
& =\min _{i}\left(v\left(x_{i} \psi_{i}^{\prime}\right)\right) .
\end{aligned}
$$

Again by Lemma 6.4 we conclude that $\boldsymbol{\psi}^{\prime}$ (in some order) is an ideal basis of $L^{\prime}$. This proves (a). Part (a) of the corollary follows now from the previous proposition.

Remark 6.8. - Corollary 6.7 (a) is false without the assumption that ( $L, Q$ ) is anisotropic. Consider, for instance, the (isotropic) quadratic form $Q(x)=x_{1}^{2}-x_{2}^{2}+4 x_{3}^{2}$ over $\mathbb{Z}_{2}$. Dividing the last vector of the standard basis by 2 we obtain the quadratic form $Q^{\prime}(x)=x_{1}^{2}-x_{2}^{2}+x_{3}^{2}$. According to Proposition 4.2(b), the Gross-Keating invariants of $Q$ are ( $0,2,2$ ), while the invariants of $Q^{\prime}$ are $(0,1,1)$.

Proposition 6.9. - Let $(L, Q)$ be an anisotropic free quadratic module over $\mathbb{Z}_{\ell}$. Then every optimal basis is an ideal basis.

The proof of this proposition uses the following lemma.
Lemma 6.10. - Let $\left(a_{1}, \ldots, a_{n}\right)$ be the Gross-Keating invariants of $(L, Q)$, and let $\boldsymbol{\psi}$ be an optimal basis. Then $v\left(\psi_{i}\right)=a_{i}$.

Proof. - Let $\boldsymbol{\psi}$ be an optimal basis and suppose that $v\left(\psi_{i}\right)>a_{i}$, for some $i$. It follows from the definition of the Gross-Keating invariants (Definition 1.2) that there exists a $j \neq i$ such that

$$
\operatorname{ord}\left(b_{i j}\right)=\left(a_{i}+a_{j}\right) / 2
$$

In particular, we have that $a_{i} \equiv a_{j} \bmod 2$. Lemma 5.3 implies therefore that $a_{k} \not \equiv$ $a_{i} \bmod 2$ for all $k \neq i, j$, since $(L, Q)$ is anisotropic. (The case that $n=4$ easily reduces to the case that $n=3$ by using the existence of an ideal basis.)

Consider the restriction $Q_{1}$ of $Q$ to $L_{1}=\left\langle\psi_{i}, \psi_{j}\right\rangle$. We distinguish three cases. First suppose that $a_{i}=a_{j}$. Then ( $\left.L_{1}, Q_{1}\right)$ is isotropic by Lemma 3.2.(b).

Next we suppose that $a_{i}<a_{j}$. Then $i<j$. We have already seen that $a_{k} \not \equiv$ $a_{i} \bmod 2$ for all $k \neq i, j$. Renumbering the indices, if necessary, we may assume that $a_{i}<a_{i+1}$ and $a_{j-1}<a_{j}$. Define $\left(\tilde{a}_{i}\right)$ by $\tilde{a}_{i}=a_{i}+1$ and $\tilde{a}_{j}=a_{j}-1$, and $\tilde{a}_{k}=a_{k}$ for all $k \neq i, j$. Then $\left(\tilde{a}_{k}\right) \in S(\psi)$. This contradicts the definition of the Gross-Keating invariants.

Finally, we suppose that $a_{i}>a_{j}$. Then $i>j$. If $v\left(\psi_{j}\right)>a_{j}$, we interchange $i$ and $j$ and obtain a contradiction by the previous case. Therefore $v\left(\psi_{j}\right)=a_{j}$. Since $a_{i} \equiv a_{j} \bmod 2$, Lemma 3.2.(b) implies that $L_{1}$ is isotropic. This gives a contradiction. We conclude that $v\left(\psi_{i}\right)=a_{i}$ for all $i$.

Proof of Proposition 6.9. - Let $\boldsymbol{\psi}$ be an optimal basis which is not ideal. Lemma 6.10 implies that $v\left(\psi_{i}\right)=a_{i}$ for all $i$. Let $k$ be minimal such that there exists a $\varphi=\sum_{i=1}^{k} x_{i} \psi_{i} \in L$ with $v(\varphi) \neq \min _{i}\left(x_{i} \psi_{i}\right)$. Lemma 6.4 implies that $k$ exists. It follows from the triangle inequality that $v(\varphi)>\min _{i}\left(x_{i} \psi_{i}\right)$. Write $\tilde{\varphi}=\sum_{i=1}^{k-1} x_{i} \psi_{i}$. The choice of $k$ implies that $v(\tilde{\varphi})=\min _{i<k} v\left(x_{i} \psi_{i}\right)$. Since $v(\varphi)=v\left(\tilde{\varphi}+x_{k} \psi_{k}\right)$, we conclude from Lemma 6.1 that $v(\tilde{\varphi})=v\left(x_{k} \psi_{k}\right)$. This implies that

$$
\begin{equation*}
2 \operatorname{ord}\left(x_{i}\right)+a_{i} \geq 2 \operatorname{ord}\left(x_{k}\right)+a_{k} . \tag{6.8}
\end{equation*}
$$

In particular, $\operatorname{ord}\left(x_{i}\right) \geq \operatorname{ord}\left(x_{k}\right)$, for all $i$. Therefore it is no restriction to assume that $x_{k}$ is a unit.

We define a new basis $\varphi=\left(\varphi_{i}\right)$ by $\varphi_{i}=\psi_{i}$ if $i \neq k$ and $\varphi_{k}=\varphi$. Write

$$
\tilde{Q}\left(\sum_{i} y_{i} \varphi_{i}\right)=\sum_{i \leq j} \tilde{b}_{i j} y_{i} y_{j}
$$

One computes that

$$
\tilde{b}_{j k}= \begin{cases}2 x_{j} b_{j j}+\sum_{i \neq j} b_{i j} x_{i} & \text { for } j<k, \\ \sum_{i} b_{i j} x_{i} & \text { for } j>k .\end{cases}
$$

Equation (6.8) implies that $\operatorname{ord}\left(\tilde{b}_{j k}\right) \geq\left(a_{j}+a_{k}\right) / 2$. Therefore $\boldsymbol{\varphi}$ is again an optimal basis. But $v\left(\varphi_{k}\right)=v(\varphi)>\min _{i} v\left(x_{i} \psi_{i}\right)=v\left(x_{k} \psi_{k}\right)=a_{k}$. This contradicts Lemma 6.10.

Lemma 6.11. - Let $M \subset L$ be a sublattice, i.e., a sub- $\mathbb{Z}_{\ell}$-module of rank $n$. Let $b_{1}, \ldots, b_{n}$ be the Gross-Keating invariants of $\left(M,\left.Q\right|_{M}\right)$. Then $b_{i} \geq a_{i}$.

Proof. - We choose ideal bases $\left(\psi_{1}, \ldots, \psi_{n}\right)$ for $L$ and $\left(\varphi_{1}, \ldots, \varphi_{n}\right)$ for $M$. Then $a_{i}=v\left(\psi_{i}\right)$ and $b_{i}=v\left(\varphi_{i}\right)$. Let us fix an index $i \in\{1, \ldots, n\}$ and show $b_{i} \geq a_{i}$. For an element $\psi=\sum_{j} x_{j} \psi_{j}$ of $L$, we set $\psi^{\prime}:=\sum_{j<i} x_{j} \psi_{j}$ and $\psi^{\prime \prime}:=\sum_{j \geq i} x_{j} \psi_{j}$. Then $\psi=\psi^{\prime}+\psi^{\prime \prime}$ and $v\left(\psi^{\prime \prime}\right) \geq a_{i}$. Since the vectors $\varphi_{1}^{\prime}, \ldots, \varphi_{i}^{\prime}$ lie in a subspace of rank $i-1$, there exist $x_{1}, \ldots, x_{i} \in \mathbb{Z}_{\ell}$, not all zero, such that $\sum_{j \leq i} x_{j} \varphi_{j}^{\prime}=0$. Then

$$
\sum_{j \leq i} x_{j} \varphi_{j}=\sum_{j \leq i} x_{j} \varphi_{j}^{\prime \prime}
$$

Applying Lemma 6.4 (6.4) to the left hand side and the triangle inequality (6.2) to the right hand side, we conclude that

$$
\min _{j \leq i}\left(b_{j}+2 \operatorname{ord}_{\ell}\left(x_{j}\right)\right) \geq \min _{j \leq i}\left(v\left(\varphi_{j}^{\prime \prime}\right)+\operatorname{ord}_{\ell}\left(x_{j}\right)\right) \geq \min _{j \leq i}\left(a_{i}+2 \operatorname{ord}_{\ell}\left(x_{j}\right)\right) .
$$

For the index $j$ for which $\operatorname{ord}_{\ell}\left(x_{j}\right)$ takes its minimal value we get $a_{i} \leq b_{j} \leq b_{i}$. This proves the lemma.

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[^7]
# 13. DEFORMATIONS OF ISOGENIES OF FORMAL GROUPS 

by<br>Michael Rapoport


#### Abstract

Let $\left(f_{1}, f_{2}, f_{3}\right): E \rightarrow E^{\prime}$ be a triple of isogenies between supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. We determine when the locus of deformation of ( $f_{1}, f_{2}, f_{3}$ ) inside the universal deformation space of $\left(E, E^{\prime}\right)$ is an Artin scheme, and in this case we give a formula for its length. These results are due to Gross and Keating.


Résumé (Déformations d'isogénies de groupes formels). - Soit $\left(f_{1}, f_{2}, f_{3}\right): E \rightarrow E^{\prime}$ un triplet d'isogénies entre des courbes elliptiques supersingulières sur $\overline{\mathbb{F}}_{p}$. Nous donnons un critère pour le lieu de déformation de $\left(f_{1}, f_{2}, f_{3}\right)$ dans l'espace de déformations universel de ( $E, E^{\prime}$ ) d'être un schéma artinien, et nous donnons dans ce cas une formule pour sa longueur. Ces résultats sont dûs à Gross et Keating.

Let $A$ and $A^{\prime}$ be abelian varieties of the same dimension $n$ over $\overline{\mathbb{F}}_{p}$. The universal deformation space $\mathcal{M}$ of the pair $A, A^{\prime}$ is the formal spectrum of a power series ring in $2 n^{2}$ variables over $W\left(\overline{\mathbb{F}}_{p}\right)$. Given an isogeny $f: A \rightarrow A^{\prime}$ one may pose the problem of determining the maximal locus inside $\mathcal{M}$, where $f$ can be deformed. More generally, given an $r$-tuple $f_{1}, \ldots, f_{r}$ of isogenies from $A$ to $A^{\prime}$, one may ask for the maximal locus inside $\mathcal{M}$ where $f_{1}, \ldots, f_{r}$ deform. And, one may ask when this maximal locus is the spectrum of a local Artin ring, and if so, to give a formula for its length.

These questions are very difficult and it even seems likely that no systematic answers exist in general. In this chapter we consider the case $n=1$, i.e., when $A$ and $A^{\prime}$ are elliptic curves. More precisely, we present the solution due to Gross and Keating [GK] to this problem when $A$ and $A^{\prime}$ are supersingular elliptic curves. Their proof is a clever application of results on quasi-canonical liftings and their endomorphisms. Unfortunately, some parts of their proof are not so easy to implement in the case $p=2$, which requires special attention. In fact, I only managed to deal with the case $p=2$ by making use of the classification of quadratic forms over $\mathbb{Z}_{2}$, comp. [B], and

2000 Mathematics Subject Classification. - 11F32, 11G15, 14L05.
Key words and phrases. - Formal group, quasi-canonical lifting, Kummer congruence, Gross-Keating invariants.
using a case-by-case analysis. Fortunately, S. Wewers afterwards found a uniform argument for this part of the proof which makes use of deeper properties of anisotropic quadratic forms over $\mathbb{Z}_{2}$. This proof is presented in the next chapter. We decided to present both proofs because the more pedestrian approach here gives insight into the subtleties of the Gross-Keating invariants in the case $p=2$.

Let us comment on the general problem above in another example, the case of ordinary elliptic curves, comp. [ $\mathbf{M e} \mathbf{2}]$. The case when $A$ and $A^{\prime}$ are ordinary elliptic curves has been known for a long time and is part of the Serre-Tate theory of canonical coordinates, comp. [Mes, Appendix]. Let $A$ and $A^{\prime}$ be ordinary elliptic curves and fix isomorphisms

$$
A\left[p^{\infty}\right]^{\mathrm{et}} \cong \mathbb{Q}_{p} / \mathbb{Z}_{p}, A^{\prime}\left[p^{\infty}\right]^{\mathrm{et}} \cong \mathbb{Q}_{p} / \mathbb{Z}_{p}
$$

which then induce, via the canonical principal polarization, isomorphisms

$$
A\left[p^{\infty}\right]^{0}=\widehat{\mathbb{G}}_{m}, A^{\prime}\left[p^{\infty}\right]^{0}=\widehat{\mathbb{G}}_{m}
$$

The isogeny $f: A \rightarrow A^{\prime}$ determines

$$
\left(z_{0}, z_{1}\right) \in \mathbb{Z}_{p}^{2}
$$

where $f$ is given by multiplication by $z_{1}$ on the étale part and by multiplication by $z_{0}$ on the connected part of $A\left[p^{\infty}\right]$. On the other hand, we have

$$
\left.\mathcal{M}=\operatorname{Spf} W\left(\overline{\mathbb{F}}_{p}\right)\left[t, t^{\prime}\right]\right]
$$

(Serre-Tate canonical coordinates). Then setting $q=1+t, q^{\prime}=1+t^{\prime}$, the locus inside $\mathcal{M}$ where $f$ deforms is defined by the equation

$$
q^{z_{1}}=q^{\prime z_{0}}
$$

cf. [Mes, Appendix, 3.3], comp. also [Me2, Example 2.3]. On the other hand, it is easy to see that, for any $r$-tuple of isogenies $f_{1}, \ldots, f_{r}: A \rightarrow A^{\prime}$, the locus where $f_{1}, \ldots, f_{r}$ deform is never of finite length, comp. [Go2, proof of Prop. 3.2]. These remarks show that already the case $n=1$ in the above-mentioned general problem defies a uniform solution.

I wish to thank I. Bouw, U. Görtz, Ch. Kaiser, S. Kudla, S. Wewers and Th. Zink for their help in the preparation of this manuscript, and the referee for his remarks.

## 1. Statement of the result

Let $E$ and $E^{\prime}$ be supersingular elliptic curves over $\overline{\mathbb{F}}_{p}$. Denoting by $W$ the ring of Witt vectors of $\mathbb{F}_{p}$, the ring

$$
R=W\left[\left[t, t^{\prime}\right]\right]
$$

is the universal deformation ring of the pair $E, E^{\prime}$. Let $\mathbb{E}, \mathbb{E}^{\prime}$ be the universal deformation of $E, E^{\prime}$ over $R$. Let $f_{1}, f_{2}, f_{3}: E \rightarrow E^{\prime}$ be a triple of isogenies. The
locus inside Spf $R$ to which $f_{1}, f_{2}, f_{3}$ deform is a closed formal subscheme. Let $I=$ minimal ideal in $R$ such that $f_{1}, f_{2}, f_{3}: E \longrightarrow E^{\prime}$ lift to isogenies $\mathbb{E} \longrightarrow \mathbb{E}^{\prime}(\bmod I)$. The problem in this chapter is: Determine

$$
\alpha\left(f_{1}, f_{2}, f_{3}\right)=\lg _{W} R / I
$$

(in particular, determine when this length is finite).
This problem reduces to a problem on formal groups, as follows. Let $\Gamma=\hat{\mathbb{E}}$ resp. $\Gamma^{\prime}=\hat{\mathbb{E}}^{\prime}$ be the formal group over $R$ corresponding to $\mathbb{E}$ resp. $\mathbb{E}^{\prime}$. By the Serre-Tate theorem we have
$I=$ minimal ideal in $R$ such that $\hat{f}_{1}, \hat{f}_{2}, \hat{f}_{3}: \hat{E} \longrightarrow \hat{E}^{\prime}$ lift to isogenies $\Gamma \longrightarrow \Gamma^{\prime}(\bmod I)$. Now $\hat{E}$ and $\hat{E}^{\prime}$ can both be identified with the formal group $G$ of dimension 1 and height 2 over $\overline{\mathbb{F}}_{p}$ (which is unique up to isomorphism). In this way $\hat{f}_{1}, \hat{f}_{2}, \hat{f}_{3}$ become non-zero elements of $\operatorname{End}(G)=\mathcal{O}_{D}$. Here $D$ denotes the quaternion division algebra over $\mathbb{Q}_{p}$.

On $\operatorname{Hom}\left(E, E^{\prime}\right)$ we have the quadratic form induced by the canonical principal polarization,

$$
Q(f)={ }^{t} f \circ f=\operatorname{deg} f
$$

This $\mathbb{Z}$-valued quadratic form is induced by the $\mathbb{Z}_{p}$-valued quadratic form

$$
Q(x)=x \cdot{ }^{\iota} x
$$

under the inclusion $\operatorname{Hom}\left(E, E^{\prime}\right) \subset \operatorname{End}(G)$. Here $x \mapsto^{\iota} x$ denotes the main involution on $D$ characterized by (reduced trace)

$$
\operatorname{tr}(x)=x+{ }^{\dagger} x
$$

We also write $Q(x)=\mathrm{Nm}(x)$ (reduced norm).
Let $L=\mathbb{Z}_{p} \hat{f}_{1}+\mathbb{Z}_{p} \hat{f}_{2}+\mathbb{Z}_{p} \hat{f}_{3}$ be the $\mathbb{Z}_{p}$-submodule of $\mathcal{O}_{D}$, with the quadratic form $Q$ obtained by restriction. Then

$$
I=\text { minimal ideal in } R \text { such that } L \subset \operatorname{Hom}_{R / I}\left(\Gamma, \Gamma^{\prime}\right)
$$

Assume that $(L, Q)$ is non-degenerate, i.e., $L$ is of rank 3. Then to $(L, Q)$ are associated integers $0 \leq a_{1} \leq a_{2} \leq a_{3}$, the Gross-Keating invariants. Recall ([B, section 2]) that if $p \neq 2$ these invariants are characterized by the fact that in a suitable basis $e_{1}, e_{2}, e_{3}$ of $L$ the matrix $T=\frac{1}{2}\left(\left(e_{i}, e_{j}\right)\right)_{i, j}$ is equal to

$$
\begin{equation*}
T=\operatorname{diag}\left(u_{1} p^{a_{1}}, u_{2} p^{a_{2}}, u_{3} p^{a_{3}}\right) \text { with } u_{1}, u_{2}, u_{3} \in \mathbb{Z}_{p}^{\times} \tag{1.1}
\end{equation*}
$$

Here $(x, y)=Q(x+y)-Q(x)-Q(y)$ is the bilinear form associated to the quadratic form $Q$.

Theorem 1.1. - The length of $R / I$ is finite if and only if $(L, Q)$ is non-degenerate. In this case, $\lg _{W} R / I$ only depends on the Gross-Keating invariants $\left(a_{1}, a_{2}, a_{3}\right)$ and equals $\alpha(Q)$ where

$$
\begin{aligned}
\alpha(Q)= & \sum_{i=0}^{a_{1}-1}(i+1)\left(a_{1}+a_{2}+a_{3}-3 i\right) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-2\right) / 2}\left(a_{1}+1\right)\left(2 a_{1}+a_{2}+a_{3}-4 i\right) p^{i} \\
& +\frac{a_{1}+1}{2}\left(a_{3}-a_{2}+1\right) p^{\left(a_{1}+a_{2}\right) / 2}, \text { if } a_{1} \equiv a_{2}(\bmod 2) \\
\alpha(Q)= & \sum_{i=0}^{a_{1}-1}(i+1)\left(a_{1}+a_{2}+a_{3}-3 i\right) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-1\right) / 2}\left(a_{1}+1\right)\left(2 a_{1}+a_{2}+a_{3}-4 i\right) p^{i}, \\
& \text { if } a_{1} \not \equiv a_{2}(\bmod 2)
\end{aligned}
$$

Remark 1.2. - Recall from [B, Lemma 5.3] that, since $(L, Q)$ is anisotropic, not all $a_{1}, a_{2}, a_{3}$ have the same parity. Hence the RHS of the formulas above is an integer in all cases.

Remark 1.3. - The formulas above imply that the length of $R / I$ only depends on the isomorphism class of the quadratic module $L$. This can be seen in an a priori way as follows.

First of all, there is an action of $\left(D^{\times}\right)^{2}$ on the universal deformation ring $R$, given by changing the identification of the special fibers of $\Gamma, \Gamma^{\prime}$ with $G, G$ by a pair of automorphisms of $G$. More precisely, an element $d \in D^{\times}$defines a quasi-isogeny of $G$, as the composition $\mathrm{Frob}^{-v} \circ d$. Here Frob denotes the Frobenius endomorphism and $v=v(d)$ is the valuation of $d$. Since this is a quasi-isogeny of height 0 , it is an automorphism of $G$. Note however, that this is only a semi-linear automorphism, and therefore also the induced automorphism by $\left(d_{1}, d_{2}\right) \in\left(D^{\times}\right)^{2}$ on $R$ is only semi-linear.

It follows that for $\left(d_{1}, d_{2}\right) \in\left(D^{\times}\right)^{2}$ with $v\left(d_{1}\right)=v\left(d_{2}\right)$, the length of the deformation ring $R / I$ for $L=\mathbb{Z}_{p} \hat{f}_{1}+\mathbb{Z}_{p} \hat{f}_{2}+\mathbb{Z}_{p} \hat{f}_{3}$ is equal to the length of the deformation ring $R / I^{\prime}$ for $L^{\prime}=\mathbb{Z}_{p} \hat{f}_{1}^{\prime}+\mathbb{Z}_{p} \hat{f}_{2}^{\prime}+\mathbb{Z}_{p} \hat{f}_{3}^{\prime}$, where $\hat{f}_{i}^{\prime}=d_{1} \hat{f}_{i} d_{2}^{-1}$. Hence it suffices to show that for any two isometric ternary lattices $L$ and $L^{\prime}$ in $\mathcal{O}_{D}$, there exists $\left(d_{1}, d_{2}\right) \in\left(D^{\times}\right)^{2}$ with $v\left(d_{1}\right)=v\left(d_{2}\right)$ and $L^{\prime}=d_{1} L d_{2}^{-1}$.

Fix a nondegenerate ternary form $Q$ over $\mathbb{Z}_{p}$. We want to show that for any two isometries $\sigma, \sigma^{\prime}$ from $Q$ to $\mathcal{O}_{D}$, there exists $\left(d_{1}, d_{2}\right) \in\left(D^{\times}\right)^{2}$ as above with $L^{\prime}=$ $d_{1} L d_{2}^{-1}$, where $L$ resp. $L^{\prime}$ denotes the image of $\sigma$, resp. $\sigma^{\prime}$. By [ $\mathbf{W d 1}$, Lemma 1.6], we may identify $\mathrm{SO}(D, \mathrm{Nm})$ with the group

$$
\left\{\left(d_{1}, d_{2}\right) \in\left(D^{\times}\right)^{2} \mid \operatorname{Nm}\left(d_{1}\right)=\operatorname{Nm}\left(d_{2}\right)\right\} / \mathbb{Q}_{p}^{\times}
$$

By [Wd2, 1.3], the group $\mathrm{SO}(D, \mathrm{Nm})$ acts simply transitively on the set of isometries $\sigma$, hence there exists a unique $\left(d_{1}, d_{2}\right) \in \mathrm{SO}(D, \mathrm{Nm})$ with $\sigma^{\prime}=d_{1} \sigma d_{2}^{-1}$. The pair $\left(d_{1}, d_{2}\right)$ has the required properties.

To start the proof of Theorem 1.1, we first recall the following proposition.

Proposition 1.4. - Let $\psi \in \operatorname{End}(G)$ be an isogeny, i.e., $\psi \neq 0$. Let $J$ be the minimal ideal in $R=W\left[\left[t, t^{\prime}\right]\right]$ such that $\psi$ lifts to an isogeny $\Gamma \rightarrow \Gamma^{\prime}(\bmod J)$. Then the closed formal subscheme $\mathcal{T}$ of $\mathcal{S}=\operatorname{Spf} R$ is a relative divisor over $\operatorname{Spf} W$. In other words, $J$ is generated by an element which is neither a unit nor divisible by $p$.

Proof. - This is the special case of [Ww1, Prop. 5.1], where (in the notation used there) $K=\mathbb{Q}_{p}$. A different proof that $\mathcal{T}$ is a divisor is (at least implicitly) contained in $[\mathbf{Z}$, section 2.5].

Let us prove the first statement of Theorem 1.1. If $(L, Q)$ is degenerate, then $L$ is generated by two elements. Hence the deformation locus is by Proposition 1.4 the intersection of two divisors on a regular 3-dimensional formal scheme and therefore cannot be of finite length. Now assume that $(L, Q)$ is non-degenerate. Now $\operatorname{Hom}\left(E, E^{\prime}\right) \otimes \mathbb{Z}_{p}=\operatorname{End}(G)$, so we find isogenies $f_{1}, f_{2}, f_{3}: E \rightarrow E^{\prime}$ with $\mathbb{Z}_{p}$-span equal to $L$. Let $T=\operatorname{Spec} W\left[\left[t, t^{\prime}\right] / J\right.$. Then $f_{1}, f_{2}, f_{3}$ deform to isogenies from $\mathbb{E}_{T}$ to $\mathbb{E}_{T}^{\prime}$. Hence at any point $t$ of $T$ we have $\operatorname{rg} \operatorname{Hom}\left(\mathbb{E}_{t}, \mathbb{E}_{t}^{\prime}\right)>2$, hence the elliptic curves $\mathbb{E}_{t}$ and $\mathbb{E}_{t}^{\prime}$ are supersingular. Since supersingular points are isolated in the moduli scheme, it follows that $T$ is an Artin scheme, as was to be shown.

From now on we assume that $(L, Q)$ is non-degenerate. Let $\psi_{1}, \psi_{2}, \psi_{3}$ be an optimal basis of $L$. If $p \neq 2$, this means that the matrix of the bilinear form $Q$ in terms of this basis is diagonal as in (1.1).

Corollary 1.5. - Let $\mathcal{T}_{i} \subset \mathcal{S}$ be the locus, defined by the ideal $I_{i}$ in $R$, where $\psi_{i}$ lifts to an isogeny $\Gamma \rightarrow \Gamma^{\prime}\left(\bmod I_{i}\right)$. Then

$$
\lg _{W} R / I=\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}
$$

Here on the RHS there appears the intersection product of divisors on a regular scheme, defined by the Samuel multiplicity or via the Koszul complex of the equations $g_{i}$ of $I_{i}$,

$$
\chi\left(\left(g_{1}, g_{2}, g_{3}\right)\right)=\sum(-1)^{i} \lg \left(H_{i}\left(K_{\bullet}\left(g_{1}, g_{2}, g_{3}\right)\right)\right)
$$

(comp. [F, Ex. 7.1.2]).
Proof. - By our non-degeneracy assumption, the $g_{i}$ form a regular sequence in a regular local ring.

The corollary allows us to apply the intersection calculus of divisors on a regular scheme. In particular, the RHS is multilinear in all three entries.

Theorem 1.1 will be proved by induction on $a_{1}+a_{2}+a_{3}$. It will follow from the following three propositions.

Proposition 1.6. - Let $a_{3} \leq 1$. Then

$$
\alpha(Q)= \begin{cases}1 & a_{2}=0 \\ 2 & a_{2}=1\end{cases}
$$

Hence Theorem 1.1 holds true in this case.

Proposition 1.7. - Let $\psi_{3}=p \cdot \psi_{3}^{\prime}$ with $\psi_{3}^{\prime} \in \operatorname{End}(G)$. Then

$$
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}=\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}^{\prime}\right)_{\mathcal{S}}+\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}
$$

Here $\mathcal{T}_{i}(i=1,2,3)$ resp. $\mathcal{T}_{3}^{\prime}$ denotes the deformation locus for $\psi_{i}$ resp. $\psi_{3}^{\prime}$ and $\mathcal{S}_{(p)}=$ $\mathcal{S} \times$ Spf ${ }_{W} \operatorname{Spf} \overline{\mathbb{F}}_{p}$ is the special fiber of $\mathcal{S}$.

Proposition 1.8. - If $a_{1} \equiv a_{2}(\bmod 2)$ then

$$
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}=\sum_{i=0}^{a_{1}-1} 2(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-2\right) / 2} 2\left(a_{1}+1\right) p^{i}+\left(a_{1}+1\right) p^{\left(a_{1}+a_{2}\right) / 2}
$$

If $a_{1} \not \equiv a_{2}(\bmod 2)$ then

$$
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}=\sum_{i=0}^{a_{1}-1} 2(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-1\right) / 2} 2\left(a_{1}+1\right) p^{i}
$$

These propositions indeed imply Theorem 1.1. For this recall ([B, Cor. 5.8]) that we can (and do) choose $\psi_{3}$ such that $v\left(\psi_{3}\right)=a_{3}$. Here, as elsewhere, we denote by $v$ the valuation function on $D$. Now, if $a_{3}>1$, then there exists $\psi_{3}^{\prime} \in \operatorname{End}(G)$ with $\psi_{3}=p \psi_{3}^{\prime}$.

Lemma 1.9. - Let $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ be an optimal basis of the lattice L. Let $\psi_{3}=p \psi_{3}^{\prime}$ with $\psi_{3}^{\prime} \in L$ and denote by $L^{\prime}$ the lattice generated by $\psi_{1}, \psi_{2}, \psi_{3}^{\prime}$. Then the invariants of $L^{\prime}$ are given in terms of the invariants $\left(a_{1}, a_{2}, a_{3}\right)$ of $L$ by

$$
\left(a_{1}, a_{2}, a_{3}-2\right)
$$

(in some order so that they form a weakly increasing sequence).
This is obvious for $p \neq 2$ from the characterization in (1.1). For $p=2$, the proof is given in the appendix, using the classification of quadratic forms over $\mathbb{Z}_{2}$. An alternative, more conceptual proof can be found in $[\mathbf{B}$, Cor. 6.7].

Using this lemma, the above propositions give an inductive procedure for calculating $\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}$. The formula in Theorem 1.1 follows from this calculation.

We now devote one section each to the proof of these three propositions. For Propositions 1.6 and 1.7 the case $p=2$ presents additional problems. In order not to obscure the argument, the problems arising for $p=2$ are relegated to the appendix to this chapter. In the chapter following this one, a variant of the proofs of Propositions 1.6 and 1.7 is given which avoids any case-by-case considerations.

## 2. The induction start: Proposition 1.5

Since not all $a_{i}$ have the same parity, we have $a_{1}=0$. Hence $\psi_{1}$ is an automorphism of $G$. Since $\Gamma^{\prime}$ is a universal deformation of $G$, the ideal $I_{1}$ in $W\left[t, t^{\prime}\right]$ defining the deformation locus of $\psi_{1}$ is of the form $I_{1}=\left(t^{\prime}-h(t)\right)$, for some $h \in W\left[[t]\right.$. For $I \supset I_{1}$,
it follows that $\psi_{i}$ lifts to an isogeny $\Gamma \rightarrow \Gamma^{\prime}(\bmod I)$ if and only if ${ }^{\iota} \psi_{1} \circ \psi_{i}$ lifts to an endomorphism of $\Gamma(\bmod I \cap W[t]])$. Let

$$
\varphi_{2}={ }^{\iota} \psi_{1} \circ \psi_{2}, \varphi_{3}={ }^{\iota} \psi_{1} \circ \psi_{3} \text { in } \operatorname{End}(G)
$$

We see that
$\mathcal{T}_{1} \cap \mathcal{T}_{2} \cap \mathcal{T}_{3}=$ locus in Spf $W[[t]]$ where $\varphi_{2}$ and $\varphi_{3}$ lift to endomorphisms of $\Gamma$. More precisely, for $i=2$ or $i=3$, let $J_{i}$ be the minimal ideal in $W[[t]]$ such that $\varphi_{i}$ lifts to an endomorphism of $\Gamma\left(\bmod J_{i}\right)$. Then $\mathcal{T}_{1} \cap \mathcal{T}_{2} \cap \mathcal{T}_{3}$ is isomorphic to the closed formal subscheme of Spf $W[t]]$ defined by $J_{2}+J_{3}$.

Now let $p \neq 2$. Then we have from the definition of an optimal basis

$$
\begin{align*}
{ }^{\iota} \varphi_{i} & =-\varphi_{i} \text { and } \operatorname{Nm}\left(\varphi_{i}\right)=u_{1} u_{i} p^{a_{i}} \quad, \quad i=2,3 .  \tag{2.1}\\
\varphi_{2} \varphi_{3} & =-\varphi_{3} \varphi_{2} .
\end{align*}
$$

Let $K=\mathbb{Q}_{p}\left(\sqrt{-u_{1} u_{2} p^{a_{2}}}\right)$. Since $a_{2} \leq 1$, we deduce from (2.1) that $\varphi_{2}$ generates the ring of integers $\mathcal{O}_{K}$. Hence $\Gamma\left(\bmod J_{2}\right)$ is the canonical lifting of $G$ relative to the quadratic extension $K$ of $\mathbb{Q}_{p}$, comp. [ $\mathbf{W w} \mathbf{w}$, Def. 3.1]. Applying the following lemma, we obtain

$$
\varphi_{3} \in \Pi^{a_{3}} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{a_{3}+1} \mathcal{O}_{D}\right)
$$

with $a_{3}=1$. Now applying [Ww1, Thm. 1.4], or [Vl, Thm. 2.1], we have

$$
\lg W[t t] /\left(J_{2}+J_{3}\right)= \begin{cases}\frac{a_{3}+1}{2}=1 & \text { if } a_{2}=0 \\ a_{3}+1=2 & \text { if } a_{2}=1\end{cases}
$$

Remark 2.1. - The proof shows more generally Theorem 1.1 in the case where $p \neq 2$ and $a_{1}=0$ : one appeals to [ $\left.\mathbf{V l}, \mathrm{Thm} .2 .1\right]$.

