國立臺灣大學理學院數學系
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# Department of Mathematics College of Science National Taiwan University Master Thesis 

論 290 定理
On the 290－theorem

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## 中文摘要

這篇文章中，我們研究 Manjul Bhargava 與 Jonathan Hanke 的 290 定理。圭要是透過分類一些被稱作上升子的特殊基本二次形，建立起決定任一正定整係數二次形是否為宇態二次形的準則。


#### Abstract

In this thesis, we study the 290-theorem of Manjul Bhargava and Jonathan Hanke. Via the classification of certain basic quadratic forms called escalators, we can estabs lish an efficient criterion to determine whether a positive integral quadratic form is universal.


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## 1 Introduction

The main aim of this thesis is to study a theorem of Manjul Bhargava and Jonathan Hanke that characterises all universal quadratic forms.


A quadratic form $Q\left(x_{1}, \ldots, x_{n}\right)$ takes integer values for all vectors $\left(x_{1}, \ldots, x_{n}\right)$ in $\mathbb{Z}^{n}$, if and only if $Q\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$, and if this holds it is called a integral quadratic form. An integral quadratic form $Q\left(x_{1}, \ldots, x_{n}\right)$ is said to represent an integer $z \in \mathbb{Z}$, if $Q\left(a_{1}, \ldots, a_{n}\right)=z$, for some non-zero $\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{Z}^{n}$, and $Q\left(x_{1}, \ldots, x_{n}\right)$ is positive-definite if it represents only positive integers. Finally, by a universal quadratic form we mean a positive-definite integral quadratic form that represents every natural number.

In history, many universal quadratic forms has been studied. A classical example comes from the famous Lagrange's four squares theorem which actually says the form $x^{2}+y^{2}+z^{2}+w^{2}$ is universal. Also, in [9] Ramanujan listed 54 other universal quadratic forms of type $a x^{2}+b y^{2}+c z^{2}+d w^{2}$, where $a, b, c, d$ are certain natural numbers.

Bhargava and Hanke [6] proves the following theorem which was originally conjectured by Conway in 1993.

Theorem 1. A positive-definite integral quadratic form is universal if and only if it represents the following 29 integers:

$$
\begin{gathered}
1,2,3,5,6,7,10,13,14,15,17,19,21,22,23,26,29 \\
\quad 30,31,34,35,37,42,58,93,110,145,203,290
\end{gathered}
$$

Definition 1.1. Those 29 integers listed in Theorem 1 are called critical integers.
Bhargava and Hanke [6] also proves that the 29 critical integers are minimal in the sense of the following:

Theorem 2. For each critical integer $c$, there is a non-universal positive-definite integral quadratic form represents other 28 critical integers except for $c$.

The proofs of the theorems relies on the concept of escalators of a given quadratic form $Q(x)$. To explain it, we apply the correspondence between quadratic forms and
lattices (see §2.4), roughly speaking, the lattice of an escalator is generated by the 霊 lattice of $Q(x)$ together with a vector having square norm equal to the smallest positive integer $a$ not represented by $Q(x)$ (see $\S 3$ for details). In particular, this escalater represents $a$. By continuing the escalating procedure, one will obtain escalators that represent all critical integers. Then one check if they are universal. It turns out one needs at most 7 steps of escalating to achieve the goal. However, there usually are numerous escalators of a given form. There is only one escalator of the zero form, namely the one-dimensional form $x^{2}$, which has 3 escalators, all 2-dimensional, and these 2-dimensional forms have totally 34 escalators, all 3-dimensional. These 3 -dimensional forms have 6560 escalators, all 4 -dimensional. Call these basic escalators. Fortunately, the proofs can be completed by just working on 4-dimensional basic escalators.

There is an effective (and efficient) algorithm to check if a basic 4-dimensional escalator locally represents an integer $m$ (see Lemma 4.18, 4.19). Although the HasseMinkowski theorem, or the local-global principle, does not always holds for integral forms, it is applicable to 1658 forms out of all 6560 basic 4-dimensional escalators so that one can determines if these forms represent an integer $m$ by checking the corresponding local representability.

For other basic 4-dimensional escalators, one needs to apply analytic method. By estimating the constants involved in a modified formula of Siegel (see (6.2.1)) and by using the Deligne's bounds of the Fourier coefficients of Hecke eigenforms, Bhargava and Hanke are able to find for each basic 4-dimensional escalator a lower bound $B$ such that the local-global principle holds for $m>B$. Fortunately, the lower bound $B$ is of reasonable size so that one can use machine to check if each $m \leq B$ is represented by the corresponding escalator. With the information of representability of all 4-dimensional basic escalators obtained, the theorems are then proved by a few deductions.

This thesis is organised in the following way. In §2, we review some basic facts and definitions including the correspondence between equivalence classes of integral quadratic forms and integral-square-norm lattices. Escalators are discussed in §3.

We review the local theory in $\S 4$ and the local-global principle in $\S 5$. The analytic method is studied in $\S 6$, and finally, the proofs are completed in $\S 7$.

## 2 Notation and setting

For simplicity, write $x$ for the vector $\left(x_{1}, \ldots, x_{n}\right)$ and denote $Q(x)=Q\left(x_{1}, \ldots, x_{n}\right)$. Let P be a partition of $\{1,2, \ldots, n\}$. For each $J \in P$, let $x_{J}$ denote the vector whose components are $x_{i}$ for $i \in J$. For a subset $S \subseteq P$, let $x_{S}$ be the vector whose components are $x_{i}$ for all $i \in \bigcup_{J \in S} J$.

### 2.1 The Gram matrix

Definition 2.1. The Gram matrix $A$ of a quadratic form $Q(x)$ is the unique symmetric matrix $A$ such that

$$
Q(x)=x^{T} A x, \text { for all } x .
$$

Call $D_{Q}:=\operatorname{det}(A)$ the determinant of $Q(x)$.
If $Q(x)$ is integral, then the entries of $A$ are contained in $\frac{1}{2} \mathbb{Z}$.
Definition 2.2. Let $Q(x)$ be an integral quadratic form with $A$ as Gram matrix. Let $N_{Q}$ be the minimal positive integer such that $N_{Q} \cdot(2 A)^{-1}$ is a matrix whose entries are integers and diagonal entries lie in $2 \mathbb{Z}$. Call $N_{Q}$ the level of $Q(x)$.

For example, the Gram matrix of

$$
\begin{equation*}
Q\left(x_{1}, x_{2}\right):=x_{1}^{2}+3 x_{1} x_{2}+5 x_{2}^{2} \tag{2.1.1}
\end{equation*}
$$

is

$$
A=\left(\begin{array}{cc}
1 & 3 / 2 \\
3 / 2 & 5
\end{array}\right) \quad \text { with } A^{-1}=\left(\begin{array}{cc}
20 / 11 & -6 / 11 \\
-6 / 11 & 4 / 11
\end{array}\right)
$$

while $D_{Q}=11 / 4$ and $N_{Q}=11$.

### 2.2 Integrally equivalence class

Two quadratic form $Q(x)$ and $Q^{\prime}(x)$ with Gam matrix $A$ and $A^{\prime}$ are integrally equivalent if and only if there exists some $B \in \mathrm{GL}(n, \mathbb{Z})$ such that

$$
A^{\prime}=B^{T} A B
$$

We shall refer this equivalence relation as the integral equivalence relation. Integrally equivalent quadratic forms have the same determinant. We define the determinant of an integral equivalence class as the that of any of its members.

Lemma 2.3. The number of integral equivalence classes of the same determinant is finite.

Proof. See [1, §9.2, corollary 1, p.129].

### 2.3 The Minkowski-reduced forms

Definition 2.4. Let $Q(x)$ be a positive-definite integral quadratic form in $n$ variables and let $A=\left(a_{i j}\right)$ be its Gram matrix. Call $Q(x)$ Minkowski-reduced if

$$
\begin{gathered}
Q(x) \geq a_{i i}, \text { whenever g.c.d. }\left(x_{i}, \ldots, x_{n}\right)=1, \\
a_{i, i+1} \geq 0, \text { for } i=1, \ldots, n-1
\end{gathered}
$$

The form $Q\left(x_{1}, x_{2}\right)$ in (2.1.1) is not Minkowski-reduced, because $Q(-1,1)<a_{22}$. However, $Q\left(x_{1}-x_{2}, x_{2}\right)=x_{1}^{2}+x_{1} x_{2}+3 x_{2}^{2}$ is Minkowski-reduced and is integrally equivalent to $Q\left(x_{1}, x_{2}\right)$. Indeed, each positive-definite integral quadratic form is integrally equivalent to at least one Minkowski-reduced form [2, p.27].

### 2.4 The corresponding lattices

Let $L$ be a full-rank lattice in the Euclidean space $\mathbb{R}^{n}$. Suppose $L$ is a $\mathbb{Z}$-lattice in the sense that $\|v\|^{2} \in \mathbb{Z}$, for all vector $v \in L$. Let $\varphi: \mathbb{Z}^{n} \longrightarrow L$ be an isomorphism of abelian groups. Then $Q(x):=\|\varphi(x)\|^{2}$ is a positive-definite integral quadratic form in $n$ variables. The form $Q(x)$ depends on $L$ as well as the choice of $\varphi$, while its integral
equivalence class depends only on $L$. Conversely, if $Q\left(x_{1}, \ldots x_{n}\right)$ is a positive-definite $\frac{\text { 龺 }}{}$ integral quadratic form with Gram matrix $A$, then the assignment $(u, v){ }_{4}{ }_{4} u^{T} A v$ defines an inner product on $\mathbb{R}^{n}$, and hence by the Gram Schmidt process, there is an orthogonal automorphism $\varphi$ of $\mathbb{R}^{n}$ such that $Q(x)=\|\varphi(x)\|^{2}$, for all $x \in \mathbb{Z}^{n}$. Then we see that $Q(x)$ is exactly the quadratic form induced by $\varphi$ and the lattice $L:=\varphi\left(\mathbb{Z}^{n}\right)^{\circ}$.

Thus, we have a bijection between equivalence classes and $\mathbb{Z}$-lattices with integral square norms. In this thesis, for convenience we will use the identification between a quadratic form $Q(x)$ and the lattice $L=L(Q)$ via certain implicitly chosen $\varphi$. Therefore, when we say that a lattice has certain property, it means the corresponding quadratic form has such property. For example, the lattice $\mathbb{Z}^{4}$ is universal, by Largrange's theorem.

Under the above correspondence, if $e_{1}, \ldots, e_{n}$ is the standard basis of $\mathbb{R}^{n}$ so that $f_{i}:=\varphi\left(f_{i}\right)$ is the corresponding $\mathbb{Z}$-basis of $L$, then the entries of the Gram matrix $A=\left(a_{i j}\right)$ of $Q(x)$ can be expressed as

$$
\begin{equation*}
a_{i j}=\left(f_{i}, f_{j}\right) \tag{2.4.1}
\end{equation*}
$$

Suppose $L \subset \mathbb{R}^{n}$ is a full $\mathbb{Z}$-lattice and let $f_{1}, \ldots, f_{n}$ be a $\mathbb{Z}$-basis of $L$. By (2.4.1), the lattice $2 L$ is contained in the dual lattice $L^{*}$ of $L$. Here

$$
\begin{equation*}
L^{*}:=\left\{x \in \mathbb{R}^{n} \mid(x, y) \in \mathbb{Z}, \text { for all } y \in L\right\} \tag{2.4.2}
\end{equation*}
$$

which is spanned by the dual basis $f_{1}^{*}, \ldots, f_{n}^{*}$ satisfying

$$
\begin{equation*}
\left(f_{i}^{*}, f_{j}\right)=\delta_{i, j}, \tag{2.4.3}
\end{equation*}
$$

where $\delta_{i, j}$ denotes the Kronecker symbol.

## 3 Escalation

The proofs of the main theorems are built on the concept of escalation. We also refer to the rank of a lattice as its dimension.

Definition 3.1. If a positive-definite integral quadratic form $Q(x)$ is not universal, we call the smallest positive integer that $Q(x)$ cannot represent the truant of $Q(x)$. We define an escalation of a non-universal $L$ to be any integral lattice generated $\overline{\bar{b}}$ $L$ and a vector whose square-norm is the truant of L. An escalator lattice is aldtice which is obtained by consecutive escalations of the zero-dimensional lattice.

Note that if a non-universal $L$ is a full lattice in $\mathbb{R}^{n}$, then its escalation is a full lattice in either $\mathbb{R}^{n}$ or $\mathbb{R}^{n+1}$.

Lemma 3.2. There are only finitely many escalator lattice of a given dimension $n$.
Proof. We prove by induction on $n$. The unique escalation of the one-dimensional lattice is the lattice $\mathbb{Z} \subset \mathbb{R}$ corresponding to the quadratic form $x^{2}$. Hence the case $n=1$ is proved. Suppose $L=\mathbb{Z} f+L_{1}$ is an $n$-dimensional escalation of an escalator lattice $L_{1}$ such that $\|f\|^{2}$ equals the truant $a$ of $L_{1}$.

