# COUNTING DIOPHANTINE APPROXIMATIONS 

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Dedicated to Professor Władysław Narkiewicz


#### Abstract

A recent development of the Davenport-Heilbronn method for diophantine inequalities is reexamined, and then applied to a class of problems in diophantine approximation. Among other things, an asymptotic formula is obtained for the number of solutions of the simultaneous inequalities $\left|n_{j}-\lambda_{j} n_{0}\right|<\varepsilon$ with square-free $n_{j} \in[1, N]$, whenever the positive real numbers $\lambda_{1}, \ldots, \lambda_{r}$ and 1 are linearly independent over the rationals. Keywords: Davenport-Heilbronn method, diophantine approximation, square-free numbers.


## 1. An illustrative example

The most classical problem in diophantine approximation concerns rational numbers near a given set of positive real numbers, say $\lambda_{1}, \ldots, \lambda_{r}$. One intends to solve the inequalities

$$
\begin{equation*}
\left|n_{j}-\lambda_{j} n_{0}\right|<\varepsilon \quad(1 \leq j \leq r) \tag{1.1}
\end{equation*}
$$

in natural numbers $n_{i}$. A prominent example is Dirichlet's theorem: If $\varepsilon=N^{-1 / r}$, then (1.1) has a solution with $1 \leq n_{i} \leq N$. Many variants of Dirichlet's theorem have been considered, with the variables $n_{j}$ restricted to special sets, and $\varepsilon$ a function of the "denominator" $n_{0}$, tending to zero as $n_{0}$ increases. The focus in these researches has mostly been on the existence of solutions, with $\varepsilon$ as small as possible. We mention a recent result of Dietmann [9]: For $\theta<\frac{2}{3}$ and subject to a mild condition on $\lambda_{1}, \ldots, \lambda_{r}$, there are infinitely many square-free numbers $n_{0}$ such that the inequalities $\left|n_{j}-\lambda_{j} n_{0}\right|<n_{0}^{-\theta / r}$ have a solution in integers $n_{1}, \ldots, n_{r}$. When $r=1$, Heath-Brown [14] preceded this with the stronger conclusion that for irrational $\lambda>0$ and any $\theta<\frac{2}{3}$ there are infinitely many pairs $n_{0}, n_{1}$ of square-free numbers with $\left|n_{1}-\lambda n_{0}\right|<n_{0}^{-\theta}$. Heath-Brown's work is among the few examples where not only the denominator $n_{0}$ but also the numerators $n_{j}$ with $1 \leq j \leq r$ in (1.1) fall into specified sets.

In the present communication we address the problem of counting solutions of (1.1), with the variables $n_{j} \in[1, N]$ restricted to a fairly general class of sets. When specialized to the set of square-free numbers, our result reads as follows.

Theorem 1. Let $\lambda_{1}, \ldots, \lambda_{r}$ denote positive real numbers such that $1, \lambda_{1}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. Let $Z_{\varepsilon}(N)$ denote the number of solutions of (1.1) in square-free numbers $n_{j}$ with $1 \leq n_{j} \leq N$. Then there exists a positive real number $C$ and a function $T(N)$ that increases to infinity as $N$ tends to infinity, such that

$$
Z_{\varepsilon}(N)=C N \varepsilon^{r}+O\left(N T(N)^{-1}\right)
$$

holds uniformly in $0<\varepsilon \leq 1$.
Note that $T(N)$ depends on $\lambda_{1}, \ldots, \lambda_{r}$. This is necessarily so, see $\S 8$.
Some of the tools in this paper are easier to describe in the simplest case $r=1$. Even then Theorem 1 seems to be new but may well be part of the folklore. Traditionally one would search for small values of $n_{1}-\lambda n_{0}$, with $n_{0}$ in some special set $\mathcal{N}_{0}$, with the aid of exponential sums. Methods of Weyl and Vinogradov are standard nowadays in this family of problems, but are set up in a way that deny access to the numerators $n_{1}$. Thus we are forced to use more subtle methods originally invented by Davenport and Heilbronn [8]. This amounts to a treatment of the inequality $\left|n_{1}-\lambda n_{0}\right|<\varepsilon$, with $n_{j}$ restricted to a special set $\mathcal{N}_{j} \subset[1, N]$, very much like a binary additive problem. One writes $\Upsilon(\alpha)=\max (0,1-|\alpha| / \varepsilon)$ as a Fourier transform of a function $K(\alpha)$ (see (5.2)), and takes the formula

$$
\begin{equation*}
\sum_{\substack{n_{0} \in \mathcal{N}_{0} \\ n_{1} \in \mathcal{N}_{1}}} \Upsilon\left(n_{1}-\lambda n_{0}\right)=\int_{-\infty}^{\infty} F_{1}(\alpha) F_{0}(-\lambda \alpha) K(\alpha) d \alpha \tag{1.2}
\end{equation*}
$$

with

$$
F_{j}(\alpha)=\sum_{n_{j} \in \mathcal{N}_{j}} e(\alpha n)
$$

as the starting point. A suitable neighbourhood of the origin contributes a term of size $\varepsilon N^{-1} \# \mathcal{N}_{0} \# \mathcal{N}_{1}$ to (1.2) whereas large $\alpha$ make a negligible contribution due to the decay of $K(\alpha)$. Only very recently Bentkus and Götze [2] developed a method that shows

$$
\begin{equation*}
F_{1}(\alpha) F_{0}(-\lambda \alpha)=o\left(\# \mathcal{N}_{0} \# \mathcal{N}_{1}\right) \tag{1.3}
\end{equation*}
$$

uniformly for a rather large range of $\alpha$, provided $\lambda$ is irrational and one has some Weyl type bounds for $F_{j}(\alpha)$ at hand. Freeman [11] reworked their ideas in a context closer to ours, and we refer to the underlying idea as the Bentkus-Götze-Freeman device. No existing version of it seems to suit our needs. Therefore, in §3, we again review the basic principles and demonstrate in Theorem 3 that the essence of the Bentkus-Götze-Freeman device can be freed from any reference to exponential
sums. In the new form, the device is ready for immediate adaption to any concrete situation of the shape (1.3), and multi-dimensional variants thereof.

As is familiar to practioneers in the additive theory of numbers, inequalities like (1.3) alone are insufficient to extract an asymptotic formula from the integral (1.2). In fact one needs to break square root cancellation for an individual exponential sum $F_{j}(\alpha)$, at least in mean square over suitable minor arcs. In [5] we classified the sets $\mathcal{N}_{j}$ where this is feasible. It turns out that multiplicative sequences like the square-free numbers featured in Theorem 1, form only a subclass of those sequences where mean square estimates of the required strength are indeed available. We summarize the consequences of [5] in §4, and use this and our version of the Bentkus-Götze-Freeman device to derive an asymptotic formula for a multidimensional analogue of (1.2) in $\S \S 5$ and 6 . The main theorem is formulated in $\S 7$, and contains Theorem 1 as a special case. In a sense, $\S \S 5-7$ may be viewed as a variation on the work in [5], for diophantine approximations rather than equations. However, the convolution integral treatment in $\S 6$ is a major digression from the iterative approach in [5] that should also prove useful elsewhere.

The notation used in this paper is mostly standard and otherwise introduced when appropriate. Frequently vector notation is used, and printed in bold. The dimension is often, but not always $r$. For $\boldsymbol{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{r}\right) \in \mathbb{R}^{r}$ we write $\|\boldsymbol{\alpha}\|=$ $\max \left|\alpha_{j}\right|$, and $\boldsymbol{\alpha} \boldsymbol{\beta}=\alpha_{1} \beta_{1}+\ldots+\alpha_{r} \beta_{r}$ is the usual scalar product.

## 2. Normalized matrices

The Fourier transform method of Davenport and Heilbronn has been used extensively to study linear systems

$$
\begin{equation*}
\left|\lambda_{i 1} x_{1}+\lambda_{i 2} x_{2}+\ldots+\lambda_{i s} x_{s}-\mu_{i}\right|<\varepsilon \quad(1 \leq i \leq r), \tag{2.1}
\end{equation*}
$$

where $\varepsilon$ is a (small) positive real number, $\Lambda=\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq r \\ 1 \leq j \leq s}}$ is a real matrix of rank $r$, $1 \leq j \leq s$ the $\mu_{i}$ are real numbers, and the variables $x_{j}$ vary over some specific set of integers (typically, the $k$-th powers, the primes, or as in our Theorem 1, the square-free numbers). If for some $i$ all $\lambda_{i j}$ are integers, then for small $\varepsilon$ the corresponding inequality reduces to an equation. It may also be that equations are "hidden" in the system (2.1), in that an inequality has integral coefficients only after suitable linear combinations of the forms in (2.1) have replaced the original ones. Freeman [12], section 2, has discussed this theme in detail, and we briefly describe his ideas here. Let $z(\Lambda)$ denote the number of rows of $\Lambda$ that have integer entries only, and let

$$
R=R(\Lambda)=\max _{\mathrm{A} \in G L_{r}(\mathbb{R})} z(\mathrm{~A} \Lambda) .
$$

We call $R$ the integral rank of $\Lambda$; this is the number of "hidden equations" in (2.1), and we proceed to make them "visible". There is then a non-singular matrix A such that $\mathrm{A} \Lambda$ has the first $R$ rows integral. This partially motivates the following definition.

A real matrix $\Lambda=\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq r \\ 1 \leq j \leq s}}$ of rank $r$ and integral rank $R$ is called normalized if the first $R$ rows are integral, and the submatrices $\left(\lambda_{i j}\right)_{1 \leq i, j \leq r}$ and $\left(\lambda_{i j}\right)_{1 \leq i, j \leq R}$ are non-singular.

Note that a matrix of rank $r$ has an $r \times r$-submatrix that is non-singular. Hence, a renumbering of the columns of $\Lambda$ suffices to make $\left(\lambda_{i j}\right)_{1 \leq i, j \leq r}$ non-singular. This corresponds to a renumbering of variables in (2.1), and does not significantly alter the original problem. Now consider the first $R$ rows of the first columns, $\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq R \\ 1 \leq j \leq r}}$. This matrix must contain an $R \times R$-submatrix that is non-singular, so another renumbering of the first $r$ columns guarantees that $\left(\lambda_{i j}\right)_{1 \leq i, j \leq r}$ is nonsingular. Thus, we have proved the following lemma.

Lemma 2. Let $\Lambda=\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq r \\ 1 \leq j \leq s}}$ be a real matrix of rank $r$ and integral rank $R$. Then, after a suitable reordering of the columns of $\Lambda$, there is a non-singular $r \times r$-matrix A such that A $\Lambda$ is normalized.

Let $\Lambda$ be a normalized matrix with rows $\mathbf{c}_{1}, \ldots, \mathbf{c}_{r}$, and suppose that $\alpha_{1}, \ldots$, $\alpha_{r} \in \mathbb{R}$ are chosen such that

$$
\alpha_{1} \mathbf{c}_{1}+\ldots+\alpha_{r} \mathbf{c}_{r} \in \mathbb{Q}^{s} .
$$

Then

$$
\alpha_{R+1}=\ldots=\alpha_{r}=0,
$$

for otherwise the integral rank would be at least $R+1$. We shall use this fact below, in the proof of Theorem 3 .

