# On Zaremba's conjecture 

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#### Abstract

Zaremba's 1971 conjecture predicts that every integer appears as the denominator of a finite continued fraction whose partial quotients are bounded by an absolute constant. We confirm this conjecture for a set of density one.


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## 1. Introduction

1.1. Statements of the main theorems. For a fixed finite set $\mathcal{A} \subset \mathbb{N}$, which we call an alphabet, let $\mathfrak{C}_{\mathcal{A}}$ denote the collection of all $x \in(0,1)$ whose continued fraction expansion

$$
x=\left[a_{1}, a_{2}, \ldots, a_{k}, \ldots\right]=\frac{1}{a_{1}+\frac{1}{a_{2}+\ddots+\frac{1}{a_{k}+\ddots}}},
$$

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has all partial quotients $a_{j}$ belonging to the alphabet $\mathcal{A}$. Any $x \in \mathfrak{C}_{\mathcal{A}}$ is uniformly badly approximable, in the sense that its partial quotients $a_{j}$ are all bounded by

$$
A:=\max \mathcal{A}
$$

When $A$ is an absolute constant, we call such a number absolutely Diophantine (of height $A$ ). That is, when we speak of numbers being absolutely Diophantine, the height $A$ is fixed in advance.

Let $\mathfrak{R}_{\mathcal{A}}$ denote the set of partial convergents to $\mathfrak{C}_{\mathcal{A}}$, that is,

$$
\mathfrak{R}_{\mathcal{A}}:=\left\{\frac{b}{d}=\left[a_{1}, a_{2}, \ldots, a_{k}\right]: 0<b<d,(b, d)=1, \text { and } \forall j, a_{j} \in \mathcal{A}\right\},
$$

and let $\mathfrak{D}_{\mathcal{A}} \subset \mathbb{N}$ be the set of denominators of fractions in $\mathfrak{\Re}_{\mathcal{A}}$,

$$
\mathfrak{D}_{\mathcal{A}}:=\left\{d \in \mathbb{N}: \exists(b, d)=1 \text { with } \frac{b}{d} \in \mathfrak{R}_{\mathcal{A}}\right\} .
$$

In 1971, S. K. Zaremba formulated the following assertion.
Conjecture 1.1 (Zaremba [Zar72, p. 76]). Every positive integer is the denominator of a reduced absolutely Diophantine fraction.

That is, the conjecture predicts the existence of some integer $A>1$ so that

$$
\mathfrak{D}_{\{1,2, \ldots, A\}}=\mathbb{N} .
$$

In fact, Zaremba suggested that $A=5$ may already be sufficient.
Zaremba's conjecture has important applications to numerical integration and pseudorandom number generation, producing collections of points of optimal discrepancy; see, e.g., the surveys [Nie78] and [Kon13]. Our main result is the following

Theorem 1.2. Almost every positive integer is the denominator of a reduced absolutely Diophantine fraction. That is, there exists an $A>1$ so that

$$
\frac{1}{N} \#\left(\mathfrak{D}_{\{1,2, \ldots, A\}} \cap[1, N]\right) \rightarrow 1
$$

as $N \rightarrow \infty$. In particular, $A=50$ suffices.
A more refined conjecture was stated by Hensley in 1996. The set $\mathfrak{C}_{\mathcal{A}} \subset$ $(0,1)$ is a Cantor-like fractal; let

$$
\delta_{\mathcal{A}}:=\operatorname{H} \cdot \operatorname{dim}\left(\mathfrak{C}_{\mathcal{A}}\right) \in[0,1]
$$

be its Hausdorff dimension. This dimension can be 0 only if $|\mathcal{A}|=1$; since we assume $\mathcal{A}$ is finite, $\delta_{\mathcal{A}}<1$. Allowing a finite number of exceptions in Zaremba's conjecture, Hensley asserts the following

Conjecture 1.3 (Hensley [Hen96, Conj. 3, p. 16]). The set of denominators $\mathfrak{D}_{\mathcal{A}}$ contains every sufficiently large integer if and only if the corresponding dimension $\delta_{\mathcal{A}}$ exceeds $1 / 2$.

As stated, Hensley's conjecture is false. For example, consider the alphabet $\mathcal{A}=\{2,4,6,8,10\}$. By implementing an algorithm due to Jenkinson and Pollicott [JP01], ${ }^{1}$ we have estimated its dimension to be $\delta_{\mathcal{A}} \approx 0.517>1 / 2$. Nevertheless, arbitrarily large numbers are missing from $\mathfrak{D}_{\mathcal{A}}$. Indeed, it is elementary to verify that

$$
\begin{equation*}
\mathfrak{D}_{\mathcal{A}}(\bmod 4) \equiv\{0,1,2\} ; \tag{1.4}
\end{equation*}
$$

see Remark 1.33.
We propose the following alternative to Hensley's conjecture, borrowing language from Hilbert's $11^{\text {th }}$ problem on representations of numbers by quadratic forms. We call an integer $d$ admissible (for $\mathcal{A}$ ) if it passes all finite local obstructions:

$$
\begin{equation*}
\forall q>1, d \in \mathfrak{D}_{\mathcal{A}}(\bmod q) . \tag{1.5}
\end{equation*}
$$

Remark 1.6. Admissibility can be checked using only one modulus $q=$ $q(\mathcal{A})$; see Remark 1.33.

Let $\mathfrak{A}_{\mathcal{A}}$ denote the set of all admissible numbers,

$$
\mathfrak{A}_{\mathcal{A}}:=\{d \in \mathbb{Z}:(1.5) \text { holds }\} .
$$

We say $d$ is represented (by $\mathcal{A}$ ) if $d \in \mathfrak{D}_{\mathcal{A}}$. The multiplicity of a denominator $d$ is the number of coprime numerators $0<b<d$ with $b / d \in \mathfrak{R}_{\mathcal{A}}$. Clearly $d$ is represented if and only if its multiplicity is positive.

Conjecture 1.7. If the dimension $\delta_{\mathcal{A}}$ exceeds $1 / 2$, then the set of denominators $\mathfrak{D}_{\mathcal{A}}$ contains every sufficiently large admissible integer.

We interpret this conjecture as a local-global principle, where the dimension condition and "sufficiently large" are local obstructions at infinity. Theorem 1.2 follows from the following more refined approximation to Conjecture 1.7.

Theorem 1.8. There exists some $\delta_{0}<1$ (see (1.21)) so that, if the dimension $\delta_{\mathcal{A}}$ exceeds $\delta_{0}$, then the set of denominators $\mathfrak{D}_{\mathcal{A}}$ contains almost every admissible integer. More precisely, there is a constant $c=c(\mathcal{A})>0$ so that

$$
\begin{equation*}
\frac{\#\left(\mathfrak{D}_{\mathcal{A}} \cap[N / 2, N]\right)}{\#\left(\mathfrak{A}_{\mathcal{A}} \cap[N / 2, N]\right)}=1+O\left(e^{-c \sqrt{\log N}}\right), \tag{1.9}
\end{equation*}
$$

[^0]as $N \rightarrow \infty$. Furthermore, each d produced above appears with multiplicity
\[

$$
\begin{equation*}
\text { multiplicity }(d) \gg N^{2 \delta_{\mathcal{A}}-\frac{1001}{1000}} . \tag{1.10}
\end{equation*}
$$

\]

The constant $c$ in (1.9) is effectively computable, and the implied constants above depend only on $\mathcal{A}$.

Some remarks are in order.
Remark 1.11. There exist alphabets $\mathcal{A}$ with $\delta_{\mathcal{A}}$ arbitrarily close to 1 , so Theorem 1.8 is not vacuous. Indeed, Hensley [Hen92] gives the asymptotic expansion

$$
\begin{equation*}
\delta_{\{1,2, \ldots, A\}}=1-\frac{6}{\pi^{2} A}-\frac{72 \log A}{\pi^{4} A^{2}}+O\left(\frac{1}{A^{2}}\right) . \tag{1.12}
\end{equation*}
$$

Remark 1.13. The number $1 / 2$ in Conjectures 1.3 and 1.7 cannot be reduced. Hensley [Hen89] showed that the truncated set of rationals

$$
\mathfrak{R}_{\mathcal{A}}(N):=\left\{\frac{b}{d} \in \mathfrak{R}_{\mathcal{A}}:(b, d)=1,0<b<d<N\right\}
$$

has cardinality

$$
\begin{equation*}
\# \mathfrak{R}_{\mathcal{A}}(N) \asymp N^{2 \delta_{\mathcal{A}}}, \tag{1.14}
\end{equation*}
$$

whence it follows immediately that

$$
\#\left(\mathfrak{D}_{\mathcal{A}} \cap[1, N]\right) \ll N^{2 \delta_{\mathcal{A}}} .
$$

Thus if $\delta_{\mathcal{A}}<1 / 2$, then certainly $\mathfrak{D}_{\mathcal{A}}$ is too thin a subset of the integers to contain even one admissible arithmetic progression.

Remark 1.15. The best previously known estimate,

$$
\begin{equation*}
\#\left(\mathfrak{D}_{\mathcal{A}} \cap[1, N]\right) \gg N^{\delta_{\mathcal{A}}}, \tag{1.16}
\end{equation*}
$$

was proved by Hensley [Hen06, Th. 3.2] and follows easily from his estimate (1.14). In particular, as long as $|\mathcal{A}|>1$, the set $\mathfrak{D}_{\mathcal{A}}$ grows at least polynomially. Moreover, taking $\mathcal{A}$ large so that $\delta_{\mathcal{A}}>1-\varepsilon$, one can already produce at least $N^{1-\varepsilon}$ denominators in $\mathfrak{D}_{\mathcal{A}}$ up to $N$.

Remark 1.17. We explain in Remark 1.33 below that for any $A \geq 2$, the alphabet $\{1,2, \ldots, A\}$ has no finite local obstructions; that is, $\mathfrak{A}_{\{1,2, \ldots, A\}}=\mathbb{Z}$. This is why the statement of Theorem 1.2 needs no mention of admissibility. Moreover, the dimension $\delta_{\{1,2\}}$ is known [Goo41], [Bum85], [JP01] to be

$$
\begin{equation*}
\delta_{\{1,2\}} \approx 0.531 \cdots, \tag{1.18}
\end{equation*}
$$

which obviously exceeds $1 / 2$. Conjecture 1.7 then implies that $\mathfrak{D}_{\{1,2\}}$ already contains every sufficiently large number, as was conjectured by Hensley [Hen96].

Remark 1.19. An earlier version ${ }^{2}$ of this paper also proved two weaker results made obsolete by Theorem 1.2, namely that for sufficiently large $\delta_{\mathcal{A}}$,
(i) $\mathfrak{D}_{\mathcal{A}}$ contains a positive proportion of numbers, and
(ii) that $\mathfrak{D}_{\mathcal{A}}$ contains almost every admissible number, without giving the rate in (1.9).
At the request of the referee to shorten the paper, we have removed these intermediary results (and of course the methods used to obtain them). We invite the interested reader to peruse the original arXiv posting for the details. Note also that some results of this paper have been announced in [BK11].

Remark 1.20. The value of $\delta_{0}$ in Theorem 1.8 coming from our proof is

$$
\begin{equation*}
\delta_{0}=307 / 312 \approx 0.984 \tag{1.21}
\end{equation*}
$$

We have made no effort to optimize this quantity, as can surely be done with a modicum of effort. In fact, Frolenkov and Kan [FK13] have since sharpened our method to prove the weaker statement ${ }^{3}$ that $\mathfrak{D}_{\mathcal{A}}$ contains a positive proportion of numbers whenever $\delta_{\mathcal{A}}>\delta_{0}$ with the improved range $\delta_{0}=5 / 6 \approx 0.833$. It does not seem likely that these methods can achieve the full range $\delta_{0}=1 / 2$ without significant new ideas. We have estimated the dimension $\delta_{\mathcal{A}}$ corresponding to the alphabet $\mathcal{A}=\{1,2, \ldots, 49,50\}$ to be about 0.986 , exceeding (1.21), whereas the alphabet $\mathcal{A}=\{1,2,3,4,5\}$ is known [Jen04] to have dimension $\delta_{\mathcal{A}}>5 / 6$.

Remark 1.22. The rate in (1.9) can be improved to a power savings, that is, an error term of the form $O\left(N^{-c}\right)$; see [Bou12].

Although our main result requires large dimension, we are also able to sharpen the best previously known estimate (1.16) in the full range $\delta_{\mathcal{A}}>1 / 2$.

Theorem 1.23. Write $\delta$ for $\delta_{\mathcal{A}}$. Then for any $\varepsilon>0$,

$$
\begin{equation*}
\#\left(\mathfrak{D}_{\mathcal{A}} \cap[1, N]\right) \ggg{ }_{\varepsilon} N^{\delta+\frac{(2 \delta-1)(1-\delta)}{5-\delta}-\varepsilon} \tag{1.24}
\end{equation*}
$$

as $N \rightarrow \infty$. This bound improves on (1.16), as long as $\delta>1 / 2$.
Remark 1.25. The improvement here is quite modest: for $\mathcal{A}=\{1,2\}$, the exponent $\delta_{\mathcal{A}} \approx 0.531$ in (1.18) and (1.16) is replaced in (1.24) by 0.537 . We have again made no attempt to optimize the exponent in (1.24), seeking just any power gain.

We state the multiplicity bound (1.10) to give another application to pseudorandom numbers. Specifically, in the (homogeneous) linear congruential

[^1]method, optimal conditions require a prime $d$ and a primitive root $b(\bmod d)$ so that the fraction $b / d$ is absolutely Diophantine (see [Kon13]). Then the pseudorandom map with modulus $d$ and multiplier $b$, that is, $x \mapsto b x(\bmod d)$, has asymptotically optimal serial correlation of pairs.

Theorem 1.26. There exist infinitely many primes $d$ with primitive roots $b(\bmod d)$ so that the fractions $b / d$ are absolutely Diophantine.

The number of such prime $d$ up to $N$ provided by our proof is on the order of $N(\log N)^{-2}$. Theorem 1.26 is an easy corollary of Theorem 1.8. In fact, if $\mathcal{A}=\{1,2, \ldots, A\}$ has dimension $\delta_{\mathcal{A}}$ exceeding $\delta_{0}$ as in Theorem 1.8, then the fractions $b / d$ produced in Theorem 1.26 can be taken to have all partial quotients bounded by $A+1$.
1.2. Reformulation and admissibility. It is an old and trivial (but for our purposes crucial) observation that

$$
\frac{b}{d}=\left[a_{1}, \ldots, a_{k}\right]
$$

is equivalent to

$$
\left(\begin{array}{ll}
* & b  \tag{1.27}\\
* & d
\end{array}\right)=\left(\begin{array}{cc}
0 & 1 \\
1 & a_{1}
\end{array}\right)\left(\begin{array}{cc}
0 & 1 \\
1 & a_{2}
\end{array}\right) \cdots\left(\begin{array}{cc}
0 & 1 \\
1 & a_{k}
\end{array}\right) .
$$

This observation will allow us to explain all local obstructions, as follows. In light of (1.27), let

$$
\mathcal{G}_{\mathcal{A}} \subset \mathrm{GL}(2, \mathbb{Z})
$$

be the semigroup generated by the matrices

$$
\left(\begin{array}{ll}
0 & 1  \tag{1.28}\\
1 & a
\end{array}\right)
$$

for $a \in \mathcal{A}$. Then the orbit

$$
\begin{equation*}
\mathcal{O}_{\mathcal{A}}:=\mathcal{G}_{\mathcal{A}} \cdot e_{2} \tag{1.29}
\end{equation*}
$$

of $e_{2}=(0,1)^{t}$ under $\mathcal{G}_{\mathcal{A}}$ corresponds to $\mathfrak{\Re}_{\mathcal{A}}$; that is, if $\gamma=\left(\begin{array}{cc}a & b \\ c & d\end{array}\right)$, then $\gamma \cdot e_{2}=$ $(b, d)^{t}$. Moreover, taking the inner product of this orbit with $e_{2}$ picks off the value of $d$, that is $\left\langle\gamma \cdot e_{2}, e_{2}\right\rangle=d$, and

$$
\begin{equation*}
\left\langle\mathcal{O}_{\mathcal{A}}, e_{2}\right\rangle=\left\langle\mathcal{G}_{\mathcal{A}} \cdot e_{2}, e_{2}\right\rangle \tag{1.30}
\end{equation*}
$$

is precisely $\mathfrak{D}_{\mathcal{A}}$ (with multiplicity). Zaremba's conjecture can then be reformulated as: For some finite alphabet $\mathcal{A}$,

$$
\mathbb{N} \subset\left\langle\mathcal{G}_{\mathcal{A}} \cdot e_{2}, e_{2}\right\rangle
$$

For convenience we pass from $\mathcal{G}_{\mathcal{A}}$ to its determinant one subsemigroup

$$
\begin{equation*}
\Gamma_{\mathcal{A}}=\mathcal{G}_{\mathcal{A}} \cap \mathrm{SL}_{2} \subset \mathrm{SL}_{2}(\mathbb{Z}) \tag{1.31}
\end{equation*}
$$

which is (freely and finitely) generated by the matrix products

$$
\left(\begin{array}{cc}
0 & 1  \tag{1.32}\\
1 & a
\end{array}\right) \cdot\left(\begin{array}{cc}
0 & 1 \\
1 & a^{\prime}
\end{array}\right)
$$

for $a, a^{\prime} \in \mathcal{A}$. The orbit $\mathcal{O}_{\mathcal{A}}$ is recovered as a finite union of "coset" orbits

$$
\mathcal{O}_{\mathcal{A}}=\Gamma_{\mathcal{A}} \cdot e_{2} \cup \bigcup_{a \in \mathcal{A}} \Gamma_{\mathcal{A}} \cdot\left(\begin{array}{ll}
0 & 1 \\
1 & a
\end{array}\right) e_{2} .
$$

Remark 1.33. It now follows from Strong Approximation [MVW84] and Goursat's Lemma (see the discussion in [Kon13, §2.2]) that the reduction of $\Gamma_{\mathcal{A}} \bmod q$ is all of $\mathrm{SL}_{2}(q)$ for all $q$ coprime to a certain "bad" modulus $\mathfrak{B}$. Here $\mathfrak{B}$ is effectively computable and depends only on $\Gamma_{\mathcal{A}}$, that is, on $\mathcal{A}$. Moreover, $\mathfrak{B}$ can be chosen so that for any $q \equiv 0(\mathfrak{B})$, the reduction $\Gamma_{\mathcal{A}}(\bmod q)$ is the full pre-image of $\Gamma_{\mathcal{A}}(\bmod \mathfrak{B})$ under the projection $\operatorname{map} \mathbb{Z} / q \rightarrow \mathbb{Z} / \mathfrak{B}$. From the $\bmod \mathfrak{B}$ reductions of $\Gamma_{\mathcal{A}}$, it is elementary to read off the reductions of $\mathcal{O}_{\mathcal{A}}$, and hence all finite local obstructions in $\mathfrak{D}_{\mathcal{A}}$; see Remark 1.6. Moreover, for the alphabet $\mathcal{A}=\{1,2\}$, it is easy to see that $\mathfrak{B}=1$, that is, $\Gamma_{\mathcal{A}}(\bmod q)$ is already all of $\mathrm{SL}_{2}(q)$, for all $q>1$; see Remark 1.17. Indeed, the group generated by $\Gamma_{\mathcal{A}}$ (that is, allowing inverses) is all of $\mathrm{SL}_{2}(\mathbb{Z})$, and the two have the same projections mod $q$. Finally, we note that this is precisely the phenomenon responsible for (1.4) in the failure of Hensley's Conjecture 1.3.
1.3. An overview of the key ideas. An observation which we had made in a slightly different context [BK10] is that there is a certain bilinear (in fact multilinear) structure to (1.30), making the problem amenable to the HardyLittlewood circle method via Vinogradov's techniques for estimating bilinear forms. We now outline the key steps.