Lemma 2.2. - We allow $p=2$. Let $K$ be a quadratic extension of $\mathbb{Q}_{p}$ contained in $D$, which is unramified or tamely ramified. Let $x \in \mathcal{O}_{D}$ which anticommutes with $K$, i.e., such that conjugation by $x$ induces on $K$ the non-trivial automorphism of $K$. Let $r=v(x)$. Then

$$
x \in \Pi^{r} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{r+1} \mathcal{O}_{D}\right)
$$

Here $\Pi$ denotes a uniformizer of $\mathcal{O}_{D}$.
Proof. - We distinguish cases.
Case $K / \mathbb{Q}_{p}$ unramified. - In this case we can choose a uniformizer $\Pi$ of $\mathcal{O}_{D}$ with $\Pi^{2}=p$ and anticommuting with $K$. Then

$$
\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \Pi
$$

where the first summand commutes with $K$, and the second summand anticommutes with $K$. Then

$$
\mathcal{O}_{K}+\Pi^{s} \mathcal{O}_{D}=\mathcal{O}_{K} \oplus p^{\left[\frac{s}{2}\right]} \mathcal{O}_{K} \cdot \Pi
$$

Now if $x$ anticommutes with $K$, then $r=v(x)=2 t+1$ is odd and $x \notin \mathcal{O}_{K}+\Pi^{r+1} \mathcal{O}_{D}=$ $\mathcal{O}_{K} \oplus p^{t+1} \mathcal{O}_{K} \cdot \Pi$.
Case $K / \mathbb{Q}_{p}$ tamely ramified. - In this case we can write $\mathcal{O}_{K}=\mathbb{Z}_{p}[\pi]$ with $\pi^{2}=u \cdot p$, for $u \in \mathbb{Z}_{p}^{\times}$. Then

$$
\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot j \quad, \quad j^{2}=u^{\prime} \in \mathbb{Z}_{p}^{\times} \backslash \mathbb{Z}_{p}^{\times, 2}
$$

where the first summand commutes with $K$ and the second summand anticommutes with $K$. Then

$$
\mathcal{O}_{K}+\Pi^{s} \mathcal{O}_{D}=\mathcal{O}_{K} \oplus \pi^{s} \mathcal{O}_{K} \cdot j
$$

If $x$ anticommutes with $K$, it lies in $\pi^{r} \mathcal{O}_{K} \cdot j$ but not in $\pi^{r+1} \mathcal{O}_{K} \cdot j$, hence $x \notin$ $\mathcal{O}_{K}+\Pi^{r+1} \mathcal{O}_{D}=\mathcal{O}_{K} \oplus \pi^{r+1} \mathcal{O}_{K} \cdot j$.

Remark 2.3. - In the case of wild ramification $(p=2)$ it can happen that $x$ can be corrected by an element of $\mathcal{O}_{K}$ to have higher valuation than $r=v(x)$.

## 3. The induction step: Proposition 1.6.

It suffices to prove

$$
\left(\mathcal{C} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}=\left(\mathcal{C} \cdot \mathcal{T}_{3}^{\prime}\right)_{\mathcal{S}}+\left(\mathcal{C} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}
$$

for every irreducible component $\mathcal{C}$ of $\mathcal{T}_{1} \cap \mathcal{T}_{2}$. Let
$J=$ minimal ideal in $W \llbracket t t]$ such that ${ }^{\circ} \psi_{1} \circ \psi_{2}$ lifts to an isogeny $\Gamma \longrightarrow \Gamma(\bmod J)$
$J^{\prime}=$ minimal ideal in $W\left[\left[t^{\prime}\right]\right]$ such that $\psi_{2} \circ^{\prime} \psi_{1}$ lifts to an isogeny $\Gamma^{\prime} \longrightarrow \Gamma^{\prime}\left(\bmod J^{\prime}\right)$.
We have an obvious inclusion

$$
\left.\mathcal{T}_{1} \cap \mathcal{T}_{2} \longleftrightarrow \mathcal{X} \underset{\text { Df. }}{=} \operatorname{Spf}(W[t t]] / J\right) \hat{\otimes}_{W}\left(W\left[\left[t^{\prime}\right]\right] / J^{\prime}\right)
$$

The proof of [Ww1, Prop. 5.1] shows that $J$ is generated by one element. Now ${ }^{\iota} \psi_{1} \circ \psi_{2}$ is not scalar. Hence the generator of $J$ is not divisible by $p$, because otherwise ${ }^{\iota} \psi_{1} \circ \psi_{2}$ would extend to the universal deformation of $G$ over $\overline{\mathbb{F}}_{p}[t t]$, contradicting [Vi, Thm. 1.1]. The same argument applies to $J^{\prime}$ instead of $J$. Hence all irreducible components of $\mathcal{X}$ have dimension 1 , and each irreducible component of $\mathcal{T}_{1} \cap \mathcal{T}_{2}$ is also an irreducible component of $\mathcal{X}$. We now determine the irreducible components of $\mathcal{X}$.

The endomorphisms $\varphi={ }^{\iota} \psi_{1} \circ \psi_{2}$ and $\varphi^{\prime}=\psi_{2} \circ^{\iota} \psi_{1}$ generate quadratic extensions $K=\mathbb{Q}_{p}(\varphi)$ resp. $K^{\prime}=\mathbb{Q}_{p}\left(\varphi^{\prime}\right)$ which are conjugate inside $D$.

Lemma 3.1. - The order $\mathbb{Z}_{p}[\varphi]$ in $K$ has conductor $\left[\left(a_{1}+a_{2}\right) / 2\right]$.
Proof for $p \neq 2$. - In this case the fact that the $\psi_{i}$ form an optimal basis, i.e., diagonalize the bilinear form as in (1.1), implies that

$$
\operatorname{tr}(\varphi)=0 \quad, \quad \varphi^{2}=-u_{1} u_{2} p^{a_{1}+a_{2}}
$$

Hence $\mathbb{Z}_{p}[\varphi]=\mathbb{Z}_{p}+p^{r} \mathcal{O}_{K}$, with $r=\left[\left(a_{1}+a_{2}\right) / 2\right]$.

We therefore obtain an equality of divisors on $\operatorname{Spf} W[t]]$,

$$
\operatorname{Spf} W[t]] / J=\sum_{s=0}^{\left[\left(a_{1}+a_{2}\right) / 2\right]} \mathcal{W}_{s}(\varphi)
$$

Here $\mathcal{W}_{s}(\varphi)$ is the quasicanonical locus of level $s$, with respect to the embedding of $K$ in $D$ defined by $\varphi$. Hence $\mathcal{W}_{s}(\varphi)$ is a reduced irreducible regular divisor such that the pullback of $\Gamma$ to $\mathcal{W}_{s}(\varphi)$ has as its endomorphism algebra the order $\mathcal{O}_{s}=\mathbb{Z}_{p}+p^{s} \mathcal{O}_{K}$ of conductor ${ }^{(1)} s$ in $K$. We may choose an identification

$$
\mathcal{W}_{s}(\varphi)=\operatorname{Spf} W_{s}
$$

where $W_{s}$ is the ring of integers in the ray class field extension $M_{s}$ of the completion $M$ of the maximal unramified extension of $K$ with norm group $\mathcal{O}_{s}^{\times}$.

Analogously we have

$$
\text { Spf } W\left[\left[t^{\prime}\right] / J^{\prime}=\sum_{s=0}^{\left[\left(a_{1}+a_{2}\right) / 2\right]} \mathcal{W}_{s}\left(\varphi^{\prime}\right)\right.
$$

We apply the following simple observation.
Lemma 3.2. - Let $M$ be a discretely valued field. Let $M \subset K \subset L$ be finite field extensions such that $K \otimes_{M} L=L^{|K: M|}$ (e.g. $K / M$ Galois). For each field embedding $\tau: K \rightarrow L$ with $\tau \mid M=\mathrm{id}$, let $\Gamma_{\tau}$ be the graph of the corresponding morphism $\operatorname{Spec} \mathcal{O}_{L} \rightarrow \operatorname{Spec} \mathcal{O}_{K}$. Then

$$
\operatorname{Spec} \mathcal{O}_{K} \otimes_{\mathcal{O}_{M}} \mathcal{O}_{L}=\bigcup_{\tau} \Gamma_{\tau}
$$

Proof. - Obviously, the RHS is a closed subscheme of the LHS with identical generic fibers. But the LHS is flat over $\mathcal{O}_{M}$, hence is the closure of its generic fiber.

Note that $W_{r} \subset W_{s}$ whenever $r \leq s$. The lemma implies that each irreducible component of $\mathcal{W}_{r}(\varphi) \cap \mathcal{W}_{s}\left(\varphi^{\prime}\right)$ is isomorphic to $\operatorname{Spf} W_{m}$, where $m=\max \{r, s\}$. Hence each irreducible component of $\mathcal{T}_{1} \cap \mathcal{T}_{2}$ is isomorphic to Spf $W_{s}$ for some $s$ with $0 \leq s \leq\left[\left(a_{1}+a_{2}\right) / 2\right]$.

Proposition 3.3. - Let $F_{r}, F_{s}$ be quasi-canonical liftings of $G$ of level $r, s$ (with respect to the quadratic extension $K$ of $\mathbb{Q}_{p}$ ) defined over the ring of integers $\mathcal{O}$ of a finite extension of Frac $W$. Assume that $\psi_{1}, \psi_{2}$ lift to isogenies $F_{r} \rightarrow F_{s}$ over $\mathcal{O}$. Let $I$ resp. $I^{\prime}$ be the minimal ideal in $\mathcal{O}$ such that $\psi_{3}=p \psi_{3}^{\prime}$, resp. $\psi_{3}^{\prime}$ lifts to an isogeny $F_{r} \rightarrow F_{s}(\bmod I)$ resp. $F_{r} \rightarrow F_{s}\left(\bmod I^{\prime}\right)$. Then $I=p I^{\prime}$.

[^8]Proof. - Perhaps replacing the isogenies by their duals, we may assume $r \leq s$. First assume $r=s$. All quasi-canonical liftings of level $r$ are conjugate under $\operatorname{Gal}\left(M_{r} / M\right)$. By [ $\mathbf{W w} \mathbf{w}$, Remark 3.3], there exists an isomorphism of the underlying formal groups

$$
\gamma: F_{s} \longrightarrow F_{r}
$$

such that

$$
\varphi \circ \gamma=\gamma \circ \varphi^{\prime}
$$

However, $\gamma$ is in general not an isomorphism of deformations of $G$, since $\gamma$ conjugates the subfield $K=\mathbb{Q}_{p}(\varphi)$ of $D$ into the subfield $K^{\prime}=\mathbb{Q}_{p}\left(\varphi^{\prime}\right)$, hence $\gamma$ may be a non-central element of $D$. Let

$$
\begin{equation*}
u=\operatorname{Nm}(\gamma) \in \mathbb{Z}_{p}^{\times} \tag{3.1}
\end{equation*}
$$

We set

$$
\varphi_{i}=\gamma \circ \psi_{i} \in \operatorname{End}\left(F_{r}\right), \quad i=1,2,3
$$

Then

$$
\varphi \circ \varphi_{i}=\varphi_{i} \circ \varphi \quad, \quad i=1,2
$$

Lemma 3.4. - We have $2 r \leq a_{2}$ and $2 r<a_{3}$.
Proof for $p \neq 2$. - Since $F_{r}$ is a quasi-canonical lifting of level $r$, it suffices for the first statement to show that the conductor of one of the orders $\mathbb{Z}_{p}\left[\varphi_{1}\right]$ resp. $\mathbb{Z}_{p}\left[\varphi_{2}\right]$ is at most $a_{2} / 2$. Now $v\left(\varphi_{i}\right)=a_{i}$. But $\varphi_{i}$ is not traceless. Set

$$
\varphi_{i}^{0}=\varphi_{i}-\frac{1}{2} \operatorname{tr}\left(\varphi_{i}\right) \quad, \quad i=1,2
$$

Then $\varphi_{i}^{0}$ is traceless and hence the conductor of $\mathbb{Z}_{p}\left[\varphi_{i}\right]=\mathbb{Z}_{p}\left[\varphi_{i}^{0}\right]$ is equal to $\left[v\left(\varphi_{i}^{0}\right) / 2\right]$. Hence it suffices to show

$$
\begin{equation*}
v\left(\varphi_{i}^{0}\right) \leq a_{2} \text { for } i=1 \text { or } i=2 \tag{3.2}
\end{equation*}
$$

We distinguish cases.
Case $K / \mathbb{Q}_{p}$ unramified. -- Then $a_{1}$ and $a_{2}$ are even and

$$
\varphi_{i}=\lambda_{i} p^{a_{i} / 2} \quad, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times} \quad, \quad i=1,2
$$

Then $\operatorname{tr}\left(\varphi_{i}\right)=\left(\lambda_{i}+{ }^{\iota} \lambda_{i}\right) p^{a_{i} / 2}$ and

$$
\varphi_{i}^{0}=\frac{1}{2}\left(\lambda_{i}-{ }^{\iota} \lambda_{i}\right) \cdot p^{a_{i} / 2}
$$

Hence $v\left(\varphi_{i}^{0}\right)=a_{i}$ unless the residue class $\left[\lambda_{i}\right]$ of $\lambda_{i}$ lies in $\mathbb{F}_{p}$. But since the $\psi_{i}$ diagonalize the bilinear form, we have

$$
\begin{equation*}
{ }^{\iota} \varphi_{1} \circ \varphi_{2}=-{ }^{\iota} \varphi_{2} \circ \varphi_{1} \tag{3.3}
\end{equation*}
$$

Hence not both $\left[\lambda_{1}\right]$ and $\left[\lambda_{2}\right]$ can lie in $\mathbb{F}_{p}$ whence the claim (3.2). Now if $a_{3}=2 r$, then $2 r=a_{2}=a_{3}$. Hence $a_{1}$ would have to be odd, which is impossible.

Case $K / \mathbb{Q}_{p}$ ramified. - Let $\pi \in \mathcal{O}_{K}$ be a uniformizer with ${ }^{\iota} \pi=-\pi$. Let

$$
\varphi_{i}=\lambda_{i} \pi^{a_{i}} \quad, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times} \quad, \quad i=1,2 .
$$

Then

$$
\varphi_{i}^{0}=\frac{1}{2}\left(\lambda_{i}-(-1)^{a_{i}} \cdot{ }^{\iota} \lambda_{i}\right) \cdot \pi^{a_{i}} .
$$

Hence $v\left(\varphi_{i}^{0}\right)=a_{i}$ if $a_{i}$ is odd. Now the identity (3.3) implies

$$
(-1)^{a_{1}} \cdot{ }^{\iota} \lambda_{1} \lambda_{2}=-(-1)^{a_{2}} \cdot \lambda_{1}{ }^{\iota} \lambda_{2} .
$$

Hence $a_{1}$ and $a_{2}$ have to have different parities which shows (3.2) in this case. Now if $a_{3}=2 r$, then $a_{1}<2 r$ would have to be odd which contradicts $2 r \leq v\left(\varphi_{1}^{0}\right)=a_{1}$.

Lemma 3.5. - We have $\varphi_{3} \in \Pi^{a_{3}} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{a_{3}+1} \mathcal{O}_{D}\right)$.
Proof for $p \neq 2$. - Again using that the $\psi_{i}$ diagonalize the bilinear form, we have

$$
\varphi_{3} \varphi={ }^{\iota} \varphi \varphi_{3}
$$

Since $v\left(\varphi_{3}\right)=a_{3}$, an application of Lemma 2.2 gives the result.

We now apply [Vl, Thm. 2.1]. Since $a_{3} \geq 2 r-1$, we are in the "stable range" of that result. Hence $I$ is the $n$-th power of the maximal ideal of $\mathcal{O}$, where

$$
\begin{equation*}
n=2 \cdot \frac{p^{r}-1}{p-1} \cdot\left|\mathcal{O}: W_{r}\right|+\left(\frac{a_{3}+1}{2}-r\right) \cdot|\mathcal{O}: W| \tag{3.4}
\end{equation*}
$$

Now $v\left(\varphi_{3}^{\prime}\right)=a_{3}-2$. Since $a_{3}-2 \geq 2 r-1$, we are again in the stable range and the ideal $I^{\prime}$ is the $n^{\prime}$-th power of the maximal ideal of $\mathcal{O}$, where $n^{\prime}$ is given by (3.4) with $a_{3}$ replaced by $a_{3}-2$. Hence $n-n^{\prime}=|\mathcal{O}: W|$. This proves the proposition in the case $r=s$.

To prove the general case, we use the following lemma. For the proof we refer to [Ww1, Cor. 5.3]. Note that the element $\pi_{1}$ appearing in the statement below has the same valuation as a uniformizer of $W_{s+1}$, by [ $\mathbf{W w} \mathbf{1}$, Cor. 4.8].

Lemma 3.6. - Let $r \leq s$ and let $F_{r}, F_{s}$ and $F_{s+1}$ be quasi-canonical liftings of level $r, s$, and $s+1$, all defined over $\mathcal{O}$. Let $\pi: F_{s} \rightarrow F_{s+1}$ be an isogeny of degree $p$ defined over $\mathcal{O}$ and write $\pi$ in terms of a formal parameter

$$
\pi(X)=\pi_{1} X+\pi_{2} X^{2}+\ldots \quad, \quad \pi_{i} \in \mathcal{O}
$$

Let $\psi \in \operatorname{End}(G) \backslash\{0\}$ and let $I(r, s)$ be the minimal ideal in $\mathcal{O}$, such that $\psi$ lifts to an isogeny $F_{r} \rightarrow F_{s}(\bmod I(r, s))$. Let $I(r, s+1)$ be the minimal ideal in $\mathcal{O}$, such that $\pi \circ \psi$ lifts to an isogeny $F_{r} \rightarrow F_{s+1}(\bmod I(r, s+1))$. Then

$$
I(r, s+1)=\pi_{1} I(r, s)
$$

The lemma shows that if the assertion of Proposition 3.3 holds for $\psi_{1}, \psi_{2}, \psi_{3}, \psi_{3}^{\prime}$ : $F_{r} \rightarrow F_{s}$, it holds for $\pi \circ \psi_{1}, \pi \circ \psi_{2}, \pi \circ \psi_{3}, \pi \circ \psi_{3}^{\prime}: F_{r} \rightarrow F_{s+1}$ as well (note that $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ is an optimal basis of their $\mathbb{Z}_{p}$-span if and only if $\left(\pi \circ \psi_{1}, \pi \circ \psi_{2}, \pi \circ \psi_{3}\right)$ is an optimal basis of their $\mathbb{Z}_{p}$-span). We note the following lemma.

Lemma 3.7. - Let $r \leq s$ and let $F_{r}, F_{s+1}$ be quasi-canonical liftings of level $r, s+1$ defined over $\mathcal{O}$. Then all isogenies $\psi: F_{r} \rightarrow F_{s+1}$ factor through an isogeny $F_{s} \rightarrow$ $F_{s+1}$ of degree $p$, where $F_{s}$ is a quasi-canonical lifting of level s.

Proof. - This follows from the proof of Prop. 1.1 in [Ww2]. After choosing suitable isogenies from the canonical lifting to $F_{r}$ and to $F_{s+1}$, we may assume that the Tate modules of $F_{r}$ and $F_{s+1}$ are of the form

$$
T_{r}=\left(\mathbb{Z}_{p} \cdot p^{-r}+\mathcal{O}_{K}\right) \cdot t, \quad T_{s+1}=\left(\mathbb{Z}_{p} \cdot p^{-(s+1)}+\mathcal{O}_{K}\right) \cdot t
$$

Let $F_{s}$ be defined by $T_{s}=\left(\mathbb{Z}_{p} \cdot p^{-s}+\mathcal{O}_{K}\right) \cdot t$. Then (loc. cit.),

$$
\begin{aligned}
\operatorname{Hom}\left(F_{r}, F_{s+1}\right) & =\left\{\alpha \in \mathcal{O}_{K} \mid \alpha T_{r} \subset T_{s+1}\right\} \\
& =\left\{\alpha \in \mathcal{O}_{K} \mid \alpha T_{r} \subset T_{r}\right\} \\
& =\left\{\alpha \in \mathcal{O}_{K} \mid \alpha T_{r} \subset T_{s}\right\}
\end{aligned}
$$

Therefore all isogenies $F_{r} \rightarrow F_{s+1}$ factor through $F_{r} \rightarrow F_{s}$.
Using the previous two lemmas we now prove Proposition 3.3 by induction on the difference $s-r$. Indeed, the induction step from $(r, s)$ to $(r, s+1)$ is obvious, except in the case ${ }^{(2)}$ when the result $\tilde{\psi}_{3}: F_{r} \longrightarrow F_{s}$ of dividing $\psi_{3}: F_{r} \longrightarrow F_{s+1}$ by $\pi$ is not of the form $\tilde{\psi}_{3}=p \tilde{\psi}_{3}^{\prime}$, for a suitable $\tilde{\psi}_{3}^{\prime}: F_{r} \longrightarrow F_{s}$. However, in this case we have $a_{3}=v\left(\psi_{3}\right)=2$ and hence $r=s=0$ and $v\left(\psi_{3}^{\prime}\right)=0$. In this case the ideal $I^{\prime}$ describes the locus where the quasi-canonical lifting $F_{1}$ is isomorphic to the canonical lifting $F_{0}$. By [ $\mathbf{W} \mathbf{w} \mathbf{1}$, Cor. 4.7], the ideal $I^{\prime}$ is equal to the $n$-th power of the maximal ideal of $\mathcal{O}$, where $n=e / e_{1}$ with $e$ the absolute ramification index of $\mathcal{O}$, and $e_{1}$ the absolute ramification index of $W_{1}$. By [ $\mathbf{V l}$, Thm. 2.1], the ideal $I(0,0)$ is equal to the $e$-th power of the maximal ideal of $\mathcal{O}$. On the other hand, the element $\pi_{1}$ occurring in Lemma 3.5 has valuation $e / e_{1}$ in $\mathcal{O}$, cf. [ $\mathbf{W w} \mathbf{1}$, Cor. 4.8]. Hence $I(0,1)=p I^{\prime}$, as required.

## 4. Intersection with $\mathcal{S}_{(p)}$ : Proposition 1.7.

For the proof of Proposition 1.8 we will make use of the Kummer congruence ([KM, 13.4.6]). We first recall the statement.

[^9]We denote by $\mathcal{M}$ the moduli stack of elliptic curves over Spec $\mathbb{F}_{p}$. For integers $a, b$ with $a \geq 0, b \geq 0$ and $a+b=n$, we form the fiber product stack $\mathcal{M}_{a, b}$,


Here $\Delta$ denotes the diagonal morphism and the upper horizontal morphism sends $\left(E, E^{\prime}\right)$ to $\left(E^{\left(p^{a}\right)}, E^{\prime\left(p^{b}\right)}\right)$. Here we denoted by $E^{\left(p^{a}\right)}$ the pullback of $E$ under the $a^{\text {th }}$ power of the Frobenius morphism. Then $\mathcal{M}_{a, b}$ classifies pairs $\left(E, E^{\prime}\right)$ with an isomorphism $\alpha: E^{\left(p^{a}\right)} \xrightarrow{\sim} E^{\prime\left(p^{b}\right)}$.

We consider the moduli stack $\mathcal{M}_{\left(p^{n}\right)}$ over Spec $\mathbb{F}_{p}$ classifying isogenies $E \rightarrow E^{\prime}$ of degree $p^{n}$ (in [Go2], this stack is denoted by $\mathcal{T}_{p^{n}, \mathbb{F}_{p}}$ ). We obtain a morphism

$$
\varphi_{a, b}: \mathcal{M}_{a, b} \longrightarrow \mathcal{M}_{\left(p^{n}\right)}
$$

It sends ( $E, E^{\prime}, \alpha$ ) to the composition isogeny

$$
E \xrightarrow{F^{a}} E^{\left(p^{a}\right)} \xrightarrow[\alpha]{\sim} E^{\prime\left(p^{b}\right)} \xrightarrow{t} F^{b} E^{\prime} .
$$

Letting $a, b$ vary we obtain a morphism

$$
\varphi: \coprod_{\substack{a+b=n \\ a \geq 0, b \geq 0}} \mathcal{M}_{a, b} \longrightarrow \mathcal{M}_{\left(p^{n}\right)}
$$

Theorem 4.1 ([KM, 13.4.6]). - The morphism $\varphi$ is an isomorphism outside the supersingular locus. The inverse image of a supersingular geometric point $x \in \mathcal{M}\left(\overline{\mathbb{F}}_{p}\right)$ in $\mathcal{M}_{\left(p^{n}\right)}\left(\mathbb{F}_{p}\right)$ consists of precisely one point $\tilde{x}$ and the completed local ring of $\tilde{x}$ is isomorphic to

$$
\overline{\mathbb{F}}_{p}\left[[X, Y] / \prod_{\substack{a+b=n \\ a \geq 0, b \geq 0}}\left(X^{p^{a}}-Y^{p^{b}}\right)\right.
$$

in such a way that $\mathcal{M}_{a, b}$ is defined by the equation $X^{p^{a}}-Y^{p^{b}}=0$.
Recall the ideal $I_{i}$ in $W[[t]]$ defining the divisors $\mathcal{T}_{i}$, for $i=1,2$. By the Kummer congruence there exist for $i=1$ and 2 uniformizers $t_{i}$ of $\left.\overline{\mathbb{F}}_{p}[t t]\right]$ and $t_{i}^{\prime}$ of $\overline{\mathbb{F}}_{p}\left[\left[t^{\prime}\right]\right]$ and generators $g_{i}$ of $I_{i}$ such that

$$
g_{i} \equiv\left(t_{i}-\left(t_{i}^{\prime}\right)^{p^{a_{i}}}\right) \cdot\left(t_{i}^{p}-\left(t_{i}^{\prime}\right)^{p^{a_{i}-1}}\right) \cdot \ldots \cdot\left(t_{i}^{p_{i}^{a_{i}}}-t_{i}^{\prime}\right)(\bmod p)
$$

Hence $\mathcal{T}_{i} \cap \mathcal{S}_{(p)}$ is the union of irreducible components $\mathcal{V}_{i \mu}\left(\mu=0,1, \ldots, a_{i}\right)$, where $\mathcal{V}_{i \mu}$ is the divisor in $\mathcal{S}_{(p)}=\operatorname{Spf} \overline{\mathbb{F}}_{p}\left[t t, t^{\prime}\right]$ defined by $t_{i}^{p^{\mu}}-\left(t_{i}^{\prime}\right)^{p^{a_{i}-\mu}}$. Hence

$$
\begin{equation*}
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}=\sum_{\mu=0}^{a_{1}} \sum_{\nu=0}^{a_{2}}\left(\mathcal{V}_{1 \mu} \cdot \mathcal{V}_{2 \nu}\right)_{\mathcal{S}} \tag{4.1}
\end{equation*}
$$

We write

$$
\begin{aligned}
t_{2} & =u \cdot t_{1}, \\
t_{2}^{\prime} & =u^{\prime} \cdot t_{1}^{\prime}, \\
, & \left.u^{\prime} \in \overline{\mathbb{F}}_{p}[\llbracket t]_{p}\left[t^{\prime}\right]\right]^{\times} .
\end{aligned}
$$

Lemma 4.2. - Let $a_{1} \equiv a_{2}(\bmod 2)$. Then $u(0), u^{\prime}(0) \in \mathbb{F}_{p^{2}}$ and $u(0) \neq u^{\prime}(0)^{p^{a_{2}}}$.
Lemma 4.3. - We have

$$
\left(\mathcal{V}_{1 \mu} \cdot \mathcal{V}_{2 \nu}\right)=p^{n}
$$

with $n=\min \left\{a_{1}-\mu+\nu, a_{2}-\nu+\mu\right\}$.
It is an elementary matter to use Lemma 4.3 to calculate the sum on the RHS of (4.1). The result is Proposition 1.8.

Proof of Lemma 4.3 (assuming Lemma 4.2). - We must show

$$
\begin{equation*}
\lg \overline{\mathbb{F}}_{p}\left[t, t^{\prime}\right] /\left(t^{p^{\mu}}-\left(t^{\prime}\right)^{p^{a_{1}-\mu}},(u t)^{p^{\nu}}-\left(u^{\prime} \cdot t^{\prime}\right)^{p^{a_{2}-\nu}}\right)=p^{n} . \tag{4.2}
\end{equation*}
$$

By symmetry it suffices to consider the following two cases.
Case 1: $\mu \leq a_{1}-\mu, \nu \leq a_{2}-\nu$
Case 2: $\mu \leq a_{1}-\mu, a_{2}-\nu \leq \nu$.
In case 1 the LHS of (4.2) is equal to

$$
\begin{aligned}
& \lg \overline{\mathbb{F}}_{p}\left[\llbracket t, t^{\prime}\right] /\left(t-t^{\prime p^{a_{1}-2 \mu}}\right)^{p^{\mu}},\left(u t-\left(u^{\prime} t^{\prime}\right)^{p^{a_{2}-2 \nu}}\right)^{p^{\nu}} \stackrel{(1)}{=} \\
& \quad p^{\mu+\nu} \cdot \lg \overline{\mathbb{F}}_{p}\left[t^{\prime}\right] /\left(u \cdot t^{\prime p^{a_{1}-2 \mu}}-\left(u^{\prime} t^{\prime}\right)^{p^{a_{2}-2 \nu}}\right) \stackrel{(2)}{=} p^{\mu+\nu+\min \left\{a_{1}-2 \mu, a_{2}-2 \nu\right\}}=p^{n}
\end{aligned}
$$

Here in (1) we used the formula ([Go2, Lemma 4.2])

$$
\lg _{A} B / x_{1} \ldots x_{n}=\sum_{i} \lg _{A} B / x_{i}
$$

valid for any $A$-algebra $B$ and non zero divisors $x_{1}, \ldots, x_{n}$ in $B$. In (2) we used Lemma 4.2 which implies that if $a_{1}-2 \mu=a_{2}-2 \nu$, then $u(0) \neq u^{\prime}(0)^{p^{a_{2}-2 \nu}}=u^{\prime}(0)^{p^{a_{2}}}$.

In case 2 , the LHS of (4.2) is equal to

$$
\begin{aligned}
& \left.\lg \overline{\mathbb{F}}_{p} \llbracket t, t^{\prime}\right] /\left(\left(t-t^{\prime p^{a_{1}-2 \mu}}\right)^{p^{\mu}},\left(u^{\prime} t^{\prime}-(u t)^{p^{2 \nu-a_{2}}}\right)^{p^{a_{2}-\nu}}\right)= \\
& p^{\mu} \cdot p^{a_{2}-\nu} \cdot \lg \mathbb{F}_{p}\left[t t, t^{\prime}\right] /\left(t-t^{\prime p_{1}-2 \mu}, u^{\prime} t^{\prime}-(u t)^{p^{2 \nu-a_{2}}}\right)= \\
& \left.\quad p^{a_{2}-\nu+\mu} \cdot \lg \overline{\mathbb{F}}_{p}\left[t^{\prime}\right]\right] /\left(u^{\prime} t^{\prime}-u^{p^{2 \nu-a_{2}}} \cdot t^{\prime p^{a_{1}-2 \mu+2 \nu-a_{2}}}\right) \stackrel{(3)}{=} p^{a_{2}-\nu+\mu}=p^{n} .
\end{aligned}
$$

Here in (3) we used Lemma 4.2: if $a_{1}-2 \mu+2 \nu-a_{2}=0$, then $a_{1}=2 \mu$ and $a_{2}=2 \nu$ are both even and $u^{\prime}(0) \neq u(0)^{p^{2 \nu-a_{2}}}=u(0)$.