We first consider the case where $L_{1}$ is of dimension $n-1$ and let $f_{1}, \ldots, f_{n-1}$ be a $\mathbb{Z}$-basis. Write $f=\alpha e_{n}+f^{\prime}$, where $e_{n} \in \mathbb{R}^{n}$ is the standard unit vector perpendicular to $\mathbb{R}^{n-1}$ and $f^{\prime} \in R^{n-1}$. By (2.4.1), the value

$$
\left(f^{\prime}, f_{i}\right)=\left(f, f_{i}\right)
$$

must be either an integer or a half integer. Hence by (2.4.3), the vector $f^{\prime} \in \frac{1}{2} L_{1}^{*}$. The discrete subgroup $\frac{1}{2} L_{1}^{*} \subset \mathbb{R}^{n-1}$ contains only finitely many vectors $f^{\prime}$ with $\left\|f^{\prime}\right\|^{2}<a$. For each such $f^{\prime}$, the number $\alpha$ must equal $\pm \sqrt{a-\left\|f^{\prime}\right\|^{2}}$. Hence a given $L_{1}$ can only have finitely many escalations. This together with the induction hypothesis implies that there are finitely many $n$-dimensional escalation of escalator lattice of dimension $n-1$. Finally, we note that if $L_{1}$ is $n$-dimensional with $f_{1}, \ldots, f_{n}$ as a $\mathbb{Z}$-basis and $L$ is an $n$-dimensional integral lattice containing $L_{1}$, then by (2.4.1), for every $f \in L$, the values $\left(f, f_{i}\right), i=1, \ldots, n$ are in $\frac{1}{2} \mathbb{Z}$, and hence $L \in \frac{1}{2} L_{1}^{*}$, by (2.4.3). Therefore, $L / L_{1} \subset \frac{1}{2} L_{1}^{*} / L_{1}$ which is finite, and hence there are only finitely many such $L$. In particular, there can only be finitely many $n$-dimensional escalator lattices containing $L_{1}$.

To determine escalator lattices of low dimensions, we follow the proof of Lemma 3.2. The only 1-dimensional escalator lattice is $L_{1}=\mathbb{Z} \subset R$ with with truant 2 and $L_{1}^{*}=L_{1}$. Denote $f_{1}=1 \in L_{1}$. Since no vector in $\frac{1}{2} L_{1}$ represents 2 , an escalation $L^{2}$ of $L_{1}$ must be 2-dimensional with $L=\mathbb{Z} f+L_{1}$ and $f=\alpha e_{2}+x f_{1}$. Now $x$ is contained in $\frac{1}{2} \mathbb{Z}$ satisfying $x^{2}<2$. Therefore, $x=0, \pm \frac{1}{2}, \pm 1$. These gives three nonisometric two-dimensional escalators having Minkowski-reduced matrices:

$$
\left(\begin{array}{ll}
1 & 0 \\
0 & 2
\end{array}\right),\left(\begin{array}{cc}
1 & \frac{1}{2} \\
\frac{1}{2} & 2
\end{array}\right) \text { and }\left(\begin{array}{cc}
1 & 0 \\
0 & 1
\end{array}\right)
$$

The corresponding lattices are $\mathbb{Z} \sqrt{2} e_{2}+\mathbb{Z} f_{1}, \mathbb{Z}\left(\frac{\sqrt{7}}{2} e_{2}+\frac{1}{2} f_{1}\right)+\mathbb{Z} f_{1}$, and $\mathbb{Z} e_{2}+\mathbb{Z} f_{1}$ whose truants are respectively 5,3 , and 3 . It turns out that all these three lattices have no 2-dimensional escalations. Escalating them, we obtain 34 three-dimensional nonisometric escalator lattices which have no three-dimensional escalations, and there are actually 6560 four-dimensional nonisometric escalations of these 34 lattices, [6]. We call these the basic escalators.

Lemma 3.3. Each universal positive-definite integral quadratic form must contain a universal escalator lattice. Conversely, the truant of any non-universal form is the same as the truant of some non-universal escalator lattice within it.

Proof. Let $Q$ be a positive-definite integral quadratic form and denote $L=L(Q)$. There exists a maximal sequence of escalator lattices

$$
\{0\} \subset L_{1} \subset L_{2} \subset \ldots \subset L_{k}
$$

within $L$. Since $L_{k}$ is maximal, it is either universal or having the same truant as that of $L$. Because $L_{k} \subset L$, if $L_{k}$ is universal, so is $L$.

## 4 The $\mathbb{Z}_{p}$-theory

To see if an integral quadratic form represents an integer $m$, we can first check if it represents $m$ locally. In this section, we study the local theory of quadratic forms. Let $p$ be a finite prime number. A quadratic form $Q(x)$ is called $\mathbb{Z}_{p}$-integral if and
only if $Q(x) \in \mathbb{Z}_{p}[x]$. Two $\mathbb{Z}_{p}$-integral quadratic forms $Q$ and $Q^{\prime}$ are $\mathbb{Z}_{p}$-equivalent if ${ }^{\text {B }}$ and only if their Gram matrices $A$ and $A^{\prime}$ satisfy

$$
A^{\prime}=B^{T} A B, \text { for some } B \in \mathrm{GL}\left(\mathrm{n}, \mathbb{Z}_{\mathrm{p}}\right)
$$



### 4.1 The normalized form

Definition 4.1. A quadratic form $Q(x)$ is $\mathbb{Z}_{p}$-elementary if and only if either $Q(x)$ is one-dimensional equal to $u x^{2}, u \in \mathbb{Z}_{p}^{*}$, or $p=2$ and $Q(x)$ is 2-dimensional with Gram matrix $\left(\begin{array}{ll}a_{11} & a_{12} \\ a_{12} & a_{22}\end{array}\right)$, such that $a_{11}, a_{22} \in \mathbb{Z}_{2}$ and $2 a_{12}, 2 a_{21} \in \mathbb{Z}_{2}^{*}$.
Definition 4.2. An $n$-dimensional $\mathbb{Z}_{p}$-integral quadratic form $Q(x)$ is a normalized form if and only if there is a partition $P$ of $\{1, \ldots, n\}$ such that

$$
\begin{equation*}
Q(x)=\sum_{J \in P} p^{\nu_{J}} Q_{J}\left(x_{J}\right) \tag{4.1.1}
\end{equation*}
$$

where each $Q_{J}$ is $\mathbb{Z}_{p}$-elementary and each $\nu_{J} \geq 0$.
Lemma 4.3. Every quadratic form $Q(x)$ is $\mathbb{Z}_{p}$-equivalent to a normalized form.
Proof. Let $A=\left(a_{i j}\right)$ denote the Gram matrix of $Q(x)$. We prove by induction on $n$, the dimension of $Q(x)$. If $\operatorname{dim} Q=1$, the lemma obviously holds. We first consider the case in which $p \neq 2$. Suppose $\min \left\{\operatorname{ord}_{p}\left(a_{i j}\right) \mid 1 \leq i, j \leq n\right)$ occurs at some $i=j$. We may assume it occurs at $i=j=1$. By completing the squares, we can cancel $a_{12}, \ldots, a_{1 n}$ as well as $a_{21}, \ldots, a_{n 1}$. Then

$$
\begin{equation*}
Q(x) \simeq p^{\operatorname{ord}_{p} a_{11}} \cdot x_{1}^{2} \oplus Q^{\prime} \tag{4.1.2}
\end{equation*}
$$

with $\operatorname{dim} Q^{\prime}=n-1$. Then the proof is completed by the induction hypothesis.
Suppose $\min \left\{\operatorname{ord}_{p}\left(a_{i j} \mid 1 \leq i, j \leq n\right)\right.$ does not occur at diagonal entries. We may assume it equals to $\operatorname{ord}_{p} a_{12}$. Then we add the second column and the second row into the first ones. By this process, $a_{11}$ is replaced by $a_{11}+2 a_{12}+a_{22}$. Since $\operatorname{ord}_{p}\left(a_{11}+2 a_{12}+a_{22}\right)=\operatorname{ord}_{p}\left(a_{12}\right)=\min \left\{\operatorname{ord}_{p}\left(a_{i j}\right) \mid 1 \leq i, j \leq n\right)$, the proof is reduced the previous case.

Now we consider the $p=2$ case. Similarly, if $\min \left\{\operatorname{ord}_{p}\left(a_{i j}\right) \mid 1 \leq i, j \leq n\right)$ occurs at some diagonal entry, the lemma is proved by completing the squares. Otherwise, we may assume

$$
\min \left\{\operatorname{ord}_{p}\left(a_{i j}\right) \mid 1 \leq i, j \leq n\right)=\operatorname{ord}_{p}\left(a_{12}\right)=k
$$

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Write the leading $2 \times 2$ submatrix as

$$
p^{k-1}\left(\begin{array}{cc}
p^{\alpha} b_{11} & b_{12}  \tag{4.1.3}\\
b_{12} & p^{\beta} b_{22}
\end{array}\right):=p^{k-1} E
$$

where $\alpha, \beta \geq 1$ and $\left|b_{12}\right|_{p}=\left|b_{21}\right|_{p}=1$. Then,

$$
\begin{equation*}
|\operatorname{det} E|_{p}=\left|p^{\alpha+\beta} b_{11} b_{22}-b_{12} b_{21}\right|_{p}=1 \tag{4.1.4}
\end{equation*}
$$

Since its determinant is a unit in $\mathbb{Z}_{2}$, the matrix $E$ is invertible. Write

$$
A=\left(\begin{array}{cc}
p^{k-1} E & p^{k-1} C \\
p^{k-1} C^{T} & U
\end{array}\right)
$$

for some matrices $C$ and $U$. Define $F:=-E^{-1} C$ and put

$$
B:=\left(\begin{array}{cc}
I_{2} & F \\
F^{T} & I_{n-2}
\end{array}\right)
$$

Then

$$
B^{T} A B=\left(\begin{array}{cc}
p^{k-1} E & 0 \\
0 & A^{\prime}
\end{array}\right)
$$

for some $(n-2) \times(n-2)$ matrix $A^{\prime}$. This means $Q(x) \simeq p^{k} Q_{1} \oplus Q^{\prime}$, where $Q_{1}$ has $p^{-1} E$ as its Gram matrix and $\operatorname{dim} Q^{\prime}=n-2$. By the induction hypothesis, the proof is complete.

Lemma 4.4. Let $Q(x)$ be an integral quadratic form. Then the following statements are equivalent:
(a) $p \mid N_{Q}$,
(b) $p \mid D_{2 Q}$,
(c) if $Q(x)$ is $\mathbb{Z}_{p}$-equivalent to a normalized form expressed by (4.1.1), then either some $v_{J} \geq 1$, or $p=2$ and some $Q_{J}$ is 1-dimensional.

Proof. By lemma 4.3, we may assume that $Q(x)$ is $\mathbb{Z}_{p}$-equivalent to a normalized form expressed by (4.1.1). We have $\operatorname{ord}_{p}\left(D_{2 Q}\right)=\sum_{J \in P} \lambda_{J}^{\prime}$, where

$$
\lambda_{J}^{\prime}=\left\{\begin{array}{cl}
\nu_{J}+\operatorname{ord}_{p}(2), & \text { if } \operatorname{dim} Q_{J}=1 \\
2 \nu_{J}, & \text { if } \operatorname{dim} Q_{J}=2,
\end{array}\right.
$$

and also $\operatorname{ord}_{p}\left(N_{Q}\right)=\max _{J \in P} \lambda_{i}^{\prime \prime}$, where

$$
\lambda_{i}^{\prime \prime}=\left\{\begin{array}{cc}
\nu_{J}+2 \operatorname{ord}_{p}(2), & \text { if } \operatorname{dim} Q_{J}=1 \\
\nu_{J}, & \text { if } \operatorname{dim} Q_{J}=2
\end{array}\right.
$$

### 4.2 The reduction maps

Let $Q(x)$ be a normalized form expressed by (4.1.1) and set

$$
S_{0}:=\left\{J \in P \mid \nu_{J}=0\right\}, S_{1}:=\left\{J \in P \mid \nu_{J}=1\right\}, \text { and } S_{2}:=\left\{J \in P \mid \nu_{J} \geq 2\right\} .
$$

Definition 4.5. For a given $m \in \mathbb{Z}_{p}$, define

$$
\begin{aligned}
& R_{p^{k}, Q}(m):=\left\{x \in \mathbb{Z}_{p}^{n} / p^{k} \mathbb{Z}_{p}^{n} \mid Q(x) \equiv m \quad\left(\bmod p^{k}\right)\right\}, \\
& R_{p^{k}, Q}^{\text {Eero }}(m):=\left\{x \in R_{p^{k}, Q}(m) \mid x \equiv 0 \quad(\bmod p)\right\}, \\
& R_{p^{k}, Q}^{G o o d}(m):=\left\{x \in R_{p^{k}, Q}(m) \mid p^{\nu_{J}} x_{J} \not \equiv 0 \quad(\bmod p), \text { for some } J \in P\right\}, \\
& R_{p^{k}, Q}^{\text {Bad }}(m):=\left\{x \in R_{p^{k}, Q}(m) \mid x \not \equiv 0 \text { but } x_{S_{0}} \equiv 0 \quad(\bmod p)\right\},
\end{aligned}
$$

and let $r_{p^{k}, Q(m)}, r_{p^{k}, Q(m)}^{Z e r o}, r_{p^{k}, Q(m)}^{G \operatorname{ood}}$ and $r_{p^{k}, Q(m)}^{\text {Bad }}$ denote respectively the cardinality of the sets $R_{p^{k}, Q}(m), R_{p^{k}, Q}^{Z \text { ero }}(m), R_{p^{k}, Q}^{G o o d}(m)$ and $R_{p^{k}, Q}^{\text {Bad }}(m)$.

