# On Representations of Integers by Indefinite Ternary Quadratic Forms 

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## CORE

 expectation (the product of local densities) and $0 \leqslant c \leqslant 2$. We give examples of $f$ and $q$ such that $c$ takes the values $0,1,2$. © 2001 Academic PressKey Words: ternary quadratic forms.

## 0. INTRODUCTION

Let $f$ be a nondegenerate indefinite integral-matrix quadratic form of $n$ variables:

$$
f\left(x_{1}, \ldots, x_{n}\right)=\sum_{i, j=1}^{n} a_{i j} x_{i} x_{j}, \quad a_{i j} \in \mathbf{Z}, \quad a_{i j}=a_{j i} .
$$

Let $q \in \mathbf{Z}, q \neq 0$. Let $W=\mathbf{Q}^{n}$. Consider the affine quadric $X$ in $W$ defined by the equation

$$
f\left(x_{1}, \ldots, x_{n}\right)=q .
$$

We wish to count the representations of $q$ by the quadratic form $f$, that is the integer points of $X$.
${ }^{1}$ Partially supported by the Hermann Minkowski Center for Geometry.

Since $f$ is indefinite, the set $X(\mathbf{Z})$ can be infinite. We fix a Euclidean norm $|\cdot|$ on $\mathbf{R}^{n}$. Consider the counting function

$$
N(T, X)=\#\{x \in X(\mathbf{Z}):|x| \leqslant T\}
$$

where $T \in \mathbf{R}, T>0$. We are interested in the asymptotic behavior of $N(T, X)$ as $T \rightarrow \infty$.

When $n \geqslant 4$, the counting function $N(T, X)$ can be approximated by the product of local densities. For a prime $p$ set

$$
\mu_{p}(X)=\lim _{k \rightarrow \infty} \frac{\# X\left(\mathbf{Z} / p^{k} \mathbf{Z}\right)}{\left(p^{k}\right)^{n-1}}
$$

For almost all $p$ it suffices to take $k=1$ :

$$
\mu_{p}(X)=\frac{\# X\left(\mathbf{F}_{p}\right)}{p^{n-1}} .
$$

Set $\mathcal{G}(X)=\prod_{p} \mu_{p}(X)$; this product converges absolutely (for $n \geqslant 4$ ) and is called the singular series. Set

$$
\mu_{\infty}(T, X)=\lim _{\varepsilon \rightarrow 0} \frac{\operatorname{Vol}\left\{x \in \mathbf{R}^{n}:|x| \leqslant T,|f(x)-q|<\varepsilon / 2\right\}}{\varepsilon},
$$

which is called the singular integral. For $n \geqslant 4$ the following asymptotic formula holds:

$$
N(T, X) \sim \Theta(X) \mu_{\infty}(T, X) \quad \text { as } \quad T \rightarrow \infty
$$

This follows from results of $[2,6.4]$ (which are based on analytical results of $[6,7,8]$ ). For certain non-Euclidean norms the similar result was earlier proved by the Hardy-Littlewood circle method, cf. [5] in the case $n \geqslant 5$ and [9] in the more difficult case $n=4$.

We are interested here in the case $n=3$, a ternary quadratic form. This case is beyond the range of the Hardy-Littlewood circle method. Set $D=$ $\operatorname{det}\left(a_{i j}\right)$. We assume that $-q D$ is not a square. Then the product $\mathfrak{S}(X)=$ $\Pi \mu_{p}(X)$ conditionally converges (see Sect. 1 below), but in general $N(T, X)$ is not asymptotically $\left.\mathbb{S}_{( } X\right) \mu_{\infty}(T, X)$. From results of [2] it follows that

$$
N(T, X) \sim c_{X} \circlearrowleft(X) \mu_{\infty}(T, X) \quad \text { as } \quad T \rightarrow \infty
$$

with $0 \leqslant c_{X} \leqslant 2$, see details in Section 1.5 below. We wish to know what values can $c_{X}$ take.

A case when $c_{X}=0$ was already known to Siegel, see also [2, 6.4.1]. Consider the quadratic form

$$
f_{1}\left(x_{1}, x_{2}, x_{3}\right)=-9 x_{1}^{2}+2 x_{1} x_{2}+7 x_{2}^{2}+2 x_{3}^{2},
$$

and take $q=1$. Let $X$ be defined by $f_{1}(x)=q$. Then $f_{1}$ does not represent 1 over $\mathbf{Z}$, so $N(T, X)=0$ for all $T$. On the other hand, $f_{1}$ represents 1 over $\mathbf{R}$ and over $\mathbf{Z}_{p}$ for all $p$, and $\Theta(X) \mu_{\infty}(T, X) \rightarrow \infty$ as $T \rightarrow \infty$. Thus $c_{X}=0$ (see details in Sect. 2).

We show that $c_{X}$ can take the value 2 . Recall that two integral quadratic forms $f, f^{\prime}$ are in the same genus if they are equivalent over $\mathbf{R}$ and over $\mathbf{Z}_{p}$ for every prime $p$, cf. e.g. [3].

Theorem 0.1. Let $f$ be an indefinite integral-matrix ternary quadratic form, $q \in \mathbf{Z}, q \neq 0$, and let $X$ be the affine quadric defined by the equation $f(x)=q$. Assume that $f$ represents $q$ over $\mathbf{Z}$ and that there exists a quadratic form $f^{\prime}$ in the genus of $f$ such that $f^{\prime}$ does not represent $q$ over $\mathbf{Z}$. Then $c_{X}=2$ :

$$
N(T, X) \sim 2 \mathbb{G}(X) \mu_{\infty}(T, X) \quad \text { as } \quad T \rightarrow \infty
$$

Theorem 0.1 will be proved in Section 3.