From now on, we treat $\mathcal{A}$ as fixed, dropping it from subscripts, writing $\delta=\delta_{\mathcal{A}}, \Gamma=\Gamma_{\mathcal{A}}$, etc. In light of (1.30), we would like to study the exponential sum

$$
\begin{equation*}
S_{N}(\theta):=\sum_{\substack{\gamma \in \Gamma \\\|\gamma\|<N}} e\left(\theta\left\langle\gamma e_{2}, e_{2}\right\rangle\right), \tag{1.34}
\end{equation*}
$$

where $\theta \in[0,1]$ and $\|\cdot\|$ is the Frobenius matrix norm, $\left\|\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)\right\|^{2}=a^{2}+b^{2}+$ $c^{2}+d^{2}$. Then the Fourier coefficient

$$
\begin{equation*}
R_{N}(d):=\widehat{S}_{N}(d)=\int_{0}^{1} S_{N}(\theta) e(-d \theta) d \theta=\sum_{\substack{\gamma \in \Gamma \\\|\gamma\|<N}} \mathbf{1}_{\left\{\left\langle\gamma e_{2}, e_{2}\right\rangle=d\right\}} \tag{1.35}
\end{equation*}
$$

is just the "representation number" of $d$ up to $N$, that is, its multiplicity. Of course if $R_{N}(d)>0$, then $d \in \mathfrak{D}$.

Note that by (1.14),

$$
\begin{equation*}
S_{N}(0)=\sum_{d} R_{N}(d)=\sum_{\substack{\gamma \in \Gamma \\\|\gamma\|<N}} 1 \asymp N^{2 \delta}, \tag{1.36}
\end{equation*}
$$

so if almost every $d \in[N / 2, N]$ is to be represented without much bias, it should occur with multiplicity roughly $N^{2 \delta-1}$.

Following the circle method, we decompose the integral in (1.35) into "major arcs" and "minor arcs," the former referring to frequencies $\theta$ quite near rationals with small denominators and the latter being the rest:

$$
R_{N}(d)=\left(\int_{\mathfrak{M}}+\int_{[0,1] \backslash \mathfrak{M}}\right) S_{N}(\theta) e(-d \theta) d \theta=\mathcal{M}_{N}(d)+\mathcal{E}_{N}(d)
$$

Here $\mathcal{M}_{N}$ is thought of as a "main" term and $\mathcal{E}_{N}$ is an "error" term, and the major arcs $\mathfrak{M}=\mathfrak{M}_{\mathcal{Q}}$ are given by

$$
\begin{equation*}
\mathfrak{M}_{\mathcal{Q}}=\bigcup_{q<\mathcal{Q}} \bigcup_{(a, q)=1}\left[\frac{a}{q}-\frac{\mathcal{Q}}{N}, \frac{a}{q}+\frac{\mathcal{Q}}{N}\right] \tag{1.37}
\end{equation*}
$$

where $\mathcal{Q}$ is roughly of size $e^{-c \sqrt{\log N}}$.
A key ingredient is to show that along the major arcs, $\theta=\frac{a}{q}+\beta \in \mathfrak{M}$, the function $S_{N}$ essentially splits into two pieces,

$$
\begin{equation*}
S_{N}\left(\frac{a}{q}+\beta\right) \sim \nu_{q}(a) \cdot \varpi(\beta) \tag{1.38}
\end{equation*}
$$

Here $\nu_{q}$ is a purely modular term and $\varpi$ is an archimedean one, which has the right order of magnitude on balls of certain size. It then follows that the main term $\mathcal{M}_{N}(d)$ also splits as a "singular series" $\mathfrak{S}$ times a "singular integral" $\Pi$,

$$
\mathcal{M}_{N}(d) \sim \mathfrak{S}(d) \Pi_{N}(d)
$$

where $\Pi$ gives the expected archimedean contribution, roughly

$$
\Pi_{N}(d) \gg N^{2 \delta-1}
$$

for $d \asymp N$, and the singular series $\mathfrak{S}$ controls the local obstructions. In particular, if $d \notin \mathfrak{A}$ is not admissible, then $\mathfrak{S}(d)=0$; otherwise, we have roughly that

$$
\mathfrak{S}(d) \asymp \prod_{p \nmid d}\left(1+\frac{1}{p^{2}-1}\right) \prod_{p \mid d}\left(1-\frac{1}{p+1}\right) \gg \frac{1}{\log \log d}
$$

The main ingredient in proving this is the renewal method in the thermodynamic formalism of Ruelle transfer operators (see Lalley [Lal89]), and the extension to "congruence" such established by Bourgain-Gamburd-Sarnak in [BGS11]. (We need here not just square-free but arbitrary moduli $q$, and we must also use the work of Bourgain-Varjú [BV12].)

With the major arcs controlled, if we could prove that the errors are individually bounded, $\left|\mathcal{E}_{N}(d)\right| \ll N^{2 \delta-1-\varepsilon}$, say, then we would conclude the full

Conjecture 1.7. We are not able to establish control of this quality individually, but do succeed on average, proving essentially that

$$
\begin{equation*}
\sum_{d \asymp N}\left|\mathcal{E}_{N}(d)\right|^{2} \ll N^{4 \delta-1} e^{-c \sqrt{\log N}} \tag{1.39}
\end{equation*}
$$

from which Theorem 1.8 follows by a standard argument.
Bounds of this type will follow from bounds on

$$
\begin{equation*}
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \tag{1.40}
\end{equation*}
$$

where we have decomposed the minor arcs $[0,1] \backslash \mathfrak{M}$ into the dyadic regions

$$
\begin{equation*}
W_{Q, K}:=\left\{\theta=\frac{a}{q}+\beta: q \asymp Q,(a, q)=1,|\beta| \asymp \frac{K}{N}\right\} . \tag{1.41}
\end{equation*}
$$

By Dirichlet's approximation theorem, the parameters $Q$ and $K$ vary in the range $Q<N^{1 / 2}$ and $K<\frac{N^{1 / 2}}{Q}$.

Unfortunately, we do not know how to obtain such strong bounds for the function $S_{N}$ as defined in (1.34). But taking a cue from Vinogradov (as we did in [BK10]), we work with a different function:

$$
\begin{equation*}
S_{N}(\theta)=\sum_{\substack{\gamma_{1} \in \Gamma \\\left\|\gamma_{1}\right\|<N^{1 / 2}}} \sum_{\substack{\gamma_{2} \in \Gamma \\\left\|\gamma_{2}\right\| \asymp N^{1 / 2}}} e\left(\theta\left\langle\gamma_{1} \gamma_{2} e_{2}, e_{2}\right\rangle\right) \tag{1.42}
\end{equation*}
$$

say. Since $\Gamma$ is a semigroup, this modified function, or rather its Fourier transform, continues to capture elements of $\mathfrak{D}$. Moreover, the bilinear nature of the problem, namely that $\left\langle\gamma_{1} \gamma_{2} e_{2}, e_{2}\right\rangle=\left\langle\gamma_{2} e_{2},{ }^{t} \gamma_{1} e_{2}\right\rangle$, allows us to separate variables.

It is here in the separation of variables and application of Cauchy-Schwarz that we replace the thin semigroup $\Gamma$ with all of $\mathrm{SL}_{2}(\mathbb{Z})$, a loss we can only tolerate if the dimension $\delta$ is large, at least some $\delta_{0}$. We are then lead to a more classical setting, and in certain large ranges of the pair $(Q, K)$ in (1.41), we can obtain the requisite cancelation. For slightly smaller values of $(Q, K)$, it is beneficial to decompose the sum further as

$$
S_{N}(\theta)=\sum_{\substack{\gamma_{1} \in \Gamma \\\left\|\gamma_{1}\right\| \simeq N^{1 / 2}}} \sum_{\substack{\gamma_{2} \in \Gamma \\\left\|\gamma_{2}\right\| \asymp N^{1 / 4}}} \sum_{\substack{\gamma_{3} \in \Gamma \\\left\|\gamma_{3}\right\| \simeq N^{1 / 4}}} e\left(\theta\left\langle\gamma_{1} \gamma_{2} \gamma_{3} e_{2}, e_{2}\right\rangle\right) .
$$

Continuing in this way, we handle every conceivable range of $(Q, K)$ by considering a sum of the form

$$
\begin{equation*}
S_{N}(\theta)=\sum_{\substack{\gamma_{1} \in \Gamma \\\left\|\gamma_{1}\right\| \not N^{1 / 2}}} \sum_{\substack{\gamma_{2} \in \Gamma \\\left\|\gamma_{2}\right\| \nsim N^{1 / 4}}} \cdots \sum_{\substack{\gamma_{J} \in \Gamma \\\left\|\gamma_{J}\right\| \approx N^{1 / 2}}} e\left(\theta\left\langle\gamma_{1} \gamma_{2} \cdots \gamma_{J} e_{2}, e_{2}\right\rangle\right), \tag{1.43}
\end{equation*}
$$

where $J \asymp \log \log N$, so that $\gamma_{J}$ is of large but constant size (independent of $N) .{ }^{4}$ Unfortunately, another problem has crept up: we can no longer control the size of the long product $\gamma_{1} \cdots \gamma_{J}$, which could have norm as large as $N(\log N)^{C}$.

To remedy this situation, we develop a bit of elementary linear algebra for $\Gamma$, showing that if the expanding vectors of two matrices are close, then their eigenvalues behave nearly multiplicatively. This forces us to concoct, for each $j=1, \ldots, J$, a certain special subset $\Xi_{j} \subset\left\{\gamma \in \Gamma:\|\gamma\| \asymp N^{1 / 2^{j}}\right\}$, all the elements of which have expanding eigenvectors pointing near a common direction (independent of $j$ ). We then simply use the pigeonhole principle to make sure all elements of $\Xi_{j}$ have almost the same eigenvalues. Moreover, we need to ensure that the representation $\gamma=\gamma_{1} \gamma_{2} \cdots \gamma_{J}$ in (1.43) is unique; that is, if

$$
\gamma_{1} \gamma_{2} \cdots \gamma_{J}=\gamma_{1}^{\prime} \gamma_{2}^{\prime} \cdots \gamma_{J}^{\prime}
$$

with $\gamma_{j}, \gamma_{j}^{\prime} \in \Xi_{j}$, then $\gamma_{j}=\gamma_{j}^{\prime}$ for all $j$. We do this by forcing each $\gamma_{j} \in \Xi_{j}$ to have the same size in the wordlength metric, again by pigeonhole.

Then the large product ensemble

$$
\Xi_{1} \cdot \Xi_{2} \cdots \Xi_{J}
$$

is a good substitute for $\{\gamma \in \Gamma:\|\gamma\| \asymp N\}$ to handle the minor arcs. Unfortunately, the concocted sets $\Xi_{j}$ are no longer amenable to the major arc methods! We rectify this by constructing a certain tiny set $\aleph \subset \Gamma$ with good modular/archimedean distribution properties and by prepending it to the product, forming

$$
\begin{equation*}
\Omega_{N}=\aleph \tilde{\Xi}_{1} \Xi_{2} \cdots \Xi_{J} . \tag{1.44}
\end{equation*}
$$

Here the size of $\Xi_{1}$ has been cut down a bit to $\tilde{\Xi}_{1}$ to make room for the set $\aleph$.
The "correct" definition of $S_{N}(\theta)$ is then to replace (1.34) by

$$
\begin{equation*}
S_{N}(\theta):=\sum_{\gamma \in \Omega_{N}} e\left(\theta\left\langle\gamma e_{2}, e_{2}\right\rangle\right), \tag{1.45}
\end{equation*}
$$

from which the argument follows as described above. In the end, we prove Theorem 1.8, and hence Theorem 1.2. As already mentioned, Theorem 1.26 is an easy corollary to Theorem 1.8.

Remark 1.46. One may ponder the flexibility of our methods in applications to other problems. For one, in particular, McMullen [McM09] has formulated the problem of producing many closed geodesics in a compact subset of the modular surface, defined over a fixed real quadratic number field $\mathbb{Q}(\sqrt{f})$. This is the same as producing many elements $\gamma \in \mathcal{G}$ so that $\operatorname{tr}(\gamma)^{2}-4$ has square-free part $f$. Specifically, McMullen asks whether there is a finite

[^2]alphabet $\mathcal{A}$ so that the set of traces in $\mathcal{G}_{\mathcal{A}}$ contains every sufficiently large admissible integer. Our use of Vinogradov's bilinear estimates relies crucially on the structure in (1.30) and does not apply as it stands to the problem of traces. We plan to return to this problem in the future.

The proof of Theorem 1.23 follows along completely different lines, and is inspired by the recent advances in projection theorems [Bou10]. The observation here is that the set $\mathfrak{D}$ has a certain "sum-set" structure. Namely, if $b / d \in \mathfrak{R}$ is a reduced fraction and $a \in \mathcal{A}$, then clearly

$$
\begin{equation*}
\frac{1}{a+\frac{b}{d}}=\frac{d}{b+a d} \in \mathfrak{R} . \tag{1.47}
\end{equation*}
$$

This implies that $b+a d \in \mathfrak{D}$ whenever $b / d \in \mathfrak{R}$ and $a \in \mathcal{A}$; we exploit this sum-set structure to produce the bound (1.24). We note further that the Discretized Ring Theorem [Bou03] can be used to get an exponent gain over the lower bound (1.16) even when $\delta \leq 1 / 2$.
1.4. Outline of the paper. In Section 2, we study the multiplicative properties of expanding eigen-values and -vectors for matrices in $\Gamma$. We use Section 3 to construct the main ensemble $\Omega_{N}$, reserving the construction of the leading set $\aleph$ for Section 8. The major arc analysis is carried out in Section 4, while the minor arc bounds are proved in Sections 5-6. Theorem 1.8 is then proved in Section 7, as is its corollary, Theorem 1.26. Lastly, we prove Theorem 1.23 in Section 9.

Notation. Throughout we use the following standard notation. We write $f \sim g$ to mean $f / g \rightarrow 1$. We use the Landau/Vinogradov notation $f=O(g)$ and $f \ll g$ synonymously to mean there exists an implied constant $C>0$ such that for $x$ sufficiently large, $f(x) \leq C g(x)$. Moreover, $f \asymp g$ denotes $f \ll g \ll f$. We allow the implied constants to depend at most on the fixed alphabet $\mathcal{A}$, unless otherwise specified. We also use the shorthand $e(x)=e^{2 \pi i x}$. The cardinality of a finite set $S$ is denoted both as $\# S$ and $|S|$, and the Lebesgue measure of an interval $\mathcal{I}$ is also $|\mathcal{I}|$. Throughout there are some constants $c, C>0$ which may change from line to line.

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## 2. Large matrix products

In this section, we develop some tools in large matrix products, reminiscent of the avalanche principle; see, e.g., [Bou05, Ch. 6] or [GS01, §2]. Recall that $\Gamma=\Gamma_{\mathcal{A}}$ is the semigroup generated by even words in the matrices $a \in \mathcal{A}$. An
easy induction shows that for $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in \Gamma, \gamma \neq I$, we have

$$
1 \leq a \leq \min (b, c) \leq \max (b, c)<d
$$

We use the Frobenius norm:

$$
\begin{equation*}
\|\gamma\|:=\sqrt{a^{2}+b^{2}+c^{2}+d^{2}} . \tag{2.1}
\end{equation*}
$$

Note that the trace and norm are comparable up to constants:

$$
\begin{equation*}
\frac{1}{2}\|\gamma\| \leq \operatorname{tr} \gamma \leq 2\|\gamma\|, \tag{2.2}
\end{equation*}
$$

as are the norm, sup-norm, and "second column" norm:

$$
\begin{equation*}
\|\gamma\|_{\infty}=d<\left|\gamma e_{2}\right|=\sqrt{b^{2}+d^{2}}<\|\gamma\|<2\left|\gamma e_{2}\right|<4\|\gamma\|_{\infty} . \tag{2.3}
\end{equation*}
$$

For $\gamma \in \Gamma$, let the expanding and contracting eigenvalues of $\gamma$ be $\lambda_{+}(\gamma)$ and $\lambda_{-}(\gamma)=1 / \lambda_{+}(\gamma)$, with corresponding normalized eigenvectors $v_{+}(\gamma)$ and $v_{-}(\gamma)$. Write $\lambda=\lambda_{+}$for the expanding eigenvalue, so that

$$
\lambda(\gamma)=\lambda_{+}(\gamma)=\frac{\operatorname{tr}(\gamma)+\sqrt{\operatorname{tr}(\gamma)^{2}-4}}{2} .
$$

Remark 2.4. Note that for all $\gamma \in \Gamma$, the eigenvalues are real, and $\lambda>1$ for $\gamma \neq I$. In particular, $\Gamma$ has no parabolic elements.

The goal of this section is to prove the following elementary but very useful observation.

Proposition 2.5. The eigenvalues of two matrices $\gamma, \gamma^{\prime} \in \Gamma$ with large norms behave essentially multiplicatively, subject to their expanding eigenvectors facing nearby directions. That is,

$$
\begin{equation*}
\lambda\left(\gamma \gamma^{\prime}\right)=\lambda(\gamma) \lambda\left(\gamma^{\prime}\right)\left[1+O\left(\left|v_{+}(\gamma)-v_{+}\left(\gamma^{\prime}\right)\right|+\frac{1}{\|\gamma\|^{2}}+\frac{1}{\left\|\gamma^{\prime}\right\|^{2}}\right)\right] \tag{2.6}
\end{equation*}
$$

Moreover, the expanding eigenvector of the product $\gamma \gamma^{\prime}$ faces a nearby direction to that of the first $\gamma$ (and the same in reverse):

$$
\begin{equation*}
\left|v_{+}\left(\gamma \gamma^{\prime}\right)-v_{+}(\gamma)\right| \ll \frac{1}{\|\gamma\|^{2}} \quad \text { and } \quad\left|v_{-}\left(\gamma \gamma^{\prime}\right)-v_{-}\left(\gamma^{\prime}\right)\right| \ll \frac{1}{\left\|\gamma^{\prime}\right\|^{2}} \tag{2.7}
\end{equation*}
$$

The implied constants above are absolute.
Proof. For $\gamma$ large, we have

$$
\begin{equation*}
\lambda(\gamma)=\frac{\operatorname{tr}(\gamma)+\sqrt{\operatorname{tr}(\gamma)^{2}-4}}{2}=\operatorname{tr}(\gamma)+O\left(\frac{1}{\|\gamma\|}\right) \tag{2.8}
\end{equation*}
$$

and

$$
\begin{aligned}
& v_{+}(\gamma)=\frac{\left(b, \lambda_{+}(\gamma)-a\right)}{\sqrt{b^{2}+\left(\lambda_{+}(\gamma)-a\right)^{2}}}=\frac{(b, d)}{\sqrt{b^{2}+d^{2}}}+O\left(\frac{1}{\|\gamma\|^{2}}\right), \\
& v_{-}(\gamma)=\frac{\left(d-\lambda_{-}(\gamma), c\right)}{\sqrt{\left(d-\lambda_{-}(\gamma)\right)^{2}+c^{2}}}=\frac{(-d, c)}{\sqrt{c^{2}+d^{2}}}+O\left(\frac{1}{\|\gamma\|^{2}}\right) .
\end{aligned}
$$

Note that for $\gamma$ large,

$$
\begin{equation*}
\left|\left\langle v_{+}(\gamma), v_{-}(\gamma)^{\perp}\right\rangle\right|=\frac{b c+d^{2}}{\sqrt{b^{2}+d^{2}} \sqrt{c^{2}+d^{2}}}+O\left(\frac{1}{\|\gamma\|^{2}}\right) \geq \frac{1}{2} \tag{2.9}
\end{equation*}
$$

meaning that the angle between expanding and contracting vectors does not degenerate.