We end this section by considering briefly the special case where $s=r+1$. It is that case that is relevant later. Following [1], a real $r \times(r+1)$-matrix $\Lambda$ is called highly non-singular if any $r \times r$-submatrix is non-singular. For such a matrix there exists $\mathrm{A} \in G L_{r}(\mathbb{R})$ such that

$$
\mathrm{A} \Lambda=\left(\begin{array}{cccc}
-\lambda_{1} & 1 & & 0  \tag{2.2}\\
\vdots & & \ddots & \\
-\lambda_{r} & 0 & & 1
\end{array}\right)
$$

with non-zero real numbers $\lambda_{1}, \ldots, \lambda_{r}$. Note $R(\Lambda)=R(\mathrm{~A} \Lambda)$, and that we may still apply row operations to the right hand side of (2.2) to determine the integral rank of $\Lambda$. In particular, it is immediate that $R=0$ if and only if the numbers $1, \lambda_{1}, \lambda_{2}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. This will be needed in $\S 6$.

The matrix $\Lambda$ is positive if the system of equations $\Lambda \mathbf{x}=0$ has a solution with all $x_{i}>0$. A matrix is positive if and only if all $\lambda_{i}$ in (2.2) are positive.

## 3. A technical estimate

In this section we set the scene for our version of the Bentkus-Götze-Freeman device. Some notation is required before we can formulate the result. Let $N$ denote the main parameter, and suppose that

$$
\begin{equation*}
1 \leq Q \leq \frac{1}{2} \sqrt{N} \tag{3.1}
\end{equation*}
$$

Then, the intervals

$$
\mathfrak{M}_{q, a}=\{\alpha \in \mathbb{R}:|q \alpha-a| \leq Q / N\}
$$

with $1 \leq q \leq Q, a \in \mathbb{Z}$ and $(a ; q)=1$ are pairwise disjoint, and we denote their union by $\mathfrak{M}(Q)$. Choose a function $H:[1, \infty) \rightarrow(0, \infty)$ that decreases monotonically to zero as its argument tends to $\infty$. Define a function $\Xi: \mathbb{R} \rightarrow[0, \infty)$ by

$$
\Xi(\alpha)= \begin{cases}0 & \text { for } \alpha \notin \mathfrak{M}(Q), \\ H(q+N|q \alpha-a|) & \text { for } \alpha \in \mathfrak{M}_{q, a} \cap \mathfrak{M}(Q) .\end{cases}
$$

We call $\Xi$ the weight on $\mathfrak{M}(Q)$ associated with $H$.
Now let $\Lambda=\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq r \\ 1 \leq j \leq s}}$ be a matrix with real entries, of rank $r$ and of integral rank $R \leq r$. Suppose that $\Lambda$ is normalized in the sense of section 2 . Fix a function $U(N)$ with $U(N) \rightarrow \infty$ as $N \rightarrow \inf$, and define

$$
\mathcal{M}=\left\{\boldsymbol{\alpha} \in \mathbb{R}^{r}: \alpha_{i} \in \mathfrak{M}(U)(1 \leq i \leq R),\left|\alpha_{i}\right| \leq U / N(R<i \leq r)\right\}
$$

The minor arcs in our set-up are then defined by

$$
\mathfrak{m}(T)=\left\{\boldsymbol{\alpha} \in \mathbb{R}^{r}: U(N) / N \leq\|\boldsymbol{\alpha}\| \leq T\right\} \backslash \mathcal{M} .
$$

We are now in a position to announce the main result of this section.
Theorem 3. Let $Q$ be a function of $N$ satisfying (3.1). Let $\Xi$ be a weight on $\mathfrak{M}(Q)$ associated with a decreasing function $H$, as above. Let $U(N)$ be a function with $U(N) \rightarrow \infty$ as $N \rightarrow \infty$. Let $\Lambda$ be a normalized $r \times s$-matrix with columns $\boldsymbol{\lambda}_{1}, \ldots, \boldsymbol{\lambda}_{s}$. Then, there exists a function $T(N) \rightarrow \infty$ with $N \rightarrow \infty$ and $T(N) \leq$ $U(N)$, and such that

$$
\begin{equation*}
\sup _{\boldsymbol{\alpha} \in \mathfrak{m}(T(N))} \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{1}\right) \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{2}\right) \ldots \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{s}\right) \ll T(N)^{-1} . \tag{3.2}
\end{equation*}
$$

This is the promised version of the Bentkus-Götze-Freeman device, freed from explicit references to exponential sums. It may be useful to compare it with the earlier versions by Freeman [10, 11] and Wooley [18] or Parsell [15]. These authors study Weyl sums for $k$-th powers, or for primes, defined by

$$
f(\alpha)=\sum_{x \leq N^{1 / k}} e\left(\alpha x^{k}\right), \quad g(\alpha)=\sum_{p \leq N}(\log p) e(\alpha p) .
$$

It is then shown that for irrational $\lambda \in \mathbb{R}$ and some suitable functions $U(N)$ and $T(N)$ as above, one has

$$
\begin{gather*}
\sup _{U(N) / N<|\alpha| \leq T(N)}|f(\alpha) f(\lambda \alpha)|=o\left(N^{2 / k}\right),  \tag{3.3}\\
\sup _{U(N) / N<|\alpha| \leq T(N)}|g(\alpha) g(\lambda \alpha)|=o\left(N^{2}\right) . \tag{3.4}
\end{gather*}
$$

These bounds are also available from Theorem 3, combined with some version of Weyl's inequality. In fact, one readily finds that for $\alpha \in \mathfrak{M}\left(N^{1 / k}\right)$ one has

$$
f(\alpha) \ll N^{1 / k}(q+N|q \alpha-a|)^{-1 / k}
$$

(see Vaughan [17], Theorem 4.1 and Lemma 4.6, for example), and when $\alpha \notin$ $\mathfrak{M}\left(N^{1 / k}\right)$ Weyl's inequality ([17], Lemma 2.4) gives $f(\alpha) \ll N^{\frac{1}{k}-\delta}$ for some $\delta>0$; the actual value of $\delta$ is of no relevance here. We take $H(t)=t^{-1 / k}$ to define the weight $\Xi$ on $\mathfrak{M}\left(N^{1 / k}\right)$. The two bounds for $f(\alpha)$ now combine to

$$
N^{-1 / k} f(\alpha) \ll \Xi(\alpha)+N^{-\delta}
$$

and (3.3) follows from Theorem 3 with $r=1, s=2$ and $\Lambda=(1, \lambda)$. A very similar argument can be used to establish (3.4). Multi-dimensional versions of (3.3) for smooth Weyl sums occur in work of Freeman. His important Lemma 5.5 of [12] can also be deduced from our Theorem 3 in the way indicated above. It is hoped that Theorem 3 serves as an easy-to-use reference in any future application of the Bentkus-Götze-Freeman device. In delicate uses of the circle method, ordinary and smooth Weyl sums often occur simultaneously, and in such situations, Theorem 3 is still easily applicable.

One might ask for a further generalization of Theorem 3, with different weights $\Xi_{1}, \ldots, \Xi_{s}$ on the linear forms $\boldsymbol{\alpha} \boldsymbol{\lambda}_{1}, \ldots, \boldsymbol{\alpha} \boldsymbol{\lambda}_{s}$. It will be clear from the proof below that such a generalization is indeed possible. However, as we shall see later, Theorem 3 is sufficiently flexible as it stands.

In this paper Theorem 3 is used only with $R=0$. In this case the sets $\mathcal{M}$ and $\mathfrak{m}(T)$ take the simple shape $\|\boldsymbol{\alpha}\| \leq U$ and $U<\|\boldsymbol{\alpha}\| \leq T$, respectively, and the proof considerably simplifies. However, in most applications the full strength of Theorem 3 is required. For an example and further comments on the matter we refer to our recent work with Wooley [7].

Proof of Theorem 3. As a first step, we show that for any fixed $T \geq 1$, one has

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \sup _{\boldsymbol{\alpha} \in \mathfrak{m}(T)} \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{1}\right) \ldots \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{s}\right)=0 \tag{3.5}
\end{equation*}
$$

We prove (3.5) by contradiction. If (3.5) is false, then there exist a real number $\varepsilon$ with $0<\varepsilon<1$, an increasing sequence of real numbers $N_{\nu}$ with $N_{\nu} \rightarrow \infty$ as $\nu \rightarrow \infty$, and $\boldsymbol{\alpha}_{\nu} \in \mathfrak{m}(T)$ with

$$
\Xi\left(\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{1}\right) \Xi\left(\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{2}\right) \ldots \Xi\left(\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{s}\right)>\varepsilon
$$

But $\Xi(\alpha) \leq 1$ whence $\Xi\left(\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j}\right)>\varepsilon$ for $1 \leq j \leq s$. In particular, $\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j} \in \mathfrak{M}(Q)$, so that there are unique coprime integers $q_{\nu j}, a_{\nu j}$ with $\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j} \in \mathfrak{M}_{q_{\nu j}, a_{\nu j}}$. Then

$$
H\left(q_{\nu j}+N_{\nu}\left|q_{\nu j} \boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j}-a_{\nu j}\right|\right)>\varepsilon .
$$

Since $H$ is decreasing, we write $E=H^{-1}(\varepsilon)$ and conclude that

$$
1 \leq q_{\nu j} \leq E, \quad\left|q_{\nu j} \boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j}-a_{\nu j}\right| \leq E N_{\nu}^{-1}
$$

hold for all $\nu \in \mathbb{N}$, all $1 \leq j \leq s$. From the second inequality, one has

$$
\left|a_{\nu j}\right| \leq E N_{\nu}^{-1}+r E T \max \left|\lambda_{i j}\right| .
$$

In particular, the $2 s$-tuples $\left(q_{\nu 1}, \ldots, q_{\nu s}, a_{\nu 1}, \ldots, a_{\nu s}\right)$ can take only finitely many values, so at least one of these, say $\left(q_{1}, \ldots, q_{s}, a_{1}, \ldots, a_{s}\right)$, occurs infinitely often as $\nu$ varies. Choosing a subsequence of the $N_{\nu}$ if necessary, we can suppose (after renumbering if necessary) that $q_{\nu j}=q_{j}, a_{\nu j}=a_{j}$ hold for all $\nu$, and we then have

$$
\begin{equation*}
1 \leq q_{j} \leq E, \quad\left|q_{j} \boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j}-a_{j}\right| \leq E N_{\nu}^{-1} . \tag{3.6}
\end{equation*}
$$

Since $\boldsymbol{\alpha}_{\nu} \in[-T, T]^{r}$, and $[-T, T]^{r}$ is compact, there is a subsequence that converges to a limit $\boldsymbol{\alpha} \in[-T, T]^{r}$. As before, we may then suppose that the $N_{\nu}$ have been chosen such that $\boldsymbol{\alpha}_{\nu}$ itself already converges to $\boldsymbol{\alpha}$. With $\nu \rightarrow \infty$ in (3.6) it follows that

$$
\begin{equation*}
\boldsymbol{\alpha} \boldsymbol{\lambda}_{j}=\frac{a_{j}}{q_{j}} \quad(1 \leq j \leq s) \tag{3.7}
\end{equation*}
$$

and since $\Lambda$ is normalized, this is only possible if $\alpha_{R+1}=\ldots=\alpha_{r}=0$. By (3.6) and (3.7),

$$
\left|\left(\boldsymbol{\alpha}_{\nu}-\boldsymbol{\alpha}\right) \boldsymbol{\lambda}_{j}\right|=\left|\left(\boldsymbol{\alpha}_{\nu} \boldsymbol{\lambda}_{j}-\frac{a_{j}}{q_{j}}\right)+\left(\frac{a_{j}}{q_{j}}-\boldsymbol{\alpha} \boldsymbol{\lambda}_{j}\right)\right| \leq E N_{\nu}^{-1}
$$

We use this with $1 \leq j \leq r$. Since $\Lambda$ is normalized, $\left(\boldsymbol{\lambda}_{1}, \ldots, \boldsymbol{\lambda}_{r}\right)$ is non-singular. Hence, there is a constant $C>0$, depending only on $\boldsymbol{\lambda}_{1}, \ldots, \boldsymbol{\lambda}_{r}$, such that

$$
\left\|\boldsymbol{\alpha}_{\nu}-\boldsymbol{\alpha}\right\| \leq C E N_{\nu}^{-1}
$$

For $R<j \leq r$, this yields

$$
\begin{equation*}
\left|\alpha_{\nu j}\right| \leq C E N_{\nu}^{-1} . \tag{3.8}
\end{equation*}
$$

When $R=0$, this contradicts $\boldsymbol{\alpha}_{\nu} \in \mathfrak{m}(T)$ when $\nu$ is large. This confirms (3.5) when $R=0$.