Proof of Lemma 4.2. - Let $\ell=\left(a_{2}-a_{1}\right) / 2$. Let
$I_{1}^{\prime}=$ minimal ideal in $W\left[\left[t, t^{\prime}\right]\right.$ such that $p^{\ell} \psi_{1}$ lifts to an isogeny $\Gamma \longrightarrow \Gamma^{\prime}\left(\bmod I_{1}^{\prime}\right)$.
By the Kummer congruence we can choose uniformizers $t_{1}$ of $\left.\overline{\mathbb{F}}_{p} \llbracket t\right]$ and $t_{1}^{\prime}$ of $\overline{\mathbb{F}}_{p}\left[t^{\prime} \rrbracket\right.$ and a generator $g_{1}^{\prime}$ of $I_{1}^{\prime}$ with

$$
g_{1}^{\prime} \equiv\left(t_{1}-t_{1}^{\prime p^{a_{2}}}\right) \cdot\left(t_{1}^{p}-t_{1}^{\prime p^{a_{2}-1}}\right) \cdot \ldots \cdot\left(t_{1}^{p_{2}}-t_{1}^{\prime}\right)(\bmod p) .
$$

Now $\psi_{2}=\alpha \circ\left(p^{\ell} \psi_{1}\right)$, where $\alpha \in \operatorname{Aut}(G)$. By the universal property of $\Gamma$ there exists a unique $W$-algebra homomorphism $h: W[t t] \rightarrow W[[t]$ such that $\alpha$ lifts to an isomorphism

$$
\tilde{\alpha}: \Gamma \longrightarrow h_{*}(\Gamma) .
$$

Hence $I_{2}$ is generated by $g_{2}^{\prime}$ with

$$
\begin{equation*}
g_{2}^{\prime} \equiv\left(h\left(t_{1}\right)-t_{1}^{\prime p^{a_{2}}}\right) \cdot\left(h\left(t_{1}\right)^{p}-t_{1}^{\prime p^{a_{2}-1}}\right) \cdot \ldots \cdot\left(h\left(t_{1}\right)^{p^{a_{2}}}-t_{1}^{\prime}\right)(\bmod p) . \tag{4.3}
\end{equation*}
$$

The two elements $g_{2}^{\prime}$ and $g_{2}$ differ by a unit and

$$
\begin{equation*}
g_{2} \equiv\left(u t_{1}-\left(u^{\prime} t_{1}^{\prime}\right)^{p^{a_{2}}}\right) \cdot\left(\left(u \cdot t_{1}\right)^{p}-\left(u^{\prime} t_{1}^{\prime}\right)^{p^{a_{2}-1}}\right) \cdot \ldots \cdot\left(\left(u t_{1}\right)^{p^{a_{2}}}-u^{\prime} t_{1}^{\prime}\right)(\bmod p) . \tag{4.4}
\end{equation*}
$$

The first factor on the RHS of (4.4) is irreducible and can only divide the first factor of the RHS of (4.3). Hence the first factors differ by a unit. Let

$$
\begin{equation*}
h\left(t_{1}\right) \equiv v \cdot t_{1}(\bmod p) \text { with } v \in \overline{\mathbb{F}}_{p}\left[\left[t_{1}\right]^{\times},\right. \tag{4.5}
\end{equation*}
$$

and put $c=v(0)$. Comparing coefficients we obtain

$$
c=u(0) / u^{\prime}(0)^{p^{a_{2}}} .
$$

The remaining factors on the RHS of (4.4) are not irreducible: $\left(u t_{1}\right)^{p^{\mu}}-\left(u^{\prime} t_{1}^{\prime}\right)^{p^{a_{2}-\mu}}$ is the $p^{\nu}$-th power of an irreducible element, where $\nu=\min \left\{\mu, a_{2}-\mu\right\}$. An analogous comparison of coefficients gives

$$
c=u(0) / u^{\prime}(0)^{p^{a_{2}-2 \mu}} \quad, \quad \mu=0, \ldots, a_{2} .
$$

It follows that $u^{\prime}(0) \in \mathbb{F}_{p^{2}}$ and by symmetry $u(0) \in \mathbb{F}_{p^{2}}$. It remains to show $c \neq 1$. Now $c$ is the induced action of $\alpha$ on the tangent space of the universal deformation of $G$ over $\overline{\mathbb{F}}_{p}$. And $\alpha$ is given in terms of the formal group law by

$$
\alpha(X)=\alpha_{1} X+\alpha_{2} X^{2}+\ldots, \alpha_{i} \in \overline{\mathbb{F}}_{p}
$$

Then $\alpha_{1} \in \mathbb{F}_{p^{2}}=\mathcal{O}_{D} / \Pi \mathcal{O}_{D}$ is the residue class of $\alpha$. By the lemma below, the action of $\alpha$ on the tangent space of the universal deformation space is by multiplication by $\alpha_{1} / \bar{\alpha}_{1}$. Hence $c=\alpha_{1} / \bar{\alpha}_{1}$. But $\alpha_{1} \notin \mathbb{F}_{p}$ and hence $c \neq 1$. Indeed, otherwise for any $a \in \mathbb{Z}_{p}$ with residue class $\alpha_{1}$ modulo $p$, we would have

$$
\begin{equation*}
v\left(\psi_{2}-a p^{\ell} \psi_{1}\right)>v\left(\psi_{2}\right) \tag{4.6}
\end{equation*}
$$

But the optimal basis $\psi_{1}, \psi_{2}, \psi_{3}$ may be chosen so that $\psi_{2}$ has maximal valuation in its residue class modulo $\mathbb{Z}_{p} \psi_{1}$. Indeed, if $p \neq 2$, any optimal basis has this property (otherwise an easy application of Hensel's lemma would imply that $L$ is isotropic).

If $p=2$, we take the optimal basis constructed in table 1 of the appendix. By assumption $a_{1} \equiv a_{2}(\bmod 2)$. Going through all cases in table 1, we see that this can only happen in cases A2 and B3 a). In the case A2, we have $v\left(\psi_{1}\right)>v\left(\psi_{2}\right)$ which contradicts (4.6). In the case B3 a), we get

$$
\left(\psi_{2}-a p^{\ell} \psi_{1}, \psi_{2}-a p^{\ell} \psi_{1}\right)=2^{\beta_{2}}\left(u_{1}+u_{2}+4\left(a^{2}-a\right) u_{1}\right)
$$

which has valuation $a_{2}=v\left(\psi_{2}\right)=\beta_{2}+2$, since in this case $u_{1}+u_{2} \equiv 4(\bmod 8)$.

Remark 4.4. - In fact, even for $p=2$, it is true that any optimal basis has the property that $\psi_{2}$ has maximal valuation in its residue class modulo $\mathbb{Z}_{p} \psi_{1}$. This follows from [B, Prop. 6.9].

Lemma 4.5. - Let $\alpha \in \mathcal{O}_{D}^{\times}=\operatorname{Aut}(G)$, with action on Lie $G$ given by (multiplication by) $\alpha_{1} \in \mathbb{F}_{p^{2}}$. The induced action of $\alpha$ on the tangent space of the universal deformation space of $G$ over $\overline{\mathbb{F}}_{p}$ is by multiplication by $\alpha_{1} / \bar{\alpha}_{1}$.

Here we denote by $x \mapsto \bar{x}$ the non-trivial automorphism of $\mathbb{F}_{p^{2}}$.
Proof (comp. [Z, Lemma 78]). - The tangent space can be canonically identified with

$$
\operatorname{Hom}\left(\operatorname{Lie}^{t} G, \operatorname{Lie} G\right)
$$

For $\varphi \in \operatorname{Hom}\left(\operatorname{Lie}{ }^{t} G\right.$, Lie $\left.G\right)$ we have

$$
\alpha_{*}(\varphi)=\alpha_{1} \circ \varphi \circ{ }^{t} \alpha_{1}^{-1} .
$$

Identifying ${ }^{t} G$ with $G$ replaces ${ }^{t} \alpha_{1}$ by the residue class of ${ }^{\iota} \alpha$, i.e., by $\bar{\alpha}_{1}$.

## A. Appendix: The case $p=2$

In sections 2 and 3 we made the assumption $p>2$. In this appendix we treat the case $p=2$. In this case one has to take into account the delicate theory of quadratic forms over $\mathbb{Z}_{2}$. We will proceed according to the following table. The table gives

- the normal form of the quadratic space $(L, Q)$ in terms of a suitable basis $e_{1}, e_{2}, e_{3}$ (we give the matrix $T=\left(\frac{1}{2}\left(e_{i}, e_{j}\right)\right)$ ),
- an optimal basis $\psi_{1}, \psi_{2}, \psi_{3}$,
- the Gross-Keating invariants $\left(a_{1}, a_{2}, a_{3}\right)$ of $(L, Q)$.

We go through all cases of anisotropic ternary lattices, according to the table in [Y1, appendix B], comp. also [B, Thm. 5.7].

## Table 1

A) $T=\operatorname{diag}\left(u_{1} 2^{\alpha}, 2^{\beta}\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)\right), \alpha \geq 0, \beta \geq-1, \alpha \equiv \beta \bmod 2$
(the condition $\alpha \equiv \beta \bmod 2$ is due to the anisotropy of $T$, comp. [ $\mathbf{B}$, section 5]).

1) $\alpha \leq \beta+1$. Then $\psi_{1}=e_{1}, \psi_{2}=e_{2}, \psi_{3}=e_{3}$ and

$$
G K(T)=(\alpha, \beta+1, \beta+1)
$$

2) $\alpha>\beta+1$. Then $\psi_{1}=e_{2}, \psi_{2}=e_{3}, \psi_{3}=e_{1}$ and

$$
G K(T)=(\beta+1, \beta+1, \alpha)
$$

B) $T=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}}\right)$ with $0 \leq \beta_{1} \leq \beta_{2} \leq \beta_{3}$ This matrix is anisotropic if and only if

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

cf. [Y2], or [B, section 5].

1) $\beta_{2} \not \equiv \beta_{1} \bmod 2$. Then $\psi_{1}=e_{1}, \psi_{2}=e_{2}, \psi_{3}=c_{1} e_{1}+c_{2} e_{2}+e_{3}$ for suitable $c_{1}, c_{2} \in \mathbb{Z}_{2}$, and

$$
G K(T)=\left(\beta_{1}, \beta_{2}, \beta_{3}+2\right)
$$

2) $\beta_{2} \equiv \beta_{1} \bmod 2$ and $\beta_{3} \leq \beta_{2}+1$.
a) $\beta_{3}=\beta_{2}$. Then $\psi_{1}=e_{1}, \psi_{2}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}, \psi_{3}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{3}$ and

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)
$$

b) $\beta_{3}=\beta_{2}+1$ and $u_{1} \equiv u_{2} \bmod 4$. Then $\psi_{1}=e_{1}, \psi_{2}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}$, $\psi_{3}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}+e_{3}$ and

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)
$$

c) $\beta_{3}=\beta_{2}+1$ and $u_{1} \equiv-u_{2} \bmod 4$. Then $\psi_{1}=e_{1}, \psi_{2}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+$ $e_{2}+e_{3}, \psi_{3}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}+2 e_{3}$ and

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)
$$

3) $\beta_{2} \equiv \beta_{1} \bmod 2$ and $\beta_{3} \geq \beta_{2}+2$.
a) $u_{1} \equiv-u_{2} \bmod 4$. Then $\psi_{1}=e_{1}, \psi_{2}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}, \psi_{3}=e_{3}$ and

$$
G K(T)=\left(\beta_{1}, \beta_{2}+2, \beta_{3}\right)
$$

b) $u_{1} \equiv u_{2} \bmod 4$. Then $\psi_{1}=e_{1}, \psi_{2}=2^{\frac{\beta_{2}-\beta_{1}}{2}} \cdot e_{1}+e_{2}, \psi_{3} c_{1} e_{1}+c_{2} e_{2}+e_{3}$ for suitable $c_{1}, c_{2} \in \mathbb{Z}_{2}$, and

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)
$$

A.1. The induction start. - Let $a_{3} \leq 1$, i.e., $a_{1}=0$ and $a_{3}=1$. We follow the proof of Proposition 1.6 in each of the following cases.

- $\quad T=\operatorname{diag}\left(u_{1} 2,2^{-1}\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)\right)$, hence $G K(T)=(0,0,1)$.

Then $\varphi_{2}={ }^{`} \psi_{1} \circ \psi_{2}={ }^{`} e_{2} \circ e_{3}$ and

$$
\operatorname{tr}\left(\varphi_{2}\right)=\left(e_{2}, e_{3}\right)=1 \text { and } \operatorname{Nm}\left(\varphi_{2}\right)=1
$$

Hence $K=\mathbb{Q}_{2}\left(\varphi_{2}\right)=\mathbb{Q}_{2}[X] /\left(X^{2}-X+1\right)$ is an unramified extension of $\mathbb{Q}_{2}$, and $\mathcal{O}_{K}=\mathbb{Z}_{2}\left[\varphi_{2}\right]$. Therefore $\Gamma\left(\bmod J_{2}\right)$ is a canonical lifting relative to $K$.

Now $\varphi_{3}={ }^{\iota} \psi_{1} \circ \psi_{3}={ }^{\iota} e_{2} \circ e_{1}$ and

$$
\operatorname{tr}\left(\varphi_{3}\right)=0 \text { and } \operatorname{Nm}\left(\varphi_{3}\right)=u_{1} \cdot 2
$$

Furthermore

$$
\operatorname{tr}\left(\varphi_{2} \circ{ }^{\iota} \varphi_{3}\right)=\operatorname{tr}\left({ }^{\iota} e_{3} \circ e_{2} \circ{ }^{\iota} e_{2} \circ e_{1}\right)=Q\left(e_{2}\right) \cdot \operatorname{tr}\left({ }^{\iota} e_{3} \circ e_{1}\right)=\left(e_{1}, e_{3}\right)=0 .
$$

Hence $-\varphi_{2} \circ \varphi_{3}+\varphi_{3} \circ{ }^{\iota} \varphi_{2}=0$, i.e., $\varphi_{3}$ anticommutes with $K$. Since $K / \mathbb{Q}_{2}$ is unramified, an application of Lemma 2.2 gives

$$
\varphi_{3} \in \Pi \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{2} \mathcal{O}_{D}\right)
$$

Hence, applying [Ww1, Thm. 1.4],

$$
\lg W[t]] /\left(J_{2}+J_{3}\right)=\frac{1+1}{2}=1=\frac{a_{3}+1}{2},
$$

which proves the claim in this case.

- $\quad T=\operatorname{diag}\left(u_{1},\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)\right)$, hence $G K=(0,1,1)$.

Then $\varphi_{2}={ }^{\iota} \psi_{1} \circ \psi_{2}={ }^{\iota} e_{1} \circ e_{2}$ and

$$
\operatorname{tr}\left(\varphi_{2}\right)=\left(e_{1}, e_{2}\right)=0 \text { and } \operatorname{Nm}\left(\varphi_{2}\right)=u_{1} \cdot 2
$$

Hence $K=\mathbb{Q}_{2}\left(\varphi_{2}\right)=\mathbb{Q}_{2}[X] /\left(X^{2}+u_{1} 2\right)$ is a ramified extension of $\mathbb{Q}_{2}$, and $\mathcal{O}_{K}=$ $\mathbb{Z}_{2}\left[\varphi_{2}\right]$. Therefore $\Gamma\left(\bmod J_{2}\right)$ is a canonical lifting relative to $K$.

Now $\varphi_{3}={ }^{'} \psi_{1} \circ \psi_{3}={ }^{\iota} e_{1} \circ e_{3}$ and

$$
\operatorname{tr}\left(\varphi_{3}\right)=0 \text { and } \operatorname{Nm}\left(\varphi_{3}\right)=u_{1} \cdot 2
$$

Furthermore

$$
\operatorname{tr}\left(\varphi_{2} \circ{ }^{\iota} \varphi_{3}\right)=\operatorname{tr}\left({ }^{\iota} e_{2} \circ e_{1} \circ{ }^{\iota} e_{1} \circ e_{3}\right)=u_{1} \cdot\left(e_{2}, e_{3}\right)=u_{1} \cdot 2 .
$$

Hence

$$
\begin{equation*}
\varphi_{2} \circ \varphi_{3}+\varphi_{3} \circ \varphi_{2}=-u_{1} \cdot 2 \tag{A.1.1}
\end{equation*}
$$

We use the presentation of $D$ resp. $\mathcal{O}_{D}$ from [G, Prop. 4.3]. Namely, assume that the different $\mathcal{D}$ of $K / \mathbb{Q}_{2}$ has valuation equal to $e$. Then

$$
\begin{equation*}
D=K \oplus K \cdot j \tag{A.1.2}
\end{equation*}
$$

where $j$ anticommutes with $K$ and where $j^{2} \in \mathbb{Z}_{2}^{\times}$satisfies $v\left(j^{2}-1\right)=2(e-1)$. Let $\pi$ be a uniformizer in $K$. Then $\alpha:=\pi^{-2}(1+j) \in \mathcal{O}_{D}^{\times}$and

$$
\begin{equation*}
\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \alpha \tag{A.1.3}
\end{equation*}
$$

In the case at hand the extension $K / \mathbb{Q}_{2}$ is wildly ramified, with different $\mathcal{D}$ of valuation $e=3$. Hence $v\left(j^{2}-1\right)=4$. As uniformizer $\pi$ we take $\varphi_{2}$.

Write $\varphi_{3}=a+b \alpha$. Then

$$
\begin{aligned}
\varphi_{2} \circ \varphi_{3}+\varphi_{3} \circ \varphi_{2} & =(a \pi+b \pi \alpha)+(a \pi+b \alpha \pi) \\
& =\left(a \pi+b \pi^{-1}+b \pi^{-1} j\right)+\left(a \pi+b \pi^{-1}-b \pi^{-1} j\right) \\
& =2 \cdot\left(a \pi+b \pi^{-1}\right)
\end{aligned}
$$

Comparing with (A.1.1) we get

$$
a \pi+b \pi^{-1}=-u_{1}
$$

Hence $v(b)=1$, i.e., $\varphi_{3} \in \Pi \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{2} \mathcal{O}_{D}\right)$. Applying [Ww1, Thm. 1.4], we obtain

$$
\lg W[t t] /\left(J_{2}+J_{3}\right)=1+1=2=a_{2}+a_{3}
$$

which proves the claim in this case.

- $\quad T=\operatorname{diag}\left(u_{1}, u_{2}, u_{3}\right)$, hence $G K(T)=(0,1,1)$.

Then $\varphi_{2}={ }^{\iota} \psi_{1} \circ \psi_{2}={ }^{\iota} e_{1} \circ\left(e_{1}+e_{2}\right)={ }^{\iota} e_{1} \circ e_{2}+u_{1} \cdot 1$, hence

$$
\operatorname{tr}\left(\varphi_{2}\right)=u_{1} \cdot 2, \quad \operatorname{Nm}\left(\varphi_{2}\right)=u_{1} \cdot\left(u_{2}+u_{1}\right)
$$

Hence $K=\mathbb{Q}_{2}\left(\varphi_{2}\right)=\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1} \cdot\left(u_{2}+u_{1}\right)\right)$. Since $T$ is anisotropic we have $u_{2}+u_{1} \equiv 2 \bmod 4$. Hence we are dealing with an Eisenstein polynomial and $\mathcal{O}_{K}=\mathbb{Z}_{2}\left[\varphi_{2}\right]$.

Now $\varphi_{3}={ }^{\dagger} \psi_{1} \circ \psi_{3}={ }^{\dagger} e_{1} \circ\left(e_{1}+e_{3}\right)={ }^{\dagger} e_{1} \circ e_{3}+u_{1} \cdot 1$. At this point it is advantageous to consider instead of $\varphi_{3}$ the endomorphism $\varphi_{3}^{\prime}={ }^{\imath} e_{1} \circ e_{3}$. It is obvious that the locus where $\varphi_{2}$ and $\varphi_{3}$ deform is the same as the locus where $\varphi_{2}$ and $\varphi_{3}^{\prime}$ deform. Now

$$
\operatorname{tr}\left(\varphi_{3}^{\prime}\right)=\left(e_{1}, e_{3}\right)=0 \text { and } \operatorname{Nm}\left(\varphi_{3}^{\prime}\right)=u_{1} u_{3}
$$

Furthermore

$$
\begin{aligned}
\operatorname{tr}\left({ }^{\iota} \varphi_{2} \circ \varphi_{3}^{\prime}\right) & =\operatorname{tr}\left({ }^{\iota}\left(e_{1}+e_{2}\right) \circ e_{1} \circ{ }^{\iota} e_{1} \circ e_{3}\right) \\
& =u_{1} \cdot\left(\left(e_{1}, e_{3}\right)+\left(e_{2}, e_{3}\right)\right) \\
& =0 .
\end{aligned}
$$

Hence

$$
{ }^{\iota} \varphi_{2} \circ \varphi_{3}^{\prime}-\varphi_{3}^{\prime} \circ \varphi_{2}=0
$$

Hence $\varphi_{3}^{\prime}$ anticommutes with $K$. Writing, as in the previous case, $D=K \oplus K \cdot j$ we have

$$
\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \alpha
$$

Here $\alpha=\pi^{-1}(1+j) \in \mathcal{O}_{D}$. Indeed, $\pi=\varphi_{2}$ is a uniformizer for $K$ and the different $\mathcal{D}$ has valuation $e=2$. Writing $\varphi_{3}^{\prime}=a+b \alpha$ we get

$$
a+b \alpha=\left(a+b \pi^{-1}\right)+b \pi^{-1} \cdot j .
$$

Hence $a+b \pi^{-1}=0$. Since $\varphi_{3}^{\prime} \in \mathcal{O}_{D}^{\times}$it follows that the valuation of $b$ is equal to 1 , hence

$$
\varphi_{3}^{\prime} \in \Pi \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\Pi^{2} \mathcal{O}_{D}\right)
$$

Applying now [Ww1, Thm. 1.4], we get

$$
\lg W\left[[t] /\left(J_{2}+J_{3}\right)=1+1=a_{2}+a_{3},\right.
$$

which proves the claim in this case. The induction start is now complete.
A.2. The induction step: Lemma 3.1. - In this section we prove Lemma 3.1. We go through all cases of the table.

Case A1: Here $\operatorname{tr}(\varphi)=0$ and $\operatorname{Nm}(\varphi)=u_{1} \cdot 2^{\alpha+\beta+1}$.
Since $\alpha+\beta+1$ is odd, we get $K=\mathbb{Q}_{2}\left(\sqrt{-u_{1} 2}\right)$ and $\mathcal{O}_{K}=\mathbb{Z}_{2}\left[\sqrt{-u_{1} 2}\right]$ and, since $\varphi=2^{\frac{\alpha+\beta}{2}} \cdot \pi$, where $\pi=\sqrt{-u_{1} 2}$ is a uniformizer, $\mathbb{Z}_{2}[\varphi] \mathbb{Z}_{2}+2^{\frac{\alpha+\beta}{2}} \cdot \mathcal{O}_{K}$. Hence the conductor of $\mathbb{Z}_{2}[\varphi]$ is equal to $\frac{\alpha+\beta}{2}=\left[\frac{\alpha+(\beta+1)}{2}\right]=\left[\frac{a_{1}+a_{2}}{2}\right]$.

Case A2: Here $\operatorname{tr}(\varphi)=2^{\beta+1}$ and $\operatorname{Nm}(\varphi)=2^{2(\beta+1)}$.
Hence $K=\mathbb{Q}_{2}[X] /\left(X^{2}-X+1\right)$ is an unramified extension and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\xi]$, where $\xi$ is the residue class of $X$. Then $\varphi=2^{\beta+1} \cdot \xi$ and $\mathbb{Z}_{2}[\varphi]=\mathbb{Z}_{2}+2^{\beta+1} \cdot \mathcal{O}_{K}$ has conductor $\beta+1=\left[\frac{(\beta+1)+(\beta+1)}{2}\right]\left[\frac{a_{1}+a_{2}}{2}\right]$.

Case B1: Here $\operatorname{tr}(\varphi)=0$ and $\operatorname{Nm}(\varphi)=u_{1} u_{2} \cdot 2^{\beta_{1}+\beta_{2}}$.
Since $\beta_{1}+\beta_{2}$ is odd, we have $K=\mathbb{Q}_{2}\left(\sqrt{-u_{1} u_{2} 2}\right)$ and $\mathcal{O}_{K}=\mathbb{Z}_{2}\left[\sqrt{-u_{1} u_{2} 2}\right]$. Now $\mathbb{Z}_{2}[\varphi]=\mathbb{Z}_{2}+2^{\frac{\beta_{1}+\beta_{2}-1}{2}} \cdot \mathcal{O}_{K}$ has conductor $\frac{\beta_{1}+\beta_{2}-1}{2}=\left[\frac{\beta_{1}+\beta_{2}}{2}\right]=\left[\frac{a_{1}+a_{2}}{2}\right]$.

Case B2 a): Here $\operatorname{tr}(\varphi)=u_{1} \cdot 2^{\frac{\beta_{1}+\beta_{2}}{2}+1}$ and $\operatorname{Nm}(\varphi) u_{1} \cdot 2^{\beta_{1}+\beta_{2}}\left(u_{1}+u_{2}\right)$.
Now by the anisotropy condition on $T$ we have $u_{1}+u_{2} \equiv 2 \bmod 4$, hence $K=$ $\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1}\left(u_{1}+u_{2}\right)\right)$ is defined by an Eisenstein polynomial and $\mathcal{O}_{K}=$ $\mathbb{Z}_{2}[\pi]$, where $\pi$ denotes the residue class of $X$. Then $\varphi 2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \pi$ and $\mathbb{Z}_{2}[\varphi]=\mathbb{Z}_{2}+$ $2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \mathcal{O}_{K}$ has conductor $\frac{\beta_{1}+\beta_{2}}{2}=\left[\frac{\beta_{1}+\left(\beta_{2}+1\right)}{2}\right]=\left[\frac{a_{1}+a_{2}}{2}\right]$.

Case B2 b): This is identical with the previous case.

Case B2 c): Here $\operatorname{tr}(\varphi)=u_{1} \cdot 2^{\frac{\beta_{1}+\beta_{2}}{2}+1}$ and $\operatorname{Nm}(\varphi) u_{1}^{2} \cdot 2^{\beta_{1}+\beta_{2}}+u_{1} u_{2} \cdot 2^{\beta_{1}+\beta_{2}}+u_{1} u_{3}$. $2^{\beta_{1}+\beta_{3}} u_{1} 2^{\beta_{1}+\beta_{2}}\left(2 u_{3}+u_{2}+u_{1}\right)$.

Hence $K=\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1} \cdot\left(2 u_{3}+u_{2}+u_{1}\right)\right)$, which is defined by an Eisenstein equation since $u_{1}+u_{2} \equiv 0 \bmod 4$. Hence $\mathcal{O}_{K}=\mathbb{Z}_{2}[\pi]$, where $\pi$ is the residue class of $X$ and $\varphi=2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \pi$ and $\mathbb{Z}_{2}[\varphi]=\mathbb{Z}_{2}+2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \mathcal{O}_{K}$ has conductor $\frac{\beta_{1}+\beta_{2}}{2}\left[\frac{\beta_{1}+\beta_{2}+1}{2}\right]=\left[\frac{a_{1}+a_{2}}{2}\right]$.
Case B3 a): Here $\operatorname{tr}(\varphi)=2^{\frac{\beta_{1}+\beta_{2}}{2}+1} \cdot u_{1}$ and $\operatorname{Nm} \varphi=2^{\beta_{1}+\beta_{2}} \cdot u_{1}\left(u_{1}+u_{2}\right)$.
Hence $K=\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1} \cdot\left(u_{1}+u_{2}\right)\right)$. Now since $T$ is anisotropic, it follows that $u_{1}+u_{2} \equiv 4 \bmod 8$. Hence writing $u_{1}+u_{2}=4 \eta$ with $\eta \in \mathbb{Z}_{2}^{\times}$, we have $K=\mathbb{Q}_{2}\left[X_{1}\right] /\left(X_{1}^{2}-u_{1} X_{1}+u_{1} \eta\right)$. Hence $K / \mathbb{Q}_{2}$ is unramified and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\xi]$, where $\xi$ denotes the residue class of $X_{1}$. Now $\varphi=2^{\frac{\beta_{1}+\beta_{2}}{2}+1} \cdot \xi$ and $\mathbb{Z}_{2}[\varphi]=\mathbb{Z}_{2}+2^{\frac{\beta_{1}+\beta_{2}}{2}+1} \cdot \mathcal{O}_{K}$ has conductor $\frac{\beta_{1}+\beta_{2}}{2}+1 \frac{\beta_{1}+\left(\beta_{2}+2\right)}{2}=\left[\frac{a_{1}+a_{2}}{2}\right]$.

Case B3 b): Here the trace and norm are as in the previous case, but this time $K=\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1} \cdot\left(u_{1}+u_{2}\right)\right)$ is defined by an Eisenstein polynomial. Hence $\mathcal{O}_{K}=\mathbb{Z}_{2}[\pi]$ where $\pi$ is the residue class of $X$ and $\varphi=2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \pi$ and $\mathbb{Z}_{2}[\varphi]=$ $\mathbb{Z}_{2}+2^{\frac{\beta_{1}+\beta_{2}}{2}} \cdot \mathcal{O}_{K}$ has conductor $\frac{\beta_{1}+\beta_{2}}{2} \frac{\beta_{1}+\left(\beta_{2}+1\right)}{2}=\left[\frac{a_{1}+a_{2}}{2}\right]$.

This proves the assertion in all cases.

By symmetry we also obtain that $\varphi^{\prime}=\psi_{2} \circ{ }^{\iota} \psi_{1}$ generates an order of conductor $\left[\frac{a_{1}+a_{2}}{2}\right]$ in $K^{\prime}$.
A.3. The induction step: Lemmas 3.4 and 3.5. - We first prove Lemma 3.4. We go through all cases, making use of the results in section A.2. Again we wish to bound the conductors of the orders $\mathbb{Z}_{2}\left[\varphi_{1}\right]$ resp. $\mathbb{Z}_{2}\left[\varphi_{2}\right]$.

Case A1: Here $K=\mathbb{Q}_{2}\left(\sqrt{-u_{1} 2}\right)$ and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\pi]$ with $\pi=\sqrt{-u_{1} 2}$. Then ${ }^{\iota} \pi=-\pi$ and thereby this case is like the ramified case for $p \neq 2$. We have

$$
\mathcal{O}_{K}=\mathbb{Z}_{2} \oplus \mathbb{Z}_{2} \pi
$$

the decomposition into traceful and traceless elements. In particular, $\operatorname{tr}\left(\mathcal{O}_{K}\right) \subset 2 \cdot \mathbb{Z}_{2}$. Let

$$
\varphi_{i}^{\circ}=\varphi_{i}-\frac{1}{2} \operatorname{tr}\left(\varphi_{i}\right) \quad, \quad i=1,2
$$

Then $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathbb{Z}_{2}\left[\varphi_{i}^{\circ}\right]$ has conductor $\frac{1}{2}\left(v\left(\varphi_{i}^{\circ}\right)-1\right)$. Let

$$
\varphi_{i}=\lambda_{i} \cdot \pi^{a_{i}} \quad, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times}
$$

Then

$$
\varphi_{i}^{\circ}=\frac{1}{2}\left(\lambda_{i}-(-1)^{a_{i}} \cdot{ }^{\iota} \lambda_{i}\right) \cdot \pi^{a_{i}} \quad, \quad i=1,2 .
$$

Writing $\lambda_{i}=a+b \pi$ we have $a \in \mathbb{Z}_{p}^{\times}$and

$$
\begin{aligned}
\lambda_{i}+{ }^{\iota} \lambda_{i} & =2 a \\
\lambda_{i}-{ }^{\iota} \lambda_{i} & =2 b \pi .
\end{aligned}
$$

Hence $v\left(\varphi_{i}^{\circ}\right)=a_{i}$ if $a_{i}$ is odd and $v\left(\varphi_{i}^{\circ}\right)>a_{i}$ if $a_{i}$ is even. Now according to our table, $a_{1}$ and $a_{2}$ have different parity which implies that $r \leq\left(a_{2}-1\right) / 2$. This shows the result in this case.

Case A2: Here $K=\mathbb{Q}_{2}[X] /\left(X^{2}-X+1\right)$ and $\mathcal{O}_{K} \mathbb{Z}_{2}[\xi]$, where $\xi$ is the residue class of $X$. In this case, $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathcal{O}_{K}$ or $\operatorname{tr}\left(\varphi_{i}\right) \in 2 \mathbb{Z}_{2}$. In the first case $r=0$ and the claim is obvious. Now let $\operatorname{tr}\left(\varphi_{i}\right) \in 2 \mathbb{Z}_{2}$ for $i=1$ and $i=2$, and consider

$$
\varphi_{i}^{\circ}=\varphi_{i}-\frac{1}{2} \operatorname{tr}\left(\varphi_{i}\right)
$$

Then writing $\varphi_{i}^{\circ}=a+b \xi$ we have $0=\operatorname{tr}\left(\varphi_{i}^{\circ}\right)=2 a+b$. Hence $\varphi_{i}^{\circ}=a \cdot(1-2 \xi)$ and $v\left(\varphi_{i}^{\circ}\right)=v(a)=v(b)-2$. The conductor of $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathbb{Z}_{2}\left[\varphi_{i}^{\circ}\right]$ is equal to $\frac{1}{2} v(b)=$ $\frac{1}{2} v\left(\varphi_{i}^{\circ}\right)+1$. Now

$$
\begin{aligned}
\varphi_{i} & =\lambda_{i} \cdot 2^{a_{i} / 2} \quad, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times} . \\
\varphi_{i}^{\circ} & =\frac{1}{2}\left(\lambda_{i}-{ }^{\iota} \lambda_{i}\right) \cdot 2^{a_{i} / 2} .
\end{aligned}
$$

Hence $v\left(\varphi_{i}^{\circ}\right)=a_{i}-2$ if the residue class $\left[\lambda_{i}\right]$ of $\lambda_{i}$ lies in $\mathbb{F}_{4} \backslash \mathbb{F}_{2}$, and is larger otherwise. Hence if $\left[\lambda_{i}\right] \in \mathbb{F}_{4} \backslash \mathbb{F}_{2}$, then $r \leq \frac{a_{i}}{2}$, hence $2 r \leq a_{2}$.

But not both $\left[\lambda_{1}\right],\left[\lambda_{2}\right]$ can lie in $\mathbb{F}_{2}$. Indeed,

$$
\varphi_{1} \circ \varphi_{2}={ }^{\iota} \lambda_{1} \cdot \lambda_{2} \cdot 2^{\left(a_{1}+a_{2}\right) / 2}={ }^{\iota} \lambda_{1} \lambda_{2} \cdot 2^{\beta+1} .
$$

On the other hand ${ }^{\iota} \varphi_{1} \circ \varphi_{2}=u \cdot \varphi$, where $u \in \mathbb{Z}_{2}^{\times}$is as in (3.1), and where $\varphi$ is as in the previous section. Now $\varphi=2^{\beta+1} \xi$. Taking the residues modulo $2^{\beta+1}$, we prove the claim.

Now assume $2 r=a_{2}=a_{3}$. Then $a_{1}<2 r$ has to be odd, which contradicts the fact that $a_{1}=a_{2}=\beta+1$.

Case B1: Here $K=\mathbb{Q}_{2}\left(\sqrt{-u_{1} u_{2} 2}\right)$ and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\pi]$, with $\pi=\sqrt{-u_{1} u_{2} 2}$. This case is completely analogous to case A1.