Remark 4.6. We have

$$
R_{p^{k}, Q}(m)=R_{p^{k}, Q}^{Z \text { ero }}(m) \sqcup R_{p^{k}, Q}^{G o o d}(m) \sqcup R_{p^{k}, Q}^{\text {bad }}(m) .
$$

Also, if $p \nmid m$, then $R_{p^{k}, Q}(m)=R_{p^{k}, Q}^{G o o d}(m)$; if $p \nmid N_{Q}$, then $R_{p^{k}, Q}^{\text {Bad }}(m)=\emptyset$.
Now, reduction maps of various types are in order.

### 4.2.1 The reduction map of good type

Definition 4.7. Let $\pi_{G}: R_{p^{k}, Q}^{G o o d}(m) \longrightarrow R_{p^{k-1, Q}}^{G o o d}(m)$ denote the reduction map sending $x\left(\bmod p^{k} \mathbb{Z}_{p}\right)$ into $x\left(\bmod p^{k-1} \mathbb{Z}_{p}\right)$.

Lemma 4.8. For $k \geq 2 \operatorname{ord}_{p} 2+1$, the map $\pi_{G}$ is surjective with multiplicity $p^{n+1}$, , , namely,

$$
r_{p^{k+1}, Q}^{G o o d}(m)=p^{(n-1)} \cdot r_{p^{k}, Q}^{G o o d}(m)
$$

Proof. This is a direct consequence of Hensel's Lemma [10, §I. 3 Theorem 3].

### 4.2.2 The reduction map of zero type

If $x \in R_{p^{k}, Q}^{Z \text { ero }}(m)$, then $x \mid p$ and hence $p^{2} \mid m$.
Definition 4.9. If $R_{p^{k}, Q}^{Z \text { ero }}(m)$ is nonempty, define the reduction map

$$
\pi_{Z}: R_{p^{k}, Q}^{Z \text { ero }}(m) \longrightarrow R_{p^{k-2}, Q}\left(\frac{m}{p^{2}}\right)
$$

to be the one sending $x\left(\bmod p^{k} \mathbb{Z}_{p}\right)$ into $\frac{x}{p}\left(\bmod p^{k-1} \mathbb{Z}_{p}\right)$.
Obviously, the following lemma holds.
Lemma 4.10. This map is surjective with multiplicity $p^{n}$

### 4.2.3 The reduction map of bad type

Define $Q^{\prime}(x)=\sum_{J \in P} p^{\nu_{J}^{\prime}} Q_{J}$ and $Q^{\prime \prime}(x)=\sum_{J \in P} p^{\nu_{J}^{\prime \prime}} Q_{J}$, where

$$
\nu_{J}^{\prime}=\left\{\begin{array}{ll}
\nu_{J}+1, & \text { if } J \in S_{0} ; \\
\nu_{J}-1, & \text { if } J \notin S_{0},
\end{array} \quad \nu_{J}^{\prime \prime}= \begin{cases}\nu_{J}, & \text { if } J \in S_{0} \cup S_{1} ; \\
\nu_{J}-2, & \text { if } J \in S_{2}\end{cases}\right.
$$

Denote, for $i=0,1,2$,

$$
R_{p^{k}, Q}^{x_{S_{i} \neq 0}}(m):=\left\{x \in R_{p^{k-2}, Q}(m) \mid x_{S_{i}} \not \equiv 0 \quad\left(\bmod p^{k}\right)\right\}
$$

and also

$$
\begin{array}{ll}
R_{p^{k}, Q}^{B a d, x_{S_{1}} \neq 0}(m)=\left\{x \in R_{p^{k}, Q}^{B a d}(m) \mid x_{S_{1}} \not \equiv 0\right. & (\bmod p)\}, \\
R_{p^{k}, Q}^{\text {Bad, } x_{S_{1}} \equiv 0}(m)=\left\{x \in R_{p^{k}, Q}^{B a d}(m) \mid x_{S_{1}} \equiv 0\right. & (\bmod p)\} .
\end{array}
$$

Definition 4.11. If $S_{1} \neq \emptyset$ define

$$
\pi_{B^{\prime}}: R_{p^{k}, Q}^{\text {Bad, }, x_{s_{1}} \neq 0}(m) \longrightarrow R_{p^{k-1}, Q^{\prime}}^{G o o d}\left(\frac{m}{p}\right)
$$

to be the map sending $x_{J}\left(\bmod p^{k}\right)$ into $\frac{x_{J}}{p}\left(\bmod p^{k-1}\right)$, for $J \in S_{0}$; and $x_{J}\left(\bmod p^{k}\right)$ into $x_{J}\left(\bmod p^{k-1}\right)$, otherwise.

Lemma 4.12. The map $\pi_{B^{\prime}}$ is surjective with multiplicity $p^{\# S_{1} \cup S_{2}}$.
Proof. We can choose arbitrary lift of $x_{J}\left(\bmod p^{k-1}\right)$ for $J \in S_{1} \bigcup S_{2}$.
Definition 4.13. If $S_{1}=\emptyset$ or $x_{S_{1}} \equiv 0(\bmod p)$, define

$$
\pi_{B^{\prime \prime}}: R_{p^{k}, Q}^{B a d, x_{S_{1}} \equiv 0}(m) \longrightarrow R_{p^{k-2}, Q^{\prime \prime}}^{x_{S_{2} \neq 0}}\left(\frac{m}{p^{2}}\right)
$$

to be the map sending $x_{J}\left(\bmod p^{k}\right)$ into $\frac{x_{j}}{p}\left(\bmod p^{k-2}\right)$, if $J \in S_{0} \bigcup S_{1} ; x_{J}\left(\bmod p^{k}\right)$ into $x_{J}\left(\bmod p^{k-2}\right)$, otherwise.

Similarly, the folllowing lemma holds.
Lemma 4.14. The map $\pi_{B^{\prime \prime}}$ is surjective with multiplicity $p^{\# S_{0} \cup S_{1}+2 \# S_{2}}$.

### 4.2.4 The depth

Definition 4.15. We define the depth of each type of solution of $R_{p^{k}, Q}(m)$ to be the maximal difference $k-k^{\prime}$ for any $x \in R_{p^{k}, Q}(m)$ to be mapped into $R_{p^{k^{\prime}}, Q^{*}}\left(m^{*}\right)$ under consecutive application of the maps $\pi_{G}, \pi_{Z}$, and $\pi_{B^{\prime}}$, or $\pi_{B^{\prime \prime}}$ for some $Q^{*}$ and $m^{*}$

Under these definitions, we can easily gain the depth of each type:
Lemma 4.16. The Good-, Zero-, Bad-type depths of $R_{p^{k}, Q}(m)$ are bounded above by $k-1, \operatorname{ord}_{p}(m)$, and $\operatorname{ord}_{p}\left(N_{Q}\right)+1$ respectively.

### 4.3 The representability

We say that an $n$-dimensional $\mathbb{Z}_{p}$-integral form $Q(x)$ represents a $p$-adic integer $m$, if and only if there is some non-zero $x \in \mathbb{Z}_{p}^{n}$ such that $Q(x)=m$, while when referring to " $Q(x)$ representing $m\left(\bmod p^{k}\right)$ ", we use the following definition.

Definition 4.17. An n-dimensional $\mathbb{Z}_{p}$-integral form $Q(x)$ represents a $p$-adic integer $m$ modulo $p^{k}$ if and only if there is some $x \in \mathbb{Z}_{p}^{n}, x \not \equiv 0(\bmod p)$ such that $Q(x)=m$ $\left(\bmod p^{k}\right)$. In other words, $x\left(\bmod p^{k}\right)$ is contained in $R_{p^{k}, Q}$ but not in $R_{p^{k}, Q}^{Z \text { его }}$.

For $n \geq 3$, we have the following useful results. Define

$$
R_{Q}(m):=\left\{\tilde{x} \in \mathbb{Z}_{p}^{n} \mid \quad Q(\tilde{x})=m, \tilde{x} \not \equiv 0 \quad(\bmod p)\right\}
$$

and for $i=0,1,2$,

$$
\left.\begin{array}{rl}
R_{Q}^{\tilde{x}_{S_{i}} \equiv 0}(m):=\left\{\tilde{x} \in \mathbb{Z}_{p}^{n} \mid Q(\tilde{x})=m, \tilde{x}_{S_{i}} \equiv 0\right. & (\bmod p)\}, \\
R_{Q}^{\tilde{x}_{i_{i}} \neq 0}(m): & =\left\{\tilde{x} \in \mathbb{Z}_{p}^{n} \mid Q(\tilde{x})=m, \tilde{x}_{S_{i}} \not \equiv 0\right.
\end{array}(\bmod p)\right\} .
$$

Lemma 4.18. Let $Q$ be an $n$-dimensional $\mathbb{Z}_{p}$-integral quaderatic form. If $n \geq 3$ and $p \nmid N_{Q}$, then $Q(x)$ represents every $p$-adic integer $m$.

Proof. First assume that $p \neq 2$. Then $Q(x)$ is $\mathbb{Z}_{p}$-equivalent to the normalized form $u_{1} x_{1}^{2}+\cdots+u_{n} x_{n}^{2}$, with $u_{1}, \ldots, u_{n} \in \mathbb{Z}_{p}^{*}$. By [10, $\S 1.62$ Corollary 2, p. 393 Theorem 5], $Q(x)$ represents $m(\bmod p)$, and hence $R_{p, Q}^{G o o d}(m)$ is non-empty. Then we apply Lemma 4.8.

Suppose $p=2$. By Lemma 4.4, $n$ must be even, and hence $n \geq 4$. Thus,

$$
Q(x)=Q_{1}\left(x_{1}, x_{2}\right)+Q_{2}\left(x_{3}, x_{4}\right)+Q^{\prime}\left(x^{\prime}\right),
$$

where $x^{\prime}$ is either (if $n=4$ ) trivial or equals $\left(x_{5}, \ldots x_{n}\right)$ (if $n>4$ ) and

$$
\begin{array}{ll}
Q_{1}\left(x_{1}, x_{2}\right)=2^{\alpha_{1}} a_{1} x_{1}^{2}+u_{1} x_{1} x_{2}+2^{\beta_{1}} b_{1} x_{2}^{2}, & a_{1}, b_{1}, u_{1} \in \mathbb{Z}_{2}^{*}, \\
\alpha_{1}, \beta_{1}>0 \\
Q_{2}\left(x_{3}, x_{4}\right)=2^{\alpha_{2}} a_{2} x_{3}^{2}+u_{2} x_{3} x_{4}+2^{\beta_{2}} b_{2} x_{4}^{2}, & a_{1}, b_{1}, u_{1} \in \mathbb{Z}_{2}^{*},
\end{array} \alpha_{2}, \beta_{2}>0 . ~ \$ ~ \$
$$

Because $Q_{1}(1,1) \in \mathbb{Z}_{2}^{*}$ and $Q_{1}\left(x_{1}+2, x_{2}\right) \equiv Q_{1}\left(x_{1}, x_{2}\right)+2 u_{1} x_{2}(\bmod 8)$, we see that $Q_{1}\left(x_{1}, x_{2}\right)$ represents all odd integers modulo 8. Therefore, $Q_{1}\left(x_{1}, x_{2}\right)+Q_{2}\left(x_{3}, x_{4}\right)$ represents all integers modulo 8 and so is $Q(x)$. Then apply again Lemma 4.8.

Lemma 4.19. Let $Q$ be an n-dimensional $\mathbb{Z}_{p}$-integral quaderatic form with $n \geq 3$ and $p \mid N_{Q}$. A integer $m$ is represented by $Q(x)$ over $\mathbb{Z}_{p}$ if and only if some quotient of $m$ by square factor is represented by $Q(x)$ modulo $p^{\text {ord }_{p}\left(4 N_{Q}\right)+2}$.

Proof. If $Q(x)=m$ in $\mathbb{Z}_{p}$ and $x=p^{\nu} x^{\prime}$ with $x^{\prime} \in \mathbb{Z}_{p}$ such that $x^{\prime} \not \equiv 0(\bmod p)$, then 逼 $Q\left(x^{\prime}\right) \equiv \frac{m}{p^{2 \nu}}\left(\bmod p^{\operatorname{ord}_{p}\left(4 N_{Q}\right)+2}\right)$ as desired.