Now suppose that $0<R \leq r$. For any $\mathbf{x} \in \mathbb{R}^{r}$ write $\tilde{\mathbf{x}}=\left(x_{1}, \ldots, x_{R}\right)$. Then, by (3.6) and (3.8),

$$
\left|\tilde{\boldsymbol{\alpha}}_{\nu} \tilde{\boldsymbol{\lambda}}_{j}-\frac{a_{j}}{q_{j}}\right| \leq E N_{\nu}^{-1}+\left|\sum_{i=R+1}^{r} \lambda_{i j} \alpha_{\nu i}\right| \leq r C E N_{\nu}^{-1} \max \left|\lambda_{i j}\right| .
$$

Since $\Lambda$ is normalized, the $R \times R$-matrix $\mathrm{A}=\left(\tilde{\boldsymbol{\lambda}}_{1}, \ldots, \tilde{\boldsymbol{\lambda}}_{R}\right)$ is non-singular. The previous display with $1 \leq j \leq R$ then yields

$$
\begin{equation*}
\tilde{\boldsymbol{\alpha}}_{\nu} \mathrm{A}=\left(\frac{a_{1}}{q_{1}}, \ldots, \frac{a_{R}}{q_{R}}\right)+\tilde{\boldsymbol{\beta}}_{\nu} \tag{3.9}
\end{equation*}
$$

with

$$
\left\|\tilde{\boldsymbol{\beta}}_{\nu}\right\| \leq C_{1} E N_{\nu}^{-1}
$$

here $C_{1}=r C \max \left|\lambda_{i j}\right|$ depends only on $\Lambda$. There is exactly one $\tilde{\mathbf{b}} \in \mathbb{R}^{R}$, one $\tilde{\boldsymbol{\gamma}}_{\nu} \in \mathbb{R}^{R}$ with

$$
\begin{equation*}
\tilde{\mathrm{b}} \mathrm{~A}=\left(\frac{a_{1}}{q_{1}}, \ldots, \frac{a_{R}}{q_{R}}\right), \quad \tilde{\boldsymbol{\gamma}}_{\nu} \mathrm{A}=\tilde{\boldsymbol{\beta}}_{\nu} . \tag{3.10}
\end{equation*}
$$

By (3.9), we then get

$$
\begin{equation*}
\tilde{\boldsymbol{\alpha}}_{\nu}=\tilde{\mathbf{b}}+\tilde{\gamma} \tag{3.11}
\end{equation*}
$$

By Cramer's rule, applied to the second equation in (3.10), the upper bounds on $\tilde{\boldsymbol{\beta}}_{\nu}$ imply that $\left\|\tilde{\boldsymbol{\gamma}}_{\nu}\right\| \leq C_{2} E N_{\nu}^{-1}$ where $C_{2}$ again only depends on $\Lambda$. Since A is an integral matrix (because $\Lambda$ is normalized), Cramer's rule also shows that $\tilde{\mathbf{b}}=\left(d_{1} / q, \ldots, d_{r} / q\right)$ with $d_{i} \in \mathbb{Z}, q \in \mathbb{N},\left(d_{1} ; \ldots, d_{r} ; q\right)=1$ and $q \mid q_{1} q_{2} \ldots q_{r} \operatorname{det} A$. In particular, $q \leq|\operatorname{det} \mathrm{A}| E^{R}$. By (3.11), it follows that $\alpha_{\nu i} \in \mathfrak{M}\left(C_{3} E+|\operatorname{det} \mathrm{A}| E^{R}\right)$ for $1 \leq i \leq R$, and in particular that $\alpha_{\nu i} \in \mathfrak{M}\left(U\left(N_{\nu}\right)\right)$ when $\nu$ is large. But then $\boldsymbol{\alpha}_{\nu} \in \mathcal{M}$, a contradiction. This establishes (3.5).

This far, we have followed Freeman's argument [12], pp. 2688-2691, with very few details that had to be adjusted to the slightly different context. We now proceed to derive Theorem 3 from (3.5) by an argument that closely follows Wooley [18]. In fact, by (3.5), for any $t \in \mathbb{N}$ there exists a number $N_{t} \in \mathbb{N}$ such that for all $N \geq N_{t}$ one has

$$
\begin{equation*}
\sup _{\alpha \in \mathfrak{m}(t)} \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{1}\right) \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{2}\right) \ldots \Xi\left(\boldsymbol{\alpha} \boldsymbol{\lambda}_{s}\right) \leq \frac{1}{t} \tag{3.12}
\end{equation*}
$$

Once $N_{t}$ is determined, it may be replaced by any larger integer. Hence, we may choose $N_{t}$ such that $N_{t}<N_{t+1}$ for all $t$, and that $U\left(N_{t}\right)>t$. Then $N_{t} \rightarrow \infty$ as $t \rightarrow \infty$. We define $T$ by $T(N)=t$ for $N_{t} \leq N<N_{t+1}$. Then $T$ satisfies $T(N) \leq U(N)$, is increasing, unbounded, and (3.2) follows from (3.12).

## 4. Extremal sequences and their exponential sums

This section is devoted to a class of sequences that have been investigated extensively in our recent memoir [5]. We begin by collecting the principal concepts and conclusions of that work, as this is fundamental for our later analysis.

Let $\mathcal{S} \subset \mathbb{N}$ be any subset of the natural numbers, and write

$$
\mathcal{S}(N, q, a)=\{s \in \mathcal{S}: s \leq N, s \equiv a \bmod q\} .
$$

Suppose that for any $a \in \mathbb{Z}, q \in \mathbb{N}$ there is a real number $s(q, a)$ such that the asymptotic formula

$$
\begin{equation*}
\# \mathcal{S}(N, q, a)=s(q, a) N+o(N) \quad(N \rightarrow \infty) \tag{4.1}
\end{equation*}
$$

holds; sequences $\mathcal{S}$ with this property are called distributed, and $\varrho(\mathcal{S})=s(1,0)$ is the density of $\mathcal{S}$. If a distributed sequence satisfies $\varrho(\mathcal{S})>0$ and the relation

$$
\begin{equation*}
\sum_{q=1}^{\infty} \sum_{\substack{a=1 \\(a, q)=1}}^{q}\left|\sum_{b=1}^{q} s(q, b) e\left(\frac{a b}{q}\right)\right|^{2}=\varrho(\mathcal{S}) \tag{4.2}
\end{equation*}
$$

holds, then $\mathcal{S}$ is an extremal sequence in the language of [5]. The series in (4.2) converges for any distributed sequence, and is bounded above by $\varrho(\mathcal{S})$ (this is [5], Lemma 1.1). Moreover, $\mathcal{S}$ is extremal if and only if the function

$$
\mathbf{s}(n)= \begin{cases}1 & s \in \mathcal{S}  \tag{4.3}\\ 0 & s \notin \mathcal{S}\end{cases}
$$

is the limit of periodic functions in the space of functions $f: \mathbb{N} \rightarrow \mathbb{R}$ equipped with the semi-norm

$$
\|f\|^{2}=\lim \sup N^{-1} \sum_{n \leq N}|f(n)|^{2}
$$

for this compare Puchta [16] or this author [4]. Many examples of extremal sequences may be found in [5]: any sequence that has positive density and a multiplicative indicator (4.3) is extremal ([5], Theorem 1.4). In particular, the squarefree numbers and more generally the $k$-free numbers form extremal sequences. Also, sequences that result from a convergent sieve process are extremal. More precisely, for any prime $p$ let $k(p) \in \mathbb{N}_{0}$, let $\Omega_{p} \subset\left(-\frac{1}{2} p^{k}, \frac{1}{2} p^{k}\right]$ with $\# \Omega_{p}<p^{k}$, and suppose that $\sum_{p} p^{-k(p)} \# \Omega_{p}$ converges. If there exists a constant $C>0$ such that $\Omega_{p} \subset[-C, C]$ for all primes $p$, then the sifted set

$$
\left\{s \in \mathbb{N}: \quad s \notin \Omega_{p} \bmod p^{k}\right\}
$$

has density $\prod_{p}\left(1-\frac{\# \Omega_{p}}{p^{k(p)}}\right.$ ), and is extremal ([5], Theorem 1.10). In particular, if $h_{1}<h_{2}<\ldots<h_{r}$ are natural numbers with $\#\left\{h_{j} \bmod p^{k}\right\}<p^{k}$ for all primes $p$, then the $r$-tuples of $k$-free numbers

$$
\left\{s \in \mathbb{N}: s+h_{j} \quad k \text {-free for } 1 \leq j \leq r\right\}
$$

are extremal. For further examples and other properties of extremal sequences, see Brüdern [4], [5] and Puchta [16].