Case B2 a): Here $K=\mathbb{Q}_{2}[X] /\left(X^{2}-2 u_{1} X+u_{1} \cdot\left(u_{1}+u_{2}\right)\right)$ and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\pi]$ where $\pi$ denotes the residue class of $X$. Then $\pi$ is a uniformizer satisfying an Eisenstein equation. Hence $\operatorname{tr}\left(\mathcal{O}_{K}\right) \subset 2 \mathbb{Z}_{2}$. We again consider $\varphi_{i}^{\circ}=\varphi_{i}-\frac{1}{2} \operatorname{tr}\left(\varphi_{i}\right)$. Then writing $\varphi_{i}^{\circ}=a+b \pi$ we have $0=\operatorname{tr}\left(\varphi_{i}^{\circ}\right)=2 a+2 b u_{1}=2\left(a+b u_{1}\right)$. Hence $\varphi_{i}^{\circ}=b \cdot\left(-u_{1}+\pi\right)$ and $v\left(\varphi_{i}^{\circ}\right)=v(b)$. The conductor of $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathbb{Z}_{2}\left[\varphi_{i}^{\circ}\right]$ is equal to $\frac{1}{2} v(b)=\frac{1}{2} v\left(\varphi_{i}^{\circ}\right)$. Now

$$
\varphi_{i}=\lambda_{i} \cdot \pi^{a_{i}}, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times} .
$$

Let us write

$$
2 u_{1}-\pi=\eta \cdot \pi \quad, \quad \eta \in \mathcal{O}_{K}^{\times} .
$$

Then $\eta=1+\eta_{1} \cdot \pi$ with $\eta_{1} \in \mathcal{O}_{K}^{\times}$. We have

$$
\begin{aligned}
\operatorname{tr}\left(\varphi_{i}\right) & =\lambda_{i} \cdot \pi^{a_{i}}+{ }^{\iota} \lambda_{i} \cdot(\eta \pi)^{a_{i}} \\
& =\left(\lambda_{i}+{ }^{\iota} \lambda_{i} \eta^{a_{i}}\right) \cdot \pi^{a_{i}} .
\end{aligned}
$$

Hence

$$
\varphi_{i}^{\circ}=\frac{1}{2} \cdot\left(\lambda_{i}-{ }^{\iota} \lambda_{i} \eta^{a_{i}}\right) \cdot \pi^{a_{i}} .
$$

Let $\lambda_{i} \equiv 1+\left[\lambda_{i}\right] \cdot \pi\left(\bmod \pi^{2}\right)$. Then ${ }^{\iota} \lambda_{i} \equiv 1-\left[\lambda_{i}\right] \pi\left(\bmod \pi^{2}\right)$. If $a_{i}$ is odd, we get

$$
\begin{aligned}
\lambda_{i}-{ }^{\iota} \lambda_{i} \eta^{a_{i}} & \equiv\left(1+\left[\lambda_{i}\right] \pi\right)-\left(1-\left[\lambda_{i}\right] \pi\right) \cdot\left(1+\eta_{1} \pi\right)\left(\bmod \pi^{2}\right) \\
& \equiv \eta_{1} \cdot \pi\left(\bmod \pi^{2}\right)
\end{aligned}
$$

Hence in this case $v\left(\varphi_{i}^{\circ}\right)=a_{i}-1$. We get $r \leq \frac{1}{2}\left(a_{i}-1\right)$. Since $a_{1}$ or $a_{2}$ are odd, we obtain the assertion.

Cases B2 b) and c): These cases are identical to the previous one.
Case B3 a): In this case $K=\mathbb{Q}_{2}[X] /\left(X^{2}-u_{1} X+u_{1} \eta\right)$, for some $\eta \in \mathbb{Z}_{2}^{\times}$. Hence $K / \mathbb{Q}_{2}$ is unramified and $\mathcal{O}_{K}=\mathbb{Z}_{2}[\xi]$, where $\xi$ is the residue class of $X$. This case is similar and almost identical to case A2). If $\operatorname{tr}\left(\varphi_{i}\right) \notin 2 \mathbb{Z}_{2}$, then $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathcal{O}_{K}$ and $r=0$ and the claim is obvious. If $\operatorname{tr}\left(\varphi_{i}\right) \in 2 \mathbb{Z}_{2}$ for $i=1$ and $i=2$, we consider again $\varphi_{i}^{\circ}=\varphi_{i}-\frac{1}{2} \operatorname{tr}\left(\varphi_{i}\right)$. Writing $\varphi_{i}^{\circ}=a+b \xi$ we get $0=\operatorname{tr}\left(\varphi_{i}^{\circ}\right)=2 a+b u_{1}$. Hence $\varphi_{i}^{\circ}=a\left(1-2 u_{i}^{-1} \xi\right)$ and $v\left(\varphi_{i}^{\circ}\right)=v(a)=v(b)-2$. The conductor of $\mathbb{Z}_{2}\left[\varphi_{i}\right]=\mathbb{Z}_{2}\left[\varphi_{i}^{\circ}\right]$ is equal to $\frac{1}{2} v(b)=\frac{1}{2} v\left(\varphi_{i}^{\circ}\right)+1$. Now

$$
\begin{aligned}
\varphi_{i} & =\lambda_{i} \cdot 2^{a_{i} / 2}, \quad \lambda_{i} \in \mathcal{O}_{K}^{\times} \\
\varphi_{i}^{\circ} & =\frac{1}{2}\left(\lambda_{i}-{ }^{\iota} \lambda_{i}\right) \cdot 2^{a_{i} / 2} .
\end{aligned}
$$

Hence $v\left(\varphi_{i}^{\circ}\right)=a_{i}-2$ if the residue class $\left[\lambda_{i}\right]$ of $\lambda_{i}$ lies in $\mathbb{F}_{4} \backslash \mathbb{F}_{2}$, and is larger otherwise. If $\left[\lambda_{i}\right] \in \mathbb{F}_{4} \backslash \mathbb{F}_{2}$, then $r \leq a_{i} / 2$, hence $2 r \leq a_{2}$. But not both $\left[\lambda_{1}\right],\left[\lambda_{2}\right]$ can lie in $\mathbb{F}_{2}$. Indeed,

$$
{ }^{\iota} \varphi_{1} \circ \varphi_{2}={ }^{\iota} \lambda_{1} \lambda_{2} \cdot 2^{\left(a_{1}+a_{2}\right) / 2}={ }^{\iota} \lambda_{1} \lambda_{2} \cdot 2^{\frac{\beta_{1}+\beta_{2}}{2}+1}=u \cdot \varphi_{2}=u \cdot 2^{\frac{\beta_{1}+\beta_{2}}{2}+1} \cdot \xi .
$$

Taking the residue modulo $2^{\frac{\beta_{1}+\beta_{2}}{2}+1}$, we get the claim.
Now assume $2 r=a_{2}=a_{3}$. Then $a_{1}<2 r$ has to be odd which contradicts the condition that $a_{1}=\beta_{1}$ has to have the same parity as $\beta_{2}+2=a_{2}=2 r$.

Case B3 b): This is again identical to cases B 2 a )-c).
The Lemma 3.4 is proved.

We now turn to the proof of Lemma 3.5. Again we inspect the various cases.
Case A1: We write $D=K \oplus K \cdot j$ as in (A.1.2) in section A.1, where $j$ anticommutes with $K$ and where $j^{2} \in \mathbb{Z}_{2}^{\times}$satisfies $v\left(j^{2}-1\right)=2(e-1)$, where the different $\mathcal{D}$ has
valuation $e$. Then $\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \delta \alpha$, where $\alpha=\pi^{-(e-1)} \cdot(1+j) \in \mathcal{O}_{D}^{\times}$, cf. (A.1.3). In the case at hand $e=3$, hence $\alpha=\pi^{-2}(1+j)$. Now

$$
\begin{equation*}
\varphi \circ{ }^{\iota} \varphi_{3}+{ }^{\iota} \varphi_{3} \circ \varphi 2^{\beta+1} \cdot{ }^{\iota} \tilde{e}_{1} \tag{A.3.1}
\end{equation*}
$$

where $\tilde{e}_{1}=\gamma \circ e_{1}$. This follows from $\varphi={ }^{\iota} \psi_{1} \circ \psi_{2}$ and the definitions $\psi_{1} e_{1}, \psi_{2}=e_{2}$, $\psi_{3}=e_{3}$. Now writing ${ }^{\iota} \varphi_{3}=a+b \alpha$ for suitable $a, b \in \mathcal{O}_{K}$ and writing $\varphi=2^{\delta} \cdot \pi$ with $\delta=\frac{1}{2}(\alpha+\beta)$, we get from (A.3.1)

$$
2^{\delta} \cdot \pi(a+b \alpha)+(a+b \alpha) \cdot 2^{\delta} \pi 2^{\beta+1} \cdot{ }^{\iota} \tilde{e}_{1}
$$

i.e., $2^{\delta+1}\left(a \pi+b \pi^{-1}\right)=2^{\beta+1} .{ }^{\iota} \tilde{e}_{1}$, hence

$$
\begin{equation*}
b=2^{\frac{\beta-\alpha}{2}} \cdot{ }^{\iota} \tilde{e}_{1} \pi-a \pi^{2} \tag{A.3.2}
\end{equation*}
$$

Now $v\left({ }^{\prime} \tilde{e}_{1}\right)=\alpha$, hence the first summand of the RHS of (A.3.2) has valuation $\beta+1$. Since $v\left(\varphi_{3}\right)=\beta+1$, it follows $v(b)=\beta+1=a_{3}$, which proves the claim in this case.

Case A2: Here we write $\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \Pi$ where $\Pi^{2}=2$ and where $\Pi$ anticommutes with $K$. In this case we have

$$
\varphi \circ{ }^{\iota} \varphi_{3}+{ }^{\iota} \varphi_{3} \circ \varphi 2^{\beta+1} \cdot{ }^{\iota} \tilde{e}_{1}
$$

Writing ${ }^{\iota} \varphi_{3}=a+b \Pi$ and $\varphi 2^{\beta+1} \cdot \xi$ we obtain

$$
2^{\beta+1}\left(2 a \xi+b\left(\xi+{ }^{`} \xi\right) \cdot \Pi\right)=2^{\beta+1} \cdot{ }^{\iota} \tilde{e}_{1}
$$

i.e.,

$$
2 a \xi+b \Pi={ }^{\imath} \tilde{e}_{1} .
$$

Now $v\left({ }^{\iota} \tilde{e}_{1}\right)=v\left({ }^{\iota} \varphi_{3}\right)=\alpha$. This implies $v(b \Pi)=\alpha$, hence $\varphi_{3} \in \Pi^{a_{3}} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\right.$ $\left.\Pi^{a_{3}+1} \mathcal{O}_{D}\right)$, since $a_{3}=\alpha$.
Case B1: This case is similar to case A1, except that the identity (A.3.1) is replaced by

$$
\varphi \circ{ }^{\iota} \varphi_{3}+{ }^{\iota} \varphi_{3} \circ \varphi 2 \cdot{ }^{\iota} \tilde{e}_{3} \circ \varphi .
$$

Now $\varphi=2^{\delta} \pi$ with $\delta=\left(\beta_{1}+\beta_{2}\right) / 2$. Writing as in case A1) ${ }^{\iota} \varphi_{3}=a+b \alpha$, where $\alpha=\pi^{-2}(1+j)$, we get

$$
2^{\delta+1} \cdot\left(a \pi+b \pi^{-1}\right)=2^{\delta+1} \cdot{ }^{\iota} \tilde{e}_{3} \cdot \pi
$$

i.e.,

$$
b={ }^{\iota} \tilde{e}_{3} \pi^{2}-a \pi^{2} .
$$

Now the first summand of the RHS has valuation $\beta_{3}+2$ and $v\left(\tilde{\psi}_{3}\right)=\beta_{3}+2$. Hence $v(b)=\beta_{3}+2$, which proves the claim, since $\beta_{3}+2=a_{3}$.
Case B2 a): In this case the valuation of the different is equal to 2 and hence $\alpha=$ $\pi^{-1} \cdot(1+j)$. Now

$$
\varphi \circ{ }^{\iota} \varphi_{3}-{ }^{\iota} \varphi_{3} \circ \varphi 2 \cdot{ }^{\iota} e_{2} e_{3}{ }^{\iota} \tilde{e}_{1} .
$$

Writing ${ }^{\iota} \varphi_{3}=a+b \alpha$ and $\varphi=2^{\delta} \pi$ with $\delta=\left(\beta_{1}+\beta_{2}\right) / 2$, we get

$$
\begin{equation*}
2^{\delta}\left(((\pi a+b)+b j)-\left((\pi a+b)+b \pi^{-1} \cdot{ }^{\iota} \pi j\right)\right)=2 \cdot{ }^{\iota} e_{2} e_{3}{ }^{\iota} \tilde{e}_{1} \tag{A.3.3}
\end{equation*}
$$

Therefore, since ${ }^{\iota} \pi=2 u_{1}-\pi$,

$$
2^{\delta+1} j \cdot b \cdot\left(1-u_{1} \pi^{-1}\right)=2 \cdot{ }^{\iota} e_{2} e_{3}{ }^{\iota} \tilde{e}_{1}
$$

Comparing valuations we obtain $v(b)=\beta_{3}+1=a_{3}$, which proves the assertion in this case.

Case B2 b): Here again $\alpha=\pi^{-1}(1+j)$, and the same equation (A.3.3) holds. The case is identical with the previous case.

Case B2 c): The same again.
Case B3 a): This case is similar to case A2. We write $\mathcal{O}_{D}=\mathcal{O}_{K} \oplus \mathcal{O}_{K} \cdot \Pi$ as in that case. Now

$$
\varphi \circ{ }^{\iota} \varphi_{3}-{ }^{\iota} \varphi_{3} \circ \varphi=-2 \cdot{ }^{\iota} e_{3} e_{2}{ }^{\iota} \tilde{e}_{1}
$$

We write ${ }^{\iota} \varphi_{3}=a+b \Pi$ and $\varphi=2^{\delta} \cdot \xi$ where $\delta=\frac{\beta_{1}+\beta_{2}}{2}+1$ and $\xi$ satisfies $\xi^{2}-u_{1} \xi+u_{1} \eta=$ 0 for some $\eta \in \mathbb{Z}_{2}^{\times}$. Then

$$
2^{\delta} \cdot\left((a \xi+b \xi \Pi)-\left(a \xi+b \cdot{ }^{\iota} \xi \Pi\right)\right)-2 \cdot{ }^{\iota} e_{3} e_{2}{ }^{\iota} \tilde{e}_{1}
$$

Now $\xi-{ }^{`} \xi=2 \xi-u_{1}$, hence

$$
2^{\delta} \cdot b \cdot \Pi \cdot\left(2 \xi-u_{1}\right)=-2 \cdot{ }^{\iota} e_{3} e_{2}{ }^{\iota} \tilde{e}_{1}
$$

Comparing valuations we get $v(b)=\beta_{3}-1=a_{3}-1$. Hence $\varphi_{3} \in \Pi^{a_{3}} \mathcal{O}_{D} \backslash\left(\mathcal{O}_{K}+\right.$ $\Pi^{a_{3}+1} \mathcal{O}_{D}$ ), as claimed.

Case B3 b): This case is similar to cases B2 a)-c). Again the valuation of the different is equal to 2 and $\alpha=\pi^{-1}(1+j)$. Now

$$
\varphi \circ{ }^{\iota} \varphi_{3}-{ }^{\iota} \varphi_{3} \circ \varphi 2 \cdot{ }^{\iota} e_{2} e_{3}{ }^{\iota} \tilde{e}_{1}
$$

Writing ${ }^{\iota} \varphi_{3}=a+b \alpha$ and $\varphi 2^{\delta} \cdot \pi$ with $\delta=\left(\beta_{1}+\beta_{2}\right) / 2$ as in case B2 a), we get just as in that case

$$
2^{\delta+1} \cdot j \cdot b\left(1-u_{1} \pi^{-1}\right)=2^{\iota} e_{2} e_{3}{ }^{\iota} \tilde{e}_{1}
$$

Comparing valuations we get $v(b)=\beta_{3}+1=a_{3}$, which proves the assertion in this case.
A.4. Lemma 1.9. - The proof of Lemma 1.9 for $p \neq 2$ was very easy. By contrast, the case $p=2$ is quite elaborate and uses more information than used so far on the construction of an optimal basis. We go through all cases of the table 1. It turns out that in the passage from the type $T$ of $L$ to the type $T^{\prime}$ of $L^{\prime}$ a number of things can happen, as can be read off from the following table.

Table 2

| Type $T$ | Type $T^{\prime}$ |
| :---: | :---: |
| $\begin{array}{ll} \text { A1 } & \alpha \neq \beta \\ & \alpha=\beta \end{array}$ | $\begin{aligned} & \text { B1 } \\ & \text { B2 b) } \end{aligned}$ |
| A2 | A2 |
| $\begin{array}{ll} \text { B1 } & \beta_{2} \leq \beta_{3}-2 \\ & \beta_{2}=\beta_{3}-1 \\ & \beta_{2}=\beta_{3} \end{array}$ | B1 <br> B2 b) or c) <br> A 1 or B 2 a ) |
| $\begin{array}{ll} \mathrm{B} 2 \text { a) } & \beta_{1}<\beta_{2} \\ & \beta_{1}=\beta_{2} \end{array}$ | $\begin{aligned} & \mathrm{B} 3 \mathrm{~b}) \\ & \mathrm{B} 2 \mathrm{~b}) \text { or } \mathrm{c}) \end{aligned}$ |
| $\begin{array}{ll} \mathrm{B} 2 \mathrm{~b}) & \beta_{1}<\beta_{2} \\ & \beta_{1}=\beta_{2} \end{array}$ | $\begin{aligned} & \text { B3 a) } \\ & \text { A2 } \end{aligned}$ |
| $\begin{array}{ll} \text { B2 c) } & \beta_{1}<\beta_{2} \\ & \beta_{1}=\beta_{2} \end{array}$ | $\begin{aligned} & \text { B3 a) } \\ & \text { A2 } \end{aligned}$ |
|  | $\begin{aligned} & \text { B3 a) } \\ & \text { B2 c) } \end{aligned}$ |
| B3 b) $\begin{aligned} & \beta_{3} \geq \beta_{2}+4 \\ & \beta_{3}=\beta_{2}+3 \\ & \beta_{3}=\beta_{2}+2 \end{aligned}$ | $\begin{aligned} & \text { B3 b) } \\ & \text { B2 b) } \\ & \text { B2 a) } \end{aligned}$ |

The calculations exhibit in fact not only the type of $T^{\prime}$ but also the precise normal form of $T^{\prime}$ from which one can then read off the Gross-Keating invariants of $T^{\prime}$. In all cases, the assertion of Lemma 1.9 is confirmed.

Since these calculations in the 16 cases are quite tedious, we will sometimes be brief.

Case A1: Here $G K=(\alpha, \beta+1, \beta+1)$, and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, e_{2}, e_{3}\right)$. Hence $\psi_{3}^{\prime} \frac{1}{2} e_{3}$, so

$$
T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\alpha}, 2^{\beta-1}\left(\begin{array}{ll}
4 & 1 \\
1 & 1
\end{array}\right)\right)
$$

Since

$$
2^{\beta-1}\left(\begin{array}{ll}
4 & 1 \\
1 & 1
\end{array}\right) \sim \operatorname{diag}\left(3 \cdot 2^{\beta-1}, 3 \cdot 2^{\beta-1}\right)
$$

we obtain

$$
T^{\prime} \sim \begin{cases}\operatorname{diag}\left(u_{1} \cdot 2^{\alpha}, 3 \cdot 2^{\beta-1}, 3 \cdot 2^{\beta-1}\right) & \text { if } \alpha \neq \beta \\ \operatorname{diag}\left(3 \cdot 2^{\alpha-1}, 3 \cdot 2^{\alpha-1}, u_{1} \cdot 2^{\alpha}\right) & \text { if } \alpha=\beta\end{cases}
$$

Hence if $\alpha \neq \beta$, and since $\alpha \equiv \beta \bmod 2$, then $T^{\prime}$ is of type B 1 and $G K\left(T^{\prime}\right)=(\alpha, \beta-$ $1, \beta+1)$ as asserted. If $\alpha=\beta$, then $T^{\prime}$ is of type B 2 b$)$ and $G K\left(T^{\prime}\right)=(\alpha-1, \alpha, \alpha+1)$, as asserted.

The case A2 is entirely similar.
Case B1: In this case $\operatorname{GK}(T)=\left(\beta_{1}, \beta_{2}, \beta_{3}+2\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, e_{2}, c_{1} e_{1}+\right.$ $c_{2} e_{2}+e_{3}$ ) for suitable $c_{1}, c_{2} \in \mathbb{Z}_{2}$. If $\beta_{2}<\beta_{3}$, then by [ $\mathbf{Y} 1$, proof of Lemma B.6], both coefficients $c_{1}$ and $c_{2}$ are divisible by 2 . Hence $L^{\prime}$ is generated by $\left(e_{1}, e_{2}, \frac{1}{2} e_{3}\right)$. Hence the matrix of $L^{\prime}$ in terms of this basis is

$$
T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}-2}\right)
$$

So if $\beta_{2} \leq \beta_{3}-2$, the type of $T^{\prime}$ is B 1 and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}, \beta_{3}\right)$ as asserted. If $\beta_{2}=\beta_{3}-1$, then $T^{\prime}$ is of type B 2 b$)$ or c$)$ and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}, \beta_{3}\right)$ as asserted.

If $\beta_{2}=\beta_{3}$, then by [ $\mathbf{Y 1}$, proof of Lemma B.6], we have $2 \mid c_{1}$. On the other hand, we have $2 \not \backslash c_{2}$ in this case, because otherwise the valuation of $\frac{1}{2}\left(\psi_{3}, \psi_{3}\right)$ would be $\beta_{2}<a_{3}=\beta_{2}+2$ which is impossible. Hence $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2}\left(e_{2}+e_{3}\right)$. Consider the matrix defined by the basis $e_{2}, \frac{1}{2}\left(e_{2}+e_{3}\right)$ of the lattice $\tilde{L}^{\prime}$ of rank 2 generated by $e_{2}$ and $\frac{1}{2}\left(e_{2}+e_{3}\right)$,

$$
\tilde{T}^{\prime}\left(\begin{array}{cl}
u_{2} 2^{\beta_{2}} & c_{2} u_{2} 2^{\beta_{2}-1} \\
* & \left(c_{2}^{2} u_{2}+u_{3}\right) 2^{\beta_{2}-2}
\end{array}\right) .
$$

We determine when $\tilde{T}^{\prime}$ is diagonalizable by determining the valuations of the ideals in $\mathbb{Z}_{2}$,

$$
s\left(\tilde{L}^{\prime}\right)=\frac{1}{2}\left(\tilde{L}^{\prime}, \tilde{L}^{\prime}\right), \text { resp. } n\left(\tilde{L}^{\prime}\right)=\left(Q(x), x \in \tilde{L}^{\prime}\right)
$$

Now

$$
\operatorname{ord} s(L)=\min \left\{\beta_{2}, \beta_{2}-1, \quad \operatorname{ord}\left(c_{2}^{2} u_{2}+u_{3}\right)+\beta_{2}-2\right\}=\beta_{2}-1 .
$$

And

$$
\begin{aligned}
\text { ord } n(L) & =\min \left\{\beta_{2}, \beta_{2}, \operatorname{ord}\left(c_{2}^{2} u_{2}+u_{3}\right)+\beta_{2}-2\right\} \\
& = \begin{cases}\beta_{2}-1 & \text { if } u_{2} \equiv u_{3} \bmod 4 \\
\beta_{2} & \text { if } u_{2} \equiv-u_{3} \bmod 4\end{cases}
\end{aligned}
$$

Hence, by [Y1, Prop. B.3],

$$
\tilde{T}^{\prime} \sim \begin{cases}\operatorname{diag}\left(\eta_{1} 2^{\beta_{2}-1}, \eta_{2} 2^{\beta_{2}-1}\right) & \text { if } u_{2} \equiv u_{3} \bmod 4 \\
2^{\beta_{2}-1} \cdot\left(\begin{array}{ll}
2 & 1 \\
1 & 2
\end{array}\right) & \text { if } u_{2} \equiv-u_{3} \bmod 4\end{cases}
$$

Here $\eta_{1}, \eta_{2} \in \mathbb{Z}^{\times}$. For the total matrix $T^{\prime}$ we get that if $u_{2} \equiv u_{3} \bmod 4$, then $T^{\prime} \sim \operatorname{diag}\left(u_{1} 2^{\beta_{1}}, \eta_{1} 2^{\beta_{2}-1}, \eta_{2} 2^{\beta_{2}-1}\right)$ is of type B 2 a) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}, \beta_{2}\right)$ as asserted. If $u_{2} \equiv-u_{3} \bmod 4$, then $T^{\prime} \sim \operatorname{diag}\left(u_{1} 2^{\beta_{1}}, 2^{\beta_{2}-1}\left(\begin{array}{c}2 \\ 1 \\ 1\end{array} \frac{1}{2}\right)\right)$ is of type A1 and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}, \beta_{2}\right)$ as asserted.

Case B2 a): In this case $G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{2}+1\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, 2^{\gamma} e_{1}+\right.$ $\left.e_{2}, 2^{\gamma} e_{1}+e_{3}\right)$, where $\gamma=\frac{1}{2}\left(\beta_{2}-\beta_{1}\right)$.

If $\gamma>0$, then $L^{\prime}$ is generated by the elements $e_{1}, e_{2}, \frac{1}{2} e_{3}$ and it follows that $T^{\prime}=$ $\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{3} 2^{\beta_{2}-2}, u_{2} 2^{\beta_{2}}\right)$. Now by the anisotropy condition we have

$$
\left(-1, u_{2} u_{3}\right)=\left(u_{1} u_{2}, u_{1} u_{3}\right)
$$

hence $u_{1} \equiv u_{3} \bmod 4$. Therefore $T^{\prime}$ is of type B3 b) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}-1, \beta_{2}+1\right)$, as asserted.

If $\gamma=0$, i.e., $\beta_{1}=\beta_{2}=\beta_{3}=: \beta$, then $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2}\left(e_{1}+e_{3}\right)$ and has matrix with respect to this basis equal to

$$
T^{\prime}=\left(\begin{array}{ccc}
u_{1} 2^{\beta} & 0 & u_{2} 2^{\beta-1} \\
* & u_{2} 2^{\beta} & 0 \\
* & * & \left(u_{1}+u_{3}\right) 2^{\beta-2}
\end{array}\right) .
$$

Now $u_{1} \equiv u_{3} \bmod 4$, hence by an argument similar to the one used in the case B1 when $\beta_{2}=\beta_{3}$, the lattice generated by $e_{1}, \frac{1}{2}\left(e_{1}+e_{3}\right)$ is diagonalizable to $\operatorname{diag}\left(\eta_{1} 2^{\beta-1}, \eta_{2} 2^{\beta-1}\right)$. Hence $T^{\prime} \sim \operatorname{diag}\left(\eta_{1} 2^{\beta-1}, \eta_{2} 2^{\beta-1}, u_{2} 2^{\beta}\right)$ is of type B 2 b$)$ or c) and $G K\left(T^{\prime}\right)=(\beta-1, \beta, \beta+1)$, as asserted.

Case B2 b): In this case $G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{2}+2\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, 2^{\gamma} e_{1}+\right.$ $\left.e_{2}, 2^{\gamma} e_{1}+e_{2}+e_{3}\right)$, with $\gamma=\frac{1}{2}\left(\beta_{2}-\beta_{1}\right)$.

If $\gamma>0$, then $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2}\left(e_{2}+e_{3}\right)$, and has matrix with respect to this basis equal to

$$
T^{\prime}=\left(\begin{array}{ccc}
u_{1} 2^{\beta_{1}} & 0 & 0 \\
* & u_{2} 2^{\beta_{2}} & u_{2} 2^{\beta_{2}-1} \\
* & * & \left(u_{2}+2 u_{3}\right) 2^{\beta_{2}-2}
\end{array}\right)
$$

By an argument similar to the one used in the case B 1 when $\beta_{2}=\beta_{3}$, we see that $T^{\prime} \sim \operatorname{diag}\left(u_{1} 2^{\beta_{1}}, \eta_{1} 2^{\beta_{2}-2}, \eta_{2} 2^{\beta_{2}+1}\right)$, hence $T^{\prime}$ is of type B3. We claim that $T^{\prime}$ is of type B3 a), so that $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}, \beta_{2}+1\right)$, as asserted. But $\eta_{1} \equiv-u_{2} \equiv-u_{1}$ $\bmod 4$, whence the assertion.

There still remains the case when $\gamma=0$, i.e., $\beta_{1}=\beta_{2}=: \beta$ and $\beta_{3}=\beta+1$. Then $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2}\left(e_{1}+e_{2}+e_{3}\right)$. Let $\tilde{L}^{\prime}$ be the sublattice generated by $f_{2}=\frac{1}{2}\left(e_{1}+e_{2}+e_{3}\right)$ and $f_{3}=\frac{1}{2}\left(e_{2}-e_{1}+e_{3}\right)$. Then

$$
\begin{aligned}
\frac{1}{2}\left(f_{2}, f_{2}\right)=\frac{1}{2}\left(f_{3}, f_{3}\right) & =u_{1} 2^{\beta-2}+u_{2} 2^{\beta-2}+u_{3} 2^{\beta-1} \\
& =\left(u_{1}+u_{2}+2 u_{3}\right) 2^{\beta-2} \\
& =\eta \cdot 2^{\beta}
\end{aligned}
$$

Now $\eta \in \mathbb{Z}^{\times}$. Indeed, by the anisotropy condition we have

$$
\left(-1, u_{1} u_{3}\right)\left(2, u_{1} u_{2}\right)
$$

It follows that if $u_{1} \equiv \pm u_{2} \bmod 8$, then $u_{1} \equiv u_{3} \bmod 4$ and if $u_{1} \equiv \pm 3 u_{2} \bmod 8$, then $u_{1} \equiv-u_{3} \bmod 4$. In either case $u_{1}+u_{2}+2 u_{3} \not \equiv 0 \bmod 8$. Similarly,

$$
\begin{aligned}
\frac{1}{2}\left(f_{2}, f_{3}\right) & =-u_{1} 2^{\beta-2}+u_{2} 2^{\beta-2}+u_{3} 2^{\beta-1}=\left(u_{2}-u_{1}+2 u_{3}\right) 2^{\beta-2} \\
& =\kappa \cdot 2^{\beta-1}, \quad \text { with } \kappa \in \mathbb{Z}_{2}^{\times}
\end{aligned}
$$

Now an argument similar to the one used previously shows that the quadratic space $\tilde{L}^{\prime}$ is equivalent to $2^{\beta-1}\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)$. The orthogonal complement of $\tilde{L}^{\prime}$ in $L \otimes_{\mathbb{Z}_{2}} \mathbb{Q}_{2}$ is the line

$$
\left(\tilde{L}^{\prime}\right)^{\perp}=\mathbb{Q}_{2} \cdot\left(-2 \frac{u_{3}}{u_{2}} e_{2}+e_{3}\right)
$$

Now $L^{\prime}$ is generated by $e_{1}+e_{2}$ and $f_{2}$ and $f_{3}$. Hence one easily calculates that

$$
\left(\tilde{L}^{\prime}\right)^{\perp} \cap L^{\prime}=\mathbb{Z}_{2} \cdot f
$$

where $f=-2 \frac{u_{3}}{u_{2}} e_{2}+e_{3}$. Now

$$
\frac{1}{2}(f, f)=\left(\frac{u_{3}}{u_{2}}\right)^{2} 2^{\beta+2}+u_{3} 2^{\beta+1}=\lambda \cdot 2^{\beta+1} \quad, \quad \lambda \in \mathbb{Z}_{2}^{\times}
$$

Hence $\mathbb{Z}_{2} \cdot f+\tilde{L}^{\prime}$ has valuation $(\beta+1)+2(\beta-1)$, equal to the valuation of $L^{\prime}$. Hence $L^{\prime}=\mathbb{Z}_{2} f+\tilde{L}^{\prime}$ is equivalent to $\operatorname{diag}\left(\lambda \cdot 2^{\beta+1}, 2^{\beta-1}\left(\begin{array}{cc}2 & 1 \\ 1 & 2\end{array}\right)\right)$, is of type A2 and $G K\left(T^{\prime}\right)=(\beta, \beta, \beta+1)$, as asserted.

Case B2 c): Here GK $(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{2}+2\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, 2^{\gamma} e_{1}+e_{2}+\right.$ $\left.e_{3}, 2^{\gamma} e_{1}+e_{2}+2 e_{3}\right)$, where $\gamma=\frac{1}{2}\left(\beta_{2}-\beta_{1}\right)$.

When $\gamma>0$, this is similar to previous cases with $L^{\prime}$ generated by $e_{1}, \frac{1}{2} e_{2}, e_{3}$. In this case $T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}-2}, u_{3} 2^{\beta_{2}+1}\right)$ is of type B3 a) and $G K\left(T^{\prime}\right)\left(\beta_{1}, \beta_{2}, \beta_{2}+1\right)$, as asserted.

When $\gamma=0$, i.e., $\beta_{1}=\beta_{2}=: \beta$ and $\beta_{3}=\beta+1$, then $L^{\prime}$ is generated by $e_{1}, \frac{1}{2}\left(e_{1}+\right.$ $\left.e_{2}\right), e_{3}$. Now the quadratic space generated by $e_{1}$ and $\frac{1}{2}\left(e_{1}+e_{2}\right)$ has matrix

$$
\tilde{T}^{\prime}=\left(\begin{array}{cc}
u_{1} 2^{\beta} & u_{1} 2^{\beta-1} \\
* & \left(u_{1}+u_{2}\right) 2^{\beta-2}
\end{array}\right)
$$

Now $\left(u_{1}+u_{2}\right) 2^{\beta-2}=\eta \cdot 2^{\beta}$ with $\eta \in \mathbb{Z}_{2}$. By the usual argument $\tilde{T}^{\prime} \sim 2^{\beta-1}\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)$ and hence $T^{\prime} \sim \operatorname{diag}\left(u_{3} 2^{\beta+1}, 2^{\beta-1}\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)\right)$ is of type A2 with $G K\left(T^{\prime}\right)=(\beta, \beta, \beta+1)$, as asserted.