By (2.8), it is enough to show that the traces behave essentially multiplicatively. We compute

$$
\begin{aligned}
\left|\operatorname{tr}\left(\gamma \gamma^{\prime}\right)-\operatorname{tr}(\gamma) \operatorname{tr}\left(\gamma^{\prime}\right)\right| & =\left|\left(a a^{\prime}+b c^{\prime}+c b^{\prime}+d d^{\prime}\right)-(a+d)\left(a^{\prime}+d^{\prime}\right)\right| \\
& \leq \frac{d}{d^{\prime}}\left|\frac{b c^{\prime} d^{\prime}}{d}-a^{\prime} d^{\prime}\right|+\frac{d^{\prime}}{d}\left|\frac{c b^{\prime} d}{d^{\prime}}-a d\right| \\
& \leq \frac{d}{d^{\prime}}\left(1+c^{\prime}\left|\frac{b d^{\prime}}{d}-b^{\prime}\right|\right)+\frac{d^{\prime}}{d}\left(1+c\left|\frac{b^{\prime} d}{d^{\prime}}-b\right|\right) \\
& =\frac{d}{d^{\prime}}+\frac{d^{\prime}}{d}+\left(c d^{\prime}+c^{\prime} d\right)\left|\frac{b}{d}-\frac{b^{\prime}}{d^{\prime}}\right| .
\end{aligned}
$$

We clearly have

$$
\left|\frac{b}{d}-\frac{b^{\prime}}{d^{\prime}}\right| \ll\left|v_{+}(\gamma)-v_{+}\left(\gamma^{\prime}\right)\right|+\frac{1}{\|\gamma\|^{2}}+\frac{1}{\left\|\gamma^{\prime}\right\|^{2}},
$$

and hence

$$
\left|\operatorname{tr}\left(\gamma \gamma^{\prime}\right)-\operatorname{tr}(\gamma) \operatorname{tr}\left(\gamma^{\prime}\right)\right| \ll d d^{\prime}\left(\left|v_{+}(\gamma)-v_{+}\left(\gamma^{\prime}\right)\right|+\frac{1}{\|\gamma\|^{2}}+\frac{1}{\left\|\gamma^{\prime}\right\|^{2}}\right)
$$

From this and (2.8), (2.6) follows easily. One proves (2.7) in a similar fashion.

## 3. Construction of $\Omega_{N}$

3.1. The leading term $\aleph$. In this subsection, we posit the existence and all necessary properties of the leading set $\aleph$ used in our construction of the main ensemble $\Omega_{N}$. The proof of its existence is arguably the most technical part of the whole paper, so in the interest of exposition, we postpone it to Section 8.

Once and for all, we fix a density point $x \in \mathfrak{C}$ and let

$$
\begin{equation*}
\mathfrak{v}=\mathfrak{v}_{x}:=\frac{(x, 1)}{\sqrt{1+x^{2}}} \tag{3.1}
\end{equation*}
$$

be the corresponding unit vector. We will henceforth be largely concerned with elements of $\Gamma$ whose expanding eigenvectors point in this direction.

For ease of exposition, we assume henceforth that for all $q \geq 1$, the reduction of $\Gamma$ is full,

$$
\begin{equation*}
\Gamma(\bmod q) \cong \mathrm{SL}_{2}(q) \tag{3.2}
\end{equation*}
$$

which is anyway the case for any alphabet $\mathcal{A}$ containing 1 and 2 ; see Remark 1.33. Minor modifications are needed in the general case.

For $N$ large and $\delta$ exceeding $\delta_{0}$ in (1.21), let

$$
\begin{equation*}
\mathfrak{b}:=\frac{1}{1000}\left(\delta-\delta_{0}\right)>0, \tag{3.3}
\end{equation*}
$$

and let $\alpha_{0}>0$ be a parameter to be chosen later in (8.20). Then we set

$$
\begin{equation*}
B:=N^{\mathfrak{b}} \tag{3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{Q}:=e^{\alpha_{0} \sqrt{\log N}} \tag{3.5}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathcal{U} \subset\left[\frac{1}{100} B, \frac{99}{100} B\right] \tag{3.6}
\end{equation*}
$$

be an arithmetic progression of real numbers starting with $u_{0}=\frac{1}{100} B$ having common difference

$$
\begin{equation*}
\left|u-u^{\prime}\right|=2 B / \mathcal{Q}^{5}, \tag{3.7}
\end{equation*}
$$

for $u, u^{\prime}$ consecutive terms in $\mathcal{U}$, and ending with $u>\left(\frac{99}{100}-\frac{2}{\mathcal{Q}^{5}}\right) B$. Then the cardinality of $\mathcal{U}$ is

$$
\begin{equation*}
|\mathcal{U}| \asymp \mathcal{Q}^{5} . \tag{3.8}
\end{equation*}
$$

Proposition 3.9. For each $u \in \mathcal{U}$, there are nonempty sets $\aleph_{u} \subset \Gamma$, all of the same cardinality

$$
\begin{equation*}
\left|\aleph_{u}\right|=\left|\aleph_{u^{\prime}}\right|, \tag{3.10}
\end{equation*}
$$

so that the following holds. For every $\mathfrak{a} \in \aleph_{u}$, its expanding eigenvector is restricted by

$$
\begin{equation*}
\left|v_{+}(\mathfrak{a})-\mathfrak{v}\right|<\mathcal{Q}^{-5} \tag{3.11}
\end{equation*}
$$

and its expanding eigenvalue $\lambda(\mathfrak{a})$ is restricted by

$$
\begin{equation*}
|\lambda(\mathfrak{a})-u|<\frac{B}{\mathcal{Q}^{5}} . \tag{3.12}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\frac{1}{200} B<\lambda(\mathfrak{a})<B \tag{3.13}
\end{equation*}
$$

for $N$ large. Moreover, for any $q<\mathcal{Q}$, any $\omega \in \operatorname{SL}_{2}(q)$, and any $u \in \mathcal{U}$, we have

$$
\begin{equation*}
\#\left\{\mathfrak{a} \in \aleph_{u}: \mathfrak{a} \equiv \omega(\bmod q)\right\}=\frac{\left|\aleph_{u}\right|}{\left|\operatorname{SL}_{2}(q)\right|}\left(1+O\left(\mathcal{Q}^{-4}\right)\right) \tag{3.14}
\end{equation*}
$$

where the implied constant does not depend on $q, \omega$, or $u$.
With the sets $\aleph_{u}$ as above, we define the main leading set $\aleph$ to be the union of the sets $\aleph_{u}$,

$$
\begin{equation*}
\aleph:=\bigsqcup_{u \in \mathcal{U}} \aleph_{u} . \tag{3.15}
\end{equation*}
$$

Note that the sets $\aleph_{u}$ are disjoint by (3.12) and (3.7). We repeat that the proof of Proposition 3.9 will be postponed to Section 8.
3.2. Sector counting. In this section we give the following slight refinement of Hensley's estimate (1.14), which follows directly from Lalley's methods [Lal89] (see also Theorem 8.1 and Remark 8.6).

Proposition 3.16. There is a constant $\mathfrak{c}=\mathfrak{c}(\mathcal{A})>0$ so that as long as $H<e^{\mathfrak{c} \sqrt{\log T}}$, we have

$$
\begin{equation*}
\#\left\{\gamma \in \Gamma:\|\gamma\|<T \text { and }\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H}\right\} \gg \frac{T^{2 \delta}}{H} \tag{3.17}
\end{equation*}
$$

as $T \rightarrow \infty$.
Sketch of proof. Lalley [Lal89, Th. 9] proves the asymptotic formula

$$
\begin{equation*}
\text { Left-hand side of }(3.17) \sim C \cdot T^{2 \delta} \mu(\mathcal{I}) \tag{3.18}
\end{equation*}
$$

under the assumption that $\Gamma$ is a nonelementary convex-cocompact subgroup of $\mathrm{SL}_{2}(\mathbb{R})$. Here $\mathcal{I}$ is the interval of length $1 / H$ about $\mathfrak{v}$, and $\mu$ is the $\delta$-dimensional Hausdorff measure supported on the limit set $\mathfrak{C}$, lifted (by abuse of notation) to $\mathbb{P}^{1}$ via

$$
d \mu(x, y)=d \mu(x / y),
$$

$y \neq 0$. After setting up the symbolic dynamics, the requirement that $\Gamma$ not contain parabolic elements is needed in the renewal method to make the distortion function eventually positive; see [Lal89, pp. 33, 41]. Our semigroup $\Gamma$ has no parabolic elements (see Remark 2.4), so the only difference here between a group and (free) semigroup is that for the latter, the transition matrix (see [Lal89, pp. 5, 32]) is trivial; that is, all sequences are allowed in the symbolic dynamics. The rate in (3.18) can be determined directly from Lalley's method (see [BGS11, §12]), with the error crudely estimated as

$$
\begin{equation*}
\ll T^{2 \delta} e^{-c \sqrt{\log T}} . \tag{3.19}
\end{equation*}
$$

Since $\mathfrak{v}$ in (3.1) corresponds to a density point in $\mathfrak{C}$, we have, again crudely, that

$$
\mu(\mathcal{I}) \gg \varepsilon H^{-\delta-\varepsilon} \gg H^{-1},
$$

since $\delta<1$. A sufficient condition for the main term, being bounded below by $C T^{2 \delta} / H$, to dominate the error in (3.19), is that $H<e^{\mathfrak{c} \sqrt{\log T}}$ with $\mathfrak{c}<c$.

Remark 3.20. The methods of Dolgopyat [Do198] and Naud [Nau05] might be used to prove (3.17) with $H$ as large as $T^{\varepsilon}$, but this is not needed in our applications.

With this crude estimate in hand, we proceed in the next subsection to detail our construction of the special sets $\Xi$ alluded to in Section 1.3.
3.3. The set $\Xi(M, H ; L, k)$.

Proposition 3.21. Given $M \gg 1$ and $H<e^{\mathfrak{c} \sqrt{\log M}}$, there exists some $L$ in the range

$$
\begin{equation*}
\frac{1}{4} M \leq L \leq 4 M \tag{3.22}
\end{equation*}
$$

an integer $k \asymp \log M$, and a set $\Xi=\Xi(M, H ; L, k) \subset \Gamma$ having the following properties. For all $\gamma \in \Xi$, the expanding eigenvalues are controlled to within $1 / \log L$ :

$$
\begin{equation*}
L\left(1-\frac{1}{\log L}\right)<\lambda(\gamma)<L, \tag{3.23}
\end{equation*}
$$

the expanding eigenvectors are controlled to within $1 / H$ :

$$
\begin{equation*}
\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H}, \tag{3.24}
\end{equation*}
$$

and the wordlength metric $\ell$ (in the generators (1.32) of $\Gamma$ ) is controlled exactly:

$$
\begin{equation*}
\ell(\gamma)=k . \tag{3.25}
\end{equation*}
$$

Moreover, the cardinality of $\Xi$ is controlled by

$$
\begin{equation*}
L^{2 \delta} \gg \# \Xi \gg \frac{L^{2 \delta}}{H(\log L)^{2}} . \tag{3.26}
\end{equation*}
$$

Recall again the the implied constants depend at most on $\mathcal{A}$, which is thought of as fixed throughout.

Proof. We proceed by the following algorithm.
(1) Let $S_{1} \subset \Gamma$ be the set of $\gamma \in \Gamma$ of norm controlled by $\|\gamma\| \asymp M$ and for which the expanding vector $v_{+}(\gamma)$ is within $\frac{1}{H}$ of the fixed vector $\mathfrak{v}$ :

$$
S_{1}:=\left\{\gamma \in \Gamma: \frac{M}{2}<\|\gamma\|<M,\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H}\right\} .
$$

By (3.17), we have that

$$
\# S_{1} \gg \frac{M^{2 \delta}}{H}
$$

(2) By (2.2) and (2.8), expanding eigenvalues $\lambda(\gamma)$ of $\gamma \in S_{1}$ satisfy

$$
\frac{1}{4} M \leq \lambda(\gamma) \leq 4 M
$$

Hence we can find (by pigeonhole) an $L$ in this range so that

$$
\#\left\{\gamma \in S_{1}: L\left(1-\frac{1}{\log L}\right)<\lambda(\gamma)<L\right\} \gg \frac{L^{2 \delta}}{H \log L} .
$$

Call the above set $S_{2}$; its expanding eigenvalues are all nearly of the same size.
(3) Lastly, note that the wordlength metric $\ell$ is commensurable with the archimedean one,

$$
\ell(\gamma) \asymp \log \|\gamma\|,
$$

with implied constant depending on $\mathcal{A}$. So (again by pigeonhole) we can find some $k$ such that

$$
\begin{equation*}
\#\left\{\gamma \in S_{2}: \ell(\gamma)=k\right\} \gg \frac{L^{2 \delta}}{H(\log L)^{2}} \tag{3.27}
\end{equation*}
$$

Call this set $S_{3}$; then the elements of $S_{3}$ all have the same wordlength, in addition to the previous qualities.
We rename this last set $S_{3}$ to $\Xi=\Xi(M, H ; L, k)$.
3.4. Decomposing $N$ and the ensemble $\Omega_{N}$. We return to our main parameter $N$ and decompose it dyadically as follows. Recall that we have already presupposed the construction of a set $\aleph$ in (3.15), all of whose expanding eigenvectors are within $\mathcal{Q}^{-5}$ of $\mathfrak{v}$, and with eigenvalues of size $B$; see (3.13). Recall from (3.4) that $B=N^{\mathfrak{b}}$ and that $\mathcal{Q}$ is given by (3.5).

Setup: We start by taking

$$
\begin{equation*}
M=\sqrt{N} / B=N^{1 / 2-\mathfrak{b}}, \quad H=\mathcal{Q}^{5} . \tag{3.28}
\end{equation*}
$$

The exponent $\alpha_{0}$ in the definition (3.5) of $\mathcal{Q}$ is chosen in (8.20) to be sufficiently small that $H<e^{\mathfrak{c} \sqrt{\log M}}$. Run the algorithm of the previous subsection to generate the set $\Xi(M, H ; L, k)$. By (3.22), the returned parameter $L$ satisfies

$$
L=\alpha_{1} M=\alpha_{1} N^{1 / 2-\mathfrak{b}}
$$

with

$$
\alpha_{1} \in(1 / 4,4) .
$$

Write

$$
\tilde{N}_{1}:=L=\alpha_{1} N^{1 / 2-\mathfrak{b}}, \quad N_{1}:=\alpha_{1} N^{1 / 2}=B \cdot \tilde{N}_{1},
$$

and rename the returned set to $\tilde{\Xi}_{1}=\Xi(M, H ; L, k)$, also setting

$$
\Xi_{1}:=\aleph \cdot \tilde{\Xi}_{1} .
$$

Remark 3.29. Despite the wordlength in $\aleph$ being unrestricted, the wordlength in $\tilde{\Xi}_{1}$ is fixed in (3.25). So the representation of an element in $\Xi_{1}$ as a product of ones in $\aleph$ and $\tilde{\Xi}_{1}$ is still unique.

We have crudely that

$$
\begin{equation*}
\left|\Xi_{1}\right| \geq\left|\tilde{\Xi}_{1}\right| \gg_{\varepsilon} \tilde{N}_{1}^{2 \delta-\varepsilon} \gg N^{\delta-2 \delta \mathbf{b}-\varepsilon} . \tag{3.30}
\end{equation*}
$$

(The cardinality of $\aleph$ is quite deficient relative to its norm, so we lose little from estimating trivially $|\aleph| \geq 1$.)

Step 1: Next we set

$$
M=\frac{N_{1}^{1 / 2}}{\alpha_{1}}=\frac{N^{1 / 4}}{\alpha_{1}^{1 / 2}}, \quad H=\log M
$$

and we generate another set $\Xi(M, H ; L, k)$. Define

$$
N_{2}:=L=\alpha_{2} M=\frac{\alpha_{2} N^{1 / 4}}{\alpha_{1}^{1 / 2}},
$$

with $\alpha_{2} \in(1 / 4,4)$, and rename the returned set to $\Xi_{2}$. We have

$$
\left|\Xi_{2}\right| \gg \frac{N_{2}^{2 \delta}}{\left(\log N_{2}\right)^{3}}
$$

Iterate: Start with $j=3$ and iterate up to $j=J-1$, where

$$
\begin{equation*}
2^{J-1}=c \log N \tag{3.31}
\end{equation*}
$$

Here the constant $c>0$ is absolute (independent of $N$ ), determined by (3.45). For each such $j$, set

$$
\begin{equation*}
M:=\frac{\left(N_{j-1}\right)^{1 / 2}}{\alpha_{j-1}}=\frac{N^{1 / 2^{j}}}{\alpha_{j-1}^{1 / 2} \alpha_{j-2}^{1 / 4} \cdots\left(\alpha_{1}\right)^{1 / 2^{(j-1)}}}, \quad H=\log M, \tag{3.32}
\end{equation*}
$$

and use Proposition 3.21 to generate the set $\Xi(M, H ; L, k)$. Define

$$
\begin{equation*}
N_{j}:=L=\alpha_{j} M=\frac{\alpha_{j} N^{1 / 2^{j}}}{\alpha_{j-1}^{1 / 2} \alpha_{j-2}^{1 / 4} \cdots\left(\alpha_{1}\right)^{1 / 2^{(j-1)}}}, \tag{3.33}
\end{equation*}
$$

with $\alpha_{j} \in(1 / 4,4)$, and call the returned set $\Xi_{j}$. Note that

$$
\begin{equation*}
\left|\Xi_{j}\right| \gg \frac{N_{j}^{2 \delta}}{\left(\log N_{j}\right)^{3}} \tag{3.34}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{1}{16} N^{1 / 2^{j}}<N_{j}<16 N^{1 / 2^{j}} \tag{3.35}
\end{equation*}
$$

End: For the last step, $j=J$, we set

$$
M=\frac{N_{J-1}}{\left(\alpha_{J-1}\right)^{2}}=\frac{N^{1 / 2^{(J-1)}}}{\alpha_{J-1}\left(\alpha_{J-2}\right)^{1 / 2} \cdots \alpha_{1}^{1 / 2^{(J-2)}}}, \quad H=\log M,
$$

and we generate one last set $\Xi_{J}:=\Xi(M, H ; L, k)$. Define

$$
N_{J}:=L=\frac{\alpha_{J} N^{1 / 2^{(J-1)}}}{\alpha_{J-1} \cdots \alpha_{1}^{1 / 2^{(J-2)}}} \asymp N^{1 / 2^{(J-1)}}=e^{1 / c} \ll 1,
$$

where we used (3.31). Since $\frac{1}{4}<N_{J} / M=\alpha_{J}<4$, we have

$$
\begin{equation*}
\frac{1}{4}<\frac{N_{1} N_{2} \cdots N_{J}}{N}=\frac{B \tilde{N}_{1} N_{2} \cdots N_{J}}{N}<4 \tag{3.36}
\end{equation*}
$$

We now define the main ensemble $\Omega_{N}$ by concatenating the sets $\Xi_{j}$ developed above.

$$
\begin{equation*}
\Omega_{N}:=\Xi_{1} \cdot \Xi_{2} \cdots \Xi_{J-1} \cdot \Xi_{J}=\aleph \cdot \tilde{\Xi}_{1} \cdot \Xi_{2} \cdots \Xi_{J-1} \cdot \Xi_{J} . \tag{3.37}
\end{equation*}
$$

3.5. Properties of $\Omega_{N}$. For $\gamma \in \Omega_{N}$, write

$$
\gamma=\mathfrak{a} \cdot \tilde{\xi}_{1} \xi_{2} \cdots \xi_{J}
$$

according to the decomposition (3.37). Note that by the fixed wordlength restriction (3.25), this decomposition is unique; see Remark 3.29. Recall that the expanding vectors $v_{+}$all point nearly in the direction of $\mathfrak{v}$ in (3.1).