To prove the implication in the other direction, we may assume that $Q(x)$ repre sents $m$ modulo $p^{\operatorname{ord}_{p}\left(4 N_{Q}\right)+2}$. Put $k=\operatorname{ord}_{p}\left(4 N_{Q}\right)+2$. Suppose $Q(x) \equiv m\left(\bmod p p^{k}\right)$ and $x \not \equiv 0(\bmod p)$, then $x \notin R_{p^{k}, Q}^{Z \text { ero }}(m)$.

If $x \in R_{p^{k}, Q}^{G o o d}(m)$, then we apply Lemma 4.8 to ensure the existence of some $\tilde{x} \in R_{Q}^{\tilde{x} \neq 0}(m)$ such that $\tilde{x} \equiv x\left(\bmod p^{k}\right)$. In particular, $R_{Q}(m)$ is non-empty.

If $x \in R_{p^{k}, Q}^{\text {Bad, } x_{S_{1}} \neq 0}(m)$, then $\pi_{B^{\prime}}(x) \in R_{p^{k-1}, Q^{\prime}}^{G \text { Good }}\left(\frac{m}{p}\right)$ and since $k-1 \geq 2 \operatorname{ord}_{p} 2+1$, Lemma 4.8 is applicable. Therefore, there there exists some $\tilde{t} \in R_{Q^{\prime}}^{\tilde{x}_{S_{0}} \neq 0}\left(\frac{m}{p}\right)$ such that $\tilde{t} \equiv \pi_{B^{\prime}}(x)\left(\bmod p^{k-1}\right)$. Consider the commutative diagram

where down-arrows are respectively reduction modulo $p^{k}$ and $p^{k-1}$, while $\varphi^{\prime}$ is the bijection sending $\tilde{x}$ to $\tilde{t}$ with $\tilde{t}_{S_{0}}=\frac{1}{p} \cdot \tilde{x}_{S_{0}}$ and $\tilde{t}_{S_{1} \cup S_{2}}=\tilde{x}_{S_{1} \cup S_{2}}$. The diagram shows $R_{Q}^{\tilde{x}_{S_{1}} \neq 0}(m)$, end hence $R_{Q}(m)$ is non-empty.

If $x \in R_{p^{k}, Q}^{B a d, x_{S_{1}} \equiv 0}(m)$, then we apply the reduction map $\pi_{B^{\prime \prime}}$ consecutively $k^{\prime}$ times until the image of $x$ launches either $R_{p^{k-2 k^{\prime}}, Q^{*}}^{\text {Bad, } x_{S_{1}} \neq 0}\left(\frac{m}{p^{2 k^{\prime}}}\right)$ or $R_{p^{k-2 k^{\prime}}, Q^{*}}^{\text {Good }}\left(\frac{m}{p^{2 k^{\prime}}}\right)$. Since we have $N_{Q^{*}} \leq p^{-2 k^{\prime}} N_{Q}$, Lemma 4.8 is applicable to $Q^{*}(x)$ and $\frac{m}{p^{2 k^{\prime}}}$. To reduce the proof to the previous cases, we only need to consider the bijection

$$
\varphi^{\prime \prime}: R_{Q}^{\tilde{x_{S}} \equiv 0}(m) \cap R_{Q}^{\tilde{S}_{S_{1}} \equiv 0}(m) \longrightarrow R_{Q^{\prime \prime}}\left(\frac{m}{p^{2}}\right)
$$

as well as the associated commutative diagrams:


Here, if $\tilde{t}=\varphi^{\prime \prime}(\tilde{x})$, then $\tilde{t}_{S_{0} \cup S_{1}}=\frac{1}{p} \tilde{x}_{S_{0} \cup S_{1}}$ and $\tilde{t}_{S_{2}}=\tilde{x}_{S_{2}}$.

## 5 The local-global principle

The Hasse-Minkowski theorem says a $\mathbb{Q}$-rational quadratic form represents anan tional number if and only it represents this number locally at all primes including the infinite prime. This local-global principle also holds to some extend for integral forms. In this section, we review such wonderful theory, our reference is Cassels [1]. The main result is Theorem 5.3 below. For convenience, denote

$$
\mathbb{Z}_{\infty}=\mathbb{R}
$$

Definition 5.1. Two integral quadratic forms are in the same genus if and only if they are $\mathbb{Z}_{p}$-equivalent at all primes including the infinite one. The number of integral equivalence classes in a given genus is called the class number.

Lemma 5.2. The class number of a genus is finite.

Proof. Two forms in the same genus must be of the same determinant, while the number of integral equivalence classes in the collection of forms of a given determinant is finite (Lemma 2.3).

Theorem 5.3. Let $Q(x)$ be an integral form in $n$ variables of determinant $d \neq 0$. Let $a \neq 0$ be an integer which is represented by $Q(x)$ over $\mathbb{Z}_{p}$, for $p=\infty, 2,3, \ldots$, then there exists some form $Q^{*}(x)$ in the same genus of $Q(x)$ representing $a$.

The above two theorems will be proved at the end of this section.

### 5.1 Basic results

The following results are consequence of Lemma 4.3 or linear algebra over $\mathbb{Z}$ or $\mathbb{Z}_{p}$. We omit their proofs.

Lemma 5.4. Let $p$ be a finite prime number. Suppose $f_{1}$ and $f_{2}$ are two $\mathbb{Z}_{p}$-integral forms and that $D_{f_{1}}=D_{f_{2}}$ is a unit in $\mathbb{Z}_{p}$, then these two forms are $\mathbb{Z}_{p}$-equivalent.

Lemma 5.5. Let $p$ be a finite prime number and let $c_{1}, \ldots, c_{J}$ be linearly independent elements in $\mathbb{Z}_{p}^{n}$. Then the following three conditions are all equivalent:
(i) There exist $c_{J+1}, \ldots, c_{n}$ such that $c_{1}, \ldots, c_{n}$ forms a basis for $\mathbb{Z}_{p}^{n}$.
(ii) The $n \times J$ matrix $c_{1}, \ldots, c_{J}$ contains a $J \times J$ minor whose determind $n$ is a p-adic unit.
(iii) If $a=l_{1} c_{1}+\cdots+l_{J} c_{J} \in \mathbb{Z}_{p}^{n}$ for some $l_{1}, \cdots, l_{J} \in \mathbb{Q}_{p}$, then it must imply the result $l_{1}, \cdots, l_{J} \in \mathbb{Z}_{p}$.

Lemma 5.6. Let $c_{1}, \ldots, c_{J}$ be elements in $\mathbb{Z}^{n}$. Then the following three conditions are all equivalent:
(i) There exist $c_{J+1}, \ldots, c_{n}$ such that $c_{1}, \ldots, c_{n}$ forms a basis for $\mathbb{Z}^{n}$.
(ii) The determinants of the $J \times J$ submatrices of the $n \times J$ matrix $\left(c_{1}, \ldots, c_{n}\right)$ have no common divisor larger than 1 .
(iii) If $a=l_{1} c_{1}+\cdots+l_{J} c_{J} \in \mathbb{Z}^{n}$ for some $l_{1}, \cdots, l_{J} \in \mathbb{Q}$, then it must imply the result $l_{1}, \cdots, l_{J} \in \mathbb{Z}$.

### 5.2 The approximation of $\mathbb{Z}_{p}$-forms

We first prove the following lemma on the simultaneous approximation. Let $P$ denote a given collection of finitely many prime numbers.

Lemma 5.7. Let $K>1$ be an integer and let $m_{k}^{(p)} \in \mathbb{Z}_{p}, 1 \leq k \leq K, p \in P$, be given such that for every $p \in P$,

$$
\begin{equation*}
\max _{1 \leq k \leq K}\left|m_{k}^{(p)}\right|_{p}=1 \tag{5.2.1}
\end{equation*}
$$

Then for each $\epsilon>0$, there exist $m_{k} \in \mathbb{Z}, 1 \leq k \leq K$, with

$$
\operatorname{gcd}\left(m_{1}, \ldots, m_{k}\right)=1
$$

such that

$$
\begin{equation*}
\left|m_{k}-m_{k}^{(p)}\right|_{p}<\varepsilon, \quad \text { for all } p \in P \tag{5.2.2}
\end{equation*}
$$

Proof. By the Chinese Remainder Theorem, we pick $m_{1} \neq 0$ such that (5.2.2) holds for $k=1$. Let $P^{*}$ be the set of primes $p^{*}$ which divide $m_{1}$ but are not in P. Again,
apply the Chinese Remainder Theorem, we can find an $m_{2} \in \mathbb{Z}$ such that (5.2.2) 逼 holds for $k=2$, and

$$
p^{*} \nmid m_{2}
$$

for all $p^{*} \in P^{*}$. Now, for $k>2$ we choose $m_{k} \in \mathbb{Z}$ to satisfy (5.2.2). We may assume that $\varepsilon<1$. Then (5.2.1) implies $\operatorname{gcd}\left(m_{1}, \ldots, m_{k}\right)$ is not divisible by any $p \in P$. The condition 5.2.3 ensures that it is not divisible by any other primes. Thus, the proof is complete.

For a $\mathbb{Z}_{p}$-integral quadratic form $f(x)$ in $n$-variables with $D_{f} \neq 0$, let $\mathrm{O}_{f}$ denote the subgroup of GL $\left(n, \mathbb{Z}_{p}\right)$ consisting of matrice $\tau$ such that $f(x)=f(\tau x)$. For such $\tau$, we have $(\operatorname{det} \tau)^{2}=1$, and hence $\operatorname{det} \tau= \pm 1$.

Lemma 5.8. Let $f(x)$ be a $\mathbb{Z}_{p}$-integral quadratic form in $n$-variables with $D_{f} \neq 0$. Then there exists $\tau_{f} \in \mathrm{O}_{f}$ such that $\operatorname{det} \tau_{f}=-1$.

Proof. We first consider the case where $f(x)$ is $\mathbb{Z}_{p}$-elementary. If $f(x)=x^{2}$, then take $\tau_{f}=-1$. If $\operatorname{dim} f=2$, then $D_{2 f}$ is a unit in $\mathbb{Z}_{2}^{*}$ and is $\equiv-u^{2}(\bmod 16)$, for some $u \in \mathbb{Z}_{2}^{*}$. This implies $D_{2 f}=-\mu^{2}+16$ for some $\mu \in \mathbb{Z}_{2}^{*}$. Lemma 5.4 says $2 f$ is $\mathbb{Z}_{p}$-equivalent to $2 g(x):=4 x_{1}^{2}+2 \mu x_{1} x_{2}+4 x_{2}^{2}$, and hence $f(x)$ is $\mathbb{Z}_{p}$-equivalent to $g(x)$. Take $\tau_{g}=\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$, and if $f(x)=g(T x)$, take $\tau_{f}=T^{-1} \tau_{g} T$.

In general, by Lemma 4.3, we can assume that $f(x)=p^{\nu_{J_{1}}} f_{1}\left(x_{J_{1}}\right)+f_{2}\left(x_{J_{2}}\right)$, where $f_{1}$ is $\mathbb{Z}_{p}$-elementary and $J_{1} \sqcup J_{2}=\{1, \ldots, n\}$. Take $\tau_{f}$ such that $\tau_{f} x_{J}=\tau_{f_{1}} x_{J}$ and $\tau_{f} x_{J_{2}}=x_{J_{2}}$.

Lemma 5.9. Let $P$ be a finite set of prime numbers and let $Q(x)$ be an integral form of determinant $d \neq 0$. For each $p \in P$, let $f_{p}$ be a $\mathbb{Z}_{p}$-integral form of determinant $d$ which is $\mathbb{Z}_{p}$-equivalent to $Q(x)$. Then there exists a form $Q^{*}(x)$ which is integrally equivalent to $Q(x)$ and whose coefficients are arbitrarily close $p$-adically to $f_{p}$ for each $p \in P$.

Proof. By the hypothesis, for each $p \in P$, there exists a $p$-adic integral matrix $T_{p}$ such that $f_{p}(x)=Q\left(T_{p} x\right)$. Since $D_{f_{p}}=D_{f}=d \neq 0$ and $\left(\operatorname{det} T_{p}\right)^{2}=1$, the determinant
$\operatorname{det} T_{p}= \pm 1$. By Lemma 5.8, we may assume that $\operatorname{det} T_{p}=1$. It is sufficientyto ied construct a global integral matrix $T$ with $\operatorname{det} T=1$, which is $p$-adically atbitrarily close to each $T_{p}$. For such $T$, the quadratic form $Q^{*}(x)=Q(T x)$ satisfies the required condition.