Case B3 a): In this case $G K(T)=\left(\beta_{1}, \beta_{2}+2, \beta_{3}\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, 2^{\gamma} e_{1}+e_{2}, e_{3}\right)$ with $\gamma=\frac{1}{2}\left(\beta_{2}-\beta_{1}\right)$.

Now $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2} e_{3}$ and has matrix $T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}-2}\right)$. If $\beta_{2}+2 \leq \beta_{3}-2$, then $T^{\prime}$ is of type B 3 a) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+2, \beta_{3}-2\right)$, as asserted. Let $\beta_{3}-2<\beta_{2}+2$. Since not all $G K$-invariants can have the same parity, we have $\beta_{1} \not \equiv \beta_{2} \bmod 2$. Hence $\beta_{3}=\beta_{2}+3$, and $T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{2}+1}\right)$ is of type B2 c) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+1, \beta_{2}+2\right)$, as asserted.
Case B3 b): In this case $G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)$ and $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)=\left(e_{1}, 2^{\frac{\beta_{2}-\beta_{1}}{2}} e_{1}+\right.$ $e_{2}, c_{1} e_{1}+c_{2} e_{2}+e_{3}$ ) for suitable $c_{1}, c_{2} \in \mathbb{Z}_{2}$. In this case we need to extract more information about the coefficients $c_{1}, c_{2}$ from [Y1, proof of Lemma B.8]. If $\beta_{3} \equiv \beta_{1}$ $\bmod 2$, then $c_{1}=2^{\frac{\beta_{3}-\beta_{1}}{2}}$ and $c_{2}=0$. Hence $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2} e_{3}$, hence its matrix is $T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}-2}\right)$. If $\beta_{3}-2 \geq \beta_{2}+2$, then $T^{\prime}$ is of type B 3 b) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}-1\right)$, as asserted. If $\beta_{3}=\beta_{2}+2$, then $T^{\prime}$ is of type B2 a) and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}-1\right)$, as asserted.

If $\beta_{3} \not \equiv \beta_{1} \bmod 2$, then by loc. cit., $c_{1}=2^{\frac{\beta_{3}-\beta_{1}-1}{2}}$ and $c_{2}=2^{\frac{\beta_{3}-\beta_{2}-1}{2}}$. Now $\beta_{3} \geq \beta_{2}+3$, hence $c_{1}$ and $c_{2}$ are divisible by 2 . Hence $L^{\prime}$ is generated by $e_{1}, e_{2}, \frac{1}{2} e_{3}$, and its matrix is $T^{\prime}=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}-2}\right)$. If $\beta_{3} \geq \beta_{2}+4$, then $T^{\prime}$ is of type $\mathrm{B} 3 \mathrm{~b})$ and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}-1\right)$, as asserted. If $\beta_{3}=\beta_{2}+3$, then $T^{\prime}$ is of type B 2 b$)$ and $G K\left(T^{\prime}\right)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}-1\right)$, as asserted.

Lemma 1.9 is now proved in all cases.

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## 14. AN ALTERNATIVE APPROACH USING IDEAL BASES

by
Stefan Wewers


#### Abstract

We give another approach to the proof of the Gross-Keating intersection formula. This approach is based on the concept of ideal bases in the theory of anisotropic quadratic forms over $\mathbb{Z}_{p}$, and in the case $p=2$ is drastically simpler than the proof given in the previous chapter. Résumé (Une approche alternative à l'aide des bases idéales). - On donne une autre approche à la démonstration de la formule de Gross et Keating. Cette approche est basée sur la notion de bases idéales de la théorie des formes quadratiques anisotropes sur $\mathbb{Z}_{p}$, et est plus simple que la démonstration dans le chapitre précédent pour $p=2$.


In this note we give an alternative proof of Proposition 1.5 and Proposition 1.6 of $[\mathbf{R}]$. This proof uses the concept of ideal bases introduced in Section 6 of $[\mathbf{B}]$ and thus avoids the difficulties encountered in the case $p=2$. In fact, our arguments work the same way for any $p$.

## 1. Homomorphisms between quasi-canonical lifts

1.1. Let $p$ be a prime number and $D$ the quaternion division algebra over $\mathbb{Q}_{p}$. The reduced norm gives an anisotropic $\mathbb{Q}_{p}$-valued quadratic form on $D$ which we denote by $Q$. The function $v: D^{\times} \rightarrow \mathbb{Z}, \alpha \mapsto \operatorname{ord}_{p} Q(\alpha)$, is the standard normalized valuation on $D$.

Let $\boldsymbol{\psi}=\left(\psi_{1}, \ldots, \psi_{n}\right)$ be an ordered tuple of linearly independent elements of $D$, and let $L \subset D$ be the $\mathbb{Z}_{p}$-lattice spanned by $\psi$. The restriction of $Q$ to $L$ gives $L$ the structure of an anisotropic quadratic $\mathbb{Z}_{p}$-module. We say that $\psi$ is an ideal basis of $L$ if

$$
v\left(\psi_{i}\right) \leq v\left(\psi_{j}\right) \quad \text { for all } i \leq j
$$

2000 Mathematics Subject Classification. - 14L05, 11F32.
Key words and phrases. - Formal $\mathcal{O}$-modules, quaternion algebras, modular intersection numbers.
and if

$$
v\left(\sum_{i} x_{i} \psi_{i}\right)=\min _{i} v\left(x_{i} \psi_{i}\right)
$$

for all $\left(x_{i}\right) \in \mathbb{Z}_{p}^{n}$. By $[\mathbf{B}]$, Lemma 6.4, this is equivalent to Definition 6.3 of loc. cit.. In particular, every sublattice $L \subset D$ has an ideal basis.

By [ $\mathbf{B}$, Proposition 6.6], an ideal basis is also optimal. Moreover, if $\boldsymbol{\psi}$ is ideal then the numbers $a_{i}:=v\left(\psi_{i}\right), i=1, \ldots, n$, are the Gross-Keating invariants of $L$.
1.2. Let $K \subset D$ be a subfield which is a quadratic extension of $\mathbb{Q}_{p}$. Then there exists an element $\varphi \in K$ such that

$$
\mathcal{O}_{K}=\mathbb{Z}_{p} \oplus \mathbb{Z}_{p} \cdot \varphi
$$

and such that $\varphi$ is a unit (resp. a uniformizer) if $K / \mathbb{Q}_{p}$ is unramified (resp. if $K / \mathbb{Q}_{p}$ is ramified). For such an element, we have

$$
\begin{equation*}
v(x+y \varphi)=\min \left\{2 \operatorname{ord}_{p} x, 2 \operatorname{ord}_{p} y+v(\varphi)\right\} \tag{1.1}
\end{equation*}
$$

for all $x, y \in \mathbb{Q}_{p}$. It follows that $\left(1, p^{r} \varphi\right)$ is an ideal basis of

$$
\mathcal{O}_{r}=\mathbb{Z}_{p} \oplus \mathbb{Z}_{p} \cdot p^{r} \varphi
$$

the unique order in $\mathcal{O}_{K}$ of conductor $p^{r}$, for all $r \geq 0$.
1.3. Let $G$ be the unique formal group of height 2 over $k=\overline{\mathbb{F}}_{p}$. We identify the ring of endomorphisms of $G$ with the maximal order $\mathcal{O}_{D}$ of $D$. Note that for $\psi \in \mathcal{O}_{D}$ the integer $v(\psi)$ is equal to the height of the isogeny $\psi: G \rightarrow G$.

Fix two positive integers $r, s \geq 0$, and let $F_{r}, F_{s}$ be quasi-canonical lifts of $G$ of level $r$ and $s$, with respect to the subfield $K \subset D$. We assume that $F_{r}, F_{s}$ are defined over $A$, a complete discrete valuation ring which is a finite extension of the ring of Witt vectors over $k$. We denote by

$$
H_{r, s}:=\operatorname{Hom}_{A}\left(F_{r}, F_{s}\right)
$$

the group of homomorphisms of formal groups $F_{r} \rightarrow F_{s}$. This is a free $\mathbb{Z}_{p}$-module of rank 2. It is also a right (resp. left) module under the order $\mathcal{O}_{r}=\operatorname{End}\left(F_{r}\right)$ (resp. the order $\mathcal{O}_{s}=\operatorname{End}\left(F_{s}\right)$ ).

Reducing a homomorphism $F_{r} \rightarrow F_{s}$ to the special fibre yields a $\mathbb{Z}_{p}$-linear embed$\operatorname{ding} H_{r, s} \hookrightarrow D$. Via this embedding we may consider $H_{r, s}$ as a quadratic $\mathbb{Z}_{p}$-module.

## Proposition 1.1

1. As a right $\mathcal{O}_{r}$-module, $H_{r, s}$ is free of rank 1, generated by a homomorphism $\psi_{1}: F_{r} \rightarrow F_{s}$ of height $|s-r|$.
2. The Gross-Keating invariants of $H_{r, s}$ are $(|s-r|, r+s)$ if $K / \mathbb{Q}_{p}$ is unramified and $(|s-r|, r+s+1)$ if $K / \mathbb{Q}_{p}$ is ramified.

Proof. - Replacing all isogenies by their duals, we may assume that $r \leq s$. Let $F / A$ be the canonical lift of $G$ with respect to the embedding $K \subset D$. By [Ww1, §4], we may identify $F_{r}$ with the quotient of $F$ corresponding to the superlattice $T_{r} \supset T:=\mathcal{O}_{K}$ defined by

$$
T_{r}:=\mathbb{Z}_{p} \cdot p^{-r}+\mathcal{O}_{K}
$$

(and similarly for $F_{s}$ ). By [ $\mathbf{W} \mathbf{w} \mathbf{1}$, Corollary 2.3], this presentation of $F_{r}, F_{s}$ yields an isomorphism of right $\mathcal{O}_{r}$-modules

$$
H_{r, s} \cong\left\{\alpha \in \mathcal{O}_{K} \mid \alpha T_{r} \subset T_{s}\right\}
$$

We let $\psi_{1} \in H_{r, s}$ denote the element corresponding to 1 under this isomorphism. Clearly, the height of $\psi_{1}$ equals the index of $T_{r}$ in $T_{s}$, which is $s-r$. To prove Part 1 of the proposition, it remains to show that $\alpha T_{r} \subset T_{s}$ if and only if $\alpha \in \mathcal{O}_{r}$. One direction is clear. For the other direction, fix $\alpha \in \mathcal{O}_{K}$ with $\alpha T_{r} \subset T_{s}$. In order to show that $\alpha \in \mathcal{O}_{r}$, we may add any element of $\mathbb{Z}_{p}$ to $\alpha$. Hence we may assume that $\alpha=x \varphi$, where $x \in \mathbb{Z}_{p}$ and $\varphi$ is as in Section 1.2. Our assumption implies that

$$
\alpha p^{-r}=x p^{-r} \varphi \in T_{s}=\mathbb{Z}_{p} \cdot p^{-s} \oplus \mathbb{Z}_{p} \cdot \varphi
$$

We conclude that $p^{r} \mid x$ and hence $\alpha \in \mathcal{O}_{r}$. This proves Part 1.
Set $\psi_{2}:=p^{r} \varphi \psi_{1}$. Clearly, $\left(\psi_{1}, \psi_{2}\right)$ is the basis of $H_{r, s}$ corresponding to the ideal basis $(1, \varphi)$ of $\mathcal{O}_{r}$ under the isomorphism $\mathcal{O}_{r} \cong H_{r, s}$. This isomorphism is not an isometry, but for $\psi=\alpha \cdot \psi_{1} \in H_{r, s}$, with $\alpha \in \mathcal{O}_{r}$, we have

$$
v(\psi)=v(\alpha)+(s-r)
$$

Therefore, it follows from (1.1) that $\left(\psi_{1}, \psi_{2}\right)$ is an ideal basis of $H_{r, s}$. By the choice of $\varphi \in K$ in Section 1.2, we get $v\left(\psi_{2}\right)=s+r$ (resp. $v\left(\psi_{2}\right)=s+r+1$ ) if $K / \mathbb{Q}_{p}$ is unramified (resp. ramified). This completes the proof of Part 2 of the proposition.
1.4. We choose a uniformizer $\lambda$ of the discrete valuation ring $A$. For $n \geq 0$ we set $A_{n}:=A /\left(\lambda^{n+1}\right)$. Let $H_{r, s, n}$ denote the subgroup of $\mathcal{O}_{D}$ consisting of endomorphisms $\psi: G \rightarrow G$ which lift to a homomorphism $F_{r} \otimes A_{n} \rightarrow F_{s} \otimes A_{n}$.

Given an element $\psi \in \mathcal{O}_{D}-H_{r, s}$, we define two integers,

$$
l_{r, s}(\psi):=\max \left\{v(\psi+\phi) \mid \phi \in H_{r, s}\right\}
$$

and

$$
n_{r, s}(\psi):=\max \left\{m \mid \psi \in H_{r, s, m}\right\}
$$

We let $e$ denote the absolute ramification index of the discrete valuation ring $A$.
Proposition 1.2. - There exists a constant $c_{r, s}$, only depending on $(r, s)$, such that the following holds. If $l_{r, s}(\psi) \geq r+s-1$ then

$$
n_{r, s}(\psi)=c_{r, s}+\frac{e}{2} \cdot l_{r, s}(\psi)
$$

Proof. - First we consider the case $r=s$. Then we may assume that $F_{r}=F_{s}$. This is the case studied in $[\mathbf{V l}]$. By ăloc. cit., Proposition 3.1, we have for $l_{r, s}(\psi) \geq 2 r-1$

$$
\begin{equation*}
n_{r, s}(\psi)=a(r-1)+p^{r-1}+\left(\frac{l_{r, s}(\psi)+1}{2}-r\right) e+1, \tag{1.2}
\end{equation*}
$$

where $a(k)=\left(p^{k}-1\right)(p+1) /(p-1)$. Hence the proposition is true for $r=s$.
For the general case, we may again assume that $r \leq s$. By induction on $s$, we will reduce to the case $r=s$. Suppose that the proposition is proved for some pair $(r, s)$ with $r \leq s$. Let $F_{r}, F_{s}, F_{s+1}$ be quasi-canonical lifts of level $r, s, s+1$. We want to prove the proposition for the pair $(r, s+1)$. By Proposition 1.1.1, the group $H_{s, s+1}$ is generated, as a right $\mathcal{O}_{s}$-module, by a homomorphism $\beta: F_{s} \rightarrow F_{s+1}$ of height one. Moreover, the map $\psi \mapsto \beta \psi$ is an isomorphism of $\mathbb{Z}_{p}$-modules $H_{r, s} \xrightarrow{\sim} H_{r, s+1}$.

Let $\psi \in \mathcal{O}_{D}-H_{r, s+1}$ with $l_{r, s+1}(\psi) \geq s+r$. In a first step we will assume in addition that either $r>0$ or that $l_{r, s+1}(\psi) \geq r+s+1$. It is no restriction of generality to assume that $v(\psi)=l_{r, s+1}(\psi)$. Then $v(\psi)>0$ and we can write $\psi=\beta \psi^{\prime}$, with $\psi^{\prime} \in \mathcal{O}_{D}$. It follows from the assertions made in the preceding paragraph that we have

$$
\begin{equation*}
l_{r, s+1}(\psi)=l_{r, s}\left(\psi^{\prime}\right)+1 \tag{1.3}
\end{equation*}
$$

In particular, $l_{r, s}\left(\psi^{\prime}\right) \geq r+s$. On the other hand, [ $\mathbf{W} \mathbf{w} \mathbf{1}$, Corollary 6.3], says that

$$
\begin{equation*}
n_{r, s+1}(\psi)=n_{r, s}\left(\psi^{\prime}\right)+e / e_{s+1} \tag{1.4}
\end{equation*}
$$

where we use the following notation. Let $M=K \cdot W[1 / p]$, and let $\mathcal{O}_{M}$ be its ring of integers. By $M_{s}$ we denote the ring class field of $O_{s}^{\times} \subset \mathcal{O}_{K}^{\times}$, by $\mathcal{O}_{M_{s}}$ its ring of integers, and by $e_{s}$ its absolute ramification index. Then $\mathcal{O}_{M_{s}}$ is the minimal subring of $A$ over which $F_{s}$ can be defined. So for $r>0$, the proposition follows from (1.3), (1.4) and induction.

Unfortunately, for $r=0$ the above argument proves the claim only for the weaker bound $l_{r, s} \geq r+s=s$. The problem is that for $s=1$ and $l=0$ the element $\psi$ is a unit in $\mathcal{O}_{D}$, and so we cannot divide by $\beta$ and reduce to the case $s=0$. However, the argument can be used to compute the value of the constant $c_{r, s}$. For instance, for $(r, s)=(0,0)$ we have $c_{0,0}=e / 2$ by (1.2), and so by (1.3) and (1.4) we get $c_{0,1}=e / e_{1}$. Therefore, the proposition is proved if we can show that for $l_{0,1}(\psi)=0$ we have $n=n_{0,1}(\psi)=e / e_{1}$.

Since $l_{0,1}(\psi)=0$, the endomorphism $\psi$ is an automorphism of $G$. Let $F_{r}^{\psi}$ denote the lift of $G$ obtained from $F_{r}$ by composing the isomorphism $F_{r} \otimes_{A} k \xrightarrow{\sim} G$ with $\psi$. Then $\psi$ lifts to a homomorphism $F_{r} \rightarrow F_{s}$ modulo $\lambda^{n}$ if and only if the two deformations $F_{r}^{\psi} \otimes A /\left(\lambda^{n}\right)$ and $F_{s} \otimes A /\left(\lambda^{n}\right)$ are isomorphic. This, in turn, means that $u\left(F_{r}^{\psi}\right) \equiv u\left(F_{s}\right)\left(\bmod \lambda^{n}\right)$ (here $u(F) \in A$ denotes the modulus of a lift of $G$ defined over $A$ ). By [ $\mathbf{W} \mathbf{w} \mathbf{1}$, Corollary 5.6], the valuation of $u\left(F_{r}^{\psi}\right)$ (resp. of $u\left(F_{s}\right)$ ) is equal to $e / e_{r}$ (resp. equal to $e / e_{s}$ ). Since $e_{r}=e_{0}<e_{s}=e_{1}$, the maximal value that $n$ can take is $e / e_{1}$. This is what we still had to prove.

## 2. The modular intersection number

2.1. Let $p$ be an arbitrary prime and $k=\overline{\mathbb{F}}_{p}$. Let $G$ be the (unique) formal group of height 2 over $k$. We identify $\operatorname{End}_{k}(G)$ with the maximal order $\mathcal{O}_{D}$ of the quaternion division algebra $D$ over $\mathbb{Q}_{p}$. Let $W=W(k)$ denote the ring of Witt-vectors over $k$. Let $\left(\Gamma, \Gamma^{\prime}\right)$ be the universal deformation of the pair of formal groups $(G, G)$. It is defined over the universal deformation space $\mathcal{S} \cong \operatorname{Spf} W\left[\left[t, t^{\prime}\right]\right]$.

Let $L \subset \mathcal{O}_{D}$ be a sub- $\mathbb{Z}_{p}$-module of rank 3 . We denote by $Q$ the quadratic form induced on $L$ by the reduced norm on $\mathcal{O}_{D}$. For $\psi \in L$ we define $v(\psi):=\operatorname{ord}_{p} Q(\psi)$. Choose an ideal basis $\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ of $(L, Q)$, see Section 1.1. Let $a_{i}:=v\left(\psi_{i}\right)$. The numbers $a_{1}, a_{2}, a_{3}$ are the Gross-Keating invariants of $L$.

For $i=1,2,3$, let $\mathcal{T}_{i}$ denote the closed subscheme of $\mathcal{S}$ corresponding to the ideal $I \triangleleft W\left[\left[t, t^{\prime}\right]\right]$ which is minimal for the property that $\psi_{i}$ lifts to a homomorphism $\Gamma \rightarrow \Gamma^{\prime}$ modulo $I$. The following proposition corresponds to Proposition 1.5 of $[\mathbf{R}]$.

Proposition 2.1. - If $a_{3} \leq 1$ then $a_{3}=1$ and

$$
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}= \begin{cases}1, & \text { for } a_{2}=0 \\ 2, & \text { for } a_{2}=1\end{cases}
$$

Proof. - Since $Q$ is anisotropic, the $a_{i}$ cannot have all the same parity. Therefore, $a_{1} \leq a_{2} \leq a_{3} \leq 1$ implies $a_{0}=0$ and $a_{3}=1$. In particular, $\psi_{1}$ is an automorphism of $G$. It follows that $\mathcal{T}_{1} \cong \operatorname{Spf} W[[t]]$, and that we may identify $\left.\Gamma\right|_{\mathcal{T}_{1}}$ with $\left.\Gamma^{\prime}\right|_{\mathcal{T}_{1}}$ via $\psi_{1}$. So for the rest of the proof, we assume that $\psi_{1}=1 \in \mathcal{O}_{D}$ and consider $\mathcal{T}_{2}, \mathcal{T}_{3}$ as closed subschemes of $\mathcal{S}^{\prime}=\operatorname{Spf} W[[t]]$, the universal deformation space of $G$. For $i=2,3, \mathcal{T}_{i}$ is defined by the condition that $\psi_{i}$ lifts to an endomorphism of $\Gamma$.

Let $\mathcal{O}=\mathbb{Z}_{p}\left[\psi_{2}\right] \subset \mathcal{O}_{D}$ denote the subring generated by $\psi_{2}$. Since $\left(\psi_{1}=1, \psi_{2}\right)$ is an ideal basis of $\mathcal{O}$, we have

$$
a_{2}=v\left(\psi_{2}\right)=\max \left\{v\left(x+\psi_{2}\right) \mid x \in \mathbb{Z}_{p}\right\} .
$$

If $a_{2}=0$, then it follows that $\mathcal{O}=\mathcal{O}_{K}$ is the maximal order of $K \subset D$, an unramified quadratic extension of $\mathbb{Q}_{p}$. Therefore, $\mathcal{T}_{2} \cong \operatorname{Spf} W \subset \mathcal{S}^{\prime}$ and $F:=\left.\Gamma\right|_{\tau_{2}}$ is the canonical lift corresponding to the subfield $K \subset D$. Moreover, in the notation of $\S 1.4$ we have $l=l_{0,0}\left(\psi_{3}\right)=v\left(\psi_{3}\right)=a_{3}$. It follows from [ $\mathbf{W w} \mathbf{1}$ ], Theorem 3.3 (see the proof of Proposition 1.2) that $\mathcal{T}_{3} \cap \mathcal{T}_{2} \subset \mathcal{T}_{2}$ corresponds to the ideal $\left(p^{n}\right) \triangleleft W$, with

$$
n=n_{0,0}\left(\psi_{3}\right)=\frac{l+1}{2} e=\frac{a_{3}+1}{2}=1 .
$$

This proves the proposition for $a_{2}=0$.
If $a_{2}=1$, then $\mathcal{O}=\mathcal{O}_{K}$ is also the maximal order of $K$, but $K / \mathbb{Q}_{p}$ is ramified. With the same arguments as above, it follows that $\mathcal{T}_{2} \cong \operatorname{Spf} \mathcal{O}_{M} \subset \mathcal{S}^{\prime}$ is the canonical locus corresponding to the subfield $K \subset D$. Applying again [Ww1], Theorem 3.3, we get

$$
n=n_{0,0}\left(\psi_{3}\right)=\frac{l+1}{2} e=a_{3}+1=2
$$

This proves the proposition for $a_{2}=1$.
2.2. The next proposition corresponds to Proposition 1.6 of $[\mathbf{R}]$.

Proposition 2.2. - Suppose that $\psi_{3}=p \psi_{3}^{\prime}$, for some $\psi_{3}^{\prime} \in \mathcal{O}_{D}$. Let $\mathcal{T}_{3}^{\prime} \subset \mathcal{S}$ be the closed formal subscheme corresponding to $\psi_{3}^{\prime}$ and $\mathcal{S}_{(p)} \subset \mathcal{S}$ the special fiber. Then

$$
\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}=\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{T}_{3}^{\prime}\right)_{\mathcal{S}}+\left(\mathcal{T}_{1} \cdot \mathcal{T}_{2} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}
$$

Proof. - Let $\left(F_{r}, F_{s}\right)$ be a pair of quasi-canonical lifts of $G$ of level $r$ and $s$, with respect to the same subfield $K \subset D$. The set $H_{r, s}:=\operatorname{Hom}\left(F_{r}, F_{s}\right)$ is a sub- $\mathbb{Z}_{p}$-module of $\mathcal{O}_{D}$ of rank two. We consider all pairs $\left(F_{r}, F_{s}\right)$ such that $\psi_{1}, \psi_{2} \in H_{r, s}$. Note that $\left(\psi_{1}, \psi_{2}\right)$ is, by construction, an ideal basis of its linear span in $H_{r, s}$. Therefore, Proposition 1.1.1 shows that

$$
a_{1} \geq|r-s|, \quad a_{2} \geq r+s+\epsilon
$$

where $\epsilon=0$ if $K / \mathbb{Q}_{p}$ is unramified and $\epsilon=1$ otherwise. We claim that

$$
\begin{equation*}
a_{3}=l_{r, s}\left(\psi_{3}\right):=\max \left\{v\left(\psi_{3}+\varphi\right) \mid \varphi \in H_{r, s}\right\} \tag{2.1}
\end{equation*}
$$

(this notation was already used in the previous section). Indeed, since $\psi_{1}, \psi_{2}, \psi_{3}$ is an ideal basis of $L$ we have

$$
\begin{equation*}
a_{3}=v\left(\psi_{3}\right)=\max \left\{v\left(x_{1} \psi_{1}+x_{2} \psi_{2}+\psi_{3}\right) \mid x_{1}, x_{2} \in \mathbb{Z}_{p}\right\} \tag{2.2}
\end{equation*}
$$

Therefore, the inequality ' $\leq$ ' in (2.1) follows from the inclusion $\left\langle\psi_{1}, \psi_{2}\right\rangle \subset H_{r, s}$. On the other hand, $[\mathbf{B}$, Corollary 6.7$]$, shows that (2.2) still holds if we allow $x_{1}, x_{2} \in \mathbb{Q}_{p}$. Hence the inequality ' $\geq$ ' follows from the inclusion $H_{r, s} \subset\left\langle\psi_{1}, \psi_{2}\right\rangle \otimes \mathbb{Q}_{p}$, proving the claim. We conclude that $l_{r, s}\left(\psi_{3}\right)=a_{3} \geq a_{2} \geq r+s+\epsilon$. In fact, we even have

$$
\begin{equation*}
l_{r, s}\left(\psi_{3}\right) \geq r+s+1 \tag{2.3}
\end{equation*}
$$

For if $K / \mathbb{Q}_{p}$ is unramified, then $a_{1}$ and $a_{2}$ are even and so $a_{3}$ must be odd.
By [ $\mathbf{B}]$, Corollary $6.7,\left(\psi_{1}, \psi_{2}, \psi_{3}^{\prime}\right)$ is again an ideal basis of its linear span (in some order). Therefore, we can apply the same argument to $\psi_{3}^{\prime}$. We get

$$
\begin{equation*}
l_{r, s}\left(\psi_{3}^{\prime}\right)=l_{r, s}\left(\psi_{3}\right)-2 \geq r+s-1 \tag{2.4}
\end{equation*}
$$

For $\alpha \in \mathcal{O}_{D}^{\times}$, let $F_{r}^{\alpha}$ denote the deformation of $G$ obtained by composing the identification $F_{r} \otimes k \xrightarrow{\sim} G$ with $\alpha$. Define $\mathcal{C}_{r, s}=\mathcal{C}\left(F_{r}, F_{s}\right) \subset \mathcal{S}$ as the closed subscheme where $\left.\Gamma\right|_{\mathcal{C}_{r, s}} \cong F_{r}^{\alpha}$ and $\left.\Gamma^{\prime}\right|_{\mathcal{C}_{r, s}} \cong F_{s}^{\alpha}$, for some $\alpha \in \mathcal{O}_{D}^{\times}$. It follows from the results of $[\mathbf{W w} \mathbf{1}]$ that $\mathcal{C}_{r, s} \cong \operatorname{Spf} \mathcal{O}_{M_{t}}$, where $t=\max \{s, r\}$. Moreover,

$$
\mathcal{T}_{1} \cdot \mathcal{T}_{2}=\bigcup_{\left(F_{r}, F_{s}\right)} \mathcal{C}_{r, s}
$$

is the decomposition into irreducible components. To prove the proposition it therefore suffices to show that

$$
\begin{equation*}
\left(\mathcal{C}_{r, s} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}=\left(\mathcal{C}_{r, s} \cdot \mathcal{T}_{3}^{\prime}\right)_{\mathcal{S}}+\left(\mathcal{C}_{r, s} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}} \tag{2.5}
\end{equation*}
$$

for all pairs $\left(F_{r}, F_{s}\right)$. We also may assume that $r \leq s$. Then $\left(\mathcal{C}_{r, s} \cdot \mathcal{S}_{(p)}\right)_{\mathcal{S}}=e_{s}$ is the ramification index of $\mathcal{O}_{M_{s}}$ over $W$. Moreover, in the notation of the last subsection, we have

$$
\begin{equation*}
\left(\mathcal{C}_{r, s} \cdot \mathcal{T}_{3}\right)_{\mathcal{S}}=n_{r, s}\left(\psi_{3}\right), \quad\left(\mathcal{C}_{r, s} \cdot \mathcal{T}_{3}^{\prime}\right)_{\mathcal{S}}=n_{r, s}\left(\psi_{3}^{\prime}\right) \tag{2.6}
\end{equation*}
$$

However, by (2.3), (2.4) and Proposition 1.2 we have $n_{r, s}\left(\psi_{3}\right)=n_{r, s}\left(\psi_{3}^{\prime}\right)+e_{s}$. This proves (2.5) and finishes the proof of the proposition.

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# 15. CALCULATION OF REPRESENTATION DENSITIES 

by

Torsten Wedhorn


#### Abstract

We calculate for all primes $p \geq 2$ the local representation density of a ternary quadratic form $Q$ over $\mathbb{Z}_{p}$ in a quadratic space of the form $N \perp H^{r}$, where $N$ is a quadratic space of rank $4, H$ is the hyperbolic plane, and $r \geq 0$ is any non-negative integer. Our principal tool is a formula of Katsurada. This defines a rational function $f_{Q, N}$ in $p^{-r}$. We also determine the derivative of $f_{Q, N}$ and relate it to the arithmetic intersection number of three modular correspondences.

Résumé (Calcul de densités de représentation). - On calcule, pour tous les nombres premiers $p \geq 2$, la densité de représentation locale d'une forme quadratique ternaire $Q$ sur $\mathbb{Z}_{p}$ dans un espace quadratique de la forme $N \perp H^{r}$, où $N$ est un espace quadratique de rang $4, H$ est le plan hyperbolique, et $r$ est un entier $\geq 0$. Notre outil principal est une formule de Katsurada. Elle est donnée par une fonction rationnelle $f_{Q, N}$ en $p^{-r}$. Nous déterminons également la dérivée de $f_{Q, N}$ et nous la relions au nombre d'intersection arithmétique de trois correspondances modulaires.


## Introduction

In this note we consider local representation densities of ternary quadratic spaces and derivatives of associated rational functions. These results are used in $[\mathbf{R W}]$ to relate the arithmetic intersection number of three modular correspondences $\left(\mathcal{T}_{m_{1}}\right.$. $\mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}$ ) to a Fourier coefficient of the restriction of the derivative at $s=0$ of a Siegel-Eisenstein series of genus 3 and weight 2. We also obtain an explicit formula for the integers $\beta_{l}(Q)$ which occur in [Go2].

Let $Q$ and $N$ be quadratic spaces over $\mathbb{Z}_{p}$ of rank 3 and 4 respectively, and let $H$ be the hyperbolic plane over $\mathbb{Z}_{p}$. Denote by $\alpha_{p}\left(Q, N \perp H^{r}\right)$ the local representation density, compare [Wd1, 4.3]. This is a rational function $f_{Q, N}(X)$ in $X=p^{-r}$.

2000 Mathematics Subject Classification. - 11E08, 11E12, 14G35.
Key words and phrases. - Representation density, modular correspondence.

In the first section we consider the case that $N$ is anisotropic and that $r=0$ :
(1) Let $D$ be "the" quaternion division algebra over $\mathbb{Q}_{p}$ and $N=O_{D}$ be its maximal order endowed with the reduced norm. Then we compute $\alpha_{p}(Q, N)$ for any ternary form $Q$ by a direct calculation (Theorem 1.1), following closely [GK, section 6].
The value obtained is of course 0 if $Q$ is isotropic, and for anisotropic $Q$ we will see that it does not depend on $Q$.

In general it is very difficult to compute local representation densities $\alpha_{p}(Q, N)$, and their computation has a long history. We give only a few references: For $p \neq 2$ a general explicit formula has been given by Hironaka and Sato [HS] for arbitrary quadratic spaces $Q$ and $N$ over $\mathbb{Z}_{p}$. If the rank of $Q$ is 2 , Yang has given a formula for $\alpha_{p}(Q, N)$ in the case of $p=2[\mathbf{Y 1}]$. We will use a result of Katsurada [Ka] who calculated $\alpha_{p}(Q, N)$ for arbitrary $p$ and $Q$ in the case that $N$ is an orthogonal sum of copies of the hyperbolic plane $H$.