Lemma 3.38. For any $2 \leq j_{1} \leq j_{2} \leq J$, and $\xi_{j_{1}} \in \Xi_{j_{1}}, \cdots \xi_{j_{2}} \in \Xi_{j_{2}}$, and any $\mathfrak{a} \in \aleph, \tilde{\xi}_{1} \in \tilde{\Xi}_{1}$, we have the following control on expanding eigen-vectors and -values of large products:

$$
\begin{gather*}
\left|v_{+}\left(\tilde{\xi}_{1} \cdot \xi_{2} \cdots \xi_{J}\right)-\mathfrak{v}\right| \ll \mathcal{Q}^{-5},  \tag{3.39}\\
\frac{1}{2}<\frac{\lambda\left(\xi_{j_{1}} \xi_{j_{1}+1} \cdots \xi_{j_{2}-1} \xi_{j_{2}}\right)}{N_{j_{1}} N_{j_{1}+1}^{\cdots N_{j_{2}-1} N_{j_{2}}}<2,}  \tag{3.40}\\
\frac{1}{2}<\frac{\lambda\left(\tilde{\xi}_{1} \xi_{2} \cdots \xi_{j_{2}-1} \xi_{j_{2}}\right)}{\tilde{N}_{1} N_{2} \cdots N_{j_{2}-1} N_{j_{2}}}<2, \tag{3.41}
\end{gather*}
$$

and

$$
\begin{equation*}
\frac{1}{2}<\frac{\lambda\left(\mathfrak{a} \tilde{\xi}_{1} \xi_{2} \cdots \xi_{j_{2}-1} \xi_{j_{2}}\right)}{\lambda(\mathfrak{a}) \tilde{N}_{1} \cdot N_{2} \cdots N_{j_{2}-1} N_{j_{2}}}<2 \tag{3.42}
\end{equation*}
$$

Proof. From (2.7), (3.24), and the choice of $H$ in (3.28), we have that

$$
\begin{aligned}
\left|v_{+}\left(\tilde{\xi}_{1} \cdot \xi_{2} \cdots \xi_{J}\right)-\mathfrak{v}\right| & \leq\left|v_{+}\left(\tilde{\xi}_{1} \cdot \xi_{2} \cdots \xi_{J}\right)-v_{+}\left(\tilde{\xi}_{1}\right)\right|+\left|v_{+}\left(\tilde{\xi}_{1}\right)-\mathfrak{v}\right| \\
& \ll \frac{1}{\left\|\tilde{\xi}_{1}\right\|^{2}}+\frac{1}{\mathcal{Q}^{5}},
\end{aligned}
$$

whence (3.39) follows from (3.5).
Similarly, we have for $j \in\left[j_{1}, j_{2}\right] \subset[2, J]$ that

$$
\begin{equation*}
\left|v_{+}\left(\xi_{j} \xi_{j+1} \cdots \xi_{j_{2}}\right)-\mathfrak{v}\right| \ll \frac{1}{\log N_{j}}, \tag{3.43}
\end{equation*}
$$

where we used the choice of $H$ in (3.32).

We now prove by downward induction on $j_{1}$ that

$$
\begin{align*}
\lambda\left(\xi_{j_{1}} \xi_{j_{1}+1} \cdots \xi_{j_{2}}\right)= & N_{j_{1}} N_{j_{1}+1} \cdots N_{j_{2}}  \tag{3.44}\\
& \times\left[1+O\left(\frac{1}{\log N_{j_{1}}}+\frac{1}{\log N_{j_{1}+1}}+\cdots+\frac{1}{\log N_{j_{2}}}\right)\right] .
\end{align*}
$$

If $j_{1}=j_{2}$, then (3.44) follows immediately from (3.23) and (3.33). If $j_{1}=j_{2}-1$, then from (2.6), (3.23), (3.24), and (3.33), we have

$$
\begin{aligned}
\lambda\left(\xi_{j_{2}-1} \xi_{j_{2}}\right)= & \lambda\left(\xi_{j_{2}-1}\right) \lambda\left(\xi_{j_{2}}\right) \\
& \times\left[1+O\left(\left|v_{+}\left(\xi_{j_{2}-1}\right)-v_{+}\left(\xi_{j_{2}}\right)\right|+\frac{1}{\left\|\xi_{j_{2}-1}\right\|^{2}}+\frac{1}{\left\|\xi_{j_{2}}\right\|^{2}}\right)\right] \\
= & N_{j_{2}-1} N_{j_{2}}\left[1+O\left(\frac{1}{\log N_{j_{2}-1}}+\frac{1}{\log N_{j_{2}}}\right)\right],
\end{aligned}
$$

as desired.
In general, we have by (3.43) that

$$
\begin{aligned}
& \lambda\left(\xi_{j_{1}} \xi_{j_{1}+1} \cdots \xi_{j_{2}}\right)=\lambda\left(\xi_{j_{1}}\right) \lambda\left(\xi_{j_{1}+1} \cdots \xi_{j_{2}}\right) \\
& \quad \times\left[1+O\left(\left|v_{+}\left(\xi_{j_{1}}\right)-v_{+}\left(\xi_{j_{1}+1} \cdots \xi_{j_{2}}\right)\right|+\frac{1}{\left\|\xi_{j_{1}}\right\|^{2}}+\frac{1}{\lambda\left(\xi_{j_{1}+1} \cdots \xi_{j_{2}}\right)^{2}}\right)\right] \\
& \quad= \\
& \quad N_{j_{1}} \lambda\left(\xi_{j_{1}+1} \cdots \xi_{j_{2}}\right)\left[1+O\left(\frac{1}{\log N_{j_{1}}}+\frac{1}{\log N_{j_{1}+1}}\right)\right],
\end{aligned}
$$

from which (3.44) follows by induction.
The rate in (3.44) may be replaced crudely by

$$
\begin{equation*}
\left[1+O\left(\frac{2^{J}}{\log N}\right)\right], \tag{3.45}
\end{equation*}
$$

whence (3.40) follows on taking the constant $c$ in (3.31) sufficiently small (independent of $N$ ). Estimates (3.41) and (3.42) are proved in the same way.

As a consequence of (3.42), (3.13), and (3.36), we have that for all $\gamma \in \Omega_{N}$,

$$
\begin{equation*}
\|\gamma\| \leq 2 \lambda(\gamma) \leq 16 N \tag{3.46}
\end{equation*}
$$

so indeed the norms are all controlled.
Moreover, the size of $\Omega_{N}$ is not too much smaller than (1.36). Indeed, we have from (3.30), (3.34), and (3.31) that

$$
\begin{align*}
\# \Omega_{N} & =\# \Xi_{1} \cdot \# \Xi_{2} \cdots \# \Xi_{J}  \tag{3.47}\\
& \gg \varepsilon \tilde{N}_{1}^{2 \delta-\varepsilon} \frac{\left(N_{2}\right)^{2 \delta}}{\left(\log N_{2}\right)^{3}} \cdots \frac{\left(N_{J}\right)^{2 \delta}}{\left(\log N_{J}\right)^{3}} \\
& \gg N^{2 \delta-2 \delta \mathfrak{b}-\varepsilon} .
\end{align*}
$$

It also follows that for any $j \geq 2$,

$$
\begin{equation*}
\# \Xi_{j} \cdot \# \Xi_{j+1} \cdots \# \Xi_{J} \gg\left(N_{j} N_{j+1} \cdots N_{J}\right)^{2 \delta} \cdot e^{-c(J-j) \log \log N_{j}} \tag{3.48}
\end{equation*}
$$

for an absolute constant $c>0$.
With the set $\Omega_{N}$ constructed, we define our exponential sum $S_{N}$ as in (1.45) and proceed with the circle method.

## 4. Major arcs analysis

In this section we estimate the major arcs contribution. First we use the set $\aleph$ described in Section 3.1 to prove that in the major arcs, our exponential sum $S_{N}$ in (1.45) splits as a product of modular and archimedean components, as in (1.38). Then we prove that the major arcs contribution is of the correct order of magnitude.
4.1. Splitting into modular and Archimedean aomponents. Let $\mathcal{Q}$ be as in (3.5) and $B$ as in (3.4). Recall from (1.37) that the major arcs of level $\mathcal{Q}$ are given by

$$
\mathfrak{M}_{\mathcal{Q}}=\bigsqcup_{q<\mathcal{Q}} \bigsqcup_{(a, q)=1}\left[\frac{a}{q}-\frac{\mathcal{Q}}{N}, \frac{a}{q}+\frac{\mathcal{Q}}{N}\right] .
$$

Let $\nu_{q}: \mathbb{Z} / q \mathbb{Z} \rightarrow \mathbb{C}$ record the $\bmod q$ distribution of $\mathfrak{D}$. That is, for $a \in \mathbb{Z} / q \mathbb{Z}$, set

$$
\begin{equation*}
\nu_{q}(a):=\frac{1}{\left|\mathrm{SL}_{2}(q)\right|} \sum_{\omega \in \mathrm{SL}_{2}(q)} e\left(\frac{a}{q}\left\langle\omega e_{2}, e_{2}\right\rangle\right) . \tag{4.1}
\end{equation*}
$$

Theorem 4.2. There exists a function $\varpi_{N}: \mathbb{R} / \mathbb{Z} \rightarrow \mathbb{C}$, given explicitly in (4.19), satisfying the following three conditions:
(1) The Fourier transform

$$
\widehat{\varpi}_{N}: \mathbb{Z} \rightarrow \mathbb{C}: n \mapsto \int_{0}^{1} \varpi_{N}(\theta) e(-n \theta) d \theta
$$

is real-valued and nonnegative, with

$$
\varpi_{N}(0)=\sum_{n} \widehat{\varpi}_{N}(n) \ll\left|\Omega_{N}\right| .
$$

(2) For $\frac{1}{25} N<n<\frac{1}{5} N$, we have

$$
\begin{equation*}
\widehat{\varpi}_{N}(n) \gg \frac{\left|\Omega_{N}\right|}{N} . \tag{4.4}
\end{equation*}
$$

(3) Moreover, we have on the major arcs $\theta=\frac{a}{q}+\beta \in \mathfrak{M}_{\mathcal{Q}}$ that

$$
\begin{equation*}
S_{N}\left(\frac{a}{q}+\beta\right)=\nu_{q}(a) \varpi_{N}(\beta)\left(1+O\left(\mathcal{Q}^{-4}\right)\right) \tag{4.5}
\end{equation*}
$$

Proof. We use the decomposition (3.37) in the form

$$
\begin{equation*}
\Omega_{N}=\aleph \cdot \Omega^{\prime} \tag{4.6}
\end{equation*}
$$

with

$$
\Omega^{\prime}=\tilde{\Xi}_{1} \Xi_{2} \cdots \Xi_{J}
$$

so that

$$
S_{N}(\theta)=\sum_{\mathfrak{a} \in \mathcal{K}} \sum_{\gamma \in \Omega^{\prime}} e\left(\theta\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle\right) .
$$

For $\mathfrak{a} \in \aleph$, recall from (3.13) that we have $\lambda(\mathfrak{a}) \asymp B$ and from (3.11) and (3.39) that

$$
\begin{equation*}
\left|v_{+}(\mathfrak{a})-\mathfrak{v}\right|<\mathcal{Q}^{-5}, \quad\left|v_{+}(\gamma)-\mathfrak{v}\right| \ll \mathcal{Q}^{-5} . \tag{4.7}
\end{equation*}
$$

To use the key property (3.12) in the construction of $\aleph$, we need to convert the expression $\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle$ in $S_{N}$ into one involving $\lambda(\mathfrak{a})$.

We will make regular use of the following elementary formula: For any two linearly independent vectors $v_{+}, v_{-} \in \mathbb{R}^{2}$, we can write any $w \in \mathbb{R}^{2}$ as

$$
\begin{equation*}
w=\frac{\left\langle w, v_{-}^{\perp}\right\rangle}{\left\langle v_{+}, v_{-}^{\perp}\right\rangle} v_{+}+\frac{\left\langle w, v_{+}^{\perp}\right\rangle}{\left\langle v_{-}, v_{+}^{\perp}\right\rangle} v_{-} . \tag{4.8}
\end{equation*}
$$

Here $(x, y)^{\perp}=(-y, x)$. Recalling (2.9), it easily follows that for a unit vector $w$ and any large $\xi \in \Gamma$,

$$
\begin{equation*}
\xi w=\lambda(\xi) \frac{\left\langle w, v_{-}^{\perp}(\xi)\right\rangle}{\left\langle v_{+}(\xi), v_{-}^{\perp}(\xi)\right\rangle} v_{+}(\xi)\left(1+O\left(\frac{1}{\|\xi\|^{2}}\right)\right) \tag{4.9}
\end{equation*}
$$

whence

$$
\begin{equation*}
\left\langle\xi e_{2}, e_{2}\right\rangle=\lambda(\xi) \frac{\left\langle e_{2}, v_{-}^{\perp}(\xi)\right\rangle}{\left\langle v_{+}(\xi), v_{-}^{\perp}(\xi)\right\rangle}\left\langle v_{+}(\xi), e_{2}\right\rangle\left(1+O\left(\frac{1}{\|\xi\|^{2}}\right)\right) . \tag{4.10}
\end{equation*}
$$

Applied to our present situation, we have by (4.7) that

$$
\begin{equation*}
\left\langle\gamma e_{2}, e_{2}\right\rangle=\lambda(\gamma) \frac{\left\langle e_{2}, v_{-}^{\perp}(\gamma)\right\rangle}{\left\langle\mathfrak{v}, v_{-}^{\perp}(\gamma)\right\rangle}\left\langle\mathfrak{v}, e_{2}\right\rangle\left(1+O\left(\frac{1}{\mathcal{Q}^{5}}\right)\right) \tag{4.11}
\end{equation*}
$$

and

$$
\begin{align*}
\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle & =\lambda(\mathfrak{a} \gamma) \frac{\left\langle e_{2}, v_{-}^{\perp}(\mathfrak{a} \gamma)\right\rangle}{\left\langle v_{+}(\mathfrak{a} \gamma), v_{-}^{\perp}(\mathfrak{a} \gamma)\right\rangle}\left\langle v_{+}(\mathfrak{a} \gamma), e_{2}\right\rangle\left(1+O\left(\frac{1}{N^{2}}\right)\right)  \tag{4.12}\\
& =\lambda(\mathfrak{a}) \lambda(\gamma) \frac{\left\langle e_{2}, v_{-}^{\perp}(\gamma)\right\rangle}{\left\langle\mathfrak{v}, v_{-}^{\perp}(\gamma)\right\rangle}\left\langle\mathfrak{v}, e_{2}\right\rangle\left(1+O\left(\frac{1}{\mathcal{Q}^{5}}\right)\right),
\end{align*}
$$

where we also used (2.6) and (2.7).
Comparing (4.12) and (4.11), we have that

$$
\begin{equation*}
\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle=\lambda(\mathfrak{a})\left\langle\gamma e_{2}, e_{2}\right\rangle+O\left(N / \mathcal{Q}^{5}\right) . \tag{4.13}
\end{equation*}
$$

For $\theta=\frac{a}{q}+\beta \in \mathfrak{M}_{\mathcal{Q}}$ with $|\beta|<\mathcal{Q} / N$, we insert (4.13) into $S_{N}$, giving

$$
\begin{align*}
S_{N}\left(\frac{a}{q}+\beta\right)= & \sum_{\mathfrak{a} \in \aleph} \sum_{\gamma \in \Omega^{\prime}} e\left(\frac{a}{q}\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle\right) e\left(\beta\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle\right)  \tag{4.14}\\
= & \sum_{\mathfrak{a} \in \aleph} \sum_{\gamma \in \Omega^{\prime}} e\left(\frac{a}{q}\left\langle\mathfrak{a} \gamma e_{2}, e_{2}\right\rangle\right) e\left(\beta \lambda(\mathfrak{a})\left\langle\gamma e_{2}, e_{2}\right\rangle\right)+O\left(\mathcal{Q}^{-4}|\Omega|\right) \\
= & \sum_{\gamma \in \Omega^{\prime}} \sum_{\omega \in \operatorname{SL}_{2}(q)} e\left(\frac{a}{q}\left\langle\omega \gamma e_{2}, e_{2}\right\rangle\right) \sum_{\substack{\mathfrak{a} \in \mathfrak{N} \\
\mathfrak{a} \equiv \omega(\bmod q)}} e\left(\beta \lambda(\mathfrak{a})\left\langle\gamma e_{2}, e_{2}\right\rangle\right) \\
& +O\left(\mathcal{Q}^{-4}\left|\Omega_{N}\right|\right),
\end{align*}
$$

where we decomposed the $\mathfrak{a}$ sum into residue classes $\omega$ in $\mathrm{SL}_{2}(q)$.
Next using (3.15), write the innermost sum above as

$$
\begin{align*}
& \sum_{\substack{\mathfrak{a} \in \mathcal{N} \\
\mathfrak{a} \equiv \omega(\bmod q)}} e\left(\lambda(\mathfrak{a}) \beta\left\langle\gamma e_{2}, e_{2}\right\rangle\right)  \tag{4.15}\\
& \quad=\sum_{u \in \mathcal{U}} \sum_{\substack{\mathfrak{a} \in \in \mathfrak{N}, \mathfrak{a}=\omega(\bmod q)}} e\left(\lambda(\mathfrak{a}) \beta\left\langle\gamma e_{2}, e_{2}\right\rangle\right) \\
& \quad=\sum_{u \in \mathcal{U}} e\left(\beta u\left\langle\gamma e_{2}, e_{2}\right\rangle\right)\left(\sum_{\substack{\mathfrak{a} \in \mathbb{N} u \\
\mathfrak{a} \equiv \omega(\bmod q)}} 1\right)\left(1+O\left(\mathcal{Q}^{-4}\right)\right),
\end{align*}
$$

where we applied (3.12).
By (3.14) and (3.10) (that the cardinality of $\aleph_{u}$ is the same for all $u$ ), the innermost sum is

$$
\begin{equation*}
\sum_{\substack{a \in \aleph_{u} \\ \mathfrak{a}=\omega(\bmod q)}} 1=\frac{|\aleph|}{|\mathcal{U}| \cdot\left|\mathrm{SL}_{2}(q)\right|}\left(1+O\left(\mathcal{Q}^{-4}\right)\right), \tag{4.16}
\end{equation*}
$$

where the implied constant does not depend on $u, \omega$, or $q$.
Returning to (4.14), inputting (4.16) and (4.15) gives

$$
\begin{align*}
S_{N}\left(\frac{a}{q}+\beta\right)= & \frac{1}{\left|\mathrm{SL}_{2}(q)\right|} \sum_{\omega \in \mathrm{SL}_{2}(q)} e\left(\frac{a}{q}\left\langle\omega e_{2}, e_{2}\right\rangle\right) \frac{|\mathcal{\aleph}|}{|\mathcal{U}|} \sum_{\gamma \in \Omega^{\prime}} \sum_{u \in \mathcal{U}} e\left(\beta u\left\langle\gamma e_{2}, e_{2}\right\rangle\right)  \tag{4.17}\\
& \times\left(1+O\left(\mathcal{Q}^{-4}\right)\right)
\end{align*}
$$

where we used the fact that the $\omega$ sum runs over all of $\mathrm{SL}_{2}(q)$, so is independent of $\gamma$. Note that $S_{N}$ has already split into modular and archimedean components, with the first piece being $\nu_{q}(a)$ as in (4.1).