Example 0.1.1. Let $f_{2}\left(x_{1}, x_{2}, x_{3}\right)=-x_{1}^{2}+64 x_{2}^{2}+2 x_{3}^{2}, q=1$. Then $f_{2}$ represents $1\left(f_{2}(1,0,1)=1\right)$ and the quadratic form $f_{1}$ considered above is in the genus of $f_{2}$ (cf. $[4,15.6]$ ). The form $f_{1}$ does not represent 1 . Take $|x|=\left(x_{1}^{2}+64 x_{2}^{2}+2 x_{3}^{2}\right)^{1 / 2}$. By Theorem $0.1 \quad c_{X}=2$ for the variety $X: f_{2}(x)$ $=1$. Analytic and numeric calculations give $2 \mathbb{S}(X) \mu_{\infty}(T, X) \sim 0.794 T$. On the other hand, numeric calculations give for $T=10,000$ the value $N(T, X) / T=0.8024$.

We also show that $c_{X}$ can take the value 1 .

Theorem 0.2. Let $f$ be an indefinite integral-matrix ternary quadratic form, $q \in \mathbf{Z}, q \neq 0$, and let $X$ be the affine quadric defined by the equation $f(x)=q$. Assume that $X(\mathbf{R})$ is two-sheeted (has two connected components). Then $c_{X}=1$ :

$$
N(T, X) \sim \subseteq(X) \mu_{\infty}(T, X) \quad \text { as } \quad T \rightarrow \infty .
$$

Theorem 0.2 will be proved in Section 4.

Example 0.2 .1 . Let $f_{2}$ and $|x|$ be as in Example 0.1.1, $q=-1, X: f_{2}(x)$ $=q$. Then $X(\mathbf{R})$ has two connected components, and by Theorem 0.2 $c_{X}=1$. Analytic and numeric calculations give $\mathfrak{S}(X) \mu_{\infty}(T, X) \sim 0.7065 T$. On the other hand, numeric calculations give for $T=10,000$ the value $N(T, X) / T=0.7048$.

Question 0.3. Can $c_{X}$ take values other than $0,1,2$ ?
The plan of the paper is the following. In Section 1 we describe results of [2] in the case of 2-dimensional affine quadrics. In Section 2 we treat in detail the example of $c_{X}=0$. In Section 3 we prove Theorem 0.1. In Section 4 we prove Theorem 0.2.

## 1. RESULTS OF [2] IN THE CASE OF TERNARY QUADRATIC FORMS

Let $f$ be an indefinite ternary integral-matrix quadratic form

$$
f\left(x_{1}, x_{2}, x_{3}\right)=\sum_{i, j=1}^{3} a_{i j} x_{i} x_{j}, \quad a_{i j} \in \mathbf{Z}, \quad a_{i j}=a_{j i}
$$

Let $q \in \mathbf{Z}, q \neq 0$. Let $D=\operatorname{det}\left(a_{i j}\right)$. We assume that $-q D$ is not a square.
Let $W=\mathbf{Q}^{3}$ and let $X$ denote the affine variety in $W$ defined by the equation $f(x)=q$, where $x=\left(x_{1}, x_{2}, x_{3}\right)$. We assume that $X$ has a $\mathbf{Q}$-point $x^{0}$. Set $G=\operatorname{Spin}(W, f)$, the spinor group of $f$. Then $G$ acts on $W$ on the left, and $X$ is an orbit (a homogeneous space) of $G$.

### 1.1. Rational Points in Adelic Orbits

Let $\mathbf{A}$ denote the adèle ring of $\mathbf{Q}$. The group $G(\mathbf{A})$ acts on $X(\mathbf{A})$; let $\mathcal{O}_{\mathbf{A}}$ be an orbit. We would like to know whether $\mathcal{O}_{\mathbf{A}}$ has a $\mathbf{Q}$-rational point.

Let $W^{\prime}$ denote the orthogonal complement of $x^{0}$ in $W$, and let $f^{\prime}$ denote the restriction of $f$ to $W^{\prime}$. Let $H$ be the stabilizer of $x^{0}$ in $G$, then $H=\operatorname{Spin}\left(W^{\prime}, f^{\prime}\right)$. Since $\operatorname{dim} W^{\prime}=2$, the group $H$ is a one-dimensional torus.

We have $\operatorname{det} f^{\prime}=D / q$, so up to multiplication by a square $\operatorname{det} f^{\prime}=q D$. It follows that up to multiplication by a scalar, $f^{\prime}$ is equivalent to the quadratic form $u^{2}+q D v^{2}$. Set $K=\mathbf{Q}(\sqrt{-q D})$, then $K$ is a quadratic extension of $\mathbf{Q}$, because $-q D$ is not a square. The torus $H$ is anisotropic over $\mathbf{Q}$ (because $-q D$ is not a square), and $H$ splits over $K$. Let $\mathbf{X}_{*}\left(H_{K}\right)$ denote the cocharacter group of $H_{K}, \mathbf{X}_{*}\left(H_{K}\right)=\operatorname{Hom}\left(\mathbb{G}_{m, K}, H_{K}\right)$; then $\mathbf{X}_{*}\left(H_{K}\right) \simeq \mathbf{Z}$. The non-neutral element of $\operatorname{Gal}(K / \mathbf{Q})$ acts on $\mathbf{X}_{*}\left(H_{K}\right)$ by multiplication by -1 .