Let $\mathcal{S}$ denote a distributed sequence, and put

$$
\begin{equation*}
S(\alpha)=\sum_{\substack{s \in \mathcal{S} \\ s \leq N}} e(\alpha s) \tag{4.4}
\end{equation*}
$$

By Parseval's identity and (4.1), one has

$$
\begin{equation*}
\int_{0}^{1}|S(\alpha)|^{2} d \alpha=\varrho(\mathcal{S}) N+o(N) \tag{4.5}
\end{equation*}
$$

Recall the definition of $\mathfrak{M}(Q)$ (below (3.1)), let $\mathfrak{N}(Q)=\mathfrak{M}(Q) \cap[0,1]$, and put $\mathfrak{n}(Q)=[0,1] \backslash \mathfrak{N}(Q)$. For extremal sequences $\mathcal{S}$, Theorem 1.6 of Brüdern [5] asserts that for any function $Q=Q(N)$ with $Q(N) \rightarrow \infty$ as $N \rightarrow \infty$ one has

$$
\begin{equation*}
\int_{\mathfrak{N}(Q)}|S(\alpha)|^{2} d \alpha=\varrho(\mathcal{S}) N+o(N) \tag{4.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{\mathfrak{n}(Q)}|S(\alpha)|^{2} d \alpha=o(N) . \tag{4.7}
\end{equation*}
$$

The properties (4.6) and (4.7) are equivalent, in view of (4.5), and actually characterize extremal sequences among the distributed ones, but this is not needed here. In the later sections of this paper, we shall make crucial use of (4.7) along with the following upper bound that in contrast to (4.7), is valid for distributed sequences.
Lemma 4. Let $\mathcal{S}$ denote a distributed sequence. Then there exists an increasing function $Q=Q(N)$ with $Q(N) \rightarrow \infty$ as $N \rightarrow \infty$, and a decreasing function $H$ with $H(t) \rightarrow 0$ as $t \rightarrow \infty$ such that the inequality $|S(\alpha)| \ll N \Xi(\alpha)$ holds for all $\alpha \in \mathfrak{M}(Q)$; here $S$ is the exponential sum (4.4), and $\Xi$ is the weight on $\mathfrak{M}(Q)$ associated with $H$.
Proof. For any distributed sequence $\mathcal{S}$, there exists a function $T(N)$ such that $T(N)$ and $N / T(N)$ are increasing, and $T(N) \rightarrow \infty$ as $N \rightarrow \infty$, with the property that

$$
|\# \mathcal{S}(N, q, b)-s(q, b) N| \leq N / T(N)
$$

holds uniformly in $N \geq 1, q \in \mathbb{N}, b \in \mathbb{Z}$ (this is [5], Lemma 2.2). Hence, by (4.4), whenever $(a ; q)=1$,

$$
\begin{equation*}
S\left(\frac{a}{q}\right)=\sum_{b=1}^{q} \sum_{n \in \mathcal{S}(N, q, b)} e\left(\frac{a b}{q}\right)=N G(q, a)+E(N, q) \tag{4.8}
\end{equation*}
$$

where

$$
G(q, a)=\sum_{b=1}^{q} s(q, b) e\left(\frac{a b}{q}\right),
$$

and where $|E(N, q)| \leq q N / T(N)$. By partial summation and (4.8),

$$
\begin{aligned}
S\left(\frac{a}{q}+\beta\right) & =e(N \beta) S\left(\frac{a}{q}\right)-2 \pi i \beta \int_{1}^{N} e(\beta \gamma) \sum_{\substack{n \leq \gamma \\
n \in \mathcal{S}}} e\left(\frac{a n}{q}\right) d \gamma \\
& =G(q, a) I(\beta)+e(N \beta) E(N, q)-2 \pi i \beta \int_{1}^{N} e(\beta \gamma) E(\gamma, q) d \gamma
\end{aligned}
$$

where

$$
\begin{equation*}
I(\beta)=\int_{0}^{N} e(\beta \gamma) d \gamma \tag{4.9}
\end{equation*}
$$

Since $N / T(N)$ is increasing, we have $|E(\gamma, q)| \leq q N / T(N)$ for $1 \leq \gamma \leq N$. It follows that

$$
\begin{equation*}
\left|S\left(\frac{a}{q}+\beta\right)-G(q, a) I(\beta)\right| \leq \frac{q N}{T(N)}(1+2 \pi N|\beta|) \tag{4.10}
\end{equation*}
$$

As we have remarked earlier, the sum in (4.2) converges. Hence, there is a function $h(q)$ that decreases to 0 as $q \rightarrow \infty$, and such that

$$
\sum_{\substack{a=1 \\(a, q)=1}}^{q}|G(q, a)|^{2} \leq h(q)^{2}
$$

holds for all $q$. In particular, this yields $|G(q, a)| \leq h(q)$ for all $(a, q)=1$. By partial integration, applied to (4.9), we infer from (4.10) the inequality

$$
\left|S\left(\frac{a}{q}+\beta\right)\right| \leq \frac{2 N h(q)}{1+N|\beta|}+\frac{q N}{T(N)}(1+2 \pi N|\beta|) .
$$

We may suppose that $h(q) \geq 1 / q$ for all $q \in \mathbb{N}$, for otherwise we replace $h(q)$ by $\max (h(q), 1 / q)$ without affecting the previous argument. We now extend the definition of $h$ to $[1, \infty)$ by $h(\gamma)=h([\gamma])$. Then $h$ is still a non-increasing function, and

$$
\left|S\left(\frac{a}{q}+\beta\right)\right| \leq 2 N h(q) h(1+N|\beta|)+\frac{q N}{T(N)}(1+2 \pi N|\beta|) .
$$

Since $h$ is non-increasing, the inequality $h(u) h(v) \leq h(1) h(\sqrt{u v})$ holds for any $u \geq 1, v \geq 1$. Hence, $H(t)=h(\sqrt{t})$ is non-increasing, satisfies $H(t) \rightarrow 0$ as $t \rightarrow \infty$, and one has

$$
\begin{equation*}
\left|S\left(\frac{a}{q}+\beta\right)\right| \leq 2 h(1) N H(q(1+N|\beta|))+\frac{N}{T(N)} q(1+2 \pi N|\beta|) . \tag{4.11}
\end{equation*}
$$

Now choose a function $Q(N)$ with $Q(N) \rightarrow \infty$ as $N \rightarrow \infty$, and such that $H(2 Q(N)) \geq 1 / \sqrt{T(N)}$; this is always possible since $H$ and $1 / T$ are non-increasing and tend to zero as $N \rightarrow \infty$. We may now replace $Q(N)$ with $\min (Q(N), \sqrt{T(N)})$ if necessary to ensure that $Q(N) \leq \sqrt{T(N)}$; the inequality $H(2 Q(N)) \geq 1 / \sqrt{T(N)}$ then still holds. By (4.11), it follows that for $\frac{a}{q}+\beta \in \mathfrak{M}(Q)$, one has

$$
\begin{aligned}
\left|S\left(\frac{a}{q}+\beta\right)\right| & \leq 2 h(1) N H(q(1+N \mid \beta))+8 N Q(N) T(N)^{-1} \\
& \leq 2 h(1) N H(q(1+N \mid \beta))+8 N T(N)^{-1 / 2} \\
& \leq(2 h(1)+8) H(q(1+N|\beta|)) .
\end{aligned}
$$

This establishes Lemma 4.

## 5. An auxiliary asymptotic formula

In this section we launch our attack on simultaneous diophantine approximations. The approach is based on the Fourier transform method of Davenport and Heilbronn. We work with the kernel functions

$$
\begin{equation*}
K_{\eta}(\alpha)=\eta\left(\frac{\sin \pi \eta \alpha}{\pi \eta \alpha}\right)^{2}, \quad \Upsilon_{\eta}(\alpha)=\max \left(0,1-\frac{|\alpha|}{\eta}\right) \tag{5.1}
\end{equation*}
$$

where $\eta>0$ is a real parameter. Note that $K_{\eta}$ and $\Upsilon_{\eta}$ are Fourier transforms of each other; we only require the classical formula

$$
\begin{equation*}
\Upsilon_{\eta}(\alpha)=\int_{-\infty}^{\infty} K_{\eta}(\beta) e(-\alpha \beta) d \beta \tag{5.2}
\end{equation*}
$$

Higher-dimensional analogues arise from Fubini's theorem. For any function $f: \mathbb{R} \rightarrow \mathbb{C}$ we define $f: \mathbb{R}^{r} \rightarrow \mathbb{C}$ by

$$
\begin{equation*}
f(\boldsymbol{\alpha})=f\left(\alpha_{1}\right) f\left(\alpha_{2}\right) \ldots f\left(\alpha_{r}\right) \tag{5.3}
\end{equation*}
$$

confusion should not arise from the ambiguous use of $f$ as the number of variables is usually clear from the context.

Fix the following notation for the rest of this section and the next. Let $\mathcal{S}_{0}, \mathcal{S}_{1}, \ldots, \mathcal{S}_{r}$ denote extremal sequences, of densities $\varrho_{j}=\varrho\left(\mathcal{S}_{j}\right)$. Let $\boldsymbol{\lambda}=$ $\left(\lambda_{1}, \ldots, \lambda_{r}\right) \in \mathbb{R}^{r}$ and suppose that $1, \lambda_{1}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. Put

$$
\begin{equation*}
\varrho=\varrho_{0} \varrho_{1} \ldots \varrho_{r}, \quad \kappa=\max \left(1, \lambda_{1}, \ldots, \lambda_{r}\right) \tag{5.4}
\end{equation*}
$$

Implicit constants, as well as several constructions below, depend on $\boldsymbol{\lambda}$ as well as the $\mathcal{S}_{j}$, but not on $N$ or $\eta$.

For any $\boldsymbol{\beta} \in \mathbb{R}^{r}$, we count solutions of the inhomogenous system

$$
\left|x_{i}-\lambda_{i} x_{0}-\beta_{i}\right|<\eta \quad(1 \leq i \leq r)
$$

with a certain weight attached to the variables $x_{j} \in \mathcal{S}_{j}$. Let

$$
\begin{equation*}
P_{\eta}(N, \boldsymbol{\beta})=\sum_{\substack{1 \leq x_{j} \leq N \\ x_{j} \in \mathcal{S}_{j} \\ j=0,1, \ldots, r}} \prod_{i=1}^{r} \Upsilon_{\eta}\left(x_{i}-\lambda_{i} x_{0}-\beta_{i}\right) \tag{5.5}
\end{equation*}
$$

Our goal is to establish an asymptotic formula for $P_{\eta}(N, \boldsymbol{\beta})$. By (5.2) and (5.3), if $S_{j}(\alpha)$ is the sum (4.4) for $\mathcal{S}_{j}$,

$$
\begin{equation*}
P_{\eta}(N, \boldsymbol{\beta})=\int_{\mathbb{R}^{r}} S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \tag{5.6}
\end{equation*}
$$

This Fourier integral is the traditional point of departure for the DavenportHeilbronn method. One splits the integral into three parts: a small neighbourhood
of the origin (the major arc) contributes a prospective main term, and the part where $\|\boldsymbol{\alpha}\|$ is large (the trivial region) contributes very little, due to the decay of $K_{\eta}(\boldsymbol{\alpha})$. In the intermediate region (the minor arc), one uses inequalities of Weyl's type and some argument of the type discussed in section 3 for an upper bound. There are two points where we have to digress from this well trodden path. Unlike earlier versions of the method, in our case the dissection will depend not only on the sets $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$, but also on $\boldsymbol{\lambda}$. This seems to be an intrinsic feature of the Bentkus-Götze-Freeman device, and apparently is unavoidable (compare the discussion in §8). Moreover, in our description of the method, there is no trivial region. Instead, we use ideas from our paper [3] that are based on the Poisson summation formula. It seems that some sort of refinement on the trivial region is actually necessary when the number of variables is too low for a direct use of mean value estimates of Hua's type; this is the case here.