In the second section we are interested in the following values:
(2) Let $N=H^{2}$. Then we specialize Katsurada's formula for $\alpha_{p}\left(Q, N \perp H^{r}\right)=$ $\alpha_{p}\left(Q, H^{r+2}\right)$ to the case where $Q$ is a ternary form and express it in terms of a refinement of the Gross-Keating invariants (see $[\mathbf{B}]$ ) of the ternary form $Q$. This is done in 2.11 .
(3) For $Q$ (ternary and) isotropic we specialize this formula to $r=0$ and therefore obtain $\alpha_{p}\left(Q, H^{2}\right)$ (Proposition 2.1) (for $Q$ anisotropic, $\alpha_{p}\left(Q, H^{2}\right)=0$ ).
(4) Finally we calculate for $N=H^{2}$ and for $Q$ a ternary anisotropic quadratic form the derivative $\frac{\partial}{\partial X} f_{Q, H^{2}}(X)$ at $X=1$ (see 2.16).
We remark that the values obtained in (3) and (4) depend only on the GrossKeating invariants of the ternary form $Q$ although the value in (2) depends on a refinement of these invariants.

Acknowledgements. - I am grateful to T. Yang for spotting a mistake in an earlier version of the manuscript.

## 1. Calculation of the representation density in the anisotropic case

1.1. We fix a prime number $p$, let $D$ be "the" quaternion division algebra over $\mathbb{Q}_{p}$, and denote by $N=O_{D}$ the maximal order of $D$ which we consider as a quadratic space of rank 4 over $\mathbb{Z}_{p}$ with respect to the reduced norm. Let $Q$ be any ternary quadratic form over $\mathbb{Z}_{p}$. In this section we are going to calculate the representation density $\alpha_{p}(Q, N)$.

As $N$ is an anisotropic quadratic space, $Q$ is represented by $N$ if and only if $Q$ is anisotropic. In this case the result is:

Theorem 1.1. - Let $Q$ be anisotropic. Then

$$
\alpha_{p}(Q, N)=2(p+1)\left(1+\frac{1}{p}\right)
$$

1.2. For the proof we quote the following lemma from $[\mathbf{K i}]$ Theorem 5.6.4(e):

Lemma 1.2. - For any integer $r \in \mathbb{Z}$ we have

$$
\alpha_{p}\left(p^{r} Q, N\right)=\alpha_{p}(Q, N)
$$

1.3. Proof of Theorem 1.1.- By Lemma 1.2 we can assume that the underlying $\mathbb{Z}_{p}$-module of the quadratic space $Q$ is a sublattice $\Lambda$ in $O_{D}$ such that $\Lambda \not \subset p O_{D}$.

Clearly any element of $\mathrm{O}(D, \mathrm{Nrd})$ preserves $N$ and hence $\mathrm{O}(D, \mathrm{Nrd})$ acts on

$$
\tilde{A}_{p^{r}}(Q, N):=\left\{\sigma: Q / p^{r} Q \longrightarrow N / p^{r} N \mid \operatorname{Nrd}(\sigma(x)) \equiv Q(x) \bmod p^{r}\right\}
$$

for all $r \geq 1$. By definition (see [ $\mathbf{W d 1}, 4.3]$ ) we have

$$
\alpha_{p}(Q, N)=\left(p^{r}\right)^{-6} \# \tilde{A}_{p^{r}}(Q, N)
$$

for $r$ sufficiently large.
The dual lattice of $N=O_{D}$ with respect to the pairing associated to the quadratic form is $N^{\vee}=\mathfrak{m}^{-1} \subset D$ where $\mathfrak{m}$ is the maximal ideal of $O_{D}$. We claim that the induced action of $\mathrm{SO}(D, \mathrm{Nrd})$ on

$$
\tilde{B}_{p^{r}}(Q, N):=\left\{\sigma: \mathbb{Z}_{p}^{3} \longrightarrow N / p^{r} \mathfrak{m}^{-1} N \mid \operatorname{Nrd}(\sigma(x)) \equiv Q(x) \bmod p^{r}\right\}
$$

is transitive for $r \geq 1$. For this it suffices to show that $\mathrm{SO}(D, \mathrm{Nrd})$ acts transitively on the set $M$ of all isometries $\tilde{\sigma}: Q \rightarrow N$. But by Witt's lemma, $\mathrm{O}(D, \mathrm{Nrd})$ acts transitively on $M$. For every such $\tilde{\sigma}$ the stabilizer in $\mathrm{O}(D, \mathrm{Nrd})$ is nothing but the orthogonal group of the orthogonal complement of the quadratic $\mathbb{Q}_{p}$-space generated by $\tilde{\sigma}(Q)$. As this complement is a one-dimensional space, we see that $\mathrm{SO}(D, \mathrm{Nrd})$ acts in fact simply transitively on $M$.

Using [Wd1, Lemma 1.6] we identify $\mathrm{SO}(D, \mathrm{Nrd})$ with

$$
\left\{\left(d, d^{\prime}\right) \in D^{\times} \times D^{\times} \mid \operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)\right\} / \mathbb{Q}_{p}^{\times}
$$

This group contains the subgroup of index 2

$$
G=\left\{\left(d, d^{\prime}\right) \in O_{D}^{\times} \times O_{D}^{\times} \mid \operatorname{Nrd}(d)=\operatorname{Nrd}\left(d^{\prime}\right)\right\} / \mathbb{Z}_{p}^{\times}
$$

Therefore $G$ acts with two orbits on $\tilde{B}_{p^{r}}(Q, N)$. Let $\bar{G}$ be the quotient of $G$ by the subgroup generated by

$$
\left\{\left(d, d^{\prime}\right) \in G \mid d \equiv d^{\prime} \equiv 1 \quad\left(\bmod p^{r} N^{\vee}\right)\right\}
$$

and by $1+p^{r-1} O_{D_{p}}$ diagonally embedded in $G$. Then $\bar{G}$ acts faithfully with 2 orbits on $\tilde{B}_{p^{r}}(Q, N)$. As

$$
\# \tilde{A}_{p^{r}}(Q, N)=\left(\# \tilde{B}_{p^{r}}(Q, N)\right) \cdot\left(\#\left(\mathfrak{m}^{-1} / O_{D}\right)\right)^{3}
$$

we see that

$$
\# \tilde{A}_{p^{r}}(Q, N)=2(\# \bar{G})\left(\#\left(\mathfrak{m}^{-1} / O_{D}\right)\right)^{3}=2(p+1)^{2} p^{6 r-7} p^{6}
$$

It follows that

$$
\alpha_{p}(Q, N)=p^{-6 r} 2(p+1)^{2} p^{6 r-1}=2(p+1)\left(1+\frac{1}{p}\right)
$$

## 2. Calculation of the representation density in the hyperbolic case

2.1. Again we fix a prime number $p$. For any element $a \in \mathbb{Q}_{p}^{\times}$we write $\operatorname{ord}(a) \in \mathbb{Z}$ for the $p$-adic valuation of $a$.

We denote by $H$ the quadratic space over $\mathbb{Z}_{p}$ whose underlying module is $\mathbb{Z}_{p}^{2}$ and whose matrix with respect to the standard basis is $\left(\begin{array}{cc}0 & \frac{1}{2} \\ \frac{1}{2} & 0\end{array}\right)$. This means that the quadratic form is given by $\mathbb{Z}_{p}^{2} \ni(x, y) \mapsto x y$.

Note that $H^{2} \cong\left(M_{2}\left(\mathbb{Z}_{p}\right)\right.$, det $)$.
Let $(M, Q)$ be any quadratic space over $\mathbb{Z}_{p}$ of rank 3 . In this section we will compute the representation density $\alpha_{p}\left(M, H^{r+2}\right)$. In fact, there is a polynomial $f_{M}(X) \in \mathbb{Q}[X]$ such that $f_{M}\left(p^{-r}\right)=\alpha_{p}\left(M, H^{2+r}\right)([\mathbf{K a}])$. We are interested in

$$
\begin{equation*}
f_{M}(1)=\alpha_{p}\left(M, H^{2}\right) \tag{2.1}
\end{equation*}
$$

and, for $(M, Q)$ anisotropic, in

$$
\begin{equation*}
\frac{\partial}{\partial X} f_{M}(X)_{\mid X=1} \tag{2.2}
\end{equation*}
$$

The first value is given in 2.12 and the second in 2.16 .
2.2. We use the formulas by Katsurada [Ka] but we express them in terms of the Gross-Keating invariants (cf. [B]) of the ternary space ( $M, Q$ ), an invariant $\tilde{\xi}=\tilde{\xi}(M) \in\{-1,0,1\}$, and an invariant $\eta=\eta(M) \in\{ \pm 1\}$.

The invariant $\eta$ is equal to +1 if $(M, Q)$ is isotropic and equal to -1 if $(M, Q)$ is anisotropic.

The Gross-Keating invariants consist of a tuple of integers $G K(M)=\left(a_{1}, a_{2}, a_{3}\right)$ such that $0 \leq a_{1} \leq a_{2} \leq a_{3}$. In addition, if $a_{1} \equiv a_{2} \bmod 2$ and $a_{2}<a_{3}$ there is a further invariant $\epsilon_{G K}(M) \in\{ \pm 1\}$.

In fact, we will not need the invariant $\epsilon_{G K}(M)$ directly in the sequel, as $\tilde{\xi}(T)$ is a refinement. But we remark that the final expressions for (2.1) and (2.2) depend only on $\eta(M)$ (that is, whether $(M, Q)$ is isotropic or not) and on the Gross-Keating invariants $G K(M)$ and $\epsilon_{G K}(M)$.

If $T$ is the matrix associated to $(M, Q)$ and a $\mathbb{Z}_{p}$-base of $M$, we also write $\eta(T)$, $G K(T), \epsilon_{G K}(T)$, and $\tilde{\xi}(T)$.
2.3. Recall the Hilbert symbol $(a, b)_{p} \in\{ \pm 1\}$ for $a, b \in \mathbb{Q}_{p}^{\times}$. It is uniquely determined by the following properties (where $a, b, b^{\prime} \in \mathbb{Q}_{p}^{\times}, u, v \in \mathbb{Z}_{p}^{\times}$):

$$
\begin{aligned}
(a, b)_{p} & =(b, a)_{p} \\
\left(a, b b^{\prime}\right)_{p} & =(a, b)_{p}\left(a, b^{\prime}\right)_{p} \\
(p, p)_{p} & =(-1, p)_{p}
\end{aligned}
$$

and, for $p$ odd, by

$$
\begin{aligned}
& (u, p)_{p}=\left(\frac{u}{p}\right), \\
& (u, v)_{p}=1
\end{aligned}
$$

and, for $p=2$, by

$$
\begin{aligned}
& (u, 2)_{2}= \begin{cases}+1, & \text { if } u \equiv \pm 1 \bmod 8 \\
-1, & \text { otherwise }\end{cases} \\
& (u, v)_{2}= \begin{cases}+1, & \text { if } u \text { or } v \equiv 1 \bmod 4, \\
-1, & \text { otherwise }\end{cases}
\end{aligned}
$$

2.4. For any symmetric matrix $T \in \operatorname{Sym}_{m}\left(\mathbb{Q}_{p}\right)$ we denote by $h(T)=h_{p}(T)$ the Hasse invariant of the associated quadratic space $(M, Q)$. We use the normalization in $[\mathbf{K i}]$. For $m=3$ we have

$$
h(T)= \begin{cases}(-1)^{\delta_{2 p}}, & \text { if }(M, Q) \text { is isotropic } \\ -(-1)^{\delta_{2 p}}, & \text { if }(M, Q) \text { is anisotropic }\end{cases}
$$

by [Ki, 3.5.1]. Here $\delta_{2 p}$ is the Kronecker delta.
2.5. In the next sections we recall some results from $[\mathbf{B}]$ (cf. also $[\mathbf{Y} 1]$ ). We start with the case $p>2$. In that case there exists a basis $\left(e_{i}\right)$ of $M$ such that the matrix $T=\left(t_{i j}\right)$ associated to $Q$ with respect to this basis (i.e., $\left.t_{i j}=\frac{1}{2}\left(Q\left(e_{i}+e_{j}\right)-Q\left(e_{i}\right)-Q\left(e_{j}\right)\right)\right)$ is a diagonal matrix. If we write $t_{i i}=u_{i} p^{a_{i}}$ for $a_{i} \in \mathbb{Z}$ and $u_{i} \in \mathbb{Z}_{p}^{\times}$, we can assume that $a_{1} \leq a_{2} \leq a_{3}$. Moreover, if $a_{i}=a_{i+1}$ we can assume that $u_{i+1}=1$. Then the Gross-Keating invariants are given as follows. We have

$$
G K(T)=\left(a_{1}, a_{2}, a_{3}\right)
$$

If $a_{1} \equiv a_{2} \bmod 2$ and $a_{2}<a_{3}$, we have

$$
\epsilon_{G K}(T)=\left(\frac{-u_{1} u_{2}}{p}\right)
$$

We set

$$
\tilde{\xi}(T)= \begin{cases}\left(\frac{-u_{1} u_{2}}{p}\right), & \text { if } a_{1} \equiv a_{2} \bmod 2 \\ 0, & \text { if } a_{1} \not \equiv a_{2} \bmod 2\end{cases}
$$

Finally, let $i, j \in\{1,2,3\}$ with $i \neq j$ and $a_{i} \equiv a_{j} \bmod 2$ and define $k \in\{1,2,3\}$ by $\{1,2,3\} \backslash\{i, j\}=\{k\}$. Then $T$ is isotropic if and only if $\left(-u_{i} u_{j}, p\right)_{p}=1$ or $a_{k} \equiv a_{j}$ $\bmod 2$.
2.6. Now assume that $p=2$. In the sequel $K$ will denote one of the matrices

$$
H=\left(\begin{array}{cc}
0 & \frac{1}{2} \\
\frac{1}{2} & 0
\end{array}\right) \quad \text { or } \quad Y:=\left(\begin{array}{cc}
1 & \frac{1}{2} \\
\frac{1}{2} & 1
\end{array}\right)
$$

There exists a basis $\mathcal{B}$ of $M$ such that the matrix $T$ associated to $Q$ with respect to $\mathcal{B}$ is of one of the following forms.

Either $Q$ is not diagonalizable (case A). Then we distinguish two subcases:
(A1) $T=\operatorname{diag}\left(u 2^{\alpha}, 2^{\beta} K\right)$ where $\alpha \leq \beta$ are integers and $u \in \mathbb{Z}_{2}^{\times}$. Then

$$
G K(T)=(\alpha, \beta, \beta)
$$

We set

$$
\tilde{\xi}(T)= \begin{cases}1, & \text { if } a_{1} \equiv a_{2} \bmod 2 \\ 0, & \text { if } a_{1} \not \equiv a_{2} \bmod 2\end{cases}
$$

(A2) $T=\operatorname{diag}\left(2^{\alpha} K, u 2^{\beta}\right)$ where $\alpha<\beta$ are integers and $u \in \mathbb{Z}_{2}^{\times}$. Then

$$
G K(T)=(\alpha, \alpha, \beta)
$$

In this case $\epsilon_{G K}(T)$ is defined and we have

$$
\epsilon_{G K}(T)= \begin{cases}+1 & \text { if } K=H \\ -1 & \text { if } K=Y\end{cases}
$$

We set $\tilde{\xi}(T):=\epsilon_{G K}(T)$.
In the nondiagonalizable case $\mathrm{A}, T$ is isotropic if and only if $K=H$ or $\alpha \equiv \beta$ $\bmod 2$.

Now assume that $T$ is diagonalizable over $\mathbb{Z}_{2}$ (case B), i.e., there exists a basis such that $T=\operatorname{diag}\left(u_{1} 2^{\beta_{1}}, u_{2} 2^{\beta_{2}}, u_{3} 2^{\beta_{3}}\right)$ where $0 \leq \beta_{1} \leq \beta_{2} \leq \beta_{3}$ are integers and $u_{i} \in \mathbb{Z}_{2}^{\times}$. Then there are four subcases (here our subdivision of cases is different from [R]):
(B1) $\beta_{1} \not \equiv \beta_{2} \bmod 2$. Then

$$
G K(T)=\left(\beta_{1}, \beta_{2}, \beta_{3}+2\right)
$$

We set $\tilde{\xi}(T):=0$.
(B2) $\beta_{1} \equiv \beta_{2} \bmod 2$ and $\left(u_{1} u_{2} \equiv 1 \bmod 4\right.$ or $\left.\left.\beta_{3}=\beta_{2}\right)\right)$. Then

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right) .
$$

We set $\tilde{\xi}(T):=0$.
(B3) $\beta_{1} \equiv \beta_{2} \bmod 2, \beta_{3}=\beta_{2}+1$, and $u_{1} u_{2} \equiv-1 \bmod 4$. Then

$$
G K(T)=\left(\beta_{1}, \beta_{2}+1, \beta_{3}+1\right)
$$

We set $\tilde{\xi}(T):=\left(-u_{1} u_{2}, 2\right)_{2}$ where $(,)_{2}$ denotes the Hilbert symbol.
(B4) $\beta_{1} \equiv \beta_{2} \bmod 2, \beta_{3}>\beta_{2}+1$, and $u_{1} u_{2} \equiv-1 \bmod 4$. Then

$$
G K(T)=\left(\beta_{1}, \beta_{2}+2, \beta_{3}\right) .
$$

In this case $\epsilon_{G K}(T)$ is defined and we have

$$
\epsilon_{G K}(T)=\left(-u_{1} u_{2}, 2\right)_{2} .
$$

We set $\tilde{\xi}(T):=\epsilon_{G K}(T)$.
Finally, let $i, j \in\{1,2,3\}$ with $i \neq j$ and $\beta_{i} \equiv \beta_{j} \bmod 2$ and define $k \in\{1,2,3\}$ by $\{1,2,3\} \backslash\{i, j\}=\{k\}$. Then $T$ is isotropic if and only if

$$
\left(-u_{k} u_{j},-u_{i} u_{j}\right)_{2}=\left(-u_{i} u_{j}, 2\right)_{2}^{\beta_{k}+\beta_{j}}
$$

2.7. Going through the cases in 2.5 and 2.6 we see that there are the following possibilities for the value of $\tilde{\xi}$ if $T$ is anisotropic:

- If $a_{1} \not \equiv a_{2} \bmod 2$, we either have $\tilde{\xi}=0$ or we have $\tilde{\xi}=-1$ and $a_{3}=a_{2}+1$.
- If $a_{1} \equiv a_{2} \bmod 2$, we always have $a_{2} \not \equiv a_{3} \bmod 2$ and $\tilde{\xi}=-1$.

If $T$ is isotropic, the possibilities for the value of $\tilde{\xi}$ are the following:

- If $a_{1} \not \equiv a_{2} \bmod 2$, we either have $\tilde{\xi}=0$ or we have $\tilde{\xi}=1$ and $a_{3}=a_{2}+1$.
- If $a_{1} \equiv a_{2} \bmod 2$, we either have $\tilde{\xi}=1$ or we have $\tilde{\xi}=-1$ and $a_{2} \equiv a_{3} \bmod 2$.
2.8. By $[\mathbf{K a}]$ there exists a polynomial $f_{M}(X)=f_{T}(X) \in \mathbb{Q}[X]$ such that $f_{T}\left(p^{-r}\right)=$ $\alpha_{p}\left(M, H^{2+r}\right)$. We use the formulas from $[\mathbf{K a}]$ to compute $f_{T}$. Indeed, by loc. cit. p. 417 and p. 428 we have

$$
f_{T}(X)=\tilde{\gamma}_{p}(T ; X) \tilde{F}_{p}(T ; X)
$$

with $\tilde{\gamma}_{p}(T ; X)=\gamma_{p}\left(T ; p^{-2} X\right)$ and $\tilde{F}_{p}(T ; X)=F_{p}\left(T ; p^{-2} X\right)$ where $\gamma_{p}(T ; X)$ and $F_{p}(T ; X)$ are the rational functions defined in loc. cit. p. 417 and p. 451 respectively. Thus

$$
\tilde{\gamma}_{p}(T ; X)=\left(1-p^{-2} X\right)\left(1-p^{-2} X^{2}\right)
$$

The function $\tilde{F}_{p}(T ; X)$ is more complicated. We will express it in the next sections using the Gross-Keating invariants $G K(T)$ and the invariant $\tilde{\xi}(T)$.
2.9. By [Ka] we have

$$
\begin{align*}
\tilde{F}_{p}(T ; X)= & \sum_{i=0}^{\hat{\delta}} \sum_{j=0}^{\tilde{\delta}^{\prime} / 2-i-1} p^{i+j} X^{i+2 j} \\
& +\eta p^{\left(\tilde{\delta}^{\prime}-2\right) / 2} X^{\delta-\tilde{\delta}^{\prime}+2} \sum_{i=0}^{\hat{\delta}} \sum_{j=0}^{\tilde{\delta}^{\prime} / 2-i-1} p^{-j} X^{i+2 j}  \tag{2.3}\\
& +\tilde{\xi}^{2} p^{\tilde{\delta}^{\prime} / 2} X^{\tilde{\delta}^{\prime}-\hat{\delta}} \sum_{i=0}^{\hat{\delta}} \sum_{j=0}^{\delta-2 \tilde{\delta}^{\prime}+\hat{\delta}} \tilde{\xi}^{j} X^{i+j}
\end{align*}
$$

where $\eta, \delta, \hat{\delta}$, and $\tilde{\delta}^{\prime}$ are the invariants defined on p. 450 of loc. cit. (note that in loc. cit. the definitions of $\tilde{\delta}$ and $\hat{\delta}$ have to be interchanged).
2.10. Going through all the cases in 2.5 and 2.6 one sees that $\eta, \delta, \hat{\delta}$, and $\tilde{\delta}^{\prime}$ can be expressed as follows (where $G K(T)=\left(a_{1}, a_{2}, a_{3}\right)$ are the Gross-Keating invariants):

$$
\begin{align*}
& \eta= \begin{cases}+1 & \text { if } T \text { is isotropic }, \\
-1 & \text { if } T \text { is anisotropic },\end{cases}  \tag{2.4}\\
& \delta=a_{1}+a_{2}+a_{3},  \tag{2.5}\\
& \hat{\delta}=a_{1},  \tag{2.6}\\
& \tilde{\delta}^{\prime}= \begin{cases}a_{1}+a_{2}, & \text { if } a_{1} \equiv a_{2} \bmod 2, \\
a_{1}+a_{2}+1, & \text { if } a_{1} \not \equiv a_{2} \bmod 2,\end{cases} \tag{2.7}
\end{align*}
$$

2.11. If we set

$$
\sigma:= \begin{cases}2, & \text { if } a_{1} \equiv a_{2} \bmod 2, \\ 1, & \text { if } a_{1} \not \equiv a_{2} \bmod 2,\end{cases}
$$

we can rewrite (2.3) using the invariants $\eta,\left(a_{1}, a_{2}, a_{3}\right)$, and $\tilde{\xi}$ :

$$
\begin{aligned}
\tilde{F}_{p}(T, X)= & \sum_{i=0}^{a_{1}} \sum_{j=0}^{\left(a_{1}+a_{2}-\sigma\right) / 2-i} p^{i+j} X^{i+2 j} \\
& +\eta \sum_{i=0}^{a_{1}} \sum_{j=0}^{\left(a_{1}+a_{2}-\sigma\right) / 2-i} p^{\left(a_{1}+a_{2}-\sigma\right) / 2-j} X^{a_{3}+\sigma+i+2 j} \\
& +\tilde{\xi}^{2} p^{\left(a_{1}+a_{2}-\sigma+2\right) / 2} \sum_{i=0}^{a_{1}} \sum_{j=0}^{a_{3}-a_{2}+2 \sigma-4} \tilde{\xi}^{j} X^{a_{2}-\sigma+2+i+j} .
\end{aligned}
$$

2.12. We now specialize to $r=0$, i.e., $X=1$. In that case we have

$$
\alpha_{p}\left(T, H^{2}\right)=f_{T}(1)=\left(1-p^{-2}\right)^{2} \tilde{F}_{p}(T, 1) .
$$

If we set $\beta_{p}(T):=\tilde{F}_{p}(T, 1)$, it follows from (2.8) that

$$
\begin{align*}
\beta_{p}(T)= & (1+\eta)\left(\sum_{i=0}^{a_{1}-1}(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right) p^{i}\right)  \tag{2.9}\\
& +p^{\left(a_{1}+a_{2}-\sigma+2\right) / 2}\left(a_{1}+1\right) R_{\tilde{\xi}}
\end{align*}
$$

where

$$
R_{\tilde{\xi}}= \begin{cases}0, & \text { if } \tilde{\xi}=0 \\ 0, & \text { if } \tilde{\xi}=-1 \text { and } a_{3} \not \equiv a_{2} \bmod 2 ; \\ a_{3}-a_{2}+2 \sigma-3, & \text { if } \tilde{\xi}=1 ; \\ 1, & \text { if } \tilde{\xi}=-1 \text { and } a_{3} \equiv a_{2} \bmod 2 .\end{cases}
$$

2.13. If $T$ is anisotropic we have $\alpha_{p}\left(T, H^{2}\right)=\beta_{p}(T)=0$, as a three dimensional anisotropic space cannot be represented by a four-dimensional hyperbolic space. Alternatively this follows also from (2.9): By (2.4) we have $\eta=-1$ and hence it suffices to show that $R_{\tilde{\xi}}=0$ if $T$ is anisotropic. By 2.7 we are in one of the following two cases:
(a) $\tilde{\xi}=0$;
(b) $\tilde{\xi}=-1$ and $a_{2} \not \equiv a_{3} \bmod 2$.

In both cases we have $R_{\tilde{\xi}}=0$ by definition.
2.14. If $T$ is isotropic, (2.9) gives Proposition 6.25 of [GK]:

Proposition 2.1. - Let $T$ be isotropic. Then:
(1) If $a_{1} \not \equiv a_{2} \bmod 2$, we have

$$
\beta_{p}(T)=2\left(\sum_{i=0}^{a_{1}-1}(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right) p^{i}\right) .
$$

(2) If $a_{1} \equiv a_{2} \bmod 2$ and $\tilde{\xi}=1$, we have

$$
\begin{aligned}
\beta_{p}(T)= & 2\left(\sum_{i=0}^{a_{1}-1}(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right) p^{i}\right) \\
& +\left(a_{1}+1\right)\left(a_{3}-a_{2}+1\right) p^{\left(a_{1}+a_{2}\right) / 2}
\end{aligned}
$$

(3) If $a_{1} \equiv a_{2} \bmod 2$ and $\tilde{\xi}=-1$, we have

$$
\begin{aligned}
\beta_{p}(T)= & 2\left(\sum_{i=0}^{a_{1}-1}(i+1) p^{i}+\sum_{i=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right) p^{i}\right) \\
& +\left(a_{1}+1\right) p^{\left(a_{1}+a_{2}\right) / 2} .
\end{aligned}
$$

Proof. - We have $\eta=1$, and by 2.7 we are in one of the following cases:
(a) $a_{1} \not \equiv a_{2} \bmod 2$ and $\tilde{\xi}=0$;
(b) $a_{1} \not \equiv a_{2} \bmod 2, \tilde{\xi}=1$, and $a_{3}=a_{2}+1$;
(c) $a_{1} \equiv a_{2} \bmod 2$ and $\tilde{\xi}=1$;
(d) $a_{1} \equiv a_{2} \bmod 2, \tilde{\xi}=-1$, and $a_{2} \equiv a_{3} \bmod 2$.

In case (a), we have $R_{\tilde{\xi}}=0$ by definition, and in case (b) we also have $R_{\tilde{\xi}}=$ $a_{3}-a_{2}+2 \sigma-3=0$. This proves (1).

In case (c), we have $R_{\tilde{\xi}}=a_{3}-a_{2}+1$ and therefore (2).
In case (d), we have $R_{\tilde{\xi}}=1$ which implies (3)
Corollary 2.2. - Set $\Delta(T)=\frac{1}{2} \operatorname{det}(2 T)=4 \operatorname{det}(T)$ and assume that $T$ is isotropic. Then $\beta_{p}(T)=1$ if $\operatorname{ord}_{p}(\Delta(T))=0$.

Proof. - For $p>2$ the equality $\operatorname{ord}_{p}(\Delta(T))=0$ is equivalent to $a_{1}=a_{2}=a_{3}=0$ by definition of the Gross Keating invariants (see 2.5). For $p=2$ the condition $\operatorname{ord}_{p}(\Delta)=0$ implies that we are in case (A1) of 2.6 with $\alpha=\beta=0$ and $K=H$. Therefore we have again $a_{1}=a_{2}=a_{3}=0$. Hence the corollary follows for all $p$ from Proposition 2.1.
2.15. From now on we assume that $T$ is anisotropic. We are going to calculate

$$
f_{T}^{\prime}(1)=\frac{\partial}{\partial X} f_{T}(X)_{\mid X=1} .
$$

As $T$ is anisotropic we have $\tilde{F}_{p}(T ; 1)=0$ and therefore

$$
\begin{align*}
f_{T}^{\prime}(1) & =\tilde{\gamma}_{p}(T, 1) \frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1}  \tag{2.10}\\
& =\left(1-p^{-2}\right)^{2} \frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1} . \tag{2.11}
\end{align*}
$$

Using (2.8) we see that

$$
\frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1}=F_{1}+F_{2}+F_{3} .
$$

Here

$$
\begin{aligned}
F_{1} & =\sum_{i=0}^{a_{1}} \sum_{j=0}^{\left(a_{1}+a_{2}-\sigma\right) / 2-i}(i+2 j) p^{i+j} \\
& =\sum_{l=0}^{a_{1}-1} \frac{3}{2}(l+1) l p^{l}+\sum_{l=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right)\left(2 l-\frac{a_{1}}{2}\right) p^{l},
\end{aligned}
$$

and

$$
\begin{aligned}
F_{2}= & -\sum_{i=0}^{a_{1}} \sum_{j=0}^{\left(a_{1}+a_{2}-\sigma\right) / 2-i}\left(a_{3}+\sigma+i+2 j\right) p^{\left(a_{1}+a_{2}-\sigma\right) / 2-j} \\
= & -\sum_{i=0}^{a_{1}} \sum_{j=i}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+a_{2}+a_{3}+i-2 j\right) p^{j} \\
= & -\sum_{l=0}^{a_{1}-1}(l+1)\left(a_{1}+a_{2}+a_{3}-\frac{3}{2} l\right) p^{l} \\
& -\sum_{l=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right)\left(\frac{3}{2} a_{1}+a_{2}+a_{3}-3 l\right) p^{l}
\end{aligned}
$$

and hence

$$
\begin{aligned}
F_{1}+F_{2}= & \sum_{l=0}^{a_{1}-1}(l+1)\left(3 l-a_{1}-a_{2}-a_{3}\right) p^{l} \\
& +\sum_{l=a_{1}}^{\left(a_{1}+a_{2}-\sigma\right) / 2}\left(a_{1}+1\right)\left(4 l-2 a_{1}-a_{2}-a_{3}\right) p^{l},
\end{aligned}
$$

and

$$
F_{3}=p^{\left(a_{1}+a_{2}-\sigma+2\right) / 2} \frac{a_{1}+1}{2} A_{\tilde{\xi}}
$$

with

$$
A_{\tilde{\xi}}= \begin{cases}0, & \text { if } \tilde{\xi}=0 \\ \left(a_{3}-a_{2}+2 \sigma-3\right)\left(a_{1}+a_{2}+a_{3}\right), & \text { if } \tilde{\xi}=1 ; \\ a_{2}-a_{3}-2 \sigma+3, & \text { if } \tilde{\xi}=-1, a_{2} \not \equiv a_{3} \bmod 2 \\ 3 a_{3}-a_{2}+a_{1}+4 \sigma-8, & \text { if } \tilde{\xi}=-1, a_{2} \equiv a_{3} \bmod 2\end{cases}
$$

2.16. We distinguish two cases. The first case is $a_{1} \not \equiv a_{2} \bmod 2$, i.e., $\sigma=1$. By 2.7 we either have $\tilde{\xi}=0$ and hence $A_{\tilde{\xi}}=0$ or we have $\tilde{\xi}=-1$ and $a_{3}=a_{2}+1$ and hence again $A_{\tilde{\xi}}=0$. Therefore we see that for $a_{1} \not \equiv a_{2} \bmod 2$ we have

$$
\begin{align*}
\frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1}= & \sum_{l=0}^{a_{1}-1}(l+1)\left(3 l-a_{1}-a_{2}-a_{3}\right) p^{l} \\
& +\sum_{l=a_{1}}^{\left(a_{1}+a_{2}-1\right) / 2}\left(a_{1}+1\right)\left(4 l-2 a_{1}-a_{2}-a_{3}\right) p^{l} . \tag{2.12}
\end{align*}
$$

The second case is $a_{1} \equiv a_{2} \bmod 2$, i.e., $\sigma=2$. Then we have $a_{3} \not \equiv a_{2} \bmod 2$ and hence

$$
\begin{align*}
\frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1}= & \sum_{l=0}^{a_{1}-1}(l+1)\left(3 l-a_{1}-a_{2}-a_{3}\right) p^{l} \\
& +\sum_{l=a_{1}}^{\left(a_{1}+a_{2}-2\right) / 2}\left(a_{1}+1\right)\left(4 l-2 a_{1}-a_{2}-a_{3}\right) p^{l}  \tag{2.13}\\
& +p^{\left(a_{1}+a_{2}\right) / 2} \frac{a_{1}+1}{2}\left(a_{2}-a_{3}-1\right) .
\end{align*}
$$

Therefore we see by $[\mathbf{R}$, Theorem 1.1] that in either case

$$
\frac{\partial}{\partial X} \tilde{F}_{p}(T ; X)_{\mid X=1}=-\lg \left(\mathcal{O}_{\mathcal{T}_{T}, \xi}\right)
$$

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# 16. THE CONNECTION TO EISENSTEIN SERIES 

by<br>Michael Rapoport \& Torsten Wedhorn


#### Abstract

We consider the non-singular Fourier coefficients of the special value of the derivative of a Siegel-Eisenstein series of genus 3 and weight 2 . We identify these coefficients with the arithmetic degrees of non-degenerate intersections of arithmetic modular correspondences.