Let $\mathcal{O}_{\mathbf{A}}$ be an orbit of $G(\mathbf{A})$ in $X(\mathbf{A}), \mathcal{O}_{\mathbf{A}}=\Pi \mathcal{O}_{v}$ where $\mathcal{O}_{v}$ is an orbit of $G\left(\mathbf{Q}_{v}\right)$ in $X\left(\mathbf{Q}_{v}\right), v$ runs over the places of $\mathbf{Q}$, and $\mathbf{Q}_{v}$ denotes the completion of $\mathbf{Q}$ at $v$. We define local invariants $v_{v}\left(\mathcal{O}_{v}\right)= \pm 1$. If $\mathcal{O}_{v}=G\left(\mathbf{Q}_{v}\right) \cdot x^{0}$, then we set $v_{v}\left(\mathcal{O}_{v}\right)=+1$, if not, we set $v_{v}\left(\mathcal{O}_{v}\right)=-1$. Then $v_{v}\left(\mathcal{O}_{v}\right)=+1$ for almost all $v$. We define $v\left(\mathcal{O}_{\mathbf{A}}\right)=\Pi v_{v}\left(\mathcal{O}_{v}\right)$ where $\mathcal{O}_{\mathbf{A}}=\Pi \mathcal{O}_{v}$. Note that the local invariants $v_{v}\left(\mathcal{O}_{v}\right)$ depend on the choice of the rational point $x^{0} \in X(\mathbf{Q})$; one can prove, however, that their product $v\left(\mathcal{O}_{\mathbf{A}}\right)$ does not depend on $x^{0}$.

Let $x \in X(\mathbf{A})$. We set $v(x)=v(G(\mathbf{A}) \cdot x)$. Then $v(x)$ takes values $\pm 1$; it is a locally constant function on $X(\mathbf{A})$, because the orbits of $G(\mathbf{A})$ are open in $X(\mathbf{A})$.

For $x \in X(\mathbf{A})$ define $\delta(x)=v(x)+1$. In other words, if $v(x)=-1$ then $\delta(x)=0$, and if $v(x)=+1$ then $\delta(x)=2$. Then $\delta$ is a locally constant function on $X(\mathbf{A})$.

Theorem 1.1. An orbit $\mathcal{O}_{\mathbf{A}}$ of $G(\mathbf{A})$ in $X(\mathbf{A})$ has a $\mathbf{Q}$-rational point if and only if $v\left(\mathcal{O}_{\mathbf{A}}\right)=+1$.

Below we will deduce Theorem 1.1 from [2, Theorem 3.6].

### 1.2. Proof of Theorem 1.1

For a torus $T$ over a field $k$ of characteristic 0 we define a finite abelian group $C(T)$ as follows

$$
C(T)=\left(\mathbf{X}_{*}\left(T_{k}^{-}\right)_{\text {Gal }(\bar{k} / k)}\right)_{\text {tors }}
$$

where $\bar{k}$ is a fixed algebraic closure of $k, \mathbf{X}_{*}\left(T_{k}^{-}\right)_{\operatorname{Gal}(\bar{k} / k)}$ denotes the group of coinvariants, and $(\cdot)_{\text {tors }}$ denotes the torsion subgroup. If $k$ is a number field and $k_{v}$ is the completion of $k$ at a place $v$, then we define $C_{v}(T)=$ $C\left(T_{k_{v}}\right)$. There is a canonical map $i_{v}: C_{v}(T) \rightarrow C(T)$ induced by an inclusion $\operatorname{Gal}\left(\bar{k}_{v} / k_{v}\right) \rightarrow \operatorname{Gal}(\bar{k} / k)$. These definitions were given for connected reductive groups (not only for tori) by Kottwitz [10]; see also [2, 3.4]. Kottwitz writes $A(T)$ instead of $C(T)$.

We compute $C(H)$ for our one-dimensional torus $H$ over $\mathbf{Q}$. Clearly

$$
C(H)=\left(\mathbf{X}_{*}\left(H_{K}\right)_{\text {Gal }(K / \mathbf{Q})}\right)_{\mathrm{tors}}=\mathbf{Z} / 2 \mathbf{Z} .
$$

We have $C_{v}(H)=1$ if $K \otimes \mathbf{Q}_{v}$ splits, and $C_{v}(H) \simeq \mathbf{Z} / 2 \mathbf{Z}$ if $K \otimes \mathbf{Q}_{v}$ is a field. The map $i_{v}$ is injective for any $v$.

We now define the local invariants $\kappa_{v}\left(\mathcal{O}_{v}\right)$ as in [2], where $\mathcal{O}_{v}$ is an orbit of $G\left(\mathbf{Q}_{v}\right)$ in $X\left(\mathbf{Q}_{v}\right)$. The set of orbits of $G\left(\mathbf{Q}_{v}\right)$ in $X\left(\mathbf{Q}_{v}\right)$ is in canonical bijection with $\operatorname{ker}\left[H^{1}\left(\mathbf{Q}_{v}, H\right) \rightarrow H^{1}\left(\mathbf{Q}_{v}, G\right)\right]$, cf. [13, I-5.4, Corollary 1 of Proposition 36]. Hence $\mathcal{O}_{v}$ defines a cohomology class $\xi_{v} \in H^{1}\left(\mathbf{Q}_{v}, H\right)$. The local Tate-Nakayama duality for tori defines a canonical homomorphism $\beta_{v}: H^{1}\left(\mathbf{Q}_{v}, H\right) \rightarrow C_{v}(H)$, see Kottwitz [10, Theorem 1.2]. (Kottwitz defines
the map $\beta_{v}$ in a more general setting, when $H$ is any connected reductive group over a number field.) The homomorphism $\beta_{v}$ is an isomorphism for any $v$. We set $\kappa_{v}\left(\mathcal{O}_{v}\right)=\beta_{v}\left(\xi_{v}\right)$. Note that if $\mathcal{O}_{v}=G\left(\mathbf{Q}_{v}\right) \cdot x^{0}$, then $\xi_{v}=0$ and $\kappa_{v}\left(\mathcal{O}_{v}\right)=0$; if $\mathcal{O}_{v} \neq G\left(\mathbf{Q}_{v}\right) \cdot x^{0}$, then $\xi_{v} \neq 0$ and $\kappa_{v}\left(\mathcal{O}_{v}\right)=1$.