We continue to massage the archimedean component. Fix $\gamma$ and $u$. For any $m \in \mathbb{Z}$ with

$$
\left|m-u\left\langle\gamma e_{2}, e_{2}\right\rangle\right| \leq B\left\langle\gamma e_{2}, e_{2}\right\rangle / \mathcal{Q}^{5}
$$

we clearly have

$$
e\left(\beta u\left\langle\gamma e_{2}, e_{2}\right\rangle\right)=e(\beta m)\left(1+O\left(\mathcal{Q}^{-4}\right)\right)
$$

and there are $2 B\left\langle\gamma e_{2}, e_{2}\right\rangle / \mathcal{Q}^{5}+O(1)$ integers $m$ in this range. Hence

$$
\begin{equation*}
e\left(\beta u\left\langle\gamma e_{2}, e_{2}\right\rangle\right)=\frac{\mathcal{Q}^{5}}{2 B\left\langle\gamma e_{2}, e_{2}\right\rangle} \sum_{m \in \mathbb{Z}} e(\beta m)\left(1+O\left(\mathcal{Q}^{-4}\right)\right) . \tag{4.18}
\end{equation*}
$$

Reversing the $u$ and $m$ sums and inserting (4.18) into (4.17) gives

$$
S_{N}\left(\frac{a}{q}+\beta\right)=\nu_{q}(a) \varpi_{N}(\beta)\left(1+O\left(\mathcal{Q}^{-4}\right)\right)
$$

where

$$
\begin{equation*}
\varpi_{N}(\beta):=\frac{|\aleph|}{|\mathcal{U}|} \sum_{\gamma \in \Omega^{\prime}} \frac{\mathcal{Q}^{5}}{2 B\left\langle\gamma e_{2}, e_{2}\right\rangle} \sum_{m \in \mathbb{Z}} e(\beta m) \sum_{u \in \mathcal{U}} \mathbf{1}_{\left|\frac{m}{\left\langle\gamma e_{2}, e_{2}\right\rangle}-u\right| \leq \frac{B}{\mathcal{Q}^{5}}} \tag{4.19}
\end{equation*}
$$

Hence (4.5) is satisfied. Also the Fourier transform

$$
\begin{equation*}
\widehat{\varpi}_{N}(n)=\frac{|\aleph|}{|\mathcal{U}|} \sum_{\gamma \in \Omega^{\prime}} \frac{\mathcal{Q}^{5}}{2 B\left\langle\gamma e_{2}, e_{2}\right\rangle} \sum_{u \in \mathcal{U}} \mathbf{1}_{\left|\frac{n}{\left\langle\gamma e_{2}, e_{2}\right\rangle}-u\right| \leq \frac{B}{\mathcal{Q}^{5}}} \tag{4.20}
\end{equation*}
$$

is clearly real and nonnegative, so (1) is satisfied.
Now, combining (4.11), (2.9), (2.2), (2.3), (2.8), (3.41), and (3.36), we have

$$
\begin{equation*}
\frac{1}{4} \frac{N}{B}<\left\langle\gamma e_{2}, e_{2}\right\rangle<4 \frac{N}{B} \tag{4.21}
\end{equation*}
$$

and hence for $\frac{1}{25} N<n<\frac{1}{5} N$, we have, crudely, that

$$
\frac{1}{100} B<\frac{n}{\left\langle\gamma e_{2}, e_{2}\right\rangle}<\frac{99}{100} B
$$

Hence by the spacing in (3.7) of $u \in \mathcal{U}$ in this range, the innermost sum in (4.20) is guaranteed to have at least one contribution, giving

$$
\widehat{\varpi}_{N}(n) \gg \frac{|\aleph|}{|\mathcal{U}|} \sum_{\gamma \in \Omega^{\prime}} \frac{\mathcal{Q}^{5}}{2 B\left\langle\gamma e_{2}, e_{2}\right\rangle} \gg \frac{|\aleph|\left|\Omega^{\prime}\right|}{N}=\frac{\left|\Omega_{N}\right|}{N}
$$

where we used (4.21), (3.8), and (4.6). So (4.4) is satisfied, and the proof of Theorem 4.2 is complete.
4.2. The major arcs contribution. Equipped with (4.5), it is now straightforward to produce the necessary major arcs contribution. For technical reasons, we need a smoothed cutoff, and we introduce the triangle function, $\psi$, given by

$$
\psi(x):= \begin{cases}1+x, & \text { if }-1<x<0  \tag{4.22}\\ 1-x, & \text { if } 0 \leq x<1 \\ 0, & \text { otherwise }\end{cases}
$$

It is well known that the Fourier transform is nonnegative:

$$
\begin{equation*}
\widehat{\psi}(y)=\left(\frac{\sin (\pi y)}{\pi y}\right)^{2} \tag{4.23}
\end{equation*}
$$

Let $\psi_{N}$ be the function localized at level $\mathcal{Q} / N$ near the origin:

$$
\psi_{N}(x):=\psi\left(\frac{N}{\mathcal{Q}} x\right) .
$$

Periodize $\psi_{N}$ to $\Psi_{N}$ on $\mathbb{R} / \mathbb{Z}$ :

$$
\Psi_{N}(\theta):=\sum_{m \in \mathbb{Z}} \psi_{N}(\theta+m),
$$

and put each such spike at a major arc:

$$
\begin{equation*}
\Psi_{\mathcal{Q}, N}(\theta):=\sum_{q<\mathcal{Q}} \sum_{(a, q)=1} \Psi_{N}\left(\theta-\frac{a}{q}\right) . \tag{4.24}
\end{equation*}
$$

Note that the support of $\Psi_{\mathcal{Q}, N}$ is $\mathfrak{M}_{\mathcal{Q}}$.
As in (1.35), write the representation number

$$
R_{N}(n):=\widehat{S}_{N}(n)=\int_{0}^{1} S_{N}(\theta) e(-n \theta) d \theta
$$

and decompose it into a (smoothed) major arcs contribution and an error

$$
\begin{equation*}
R_{N}(n)=\mathcal{M}_{N}(n)+\mathcal{E}_{N}(n) \tag{4.25}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{M}_{N}(n):=\int_{0}^{1} \Psi_{\mathcal{Q}, N}(\theta) S_{N}(\theta) e(-n \theta) d \theta \tag{4.26}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{E}_{N}(n):=\int_{0}^{1}\left(1-\Psi_{\mathcal{Q}, N}(\theta)\right) S_{N}(\theta) e(-n \theta) d \theta \tag{4.27}
\end{equation*}
$$

The ultimate goal of this section is to prove the following
Theorem 4.28. For $\frac{1}{20} N \leq n<\frac{1}{10} N$,

$$
\begin{equation*}
\mathcal{M}_{N}(n) \gg \frac{1}{\log \log N} \frac{\left|\Omega_{N}\right|}{N} . \tag{4.29}
\end{equation*}
$$

Proof. Fix $\frac{1}{20} N \leq n<\frac{1}{10} N$. Starting with (4.26), insert (4.24) and (4.5) (recall supp $\Psi_{\mathcal{Q}, N} \subset \mathfrak{M}_{\mathcal{Q}}$ ), and make the change of variables $\beta=\theta-a / q$ :

$$
\begin{align*}
\mathcal{M}_{N}(n)= & \sum_{q<\mathcal{Q}} \sum_{(a, q)=1} \nu_{q}(a) e\left(-n \frac{a}{q}\right)  \tag{4.30}\\
& \times \int_{0}^{1} \Psi_{N}(\beta) \varpi_{N}(\beta) e(-n \beta) d \beta \\
& +O\left(\mathcal{Q} \mathcal{Q} \frac{\mathcal{Q}}{N}\left|\Omega_{N}\right| \mathcal{Q}^{-4}\right)
\end{align*}
$$

where we used (4.3).
Note that $\mathcal{M}_{N}$ has already split (up to acceptable error) into the product of the singular series

$$
\begin{equation*}
\mathfrak{S}_{\mathcal{Q}}(n):=\sum_{q<\mathcal{Q}} \sum_{(a, q)=1} \nu_{q}(a) e\left(-n \frac{a}{q}\right) \tag{4.31}
\end{equation*}
$$

and the singular integral

$$
\begin{align*}
\Pi_{N}(n) & :=\int_{0}^{1} \Psi_{N}(\beta) \varpi_{N}(\beta) e(-n \beta) d \beta=\sum_{m \in \mathbb{Z}} \widehat{\Psi}_{N}(n-m) \widehat{\varpi}_{N}(m)  \tag{4.32}\\
& =\frac{\mathcal{Q}}{N} \sum_{m \in \mathbb{Z}} \widehat{\psi}\left(\frac{\mathcal{Q}}{N}(n-m)\right) \widehat{\varpi}_{N}(m) .
\end{align*}
$$

First we sketch an analysis of the singular series, which is standard. Insert (4.1) into (4.31):

$$
\mathfrak{S}_{\mathcal{Q}}(n)=\sum_{q<\mathcal{Q}} \frac{1}{\left|\operatorname{SL}_{2}(q)\right|} \sum_{\gamma \in \mathrm{SL}_{2}(q)} c_{q}\left(\left\langle\gamma e_{2}, e_{2}\right\rangle-n\right),
$$

where $c_{q}$ is the classical Ramanujan sum

$$
c_{q}(m)=\sum_{(a, q)=1} e(a m / q) .
$$

Recall that $c_{q}$ is multiplicative in $q$ and that $c_{q}(m)=\mu(q)$ if $(m, q)=1$. (Here $\mu$ is the Möbius function.) Hence we may extend the range of the sum $q<\mathcal{Q}$ to $q<\infty$ with a negligible error, obtaining a sum which factors into an Euler product. At each place, the contribution from prime powers is negligible. We are left to analyze

$$
\begin{align*}
\mathfrak{S}_{\mathcal{Q}}(n) & \gg \mathfrak{S}(n) \gg \prod_{p}\left(1+\frac{1}{\left|\mathrm{SL}_{2}(p)\right|} \sum_{\gamma \in \mathrm{SL}_{2}(p)} c_{p}\left(\left\langle\gamma e_{2}, e_{2}\right\rangle-n\right)\right)  \tag{4.33}\\
& =\prod_{p \nmid n}\left(1+\frac{1}{p^{2}-1}\right) \prod_{p \mid n}\left(1-\frac{1}{p+1}\right) \gg \frac{1}{\log \log n} .
\end{align*}
$$

Returning to (4.32), we now analyze the singular integral. By positivity and using $(4.23)$ that $\widehat{\psi}(y)>2 / 5$ for $|y|<1 / 2$, we have

$$
\Pi_{N}(n)=\frac{\mathcal{Q}}{N} \sum_{m \in \mathbb{Z}} \widehat{\psi}\left(\frac{\mathcal{Q}}{N}(n-m)\right) \widehat{\varpi}_{N}(m) \geq \frac{2}{5} \frac{\mathcal{Q}}{N} \sum_{|m-n|<N /(2 \mathcal{Q})} \widehat{\varpi}_{N}(m)
$$

For $N$ (and hence $\mathcal{Q}$ ) sufficiently large, the ranges $n / N \in\left[\frac{1}{20}, \frac{1}{10}\right]$ and $|m-n|<$ $N /(2 \mathcal{Q})$ force $m / N \in\left[\frac{1}{25}, \frac{1}{5}\right]$, so (4.4) applies, giving

$$
\begin{equation*}
\Pi_{N}(n) \gg \frac{\mathcal{Q}}{N} \frac{N}{2 \mathcal{Q}} \frac{\left|\Omega_{N}\right|}{N} \gg \frac{\left|\Omega_{N}\right|}{N} \tag{4.34}
\end{equation*}
$$

Inserting (4.34) and (4.33) into (4.30) gives (4.29), as claimed.

## 5. Minor arcs analysis I

We keep all the notation from the previous section. Having dealt with the main term (4.26), we are now tasked with estimating the error $\mathcal{E}_{N}$ in (4.27). As discussed in (1.39)-(1.41), the key goal is to estimate

$$
\sum_{n \in \mathbb{Z}}\left|\mathcal{E}_{N}(n)\right|^{2}=\int_{0}^{1}\left|1-\Psi_{\mathcal{Q}, N}(\theta)\right|^{2}\left|S_{N}(\theta)\right|^{2} d \theta
$$

where we applied Parseval's formula. For $\theta$ outside the major $\operatorname{arcs}, \Psi_{\mathcal{Q}, N}$ vanishes, and we decompose the above integral into regions

$$
W_{Q, K}:=\left\{\theta=\frac{a}{q}+\beta: \frac{1}{2} Q \leq q<Q,(a, q)=1, \frac{K}{2 N} \leq|\beta|<\frac{K}{N}\right\}
$$

Here the parameters $Q$ and $K$ range dyadically in

$$
\begin{equation*}
Q<N^{1 / 2}, \quad K<N^{1 / 2} / Q \tag{5.1}
\end{equation*}
$$

If $K=O(1)$, we replace the condition $\frac{1}{2} K / N \leq|\beta|<K / N$ in $W_{Q, K}$ by just $|\beta|<K / N$, and any appearances of $K$ should be replaced by 1.

In this section, we give two bounds for $S_{N}(\theta)$, similar to Theorems 5.1 and 6.1 of [BK10]. These will suffice as long as $Q$ or $K$ is large.
5.1. The bound for $K$ large. We will first bound $\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta$ by pulling out the largest value of the integrand and multiplying by the measure of the domain, which $\ll Q^{2} \frac{K}{N}$. To get the desired bound of $N^{4 \delta-1}$ (see (1.39)), we need to bound the sup norm of $S_{N}$ on $W_{Q, K}$ by a bit less than $N^{2 \delta} /\left(K^{1 / 2} Q\right)$. We will win by an extra $K^{1 / 2}$.

Proposition 5.2. Let $N, Q, K$ be as above, and write $\theta=\frac{a}{q}+\beta \in W_{Q, K}$. Then

$$
\begin{equation*}
\left|S_{N}(\theta)\right| \ll N^{2 \delta}\left(\frac{N^{1-\delta}}{K Q}\right) \tag{5.3}
\end{equation*}
$$

as $N \rightarrow \infty$.

Proof. This is a simplified version of Theorem 5.1 in [BK10]; we repeat the arguments. By (3.37), we decompose

$$
\begin{equation*}
\Omega_{N}=\Xi_{1}\left(\Xi_{2} \Xi_{3} \cdots \Xi_{J}\right)=\Xi_{1} \cdot \Omega^{\prime} \tag{5.4}
\end{equation*}
$$

Then by (2.3), (2.2), (2.8), Lemma 3.38, and (3.35), we have for $\gamma \in \Xi_{1}$ and $\omega \in \Omega^{\prime}$ that

$$
\begin{equation*}
\left.\right|^{t} \gamma e_{2}\left|,\left|\omega e_{2}\right|<50 N^{1 / 2} .\right. \tag{5.5}
\end{equation*}
$$

Note also from (3.26) that

$$
\begin{equation*}
\# \Xi_{1}, \# \Omega^{\prime} \ll N^{\delta} . \tag{5.6}
\end{equation*}
$$

Then we can rewrite $S_{N}(\theta)$ as

$$
\begin{equation*}
S_{N}(\theta)=\sum_{x \in \mathbb{Z}^{2}} \sum_{y \in \mathbb{Z}^{2}} \mu(x) \nu(y) e(\theta\langle x, y\rangle), \tag{5.7}
\end{equation*}
$$

where $\mu$ and $\nu$ are image measures in $\mathbb{Z}^{2}$ defined by

$$
\mu(x):=\sum_{\gamma \in \Xi_{1}} \mathbf{1}_{\left\{x=t^{t} \gamma \cdot e_{2}\right\}},
$$

and similarly

$$
\nu(y):=\sum_{\omega \in \Omega^{\prime}} \mathbf{1}_{\left\{y=\omega \cdot e_{2}\right\}} .
$$

The projection $\omega \mapsto \omega \cdot e_{2}$ in $\nu$ is one-to-one, since it is well known that the continued fraction of a rational number, if restricted to have even length, is unique. The map $\mu$ is also one-to-one, since $\mathcal{G}$ is preserved under transposition, ${ }^{t} g \in \mathcal{G}$ for $g \in \mathcal{G}$ (since its generators (1.28) are fixed by transposition). Hence we have

$$
\begin{equation*}
\|\mu\|_{\infty} \leq 1, \quad\|\nu\|_{\infty} \leq 1 \tag{5.8}
\end{equation*}
$$

Note that for any $y, y^{\prime} \in \operatorname{supp} \nu$, we have from (5.5) that $\left|y-y^{\prime}\right|<100 N^{1 / 2}$. Decompose $\nu$ into 100000 blocks $\nu=\sum_{\alpha} \nu^{(\alpha)}$ so that for each $\alpha$ and any $y, y^{\prime} \in \operatorname{supp} \nu^{(\alpha)}$,

$$
\begin{equation*}
\left|y-y^{\prime}\right|<\frac{1}{2} N^{1 / 2} . \tag{5.9}
\end{equation*}
$$

Write $\left|S_{N}(\theta)\right| \leq \sum_{\alpha}\left|S_{N}^{(\alpha)}(\theta)\right|$, where

$$
S_{N}^{(\alpha)}(\theta):=\sum_{x} \sum_{y} \mu(x) \nu^{(\alpha)}(y) e(\theta\langle x, y\rangle) .
$$

We will bound each such $S_{N}^{(\alpha)}$ independently of $\alpha$, so we drop the superscripts $\alpha$.
Let $\Upsilon: \mathbb{R}^{2} \rightarrow \mathbb{R}_{+}$be a smooth test function which exceeds 1 on the square $[-1,1] \times[-1,1]$ and has Fourier transform supported in a ball of radius 1 about
the origin. Apply Cauchy-Schwarz in the $x$ variable, insert $\Upsilon$, and open the squares:
$\left|S_{N}(\theta)\right| \ll\left(\sum_{x} \mu^{2}(x)\right)^{1 / 2}\left(\sum_{x} \Upsilon\left(\frac{x}{50 N^{1 / 2}}\right) \sum_{y} \nu(y) \sum_{y^{\prime}} \nu\left(y^{\prime}\right) e\left(\left\langle x, y-y^{\prime}\right\rangle \theta\right)\right)^{1 / 2}$.
The first parentheses contribute $N^{\delta / 2}$ by (5.6) and (5.8). To the last sum on $x$ apply Poisson summation, recalling the support of $\widehat{\Upsilon}$ :

$$
\begin{equation*}
\left|S_{N}(\theta)\right| \ll N^{\delta / 2}\left(\sum_{y} \nu(y) \sum_{y^{\prime}} \nu\left(y^{\prime}\right) N \mathbf{1}_{\left\{\left\|\left(y-y^{\prime}\right) \theta\right\|<\frac{1}{\left.50 N^{1 / 2}\right\}}\right.}\right)^{1 / 2} \tag{5.10}
\end{equation*}
$$

Here $\|\cdot\|$ is the distance to the nearest lattice point in $\mathbb{Z}^{2}$. For such $y, y^{\prime}, \theta$, we have

$$
\left\|\left(y-y^{\prime}\right) \frac{a}{q}\right\| \leq\left\|\left(y-y^{\prime}\right) \theta\right\|+\left|y-y^{\prime} \| \beta\right|<\frac{1}{50 N^{1 / 2}}+\frac{1}{2} N^{1 / 2} \frac{K}{N}<\frac{1}{Q},
$$

where we used (5.9) and (5.1). Then $q<Q$ forces $\left\|\left(y-y^{\prime}\right) \frac{a}{q}\right\|=0$, or

$$
y \equiv y^{\prime}(q) .
$$

This being the case, we now have

$$
\frac{1}{50 N^{1 / 2}}>\left\|\left(y-y^{\prime}\right) \theta\right\|=\left|\left(y-y^{\prime}\right) \beta\right| ;
$$

that is,

$$
\left|y-y^{\prime}\right| \ll \frac{N^{1 / 2}}{K}
$$

In summary, we have

$$
\left|S_{N}(\theta)\right| \ll N^{(\delta+1) / 2}\left(\sum_{y} \nu(y) \sum_{y^{\prime}} \mathbf{1}_{\left\{\begin{array}{c}
y \equiv y^{\prime}(q) \\
\left|y-y^{\prime}\right| \ll \frac{N^{1 / 2}}{K}
\end{array}\right\}}\right)^{1 / 2}
$$

where we used (5.8). Using $Q<\frac{N^{1 / 2}}{K}$ and the crudest bound on the $y^{\prime}$ sum gives

$$
\left|S_{N}(\theta)\right| \ll N^{(\delta+1) / 2}\left(\sum_{y} \nu(y)\left(\frac{N^{1 / 2}}{Q K}\right)^{2}\right)^{1 / 2} \ll \frac{N^{\delta+1}}{Q K}
$$

as claimed.
The bound (5.3) is already conclusive if $K$ is a bit larger than $N^{2(1-\delta)}$.
Theorem 5.11. Assume $Q<N^{1 / 2}$ and $K<N^{1 / 2} / Q$. Then

$$
\begin{equation*}
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N}\left[\frac{N^{2(1-\delta)+4 \mathfrak{b}}}{K}\right] . \tag{5.12}
\end{equation*}
$$

Proof. We bound trivially using (5.3):

$$
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \frac{K}{N} Q^{2}\left(\frac{N^{\delta+1}}{Q K}\right)^{2} \ll N^{4 \delta-1}\left[\frac{N^{2(1-\delta)}}{K}\right],
$$

and the claim follows from (3.47), on crudely using $\delta<1$.
5.2. Another bilinear forms estimate. Next we introduce the cross-section of $W_{Q, K}$ for given $\beta$ :

$$
P_{Q, \beta}:=\left\{\theta=\frac{a}{q}+\beta: \frac{1}{2} Q \leq q<Q,(a, q)=1\right\} .
$$

We will bound using (5.3), giving essentially

$$
\int_{W_{Q, K}}\left|S_{N}\right|^{2} \ll \sup \left|S_{N}\right| \frac{K}{N} \sup _{\beta} \sum_{P_{Q, \beta}}\left|S_{N}\right| \ll \frac{N^{2 \delta+}}{K Q} \frac{K}{N} \sup _{\beta} \sum_{P_{Q, \beta}}\left|S_{N}\right| .
$$

The trivial bound on $\sum_{P_{Q, \beta}}\left|S_{N}\right|$ is of course $N^{2 \delta} Q^{2}$, so we need to save a little more than a power of $Q$ to get our target bound of less than $N^{4 \delta-1}$. This is achieved by exploiting the extra structure in the $a$ and $q$ sums, as follows.