We begin by constructing the major arc. By (4.10), applied to $S_{j}$, with $q=1$ and $a=0$, there exists a function $V_{j}(N)$ with $V_{j}(N) \rightarrow \infty$ as $N \rightarrow \infty$, and such that whenever $|\alpha| \leq V_{j}(N) / N$ one has

$$
\begin{equation*}
\left|S_{j}(\alpha)-\varrho_{j} I(\alpha)\right| \leq N / V_{j}(N) \tag{5.7}
\end{equation*}
$$

Let $V=\min \left(V_{0}, \ldots, V_{r}\right)$ and $U^{*}=V^{1 /(r+1)}$. Then $U^{*}(N) \rightarrow \infty$ as $N \rightarrow \infty$. For $1 \leq U \leq U^{*}(N)$ we define

$$
\begin{equation*}
\mathfrak{K}=\{\boldsymbol{\alpha}:\|\boldsymbol{\alpha}\| \leq U / N\}, \quad \mathfrak{k}=\mathbb{R}^{r} / \mathfrak{K}=\left\{\boldsymbol{\alpha} \in \mathbb{R}^{r}:\|\boldsymbol{\alpha}\|>U / N\right\} . \tag{5.8}
\end{equation*}
$$

For $\|\boldsymbol{\alpha}\| \leq U^{*} / N$, one has $|\boldsymbol{\lambda} \boldsymbol{\alpha}| \leq r\|\boldsymbol{\lambda}\| U^{*}(N) / N \leq V(N) / N$ for large $N$, whence by (5.7), for these $\boldsymbol{\alpha}$,

$$
\begin{equation*}
S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha})=\varrho I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha})+O\left(N^{r+1} / V(N)\right) \tag{5.9}
\end{equation*}
$$

where $\varrho$ is given by (5.4) and $I(\boldsymbol{\alpha})$ is defined via (4.9) and (5.3). We multiply with the bounded function $K(\boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta})$ and integrate over $\mathfrak{K}$ to infer that

$$
\begin{aligned}
\int_{\mathfrak{K}} S_{1}\left(\alpha_{1}\right) & \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \\
& =\varrho \int_{\mathfrak{K}} I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}+O\left(N U^{r} V^{-1}\right)
\end{aligned}
$$

We now complete the integral on the right hand side to the singular integral

$$
\begin{equation*}
P_{\eta}^{*}(N, \boldsymbol{\beta})=\int_{\mathbb{R}^{r}} I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}, \tag{5.10}
\end{equation*}
$$

and then proceed to evaluate (5.10) asymptotically. The arguments are largely standard, but some care is required because only very few variables are present. By (5.8) we have

$$
\begin{align*}
& \int_{\mathfrak{K}} S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \\
& \quad=\varrho P_{\eta}^{*}(N, \boldsymbol{\beta})-\int_{\mathfrak{k}} I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}+O\left(N U^{-1}\right) . \tag{5.11}
\end{align*}
$$

For simplicity, write $\alpha_{0}=-\boldsymbol{\lambda} \boldsymbol{\alpha}$ (as a function of $\alpha_{1}, \ldots, \alpha_{r}$ ) and put

$$
J_{l}(\boldsymbol{\alpha})=\prod_{\substack{0 \leq j \leq r \\ j \neq l}} I\left(\alpha_{j}\right) .
$$

Then, by Hölder's inequality,

$$
\begin{equation*}
\int_{\mathfrak{k}}|I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha})| d \boldsymbol{\alpha} \leq \prod_{l=0}^{r}\left(\int_{\mathfrak{k}}\left|J_{l}(\boldsymbol{\alpha})\right|^{\frac{r+1}{r}} d \boldsymbol{\alpha}\right)^{\frac{1}{r+1}} \tag{5.12}
\end{equation*}
$$

But $J_{0}(\boldsymbol{\alpha})=I(\boldsymbol{\alpha})$, and $I(\alpha) \ll N(1+N|\alpha|)^{-1}$, as one finds by partial integration. Consequently,

$$
\int_{\mathfrak{k}}\left|J_{0}(\boldsymbol{\alpha})\right|^{\frac{r+1}{r}} d \boldsymbol{\alpha} \ll N^{r+1} \int_{\mathfrak{k}} \prod_{j=0}^{r}\left(1+N\left|\alpha_{j}\right|\right)^{-\frac{r+1}{r}} d \boldsymbol{\alpha} \ll N U^{-1 / r} .
$$

A cruder argument suffices for $1 \leq l \leq r$. When $l=r$, substitute $\alpha_{0}$ for $\alpha_{r}$ to infer that

$$
\begin{aligned}
\int_{\mathfrak{k}}\left|J_{r}(\boldsymbol{\alpha})\right|^{\frac{r+1}{r}} d \boldsymbol{\alpha} & \leq \int_{\mathbb{R}^{r}}\left|J_{r}(\boldsymbol{\alpha})\right|^{\frac{r+1}{r}} d \boldsymbol{\alpha} \\
& =\lambda_{r}^{-1} \int_{\mathbb{R}^{r}}\left|I\left(\alpha_{0}\right) I\left(\alpha_{1}\right) \ldots I\left(\alpha_{r-1}\right)\right|^{\frac{r+1}{r}} d \alpha_{0} \ldots d \alpha_{r-1} \\
& =\lambda_{r}^{-1}\left(\int_{-\infty}^{\infty}|I(\alpha)|^{\frac{r+1}{r}} d \alpha\right)^{r} \ll N
\end{aligned}
$$

here the implicit constant depends on $\boldsymbol{\lambda}$. By symmetry, the same bound holds with $J_{l}$ in place of $J_{r}$, and we infer from (5.12) that

$$
\begin{equation*}
\int_{\mathfrak{k}}|I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha})| d \boldsymbol{\alpha} \ll N U^{-\frac{1}{r(r+1)}} . \tag{5.13}
\end{equation*}
$$

Since $e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) \ll 1$ for $0<\eta \leq 1$, we can use (5.13) to simplify (5.11) to

$$
\begin{equation*}
\int_{\mathfrak{K}} S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}=\varrho P_{\eta}^{*}(N, \boldsymbol{\beta})+O\left(N U^{-\frac{1}{r(r+1)}}\right) \tag{5.14}
\end{equation*}
$$

It remains to evaluate $P_{\eta}^{*}(N, \boldsymbol{\beta})$, as defined in (5.10). Use the definition of $I(\alpha)$ in (4.9) to write $I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha})$ as an $r+1$-fold integral. Then, by (5.10) and 5.2,

$$
P_{\eta}^{*}(N, \boldsymbol{\beta})=\int_{0}^{N} \int_{[0, N]^{r}} \Upsilon_{\eta}\left(\beta_{1}-t_{1}+\lambda_{1} t_{0}, \ldots, \beta_{r}-t_{r}+\lambda_{r} t_{0}\right) d \mathbf{t} d t_{0}
$$

Note the formal similarity with the definition of $P_{\eta}(N, \boldsymbol{\beta})$ in (5.5), $P^{*}(N, \boldsymbol{\beta})$ is a continuous analogue that is easy to compute. We substitute $u_{j}=t_{j}-\beta_{j}-\lambda_{j} t_{0}$ in the inner integral and write $t$ for $t_{0}$. Then, since $\Upsilon_{\eta}(-u)=\Upsilon_{\eta}(u)$,

$$
\begin{equation*}
P_{\eta}^{*}(N, \boldsymbol{\beta})=\int_{0}^{N} \int_{\mathfrak{B}(t)} \Upsilon_{\eta}(\mathbf{u}) d \mathbf{u} d t \tag{5.15}
\end{equation*}
$$

where $\mathfrak{B}(t)$ is the set of $\mathbf{u} \in \mathbb{R}^{r}$ with $0 \leq u_{j}+\lambda_{j} t+\beta_{j} \leq N$ for all $j=1, \ldots, r$.

So far, all estimates have been uniform with respect to $\boldsymbol{\beta}$. It is now necessary to assume that $\boldsymbol{\beta}$ is small, whence we suppose that $\|\boldsymbol{\beta}\| \leq \sqrt{N}$. Then, for $N^{2 / 3} \leq$ $t \leq \kappa^{-1} N-N^{2 / 3}$, we have $\lambda_{j} N^{2 / 3}-N^{1 / 2} \leq \lambda_{j} t+\beta_{j} \leq N-\lambda_{j} N^{2 / 3}+N^{1 / 2}$, and hence, for large $N$, the box $\|\mathbf{u}\| \leq 1$ is contained in $\mathfrak{B}(t)$. Hence, for these $t$, and for $0<\eta \leq 1$, one has

$$
\int_{\mathfrak{B}(t)} \Upsilon_{\eta}(\mathbf{u}) d \mathbf{u}=\left(\int_{-1}^{1} \Upsilon_{\eta}(u) d u\right)^{r}=\eta^{r} .
$$

Similarly, when $t \geq \kappa^{-1} N+N^{2 / 3}$ and $N$ is large, there is at least one $j$ for which the intersection of the intervals $0 \leq u_{j}+\lambda_{j} t+\beta_{j} \leq N$ and $\left|u_{j}\right| \leq 1$ is empty. Hence, in this case, the inner integral in (5.15) vanishes. In the regions $0 \leq t \leq N^{2 / 3}$ and $\left|t-\kappa^{-1} N\right| \leq N^{2 / 3}$ it suffices to use the upper bound

$$
\int_{\mathfrak{B}(t)} \Upsilon_{\eta}(\mathbf{u}) d \mathbf{u} \leq \int_{\mathbb{R}^{r}} \Upsilon_{\eta}(\mathbf{u}) d \mathbf{u}=\eta^{r}
$$

that is actually valid for all $t \in \mathbb{R}$. Now split the outer integral into the regions $0 \leq t \leq N^{2 / 3}, N^{2 / 3}<t \leq \kappa^{-1} N-N^{2 / 3}, \kappa^{-1} N-N^{2 / 3}<t \leq \min \left(N, \kappa^{-1} N+N^{2 / 3}\right)$, larger $t \leq N$ (if any) do not contribute by the above observations. Then, on collecting together,

$$
P_{\eta}^{*}(N, \boldsymbol{\beta})=\kappa^{-1} N \eta^{r}+O\left(N^{2 / 3} \eta^{r}\right) .
$$

By (5.14) we have proved the first statement in the following lemma.
Lemma 5. Let $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$ denote extremal sequences, and let $\lambda_{1}, \ldots, \lambda_{r}$ denote positive real numbers such that $1, \lambda_{1}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. Then there exists a function $U^{*}(N)$ with $U^{*}(N) \rightarrow \infty$ as $N \rightarrow \infty$, and such that uniformly in $0<\eta \leq 1,\|\boldsymbol{\beta}\| \leq \sqrt{N}$ and $1 \leq U \leq U^{*}(N)$ one has

$$
\begin{aligned}
& \int_{\|\boldsymbol{\alpha}\| \leq U / N} S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha}) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \\
& =\varrho \kappa^{-1} \eta^{r} N+O\left(N U^{-\frac{1}{r(r+1)}}\right)
\end{aligned}
$$

and

$$
\int_{U / N \leq\|\boldsymbol{\alpha}\| \leq U^{*} / N}\left|S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha})\right| d \boldsymbol{\alpha} \ll N U^{-\frac{1}{r(r+1)}} .
$$

Here $\kappa$ and $\varrho$ are defined by (5.4).
To establish the final inequality in Lemma 5, apply the triangle inequality to (5.9) and integrate over $U / N \leq\|\boldsymbol{\alpha}\| \leq U^{*} / N$. Then, by (5.8),

$$
\begin{aligned}
& \int_{U / N \leq\|\boldsymbol{\alpha}\| \leq U^{*} / N}\left|S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha})\right| d \boldsymbol{\alpha} \\
& \ll \int_{\mathfrak{k}}|I(\boldsymbol{\alpha}) I(-\boldsymbol{\lambda} \boldsymbol{\alpha})| d \boldsymbol{\alpha}+N U^{* r} V^{-1} .
\end{aligned}
$$

But $U^{*}=V^{1 /(r+1)}$, and the required bound follows from (5.13).

## 6. Counting continued

In this section we complete the analysis of the weighted counter $P_{\eta}(N, \boldsymbol{\beta})$, but before we direct our forces to the treatment of the minor arc portion of the integral (5.6), we pause for two auxiliary identities that will frequently be used later.