We define the Kottwitz invariant $\kappa\left(\mathcal{O}_{\mathbf{A}}\right)$ of an orbit $\mathcal{O}_{\mathbf{A}}=\prod \mathcal{O}_{v}$ of $G(\mathbf{A})$ in $X(\mathbf{A})$ by $\kappa\left(\mathcal{O}_{\mathbf{A}}\right)=\sum_{v} i_{v}\left(\kappa_{v}\left(\mathcal{O}_{v}\right)\right)$. We identify $C(H)$ with $\mathbf{Z} / 2 \mathbf{Z}$, and $C_{v}(H)$ with a subgroup of $\mathbf{Z} / 2 \mathbf{Z}$. With this identifications $\kappa\left(\mathcal{O}_{\mathbf{A}}\right)=\sum \kappa_{v}\left(\mathcal{O}_{v}\right)$.

We prefer the multiplicative rather than additive notation. Instead of $\mathbf{Z} / 2 \mathbf{Z}$ we consider the group $\{+1,-1\}$, and set

$$
v_{v}\left(\mathcal{O}_{v}\right)=(-1)^{\kappa_{v}\left(\mathcal{O}_{v}\right)}, \quad v\left(\mathcal{O}_{\mathbf{A}}\right)=(-1)^{\kappa\left(\mathcal{O}_{\mathbf{A}}\right)} .
$$

Here $v_{v}\left(\mathcal{O}_{v}\right)$ and $v\left(\mathcal{O}_{\mathbf{A}}\right)$ take the values $\pm 1$. We have $v\left(\mathcal{O}_{\mathbf{A}}\right)=\Pi v_{v}\left(\mathcal{O}_{v}\right)$. Since $\kappa_{v}\left(\mathcal{O}_{v}\right)=0$ if and only if $\mathcal{O}_{v}=G\left(\mathbf{Q}_{v}\right) \cdot x^{0}$, we see that $v_{v}\left(\mathcal{O}_{v}\right)=+1$ if and only if $\mathcal{O}_{v}=G\left(\mathbf{Q}_{v}\right) \cdot x^{0}$. Hence our $v_{v}\left(\mathcal{O}_{v}\right)$ and $v\left(\mathcal{O}_{\mathbf{A}}\right)$ coincide with $v_{v}\left(\mathcal{O}_{v}\right)$ and $v\left(\mathcal{O}_{\mathbf{A}}\right)$, resp., introduced in Section 1.1.

By Theorem 3.6 of [2] an adelic orbit $\mathcal{O}_{\mathbf{A}}$ contains $\mathbf{Q}$-rational points if and only if $\kappa\left(\mathcal{O}_{\mathbf{A}}\right)=0$. With our multiplicative notation $\kappa\left(\mathcal{O}_{\mathbf{A}}\right)=0$ if and only if $v\left(\mathcal{O}_{\mathbf{A}}\right)=+1$. Thus $\mathcal{O}_{\mathbf{A}}$ contains $\mathbf{Q}$-points if and only if $v\left(\mathcal{O}_{\mathbf{A}}\right)=+1$. We have deduced Theorem 1.1 from [2, Theorem 3.6].

### 1.3. Tamagawa Measure

We define a gauge form on $X$, i.e. a regular differential form $\omega \in \Lambda^{2}(X)$ without zeroes. Recall that $X$ is defined by the equation $f(x)=q$. Choose a differential form $\mu$ of degree 2 on $W$ such that $\mu \wedge d f=d x_{1} \wedge d x_{2} \wedge d x_{3}$, where $x_{1}, x_{2}, x_{3}$ are the coordinates in $W=\mathbf{Q}^{3}$. Let $\omega=\left.\mu\right|_{X}$, the restriction of $\mu$ to $X$. Then $\omega$ is a gauge form on $X$, cf. [2, 1.3], and it does not depend on the choice of $\mu$. The gauge form $\omega$ is $G$-invariant, because there exists a $G$-invariant gauge form on $X$, cf. $[2,1.4]$, and a gauge form on $X$ is unique up to a scalar multiple, cf. [2, Corollary 1.5.4].

For any place $v$ of $\mathbf{Q}$ one associates with $\omega$ a local measure $m_{v}$ on $X\left(\mathbf{Q}_{v}\right)$, cf. $[14,2.2]$. We show how to define a Tamagawa measure on $X(\mathbf{A})$, following [2, 1.6.2].

We have by $[2,1.8 .1], \mu_{p}(X)=m_{p}\left(X\left(\mathbf{Z}_{p}\right)\right)$, where $\mu_{p}(X)$ is defined in the Introduction. By [14, Theorem 2.2.5], for almost all $p$ we have $m_{p}\left(X\left(\mathbf{Z}_{p}\right)\right)$ $=\# X\left(\mathbf{F}_{p}\right)$.