Write $T_{p}=\left(c_{1}^{(p)}, \ldots, c_{n}^{(p)}\right)$ so that the columns $c_{i}^{(p)}, i=1, \ldots, n$, form a basis for $\mathbb{Z}_{p}^{n}$, and denote

$$
\|x\|_{p}=\max _{1 \leq k \leq n}\left|x_{k}\right|_{p}, \quad \text { for } x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{Z}_{p}^{n}
$$

For $J=1, \ldots, n$, we claim that for every $\varepsilon>0$, there exist $c_{1}, \ldots, c_{J} \in \mathbb{Z}^{n}$ extendable to a basis $c_{1}, \ldots, c_{j}, b_{J+1}, \ldots, b_{n}$ of $\mathbb{Z}^{n}$ such that

$$
\begin{equation*}
\left\|c_{j}-c_{j}^{(p)}\right\|_{p}<\varepsilon, \quad \text { for } p \in P, \quad 1 \leq j \leq J \tag{5.2.4}
\end{equation*}
$$

Then the lemma is proved by taking $J=n$. We prove the claim by induction on $J$.
For $J=1$ the claim follows from Lemma 5.5, 5.7 and 5.6. Suppose the claim holds for $J$ and $c_{1}, \ldots, c_{J}, b_{J+1}, \ldots, b_{n}$. Then

$$
c_{J+1}^{(p)}=l_{1}^{(p)} c_{1}+\ldots+l_{J}^{(p)} c_{J}+m_{J+1}^{(p)} b_{J+1}+\ldots+m_{n}^{(p)} b_{n},
$$

for some $l_{1}^{(p)}, \ldots, l_{J}^{(p)}, m_{J+1}^{(p)}, \ldots, m_{n}^{(p)} \in \mathbb{Z}_{p}$. By Lemma 5.5, we have $\max _{J<j \leq n}\left|m_{j}^{(p)}\right|_{p}=1$. By Chinese Remainder Theorem, we can find $l_{j} \in \mathbb{Z}$ such that $\left|l_{i}-l_{j}^{(p)}\right|_{p}<\varepsilon$ for all $j \leq J$ and $p \in P$. Assume that $J<n-1$ so that $K:=n-J>1$ and Lemma 5.6 applicable. It ensures that we can find $m_{j} \in \mathbb{Z}$ for $J<j \leq n$ satisfying $\left|m_{j}-m_{j}^{(p)}\right|_{p}<\varepsilon$ and $\operatorname{gcd}\left(m_{J+1}, \ldots, m_{n}\right)=1$. Put

$$
c_{J+1}=l_{1} c_{1}+\ldots+l_{J-} c_{J}+m_{J+1} b_{J+1}+\ldots+m_{n} b_{n} .
$$

Then $\left\|c_{J+1}-c_{J+1}^{(p)}\right\|_{p}<\epsilon$, for all $p \in P$. Moreover, since $\operatorname{gcd}\left(m_{J+1}, \ldots, m_{n}\right)=1$, by Lemma 5.6, $c_{1}, \ldots, c_{J+1}$ can be extended to a basis of $\mathbb{Z}^{n}$. Hence the claim holds for $J+1$.

Now suppose $J=n-1$. We may assume that $b_{n}$ is chosen to have

$$
\operatorname{det}\left(c_{1}, \ldots, c_{n-1}, b_{n}\right)=1
$$

and that $c_{1}, \ldots, c_{n-1}$ are sufficiently close to $c_{1}^{(p)}, \ldots, c_{n-1}^{(p)}$ so that $c_{1}, \ldots, c_{n-1},,_{n}^{(p)}$ forms a basis of $\mathbb{Z}_{p}^{n}$ and the determinant $\operatorname{det}_{p}:=\operatorname{det}\left(c_{1}, \ldots, c_{n-1}, c_{n}^{(p)}\right)$ sufficiently close to 1. Let $\widehat{c_{n}}{ }^{(p)}:=\frac{1}{\operatorname{det}_{p}} c_{n}^{(p)}$. Then $\operatorname{det}\left(c_{1}, \ldots, c_{n-1}, \widehat{c_{n}}{ }^{(p)}\right)=1$. Hence

$$
{\widehat{c_{n}}}^{(p)}=\widehat{l}_{1}^{(p)} c_{1}+\ldots+{\widehat{l_{n-1}}}^{(p)} c_{n-1}+b_{n}
$$

By Chinese Remainder Theorem, we pick $l_{1}, \ldots, l_{n-1}$ sufficiently close to $\widehat{l}_{1}^{(p)}, \ldots,{\widehat{l_{n-1}}}^{(p)}$, such that the vector

$$
c_{n}:=l_{1} c_{1}+\ldots+l_{n-1} c_{n-1}+b_{n},
$$

satisfies

$$
\left\|c_{n}-c_{n}^{(p)}\right\|_{p} \leq \max \left\{\left\|c_{n}-{\widehat{c_{n}}}^{(p)}\right\|_{p},\left\|{\widehat{c_{n}}}^{(p)}-c_{n}^{(p)}\right\|_{p}\right\}<\varepsilon
$$

The proof is complete.
Lemma 5.10. Let $f(x)$ and $g(x)$ be $\mathbb{Z}_{p}$-integral forms with Gram matrices $F$ and $G$ respectively. If

$$
\begin{equation*}
2 F \equiv 2 G \quad\left(\bmod p^{\delta+2 \lambda+1}\right) \tag{5.2.5}
\end{equation*}
$$

where $\delta$ is defined by $\left|D_{2 F}\right|_{p}=p^{-\delta}$ and $\lambda=\delta_{2, p}$ (the Kronecker symbol), then $f(x)$ and $g(x)$ are $\mathbb{Z}_{p}$-equivalent

Proof. Set

$$
\begin{equation*}
S:=\frac{1}{2} F^{-1}(G-F)=\frac{1}{2}(2 F)^{-1}(2 G-2 F) . \tag{5.2.6}
\end{equation*}
$$

Suppose $2 G \equiv 2 F\left(\bmod p^{\mu}\right)$ with $\mu \geq \delta+2 \lambda+1$.Then

$$
\begin{equation*}
2(\operatorname{det} 2 F) S=(\operatorname{adj} 2 F)(2 G-2 F) \equiv 0 \quad\left(\bmod p^{\mu}\right) \tag{5.2.7}
\end{equation*}
$$

Since adj $2 F$ is a $\mathbb{Z}_{p}$-integral matrix, $S \equiv 0\left(\bmod p^{\mu-\delta-\lambda}\right)$. Put $F_{1}=(I+S)^{t} F(I+S)$, so that

$$
\begin{equation*}
\left(G-F_{1}\right)=-S^{t} F S=-\frac{1}{2} S^{t}(G-F) . \tag{5.2.8}
\end{equation*}
$$

Then $\left.2 G-2 F_{1}\right) \equiv 0\left(\bmod p^{2 \mu-\delta-2 \lambda}\right)$. Since $(2 \mu-\delta-2 \lambda)>\mu$, we may replace $F$ by $F_{1}$ and repeat the procedure to complete the proof.

### 5.3 The proofs of Theorem 5.3

For a finite prime number, we say an element $x \in \mathbb{Z}_{p}^{n}$ is primitive if $\|x\|_{p}=1$, and $a \in \mathbb{Z}_{p}$ is primitively represented by a quadratic form $f(x)$ if there exists a primitive element $b \in \mathbb{Z}_{p}^{n}$ with $a=f(b)$. For $p=\infty, f(x)$ primitively represents $a$, if $f(x)$ represents $a$.

Theorem 5.11. Let $n \in \mathbb{N}$ and $d \neq 0$ be given. Let $f_{p}$ be a $n$-dimensional $\mathbb{Z}_{p}$-integral form of determinant d, for all $p=\infty, 2,3 \ldots$. Suppose there exists a rational form $R(x)$ which is $\mathbb{Q}_{p}$-equivalent to $f_{p}(x)$ for each $p$. Then there is a $\mathbb{Z}$-integral form $Q(x)$ which is $\mathbb{Z}_{p}$-equivalent to $f_{p}$ for each $p$, and $\mathbb{Q}$-equivalent to $R(x)$. Moreover, if each $f_{p}(x)$ primitively represents a rational integer $a$, then $Q(x)$ represents $a$.

Proof. We prove by induction on the dimension $n$. When $n=1, f_{p}(x)=d x_{1}^{2}$ for all $p$. The theorem holds for $Q(x)=d x_{1}^{2}$. Indeed, if $f_{p}(x)$ primitively represents $a$, for all $p$, then $d=a$, and hence $Q(1)=a$ as required. For the rest of the proof, assume that $n \geq 2$ and the theorem holds for $n-1$.

Next, we claim that there is an integer $c$ primitively represented by $f_{p}$ for all $p$ over $\mathbb{Z}_{p}$ and by $R(x)$ over $\mathbb{Q}$. To show it, let $b$ be any non-zero integer such that $b=R(x)$ where $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathbb{Q}^{n}$. Let $B=\left(B_{i j}\right)$ be the gram matrix of $R(x)$ and write $B_{i j}=U_{i j} / V_{i j}$ with $U_{i j}, V_{i j} \in Z, \operatorname{gcd}\left(U_{i j}, V_{i j}\right)=1$. Also write $x_{i}=s_{i} / t_{i}$, $s_{i}, t_{i} \in \mathbb{Z}$ with $\operatorname{gcd}\left(s_{i}, t_{i}\right)=1$. Let $P$ be the set of primes dividing $2 b d, U_{i j}, V_{i j}, s_{i}$ and $t_{i}$. For $p \notin P, R(x)$ is a $\mathbb{Z}_{p}$-integral form primitively representing $b$. Since $d$ is a $p$-adic unit, by Lemma 5.4, $f_{p}$ is $\mathbb{Z}_{p}$-equivalent to $R(x)$, and hence primitively represents $b$. For $p \in P$, since $f_{p}$ is $\mathbb{Q}_{p}$-equivalent to $R(x), f_{p}$ represents $b$ over $\mathbb{Q}_{p}$. Suppose $b=f_{p}\left(x_{p}\right)$ for $x_{p} \in \mathbb{Q}_{p}^{n}$. We may choose $\beta(p) \in \mathbb{Z}$ so that $\left\|p^{\beta(p)} x_{p}\right\|_{p}=1$. Then take $c=b \prod_{p \in P} p^{2 \beta(p)}$ for the claim.

Now let $a$ be an integer primitively represented by $f_{p}$ over $\mathbb{Z}_{p}$, for all $p$. By lemma 5.5, we may assume $f_{p}(1,0, \ldots, 0)=a$ for all p. Also, by Hasse-Minkowski Theorem, $R(x)$ represents $a$ over $\mathbb{Q}$, and we may assume $R(1,0, \ldots, 0)=a$. By completing the square, we have

$$
4 a f_{p}(x)=\left(2 a x_{1}+b_{2 p} x_{2}+\ldots+b_{n p} x_{n}\right)^{2}+f_{p}^{*}\left(x_{2}, \ldots, x_{n}\right)
$$

for some $b_{2 p}, \ldots, b_{n p} \in \mathbb{Z}_{p}$, where $f_{p}^{*}(x)$ is a $\mathbb{Z}_{p}$-integral form in $(n-1)$-variables with 1 determinant $d^{*}=a^{(n-2)} d$. Similarly,

$$
4 a R(x)=\left(2 a x_{1}+c_{2} x_{2}+\ldots+b_{n} x_{n}\right)^{2}+R^{*}\left(x_{2}, \ldots, x_{n}\right)
$$

for some $c_{2}, \ldots, c_{n} \in \mathbb{Q}$ and some rational form $R^{*}(x)$ in $(n-1)$-variables. Taking $y_{1}=2 a x_{1}+b_{2 p} x_{2}+\ldots+b_{n p} x_{n}$, we can write

$$
4 a f_{p}=y_{1}^{2}+f_{p}^{*}\left(x_{2}, \ldots, x_{n}\right)
$$

Similarly, by taking $z_{1}=2 a x_{1}+c_{2} x_{2}+\ldots+b_{n} x_{n}$, we can write

$$
4 a R=z_{1}^{2}+R^{*}\left(x_{2}, \ldots x_{n}\right)
$$

Since $y_{1}^{2}$ is $\mathbb{Q}_{p}$-equivalent to $z_{1}^{2}$, by Witt's Theorem, $f_{p}^{*}\left(x_{2}, \ldots, x_{n}\right)$ is $\mathbb{Q}_{p}$-equivalent to $R^{*}\left(x_{2}, \ldots, x_{n}\right)$. By the induction hypothesis, there exists a global integral form $Q^{*}\left(x_{2}, \ldots, x_{n}\right)$ of determinant $d^{*}=D_{f_{p}^{*}}=(4 a)^{n} d / a^{2}$ which is $\mathbb{Z}_{p}$-equivalent to $f_{p}^{*}\left(x_{2}, \ldots, x_{n}\right)$ for all p and $\mathbb{Q}$-equivalent to $R^{*}$. By Chinese Remainder Theorem, we can find $b_{2}, \ldots, b_{n} \in \mathbb{Z}$ arbitrarily close p-adically to $b_{2 p}, \ldots, b_{n p}$ for all $p \mid 2 a d$. Moreover, by Lemma 5.9, we can replace $Q^{*}$ by an equivalent form that is p-adically sufficiently close to $f_{p}^{*}$ for all $p \mid 2 a d$ so that the quadratic form $Q(x)$ defined by

$$
\begin{equation*}
4 a Q(x)=\left(2 a x_{1}+b_{2} x_{2}+\ldots+b_{n} x_{n}\right)^{2}+Q^{*}\left(x_{2}, \ldots, x_{n}\right) \tag{5.3.1}
\end{equation*}
$$

is an integral form. Since $D_{Q^{*}}=d^{*}, D_{Q}=d$. Obviously, $Q(1,0, \ldots, 0)=a$. Since $Q(x)$ is sufficiently close p-adically to $f_{p}$ for $p \mid 2 d$, by lemma $5.10, Q(x)$ and $f_{p}$ are $\mathbb{Z}_{p}$-equivalent for these p . For $p \nmid 2 d, p \neq \infty$, by Lemma 5.4, $Q(x)$ is automatically $\mathbb{Z}_{p}$-equivalent to $f_{p}$. For $p=\infty, Q(x)$ also $\mathbb{Z}_{p}$-equivalent to $f_{p}(x)$, because they have the same signature. Therefore, by Hasse-Minkowski Theorem, $Q(x)$ is $\mathbb{Q}$-equivalent to $R(x)$. The proof is complete.