Résumé (Relation avec les séries d'Eisenstein). - Nous identifions les coefficients de Fourier non-dégénerés d'une valeur spéciale de la dérivée d'une série de Siegel-Eisenstein de genre 3 et de poids 2 avec les degrés arithmétiques des intersections de correspondances modulaires arithmétiques.


## Introduction

In a previous chapter [Go2] an expression was obtained for the arithmetic intersection number of three modular correspondences $\left(\mathcal{T}_{m_{1}} \cdot \mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}\right)$, when their intersection is of dimension 0 . This expression is quite complicated, and involves certain local representation densities $\beta_{\ell}(Q)$ of quadratic forms and a local intersection multiplicity $\alpha_{p}(Q)$. It is this expression that is the main result of $[\mathbf{G K}]$. However, already in the introduction to their paper, Gross and Keating mention that computations of S. Kudla and D. Zagier strongly suggest that the arithmetic intersection number $\left(\mathcal{T}_{m_{1}} \cdot \mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}\right)$ agrees (up to a constant) with a Fourier coefficient of the restriction of the derivative at $s=0$ of a Siegel-Eisenstein series of genus 3 and weight 2 .

In the intervening years since the publication of [GK], Kudla has vastly advanced this idea and has in particular proved the analogue of this statement for the intersection of two Hecke correspondences on Shimura curves [Ku3]. In fact, Kudla has proposed a whole program which postulates a relation between special values of

[^10]derivatives of Siegel-Eisenstein series and arithmetic intersection numbers of special cycles on Shimura varieties for orthogonal groups, comp. [Ku4].

The purpose of the present chapter is to sketch these ideas of Kudla and to derive from Kudla's various papers on the subject the statement alluded to in the introduction of $[\mathbf{G K}]$. We stress that what we have done here is simply a task of compilation, since we do not (and cannot) claim to have mastered the automorphic side of the statement in question. We use the results of Katsurada [Ka] on local representation densities of quadratic forms, valid even for $p=2$, to relate the local intersection multiplicities to the derivatives of certain local Whittaker functions, comp. [Wd2]. For $p \neq 2$ the corresponding calculations of representation densities are much older and are based on results of Kitaoka [Kit].

We thank S. Kudla for his help with this chapter.

## 1. Decomposition of the intersections of modular correspondences

1.1. To $m \in \mathbb{Z}_{>0}$ we have associated the Deligne-Mumford stack which parametrizes the category of isogenies of degree $m$ between elliptic curves,

$$
\mathcal{T}_{m}(S)=\left\{f: E \longrightarrow E^{\prime} \mid \operatorname{deg}(f)=m\right\}
$$

Here $E$ and $E^{\prime}$ are elliptic curves over $S$. Then $\mathcal{T}_{m}$ maps by a finite unramified morphism to the stack $\mathcal{M}^{(2)}=\mathcal{M} \times_{\text {Spec } \mathbb{Z}} \mathcal{M}$ parametrizing pairs ( $E, E^{\prime}$ ) of elliptic curves.

Let now $m_{1}, m_{2}, m_{3} \in \mathbb{Z}_{>0}$ and consider

$$
\mathcal{T}\left(m_{1}, m_{2}, m_{3}\right)=\left\{\mathbf{f}=\left(f_{1}, f_{2}, f_{3}\right) \mid f_{i}: E \longrightarrow E^{\prime}, \operatorname{deg} f_{i}=m_{i}\right\},
$$

the fiber product of $\mathcal{T}_{m_{1}}, \mathcal{T}_{m_{2}}, \mathcal{T}_{m_{3}}$ over $\mathcal{M}^{(2)}$. Denoting by $Q$ the degree quadratic form on $\operatorname{Hom}\left(E, E^{\prime}\right)$, we obtain a disjoint sum decomposition,

$$
\begin{equation*}
\mathcal{T}\left(m_{1}, m_{2}, m_{3}\right)=\coprod_{T} \mathcal{I}_{T} . \tag{1.1}
\end{equation*}
$$

Here

$$
\mathcal{T}_{T}(S)=\left\{\mathbf{f} \in \operatorname{Hom}_{S}\left(E, E^{\prime}\right)^{3} \left\lvert\, \frac{1}{2}(\mathbf{f}, \mathbf{f})=T\right.\right\}
$$

where $(\mathbf{f}, \mathbf{f})$ denotes the matrix $\left(a_{i j}\right)$ with $a_{i j}=\left(f_{i}, f_{j}\right)=Q\left(f_{i}+f_{j}\right)-Q\left(f_{i}\right)-Q\left(f_{j}\right)$. Note that, due to the positive definiteness of $Q$, the index set in $(1.1)$ is $\operatorname{Sym}_{3}(\mathbb{Z})_{\geq 0}^{\vee}$, the set of positive semi-definite half-integral matrices.

Lemma 1.1. - Let $T \in \operatorname{Sym}_{3}(\mathbb{Z})_{>0}^{\vee}$, i.e., $T$ is positive definite. Then there exists a unique prime number $p$ such that $\mathcal{I}_{T}$ is a finite scheme with support lying over the supersingular locus of $\mathcal{M}_{p}^{(2)}=\mathcal{M}^{(2)} \otimes_{\mathbb{Z}} \mathbb{F}_{p}$.
Proof. - Let $\left(E, E^{\prime}\right) \in \mathcal{M}^{(2)}$ be in the image of $\mathcal{T}_{T}$. Since $\operatorname{Hom}\left(E, E^{\prime}\right)$ has rank at least 3 , it follows that $\left(E, E^{\prime}\right)$ has to be a pair of supersingular elliptic curves in some positive characteristic $p$. To see that $p$ is uniquely determined by $T$, note that $T$ is
represented by the quadratic space over $\mathbb{Q}$ corresponding to the definite quaternion algebra ramified in $p$. However, by $[\mathbf{K u} 3$, Prop. 1.3], there is only one quadratic space with fixed discriminant which represents $T$.
1.2. In this chapter we consider, for $T \in \operatorname{Sym}_{3}(\mathbb{Z})_{>0}^{\vee}$, the number

$$
\widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right)=\lg \left(\mathcal{T}_{T}\right) \cdot \log p
$$

where $p$ is the unique prime in the statement of Lemma 1.1, and where

$$
\lg \left(\mathcal{T}_{T}\right)=\sum_{\xi \in \mathcal{T}_{T}\left(\mathbb{F}_{p}\right)} e_{\xi}^{-1} \cdot \lg \left(\mathcal{O}_{\mathcal{T}_{T}, \xi}\right)
$$

with $e_{\xi}=|\operatorname{Aut}(\xi)|$. Our aim is to compare $\widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right)$ with the $T^{\text {th }}$ Fourier coefficient of a certain Siegel-Eisenstein series of genus 3 and weight 2 .

We first define a class of Eisenstein series, among which will be the one appearing in our main theorem.

## 2. Eisenstein series and the main theorem

2.1. Let $B$ be a quaternion algebra over $\mathbb{Q}$. We denote by $V=V_{B}$ the quadratic space defined by $B$, i.e., $B$ with its norm form $Q$. We note that the idèle class character usually associated to a quadratic space, $x \mapsto\left(x,(-1)^{n(n-1) / 2} \operatorname{det}(V)\right)_{\mathbb{Q}}$ is in this case the trivial character $\chi_{0}(4 \mid n$, and $\operatorname{det}(V)$ is a square $)$. Let $H=\mathrm{O}(V)$ be the associated orthogonal group. Let $W=\mathbb{Q}^{6}$, with standard symplectic form $\langle$,$\rangle whose$ matrix with respect to the standard basis is given by $\left(\begin{array}{cc}0 & I_{3} \\ -I_{3} & 0\end{array}\right)$. Let $G=\operatorname{Sp}(W)=\operatorname{Sp}_{6}$, and denote by $P=M . N$ the Siegel parabolic subgroup, with

$$
\begin{aligned}
M & =\left\{\left.m(a)=\left(\begin{array}{cc}
a & 0 \\
0 & { }^{t} a^{-1}
\end{array}\right) \right\rvert\, a \in \mathrm{GL}_{3}\right\}, \\
N & =\left\{\left.n(b)=\left(\begin{array}{cc}
1 & b \\
0 & 1
\end{array}\right) \right\rvert\, b \in \operatorname{Sym}_{3}\right\}
\end{aligned}
$$

Let $K=K_{\infty} \cdot K_{f}=\prod_{v} K_{v}$ be the maximal compact subgroup of $G(\mathbb{A})$ with

$$
K_{v}= \begin{cases}\operatorname{Sp}_{6}\left(\mathbb{Z}_{p}\right), & \text { if } v=p<\infty  \tag{2.1}\\
\left\{\left.\left(\begin{array}{cc}
a & b \\
-b & a
\end{array}\right) \right\rvert\, a+i b \in \mathrm{U}_{3}(\mathbb{R})\right\}, & \text { if } v=\infty\end{cases}
$$

We have the Weil representation $\omega$ of $G(\mathbb{A}) \times H(\mathbb{A})$ (for the standard additive character $\psi$ of $\mathbb{A}$ with archimedean component $\psi_{\infty}(x)=\exp (2 \pi i x)$ and of conductor zero at all non-archimedean places) on the Schwartz space $\mathcal{S}\left(V(\mathbb{A})^{3}\right)$ (the action of the elements $P(\mathbb{A}) \times H(\mathbb{A})$ are given by simple formulae $[\mathbf{W e}]$, comp. also (4.1) and (4.2) below). In the local version at a place $v$, we have a representation $\omega_{v}$ of $G\left(\mathbb{Q}_{v}\right) \times H\left(\mathbb{Q}_{v}\right)$ on $\mathcal{S}\left(V\left(\mathbb{Q}_{v}\right)^{3}\right)$.

We have the Iwasawa decomposition

$$
G(\mathbb{A})=P(\mathbb{A}) K=N(\mathbb{A}) M(\mathbb{A}) K
$$

If $g=n m(a) k \in G(\mathbb{A})$, then

$$
|a(g)|=|\operatorname{det}(a)|_{\mathbb{A}}
$$

is well-defined. For a character $\chi$ of $\mathbb{A}^{\times} / \mathbb{Q}^{\times}$, we have the induced representation of $G(\mathbb{A})$, corresponding to $s \in \mathbb{C}$,

$$
\begin{aligned}
& I(s, \chi)=\{\Phi: G(\mathbb{A}) \rightarrow \mathbb{C} K \text {-finite function } \mid \\
&\left.\Phi(n m(a) g)=\chi(\operatorname{det}(a)) \cdot|a(g)|^{s+2} \cdot \Phi(g)\right\}
\end{aligned}
$$

For $\varphi \in \mathcal{S}\left(V(\mathbb{A})^{3}\right)$, we set

$$
\Phi(g, s)=(\omega(g) \varphi)(0) \cdot|a(g)|^{s}
$$

In this way, we obtain an intertwining map

$$
\begin{equation*}
\mathcal{S}\left(V(\mathbb{A})^{3}\right) \longrightarrow I\left(0, \chi_{0}\right), \quad \varphi \longmapsto \Phi(g, 0) \tag{2.2}
\end{equation*}
$$

Note that $|a(g)|$ is a right $K$-invariant function on $G(\mathbb{A})$, so $\Phi(g, s)$ is a standard section of the induced representation, i.e., its restriction to $K$ is independent of $s$. We will also need the local version $I\left(s, \chi_{v}\right)$ of the induced representation at a place $v$ and the $G\left(\mathbb{Q}_{v}\right)$-equivariant map

$$
\begin{equation*}
\mathcal{S}\left(V_{v}^{3}\right) \longrightarrow I\left(0, \chi_{0, v}\right) . \tag{2.3}
\end{equation*}
$$

2.2. Returning to the global situation, we consider the Eisenstein series associated to $\varphi \in \mathcal{S}\left(V(\mathbb{A})^{3}\right)$,

$$
E(g, s, \Phi)=\sum_{\gamma \in P(\mathbb{Q}) \backslash G(\mathbb{Q})} \Phi(\gamma g, s) .
$$

This series is absolutely convergent for $\operatorname{Re}(s)>2$, and defines an automorphic form. It has a meromorphic continuation and a functional equation with $s=0$ as its center of symmetry.

We will now make a specific choice of $\Phi$ which will define an incoherent Eisenstein series. Let $B=\mathrm{M}_{2}(\mathbb{Q})$ and let $V\left(\mathbb{Z}_{p}\right)=\mathrm{M}_{2}\left(\mathbb{Z}_{p}\right)$ for any $p$. We let $\varphi_{f}=\otimes \varphi_{p}=$ $\otimes \operatorname{char} V\left(\mathbb{Z}_{p}\right)$, and let $\Phi_{f}=\otimes \Phi_{p}$ be the corresponding factorizable standard section. For $\Phi_{\infty}$ we take the standard section uniquely determined by

$$
\Phi_{\infty}(k, 0)=\operatorname{det}(\underline{k})^{2}
$$

where $k \in K_{\infty}$ is the image of $\underline{k} \in \mathrm{U}_{3}(\mathbb{R})$ under the natural identification in (2.1). Then by $[\mathbf{K u 3},(7.13)], \Phi_{\infty}$ is the image of the Gaussian $\varphi_{\infty}$ under the local map (2.3) for $v=\infty$, where the local quadratic space is $V_{\infty}^{+}$, the positive-definite quadratic space corresponding to the Hamilton quaternion algebra over $\mathbb{R}$, and where

$$
\begin{equation*}
\varphi_{\infty}(x)=\exp (-\pi \operatorname{tr}(x, x)), \quad x \in\left(V_{\infty}^{+}\right)^{3} \tag{2.4}
\end{equation*}
$$

Since $V_{\infty}^{+} \otimes V\left(\mathbb{A}_{f}\right)$ does not correspond to a quaternion algebra over $\mathbb{Q}$, the standard section $\Phi=\Phi_{\infty} \otimes \Phi_{f}$ is incoherent in the sense of loc. cit., and hence (loc. cit., Theorem 2.2),

$$
E(g, 0, \Phi) \equiv 0
$$

Consider the Fourier expansion of $E(g, s, \Phi)$,

$$
E(g, s, \Phi)=\sum_{T \in \operatorname{Sym}_{3}(\mathbb{Q})} E_{T}(g, s, \Phi),
$$

where

$$
E_{T}(g, s, \Phi)=\int_{N(\mathbb{Q}) \backslash N(\mathbb{A})} E(n g, s, \Phi) \cdot \psi_{T}(n)^{-1} d n
$$

with

$$
\begin{equation*}
\psi_{T}(n(b))=\psi(\operatorname{tr}(T b)), \quad b \in \operatorname{Sym}_{3}(\mathbb{A}) . \tag{2.5}
\end{equation*}
$$

For $T \in \operatorname{Sym}_{3}(\mathbb{Q})$ with $\operatorname{det}(T) \neq 0$, the Fourier coefficient has an explicit expression as a product

$$
\begin{equation*}
E_{T}(g, s, \Phi)=\prod_{v} W_{T, v}\left(g_{v}, s, \Phi_{v}\right) \tag{2.6}
\end{equation*}
$$

see $[\mathbf{K u 3},(4.4)]$. Here $W_{T, v}\left(g_{v}, s, \Phi_{v}\right)$ is the local Whittaker function, cf. section 5. The local Whittaker functions are entire (cf. [Ku3, (4.2) and (4.3)]), and the product (2.6) is absolutely convergent and holomorphic in $s=0$. More precisely, for $\operatorname{Re}(s)>2$, the identity (2.6) holds and for almost all places $p$, the local factor at $p$ on the right hand side equals $\zeta_{p}(s+2)^{-1} \cdot \zeta_{p}(2 s+2)^{-1}=\left(1-p^{-s-2}\right) \cdot\left(1-p^{-2 s-2}\right)$, and for all places the local factor is an entire function.

### 2.3. For $T \in \operatorname{Sym}_{3}(\mathbb{Q})_{>0}$, let

$$
\operatorname{Diff}(T, V)=\left\{p \mid T \text { not represented by } V\left(\mathbb{Q}_{p}\right)\right\}
$$

Then the cardinality $|\operatorname{Diff}(T, V)|$ is odd, cf. [Ku3, Cor. 5.2]. Moreover we have $W_{T, p}\left(g_{p}, 0, \Phi_{p}\right) \equiv 0$ for $p \in \operatorname{Diff}(T, V)$, cf. [Ku3, Prop. 1.4]. On the other hand, $W_{T, \infty}\left(g_{\infty}, 0, \Phi_{\infty}\right) \neq 0$, cf. [Ku3, Prop. 9.5]. Hence

$$
\underset{s=0}{\operatorname{ord}} E_{T}(g, s, \Phi) \geq|\operatorname{Diff}(T, V)| .
$$

In particular, if $E_{T}^{\prime}(g, 0, \Phi) \neq 0$, then $\operatorname{Diff}(T, V)=\{p\}$ for a unique prime number $p$.
2.4. We may now formulate our main theorem.

Theorem 2.1. - Let $V=\mathrm{M}_{2}(\mathbb{Q})$ and let $\Phi=\Phi_{\infty} \otimes \Phi_{f}$ be the incoherent standard section as above. Let $T \in \operatorname{Sym}_{3}(\mathbb{Q})_{>0}$ with $\operatorname{Diff}(T, V)=\{p\}$.
(i) If $T \notin \operatorname{Sym}_{3}(\mathbb{Z})^{\vee}$, then $\mathcal{T}_{T}=\varnothing$ and $\widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right)=0$ and $E_{T}^{\prime}(g, 0, \Phi) \equiv 0$.
(ii) Let $T \in \operatorname{Sym}_{3}(\mathbb{Z})^{\vee}$. Then $\mathcal{I}_{T}$ has support in characteristic $p$. For $g=$ $\left(g_{\infty}, e, e \ldots\right) \in G(\mathbb{A})$ with

$$
g_{\infty}=\left(\begin{array}{ll}
1 & x \\
& 1
\end{array}\right)\left(\begin{array}{ll}
y^{1 / 2} & \\
& y^{-1 / 2}
\end{array}\right), \quad x, y \in \operatorname{Sym}_{3}(\mathbb{R}), y>0
$$

let $\tau=g_{\infty} \cdot i 1_{3}=x+i y \in \mathfrak{H}_{3}$. Then

$$
\operatorname{det}(y) \widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right) \cdot q^{T}=\kappa \cdot E_{T}^{\prime}(g, 0, \Phi)
$$

where $q^{T}=\exp (2 \pi i \operatorname{tr}(T \tau))$ and where the negative constant $\kappa$ is independent of $T$.

Here $\mathfrak{H}_{3}=\left\{\tau \in \operatorname{Sym}_{3}(\mathbb{C}) \mid \operatorname{Im}(\tau)>0\right\}$ is the Siegel upper half space.
The proof of the theorem consists in calculating explicitly both sides of the identity. The first assertion of (i) is obvious and the second is a consequence of section 5 below, where the local Whittaker functions are related to local representation densities (see Proposition 5.2 below). The proof of (ii) will be reduced in section 3 to a statement about local Whittaker functions which will be taken up in sections 4 and 5 .
2.5. In the rest of this section we relate the adelic Eisenstein series to the classical Siegel-Eisenstein series, following [Ku1, section IV.2]. By strong approximation,

$$
G(\mathbb{A})=G(\mathbb{Q}) G(\mathbb{R}) K
$$

By our choice of $\Phi$, which is right $K_{f}$-invariant, the Eisenstein series $E(g, s, \Phi)$ is determined by its restriction to $G(\mathbb{R})$ (embedded via $g_{\infty} \mapsto\left(g_{\infty}, e, e \ldots\right.$ ) in $G(\mathbb{A})$ ).

We have

$$
G(\mathbb{Z})=G(\mathbb{Q}) \cap\left(G(\mathbb{R}) \cdot K_{f}\right) .
$$

Also, $P(\mathbb{Q}) \backslash G(\mathbb{Q})=P(\mathbb{Z}) \backslash G(\mathbb{Z})$, hence for $g=g_{\infty}$,

$$
\begin{align*}
E(g, s, \Phi) & =\sum_{\gamma \in P(\mathbb{Q}) \backslash G(\mathbb{Q})} \Phi_{\infty}\left(\gamma g_{\infty}, s\right) \cdot \Phi_{f}(\gamma, s)  \tag{2.7}\\
& =\sum_{\gamma \in P(\mathbb{Z}) \backslash G(\mathbb{Z})} \Phi_{\infty}\left(\gamma g_{\infty}, s\right) .
\end{align*}
$$

For our choice of $\Phi_{\infty}$ and of $g_{\infty}=\left(\begin{array}{ll}1 & x \\ & 1\end{array}\right)\left(\begin{array}{ll}y^{1 / 2} & \\ & y^{-1 / 2}\end{array}\right)$, we have

$$
\Phi_{\infty}\left(\gamma g_{\infty}, s\right)=\operatorname{det}(y)^{\frac{s}{2}+1} \cdot \operatorname{det}(c \tau+d)^{-2} \cdot|\operatorname{det}(c \tau+d)|^{-s},
$$

where

$$
\gamma=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in \operatorname{Sp}_{6}(\mathbb{Z})
$$

Inserting this expression into the sum (2.7), one obtains from [Ku1, IV.2.23], (for $\ell=\rho_{n}=2$ ),

$$
\begin{equation*}
E(g, s, \Phi)=\operatorname{det}(y) \cdot E_{\text {class }}(\tau, s), \tag{2.8}
\end{equation*}
$$

where

$$
E_{\text {class }}(\tau, s)=\operatorname{det}(y)^{s / 2} \sum_{(c, d)} \operatorname{det}(c \tau+d)^{-2} \cdot|\operatorname{det}(c \tau+d)|^{-s}
$$

is the classical Siegel Eisenstein series (the sum here ranges over a complete set of representatives of the equivalence classes of pairs of co-prime symmetric integer matrices).
2.6. Using the comparison (2.8) between the adelic and the classical Eisenstein series, we obtain from Theorem 2.1 the following statement. We consider the Fourier expansion of the classical Eisenstein series,

$$
E_{\text {class }}(\tau, s)=\sum_{T \in \operatorname{Sym}_{3}(\mathbb{Z})^{\vee}} c(T, y, s) q^{T}
$$

Here $\tau=x+i y \in \mathfrak{H}_{3}$ and $q^{T}=\exp (2 \pi i \operatorname{tr}(T \tau))$.
Theorem 2.2. - Let $T \in \operatorname{Sym}_{3}(\mathbb{Z})_{>0}^{\vee}$.
(1) Then $c^{\prime}(T)=\left(\frac{\partial}{\partial s} c(T, y, s)\right)_{\mid s=0}$ is independent of $y$.
(2) If $\operatorname{Diff}(T, V)=\{p\}$, then $\mathcal{T}_{T}$ has support in characteristic $p$ and

$$
\widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right)=\kappa \cdot c^{\prime}(T)
$$

for a negative constant $\kappa$ independent of $T$.
Corollary 2.3. - Assume that there is no positive definite binary quadratic form over $\mathbb{Z}$ which represents $m_{1}, m_{2}$ and $m_{3}$, so that the divisors $\mathcal{T}_{m_{1}}, \mathcal{T}_{m_{2}}$, and $\mathcal{T}_{m_{3}}$ intersect in dimension 0 , cf. [Go2, Prop. 3.2]. Then there exists a constant $\kappa$ independent of $\left(m_{1}, m_{2}, m_{3}\right)$ such that

$$
\left(\mathcal{T}_{m_{1}} \cdot \mathcal{T}_{m_{2}} \cdot \mathcal{T}_{m_{3}}\right)=\kappa \cdot \sum_{\substack{T \in \operatorname{Sym}_{3}(\mathbb{Z})_{>0}^{\vee} \\ \operatorname{diag}(T)=\left(m_{1}, m_{2}, m_{3}\right)}} c^{\prime}(T)
$$

Proof. - The hypothesis implies that in the disjoint sum (1.1) only positive definite $T \in \operatorname{Sym}_{3}(\mathbb{Z})^{\vee}$ occur as indices, comp. [Go2, Prop. 3.5]. Therefore the assertion follows from Theorem 2.2.

## 3. Use of the Siegel-Weil formula

3.1. Let $\tilde{V}$ be the quadratic space associated to a quaternion algebra $\tilde{B}$ over $\mathbb{Q}$. For $\tilde{\varphi} \in \mathcal{S}\left(\tilde{V}(\mathbb{A})^{3}\right)$, there is the theta series

$$
\theta(g, h ; \tilde{\varphi})=\sum_{x \in \tilde{V}(\mathbb{Q})^{3}}(\omega(g) \tilde{\varphi})\left(h^{-1} x\right)
$$

and the corresponding theta integral over the orthogonal group $\tilde{H}=\mathrm{O}(\tilde{V})$ associated to $\tilde{V}$,

$$
I(g ; \tilde{\varphi})=\int_{\tilde{H}(\mathbb{Q}) \backslash \tilde{H}(\mathbb{A})} \theta(g, h ; \tilde{\varphi}) d h .
$$

Here the Haar measure $d h$ is normalized so that

$$
\operatorname{vol}(\tilde{H}(\mathbb{Q}) \backslash \tilde{H}(\mathbb{A}))=1
$$

We will only consider the case in which the quadratic space $\tilde{V}$ is anisotropic. If $\tilde{\varphi}$ is $K$-finite, then $I(g ; \tilde{\varphi})$ is an automorphic form on $G(\mathbb{A})$. The Siegel-Weil formula $[\mathbf{K R}]$ states that, if $\tilde{\varphi}$ gives rise to $\tilde{\Phi}$ via the map (2.2), then

$$
\begin{equation*}
E(g, 0, \tilde{\Phi})=2 \cdot I(g ; \tilde{\varphi}) \tag{3.1}
\end{equation*}
$$

Let $T \in \operatorname{Sym}_{3}(\mathbb{Q})$ with $\operatorname{det}(T) \neq 0$. Then the $T^{\text {th }}$-Fourier coefficient of $I(g ; \tilde{\varphi})$ is equal to ([KR, (6.21)])

$$
I_{T}(g ; \tilde{\varphi})=\int_{\tilde{H}(\mathbb{Q}) \backslash \tilde{H}(\mathbb{A})} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}}(\omega(g) \tilde{\varphi})\left(h^{-1} x\right) d h
$$

where

$$
\tilde{V}(\mathbb{Q})_{T}^{3}=\left\{x \in \tilde{V}(\mathbb{Q})^{3} \left\lvert\, \frac{1}{2}(x, x)=T\right.\right\} .
$$

3.2. We now return to the situation considered in Theorem 2.1. Let $V=\mathrm{M}_{2}(\mathbb{Q})$ and let $\Phi$ be the standard section defined in the previous section. We also fix $T \in$ $\operatorname{Sym}_{3}(\mathbb{Q})_{>0}$ with $\operatorname{Diff}(T, V)=\{p\}$. Let $\tilde{V}$ be the quadratic space associated to the definite quaternion algebra $\tilde{B}=B^{(p)}$ ramified at $p$, and unramified at all other finite primes. Note that $\tilde{V}(\mathbb{R})=V_{\infty}^{+}$. We consider the standard section $\tilde{\Phi}$ which is the image of $\tilde{\varphi}=\tilde{\varphi}_{\infty} \otimes \tilde{\varphi}_{f}^{p} \otimes \tilde{\varphi}_{p}$ under

$$
\mathcal{S}\left(\tilde{V}(\mathbb{A})^{3}\right) \longrightarrow I\left(0, \chi_{0}\right),
$$

where $\tilde{\varphi}_{f}^{p}=\varphi_{f}^{p}$, where $\tilde{\varphi}_{\infty}=\varphi_{\infty}$ is the Gaussian (2.4) and where $\tilde{\varphi}_{p}=\operatorname{char} \tilde{V}\left(\mathbb{Z}_{p}\right)^{3}$, with $\tilde{V}\left(\mathbb{Z}_{p}\right)$ the maximal order of the division algebra $B^{(p)} \otimes \mathbb{Q}_{p}$. Hence $\tilde{\Phi}_{\infty}=\Phi_{\infty}$, $\tilde{\Phi}_{f}^{p}=\Phi_{f}^{p}$ and $\tilde{\Phi}$ is a coherent standard section. Comparing the expressions (2.6) for the Fourier coefficients of $E(g, s, \Phi)$ and $E(g, s, \tilde{\Phi})$, we can write, for $g=g_{\infty} \in G(\mathbb{R})$,

$$
E_{T}^{\prime}(g, 0, \Phi)=\frac{W_{T, p}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T, p}\left(e, 0, \tilde{\Phi}_{p}\right)} \cdot E_{T}(g, 0, \tilde{\Phi})
$$

We refer to Corollary 5.3 below for a proof of the fact that the denominator here is nonzero. Using the Siegel-Weil formula (3.1) for the anisotropic quadratic space $\tilde{V}$, we can rewrite this as

$$
\begin{equation*}
E_{T}^{\prime}(g, 0, \Phi)=2 \cdot \frac{W_{T, p}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T, p}\left(e, 0, \tilde{\Phi}_{p}\right)} \cdot I_{T}(g ; \tilde{\varphi}) \tag{3.2}
\end{equation*}
$$

Now the function $\tilde{\varphi}_{\infty}$ is invariant under $\tilde{H}(\mathbb{R})$. For

$$
g_{\infty}=\left(\begin{array}{ll}
1 & x \\
& 1
\end{array}\right)\left(\begin{array}{ll}
y^{1 / 2} & \\
& y^{-1 / 2}
\end{array}\right), \quad x, y \in \operatorname{Sym}_{3}(\mathbb{R}), y>0
$$

the value of $\omega\left(g_{\infty}\right) \tilde{\varphi}_{\infty}$ at $t \in \tilde{V}(\mathbb{R})^{3}$ with $\frac{1}{2}(t, t)=T$ is equal to

$$
\left(\omega\left(g_{\infty}\right) \tilde{\varphi}_{\infty}\right)(t)=\exp (2 \pi i \operatorname{tr}(T \tau)) \operatorname{det}(y)
$$

Since $\tilde{H}(\mathbb{R})=\mathrm{O}(\tilde{V}(\mathbb{R}))$ is compact, we may write using the product measure $d h=$ $d_{\infty} h \times d_{f} h$,

$$
\begin{equation*}
2 \cdot I_{T}(g ; \tilde{\varphi})=2 \operatorname{det}(y) \cdot q^{T} \cdot \operatorname{vol}\left(\tilde{H}(\mathbb{R}), d_{\infty} h\right) \cdot I_{T}\left(\tilde{\varphi}_{f}\right) \tag{3.3}
\end{equation*}
$$

where

$$
I_{T}\left(\tilde{\varphi}_{f}\right)=\int_{\tilde{H}(\mathbb{Q}) \backslash \tilde{H}\left(\mathbb{A}_{f}\right)} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}} \tilde{\varphi}_{f}\left(h^{-1} x\right) d_{f} h
$$

### 3.3. Let

$$
\tilde{H}^{\prime}=\left\{\tilde{g}=\left(g, g^{\prime}\right) \in \tilde{B}^{\times} \times \tilde{B}^{\times} \mid \operatorname{Nm}(g)=\operatorname{Nm}\left(g^{\prime}\right)\right\}
$$

Then $\tilde{H}^{\prime}$ acts on $\tilde{V}$ via

$$
\tilde{g} \cdot x=\left(g, g^{\prime}\right) \cdot x=g^{\prime} x g^{-1}
$$

This induces an exact sequence, where $\mathbb{G}_{m}$ lies in the center of $\tilde{H}^{\prime}$, cf. [ $\mathbf{W d} \mathbf{d}$, Lemma 1.6],

$$
\begin{equation*}
1 \rightarrow \mathbb{G}_{m} \rightarrow \tilde{H}^{\prime} \xrightarrow{\mathrm{pr}} \mathrm{SO}(\tilde{V}) \rightarrow 1 \tag{3.4}
\end{equation*}
$$

We fix the Haar measure on $\tilde{H}^{\prime}(\mathbb{A})$ such that the measure induced by the exact sequence (3.4) on $\operatorname{SO}(\tilde{V})(\mathbb{A})$ is the Tamagawa measure, and with the standard Haar measure on the central idele group $\mathbb{A}^{\times}$which is the product of the local measures $\lambda_{\ell} \frac{d x_{\ell}}{\left|x_{\ell}\right|}$ with convergence factors $\lambda_{\ell}=1-\ell^{-1}$, so that $\operatorname{vol}\left(\hat{\mathbb{Z}}^{\times}\right)=1$. Let

$$
\tilde{K}^{\prime}=\tilde{H}^{\prime}\left(\mathbb{A}_{f}\right) \cap\left(\left(O_{\tilde{B}} \otimes \hat{\mathbb{Z}}\right)^{\times} \times\left(O_{\tilde{B}} \otimes \hat{\mathbb{Z}}\right)^{\times}\right)
$$

Proposition 3.1. - Let

$$
\left|\mathcal{T}_{T}\right|=\sum_{\xi \in \mathcal{T}_{T}\left(\mathbb{F}_{p}\right)} e_{\xi}^{-1}
$$

with $e_{\xi}=|\operatorname{Aut}(\xi)|$. Then

$$
\left|\mathcal{T}_{T}\right|=\kappa_{1} \cdot I_{T}\left(\tilde{\varphi}_{f}\right)
$$

where $\kappa_{1}=2 \operatorname{vol}\left(\tilde{K}^{\prime}\right)^{-1}$.
Proof. - We choose a finite set of double coset representatives $h_{j} \in \tilde{H}^{\prime}\left(\mathbb{A}_{f}\right)$ such that

$$
\tilde{H}^{\prime}\left(\mathbb{A}_{f}\right)=\coprod_{j} \tilde{H}^{\prime}(\mathbb{Q}) h_{j} \tilde{K}^{\prime}
$$