We compute $\# X\left(\mathbf{F}_{p}\right)$. The group $\operatorname{SO}(f)\left(\mathbf{F}_{p}\right)$ acts on $X\left(\mathbf{F}_{p}\right)$ with stabilizer $\operatorname{SO}\left(f^{\prime}\right)\left(\mathbf{F}_{p}\right)$, where $\operatorname{SO}\left(f^{\prime}\right)\left(\mathbf{F}_{p}\right)$ is defined for almost all $p$. This action is transitive by Witt's theorem. Thus we obtain that $\# X\left(\mathbf{F}_{p}\right)=$ $\# \operatorname{SO}(f)\left(\mathbf{F}_{p}\right) / \# \operatorname{SO}\left(f^{\prime}\right)\left(\mathbf{F}_{p}\right)$. By [1, III-6],

$$
\# \mathrm{SO}(f)\left(\mathbf{F}_{p}\right)=p\left(p^{2}-1\right), \quad \# \mathrm{SO}\left(f^{\prime}\right)\left(\mathbf{F}_{p}\right)=p-\chi(p)
$$

where $\chi(p)=-1$ if $f^{\prime} \bmod p$ does not represent 0 , and $\chi(p)=+1$ if $f^{\prime} \bmod p$ represents 0 . We have $\chi(p)=\left(\frac{-q D}{p}\right)$. We obtain for $p \nmid q D$

$$
\# X\left(\mathbf{F}_{p}\right)=\frac{p\left(p^{2}-1\right)}{p-\chi(p)}, \quad \mu_{p}(X)=\frac{\# X\left(\mathbf{F}_{p}\right)}{p^{2}}=\frac{1-1 / p^{2}}{1-\chi(p) / p} .
$$

For $p \mid q D$ set $\chi(p)=0$. We define

$$
L_{p}(s, \chi)=\left(1-\chi(p) p^{-s}\right)^{-1}, \quad L(s, \chi)=\prod_{p} L_{p}(s, \chi),
$$

where $s$ is a complex variable. We set

$$
\lambda_{p}=L_{p}(1, \chi)^{-1}=1-\frac{\chi(p)}{p}, \quad r=L(1, \chi)^{-1} .
$$

Then the product $\prod_{p}\left(\lambda_{p}^{-1} \mu_{p}\right)$ converges absolutely, hence the family $\left(\lambda_{p}\right)$ is a family of convergence factors in the sense of [14, 2.3]. We define, as in $[2,1.6 .2]$, the measures

$$
m_{f}=r^{-1} \prod_{p}\left(\lambda_{p}^{-1} m_{p}\right), \quad m=m_{\infty} m_{f},
$$

then $m_{f}$ is a measure on $X\left(\mathbf{A}_{f}\right)$ (where $\mathbf{A}_{f}$ is the ring of finite adèles) and $m$ is a measure on $X(\mathbf{A})$. We call $m$ the Tamagawa measure on $X(\mathbf{A})$.

### 1.4. Counting Integer Points

For $T>0$ set $X(\mathbf{R})^{T}=\{x \in X(\mathbf{R}):|x| \leqslant T\}$.
Theorem 1.2.

$$
N(T, X) \sim \int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m .
$$

In other words,

$$
\begin{equation*}
N(T, X) \sim 2 m\left(\left\{x \in X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}}): v(x)=+1\right\}\right) . \tag{1}
\end{equation*}
$$

Theorem 1.2 follows from [2, Theorem 5.3] (cf. [2, 6.4] and [2, Definition 2.3]).

For comparison note that

$$
\begin{equation*}
m\left(X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})\right)=m_{\infty}\left(X(\mathbf{R})^{T}\right) m_{f}(X(\hat{\mathbf{Z}}))=\mu_{\infty}(T, X) \mathbb{S}_{(X)} \tag{2}
\end{equation*}
$$

cf. $[2,1.8]$.
The following lemma will be used in the proof of Theorem 0.1.

Lemma 1.3. Assume that there exists $y \in X(\mathbf{R} \times \hat{\mathbf{Z}})$ such that $v(y)=+1$. Then the set $X(\mathbf{Z})$ is infinite.

Proof. Since $v$ is a locally constant function on $X(\mathbf{A})$, there exists a nonempty open subset $\mathscr{U}_{f} \in X(\hat{\mathbf{Z}})$ and an orbit $\mathscr{U}_{\infty}$ of $G(\mathbf{R})$ in $X(\mathbf{R})$ such that $v(x)=+1$ for all $x \in U_{\infty} \times \mathscr{U}_{f}$. Set $\mathscr{U}_{\infty}^{T}=\left\{x \in \mathscr{U}_{\infty}:|x| \leqslant T\right\}$, then $m_{\infty}\left(\mathscr{U}_{\infty}^{T}\right) \rightarrow \infty$ as $T \rightarrow \infty$. We have

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m \geqslant \int_{\mathscr{U}_{\infty}^{T} \times \mathscr{U}_{f}} \delta(x) d m=2 m_{\infty}\left(\mathscr{U}_{\infty}^{T}\right) m_{f}\left(\mathscr{U}_{f}\right) .
$$

Since $2 m_{\infty}\left(\mathscr{U}_{\infty}^{T}\right) m_{f}\left(\mathscr{U}_{f}\right) \rightarrow \infty$ as $T \rightarrow \infty$, we see that

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m \rightarrow \infty \quad \text { as } \quad T \rightarrow \infty
$$

and by Theorem $1.2 N(T, X) \rightarrow \infty$. Hence $X(\mathbf{Z})$ is infinite.

### 1.5. The Constant $c_{X}$

Here we prove the following result:
Proposition 1.4.

$$
N(T, X) \sim c_{X} \Xi(X) \mu_{\infty}(T, X) \quad \text { as } \quad T \rightarrow \infty
$$

with some constant $c_{X}, 0 \leqslant c_{X} \leqslant 2$.
Proof. If $X(\mathbf{R})$ has two connected components, then by Theorem 0.2 (which we will prove in Section 4 below), $N(T, X) \sim \subseteq(X) \mu_{\infty}(T, X)$, so the proposition holds with $c_{X}=1$.