The proof of Theorem 5.3. If $p \nmid a, Q(x)$ must primitively represents $a$. Otherwise, we assume $Q\left(x_{p}\right)=a$ over $\mathbb{Z}_{p}$. Then $\left\|x_{p} p^{\beta(p)}\right\|_{p}=1$ for some non-positive integers $\beta(p)$. We set $\hat{a}=a \prod_{p \mid a} p^{2 \beta(p)}$. Since

$$
\begin{equation*}
|\hat{a}|_{p}=|a| \cdot\left|\prod_{p \mid a} p^{2 \beta(p)}\right|=\left|f\left(x_{p}\right)\right| \cdot \frac{1}{p^{2 \beta(p)}} \leq 1 \tag{5.3.2}
\end{equation*}
$$

so $\hat{a} \in \mathbb{Z}$. It is clear that $Q(x)$ primitively represents $\hat{a}$. Theorem 5.11 applied to the case where $f_{p}=Q(x)$ and $R(x)=Q(x)$ says there exists in the genus of $Q(x)$ a $\mathbb{Z}$-integral form $Q^{*}(x)$ that globally represents $\hat{a}$. Assume $Q^{*}(\xi)=\hat{a}$ for somé $\xi \in \mathbb{Z}^{n}$. Then

$$
Q^{*}\left(\xi \cdot \prod_{p \mid a} p^{-\beta(p)}\right)=a
$$

(5.3.3)

### 5.4 Example and conclusion

In some cases, we can apply the arithmetic method discussed in this section and the section before to find the set of integers which are represented by a four-dimensional escalator $L_{4}$. Suppose $L_{4}$ is the escalator of a three-dimensional sublattice $L_{3}$, which is unique in its genus. By Theorem 5.3, such $L_{3}$ represents all integers which is locally represented over $\mathbb{Z}_{p}$. We can use the information of $L_{3}$ to show that the direct sum of $L_{3}$ with orthogonal complement in $L_{4}$ represents all sufficiently large integers $n$ locally represented by it. Thus, we only need to check the representability of small $n$ to determine the set of integers which are represented by the four-dimensional escalator $L_{4}$.

For example, we consider the quadratic form $x^{2}+y^{2}+z^{2}$ and its corresponding lattice $L_{3}$. It is well known that a natural number $n$ is represented by $L_{3}$ if and only if $n$ is not of the form $4^{a}(8 k+7)$ [5]. Now $L_{3}^{*}=L_{3}$, and hence $L_{4}$ is generated by $L_{3}$ together with a vector $f=\alpha e_{4}+f^{\prime}$, such that $\|f\|^{2}=7$ and $f^{\prime} \in \frac{1}{2} L_{3}$. If $f^{\prime} \in L_{3}$, then $L_{4}=L_{3} \oplus \alpha e_{4}=: L_{3} \oplus[m]$, where $m=\alpha^{2}=1,2,3,4,5,6$. If $f^{\prime} \notin L_{3}$, then $L_{4}$ contains $L_{3} \oplus 2 \alpha e_{4}:=L_{3} \oplus[m]$, where $m=4 \alpha^{2}<28$. In this case we can check that $4 \nmid m$ and $m \neq 21$. In both cases, we have $L_{4}$ conatins $L_{3} \oplus[m]$ with $m \leq 27$ and $m \not \equiv 0(\bmod 8)$. If $L_{4}$ were not universal and let $u$ be the smallest natural number missed by $L_{4}$, then $u$ is also missed by $L_{3}$. By the minimality, $u \equiv 7(\bmod 8)$. Since $L_{4}$ represents $7, u$ is at least 15 .

If $m \not \equiv 3,7(\bmod 8)$, then $u-m$ is not of the form $4^{a}(8 k+7)$. If $u-m>0$, then it is represented by $L_{3}$, and hence $u$ would be represented by $L_{4}$, which is absurd.

Therefore, $u \leq m$. Thus, if $m \leq 14$, then $L_{4}$ is universal. One can check that in this case all possible $L_{4}$ are universal.

If $m \equiv 3,7(\bmod 8)$, then $u-4 m$ is not of the form $4^{a}(8 k+7)$ and is represented by $L_{3}$ in case $u>4 m$. This implies $u \leq 4 m$. Thus, if $m=3$, then $L_{4}$ is universal. One can check that all possible $L_{4}$ in this case are also universal.

Among 34 three-dimensional basic escalator, 20 are unique in their genus. Moreover, the escalations of 17 of these 20 escalator lattices can be handled by this process. We can apply such arithmetic method to determine the sets of integers represented by 1658 of the 6560 basic four-dimensional escalators [6].

## 6 Analytic method

For basic 4-dimensional escalators on which the arithmetic method does not work well, one can apply the analytic method discussed in this section. The aim is to find for each form a lower bound of reasonable size such that every positive integer greater than this bound is represented by the given form if and only if it is locally represented by the form.

### 6.1 The theta function associated to 4-dimensional escalators

For a basic four-dimensional escalator $Q(x)$, we define the theta function associated with $Q(x)$ as

$$
\Theta_{Q}(z):=\sum_{x \in \mathbb{Z}^{4}} e^{2 \pi i Q(x) z}=\sum_{m \geq 0} r_{Q}(m) e^{2 \pi i m z}
$$

The theta function $\Theta_{Q}(z)$ is a modular form [6] of weight 2, and hence there exist an Eisenstein series $E(z)$ and a cusp form $f(x)$ such that

$$
\Theta_{Q}(z)=E(z)+f(z) .
$$

We observe that $r_{Q}(m)>0$ if and only if $Q(x)$ represents $m$. We have to compare the growth rates of Fourier coefficients $a_{E}(m)$ of $E(x)$ and $a_{f}(m)$ of $f(x)$ to establish an effective criterion for $r_{Q}(m)>0$ for $m$ sufficiently large by checking if $m$ is locally represented by $Q(x)$.

### 6.2 Fourier coefficients of Eisenstein series E(z)

Let $a_{E}(m)$ be as before. According to Segel's theory[8], we have

$$
a_{E}(m)=\prod_{p=\infty, 2,3, \ldots} \beta_{p}(m)
$$


where

$$
\beta_{p}(m)= \begin{cases}\lim _{k \rightarrow \infty} \frac{r_{p^{k}, Q}(m)}{p^{(n-1) k}}, & \text { if } p \neq \infty  \tag{6.2.2}\\ \mu\left(\left\{x \in \mathbb{R}^{4} \mid Q(x)=m\right\}\right), & \text { if } p=\infty\end{cases}
$$

and $\mu$ denotes the three-dimensional Lebesgue measure. By direct computation, we get

$$
\begin{equation*}
\beta_{\infty}(m)=\frac{2 \omega_{4} m}{\sqrt{\operatorname{det}(Q)}} \tag{6.2.3}
\end{equation*}
$$

where $\omega_{4}$ is the measure of four-dimensional unit ball. However, the evaluation of $a_{E}(m)$ needs more discussion.

Definition 6.1. The discriminant of a non-singular quadratic form $f(x)$ in a field $F$ is defined to be $d_{F} f:=D_{f} \cdot\left(F^{*}\right)^{2} \in F^{*} /\left(F^{*}\right)^{2}$.

Definition 6.2. A binary non-singular quadratic form $f(x)$ over a field $F$ is called a hyperbolic plane if $f \simeq 2 x y$, and an even-dimensional non-singular quadratic form is called hyperbolic if it is $F$-equivalent to a direct sum of some hyperbolic planes.

The proofs of the following two basic lemmas can be found in[5]
Lemma 6.3. Let $f(x)$ be a non-singular binary quadratic form in $\mathbb{F}_{p}, p \neq 2$, then the following are all equivalent:
(i) $f(x)$ is a hyperbolic plane.
(ii) $f(x)$ is isotropic, that is, there is some $a \in \mathbb{F}_{p}^{*}$ such that $f(a)=0$.
(iii) $d_{F} f=-1$.

Lemma 6.4. Let $f(x)$ and $g(x)$ be two non-singular quadratic forms over $\mathbb{F}_{p}, p \neq 2$, then $f(x)$ is $\mathbb{F}_{p}$-equivalent to $g(x)$ if and only if $\operatorname{dim} f(x)=\operatorname{dim} g(x)$ and $d_{\mathbb{F}_{p}} f=d_{\mathbb{F}_{p}} g$.

Now，we compute $r_{p, Q}(m)$ as follow．
Lemma 6．5．Let $Q(x)$ be positive integral quadratic form of dimension $2 k$ and het $p \neq 2$ be a prime number not dividing $N_{Q}$ ．For an integer $m \not \equiv 0(\bmod p)$ ，we haveg

$$
r_{p, Q}(m)=\frac{1}{p-1}\left(p^{2 k}-r_{p, Q}(0)\right) .
$$

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Proof．Put $Q^{(m)}(x)=m Q(x)$ ．Since the discriminants of $d_{\mathbb{F}_{p}} Q^{(m)}=d_{\mathbb{F}_{p}} Q$ ，Lemma 6.4 says $Q^{(m)}(x) \simeq Q(x)$ over $\mathbb{F}_{p}$ ．Therefore，the two sets $\left\{a \in \mathbb{F}_{p}^{2 k} \mid Q(a)=1\right\}$ and $\left\{a \in \mathbb{F}_{p}^{2 k} \mid Q(a)=m\right\}$ have an one－to－one correspondence，hence $r_{p, Q}(m)=r_{p, Q}(1)$ ． By counting，we obtain the desired equality．

Lemma 6．6．Let $Q(x)$ be a four－dimensional positive integral quadratic form with $p \neq 2$ and $p \nmid N_{Q}$ ．The following holds：
（i）If $Q(x)$ is hyperbolic，then $r_{p, Q}(m)=\left\{\begin{array}{cl}p^{3}+p(p-1), & \text { for } m \equiv 0(\bmod p) \text { ；} \\ p^{3}-p, & \text { for } m \not \equiv 0(\bmod p) .\end{array}\right.$
（ii）If $Q(x)$ is not hyperbolic，then $r_{p, Q}(m)=\left\{\begin{array}{cl}p^{3}-p(p-1), & \text { for } m \equiv 0(\bmod p) \text { ；} \\ p^{3}+p, & \text { for } m \not \equiv 0(\bmod p) \text { ．}\end{array}\right.$ Proof．To prove（i），write $Q=Q_{1} \oplus Q_{2}$ with $Q_{1} \simeq Q_{2} \simeq\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$ ．Then

$$
\begin{aligned}
r_{p, Q}(0) & =r_{p, Q_{1}}(0) \cdot r_{p, Q_{2}}(0)+\left(r_{p, Q_{1}}(1)\right)^{2} \cdot(p-1) \\
& =(2 p-1)^{2}+(p-1)^{3} \\
& =p^{3}+p(p-1) .
\end{aligned}
$$

By Lemma 6．5，$r_{p, Q}(m)=p^{3}-p$ if $m \not \equiv 0(\bmod p)$ ．
For（ii），by Lemma 6．4，we can write $Q=Q_{1} \oplus Q_{2}$ with $Q_{1} \simeq\left(\begin{array}{cc}0 & 1 \\ 1 & 0\end{array}\right)$ and $Q_{2}$ is not hyperbolic，and hence not isotropic，by Lemma 6．3．Thus，

$$
\begin{aligned}
r_{p, Q}(0) & =1 \cdot r_{p, Q_{1}}(0)+\left(p^{2}-1\right) r_{p, Q_{1}}(1) \\
& =1 \cdot(2 p-1)+\left(p^{2}-1\right)(p-1) \\
& =p^{3}-p(p-1) .
\end{aligned}
$$

By Lemma 6.5, $r_{p, Q}(m)=p^{3}+p$ if $m \not \equiv 0(\bmod p)$.