Let $S(\alpha)$ denote any exponential sum, as in (4.4). Let $\nu, \mu, \mu^{\prime}$ denote real numbers with $\nu \neq 0$. Then, provided only that $|\nu| \geq \eta>0$, we deduce from (5.2) the convolution formula
$\int_{-\infty}^{\infty} S(\nu \alpha+\mu) S\left(-\nu \alpha-\mu^{\prime}\right) K_{\eta}(\alpha) d \alpha=\sum_{\substack{s, s^{\prime} \leq N \\ s, s^{\prime} \in \mathcal{S}}} \Upsilon_{\eta}\left(\nu\left(s-s^{\prime}\right)\right) e\left(\mu s-\mu^{\prime} s^{\prime}\right)=S\left(\mu-\mu^{\prime}\right)$.
Our second tool concerns functions $f \in L^{1}(\mathbb{R} / \mathbb{Z})$ that we view as functions of period 1 on the real line. Then, by Lemma 1 of Brüdern and Perelli [6], or Brüdern [3], eqn. (4), for any function $K \in L^{1}(\mathbb{R})$ that has a compactly supported Fourier transform $\hat{K}$ one has

$$
\begin{equation*}
\int_{-\infty}^{\infty} f(\alpha) K(\alpha) e(-\alpha \beta) d \alpha=\sum_{n=-\infty}^{\infty} \hat{K}(\beta-n) \int_{0}^{1} f(\alpha) e(-\alpha n) d \alpha \tag{6.2}
\end{equation*}
$$

this is readily proved by Poisson summation ([3], p. 63). In particular, we have

$$
\int_{-\infty}^{\infty} f(\alpha) K_{\eta}(\alpha) d \alpha=\int_{0}^{1} f(\alpha) d \alpha \quad(0<\eta \leq 1) .
$$

Two special cases are of interest to us. With $S(\alpha)$ as before and $\mathfrak{m}(Q)$ and $\mathfrak{n}(Q)=$ $\mathfrak{m}(Q) \cap[0,1]$ as in section 3 , one has

$$
\begin{equation*}
\int_{\mathfrak{m}(Q)}|S(\alpha)|^{2} K_{\eta}(\alpha) d \alpha=\int_{\mathfrak{n}(Q)}|S(\alpha)|^{2} d \alpha \tag{6.3}
\end{equation*}
$$

To see this, it suffices to take $f(\alpha)=|S(\alpha)|^{2}$ when $\alpha \in \mathfrak{m}(Q)$, and $f(\alpha)=0$ otherwise, in the preceding identity. Similarly, one finds

$$
\begin{equation*}
\int_{\mathfrak{M}(Q)} K_{\eta}(\alpha) d \alpha=\int_{\mathfrak{N}(Q)} d \alpha \ll Q^{2} / N . \tag{6.4}
\end{equation*}
$$

We return to the main theme, and use the notation from section 5 throughout. The objective is to estimate the quantity

$$
\begin{equation*}
\int_{\mathfrak{k}} S_{0}\left(\alpha_{0}\right) S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right) e(-\boldsymbol{\alpha} \boldsymbol{\beta}) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}=\mathcal{I}(\mathfrak{k}) \tag{6.5}
\end{equation*}
$$

say; here $\alpha_{0}=-\boldsymbol{\lambda} \boldsymbol{\alpha}$ as on earlier occasions. Some preparation is required before we can bring Theorem 3 into play. First we apply Lemma 4 to each of the sequences $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$. This supplies functions $H_{j}(t)$ that decrease to 0 as
$t \rightarrow \infty$, and increasing unbounded functions $Q_{j}(N)$, such that $\left|S_{j}(\alpha)\right| \ll N \Xi_{j}(\alpha)$ holds for $\alpha \in \mathfrak{M}\left(Q_{j}(N)\right)$; here $\Xi_{j}$ is the weight associated to $H_{j}$. Now take $H=\max \left(H_{0}, H_{1}, \ldots, H_{r}\right)$ and $Q=\min \left(Q_{0}, \ldots, Q_{r}\right)$. Then

$$
\begin{equation*}
S_{j}(\alpha) \ll N \Xi(\alpha) \text { for } \alpha \in \mathfrak{M}(Q) \tag{6.6}
\end{equation*}
$$

holds for all $0 \leq j \leq r$. With this choice of $Q(N)$ and $U^{*}(N)$ in the role of $U$, apply Theorem 3. By the comment at the end of section 2 , the relevant coefficient matrix has integral rank 0 , and there exists a function $T(N)$, tending to $\infty$ with $N$, with $T(N) \leq U^{*}(N)$ and such that

$$
\begin{equation*}
\sup _{U^{*}(N) / N \leq\|\boldsymbol{\alpha}\| \leq T(N)} \Xi\left(\alpha_{0}\right) \Xi\left(\alpha_{1}\right) \ldots \Xi\left(\alpha_{r}\right) \ll T(N)^{-1} . \tag{6.7}
\end{equation*}
$$

The set $\mathfrak{k}$ is to be dissected into $r+2$ subsets that must be considered separately. We need to define various functions before we can describe this in detail. First, we apply (4.7) in combination with (6.3) to $S=S_{j}$. This shows that for any function $V_{j}$ with $V_{j}(N) \rightarrow \infty$ as $N \rightarrow \infty$ there exists a function $U_{j}$ that also increases to infinity, and such that

$$
\begin{equation*}
\int_{\mathfrak{m}\left(V_{j}\right)}\left|S_{j}(\alpha)\right| K_{\eta}(\alpha) d \alpha \ll N U_{j}^{-1} \tag{6.8}
\end{equation*}
$$

Here we may suppose that $U_{j} \leq V_{j}$ for otherwise $U_{j}$ can be replaced by $\min \left(U_{j}, V_{j}\right)$. Now choose functions $V_{0}, V_{r}, V_{r-1}, \ldots, V_{1}$ as follows: let

$$
V=V_{0}=\min \left(Q, T^{1 / 4 r}\right) .
$$

This defines a function $U_{0}$ with $U_{0}(N) \rightarrow \infty$ as $N \rightarrow \infty$, and such that (6.8) holds for $j=0$. Next, put $V_{r}=U_{0}^{1 / 5 r}$ and choose an increasing, unbounded function $U_{r}$ such that (6.8) holds for $j=r$. Proceed likewise and recursively put

$$
\begin{equation*}
V_{j}=U_{j+1}^{1 / 5 j} \quad(1 \leq j \leq r-1) . \tag{6.9}
\end{equation*}
$$

Note that $U_{j} \leq V_{j}$ implies the inequalities

$$
V_{j} \leq V_{j+1}^{1 / 5 j} \quad(1 \leq j \leq r-1), \quad V_{r}^{1 / 5 r} \leq V_{0}=V
$$

The functions $V_{j}$ now being defined, we choose $U(N)=\min \left(U_{0}(N), V_{1}(N)\right)$ in the work of section 5 . This fixes the definition of $\mathfrak{k}$, and Lemma 5 is applicable. Now write

$$
\begin{aligned}
\mathfrak{h}_{l} & =\left\{\boldsymbol{\alpha} \in \mathfrak{k}: \alpha_{l} \in \mathfrak{m}\left(V_{l}\right), \alpha_{j} \in \mathfrak{M}\left(V_{j}\right) \text { for } 1 \leq j \leq l-1\right\}, \\
\mathfrak{h}^{*} & =\left\{\boldsymbol{\alpha} \in \mathfrak{k}: \alpha_{0} \in \mathfrak{m}\left(V_{0}\right), \alpha_{j} \in \mathfrak{M}\left(V_{j}\right) \text { for } 1 \leq j \leq r\right\}, \\
\mathfrak{H} & =\left\{\boldsymbol{\alpha} \in \mathfrak{k}: \alpha_{j} \in \mathfrak{M}\left(V_{j}\right) \text { for } 0 \leq j \leq r\right\} .
\end{aligned}
$$

Then $\mathfrak{k}$ is the disjoint union of $\mathfrak{h}_{1}, \ldots, \mathfrak{h}_{r}, \mathfrak{h}^{*}$ and $\mathfrak{H}$. Accordingly, by an appropriate interpretation of (6.5),

$$
\begin{equation*}
\mathcal{I}(\mathfrak{k})=\mathcal{I}(\mathfrak{H})+\mathcal{I}\left(\mathfrak{h}^{*}\right)+\sum_{l=1}^{r} \mathcal{I}\left(\mathfrak{h}_{l}\right) . \tag{6.10}
\end{equation*}
$$

The integral $\mathcal{I}(\mathfrak{H})$ is estimated in three steps. We consider the subsets

$$
\mathfrak{H}^{*}=\left\{\boldsymbol{\alpha} \in \mathfrak{H}: U^{*} / N<\|\boldsymbol{\alpha}\| \leq T\right\}, \quad \mathfrak{H}_{*}=\{\boldsymbol{\alpha} \in \mathfrak{H}:\|\boldsymbol{\alpha}\|>T\}
$$

and then first observe that by Lemma 5 we have

$$
\mathcal{I}(\mathfrak{H}) \ll\left|\mathcal{I}\left(\mathfrak{H}^{*}\right)\right|+\left|\mathcal{I}\left(\mathfrak{H}_{*}\right)\right|+N U^{-1 / r(r+1)}
$$

Note that $\mathfrak{H} \subset \mathfrak{M}^{r}$ where $\mathfrak{M}=\mathfrak{M}(V)$. By (6.5) and (6.6), then (6.7) followed by (6.4), we find that

$$
\begin{aligned}
\mathcal{I}\left(\mathfrak{H}^{*}\right) & \ll N^{r+1} \int_{\mathfrak{H}^{*}} \Xi\left(\alpha_{0}\right) \Xi\left(\alpha_{1}\right) \ldots \Xi\left(\alpha_{r}\right) K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \\
& \ll N^{r+1} T^{-1} \int_{\mathfrak{M}^{r}} K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \ll N T^{-1} V^{2 r} \ll N T^{-1 / 2} .
\end{aligned}
$$

For $\mathfrak{H}_{*}$ we apply brute force, estimate all $S_{j}\left(\alpha_{j}\right)$ trivially and then use symmetry in $\boldsymbol{\alpha}$ to confirm the bounds

$$
\mathcal{I}\left(\mathfrak{H}_{*}\right) \ll N^{r+1} \int_{\mathfrak{H}_{*}} K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha} \ll N^{r+1}\left(\int_{\mathfrak{M}} K_{\eta}(\alpha) d \alpha\right)^{r-1}\left(\int_{\mathfrak{L}(T)} K_{\eta}(\alpha) d \alpha\right)
$$

where $\mathfrak{L}(T)=\mathfrak{M} \cap[T, \infty)$. But straightforward estimates based on (5.1) give

$$
\int_{\mathfrak{L}(T)} K_{\eta}(\alpha) d \alpha \ll \sum_{m=[T]}^{\infty} m^{-2} \int_{\mathfrak{N}} d \alpha \ll T^{-1} V^{2} N^{-1}
$$

so that (6.4) confirms the bound $\mathcal{I}\left(\mathfrak{H}_{*}\right) \ll N V^{2 r} T^{-1}$. This combines to

$$
\begin{equation*}
\mathcal{I}(\mathfrak{H}) \ll N T^{-1 / 2}+N U^{-1 / r(r+1)} \tag{6.11}
\end{equation*}
$$