If $X(\mathbf{R})$ has one connected component, then $X(\mathbf{R})$ consists of one $G(\mathbf{R})$-orbit and $v_{\infty}(X(\mathbf{R}))=+1$. For an orbit $\mathcal{O}_{f}=\prod \mathcal{O}_{p}$ of $G\left(\mathbf{A}_{f}\right)$ in $X\left(\mathbf{A}_{f}\right)$ we set $v_{f}\left(\mathcal{O}_{f}\right)=\prod_{p} v_{p}\left(\mathcal{O}_{p}\right)$. We regard $v_{f}$ as a locally constant function on $X\left(\mathbf{A}_{f}\right)$ taking the values $\pm 1$. Define $X(\hat{\mathbf{Z}})_{+}=\left\{x_{f} \in X(\hat{\mathbf{Z}}): v_{f}\left(x_{f}\right)=+1\right\}$. We have

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m=2 m_{\infty}\left(X(\mathbf{R})^{T}\right) m_{f}\left(X(\hat{\mathbf{Z}})_{+}\right) .
$$

Set $c_{X}=2 m_{f}\left(X(\hat{\mathbf{Z}})_{+}\right) / m_{f}(X(\hat{\mathbf{Z}}))$, then $0 \leqslant c_{X} \leqslant 2$ and

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m=c_{X} m_{\infty}\left(X(\mathbf{R})^{T}\right) m_{f}(X(\hat{\mathbf{Z}}))=c_{X} \mu_{\infty}(T, X) \Theta(X) .
$$

Using Theorem 1.2, we see that

$$
N(T, X) \sim c_{X} \mu_{\infty}(T, X) \mathbb{S}_{(X)} \quad \text { as } \quad T \rightarrow \infty .
$$

## 2. AN EXAMPLE OF $c_{X}=0$

Let

$$
f_{1}\left(x_{1}, x_{2}, x_{3}\right)=-9 x_{1}^{2}+2 x_{1} x_{2}+7 x_{2}^{2}+2 x_{3}^{2}, \quad q=1 .
$$

This example was mentioned in [2, 6.4.1]. Here we provide a detailed exposition.

Consider the variety $X$ defined by the equation $f_{1}(x)=q$. We have $f_{1}\left(-\frac{1}{2}, \frac{1}{2}, 1\right)=1$. It follows that $f_{1}$ represents 1 over $\mathbf{R}$ and over $\mathbf{Z}_{p}$ for $p>2$.

We have $f_{1}(4,1,1)=-127 \equiv 1\left(\bmod 2^{7}\right)$. We prove that $f_{1}$ represents 1 over $\mathbf{Z}_{2}$. Define a polynomial of one variable $F(Y)=f_{1}(4,1, Y)-1$, $F \in \mathbf{Z}_{2}[Y]$. Then $F(1)=-2^{7},|F(1)|_{2}=2^{-7}, F^{\prime}(Y)=4 Y,\left|F^{\prime}(1)^{2}\right|_{2}=2^{-4}$, $|F(1)|_{2}<\left|F^{\prime}(1)^{2}\right|_{2}$. By Hensel's lemma (cf. [11, II-§2, Proposition 2]) $F$ has a root in $\mathbf{Z}_{2}$. Thus $f_{1}$ represents 1 over $\mathbf{Z}_{2}$.

Now we prove that $f_{1}$ does not represent 1 over $\mathbf{Z}$. I know the following elementary proof from D. Zagier.

We prove the assertion by contradiction. Assume on the contrary that

$$
-9 x_{1}^{2}+2 x_{1} x_{2}+7 x_{2}^{2}+2 x_{3}^{2}=1 \quad \text { for some } \quad x_{1}, x_{2}, x_{3} \in \mathbf{Z}
$$

We may write this equation as follows:

$$
2 x_{3}^{2}-1=\left(x_{1}-x_{2}\right)^{2}+8\left(x_{1}-x_{2}\right)\left(x_{1}+x_{2}\right) .
$$

The left hand side is odd, hence $x_{1}-x_{2}$ is odd and therefore $x_{1}+x_{2}$ is odd. We have $\left(x_{1}-x_{2}\right)^{2} \equiv 1(\bmod 8)$. Hence the right hand side is congruent to $1(\bmod 8)$. We see that $x_{3}$ is odd, hence $2 x_{3}^{2}-1 \equiv 1(\bmod 16)$. But

$$
8\left(x_{1}-x_{2}\right)\left(x_{1}+x_{2}\right) \equiv 8 \quad(\bmod 16)
$$

It follows that

$$
\begin{aligned}
\left(x_{1}-x_{2}\right)^{2} & \equiv 9 \\
x_{1}-x_{2} & \equiv \pm 3 \quad(\bmod 16) \\
& (\bmod 8) .
\end{aligned}
$$

Therefore $x_{1}-x_{2}$ must have a prime factor $p \equiv \pm 3(\bmod 8)$. Hence $2 x_{3}^{2}-1$ has a prime factor $p \equiv \pm 3(\bmod 8)$. On the other hand, if $p \mid\left(2 x_{3}^{2}-1\right)$, then

$$
2 x_{3}^{2} \equiv 1 \quad(\bmod p)
$$

and 2 is a square modulo $p,\left(\frac{2}{p}\right)=1$. By the quadratic reciprocity law $p \equiv \pm 1(\bmod 8)$. Contradiction. We have proved that $f_{1}$ does not represent 1 over $\mathbf{Z}$, hence $N(T, X)=0$ for all $T$.