Let $\chi$ denote the Dirichlet character such that $\chi(p):=\left(\frac{D_{Q}}{p}\right)$, for $p \nmid N_{Q}$ and let $\mathscr{L}(s, \chi)$ denote the Dirichlet $L$-function. Namely,

$$
\mathscr{L}(s, \chi):=\prod_{p}(1-\chi(p))^{-s}, \quad s \in \mathbb{C}
$$

where in the product $p$ runs through all finite primes.
Corollary 6.7. For $p \nmid 2 N_{Q}$, we have

$$
r_{p, Q}=\left\{\begin{array}{cc}
p^{3}+p(p-1) \chi(p), & \text { if } m \equiv 0 \quad(\bmod p) \\
p^{3}-p \chi(p), & \text { if } m \not \equiv 0
\end{array}(\bmod p) .\right.
$$

Proof. As in the above proof, write $Q=Q_{1} \oplus Q_{2}$ so that $Q$ is hyperbolic if and only if $Q_{2} \simeq Q_{1}$. By Lemma 6.4, $Q_{2} \simeq Q_{1}$ if and only if $d_{\mathbb{F}_{p}}\left(Q_{2}\right)=d_{\mid F_{p}}\left(Q_{1}\right)$, or equivalently, $d_{\mathbb{F}_{p}}(Q)=-1$.

Definition 6.8. Let $p$ be a finite prime. We say that an integer $m$ is $p$-stable with respect to an integral form $Q(x)$, if it is locally represented by $Q(x)$ at $p$ and for each $k \gg 1$,

$$
r_{p^{k}, Q}^{G o o d}\left(m p^{2 \nu}\right)+r_{p^{k}, Q}^{B a d}\left(m p^{2 \nu}\right)
$$

is constant for all $\nu \geq 1$. Furthermore, we define $\operatorname{Stable}(m)$ to be the set containing all primes at which $m$ is stable.

Lemma 6.9. Let $Q(x)$ be a four dimensional integral form of level $N$ and let $m$ be an integer. Then $m$ is $p$-stable for every $p \nmid 2 N$.

Proof. Since $p \nmid D_{Q}, r_{p^{k}, Q}^{B a d}\left(m p^{2 \nu}\right)=0$ for each $\nu \geq 1$. Because $r_{p, Q}^{G o o d}\left(m p^{\nu}\right)=r_{p, Q}(0)$, the lemma is a consequence of Corollary 6.7 and Lemma 4.8.

We say that an integer $m$ is supported on a set $S$ of finite $\operatorname{primes}^{\text {if }} \operatorname{ord}_{p}(m)=0$ whenever $p \notin S$. This means all prime factors of $m$ lie in $S$.

Lemma 6.10. Let $Q(x)$ be a $2 n$-dimensional integral form of level $N$ and let mebe an integer. Suppose $m=\left(m^{\prime}\right)^{2} \cdot t$ and $p \in \operatorname{Stable}(t)$ with $\nu:=\operatorname{ord}_{p}\left(m^{\prime}\right) \geq 1$. . Then

$$
\frac{\beta(m)}{\beta(t)} \geq C_{p}(t)
$$

where

$$
C_{p}(t):=\min \left\{\frac{p^{n-2}}{p^{n-2}-1} \cdot \frac{\beta_{p}^{G o o d}\left(p^{2} t\right)+\beta_{p}^{\text {Bad }}\left(p^{2} t\right)}{\beta_{p}(t)}, 1\right\}
$$

Proof. For a $\mu \in \mathbb{Z}_{p}^{*}$, the form $Q(\mu x)$ is $\mathbb{Z}_{p}$-equivalent to $Q(x)$. Therefore, without loss of generality, we may assume $m^{\prime}=p^{\nu}$. For simplicity, write $r_{p^{k}}^{\operatorname{Good} \cup \operatorname{Bad}}(m)$ for $r_{p^{k}, Q}^{G o o d}(m)+r_{p^{k}, Q}^{B a d}(m)$. By Lemma 4.10, we have the recursive formula

$$
\begin{aligned}
r_{p^{k}}(m) & =r_{p^{k}}^{G o o d \cup B a d}(m)+r_{p^{k}, Q}^{\text {Zero }}(m) \\
& =r_{p^{k}}^{\text {Good } \cup \text { Bad }}(m)+p^{n} \cdot r_{p^{k-2}}\left(\frac{m}{p^{2}}\right),
\end{aligned}
$$

and hence

$$
r_{p^{k}}(m)=\sum_{i=0}^{\nu-1} p^{n i} r_{p^{k}}^{G o o d \cup B a d}\left(\frac{m}{p^{2 i}}\right)+p^{n \nu} r_{p^{k-2 \nu}, Q}(t)
$$

Put $K=\beta_{p}^{\text {Good } \cup \text { Bad }}\left(p^{2} t\right)$. Since $p \in \operatorname{Stable}(t)$, letting $k \rightarrow \infty$, we get

$$
\beta_{p}(m)=K \cdot \frac{\left(\frac{1}{p^{n-2}}\right)^{v}-1}{\frac{1}{p^{n-2}}-1}+\frac{\beta_{p}(t)}{p^{(n-2) v}} .
$$

Therefore, by setting

$$
L:=\frac{p^{n-2}}{p^{n-2}-1} \cdot \frac{K}{\beta_{p}(t)},
$$

we obtain

$$
\begin{aligned}
\frac{\beta_{p}(m)}{\beta_{p}(t)} & =\frac{1}{p^{(n-2) \nu}}+\frac{K}{\beta_{p}(t)} \cdot \frac{\left(\frac{1}{p^{n-2}}\right)^{\nu}-1}{\frac{1}{p^{n-2}}-1} \\
& =\frac{1}{p^{(n-2) \nu}}\left((1-L)+L \cdot p^{(n-2) \nu}\right)
\end{aligned}
$$

Therefore, we have the desired

$$
\frac{\beta_{p}(m)}{\beta_{p}(t)}= \begin{cases}\frac{1-L}{p^{(n-2) \nu}}+L \geq L, & \text { if } 1 \geq L \\ \frac{(L-1)\left(p^{(n-2) \nu}-1\right)}{p^{(n-2) \nu}}+1 \geq 1, & \text { if } L \geq 1\end{cases}
$$

Lemma 6.11. Let notation be as in Lemma 6.9. Suppose $n=2$ and $p \nmid 2 N$. If if $p \nmid t$ and $\chi(p)=-1$, then $C_{p}(t) \geq \frac{p}{p+1}$; otherwise $C_{p}(t)=1$.

Proof. By Lemma 5.4, we have $Q(x) \simeq x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+D_{Q} x_{4}^{2}$ over $\mathbb{Z}_{p}$. Consequently, $r_{p, Q}(a)=r_{p, Q}^{G o o d}(a)$ for every $a$. If $p \mid t$, then since $r_{p, Q}\left(p^{2} t\right)=r_{p, Q}(t)=r_{p, Q}(0)$, by Lemma 4.8, we have

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

Hence $C_{p}(t)=1$. If $p \nmid t$, then $r_{p, Q}\left(p^{2} t\right)=r_{P, Q}(0)$ and $r_{p, Q}(t)=r_{p, Q}(1)$. By Corollary 6.7 and Lemma 4.8, we have

$$
\begin{aligned}
\frac{p^{2}}{p^{2}-1} \cdot \frac{\beta_{p}\left(p^{2} T^{\prime}\right)}{\beta_{p}\left(T^{\prime}\right)} & =\frac{p^{2}}{p^{2}-1} \cdot \frac{p^{3}+p(p-1) \chi(p)}{p^{3}-p \chi(p)} \\
& = \begin{cases}\frac{p^{2}}{p^{2}-1} \cdot \frac{p^{3}+p^{2}-p}{p^{3}-p}>1, & \text { if } \chi(p)=1 ; \\
\frac{p^{2}}{p^{2}-1} \cdot \frac{p^{3}-p^{2}+p}{p^{3}+p}>\frac{p}{p+1}, & \text { if } \chi(p)=-1 .\end{cases}
\end{aligned}
$$

The following is the main theorem of this section. By Lemma 6.9, we can write $m=\left(m^{\prime}\right)^{2} \cdot t$ such that $m^{\prime}$ supported on Stable $(t)$ and $\operatorname{ord}_{p}(t) \leq 1$, for $p \nmid 2 N$. There could be more than one choice of $t^{\prime}$, let $T$ denote the set of all possible choices. Let $C_{p}(t)$ be as in Lemma 6.9 and denote

$$
\delta_{2}= \begin{cases}C_{2}(t), & \text { if } 2 \mid m^{\prime} \\ 1, & \text { otherwise }\end{cases}
$$

Theorem 6.12. Let $Q(x)$ be a four-dimensional basic escalator of level $N$ and determinant $D$ and let $m$ be locally represented by $Q(x)$. Then we have the following estimation

$$
a_{E}(m) \geq C_{E} \cdot m \cdot \prod_{p \mid m, p \nmid 2 N, \chi(p)=-1} \frac{p-1}{p+1},
$$

where

$$
C_{E}=\min _{t \in T}\left\{\frac{2 \omega_{4} \cdot D^{1 / 2}}{\mathscr{L}(2, \chi)} \cdot \prod_{p \mid 2 N} \frac{\beta_{p}(t)}{1-\frac{\chi(p)}{p^{2}}} \cdot \prod_{\substack{p \mid N \\ p \in \text { Stable }(t)}} C_{p}(t) \cdot \delta_{2}\right\} .
$$

Proof. Obviously, for $p \nmid m^{\prime}, \beta_{p}(m)=\beta_{p}\left(t^{\prime}\right)$. On the other hand, if $p \mid m^{\prime}$, with $\operatorname{ord}_{p}\left(m^{\prime}\right)=\nu \geq 1$, then since $p \in \operatorname{Stable}\left(t^{\prime}\right)$, by Lemma 6.9, $\frac{\beta(m)}{\beta(t)} \geq C_{p}(t)$. Thus by (6.2.1), (6.2.2) and (6.2.3),

$$
a_{E}(m) \geq a_{E}(t) \cdot\left(m^{\prime}\right)^{2} \cdot \prod_{p \mid m^{\prime}} C_{p}(t)
$$

Lemma 6.11 says

$$
\begin{equation*}
\prod_{p \mid m^{\prime}} C_{p}(t) \geq \prod_{\substack{p \nmid 2 N t \\ \chi(p)=-1}} \frac{p}{p+1} \cdot \prod_{\substack{p \mid N \\ p \in \operatorname{Stable}(t)}} C_{p}(t) \cdot \delta_{2} \tag{6.2.5}
\end{equation*}
$$

To estimate $a_{E}(t)$, we first note that by Lemma 6.6,

$$
\begin{aligned}
\prod_{p \nmid 2 N} \beta_{p}(t) & =\prod_{p \nmid 2 N t} \beta_{p}(t) \cdot \prod_{p \nmid 2 N, p \mid t} \beta_{p}(t) \\
& =\prod_{p \nmid 2 N t}\left(1-\frac{\chi(p)}{p^{2}}\right) \cdot \prod_{p \nmid 2 N, p \mid t} \frac{p^{3}+p(p-1) \chi(p)}{p^{3}} \\
& \geq \prod_{p \nmid 2 N t}\left(1-\frac{\chi(p)}{p^{2}}\right) \cdot \prod_{p \nmid 2 N, p \mid t} \frac{(p+\chi(p))\left(p^{2}-\chi(p)\right)}{p^{3}} \\
& =\prod_{p \nmid 2 N}\left(1-\frac{\chi(p)}{p^{2}}\right) \cdot \prod_{p \nmid 2 N, p \mid t} \frac{(p+\chi(p))\left(p^{2}-\chi(p)\right)}{p\left(p^{2}-\chi(p)\right)} \\
& \geq \prod_{p \nmid 2 N}\left(1-\frac{\chi(p)}{p^{2}}\right) \cdot \prod_{\substack{p \nmid 2 N, p \mid t, \chi(p)=-1}}\left(1-\frac{1}{p}\right) \\
& =\frac{1}{\mathscr{L}(2, \chi)} \cdot \prod_{p \mid 2 N}^{1-\frac{\chi(p)}{p^{2}}} \prod_{\substack{p \nmid 2 N, p \mid t, \chi(p)=-1}}\left(1-\frac{1}{p}\right) .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
a_{E}(t) & =\beta_{\infty}(t) \cdot \prod_{p \mid 2 N} \beta_{p}(t) \cdot \prod_{p \nmid 2 N} \beta_{p}(t) \\
& \geq \frac{\beta_{\infty}(t)}{\mathscr{L}(2, \chi)} \cdot \prod_{p \mid 2 N} \frac{\beta_{p}(t)}{1-\frac{\chi(p)}{p^{2}}} \cdot \prod_{\substack{p \nmid 2 N, p \mid t, \chi(p)=-1}}\left(1-\frac{1}{p}\right) .
\end{aligned}
$$

Thus, by (6.2.3), (6.2.4) and (6.2.5)

$$
a_{E}(m) \geq\left(\frac{2 \omega_{4} \cdot D^{1 / 2}}{\mathscr{L}(2, \chi)} \cdot \prod_{p \mid 2 N} \frac{\beta_{p}(t)}{1-\frac{\chi(p)}{p^{2}}} \cdot \prod_{p \mid N, p \in \text { stable }(t)} C_{p}(t) \cdot \delta_{2}\right) \cdot m \cdot \prod_{\substack{p \mid p, p+p+N, \chi(p)=-1}} \frac{p-1}{p+1} .
$$

Then take minimum for $t \in T$ and complete the proof.