Proposition 5.13. Let the notation be as above. Then for all $\varepsilon>0$,

$$
\begin{equation*}
\sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right|<_{\varepsilon} N^{2 \delta} Q^{2} N^{1-\delta+\varepsilon}\left[\frac{1}{Q^{3 / 2}}+\frac{1}{Q N^{1 / 8}}\right] . \tag{5.14}
\end{equation*}
$$

Proof. The proof is nearly identical to that of Theorem 6.1 in [BK10], but we reproduce it for the reader's convenience. We again use (3.37) to decompose $\Omega_{N}$ into pieces, now grouping by

$$
\Omega_{N}=\left(\Xi_{1} \Xi_{2}\right)\left(\Xi_{3} \cdots \Xi_{J}\right)=\Omega^{\prime} \cdot \Omega^{\prime \prime} .
$$

As before, we have for $\gamma \in \Omega^{\prime}$ and $\omega \in \Omega^{\prime \prime}$ that

$$
\begin{equation*}
\left|{ }^{t} \gamma e_{2}\right|<300 N^{3 / 4} \quad \text { and } \quad\left|\omega e_{2}\right|<2000 N^{1 / 4} \tag{5.15}
\end{equation*}
$$

Also from (3.26), we have

$$
\begin{equation*}
\# \Omega^{\prime} \ll N^{3 \delta / 2} \quad \text { and } \quad \# \Omega^{\prime \prime} \ll N^{\delta / 2} \tag{5.16}
\end{equation*}
$$

Again we define the measures $\mu$ and $\nu$ on $\mathbb{Z}^{2}$ by

$$
\begin{aligned}
\mu(x) & :=\sum_{\gamma \in \Omega^{\prime}} \mathbf{1}_{\left\{x=t \gamma e_{2}\right\}}, \\
\nu(y) & :=\sum_{\omega \in \Omega^{\prime \prime}} \mathbf{1}_{\left\{y=\omega e_{2}\right\}},
\end{aligned}
$$

with $\mu, \nu \leq 1$. For any two elements $y, y^{\prime}$ in the support of $\nu$, we have $\left|y-y^{\prime}\right|<$ $4000 N^{1 / 4}$. Hence we again decompose $\nu$ into $O(1)$ pieces, $\nu=\sum_{\alpha} \nu^{(\alpha)}$, so as
to make the difference

$$
\begin{equation*}
\left|y-y^{\prime}\right|<\frac{1}{10000} N^{1 / 4} \tag{5.17}
\end{equation*}
$$

for $y, y^{\prime}$ in the support of $\nu^{(\alpha)}$. Writing

$$
S_{N}^{(\alpha)}(\theta)=\sum_{x} \sum_{y} \mu(x) \nu^{(\alpha)}(y) e(\theta\langle x, y\rangle),
$$

and dropping the superscripts $\alpha$, we proceed to bound

$$
\begin{aligned}
\sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right| & =\sum_{q \asymp Q} \sum_{(a, q)=1} \zeta(\theta) S_{N}(\theta) \\
& =\sum_{q \asymp Q} \sum_{(a, q)=1} \zeta(\theta) \sum_{x} \sum_{y} \mu(x) \nu(y) e(\theta\langle x, y\rangle),
\end{aligned}
$$

where $\zeta$ has modulus 1. Recall the bump function $\Upsilon$ which is at least one on $[-1,1]^{2}$; assume now that its Fourier transform is supported in a ball of radius $1 / 40$ about the origin. Apply Cauchy-Schwarz in the $x$ sum and (5.16), insert the function $\Upsilon$, reverse orders, and apply Poisson summation:

$$
\begin{align*}
& \sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right|  \tag{5.18}\\
& \ll N^{3 \delta / 4}\left(\sum_{x} \Upsilon\left(\frac{x}{300 N^{3 / 4}}\right)\left|\sum_{q \asymp Q} \sum_{(a, q)=1} \zeta(\theta) \sum_{y} \nu(y) e(\theta\langle x, y\rangle)\right|^{2}\right)^{1 / 2} \\
& \ll N^{3(\delta+1) / 4} \mathcal{X}^{1 / 2},
\end{align*}
$$

where

$$
\begin{equation*}
\mathcal{X}=\mathcal{X}_{Q, \beta}:=\sum_{q} \sum_{q^{\prime}} \sum_{a} \sum_{a^{\prime}} \sum_{y} \sum_{y^{\prime}} \nu(y) \nu\left(y^{\prime}\right) \mathbf{1}_{\left\{\left\|y \theta-y^{\prime} \theta^{\prime}\right\|<\frac{1}{12000 N^{3 / 4}}\right\}} . \tag{5.19}
\end{equation*}
$$

Here $\theta^{\prime}=\frac{a^{\prime}}{q^{\prime}}+\beta$; note that $\beta$ is the same for $\theta$ and $\theta^{\prime}$.
Write $y=\left(y_{1}, y_{2}\right)$ and the same with $y^{\prime}$. Consider the innermost condition in (5.19):

$$
\begin{equation*}
\left\|y_{1} \theta-y_{1}^{\prime} \theta^{\prime}\right\|,\left\|y_{2} \theta-y_{2}^{\prime} \theta^{\prime}\right\|<\frac{1}{12000 N^{3 / 4}} \tag{5.20}
\end{equation*}
$$

Recall that $y=\gamma e_{2}$ for some (nonidentity) $\gamma \in \Gamma$, and the same for $y^{\prime}$; hence we have

$$
y_{1} y_{2} y_{1}^{\prime} y_{2}^{\prime} \neq 0
$$

Also note using (5.20), (5.17), and $|\beta|<K / N<1 /\left(N^{1 / 2} Q\right)$ that

$$
\begin{equation*}
\left\|y_{1} \frac{a}{q}-y_{1}^{\prime} \frac{a^{\prime}}{q^{\prime}}\right\| \leq\left\|y_{1} \theta-y_{1}^{\prime} \theta^{\prime}\right\|+\left|\left(y_{1}-y_{1}^{\prime}\right) \beta\right|<\frac{1}{12000 N^{3 / 4}}+\frac{N^{1 / 4}}{10000 N^{1 / 2} Q}, \tag{5.21}
\end{equation*}
$$

and similarly with $y_{2}, y_{2}^{\prime}$.

Let $Y:=\left(\begin{array}{ll}y_{1} & y_{1}^{\prime} \\ y_{2} & y_{2}^{\prime}\end{array}\right)$, so that

$$
\begin{equation*}
\mathcal{Y}:=\operatorname{det}(Y)=y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2} . \tag{5.22}
\end{equation*}
$$

Observe then by (5.21), (5.15), and $Q<N^{1 / 2}$ that

$$
\begin{aligned}
\left\|\mathcal{Y} \frac{a}{q}\right\| & \leq\left\|y_{2}^{\prime}\left(y_{1} \frac{a}{q}-y_{1}^{\prime} \frac{a^{\prime}}{q^{\prime}}\right)\right\|+\left\|y_{1}^{\prime}\left(y_{2}^{\prime} \frac{a^{\prime}}{q^{\prime}}-y_{2} \frac{a}{q}\right)\right\| \\
& <2000 N^{1 / 4}\left(\frac{1}{12000 N^{3 / 4}}+\frac{N^{1 / 4}}{10000 N^{1 / 2} Q}\right) \times 2 \\
& <\frac{1}{Q} .
\end{aligned}
$$

Of course this forces $\mathcal{Y} \equiv 0(\bmod q)$. The same argument gives $\mathcal{Y} \equiv 0\left(\bmod q^{\prime}\right)$, and hence we have

$$
\begin{equation*}
\mathcal{Y} \equiv 0(\bmod \mathfrak{q}), \tag{5.23}
\end{equation*}
$$

where $\frac{1}{2} Q \leq \mathfrak{q}<Q^{2}$ is the least common multiple of $q$ and $q^{\prime}$.
Decompose $\mathcal{X}$ in (5.19) as $\mathcal{X}=\mathcal{X}_{1}+\mathcal{X}_{2}$ according to whether $\mathcal{Y}=0$ or not; we handle the two contributions separately. We will prove the following two lemmata.

Lemma 5.24. For any $\varepsilon>0$,

$$
\mathcal{X}_{1} \ll_{\varepsilon} N^{\delta / 2+\varepsilon} Q^{4}\left[\frac{1}{N^{3 / 4}}+Q^{-2}\right] .
$$

Lemma 5.25. For any $\varepsilon>0$,

$$
\mathcal{X}_{2}<_{\varepsilon} N^{\delta+\varepsilon} Q .
$$

We momentarily postpone the proofs of these two lemmata, first using them to finish the proof of Proposition 5.13. Returning to (5.18), we have

$$
\sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right|<_{\varepsilon} N^{3(\delta+1) / 4+\varepsilon}\left[N^{\delta / 2} Q^{4}\left(\frac{1}{N^{3 / 4}}+Q^{-2}\right)+N^{\delta} Q\right]^{1 / 2}
$$

from which the claim follows using $Q<N^{1 / 2}$.
Now we establish the lemmata separately.

### 5.2.1. Bounding $\mathcal{X}_{2}$ : the case $\mathcal{Y} \neq 0$.

Proof of Lemma 5.25. Note from (5.15), (5.17), and (5.22) that

$$
|\mathcal{Y}| \leq\left|y_{1}\left(y_{2}^{\prime}-y_{2}\right)\right|+\left|\left(y_{1}-y_{1}^{\prime}\right) y_{2}\right|<2000 N^{1 / 4} \frac{1}{10000} N^{1 / 4} \times 2<N^{1 / 2} .
$$

Since $\mathfrak{q} \mid \mathcal{Y}$ and $\mathcal{Y} \neq 0$, we have

$$
\mathfrak{q} \leq \min \left(Q^{2}, N^{1 / 2}\right) \leq Q N^{1 / 4} .
$$

Then (5.21) and $Q<N^{1 / 2}$ forces

$$
\begin{equation*}
y_{1} \frac{a}{q}-y_{1}^{\prime} \frac{a^{\prime}}{q^{\prime}} \equiv 0(\bmod 1), \tag{5.26}
\end{equation*}
$$

and the same holds for $y_{2}, y_{2}^{\prime}$. Let $\tilde{q}:=\left(q, q^{\prime}\right)$ and $q=q_{1} \tilde{q}, q^{\prime}=q_{1}^{\prime} \tilde{q}$ so that $\mathfrak{q}=q_{1} q_{1}^{\prime} \tilde{q}$. Then (5.26) becomes

$$
y_{1} a q_{1}^{\prime} \equiv y_{1}^{\prime} a^{\prime} q_{1}(\bmod \mathfrak{q})
$$

and the same for $y_{2}, y_{2}^{\prime}$. Recall $a$ and $q$ are coprime, as are $a^{\prime}$ and $q^{\prime}$. It then follows that $q_{1} \mid y_{1}$, and similarly, $q_{1} \mid y_{2}$. But since the coordinates of $y$ are coprime, $\left(y_{1}, y_{2}\right)=1$, this forces $q_{1}=1$. The same argument applies to $q_{1}^{\prime}$, so we have $q=q^{\prime}=\mathfrak{q}$. Then (5.26) now reads

$$
\begin{equation*}
y_{1} a \equiv y_{1}^{\prime} a^{\prime}(\bmod q), \tag{5.27}
\end{equation*}
$$

and similarly for $y_{2}, y_{2}^{\prime}$.
Hence, once we fix $y, y^{\prime} \in \Omega^{\prime \prime} e_{2}$, the value of $\mathcal{Y}$ is determined, and $q \mid \mathcal{Y}$ leaves at most $N^{\varepsilon}$ choices for $q$. Then there are at most $Q$ choices for $a$, from which $a^{\prime}$ is determined by (5.27) (again using the coprimality of the coordinates of $y$ and $y^{\prime}$ ).

Then using (5.16), $\mathcal{X}_{2}$ is bounded by

$$
\begin{align*}
\mathcal{X}_{2} & \ll \sum_{y} \nu(y) \sum_{y^{\prime}} \nu\left(y^{\prime}\right) \sum_{\substack{q \left\lvert\, \mathcal{Y} \\
\frac{1}{2} Q \leq q<Q\right.}} \sum_{a(\bmod q)} 1  \tag{5.28}\\
& \ll \varepsilon\left(N^{\delta / 2}\right)^{2} N^{\varepsilon} Q,
\end{align*}
$$

as claimed.
5.2.2. Bounding $\mathcal{X}_{1}$ : the case $\mathcal{Y}=0$.

Proof of Lemma 5.24. The condition $\mathcal{Y}=0$ implies $y_{1} / y_{2}=y_{1}^{\prime} / y_{2}^{\prime}$. Recall that rationals have unique continued fraction expansions (of even length), and thus $y=y^{\prime}$. The bottom line savings from this fact is at most $N^{1 / 4}$, whereas we need to save a bit more than $Q$, which can be as large as $N^{1 / 2}$.

Let $N^{\prime}:=\frac{1}{12000} N^{3 / 4}$. The condition (5.20) then becomes

$$
\begin{equation*}
\left\|y_{1}\left(\frac{a}{q}-\frac{a^{\prime}}{q^{\prime}}\right)\right\|<\frac{1}{N^{\prime}} . \tag{5.29}
\end{equation*}
$$

Let $\left(y_{1}, q\right)=v$ and $\left(y_{1}, q^{\prime}\right)=v^{\prime}$ with $q=v r$. Assume without loss of generality that $v \leq v^{\prime}$. Fix $y$ (for which there are $N^{\delta / 2}$ choices) and $v, v^{\prime} \mid y_{1}$ (at most $N^{\varepsilon}$ choices). There are $\ll Q / v^{\prime}$ choices for $q^{\prime} \equiv 0\left(\bmod v^{\prime}\right)$, and then $\ll Q$ choices for $\left(a^{\prime}, q^{\prime}\right)=1$. Write $\psi$ for $y_{1} a^{\prime} / q^{\prime}(\bmod 1)$, which is now fixed, and write $y_{1}=v z$ with $(z, r)=1$. Then (5.29) becomes

$$
\left\|z \frac{a}{r}-\psi\right\|<\frac{1}{N^{\prime}} .
$$

Let $\mathcal{U}_{z}$ be the set of possible fractions $\frac{z a}{r}(\bmod 1)$ as $r$ varies in $Q /(2 v) \leq r<$ $Q / v$, and $a$ ranges up to $Q$ subject to $(a, v r)=1$. Note that distinct points
$u \in \mathcal{U}_{z}$ are separated by a distance of at least $v^{2} / Q^{2}$. Hence the size of the intersection of $\mathcal{U}_{z}$ with the interval

$$
\left[\psi-\frac{1}{N^{\prime}}, \psi+\frac{1}{N^{\prime}}\right]
$$

contains at most $\frac{Q^{2}}{v^{2} N^{\prime}}+1$ points. Once $u=f / r \in \mathcal{U}_{z}$ is determined, so is its denominator; that is, $r$ is determined. Also $a(\bmod r)$ is determined (to be $f$ ), hence $a(\bmod q)$ has $v$ possible values (recall $q=r v)$.

In summary, we use (5.16) again to bound $\mathcal{X}_{1}$ by

$$
\begin{aligned}
\mathcal{X}_{1} & \ll \sum_{y} \nu(y) \sum_{\substack{v, v^{\prime} \mid y_{1} \\
v \leq v^{\prime}}} \sum_{\substack{q^{\prime} \equiv 0\left(\bmod v^{\prime}\right)}} \sum_{\substack{\left.a<q \\
a^{\prime}, q^{\prime}\right)=1}} 1 \\
& \ll \sum_{y} \nu(y) \sum_{\substack{v, v^{\prime} \mid y_{1} \\
v \leq v^{\prime}}} \frac{Q}{v^{\prime}} Q\left(\frac{Q^{2}}{v^{2} N^{\prime}}+1\right) v \\
& \ll \varepsilon \mathcal{U}_{z} \cap \psi-\frac{1}{q-r N^{\prime}}, \psi+\frac{1}{\left.N^{\prime}\right]} \\
& N^{\delta / 2} N^{\varepsilon} Q^{2}\left(\frac{Q^{2}}{N^{3 / 4}}+1\right),
\end{aligned}
$$

as claimed.
With the lemmata established, we have completed the proof of Propostion 5.13.
5.3. The bound for $Q$ large. Lastly, we input this bound to get another bound on the main integral, one which is favorable as long as $Q$ is a bit bigger than $N^{4(1-\delta)}$.

Theorem 5.30. Assume that $Q<N^{1 / 2}$ and $K Q<N^{1 / 2}$. Then

$$
\begin{equation*}
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N} N^{2(1-\delta)} N^{4 \mathfrak{b}}\left(\frac{1}{Q^{1 / 2}}+\frac{1}{N^{1 / 8}}\right) . \tag{5.31}
\end{equation*}
$$

Proof. Write

$$
\begin{aligned}
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta & \ll \sup _{\theta \in W_{Q, K}}\left|S_{N}(\theta)\right| \cdot \frac{K}{N} \sup _{|\beta| \asymp \frac{K}{N}} \sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right| \\
& \ll{ }_{\varepsilon} N^{2 \delta}\left(\frac{N^{1-\delta}}{K Q}\right) \cdot \frac{K}{N}\left(N^{2 \delta} Q^{2} N^{1-\delta+\varepsilon}\left[\frac{1}{Q^{3 / 2}}+\frac{1}{Q N^{1 / 8}}\right]\right) \\
& \ll N^{4 \delta-1} N^{2(1-\delta)+\varepsilon}\left(\frac{1}{Q^{1 / 2}}+\frac{1}{N^{1 / 8}}\right),
\end{aligned}
$$

where we used (5.3) and (5.14). The claim follows from (3.47), again crudely using $\delta<1$.

It remains to handle the regions when both $K$ and $Q$ are very small, less than $N^{\varepsilon}$ for $\varepsilon$ near zero.

## 6. Minor arcs analysis II

We now push the methods of the previous section down to the level of $Q$ and $K$ being of constant size. We again do this in two stages. But first we record the following counting bound.