On the other hand,

$$
\mathfrak{\Im}(X) \mu_{\infty}(T, X)=m_{f}(X(\hat{\mathbf{Z}})) m_{\infty}\left(X(\mathbf{R})^{T}\right) .
$$

Since $X(\hat{\mathbf{Z}})$ is a nonempty open subset in $X\left(\mathbf{A}_{f}\right), m_{f}(X(\hat{\mathbf{Z}}))>0$. Now $m_{\infty}\left(X(\mathbf{R})^{T}\right) \rightarrow \infty$ as $T \rightarrow \infty$. Hence $\Theta(X) \mu_{\infty}(T, X) \rightarrow \infty$ as $T \rightarrow \infty$, and thus $c_{X}=0$.

## 3. PROOF OF THEOREM 0.1

Lemma 3.1. Let $k$ be a field of characteristic different from 2, and let $V$ be a finite-dimensional vector space over $k$. Let $f$ be a non-degenerate quadratic form on $V$. Let $u \in \operatorname{GL}(V)(k), f^{\prime}=u^{*} f$. Then the map $y \mapsto u y: V$ $\rightarrow V$ takes the orbits of $\operatorname{Spin}(f)(k)$ in $V$ to the orbits of $\operatorname{Spin}\left(f^{\prime}\right)(k)$.

Proof. Let $x \in V, f(x) \neq 0$. The reflection (symmetry) $r_{x}=r_{f, x}: V \rightarrow V$ is defined by

$$
r_{x}(y)=y-\frac{2 B(x, y)}{f(x)} x, \quad y \in V,
$$

where $B$ is the symmetric bilinear form on $V$ associated with $f$. Every $s \in \operatorname{SO}(f)(k)$ can be written as

$$
\begin{equation*}
s=r_{x_{1}} \cdots r_{x_{l}} \tag{3}
\end{equation*}
$$

cf. [12, Theorem 43:3]. The spinor norm $\theta(s)$ of $s$ is defined by

$$
\theta(s)=f\left(x_{1}\right) \cdots f\left(x_{l}\right) \quad\left(\bmod k^{* 2}\right) \in k^{*} / k^{* 2}
$$

and it does not depend on the choice of the representation given by (3), cf. [12, §55]. Let $\Theta(f)$ denote the image of $\operatorname{Spin}(f)(k)$ in $\operatorname{SO}(f)(k)$. Then $s \in \mathrm{SO}(f)(k)$ is contained in $\Theta(f)$ if and only if $\theta(s)=1$, cf. [13, III-3.2] or [3, Chap. 10, Theorem 3.3].

Now let $u, f^{\prime}$ be as above. Then $r_{f^{\prime}, u x}=u r_{f, x} u^{-1}, f^{\prime}(u x)=f(x)$, and so $\theta_{f^{\prime}}\left(u s u^{-1}\right)=\theta_{f}(s)$. We conclude that $u \Theta(f) u^{-1}=\Theta\left(f^{\prime}\right)$ and that the map $y \mapsto u y$ takes the orbits of $\Theta(f)$ in $V$ to the orbits of $\Theta\left(f^{\prime}\right)$.

Let $f, f^{\prime}$ be integral-matrix quadratic forms on $\mathbf{Z}^{n}$ and assume that $f^{\prime}$ is in the genus of $f$. Then there exists $u \in \mathrm{GL}_{n}(\mathbf{R} \times \hat{\mathbf{Z}})$ such that $f^{\prime}(x)=$ $f\left(u^{-1} x\right)$ for $x \in \mathbf{A}^{n}$. Let $q \in \mathbf{Z}, q \neq 0$. Let $X$ denote the affine quadric $f(x)=q$, and $X^{\prime}$ denote the quadric $f^{\prime}(x)=q$.

Lemma 3.2. The map $x \mapsto u x: \mathbf{A}^{n} \rightarrow \mathbf{A}^{n}$ takes $X(\mathbf{R} \times \hat{\mathbf{Z}})$ to $X^{\prime}(\mathbf{R} \times \hat{\mathbf{Z}})$ and takes orbits of $\operatorname{Spin}(f)(\mathbf{A})$ in $X(\mathbf{A})$ to orbits of $\operatorname{Spin}\left(f^{\prime}\right)(\mathbf{A})$ in $X^{\prime}(\mathbf{A})$.

Proof. Let $A$ denote the matrix of $f$, and $A^{\prime}$ denote the matrix of $f^{\prime}$. We have

$$
\left(u^{-1}\right)^{t} A u^{-1}=A^{\prime}, \quad A=u^{t} A^{\prime} u .
$$

The variety $X$ is defined by the equation $x^{t} A x=q$, and $X^{\prime}$ is defined by $x^{t} A^{\prime} x=q$. One can easily check that the map $x \mapsto u x$ takes $X(\mathbf{R} \times \hat{\mathbf{Z}})$ to $X^{\prime}(\mathbf{R} \times \hat{\mathbf{Z}})$ and $X(\mathbf{A})$ to $X^{\prime}(\mathbf{A})$.

In order to prove that the map $x \mapsto u x: X(\mathbf{A}) \rightarrow X^{\prime}(\mathbf{A})$ takes the orbits of $\operatorname{Spin}(f)(\mathbf{A})$ to the orbits of $\operatorname{Spin}\left(f^{\prime}\right)(\mathbf{A})$, it suffices to prove that the map $x \mapsto u_{v} x: X\left(\mathbf{Q}_{v}\right) \rightarrow X^{\prime}\left(\mathbf{Q}_{v}\right)$ takes the orbits of $\operatorname{Spin}(f)\left(\mathbf{Q}_{v}\right)$ to the orbits of $\operatorname{Spin}\left(f^{\prime}\right)\left(\mathbf{Q}_{v}\right)$ for every $v$, where $u_{v}$ is the $v$-component of $u$. This last assertion follows from Lemma 3.1.