We now turn our attention to the set $\mathfrak{h}_{l}$ with $1 \leq l \leq r$. Note that $U(N) \leq$ $V_{j}(N)$ for all $j=0, \ldots, r$. In particular, when $1 \leq l \leq r$ and $\boldsymbol{\alpha} \in \mathbb{R}^{r}$ with $\alpha_{l} \in \mathfrak{m}\left(V_{l}\right)$, then $\|\boldsymbol{\alpha}\| \geq V_{l} / N$, and hence $\boldsymbol{\alpha} \in \mathfrak{k}$. This shows that $\mathfrak{h}_{l}$ is a cartesian product. It is convenient to introduce the following notation: for $1 \leq l \leq r$ and $\boldsymbol{\alpha} \in \mathbb{R}^{s}$, let $\boldsymbol{\alpha}_{l}=\left(\alpha_{1}, \ldots, \alpha_{l-1}, \alpha_{l+1}, \ldots, \alpha_{r}\right)$, and put

$$
\mathfrak{F}_{l}=\left\{\boldsymbol{\alpha}_{l} \in \mathbb{R}^{r-1}: \alpha_{j} \in \mathfrak{M}\left(V_{j}\right)(1 \leq j \leq l-1)\right\}
$$

Then, by (6.5), (5.3) and Fubini's theorem,

$$
\begin{equation*}
\mathcal{I}\left(\mathfrak{h}_{l}\right)=\int_{\mathfrak{m}\left(V_{l}\right)} S_{l}\left(\alpha_{l}\right) e\left(-\alpha_{l} \beta_{l}\right) F_{l}\left(\alpha_{l}\right) K_{\eta}\left(\alpha_{l}\right) d \alpha_{l} \tag{6.12}
\end{equation*}
$$

where

$$
\begin{equation*}
F_{l}\left(\alpha_{l}\right)=\int_{\tilde{F}_{l}} \prod_{\substack{j=0 \\ j \neq l}}^{r} S_{j}\left(\alpha_{j}\right) e\left(-\boldsymbol{\alpha}_{l} \boldsymbol{\beta}_{l}\right) K_{\eta}\left(\boldsymbol{\alpha}_{l}\right) d \boldsymbol{\alpha}_{l}, \tag{6.13}
\end{equation*}
$$

here we continue to use the abbreviation $\alpha_{0}=-\boldsymbol{\lambda} \boldsymbol{\alpha}$. By Cauchy's inequality,

$$
\begin{equation*}
\mathcal{I}\left(\mathfrak{h}_{l}\right) \leq\left(\int_{\mathfrak{m}\left(V_{l}\right)}\left|S_{l}(\alpha)\right|^{2} K_{\eta}(\alpha) d \alpha\right)^{\frac{1}{2}} \mathcal{J}_{l}^{\frac{1}{2}} \tag{6.14}
\end{equation*}
$$

where

$$
\mathcal{J}_{l}=\int_{-\infty}^{\infty}\left|F_{l}(\alpha)\right|^{2} K_{\eta}(\alpha) d \alpha
$$

We expand this last integral. On writing

$$
\Phi_{l}\left(\boldsymbol{\alpha}_{l}\right)=\prod_{\substack{j=1 \\ j \neq l}}^{r} S_{j}\left(\alpha_{j}\right)
$$

and

$$
\Psi_{l}\left(\boldsymbol{\alpha}_{l}, \boldsymbol{\alpha}_{l}^{\prime}\right)=\int_{-\infty}^{\infty} S_{0}\left(-\lambda_{l} \alpha_{l}-\boldsymbol{\lambda}_{l} \boldsymbol{\alpha}_{l}\right) S_{0}\left(\lambda_{l} \alpha_{l}+\boldsymbol{\lambda}_{l} \boldsymbol{\alpha}_{l}^{\prime}\right) K_{\eta}\left(\alpha_{l}\right) d \alpha_{l},
$$

one finds from (6.13) and Fubini's theorem that

$$
\mathcal{J}_{l}=\int_{\mathfrak{F}_{l}} \int_{\mathfrak{F}_{l}} \Phi_{l}\left(\boldsymbol{\alpha}_{l}\right) \Phi_{l}\left(-\boldsymbol{\alpha}_{l}^{\prime}\right) e\left(\boldsymbol{\beta}_{l}\left(\boldsymbol{\alpha}_{l}^{\prime}-\boldsymbol{\alpha}_{l}\right)\right) \Psi_{l}\left(\boldsymbol{\alpha}_{l}, \boldsymbol{\alpha}_{l}^{\prime}\right) K_{\eta}\left(\boldsymbol{\alpha}_{l}\right) K_{\eta}\left(\boldsymbol{\alpha}_{l}^{\prime}\right) d \boldsymbol{\alpha}_{l} d \boldsymbol{\alpha}_{l}^{\prime}
$$

We suppose that $0<\eta \leq \min \lambda_{j}$. Then, by (6.1), $\Psi_{l}\left(\boldsymbol{\alpha}_{l}, \boldsymbol{\alpha}_{l}^{\prime}\right)=S_{0}\left(\boldsymbol{\lambda}_{l}\left(\boldsymbol{\alpha}_{l}^{\prime}-\boldsymbol{\alpha}_{l}\right)\right)$ and
$\mathcal{J}_{l}=\int_{\widetilde{\mathcal{F}}_{l}} \int_{\widetilde{\mathcal{F}}_{l}} \Phi_{l}\left(\boldsymbol{\alpha}_{l}\right) \Phi_{l}\left(-\boldsymbol{\alpha}_{l}^{\prime}\right) S_{0}\left(\boldsymbol{\lambda}_{l}\left(\boldsymbol{\alpha}_{l}^{\prime}-\boldsymbol{\alpha}_{l}\right)\right) e\left(\boldsymbol{\beta}_{l}\left(\boldsymbol{\alpha}_{l}^{\prime}-\boldsymbol{\alpha}_{l}\right)\right) K_{\eta}\left(\boldsymbol{\alpha}_{l}\right) K_{\eta}\left(\boldsymbol{\alpha}_{l}^{\prime}\right) d \boldsymbol{\alpha}_{l} d \boldsymbol{\alpha}_{l}^{\prime}$.

The case $l=1$ is particularly simple. Here $\mathfrak{F}_{1}=\mathbb{R}^{r-1}$. Hence, by (5.3) and a consideration of the underlying diophantine problem, we see that when $l=1$, the integral in (6.15) does not exceed the number of solutions of the system of inequalities

$$
\begin{equation*}
\left|\lambda_{j} w-x_{j}\right|<\eta, \quad\left|\lambda_{j} w-y_{j}\right|<\eta \quad(2 \leq j \leq r) \tag{6.16}
\end{equation*}
$$

in integers $w, x_{j}, y_{j}$, all in the interval $[1, N]$. For $\eta \leq 1$, the value of $w$ fixes $x_{j}$ and $y_{j}$ within $O(1)$ possibilities. Hence, $\mathcal{J}_{1} \ll N$, and by (6.14) and (6.8), it follows that

$$
\begin{equation*}
\mathcal{I}\left(\mathfrak{h}_{1}\right) \ll N U_{1}^{-1 / 2} \tag{6.17}
\end{equation*}
$$

The cases $2 \leq l \leq r$ can be treated along similar lines. We return to (6.15) and carry out all integrations against $\alpha_{l+1}, \ldots, \alpha_{r}, \alpha_{l+1}^{\prime}, \ldots, \alpha_{r}^{\prime}$, with the aid of 5.3. On writing $\tilde{\boldsymbol{\alpha}}=\left(\alpha_{1}, \ldots \alpha_{l-1}\right)$,

$$
\mathfrak{G}_{l}=\mathfrak{M}\left(V_{1}\right) \times \ldots \times \mathfrak{M}\left(V_{l-1}\right)
$$

and

$$
\Theta_{l}\left(\tilde{\boldsymbol{\alpha}}, \tilde{\boldsymbol{\alpha}}^{\prime}\right)=\sum_{w, x_{j}, y_{j}} \prod_{j=l+1}^{r} \Upsilon_{\eta}\left(x_{j}-\lambda_{j} w-\beta_{j}\right) e\left(w \sum_{i=1}^{l-1} \lambda_{i}\left(\alpha_{i}^{\prime}-\alpha_{i}\right)\right),
$$

where the summation is over all $w \in \mathcal{S}_{0}, x_{j}, y_{j} \in \mathcal{S}$ for $1 \leq j \leq l-1$ that satisfy (6.16) and lie in the interval $[1, N]$, one finds that 6.15 reduces to

$$
\mathcal{J}_{l}=\int_{\mathfrak{G}_{l}} \int_{\mathfrak{G}_{l}} \Theta_{l}\left(\tilde{\boldsymbol{\alpha}}, \tilde{\boldsymbol{\alpha}}^{\prime}\right) \prod_{j=1}^{l-1} S_{j}\left(\alpha_{j}\right) S_{j}\left(-\alpha_{j}^{\prime}\right) e\left(\beta_{j}\left(\alpha_{j}^{\prime}-\alpha_{j}\right)\right) K_{\eta}(\tilde{\boldsymbol{\alpha}}) K_{\eta}\left(\tilde{\boldsymbol{\alpha}}^{\prime}\right) d \tilde{\boldsymbol{\alpha}} d \tilde{\boldsymbol{\alpha}}^{\prime}
$$

All remaining exponential sums will now be estimated trivially. As we saw in the case $l=1$, there are only $O(N)$ tuples $w, x_{j}, y_{j}$ in the summation conditions for the sum defining $\Theta$. Hence, $\Theta \ll N$, and therefore, we may conclude from (6.4) that

$$
\mathcal{J}_{l} \ll N^{2 l-1} \prod_{j=1}^{l-1}\left(\int_{\mathfrak{M}\left(V_{j}\right)} K_{\eta}\left(\alpha_{j}\right) d \alpha_{j}\right)^{2} \ll N\left(V_{1} \ldots V_{l-1}\right)^{4}
$$

By (6.14) and (6.8), it follows that for $2 \leq l \leq r$ one has

$$
\begin{equation*}
\mathcal{I}\left(\mathfrak{h}_{l}\right) \ll N\left(V_{1} \ldots V_{l-1}\right)^{2} U_{l}^{-1 / 2} \ll N U_{l}^{-1 / 10} \tag{6.18}
\end{equation*}
$$

Finally, we consider $\mathfrak{h}^{*}$. Here we apply Schwarz's inequality directly to (6.5) to obtain

$$
\mathcal{I}\left(\mathfrak{h}^{*}\right) \leq\left(\mathcal{I}_{1} \mathcal{I}_{2}\right)^{1 / 2}
$$

with

$$
\mathcal{I}_{1}=\int_{\mathfrak{h}^{*}}\left|S_{0}(-\boldsymbol{\lambda} \boldsymbol{\alpha})\right|^{2} K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}, \quad \mathcal{I}_{2}=\int_{\mathbb{R}^{r}}\left|S_{1}\left(\alpha_{1}\right) \ldots S_{r}\left(\alpha_{r}\right)\right|^{2} K_{\eta}(\boldsymbol{\alpha}) d \boldsymbol{\alpha}
$$