Lemma 6.1. For $(q K)^{13 / 5}<Y<X$, and vectors $\eta, \eta^{\prime} \in \mathbb{Z}^{2}$ having coprime coordinates with $|\eta| \asymp X / Y$ and $\left|\eta^{\prime}\right| \asymp Y$,

$$
\#\left\{\gamma \in \mathrm{SL}_{2}(\mathbb{Z}):\|\gamma\| \asymp Y,\left|\gamma \eta-\eta^{\prime}\right|<\frac{X}{Y K}, \text { and } \gamma \eta \equiv \eta^{\prime}(\bmod q)\right\} \ll\left(\frac{Y}{q K}\right)^{2}
$$

The implied constant is absolute, depending on the implied constants above.
Sketch of proof. Write $G(\mathbb{Z})=\mathrm{SL}_{2}(\mathbb{Z})$, and let $G_{\eta}(q)$ be the stabilizer of $\eta \bmod q$ :

$$
G_{\eta}(q):=\{\gamma \in G(\mathbb{Z}): \gamma \eta \equiv \eta(q)\}
$$

Then $G(\mathbb{Z}) \cong\left(G(\mathbb{Z}) / G_{\eta}(q)\right) \times G_{\eta}(q)$. Let $R=R_{Y, K}$ denote the region

$$
R:=\left\{g \in \mathrm{SL}_{2}(\mathbb{R}):\|g\| \asymp Y,\left|g \eta-\eta^{\prime}\right|<X /(Y K)\right\}
$$

The methods in [Goo83] (see also [BKS10]) give an estimate of the form

$$
\begin{aligned}
\sum_{\gamma \in G(\mathbb{Z})} \mathbf{1}_{\{\gamma \in R\}} \mathbf{1}_{\gamma \eta \equiv \eta^{\prime}(\bmod q)} & =\sum_{\omega \in G(\mathbb{Z}) / G_{\eta}(q)} \mathbf{1}_{\omega \eta \equiv \eta^{\prime}(\bmod q)} \sum_{\gamma^{\prime} \in G_{\eta}(q)} \mathbf{1}_{\left\{\omega \gamma^{\prime} \in R\right\}} \\
& \ll \varepsilon \sum_{\omega \in G(\mathbb{Z}) / G_{\eta}(q)} \mathbf{1}_{\omega \eta \equiv \eta^{\prime}(\bmod q)}\left(\left(\frac{Y}{q K}\right)^{2}+Y^{2 \Theta+\varepsilon}\right) \\
& \ll\left(\frac{Y}{q K}\right)^{2}+Y^{2 \Theta+\varepsilon},
\end{aligned}
$$

where $\Theta=1 / 2+7 / 64$ is the best known bound towards the Ramanujan conjectures [Kim03]. (We apply the argument to a smoothed sum.) The first term dominates as long as $(q K)^{2}<Y^{25 / 32-\varepsilon}$, and the claim follows using $64 / 25+\varepsilon<65 / 25=13 / 5$.

Remark 6.2. Recall that we are not interested here in optimizing the final value of $\delta_{0}$ in Theorem 1.8 , so we allow ourselves to be a bit crude in the above for the sake of exposition.
6.1. The bound for $K$ at least a small power of $Q$. We return to the approach of Section 5.1, that is just bounding the sup norm and needing to win more than $K^{1 / 2} Q$ off the trivial bound. Now we use the fact that $K Q$ is quite small to improve on the trivial bound by $(K Q)^{1-\varepsilon}$ with $\varepsilon$ small (depending on the distance from $\delta$ to 1 ). Then we will have, roughly

$$
\int_{W_{Q, K}}\left|S_{N}\right|^{2} \ll Q^{2} \frac{K}{N}\left(\frac{\# \Omega_{N}}{(K Q)^{1-\varepsilon}}\right)^{2}=\frac{\left(\# \Omega_{N}\right)^{2}}{N}\left(\frac{Q^{\varepsilon}}{K^{1-\varepsilon}}\right)
$$

which is a savings as long as $K$ is at least a small power of $Q$. Note now for $K$ and $Q$ small that we must be careful with the loss in the size of $\Omega_{N}$ in the lower bound (3.47). We make all of this precise below.

Proposition 6.3. Assume $\theta \in W_{Q, K}$ with

$$
\begin{equation*}
1 \ll K Q<N^{5 / 52} \tag{6.4}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left|S_{N}(\theta)\right| \ll \# \Omega_{N}\left(\frac{e^{c(\log \log (K Q))^{2}}}{(K Q)^{1-(1-\delta) 52 / 5}}\right) . \tag{6.5}
\end{equation*}
$$

Proof. Recalling (3.35) that $N_{j} \asymp N^{1 / 2^{j}}$, we can find a $1 \leq j \leq J$ so that

$$
\begin{equation*}
\frac{1}{100}(Q K)^{13 / 5}<N_{j}<(Q K)^{26 / 5} \tag{6.6}
\end{equation*}
$$

say. Here we used (6.4) that $(Q K)^{26 / 5}<N^{1 / 2}$.
Define the sets

$$
\begin{align*}
& \Omega^{(1)}:=\Xi_{1} \Xi_{2} \cdots \Xi_{j-1},  \tag{6.7}\\
& \Omega^{(2)}:=\Xi_{j}, \\
& \Omega^{(3)}:=\Xi_{j+1} \Xi_{j+2} \cdots \Xi_{J} .
\end{align*}
$$

Hence for $g_{i} \in \Omega^{(i)}$,

$$
\begin{align*}
& \lambda\left(g_{3}\right) \sim N_{j+1} N_{j+2} \cdots N_{J}=: M  \tag{6.8}\\
& \lambda\left(g_{2}\right) \sim N_{j}  \tag{6.9}\\
& \lambda\left(g_{1}\right) \asymp \frac{N}{M N_{j}} . \tag{6.10}
\end{align*}
$$

Note that

$$
\begin{equation*}
\frac{N_{j}}{\log N_{j}} \ll M \ll N_{j} \log N_{j} \tag{6.11}
\end{equation*}
$$

and that from (3.48) and (3.26) we have

$$
\begin{equation*}
\left|\Omega^{(3)}\right| \gg \frac{M^{2 \delta}}{e^{c(\log \log M)^{2}}}, \quad\left|\Omega^{(2)}\right| \gg \frac{\left(N_{j}\right)^{2 \delta}}{\left(\log N_{j}\right)^{3}} . \tag{6.12}
\end{equation*}
$$

In the above, we used that $J-j \asymp \log \log M$.
Estimate

$$
\begin{equation*}
\left|S_{N}(\theta)\right| \ll \sum_{g_{1} \in \Omega^{(1)}} \sum_{g_{3} \in \Omega^{(3)}}\left|\sum_{g_{2} \in \Omega^{(2)}} e\left(\left\langle g_{3} e_{2},{ }^{t} g_{2}^{t} g_{1} e_{2}\right\rangle \theta\right)\right| . \tag{6.13}
\end{equation*}
$$

Fix $g_{1}$ and set $\eta={ }^{t} g_{1} e_{2}$. Note that

$$
\begin{equation*}
|\eta| \asymp \frac{N}{M N_{j}} . \tag{6.14}
\end{equation*}
$$

Estimate as in (5.10):

$$
\begin{align*}
& \sum_{g_{3} \in \Omega^{(3)}}\left|\sum_{g_{2} \in \Omega^{(2)}} e\left(\left\langle g_{3} e_{2},{ }^{t} g_{2} \eta\right\rangle \theta\right)\right|  \tag{6.15}\\
& \ll\left(\# \Omega^{(3)}\right)^{1 / 2} M\left[\#\left\{\left(g, g^{\prime}\right) \in{ }^{t} \Omega^{(2)} \times{ }^{t} \Omega^{(2)}: \begin{array}{l}
\left\|\left\langle\left(g-g^{\prime}\right) \eta, e_{1}\right\rangle\right\| \ll \frac{1}{M} \\
\left\|\left\langle\left(g-g^{\prime}\right) \eta, e_{2}\right\rangle\right\| \ll \frac{1}{M}
\end{array}\right\}\right]^{1 / 2},
\end{align*}
$$

where we extended the sum over $g_{3}$ to $g_{3} e_{2} \in\left\{z \in \mathbb{Z}^{2}:|z| \ll M\right\}$. Write

$$
\begin{equation*}
\left\|\left\langle\left(g-g^{\prime}\right) \eta, e_{i}\right\rangle \frac{a}{q}\right\|=\left\|\left\langle\left(g-g^{\prime}\right) \eta, e_{i}\right\rangle \theta\right\|+\left|\left\langle\left(g-g^{\prime}\right) \eta, e_{i}\right\rangle \beta\right|, \tag{6.16}
\end{equation*}
$$

where

$$
\left|\left\langle\left(g-g^{\prime}\right) \eta, e_{i}\right\rangle \beta\right| \ll N_{j} \frac{N}{M N_{j}} \frac{K}{N}=\frac{K}{M}
$$

From (6.6) and (6.11) we clearly have $\frac{K}{M}<\frac{1}{Q}$, so (6.16) forces

$$
\begin{equation*}
\left(g-g^{\prime}\right) \eta \equiv 0(q) \tag{6.17}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\left(g-g^{\prime}\right) \eta\right| \ll \frac{1}{M|\beta|} \ll \frac{N}{K M} . \tag{6.18}
\end{equation*}
$$

Fix $g^{\prime}$, and enlarge $g \in{ }^{t} \Omega^{(2)}$ to $\left\{g \in \mathrm{SL}_{2}(\mathbb{Z}):\|g\| \ll N_{j}\right\}$. Applying Lemma 6.1 with $\eta^{\prime}=g^{\prime} \eta, X=N / M$, and $Y=N_{j}$, the $g$ cardinality contributes

$$
\ll\left(\frac{N_{j}}{K Q}\right)^{2} .
$$

Thus we have by (6.12) and (6.11) that

$$
\begin{aligned}
(6.15) & \ll\left(\# \Omega^{(3)} \cdot \# \Omega^{(2)}\right)^{1 / 2} \frac{N_{j}^{2}}{K Q} \\
& \ll \# \Omega^{(3)} \cdot \# \Omega^{(2)} \frac{\left(M N_{j}\right)^{1-\delta} e^{c(\log \log M)^{2}}\left(\log N_{j}\right)^{3}}{K Q} .
\end{aligned}
$$

Hence by (6.6) and (6.11),

$$
(6.13) \ll \# \Omega_{N} \frac{(K Q)^{(1-\delta) 52 / 5} e^{c(\log \log (K Q))^{2}}}{K Q},
$$

as claimed.
Inserting this bound into the main integral and estimating trivially gives
Theorem 6.19. Assuming (6.4),

$$
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N} \frac{Q^{(1-\delta) 104 / 5} e^{c(\log \log (K Q))^{2}}}{K^{1-(1-\delta) 104 / 5}} .
$$

This bound is sufficient for our purposes unless $K$ is much less than

$$
\begin{equation*}
Q^{\frac{104 / 5(1-\delta)}{1-104 / 5(1-\delta)}} \approx Q^{\varepsilon} . \tag{6.20}
\end{equation*}
$$

6.2. The bound for $K$ even smaller. In this last section, we give the final bound for minor arcs, which we apply to the remaining range of $K$. Recall the approach of Section 5.2: we bound $\int_{W_{Q, K}}\left|S_{N}\right|^{2}$ by the sup norm times $K / N$ times $\sum_{P_{\beta, Q}}\left|S_{N}\right|$. The sup norm has already gained almost $K Q$, so we need to gain more than a power of $Q$ off of the last summation. We proceed as follows.

Proposition 6.21. Recall the cross section $P_{Q, \beta}$ for a given $|\beta| \asymp \frac{K}{N}$ :

$$
P_{Q, \beta}:=\left\{\theta=\frac{a}{q}+\beta: q \asymp Q,(a, q)=1\right\} .
$$

Then assuming (6.4), we have

$$
\begin{equation*}
\sum_{\theta \in P_{Q, \beta}}\left|S_{N}(\theta)\right| \ll \# \Omega_{N} Q^{2}\left(\frac{(K Q)^{(1-\delta) 52 / 5} e^{c(\log \log (K Q))^{2}}}{Q^{3 / 2}}\right) \tag{6.22}
\end{equation*}
$$

Proof. This argument is similar to Proposition 5.13, and we sketch the proof. Using the same decomposition (6.7), we follow (5.18) and bound the left-hand side of (6.22) by

$$
\begin{align*}
\ll & \sum_{g_{1} \in \Omega^{(1)}}\left(\# \Omega^{(3)}\right)^{1 / 2} M  \tag{6.23}\\
& \times\left[\#\left\{\left(\theta, \theta^{\prime}, g, g^{\prime}\right) \in P_{\beta} \times P_{\beta} \times{ }^{t} \Omega^{(2)} \times{ }^{t} \Omega^{(2)}:\left\|\left(g \theta-g^{\prime} \theta^{\prime}\right) \eta\right\| \ll \frac{1}{M}\right\}\right]^{1 / 2},
\end{align*}
$$

where $\eta={ }^{t} g_{1} e_{2}$. The innermost condition guarantees $q=q^{\prime}$ and

$$
a(g \eta) \equiv a^{\prime}\left(g^{\prime} \eta\right)(\bmod q)
$$

The number of choices for $g^{\prime}$ given $g, q, a$, and $a^{\prime}$ is

$$
\ll\left(\frac{N_{j}}{Q}\right)^{2},
$$

hence

$$
\begin{aligned}
(6.23) & \ll \# \Omega^{(1)}\left(\# \Omega^{(3)}\right)^{1 / 2} M\left[Q Q^{2} \# \Omega^{(2)}\left(\frac{N_{j}}{Q}\right)^{2}\right]^{1 / 2} \\
& \ll \# \Omega_{N} Q^{2}\left(\frac{(K Q)^{(1-\delta) 52 / 5} e^{c(\log \log (K Q))^{2}}}{Q^{3 / 2}}\right),
\end{aligned}
$$

as claimed.

Using (6.22) and (6.5), we now have the bound

$$
\begin{aligned}
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll & \frac{K}{N} \# \Omega_{N}\left(\frac{e^{c(\log \log (K Q))^{2}}}{(K Q)^{1-(1-\delta) 52 / 5}}\right) \\
& \times \# \Omega_{N} Q^{2}\left(\frac{(K Q)^{(1-\delta) 52 / 5} e^{c(\log \log (K Q))^{2}}}{Q^{3 / 2}}\right),
\end{aligned}
$$

from which we immediately have
Theorem 6.24. Assuming (6.4),

$$
\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N}\left(\frac{(K Q)^{(1-\delta) 104 / 5} e^{c(\log \log (K Q))^{2}}}{Q^{1 / 2}}\right)
$$

## 7. Proofs of Theorems 1.8 and 1.25

Keeping all the notation of previous sections, we now prove the main minor arcs estimate, analogous to (1.39), before completing a proof of Theorem 1.8.

Theorem 7.1. Assume

$$
\begin{equation*}
\delta>\delta_{0}, \tag{7.2}
\end{equation*}
$$

with $\delta_{0}$ given by (1.21). Then for some $c>0$,

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}}\left|\mathcal{E}_{N}(n)\right|^{2} \ll \frac{\left|\Omega_{N}\right|^{2}}{N} \mathcal{Q}^{-c} \tag{7.3}
\end{equation*}
$$

Proof. By Parseval, we have

$$
\sum_{n \in \mathbb{Z}}\left|\mathcal{E}_{N}(n)\right|^{2}=\int_{0}^{1}\left|1-\Psi_{\mathcal{Q}, N}(\theta)\right|^{2}\left|S_{N}(\theta)\right|^{2} d \theta=\int_{\mathfrak{M}_{\mathcal{Q}}}+\int_{\mathfrak{m}},
$$

where we decomposed the integral into the major arcs $\mathfrak{M}_{\mathcal{Q}}$ and the complementary minor arcs $\mathfrak{m}=[0,1] \backslash \mathfrak{M}_{\mathcal{Q}}$.

On the major arcs, note from (4.22) that $1-\psi(x)=|x|$ on $[-1,1]$. Then using (6.5) with $K \asymp N|\beta|$ gives

$$
\begin{aligned}
\int_{\mathfrak{M}_{\mathcal{Q}}} & \ll \sum_{q<\mathcal{Q}} \sum_{(a, q)=1} \int_{|\beta|<\mathcal{Q} / N}\left|\frac{N}{\mathcal{Q}} \beta\right|^{2}\left(\left|\Omega_{N}\right|\left(\frac{1}{(N|\beta| \mathcal{Q})^{1-c}}\right)\right)^{2} d \beta \\
& \ll \frac{\left|\Omega_{N}\right|^{2}}{N} \frac{1}{\mathcal{Q}^{1-4 c}} .
\end{aligned}
$$

Here $0<c<(1-\delta) 52 / 5<1 / 4$ by (7.2), so renaming the constant $c>0$, we are done with the major arcs.

Decompose the minor arcs $\mathfrak{m}$ into dyadic regions

$$
\int_{\mathfrak{m}}\left|S_{N}(\theta)\right|^{2} d \theta \ll \sum_{\substack{Q<N^{1 / 2} \\ \text { dyadic }}} \sum_{\substack{K<\frac{N^{1 / 2}}{Q} \\ \text { dyadic }}} \mathcal{I}_{Q, K},
$$

where at least one of $Q$ or $K$ exceeds $\mathcal{Q}$, and

$$
\mathcal{I}(Q, K):=\int_{W_{Q, K}}\left|S_{N}(\theta)\right|^{2} d \theta
$$

Write $Q=N^{\alpha}, K=N^{\kappa}$, with the parameters ( $\alpha, \kappa$ ) ranging in

$$
\begin{equation*}
0 \leq \alpha<1 / 2 \text { and } 0 \leq \kappa<1 / 2-\alpha . \tag{7.4}
\end{equation*}
$$

It will be convenient to define

$$
\begin{equation*}
\eta:=(1-\delta) 104 / 5 . \tag{7.5}
\end{equation*}
$$

Assume that $1-\delta<5 / 208$, so that $0<\eta<1 / 2$. We break the summation into the following four ranges:

$$
\begin{aligned}
& \mathcal{R}_{1}:=\{(\alpha, \kappa): \kappa>2(1-\delta)+4 \mathfrak{b}\}, \\
& \mathcal{R}_{2}:=\{(\alpha, \kappa): \alpha>4(1-\delta)+8 \mathfrak{b}\}, \\
& \mathcal{R}_{3}:=\{(\alpha, \kappa): \eta(\alpha+\kappa)<\kappa \text { and } \alpha+\kappa<5 / 52\}, \\
& \mathcal{R}_{4}:=\left\{(\alpha, \kappa): \eta(\alpha+\kappa)<\frac{1}{2} \alpha \text { and } \alpha+\kappa<5 / 52\right\} .
\end{aligned}
$$

We need to show that these four regions cover the entire range (7.4). Using (3.3) and (7.2) with (1.21) guarantees that the regions $\mathcal{R}_{1}$ and $\mathcal{R}_{2}$ certainly cover the range $\alpha+\kappa \geq 5 / 52$. In the complimentary range, $\mathcal{R}_{3}$ and $\mathcal{R}_{4}$ give two regions: the region below the line through the origin with slope $\eta /(1-\eta)$, and the region above the line through the origin with slope $(1 / 2-\eta) / \eta$. These two regions overlap when the slopes overlap, that is, when $\eta<1 / 3$. Then (7.5) explains the value of $\delta_{0}$ in (1.21).

Since the four regions cover the full range (7.4), we now just collect the results of the previous two sections. In the range $\mathcal{R}_{1}$, we apply Theorem 5.11, getting

$$
\begin{equation*}
\mathcal{I}(Q, K) \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N} K^{-c} \tag{7.6}
\end{equation*}
$$

for some $c>0$. In $\mathcal{R}_{2}$, we apply Theorem 5.30, getting

$$
\begin{equation*}
\mathcal{I}(Q, K) \ll \frac{\left(\# \Omega_{N}\right)^{2}}{N} Q^{-c} \tag{7.7}
\end{equation*}
$$

In the range $\mathcal{R}_{3}$ with $K>\mathcal{Q}$, Theorem 6.19 gives (7.6), and in $\mathcal{R}_{4}$ with $Q>\mathcal{Q}$, Theorem 6.24 gives (7.7). Combining these estimates completes the proof of (7.3).

It is now standard to derive Theorem 1.8 from (4.29) and (7.3).