Since each double coset $\tilde{H}^{\prime}(\mathbb{Q}) h_{j} \tilde{K}^{\prime}$ is stable under $\widehat{\mathbb{Z}}^{\times} \mathbb{Q}^{\times}=\mathbb{A}_{f}^{\times}$, we obtain a disjoint decomposition,

$$
\mathrm{SO}(\tilde{V})\left(\mathbb{A}_{f}\right)=\coprod_{j} \mathrm{SO}(\tilde{V})(\mathbb{Q}) \operatorname{pr}\left(h_{j}\right) \operatorname{pr}\left(\tilde{K}^{\prime}\right)
$$

Let

$$
\tilde{\Gamma}_{j}^{\prime}=\tilde{H}^{\prime}(\mathbb{Q}) \cap h_{j} \tilde{K}^{\prime} h_{j}^{-1}
$$

Note that $\operatorname{vol}(\operatorname{SO}(\tilde{V})(\mathbb{Q}) \backslash \operatorname{SO}(\tilde{V})(\mathbb{A}))=2$. We have

$$
\tilde{H}\left(\mathbb{A}_{f}\right)=\operatorname{SO}(\tilde{V})\left(\mathbb{A}_{f}\right) \rtimes \mu_{2}\left(\mathbb{A}_{f}\right)
$$

Hence

$$
\begin{aligned}
1 & =\operatorname{vol}(\mathrm{O}(\tilde{V})(\mathbb{Q}) \backslash \mathrm{O}(\tilde{V})(\mathbb{A})) \\
& =\frac{1}{2} \operatorname{vol}(\mathrm{SO}(\tilde{V})(\mathbb{Q}) \backslash \mathrm{O}(\tilde{V})(\mathbb{A})) \\
& =\frac{1}{2} \operatorname{vol}(\mathrm{SO}(\tilde{V})(\mathbb{Q}) \backslash \mathrm{SO}(\tilde{V})(\mathbb{A})) \operatorname{vol}\left(\mu_{2}(\mathbb{A})\right) \\
& =\operatorname{vol}\left(\mu_{2}(\mathbb{A})\right)
\end{aligned}
$$

and therefore

$$
\operatorname{vol}\left(\mu_{2}(\mathbb{Q}) \backslash \mu_{2}(\mathbb{A})\right)=\frac{1}{2}
$$

Let us normalize the Haar measure on $\mu_{2}(\mathbb{R})$ by $\operatorname{vol}\left(\mu_{2}(\mathbb{R})\right)=1$. Then we get $\operatorname{vol}\left(\mu_{2}(\mathbb{Q}) \backslash \mu_{2}\left(\mathbb{A}_{f}\right)\right)=\frac{1}{2}$. Then we obtain as in $[\mathbf{K u} \mathbf{3},(7.28)]$,

$$
\begin{aligned}
I_{T}\left(\tilde{\varphi}_{f}\right) & =\int_{\operatorname{SO}(\tilde{V})(\mathbb{Q}) \backslash \operatorname{SO}(\tilde{V})\left(\mathbb{A}_{f}\right)} \int_{\mu_{2}(\mathbb{Q}) \backslash \mu_{2}\left(\mathbb{A}_{f}\right)} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}} \tilde{\varphi}_{f}\left(h^{-1} c x\right) d_{f} h d c \\
& =\frac{1}{2} \int_{\operatorname{SO}(\tilde{V})(\mathbb{Q}) \backslash \operatorname{SO}(\tilde{V})\left(\mathbb{A}_{f}\right)} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}} \tilde{\varphi}_{f}\left(h^{-1} x\right) d_{f} h \\
& =\frac{1}{2} \sum_{j} \int_{\operatorname{SO}(\tilde{V})(\mathbb{Q}) \backslash \operatorname{SO}(\tilde{V})(\mathbb{Q}) \operatorname{pr}\left(h_{j}\right) \operatorname{pr}\left(\tilde{K}^{\prime}\right)} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}} \tilde{\varphi}_{f}\left(h^{-1} x\right) d_{f} h \\
& =\frac{1}{2} \cdot \operatorname{vol}\left(\operatorname{pr}\left(\tilde{K}^{\prime}\right)\right) \cdot \sum_{j} \sum_{x \in \tilde{V}(\mathbb{Q})_{T}^{3}} \frac{1}{\left|\tilde{\Gamma}_{j, x}\right|} \cdot \tilde{\varphi}_{f}\left(h_{j}^{-1} x\right) .
\end{aligned}
$$

Here $\tilde{\Gamma}_{j, x}$ is the image of $\tilde{\Gamma}_{j, x}^{\prime}$ in $\mathbb{Q}^{\times} \backslash \tilde{H}^{\prime}(\mathbb{Q})=\operatorname{SO}(\tilde{V})(\mathbb{Q})$. Therefore we have $\left|\tilde{\Gamma}_{j, x}\right|=$ $\frac{1}{2} \cdot\left|\tilde{\Gamma}_{j, x}^{\prime}\right|$. Note that $\tilde{\Gamma}_{j, x}$ is trivial since $x$ spans a three-dimensional subspace of the 4-dimensional space $\tilde{V}$.

To make the connection with $\mathcal{T}_{T}$, note that the supersingular locus of $\mathcal{M}_{p}^{(2)}$ can be written as a double coset space (cf. $[\mathbf{M i}, 6]$ ),

$$
\left(\mathcal{M}^{(2)}\right)^{s s}=\tilde{H}^{\prime}(\mathbb{Q}) \backslash \tilde{H}^{\prime}\left(\mathbb{A}_{f}\right) / \tilde{K}^{\prime}
$$

Here we chose $\left(E_{0}, E_{0}\right)$ as a base point, such that $\tilde{K}^{\prime}$ is the stabilizer of the Tate module $\hat{T}\left(E_{0}\right) \times \hat{T}\left(E_{0}\right)$ (completed by the Dieudonné module at $p$ ). To $\tilde{g}=\left(g, g^{\prime}\right) \in$ $\tilde{H}^{\prime}\left(\mathbb{A}_{f}\right)$ corresponds $E_{g} \times E_{g^{\prime}}$ with the diagonal isogeny,

$$
\left(g, g^{\prime}\right): E_{0} \times E_{0} \longrightarrow E_{g} \times E_{g^{\prime}}
$$

The lattice $\operatorname{Hom}\left(E_{g}, E_{g^{\prime}}\right)$ in $\tilde{V}(\mathbb{Q})=\operatorname{Hom}\left(E_{0}, E_{0}\right) \otimes \mathbb{Q}$ is given by

$$
\begin{aligned}
\operatorname{Hom}\left(E_{g}, E_{g^{\prime}}\right) & =\left\{y \in \tilde{B} \mid y g\left(\hat{T}\left(E_{0}\right)\right) \subset g^{\prime} \hat{T}\left(E_{0}\right)\right\} \\
& =\left\{y \in \tilde{B} \mid g^{\prime-1} y g \in \tilde{V}(\hat{\mathbb{Z}})\right\} \\
& =\left\{y \in \tilde{B} \mid \tilde{g}^{-1} y \in \tilde{V}(\hat{\mathbb{Z}})\right\}
\end{aligned}
$$

Hence we obtain

$$
\begin{aligned}
\left|\mathcal{T}_{T}\right| & =\sum_{[y, \tilde{g}] \in \tilde{H}^{\prime}(\mathbb{Q}) \backslash\left(\tilde{V}^{3}(\mathbb{Q})_{T} \times \tilde{H}^{\prime}\left(\mathbb{A}_{f}\right) / \tilde{K}\right)} \tilde{\varphi}_{f}\left(\tilde{g}^{-1} \cdot y\right) \\
& =\sum_{j} \sum_{x \in \tilde{V}^{3}(\mathbb{Q})_{T}} \tilde{\varphi}_{f}\left(h_{j}^{-1} \cdot x\right) \\
& =2 \cdot \operatorname{vol}\left(\operatorname{pr}\left(\tilde{K}^{\prime}\right)\right)^{-1} \cdot I_{T}\left(\tilde{\varphi}_{f}\right)
\end{aligned}
$$

Since $\operatorname{vol}\left(\tilde{K}^{\prime}\right)=\operatorname{vol}\left(\operatorname{pr}\left(\tilde{K}^{\prime}\right)\right)$, the result follows.
3.4. The next result will be proved in section 5.6.

Theorem 3.2. - The lengths of the local rings $\mathcal{O}_{\mathcal{T}_{T, \xi}}$ at all points $\xi \in \mathcal{T}_{T}\left(\overline{\mathbb{F}}_{p}\right)$ are all equal to

$$
\lg \left(\mathcal{O}_{\mathcal{T}_{T, \xi}}\right)=-\frac{2}{(p-1)^{2}} \cdot \frac{W_{T, p}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T, p}\left(e, 0, \tilde{\Phi}_{p}\right)} \cdot(\log p)^{-1}
$$

3.5. We will now prove Theorem 2.1 using Theorem 3.2. Let

$$
\begin{aligned}
& H^{\prime}=\left\{\tilde{g}=\left(g, g^{\prime}\right) \in \mathrm{GL}_{2} \times \mathrm{GL}_{2} \mid \operatorname{det}(g)=\operatorname{det}\left(g^{\prime}\right)\right\} \\
& K^{\prime}=H^{\prime}\left(\mathbb{A}_{f}\right) \cap\left(\mathrm{GL}_{2}(\hat{\mathbb{Z}}) \otimes \mathrm{GL}_{2}(\hat{\mathbb{Z}})\right)
\end{aligned}
$$

Then $\tilde{H}^{\prime}$ is an inner form of $H^{\prime}$.
We now fix Haar measures on $\tilde{H}^{\prime}(\mathbb{A})$ and on $H^{\prime}(\mathbb{A})$ following $[\mathbf{K u 3}$, p. 573]. More precisely, in loc. cit. Kudla defines for any quaternion algebra $B$ over $\mathbb{Q}$ a Haar measure on $(B \otimes \mathbb{A})^{\times}$which is decomposed, i.e., the explicit product of local Haar measures on $\left(B \otimes \mathbb{Q}_{v}\right)^{\times}$for all places $v$. By our fixed choice of Haar measure on $\mathbb{A}^{\times}$, we therefore also obtain a decomposed Haar measure on $H(B)^{\prime}(\mathbb{A})$, where

$$
H(B)^{\prime}=\left\{\tilde{g}=\left(g, g^{\prime}\right) \in B^{\times} \times B^{\times} \mid \operatorname{Nm}(g)=\operatorname{Nm}\left(g^{\prime}\right)\right\}
$$

By loc. cit., the induced Haar measure on $S O(V(B))(\mathbb{A})$ is the Tamagawa measure, as required above.

We apply this construction to $B=\mathrm{M}_{2}(\mathbb{Q})$ and to $B=\tilde{B}=B^{(p)}$, the definite quaternion algebra, ramified at $p$ and unramified at all other finite places. Then we have for these Haar measures (comp. [Ku3, Lemma 14.10]),

$$
\frac{\operatorname{vol}\left(K_{p}^{\prime}\right)}{\operatorname{vol}\left(\tilde{K}_{p}^{\prime}\right)}=(p-1)^{2}
$$

and

$$
\frac{\operatorname{vol}\left(K^{\prime}\right)}{\operatorname{vol}\left(\tilde{K}^{\prime}\right)}=(p-1)^{2} .
$$

Hence

$$
\begin{aligned}
q^{T} \cdot \widehat{\operatorname{deg}}\left(\mathcal{T}_{T}\right) & =q^{T} \lg \left(\mathcal{O}_{\mathcal{T}_{T}}\right) \cdot \log p \\
& =q^{T} \lg \left(\mathcal{O}_{\mathcal{T}_{T, \xi}}\right) \cdot\left|\mathcal{T}_{T}\right| \cdot \log p \\
& =-\frac{2}{(p-1)^{2} \cdot \operatorname{vol}\left(\tilde{K}^{\prime}\right)} \cdot 2 \cdot \frac{W_{T, p}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T, p}\left(e, 0, \tilde{\Phi}_{p}\right)} \cdot q^{T} \cdot I_{T}\left(\tilde{\varphi}_{f}\right) \\
& =-\frac{2}{\operatorname{vol}\left(K^{\prime}\right)} \cdot E_{T}^{\prime}(g, 0, \Phi) \operatorname{det}(y)^{-1} v^{-1}
\end{aligned}
$$

where we used (3.2) and (3.3) in the last step, and where $v=\operatorname{vol}\left(\tilde{H}(\mathbb{R}), d_{\infty} h\right)$. This proves the main theorem with the negative constant $\kappa=-\frac{2}{\operatorname{vol}\left(K^{\prime}\right)} \cdot v^{-1}$.

## 4. The Weil representation

4.1. The remainder of this chapter is devoted to the proof of Theorem 3.2. This is a purely local statement.

We fix a prime number $p$ and change our notation: We replace $V$ by $V \otimes \mathbb{Q}_{p}, G$ by $G \otimes \mathbb{Q}_{p}, \psi$ by its localization $\psi_{p}$ (of conductor zero), etc. At the same time we consider a more general situation.
4.2. Instead of the quadratic space associated to a quaternion algebra, we now let $V$ be any $\mathbb{Q}_{p}$-vector space and (, ) a symmetric nondegenerate bilinear form on $V$. Then $Q(x)=\frac{1}{2}(x, x)$ is a quadratic form on $V$.

We assume that $n:=\operatorname{dim}(V)$ is even. In fact, we will only need the case $V=B \perp$ $H^{r}$ where $B$ is a quaternion algebra over $\mathbb{Q}_{p}$ endowed with the reduced norm, and where $H^{r}$ is the orthogonal sum of $r$ copies of the hyperbolic plane $H$.

We denote by $\operatorname{det}(V)$ the image in $\mathbb{Q}_{p}^{\times} /\left(\mathbb{Q}_{p}^{\times}\right)^{2}$ of the determinant of the matrix $\left(\left(v_{i}, v_{j}\right)\right)_{i j}$ where $\left(v_{1}, \ldots, v_{n}\right)$ is some basis of $V$. As in 2.1 we have the quadratic character $\chi_{V}$ of $\mathbb{Q}_{p}^{\times}$associated to $V$ given by

$$
\chi_{V}(x)=\left(x,(-1)^{n(n-1) / 2} \operatorname{det}(V)\right)_{p}=\left(x,(-1)^{n / 2} \operatorname{det}\left(\left(v_{i}, v_{j}\right)\right)_{i j}\right)_{p},
$$

where $(,)_{p}$ denotes the Hilbert symbol.
4.3. Let $(W,\langle\rangle$,$) be the space \mathbb{Q}_{p}^{2 m}$ endowed with the standard symplectic form whose matrix with respect to the standard basis is given by $\left(\begin{array}{cc}0 & I_{m} \\ -I_{m} & 0\end{array}\right)$. We consider $W$ as vector space of row vectors, in particular the canonical $\mathrm{GL}_{2 m}$-action is from the right. To prove Theorem 3.2 we will need only the case $m=3$.

As in 2.1 we denote by $P=M N$ the Siegel parabolic subgroup of $G=\operatorname{Sp}_{2 m}\left(\mathbb{Q}_{p}\right)$ over $\mathbb{Q}_{p}$ where

$$
\begin{aligned}
& M=\left\{\left.m(a)=\left(\begin{array}{cc}
a & 0 \\
0 & { }^{t} a^{-1}
\end{array}\right) \right\rvert\, a \in \operatorname{GL}_{m}\left(\mathbb{Q}_{p}\right)\right\}, \\
& N=\left\{\left.n(b)=\left(\begin{array}{cc}
I_{m} & b \\
0 & I_{m}
\end{array}\right) \right\rvert\, b \in \operatorname{Sym}_{m}\left(\mathbb{Q}_{p}\right)\right\}
\end{aligned}
$$

Let $K=\operatorname{Sp}_{2 m}\left(\mathbb{Z}_{p}\right) \subset G$ the standard maximal compact subgroup and set

$$
w:=\left(\begin{array}{cc}
0 & I_{m} \\
-I_{m} & 0
\end{array}\right) \in G
$$

4.4. In the sequel we let $a \in \mathrm{GL}_{m}$ act on $V^{m}=V \otimes \mathbb{Q}_{p}^{m}$ via right multiplication, which we denote by $x \mapsto x a$.

Moreover for $x, y \in V^{m}$ we set

$$
(x, y):=\left(\left(x_{i}, y_{j}\right)\right)_{i j} \in \operatorname{Sym}_{m}\left(\mathbb{Q}_{p}\right) .
$$

4.5. Associated to the quadratic space $V$ and the fixed additive character $\psi$ there is a Weil representation $\omega_{V}$ of $G$ on the vector space $\mathcal{S}\left(V^{m}\right)$ of Schwartz functions on $V^{m}$. For $g=\left(\begin{array}{cc}a & b \\ c & d\end{array}\right) \in G, \varphi \in \mathcal{S}\left(V^{m}\right)$, and $x \in V^{m}$ we have by [Ku2, Prop. 4.3] (cf. also [Rao, Lemma 3.2], and [We]),

$$
\begin{aligned}
\left(\omega_{V}(g)(\varphi)\right)(x) & =\gamma(V, \psi, g) \\
& \cdot \int_{V^{m} / \operatorname{Ker}(c)} \psi\left(\operatorname{tr}\left(\frac{1}{2}(x a, x b)+(x b, y c)+\frac{1}{2}(y c, y d)\right)\right) \varphi(x a+y c) d_{g} y
\end{aligned}
$$

where $\gamma(V, \psi, g)$ is a certain $8^{\text {th }}$ root of unity depending on $V, \psi$, and $g$ such that $\gamma(V, \psi, e)=1$ and where $d_{g} y$ is a suitable Haar measure. We make this more explicit in three special cases:

$$
\begin{align*}
\left(\omega_{V}(m(a)) \varphi\right)(x) & =\chi_{V}(\operatorname{det} a)|\operatorname{det} a|^{n / 2} \varphi(x a)  \tag{4.1}\\
\left(\omega_{V}(n(b)) \varphi\right)(x) & =\psi\left(\frac{1}{2} \operatorname{tr}((x, x) b)\right) \varphi(x)  \tag{4.2}\\
\left.\left(\omega_{V}\left(w^{-1}\right)\right) \varphi\right)(x) & =\gamma(V) \int_{V^{m}} \psi(-\operatorname{tr}((x, y))) \varphi(y) d y \tag{4.3}
\end{align*}
$$

where in (4.3) $d y$ is the Haar measure on $V^{m}$ which is self dual for Fourier transform and where $\gamma(V)=\gamma\left(V, \psi, w^{-1}\right)$ is the $8^{\text {th }}$ root of unity explicitly given in [Ku3, A.4].

## 5. Local Whittaker functions and representation densities

5.1. We keep the notation of section 4 and assume from now on that $m=3$ and hence $G=\mathrm{Sp}_{6}\left(\mathbb{Q}_{p}\right)$, and $n=4$.

For $s \in \mathbb{C}$ let $I\left(s, \chi_{V}\right)$ be the degenerate principal series representation of $G$ induced from $P$, i.e., $I\left(s, \chi_{V}\right)$ consists of $K$-finite functions $\Phi(\cdot, s): G \rightarrow \mathbb{C}$ such that

$$
\Phi(n m(a) g, s)=\chi_{V}(\operatorname{det} a)|\operatorname{det} a|^{s+2} \Phi(g, s)
$$

for all $n \in N, a \in \mathrm{GL}_{3}\left(\mathbb{Q}_{p}\right)$, and $g \in G$.
We also set for $T \in \operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)$, as in (2.5),

$$
\psi_{T}: N \longrightarrow \mathbb{C}^{\times}, \quad \psi_{T}(n(b))=\psi(\operatorname{tr}(T b)) .
$$

5.2. For $s \in \mathbb{C}, \Phi \in I\left(s, \chi_{V}\right), T \in \operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)$ with $\operatorname{det}(T) \neq 0$, and $g \in G$ we define the local Whittaker function by

$$
W_{T}(g, s, \Phi)=\int_{N} \Phi\left(w^{-1} n(b) g, s\right) \psi_{T}(n(b))^{-1} d b
$$

where $d b$ is the Haar measure on $\operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)$ which is selfdual with respect to the pairing

$$
\psi_{N}: \operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right) \times \operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right) \longrightarrow \mathbb{C}, \quad\left(b, b^{\prime}\right) \longmapsto \psi\left(\operatorname{tr}\left(b b^{\prime}\right)\right)
$$

As the conductor of $\psi$ is zero, we have

$$
\begin{equation*}
\left\{b \in \operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right) \mid \psi_{N}\left(b, b^{\prime}\right)=1 \text { for all } b^{\prime} \in \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)\right\}=\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee} \tag{5.1}
\end{equation*}
$$

Therefore

$$
\operatorname{vol}_{d b}\left(\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)\right) \operatorname{vol}_{d b}\left(\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}\right)=1
$$

As the index of $\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)$ in $\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}$ is $2^{3 \delta_{2 p}}$, we obtain

$$
\begin{equation*}
\operatorname{vol}_{d b}\left(\operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)\right)=2^{-(3 / 2) \delta_{2 p}} \tag{5.2}
\end{equation*}
$$

It is known that $W_{T}(g, s, \Phi)$ converges for $\operatorname{Re}(s)>2$ and admits a holomorphic continuation to the entire complex plane, if $\Phi$ is standard, i.e., if its restriction to $K$ is independent of $s[\mathbf{K u} 3$, Prop. 1.4]. Moreover, we will see in Proposition 5.2 below that $W_{T}(e, s, \Phi)$ is a polynomial in $p^{-s}$.
5.3. For $\varphi \in \mathcal{S}\left(V^{3}\right)$ we set

$$
\Phi(g, s)=(\omega(g) \varphi)(0) \cdot|a(g)|^{s}
$$

It follows from (4.1) and (4.2) that $\Phi(g, s) \in I\left(s, \chi_{V}\right)$. In this way, we obtain a $G$-equivariant map similar to (2.2),

$$
\mathcal{S}\left(V^{3}\right) \longrightarrow I\left(0, \chi_{V}\right), \quad \varphi \longmapsto \Phi(g, 0)
$$

5.4. For $r \geq 0$ we denote by the $H_{r}$ the quadratic space $\mathbb{Q}_{p}^{2 r}$ whose associated bilinear form has the matrix $\frac{1}{2}\left(\begin{array}{ccc}0 & I_{r} \\ I_{r} & 0\end{array}\right)$ with respect to the standard basis, and set

$$
V_{r}=V \perp H_{r}
$$

It is known [Ku3, Lemma A.2] that $\omega_{V_{r}}=\omega_{V} \otimes \omega_{H_{r}}$ as representations of $G$ on $\mathcal{S}\left(V_{r}^{3}\right)=\mathcal{S}\left(V^{3}\right) \otimes \mathcal{S}\left(H_{r}^{3}\right)$.

We also recall Lemma A. 3 from [Ku3] (see also [Ral, Remark II.3.2]):

Lemma 5.1. - Let $\varphi_{r}^{0} \in \mathcal{S}\left(H_{r}^{3}\right)$ be the characteristic function of $\mathrm{M}_{2 r, 3}\left(\mathbb{Z}_{p}\right)$ and $\varphi \in$ $\mathcal{S}\left(V^{3}\right)$ with associated $\Phi(g, s) \in I\left(s, \chi_{V}\right)$. Set $\varphi^{[r]}=\varphi \otimes \varphi_{r}^{0} \in \mathcal{S}\left(V_{r}^{3}\right)=\mathcal{S}\left(V^{3}\right) \otimes \mathcal{S}\left(H_{r}^{3}\right)$. Then we have for all $g \in G$ and $r \geq 0$

$$
\Phi(g, r)=\left(\omega_{V_{r}}(g) \varphi^{[r]}\right)(0)
$$

5.5. We fix a $\mathbb{Z}_{p}$-lattice $L$ of $V$ such that (, ) is integral on $L$. Choose a $\mathbb{Z}_{p}$-basis of $L$ and let $S_{r}$ be the matrix associated to the quadratic form on $V_{r}=V \oplus H_{r}$ with respect to the chosen basis of $L$ and the standard basis of $H_{r}$. In particular, the matrix of the bilinear form (, ) with respect to the chosen base of $L$ equals $2 S_{0}$.

Let $\varphi \in \mathcal{S}\left(V^{3}\right)$ be the characteristic function of $L^{3}$ with associated $\Phi=\Phi(g, s) \in$ $I\left(s, \chi_{V}\right)$. Then the local Whittaker function $W_{T}(e, s, \Phi)$ interpolates the local representation densities:

Proposition 5.2. - For all integers $r \geq 0$ we have

$$
W_{T}(e, r, \Phi)=2^{-(15 / 2) \delta_{2 p}}\left|\operatorname{det} S_{0}\right|^{3 / 2} \gamma(V) \alpha_{p}\left(T, S_{r}\right)
$$

where we denote by $\alpha_{p}($,$) the local representation density as normalized in [\mathbf{W d 1}$, (4.4)]. In particular, $W_{T}(e, s, \Phi)$ is a polynomial in $X=p^{-s}$.

Proof. - The right hand side is a polynomial in $X=p^{-r}[\mathbf{K i t}]$ and the left hand side is an entire function in $r$. Hence it suffices to show the identity for $r>2$. Now we have

$$
\begin{aligned}
& W_{T}(e, r, \Phi)=\int_{\operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)} \Phi\left(w^{-1} n(b), r\right) \psi(-\operatorname{tr}(T b)) d b \\
& \stackrel{(5.1)}{=} \int_{\operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)}\left(\omega_{V_{r}}\left(w^{-1} n(b)\right) \varphi^{[r]}\right)(0) \psi(-\operatorname{tr}(T b)) d b \\
& \stackrel{(4.3)}{=} \int_{\operatorname{Sym}_{3}\left(\mathbb{Q}_{p}\right)} \gamma(V) \int_{V_{r}^{3}} \psi\left(\frac{1}{2} \operatorname{tr}(b(y, y))\right) \cdot \varphi^{[r]}(y) d y \psi(-\operatorname{tr}(T b)) d b \\
& =\gamma(V) \lim _{t \rightarrow \infty} \int_{p^{-t}} \int_{\mathrm{Sym}_{3}\left(\mathbb{Z}_{p}\right)} \int_{V_{r}^{3}} \psi\left(\operatorname{tr}\left(b\left(\frac{1}{2}(y, y)-T\right)\right)\right) \cdot \varphi^{[r]}(y) d y d b \\
& \stackrel{(5.1)}{=} \gamma(V) \lim _{t \rightarrow \infty} \operatorname{vol}_{d b}\left(p^{-t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)\right) \cdot \int_{\substack{y \in V_{r}^{3} \\
\frac{1}{2}(y, y)-T \in p^{t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}}} \varphi^{[r]}(y) d y \\
& \stackrel{(5.2)}{=} \gamma(V) \lim _{t \rightarrow \infty} 2^{-(3 / 2) \delta_{2 p}} p^{6 t} d y \text {. } \\
& y \in \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p}\right) \\
& { }^{t} y S_{r} y-T \in p^{t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}
\end{aligned}
$$

Now $\left\{\left.y \in \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p}\right)\right|^{t} y S_{r} y-T \in p^{t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}\right\}$ is a union of

$$
\tilde{A}_{p^{t}}\left(T, S_{r}\right):=\#\left\{\left.y \in \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p} / 2 p^{t} \mathbb{Z}_{p}\right)\right|^{t} y S_{r} y-T \in p^{t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}\right\}
$$

cosets for $2 p^{t} \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p}\right)$. Moreover, by the definition of $d y$ (4.3) we have

$$
\begin{aligned}
\operatorname{vol}_{d y}\left(\mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p}\right)\right) & =\left|\operatorname{det} 2 S_{r}\right|^{3 / 2} \\
& =\left|\operatorname{det} 2 S_{0}\right|^{3 / 2} \\
& =2^{-6 \delta_{2 p}}\left|\operatorname{det} S_{0}\right|^{3 / 2}
\end{aligned}
$$

and hence

$$
\operatorname{vol}_{d y}\left(2 p^{t} \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p}\right)\right)=2^{-6 \delta_{2 p} p}\left|\operatorname{det} S_{0}\right|^{3 / 2} 2^{-3(4+2 r) \delta_{2 p}} p^{-t 3(4+2 r)}
$$

Therefore $W_{T}(e, r, \Phi)$ is equal to

$$
\gamma(V) 2^{-6 \delta_{2 p}}\left|\operatorname{det} S_{0}\right|^{3 / 2} 2^{(-(3 / 2)-3(4+2 r)) \delta_{2 p}} \lim _{t \rightarrow \infty} p^{6 t-t 3(4+2 r)} \tilde{A}_{p^{t}}\left(T, S_{r}\right)
$$

Now we have

$$
\tilde{A}_{p^{t}}\left(T, S_{r}\right)=2^{3(4+2 r) \delta_{2 p}} A_{p^{t}}\left(T, S_{r}\right)
$$

with

$$
A_{p^{t}}\left(T, S_{r}\right)=\#\left\{\left.y \in \mathrm{M}_{2 r+4,3}\left(\mathbb{Z}_{p} / p^{t} \mathbb{Z}_{p}\right)\right|^{t} y S_{r} y-T \in p^{t} \operatorname{Sym}_{3}\left(\mathbb{Z}_{p}\right)^{\vee}\right\}
$$

By definition we have

$$
\alpha_{p}\left(T, S_{r}\right)=\lim _{t \rightarrow \infty} p^{6 t-t 3(4+2 r)} A_{p^{t}}\left(T, S_{r}\right)
$$

and this proves the proposition.
Corollary 5.3. - For $\Phi$ as in Proposition 5.2, $W_{T}(e, 0, \Phi) \neq 0$ if and only if $T$ is represented by $S_{0}$.
5.6. We will now prove Theorem 3.2.

As $\alpha_{p}\left(T, S_{r}\right)$ is a rational function in $X=p^{-r}$, it follows from Proposition 5.2 that

$$
\begin{equation*}
W_{T}^{\prime}(e, 0, \Phi)=-\log (p) 2^{-(15 / 2) \delta_{2 p}}\left|\operatorname{det} S_{0}\right|^{3 / 2} \gamma(V) \frac{\partial}{\partial X} \alpha_{p}\left(T, S_{r}\right)_{\mid X=1} \tag{5.3}
\end{equation*}
$$

Let $D$ be the division quaternion algebra over $\mathbb{Q}_{p}$ and denote by $\mathcal{O}_{D}$ its maximal order. We denote by $S=S_{0}$ (resp. $\tilde{S}=\tilde{S}_{0}$ ) the matrix associated to the quadratic space $V=\mathrm{M}_{2}\left(\mathbb{Z}_{p}\right)$ (resp. $\left.\tilde{V}=\mathcal{O}_{D}\right)$ endowed with the reduced Norm. Then we have (see e.g., [Wd1, (4.5) and (4.6)])

$$
\begin{aligned}
\left|\operatorname{det}\left(S_{0}\right)\right| & =2^{4 \delta_{2 p}}, \\
\left|\operatorname{det}\left(\tilde{S}_{0}\right)\right| & =2^{4 \delta_{2 p}} p^{-2} .
\end{aligned}
$$

Moreover we have by the explicit formulas in the Appendix of [Ku3]

$$
\gamma(V)=-\gamma(\tilde{V})
$$

Using the notation of Theorem 3.2, we therefore have by Proposition 5.2 and (5.3)

$$
\begin{equation*}
\frac{W_{T, p}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T, p}\left(e, 0, \tilde{\Phi}_{p}\right)}(\log p)^{-1}=p^{3} \frac{\frac{\partial}{\partial X} \alpha_{p}\left(T, S_{r}\right)_{\mid X=1}}{\alpha_{p}\left(T, \tilde{S}_{0}\right)} \tag{5.4}
\end{equation*}
$$

But now by [ $\mathbf{W d 2}$, Theorem 1.1 and 2.16] we have

$$
\begin{equation*}
\alpha_{p}\left(T, \tilde{S}_{0}\right)=2(p+1)^{2} p^{-1} \tag{5.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial}{\partial X} \alpha_{p}\left(T, S_{r}\right)_{\mid X=1}=-p^{-4}\left(p^{2}-1\right)^{2} \lg \left(\mathcal{O}_{\mathcal{T}_{T}, \xi}\right) \tag{5.6}
\end{equation*}
$$

Therefore we have

$$
\begin{aligned}
& \lg \left(\mathcal{O}_{T_{T}, \xi}\right) \log (p) \stackrel{(5.5)}{=}-\frac{p^{4}}{\left(p^{2}-1\right)^{2}} \frac{2(p+1)^{2}}{p} \frac{\frac{\partial}{\partial X} \alpha_{p}\left(T, S_{r}\right)_{\mid X=1}}{\alpha_{p}\left(T, \tilde{S}_{0}\right)} \\
& \stackrel{(5.4)}{=}-\frac{2}{(p-1)^{2}} \frac{W_{T}^{\prime}\left(e, 0, \Phi_{p}\right)}{W_{T}\left(e, 0, \tilde{\Phi}_{p}\right)}
\end{aligned}
$$

which proves the theorem.

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[^0]:    ${ }^{(1)}$ Acronym for Arithmetische Geometrie Oberseminar.

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[^8]:    ${ }^{(1)}$ It is more traditional to attribute the conductor $p^{s}$ to this order.

[^9]:    ${ }^{(2)}$ I thank S. Wewers for pointing out this possibility, which I had overlooked.

[^10]:    2000 Mathematics Subject Classification. - 11E08, 11F30, 11F32, 11G18.
    Key words and phrases. - Eisenstein series, local Whittaker functions, Siegel-Weil formula, local representation densities.

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