Proposition 3.3. Let $f^{\prime}$ and $q$ be as in Theorem 0.1, in particular $f^{\prime}$ represents $q$ over $\mathbf{Z}_{v}$ for any $v$ (we set $\mathbf{Z}_{\infty}=\mathbf{R}$ ), but not over $\mathbf{Z}$. Let $X^{\prime}$ be the quadric defined by $f^{\prime}(x)=q v$. Then $X^{\prime}(\mathbf{R} \times \hat{\mathbf{Z}})$ is contained in one orbit of $\operatorname{Spin}\left(f^{\prime}\right)(\mathbf{A})$.

Proof. Set $G^{\prime}=\operatorname{Spin}\left(f^{\prime}\right)$. We prove that $X^{\prime}\left(\mathbf{Z}_{v}\right)$ is contained in one orbit of $G^{\prime}\left(\mathbf{Q}_{v}\right)$ for every $v$ by contradiction. Assume on the contrary that for some $v$ the set $X^{\prime}\left(\mathbf{Z}_{v}\right)$ has nontrivial intersection with two orbits of $G^{\prime}\left(\mathbf{Q}_{v}\right)$. Then $v_{v}$ takes both values +1 and -1 on $X^{\prime}\left(\mathbf{Z}_{v}\right)$. It follows that $v$ takes both values +1 and -1 on $X^{\prime}(\mathbf{R} \times \hat{\mathbf{Z}})$. Hence by Lemma $1.3 X^{\prime}$ has infinitely many $\mathbf{Z}$-points. This contradicts to the assumption that $f^{\prime}$ does not represent $q$ over $\mathbf{Z}$.

Proof of Theorem 0.1. Let $u \in \mathrm{GL}_{3}(\mathbf{R} \times \hat{\mathbf{Z}})$ be such that $f^{\prime}(x)=f\left(u^{-1} x\right)$. Let $X, X^{\prime}$ be as above, in particular $X^{\prime}$ has no $\mathbf{Z}$-points. By Proposition 3.3 $X^{\prime}(\mathbf{R} \times \hat{\mathbf{Z}})$ is contained in one orbit of $\operatorname{Spin}\left(f^{\prime}\right)(\mathbf{A})$. It follows from Lemma 3.2 that $X(\mathbf{R} \times \hat{\mathbf{Z}})$ is contained in one orbit of $\operatorname{Spin}(f)(\mathbf{A})$. Since $f$ represents $q$ over $\mathbf{Z}$, this orbit has $\mathbf{Q}$-rational points, and $v$ equals +1 on
$X(\mathbf{R} \times \hat{\mathbf{Z}})$. Thus $\delta$ equals 2 on $X(\mathbf{R} \times \hat{\mathbf{Z}})$, and by Formulas (1) and (2) of Section $1.4 N(T, X) \sim 2 \mathbb{S}(X) \mu_{\infty}(T, X)$.

## 4. PROOF OF THEOREM 0.2

We prove Theorem 0.2. We define an involution $\tau_{\infty}$ of $X(\mathbf{R})$ by $\tau_{\infty}(x)=$ $-x, x \in X(\mathbf{R}) \subset \mathbf{R}^{3}$. Since $f(x)=f(-x), \tau_{\infty}$ is well defined, i.e takes $X(\mathbf{R})$ to itself. Since $|-x|=|x|, \tau_{\infty}$ takes $X(\mathbf{R})^{T}$ to itself. We define an involution $\tau$ of $X(\mathbf{A})$ by defining $\tau$ as $\tau_{\infty}$ on $X(\mathbf{R})$ and as 1 on $X\left(\mathbf{Q}_{p}\right)$ for all prime $p$. Then $\tau$ respects the Tamagawa measure $m$ on $X(\mathbf{A})$.

By assumption $X(\mathbf{R})$ has two connected components. These are the two orbits of $\operatorname{Spin}(f)(\mathbf{R})$. The involution $\tau_{\infty}$ of $X(\mathbf{R})$ interchanges these two orbits. Thus we have

$$
\begin{align*}
v_{\infty}\left(\tau_{\infty}\left(x_{\infty}\right)\right) & =-v_{\infty}\left(x_{\infty}\right) & & \text { for all } \tag{4}
\end{align*} \quad x_{\infty} \in X(\mathbf{R})
$$

Let $X(\mathbf{R})_{1}$ and $X(\mathbf{R})_{2}$ be the two connected components of $X(\mathbf{R})$. Set

$$
X(\mathbf{R})_{1}^{T}=X(\mathbf{R})_{1} \cap X(\mathbf{R})^{T}, \quad X(\mathbf{R})_{2}^{T}=X(\mathbf{R})_{2} \cap X(\mathbf{R})^{T}
$$

Then $\tau$ interchanges $X(\mathbf{R})_{1}^{T} \times X(\hat{\mathbf{Z}})$ and $X(\mathbf{R})_{2}^{T} \times X(\hat{\mathbf{Z}})$. From Formula (5) in this section we have

$$
\int_{X(\mathbf{R})_{1}^{T} \times X(\mathbf{Z})} v(x) d m=-\int_{X(\mathbf{R})_{2}^{T} \times X(\hat{\mathbf{Z}})} v(x) d m,
$$

hence

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} v(x) d m=0 .
$$

Since $\delta(x)=v(x)+1$, we obtain

$$
\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m=\int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} d m=m\left(X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})\right),
$$

and $m\left(X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})\right)=\mathfrak{S}(X) \mu_{\infty}(T, X)$. By Theorem 1.2

$$
N(T, X) \sim \int_{X(\mathbf{R})^{T} \times X(\hat{\mathbf{Z}})} \delta(x) d m .
$$

Thus $N(T, X) \sim \Theta(X) \mu_{\infty}(T, X)$ as $T \rightarrow \infty$, i.e. $c_{X}=1$.

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