### 6.3 Fourier coefficients of the cusp form $f(z)$

In order to estimate the growth rate of the Fourier coefficients $a_{f}(m)$ of the cusp. form $f(z)$, we apply general theory of Hecke eigenforms to write $f(z)$ as a linear combination of Hecke eigenforms

$$
f(z)=\sum_{i=1}^{r} \gamma_{i} g_{i}(z)
$$

where $\gamma_{i} \in \mathbb{C}$. Suppose we write $g_{i}(z)=\sum_{m \geq 0} b_{i}(m) e^{2 \pi i m z}$, then $a_{f}(m)=\sum_{i=1}^{r} \gamma_{i} b_{i}(m)$. By the Deligne's bound [7] on Hecke eigenforms, we obtain the estimation

$$
\left|b_{i}(m)\right| \leq \tau(m) \sqrt{m},
$$

where $\tau(m)$ is the number of positive divisors of $m$; hence,

$$
\begin{equation*}
\left|a_{f}(m)\right| \leq C_{f} \tau(m) \sqrt{m}, \tag{6.3.1}
\end{equation*}
$$

where $C_{f}=\sum_{i=1}^{r}\left|\gamma_{i}\right|$.

### 6.4 The criterion of representability

Let $Q(x)$ be a basic escalator. Combining the bounds of the Fourier coefficients of $E(z)$ and $f(z)$, we know that any number $m$ locally represented by $Q(x)$ and satisfying

$$
\begin{equation*}
\frac{\sqrt{m}}{\tau(m)} \prod_{\substack{p+N, p \mid m, \chi(p)=-1}} \frac{p-1}{p+1}>\frac{C_{f}}{C_{E}} \tag{6.4.1}
\end{equation*}
$$

is represented by $Q$.

Definition 6.13. For a given basic escalator $Q(x)$, an positive integer $m$ is called an eligible integer of the form if $m$ is locally represented by $Q(x)$ but fails to satisfy the inequity 6.4.1. Furthermore, define

$$
B(m)=\frac{\sqrt{m}}{\tau(m)} \prod_{\substack{p \nmid N, p \mid m, \chi(p)=-1}} \frac{p-1}{p+1} .
$$

In order to determine the set of integers represented by a basic quaternary escalator $Q(x)$, it is necessary to handle the eligible integers. Thus, we require more information about the function $B(m)$.

Lemma 6.14. $B(m)$ is a multiplicative function, and $B(p)>1$ for each prime $p>$
Moreover, for any positive $m$ and prime $p$,

$$
B\left(m p^{\nu}\right)>B(m) \text { if } \begin{cases}\nu \geq 1, & \text { for } p \geq 11  \tag{6.4.2}\\ \nu \geq 2, & \text { for } p=5,7 \\ \nu \geq 5, & \text { for } p=3 \\ \nu \geq 11, & \text { for } p=2\end{cases}
$$

Proof. First, $B(m)$ is obviously multiplicative and by direct computation, the inequality $B(p)>1$ is obtained for $p>7$. We need to check the inequality 6.4.2. Suppose $m=m_{1} p^{\nu_{1}}$ with $p \nmid m_{1}$, then

$$
\begin{aligned}
B\left(m p^{\nu}\right) & =\frac{p^{\frac{\nu}{2}} \sqrt{m_{1} p^{\nu_{1}}}}{\left(1+\nu+\nu_{1}\right) \tau\left(m_{1}\right)} \prod_{\substack{q \mid m_{1} p, q \uparrow N_{Q}, \chi(q)=-1}} \frac{q-1}{q+1} \\
& \geq\left(\frac{1+\nu_{1}}{1+\nu+\nu_{1}}\right)\left(\frac{p^{\frac{\nu}{2}}(p-1)}{p+1}\right) \frac{\sqrt{m_{1} p^{\nu_{1}}}}{\left(1+\nu_{1}\right) \tau\left(m_{1}\right)} \prod_{\substack{q \nmid N_{Q}, q \mid m_{1},, \chi(q)=-1}} \frac{q-1}{q+1} \\
& \geq \frac{1+\nu_{1}}{1+\nu+\nu_{1}} \cdot \frac{p^{\frac{\nu}{2}}(p-1)}{p+1} B(m) \\
& \geq \frac{p^{\frac{\nu}{2}}}{1+\nu} \cdot \frac{p-1}{p+1} B(m) .
\end{aligned}
$$

Therefore, the proof is completed by solving the inequality

$$
\frac{p^{\frac{\nu}{2}}}{1+\nu} \cdot \frac{p-1}{p+1}>1 .
$$

Corollary 6.15. For a basic escalator $Q(x)$, the set of all eligible integers of $Q(x)$ is finite.

Proof. Since $B(m)$ is a multiplicative function and $B(p) \leq \frac{C_{f}}{C_{E}}$ holds only for finitely many $p$, the assertion follows from Lemma 6.14 and its proof.

Since the number of eligible integers is finite (although maybe very large), the set of positive integers represented by the form $Q(x)$ could be determined by computation [ $6, \S 4.3]$. In the Bhargave and Hanke's paper [6], they present a quaternary basic form $Q(\tilde{x})=x^{2}+2 y^{2}+4 z^{2}+31 w^{2}+y z-y w+3 z w$. The level $N_{Q}=3774, \chi(p)=\left(\frac{104}{p}\right), C_{E}=$ $\frac{36}{125}$, and $C_{f} \approx 2331.99$, giving the largest bound $\frac{C_{f}}{C_{E}}<8100.65$ of 6560 quaternaries. In this case, there are 28 billion square-free eligible integers.

## 7 Proofs of the main theorems

### 7.1 Summary

By checking with machine using the above arithmetic and analytic methods, we classify all basic four-dimensional escalator lattices $L$ into three types:

Type I: $L$ is universal.
Type II : $L$ is not universal but miss at most three positive integers, every missed number is a critical integer.

Type III : Otherwise.
We find that among all basic four-dimensional escalator lattices, there are 6402 lattices of Type I, 153 one of type II, and 5 ones of Type III.

### 7.2 The $10-14$ switch

The lattice escalators of Type III all have truants 14 and are given by

$$
\left(\begin{array}{cccc}
1 & 0 & -1 / 2 & -3 \\
0 & 2 & 1 & 0 \\
-1 / 2 & 1 & 5 & 1 \\
-3 & 0 & 1 & 10
\end{array}\right),\left(\begin{array}{cccc}
1 & 0 & -1 / 2 & -2 \\
0 & 2 & 1 & -2 \\
-1 / 2 & 1 & 5 & 3 \\
-2 & -2 & 3 & 10
\end{array}\right),\left(\begin{array}{cccc}
1 & 0 & -1 / 2 & -2 \\
0 & 2 & 1 & -2 \\
-1 / 2 & 1 & 5 & 1 \\
-2 & -2 & 1 & 10
\end{array}\right)
$$

$$
\left(\begin{array}{cccc}
1 & 0 & -1 / 2 & -1 \\
0 & 2 & 1 & 0 \\
-1 / 2 & 1 & 5 & 3 \\
-1 & 0 & 3 & 10
\end{array}\right) \text { and }\left(\begin{array}{cccc}
1 & 0 & -1 / 2 & -1 \\
0 & 2 & 1 & 0 \\
-1 / 2 & 1 & 5 & 2 \\
-1 & 0 & 2 & 10
\end{array}\right)
$$

These five lattices are escalated from the three-dimensional escalator $L_{3}$ given by $\left(\begin{array}{ccc}1 & 0 & 1 / 2 \\ 0 & 2 & 1 \\ 1 / 2 & 1 & 5\end{array}\right)$, whose truant is 10 .

Now instead of escalating $L_{3}$ by adding a new vector of square norm 10, we add a vector of square norm 14 to it. There are 330 lattices obtained in this way, they are all 4-dimensional. Call them the auxiliary quaternaries.

Lemma 7.1. The escalator of a basic lattice $L$ of Type III contains an auxiliary quaternary.

Proof. Let $M=L+\mathbb{Z} e$, with $\|e\|^{2}=14$. Then $L^{\prime}=L_{3}+\mathbb{Z} e$ is an auxiliary quaternary contained in $M$.

One finds that 226 of these 330 auxiliary quaternaries recurred in the 6555 basic escalator lattices of Type I and II. By machine computation, we can classify the remaining 104 forms $L$ into three types:

Type I: $L$ is universal.
Type II: $L$ not universal but miss at most three positive integers, every missed number is a critical integer.
Type IV : $L$ represents all integers except for those of forms $10 n^{2}$ and $13 n^{2}$.

### 7.3 The proof of Theorem 1

So far, all basic quaternary escalators and auxiliary quaternaries are classified.

Definition 7.2. Let $C_{290}$ denote the collection of all integers missed by at least one of all basic quaternary escalators and auxiliary quaternaries.

The result of computation tells the following:
Lemma 7.3. $C_{290}$ is exactly the set of all critical integers listed in Theorem 1.
Lemma 7.4. The zero lattice can be escalated at most 7 times.

Proof. If a basic 4-dimensional escalator is of Type II, then it becomes universal after at most three escalating. If it is of Type III, then any of its escalator $L^{\prime}$ contains an auxiliary quaternary. If this auxiliary quaternary, which misses 10 , is of Type II, then $L^{\prime}$, which represents 10 , misses at most two integers, and hence becomes universal after at most two escalations. If this auxiliary quaternary, which represents 10 , is of Type II, then $L^{\prime}$ is exactly this auxiliary, and hence becomes universal at most 3 escalations. If the auxiliary quaternary is of Type IV, then it is universal.

Proof of Theorem 1. First, let $Q$ be a positive-definite integral quadratic form, then there exists a maximal sequence of escalator lattices

$$
\{0\} \subset L_{1} \subset L_{2} \subset \ldots \subset L_{k}
$$

within $L(Q)(\mathrm{k} \leq 7)$. If $L(Q)$ is universal, then $L_{k}$ is also universal. Conversely, suppose $L(Q)$ is not universal, then $\operatorname{truant}(L)=\operatorname{truant}\left(L_{k}\right)$; otherwise, if $\operatorname{truant}(L) \geq$ truant $\left(L_{k}\right)$, then one can obtain a escalator $L_{k+1}$ within $L(Q)$ properly containing $L_{k}$.

### 7.4 The proof of Theorem 2

Now we complete the prove of Theorem 2.
The proof of Theorem 2. First, we know that all critical integers except 203 and 290 arise as the truants of $1,2,3,4$ dimensional escalator lattices. For 203 and 290, we consider the special escalator lattice

$$
L_{145}:=\left(\begin{array}{ccc}
1 & 0 & -1 / 2 \\
0 & 2 & -1 / 2 \\
-1 / 2 & -1 / 2 & 4
\end{array}\right) \oplus[29] .
$$

$L_{145}$ has 145 as its truant, and its escalation

$$
\left(\begin{array}{ccc}
1 & 0 & -1 / 2 \\
0 & 2 & -1 / 2 \\
-1 / 2 & -1 / 2 & 4
\end{array}\right) \oplus[29] \oplus[145]
$$

has the truant 290. Remarkably, $L_{145}$ is the only basic four-dimensional escalator lattice dose not represent 203. We set a lattice $L_{203}:=L_{145} \oplus$ [58]. Since the only square multiples of 58 less than 203 are 0 and 58 itself, $L_{203}$ represents

$$
\{0,1, \ldots, 144\} \bigcup\{0+58,1+58, \ldots, 144+58\}
$$

but dose not represent 203. This is, 203 is the truant of $L_{203}$. Since $L_{203}$ represents 145 , we may pick a vector $v$ of norm 145 in the lattice $L_{203}$, and set $L_{203}^{\prime}$ as the lattice generated by $L_{145}$ and the vector $v$. Because $L_{145}$ cannot represent $145, v$ is form of $(*, *, *, *, 1)$, which implies $L_{203}=L_{203}^{\prime}$.

So far, we have proven that each critical integer is the truant of some escalator lattice. For any critical integer $m$, let $L$ be an escalator lattice whose truant is $m$, and define $L_{m}=L \oplus[m+1]^{\oplus 4} \oplus[m x+1]$. We claim that $L_{m}$ represents all nonzero integers except for $m$.

By Lagrange's four theorem, $[m+1]^{\oplus 4}$ represents all integers of the form $(m+1) n$. Every integer $a$ can be expressed as $a=(m+1) q+r$, where $0 \leq r \leq m$. Therefore, $L_{m}$ represent all nonzero integers except for $m$, since $L$ represents all integers in the interval $[0, m-1]$ while $[2 m+1]$ represents $2 m+1$. Thus, the proof is complete.

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