The second integral factorizes into

$$
\int_{-\infty}^{\infty}\left|S_{j}(\alpha)\right|^{2} K_{\eta}(\alpha) d \alpha=\#\left\{n \leq N: n \in \mathcal{S}_{j}\right\} \ll N
$$

whence $\mathcal{I}_{2} \ll N^{r}$. For $\mathcal{I}_{1}$, we substitute $\alpha_{0}=-\boldsymbol{\lambda} \boldsymbol{\alpha}$ for $\alpha_{r}$ as a new variable of integration. Then, with $\boldsymbol{\alpha}_{r}=\left(\alpha_{1}, \ldots, \alpha_{r-1}\right)$ as before,

$$
\mathcal{I}_{1} \leq \frac{1}{\lambda_{r}} \int_{\mathfrak{G}_{r}} \int_{\mathfrak{m}\left(V_{0}\right)}\left|S_{0}\left(\alpha_{0}\right)\right|^{2} K_{\eta}\left(\alpha_{r}\right) K_{\eta}\left(\boldsymbol{\alpha}_{r}\right) d \alpha_{0} d \boldsymbol{\alpha}_{r}
$$

Here $\alpha_{r}$ is the function of $\alpha_{0}$ and $\boldsymbol{\alpha}_{r}$ defined via $\alpha_{0}=-\lambda_{r} \alpha_{r}-\boldsymbol{\lambda}_{r} \boldsymbol{\alpha}_{r}$. We can rewrite the inner integral as

$$
\int_{\mathfrak{m}\left(V_{0}\right)}\left|S_{0}\left(\alpha_{0}\right)\right|^{2} K_{\eta}\left(\frac{\alpha_{0}+\boldsymbol{\lambda}_{r} \boldsymbol{\alpha}_{r}}{\lambda_{r}}\right) d \alpha_{0} .
$$

The Fourier transform of $\alpha \mapsto K_{\eta}\left(\lambda_{r}^{-1}\left(\alpha+\boldsymbol{\lambda}_{r} \boldsymbol{\alpha}_{r}\right)\right)$ is $\beta \mapsto \Upsilon_{\eta}\left(\lambda_{r} \beta\right) \lambda_{r} e\left(\beta \boldsymbol{\lambda}_{r} \boldsymbol{\alpha}_{r}\right)$. Hence, by (6.2), when $\eta \leq \lambda_{r}$,

$$
\int_{\mathfrak{m}\left(V_{0}\right)}\left|S_{0}\left(\alpha_{0}\right)\right|^{2} K_{\eta}\left(\frac{\alpha_{0}+\boldsymbol{\lambda}_{r} \boldsymbol{\alpha}_{r}}{\lambda_{r}}\right) d \alpha_{0}=\lambda_{r} \int_{\mathfrak{n}\left(V_{0}\right)}\left|S_{0}(\alpha)\right|^{2} d \alpha
$$

When combined with (6.3) and (6.4), it follows that

$$
\mathcal{I}_{1} \ll N U_{0}^{-1} \int_{\mathfrak{G}_{r}} K_{\eta}\left(\boldsymbol{\alpha}_{r}\right) d \boldsymbol{\alpha}_{r} \ll N^{2-r} U_{0}^{-1}\left(V_{1} \ldots V_{r-1}\right)^{2} \ll N^{2-r} U_{0}^{-1 / 2}
$$

Collecting together, we deduce that $\mathcal{I}\left(\mathfrak{h}^{*}\right) \ll N U_{0}^{-1 / 4}$. By (6.10), (6.11), (6.17) and (6.18) we finally confirm the bound $\mathcal{I}(\mathfrak{k})=o(N)$ as $N \mapsto \infty$ When combined with (5.6) and Lemma 5 , we have reached the following conclusion.

Theorem 6. Let $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$ denote extremal sequences. Let $\lambda_{1}, \ldots, \lambda_{r}$ be positive real numbers such that $1, \lambda_{1}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. Define $P_{\eta}(N, \boldsymbol{\beta})$ by (5.5). There exists a function $U(N)$ with $U(N) \rightarrow \infty$ as $N \rightarrow \infty$, and such that

$$
P_{\eta}(N, \boldsymbol{\beta})=\varrho \kappa^{-1} \eta^{r} N+O\left(N U^{-1}\right)
$$

holds uniformly in $0<\eta \leq \min \lambda_{j}$ and $|\boldsymbol{\beta}| \leq \sqrt{N}$. Here $\varrho$ and $\kappa$ are defined by (5.4).

We close this section with a brief comment on the method. The kernel in our Fourier approach factorizes, as in (5.3). It seems very important to use such a kernel. This property has been used ubiquitously, but most notably in (6.12). The second factor in (6.14) reduces, at least in part, to a diophantine system of linear inequalities. This provides a crucial saving over more traditional methods of estimation which would require $2 r+1$ variables for $r$ inequalities, at the very least. The price we have to pay for this is rather low; it is the tedious and unsymmetric construction of the functions $V_{j}$ in (6.9), with consequential asymmetries in the sets $\mathfrak{h}_{1}, \ldots, \mathfrak{h}_{l}$ later. One can use the techniques of this section also in the context of diophantine equations. For example, one can work along the lines of sections 5 and 6 of this paper to give an alternative, more direct proof of the $k$-tuple-Theorem 1.7 from our recent memoir [5].

## 7. The main theorem

The weights can be removed from the counting function $P_{\eta}(N, \boldsymbol{\beta})$ in Theorem 6, and we may consider more general linear forms than those covered by the methods
of sections 5 and 6 . With this in view, let $\Lambda=\left(\lambda_{i j}\right)_{\substack{1 \leq i \leq r \\ 0 \leq j \leq r}}$ denote a real matrix, and fix a choice of extremal sequences $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$. For real numbers $\tau_{j}>0$ define $Z_{\boldsymbol{\tau}}(N, \boldsymbol{\mu})$ as the number of solutions of the system

$$
\begin{equation*}
\left|\sum_{j=0}^{r} \lambda_{i j} x_{j}-\mu_{i}\right| \leq \tau_{i} \quad(1 \leq i \leq r) \tag{7.1}
\end{equation*}
$$

with $x_{j} \in \mathcal{S}_{j}$ and $1 \leq x_{j} \leq N$ for $0 \leq j \leq r$.
Theorem 7. Suppose that the real $r \times(r+1)$-matrix $\Lambda$ is highly non-singular, positive and of integral rank 0 . Then there exists a constant $c>0$, depending only on $\Lambda$, and such that the asymptotic formula

$$
Z_{\boldsymbol{\tau}}(N, \boldsymbol{\mu})=c \varrho N \tau_{1} \tau_{2} \ldots \tau_{r}+o(N)
$$

holds uniformly for all $|\boldsymbol{\mu}| \leq \sqrt{N}$ and for all positive $\tau_{j}$ that are sufficiently small in terms of $\Lambda$.

Proof. Choose a matrix $\mathrm{A} \in G L_{r}(\mathbb{R})$ such that $\mathrm{A} \Lambda$ takes the shape (2.2); then all $\lambda_{j}$ are positive, and by the remark at the end of section 2 , the numbers $1, \lambda_{1}, \ldots, \lambda_{r}$ are linearly independent over $\mathbb{Q}$. Let

$$
\mathcal{D}=\left\{\boldsymbol{\alpha} \in \mathbb{R}^{r}:\left|\alpha_{j}-\mu_{j}\right| \leq \tau_{j} \quad(1 \leq j \leq r)\right\}
$$

Then $A \mathcal{D}$ is a parallelotope, and the inequalities (7.1) are equivalent with the condition that

$$
\left(\lambda_{1} x_{0}-x_{1}, \lambda_{2} x_{0}-x_{2}, \ldots, \lambda_{r} x_{0}-x_{r}\right) \in \mathrm{A} \mathcal{D} .
$$

Choose $U(N)$ in accordance with Theorem 6, and let $\eta=U^{-1 / 4 r}$. Suppose that $\tau_{j} \geq \eta^{1 / 4}$ holds for all $j$. Then put

$$
\mathcal{D}^{+}=\left\{\boldsymbol{\alpha}:\left|\alpha_{j}-\mu_{j}\right| \leq \tau_{j}+\sqrt{\eta}\right\}, \quad \mathcal{D}^{-}=\left\{\boldsymbol{\alpha}:\left|\alpha_{j}-\mu_{j}\right| \leq \tau_{j}-\sqrt{\eta}\right\}
$$

Now consider the set of lattice points $\mathcal{B}^{ \pm}=\eta \mathbb{Z}^{r} \cap \mathrm{AD} \mathcal{D}^{ \pm}$. The volume of $\mathcal{D}^{ \pm}$is $\tau_{1} \tau_{2} \ldots \tau_{r}\left(1+O\left(\eta^{1 / 4}\right)\right)$, and hence the volume $\mathrm{A} \mathcal{D}^{ \pm}$equals $|\operatorname{det} \mathrm{A}| \tau_{1} \tau_{2} \ldots \tau_{r}(1+$ $\left.O\left(\eta^{1 / 4}\right)\right)$. By standard lattice point arguments, we find that

$$
\begin{equation*}
\# \mathcal{B}^{ \pm}=|\operatorname{det} A| \tau_{1} \ldots \tau_{r} \eta^{-r}\left(1+O\left(\eta^{1 / 4}\right)\right) \tag{7.2}
\end{equation*}
$$

Now write

$$
Z^{ \pm}=\sum_{\boldsymbol{\beta} \in \mathcal{B}^{ \pm}} P_{\eta}(N, \boldsymbol{\beta})
$$

By (5.1) and (5.5), one readily confirms that whenever $N$ is large one has

$$
Z^{-} \leq Z_{\boldsymbol{\tau}}(N, \boldsymbol{\beta}) \leq Z^{+}
$$

By Theorem 6 and (7.2), one finds that for some constant $c$ one has

$$
Z^{ \pm}=\varrho c \tau_{1} \ldots \tau_{r} N+O\left(N U^{-3 / 4}\right)+O\left(N \tau_{1} \tau_{2} \ldots \tau_{r} \eta^{1 / 4}\right)
$$

Theorem 7 is now immediate. This also proves Theorem 1 because the square-free numbers form an extremal sequence.

It is perhaps of interest that Theorem 7 may be generalized to the case of nonzero integral rank. In such a situation there are hidden diophantine equations in the system (7.1), as we pointed out in $\S 2$. The asymptotic formula for $Z_{\boldsymbol{\tau}}(N, \boldsymbol{\mu})$ will reflect this with the appearance of a suitable singular series that, however, is not necessarily positive. When $\Lambda$ is of integral rank $r$, then the desired formula is actually contained in [5], chapter 5, and the intermediate cases can be handled by the methods of this paper and the intersection principles of [5], chapter 4.

## 8. Outroduction

Suppose that $0<\lambda<\frac{1}{2}$ is an irrational number, and let $Z_{\tau}(N)$ denote the number of solutions of $|m-\lambda n|<\tau$ with $1 \leq n, m \leq N$. This is the simplest situation where Theorem 7 is applicable. The work of Hardy and Littlewood [13], section 4.4 is readily modified to show that in the current context the error term in the formula $Z_{\tau}(N)=2 \tau N+o(N)$ is sharp: for any function $T(N)$ increasing to infinity with $N$, there exists an irrational $\lambda \in\left(0, \frac{1}{2}\right)$ such that $\left|Z_{\tau}(N)-2 \tau N\right|>N / T(N)$ holds on a sequences of values for $N$ tending to infinity. In particular, the error term in Theorem 7 crucially depends on $\lambda$. One would expect that this is so for any $r$, and any choice of extremal sequences $\mathcal{S}_{0}, \ldots, \mathcal{S}_{r}$. However, a more thorough discussion of this matter requires another development of our methods that we intend to communicate on a future occasion.

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