Proof of Theorem 1.8, assuming Proposition 3.9. In light of (4.29), the proof of which uses Proposition 3.9, we have for $n \asymp N$ that

$$
\begin{aligned}
\mathcal{M}_{N}(n) & \gg\left|\Omega_{N}\right| /(N \log \log N) \\
& \gg N^{2 \delta-1-1 / 1000},
\end{aligned}
$$

where we crudely used (3.47) and (3.3). Hence we expect the same for $R_{N}(n)$. If this is not the case, it means that

$$
\left|\mathcal{E}_{N}(n)\right|=\left|R_{N}(n)-\mathcal{M}_{N}(n)\right| \gg \frac{1}{\log \log N} \frac{\left|\Omega_{N}\right|}{N} .
$$

Let $\mathfrak{E}(N)$ denote the set of $n \asymp N$ which have a small representation number $R_{N}(n)$ :

$$
\mathfrak{E}(N):=\left\{\frac{1}{20} N \leq n<\frac{1}{10} N: R_{N}(n)<\frac{1}{2} \mathcal{M}_{N}(n)\right\} .
$$

Then assuming (7.2), we have

$$
\begin{aligned}
\# \mathfrak{E}(N) & \ll \sum_{\frac{1}{20} N \leq n<\frac{1}{10} N} \mathbf{1}_{\left\{\left|\mathcal{E}_{N}(n)\right| \gg \frac{\left|\Omega_{N}\right|}{N \log \log N}\right\}} \\
& \ll \frac{N^{2}(\log \log N)^{2}}{\left|\Omega_{N}\right|^{2}} \sum_{n}\left|\mathcal{E}_{N}(n)\right|^{2} \\
& \ll \frac{N^{2}(\log \log N)^{2}}{\left|\Omega_{N}\right|^{2}} \frac{\left|\Omega_{N}\right|^{2}}{N} \mathcal{Q}^{-c} \ll N e^{-c \sqrt{\log N}},
\end{aligned}
$$

using (7.3) and (3.5).
This completes the proof of Theorem 1.8, modulo the construction of the leading set $\aleph$, which is taken up in the next section. First we give a quick
7.1. Proof of Theorem 1.25. Let $\mathcal{P}=\mathcal{P}_{N}$ be the set of primes $p$ up to $N$ which are $3(\bmod 4)$, so that $(p-1) / 2$ is a 10 -almost-prime; that is,

$$
\mathcal{P}:=\left\{p<N: p \equiv 3(\bmod 4), \text { and } m \mid(p-1) / 2 \Longrightarrow m>N^{1 / 10}\right\} .
$$

A standard sieve argument shows that $\mathcal{P}$ has cardinality $\gg \frac{N}{(\log N)^{2}}$. By (1.9), the cardinality $N e^{-c \sqrt{\log N}}$ of the exceptional set is much smaller, and so $\mathfrak{D}_{\mathcal{A}}(N) \cap \mathcal{P}$ is unbounded in $N$ for $\delta_{\mathcal{A}}>\delta_{0}$.

By (1.10), each $p=d$ in the intersection appears with multiplicity at least $N^{2 \delta-1.001}>N^{10 / 11}$, say. That is, there are distinct $b_{1}, \ldots, b_{L}$ so that $b_{j} / d \in$ $\mathfrak{R}_{\mathcal{A}}, j=1, \ldots, L$, and $L>N^{10 / 11}$. Let $r$ be any primitive root $\bmod d$. For $j=1, \ldots, L$, let $k_{j}$ be defined by $b_{j} \equiv r^{k_{j}}(\bmod d)$, and let $K=\left\{k_{1}, \ldots, k_{L}\right\}$. Of course $b_{j}$ is a primitive root $\bmod d$ if and only if $\left(k_{j}, d-1\right)=1$.

Consider the subset $K^{\prime}$ of $k \in K$ for which $(k, d-1)>2$. Since $d \in \mathcal{P}$, each such $k$ has a prime factor of size $N^{1 / 10}$, and hence the cardinality of $K^{\prime}$
is $\ll N^{9 / 10}$. This is less than the cardinality of $K$, so we may safely discard $K^{\prime}$ from $K$, leaving a nonempty set $K^{\prime \prime}$.

Consider $b \equiv r^{k}(\bmod d)$ with $k \in K^{\prime \prime}$. If $(k, d-1)=1$, we are done, since $b$ is a root $\bmod d$ and $b / d \in \mathfrak{R}_{\mathcal{A}}$. The only other possibility is $(k, d-1)=2$, whence $b$ is a square $\bmod d$. Set $b^{\prime}:=d-b$, so $b^{\prime} \equiv-r^{k}(\bmod d)$; since $d \equiv 3(\bmod 4), b^{\prime}$ is now a primitive root. It is elementary to verify that $b / d \in \mathfrak{R}_{\{1,2, \ldots, A\}}$ implies that $b^{\prime} / d=1-b / d \in \mathfrak{R}_{\{1,2, \ldots, A+1\}}$. That is, these quotients are still absolutely Diophantine, completing the proof.

## 8. Construction of $\aleph$

In this section, we arrange the special leading set $\aleph$ in the ensemble $\Omega_{N}$ as described in Section 3.1. We need two pieces of background, using Section 8.1 to extract some modular/archimedean counting statements from [BGS11], and spending Section 8.2 proving a certain "randomness extraction argument." Finally, we proceed in Section 8.3 to construct «, thereby proving Proposition 3.9 and finalizing the proof of Theorem 1.8.
8.1. Congruence counting theorems. Recall from Section 3.2 that $\mu$ is the $\delta$-dimensional Hausdorff measure supported on the limit set $\mathfrak{C}$, lifted to $\mathbb{P}^{1}$. Extending the work of Lalley's [Lal89] to the congruence setting, Bourgain-Gamburd-Sarnak [BGS11, Th. 1.5] proved the following theorem, adapted to our present context.

Theorem 8.1 ([BGS11]). Let $\Gamma=\Gamma_{\mathcal{A}}$ be the semigroup in (1.31). There exist an integer

$$
\begin{equation*}
\mathfrak{B}=\mathfrak{B}(\mathcal{A}) \geq 1 \tag{8.2}
\end{equation*}
$$

and a constant

$$
\begin{equation*}
\mathfrak{c}=\mathfrak{c}(\mathcal{A})>0 \tag{8.3}
\end{equation*}
$$

so that the following holds. For any $(q, \mathfrak{B})=1$, any $\omega \in \operatorname{SL}_{2}(q)$, any $\gamma_{0} \in \Gamma$, and parameters $T, H \rightarrow \infty$ with $H<e^{c \sqrt{\log T}}$, there is a constant $C\left(\gamma_{0}\right)>0$ so that

$$
\begin{align*}
& \#\left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod q),\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H}, \text { and } \frac{\left\|\gamma \gamma_{0}\right\|}{\left\|\gamma_{0}\right\|} \leq T\right\}  \tag{8.4}\\
& =C\left(\gamma_{0}\right) \cdot T^{2 \delta} \frac{\mu(\mathcal{I})}{\left|\mathrm{SL}_{2}(q)\right|}+O\left(T^{2 \delta} e^{-\mathfrak{c} \sqrt{\log T}}\right), \quad \text { as } T \rightarrow \infty
\end{align*}
$$

Here $\mathcal{I}$ is the interval of length $1 / H$ about $\mathfrak{v}$, and the implied constant does not depend on $T, H, q, \omega$, or $\gamma_{0}$.

With the same conditions as above, except for a modulus $q$ with $\mathfrak{B} \mid q$, we have

$$
\begin{align*}
\# & \left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod q), v_{+}(\gamma) \in \mathcal{I}, \text { and } \frac{\left\|\gamma \gamma_{0}\right\|}{\left\|\gamma_{0}\right\|} \leq T\right\}  \tag{8.5}\\
= & \frac{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|}{\left|\mathrm{SL}_{2}(q)\right|} \cdot \#\left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod \mathfrak{B}), v_{+}(\gamma) \in \mathcal{I}, \text { and } \frac{\left\|\gamma \gamma_{0}\right\|}{\left\|\gamma_{0}\right\|} \leq T\right\} \\
& +O\left(T^{2 \delta} e^{-c \sqrt{\log T}}\right) .
\end{align*}
$$

Remark 8.6. Recall that throughout, the appearance of constants $c$ and $C$ may change from line to line. The special constant $\mathfrak{c}$ in (8.3) is in contradistinction with this principle, being the same constant wherever it appears. Moreover, Proposition 3.16 follows immediately from (8.4), so the constant $\mathfrak{c}$ appearing there can be taken to be the same as the one here.

Remark 8.7. In Theorem 8.1, we have stated the result only for the extreme cases $(\mathfrak{B}, q)=1$ and $\mathfrak{B} \mid q$. Of course intermediate cases can be obtained by summing over suitable arithmetic progressions. This introduces no extra error since the number of terms is $<_{\mathfrak{B}} 1$, and our implied constants may depend on the fixed alphabet $\mathcal{A}$ (recall $\mathfrak{B}$ depends only on $\mathcal{A}$ ).

Remark 8.8. The proof of Theorem 8.1 is the same as that of [BGS11, Th. 1.5], with the following caveats. The latter was proved under the further assumptions that the modulus $q$ is square-free and that $\Gamma$ is a convexcocompact subgroup of $\mathrm{SL}_{2}(\mathbb{Z})$. As discussed in Section 3.2, the proof also works when the group is replaced by our free semigroup $\Gamma$ with no parabolic elements. As for the level, taking $q$ square-free was enough for the sieving purposes in [BGS11], but for the circle method used here, we must take arbitrary $q$. The main ingredient in analyzing the modular aspect (see [BGS11, Lemma 7.2]) was the spectral gap (expansion property) proved using methods of additive combinatorics and an $L^{2}$-flattening lemma in [BG08], [BGS10], again for square-free $q$. The relevant results have since been established for arbitrary modulus; see [BV12] and [GV12, Rem. 30]. With this input, [BGS11, Prop. 12.1] holds for arbitrary $q$, from which one can derive Theorem 8.1.

Remark 8.9. A final caveat is that there is a minor inaccuracy in the statement of [BGS11, Th. 1.5], in that one must "unsmoothe" the smoothed estimate given in [BGS11, Prop. 12.1] to obtain [BGS11, Th. 1.5]. Applying a Tauberian theorem, the removal of smooth weights results in a loss in the stated error term where the smooth error term $T^{-c / \log \log T}$ must be replaced in the sharp cutoff by $e^{-c \sqrt{\log T}}$; hence the error terms stated in Theorem 8.1.

The condition $\left\|\gamma \gamma_{0}\right\| /\left\|\gamma_{0}\right\|<T$ arises naturally in the above through the renewal method; in fact, this condition is essentially equivalent to

$$
d_{\mathbb{H}}\left(\gamma \gamma_{0} i, i\right)-d_{\mathbb{H}}\left(\gamma_{0} i, i\right)<C \log T,
$$

where $d_{\mathbb{H}}(\cdot, \cdot)$ denotes hyperbolic distance in the upper half-plane $\mathbb{H}$. As in Section 3.2 , one typically sets $\gamma_{0}=I$, but we will use $\gamma_{0}$ for a different purpose. Namely, we will need to control both the expanding direction $v_{+}(\gamma)$, and its expanding eigenvalue $\lambda(\gamma)$, but taking $\gamma_{0}=I$ gives us control only on the norm $\|\gamma\|$ (which can be off by a constant from the eigenvalue). So we instead do the following.

It is easy to see from (4.10) that

$$
\begin{equation*}
\|\gamma\|=\frac{\lambda(\gamma)}{\left|\left\langle v_{+}(\gamma), v_{-}^{\perp}(\gamma)\right\rangle\right|}\left(1+O\left(\frac{1}{\|\gamma\|^{2}}\right)\right) . \tag{8.10}
\end{equation*}
$$

Assume now that both $\gamma$ and $\gamma_{0}$ lie in $\mathcal{I}$, the interval of length $1 / H$ about $\mathfrak{v}$, and assume $\|\gamma\|,\left\|\gamma_{0}\right\|>H$. Then applying (8.10) to $\gamma_{0}$ and $\gamma \gamma_{0}$ gives

$$
\left\|\gamma_{0}\right\|=\frac{\lambda\left(\gamma_{0}\right)}{\left|\left\langle\mathfrak{v}, v_{-}^{\perp}\left(\gamma_{0}\right)\right\rangle\right|}\left(1+O\left(\frac{1}{H}\right)\right)
$$

and

$$
\begin{aligned}
\left\|\gamma \gamma_{0}\right\| & =\frac{\lambda\left(\gamma \gamma_{0}\right)}{\left|\left\langle v_{+}\left(\gamma \gamma_{0}\right), v_{-}^{\perp}\left(\gamma \gamma_{0}\right)\right\rangle\right|}\left(1+O\left(\frac{1}{\left\|\gamma \gamma_{0}\right\|^{2}}\right)\right) \\
& =\frac{\lambda(\gamma) \lambda\left(\gamma_{0}\right)}{\left|\left\langle\mathfrak{v}, v_{-}^{\perp}\left(\gamma_{0}\right)\right\rangle\right|}\left(1+O\left(\frac{1}{H}\right)\right),
\end{aligned}
$$

where we used (2.6) and (2.7). On dividing, we obtain

$$
\frac{\left\|\gamma \gamma_{0}\right\|}{\left\|\gamma_{0}\right\|}=\lambda(\gamma)\left(1+O\left(\frac{1}{H}\right)\right) .
$$

We can thus convert statements restricting norms into ones controlling eigenvalues, without losing constants.

The constant $C\left(\gamma_{0}\right)$ in (8.4) approaches a constant $C(\mathfrak{v})>0$ as $\left\|\gamma_{0}\right\| \rightarrow \infty$ with $v_{+}\left(\gamma_{0}\right) \rightarrow \mathfrak{v}$; indeed, $C\left(\gamma_{0}\right)$ is obtained by evaluating a certain Gibbs measure. (See [BGS11, §10] or [Lal89, (2.5)], where his $x$ plays the role of our $\gamma_{0}$.) Hence (8.4) can be replaced by

$$
\begin{aligned}
& \#\left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod q),\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H}, \text { and } \lambda(\gamma) \leq T\right\} \\
& =C(\mathfrak{v}) \cdot T^{2 \delta} \frac{\mu(\mathcal{I})}{\left|\mathrm{SL}_{2}(q)\right|}\left(1+O\left(\frac{1}{H}\right)\right)+O\left(T^{2 \delta} e^{-\mathfrak{c} \sqrt{\log T}}\right)
\end{aligned}
$$

and a similar expression analogous to (8.5).

Finally, we can restrict $\lambda(\gamma)$ to a smaller range, $\lambda(\gamma)=T\left(1+O\left(1 / H_{1}\right)\right)$, for a smaller parameter $H_{1}$; in applications we take $H_{1}=H^{1 / 2}$. In summary, we have the following

Corollary 8.11. With notation as above, we have for any $T, H, H_{1} \rightarrow \infty$ with $H_{1}=o(H)$ and $H<e^{c \sqrt{\log T}}$, and any $(q, \mathfrak{B})=1, \omega \in \mathrm{SL}_{2}(q)$, that

$$
\begin{align*}
& \#\left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod q),\left|v_{+}(\gamma)-\mathfrak{v}\right|<\frac{1}{H},|\lambda(\gamma)-T|<\frac{T}{H_{1}}\right\}  \tag{8.12}\\
& =C(\mathfrak{v}) \cdot \frac{T^{2 \delta}}{H_{1}} \frac{\mu(\mathcal{I})}{\left|\operatorname{SL}_{2}(q)\right|}\left(1+O\left(\frac{1}{H_{1}}+\frac{H_{1}}{H}\right)\right)+O\left(T^{2 \delta} e^{-\mathfrak{v} \sqrt{\log T}}\right) .
\end{align*}
$$

The implied constants are independent of $T, H, H_{1}, q$, and $\omega$.
For a modulus $q \equiv 0(\bmod \mathfrak{B})$, we have

$$
\begin{align*}
\#\{ & \left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod q), v_{+}(\gamma) \in \mathcal{I},|\lambda(\gamma)-T|<\frac{T}{H^{\prime}}\right\}  \tag{8.13}\\
= & \frac{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|}{\left|\mathrm{SL}_{2}(q)\right|} \cdot \#\left\{\gamma \in \Gamma: \gamma \equiv \omega(\bmod \mathfrak{B}), v_{+}(\gamma) \in \mathcal{I},|\lambda(\gamma)-T|<\frac{T}{H^{\prime}}\right\} \\
& \times\left(1+O\left(\frac{1}{H_{1}}+\frac{H_{1}}{H}\right)\right)+O\left(T^{2 \delta} e^{-\mathfrak{c} \sqrt{\log T}}\right) .
\end{align*}
$$

8.2. A randomness extraction argument. Corollary 8.11 gives us good modular/achimedean control away from the modulus $\mathfrak{B}$, but we need $\aleph$ to have good distribution properties for all moduli. So in the next subsection we will concoct certain special sets engineered to have good equidistribution mod $\mathfrak{B}$. But in so doing, we will potentially ruin the distribution away from $\mathfrak{B}$. To recover this distribution, we apply a certain more-or-less standard "randomness extraction" argument, which states roughly that if a large set has good modular distribution, then so does a sufficiently large random subset of it. We will need to have the flexibility to stay away from a modulus $q_{0}$, which in applications is either 1 or $\mathfrak{B}$.

Lemma 8.14. Let $\mu=\mu_{S}$ be the normalized (probability) measure of a finite subset $S \subset \mathrm{SL}(2, \mathbb{Z})$,

$$
\mu(\gamma)=\frac{1}{|S|} \sum_{s \in S} \mathbf{1}_{\{s=\gamma\}}
$$

and fix $\eta>0$. Let $q_{0}<Q$ be a fixed modulus, let $\omega_{0} \in \mathrm{SL}_{2}\left(q_{0}\right)$ be a fixed element, and let $\mathfrak{Q}=\mathfrak{Q}_{q_{0}} \subset[1, Q]$ be the set of moduli $q<Q$ with $q_{0} \mid q$. Assume that for all $q \in \mathfrak{Q}$ and all $\omega \in \operatorname{SL}_{2}(q)$ with $\omega \equiv \omega_{0}\left(\bmod q_{0}\right)$, the projection

$$
\pi_{q}[\mu](\omega)=\sum_{\gamma \equiv \omega(\bmod q)} \mu(\gamma)
$$

is near the uniform measure on $\mathrm{SL}_{2}(q)$ conditioned on being $\equiv \omega_{0}\left(\bmod q_{0}\right)$, (8.15)

$$
\left|\left|\pi_{q}[\mu]-\frac{\left|\mathrm{SL}_{2}\left(q_{0}\right)\right|}{\left|\mathrm{SL}_{2}(q)\right|}\right|_{L^{\infty}}\right|_{\equiv \omega_{0}\left(\bmod q_{0}\right)}=\max _{\substack{\omega \in \operatorname{SL}, 2, q) \\ \omega \equiv \omega_{0}\left(\bmod q_{0}\right)}}\left|\pi_{q}[\mu](\omega)-\frac{\left|\mathrm{SL}_{2}\left(q_{0}\right)\right|}{\left|\mathrm{SL}_{2}(q)\right|}\right|<\eta .
$$

Then for any $T$ with

$$
\begin{equation*}
\eta^{-2} \log Q<T=o\left(|S|^{1 / 2}\right) \tag{8.16}
\end{equation*}
$$

there exist $T$ distinct points $\gamma_{1}, \ldots, \gamma_{T} \in S=\operatorname{supp} \mu$ such that the probability measure $\nu=\nu_{T, \gamma_{1}, \ldots, \gamma_{T}}$ defined by

$$
\begin{equation*}
\nu=\frac{1}{T}\left(\mathbf{1}_{\gamma_{1}}+\cdots+\mathbf{1}_{\gamma_{T}}\right) \tag{8.17}
\end{equation*}
$$

has the same property. That is, for all $q \in \mathfrak{Q}$, projection $\pi_{q}[\nu]$ is also nearly uniform:

$$
\begin{equation*}
\max _{q \in \mathfrak{Z}}\left(\left.\left\|\pi_{q}[\nu]-\frac{\left|\mathrm{SL}_{2}\left(q_{0}\right)\right|}{\left|\mathrm{SL}_{2}(q)\right|}\right\|_{L^{\infty}}\right|_{\equiv \omega_{0}\left(\bmod q_{0}\right)}\right) \ll \eta . \tag{8.18}
\end{equation*}
$$

The implied constant above is absolute.
Proof. This is a standard argument, so we give a sketch. Take $\nu$ as in (8.17). Let $\mathcal{D}$ be the expectation with respect to $\mu$ of the left-hand side of (8.18),

$$
\mathcal{D}:=\sum_{\gamma \in \mathrm{SL}_{2}(\mathbb{Z})^{T}} \max _{q \in \mathfrak{Q}} \max _{\substack{\omega \in \mathrm{SL}_{2}(q) \\ \omega \equiv \omega_{0}\left(\bmod q_{0}\right)}}\left|\frac{1}{T} \sum_{j=1}^{T} \mathbf{1}_{\left\{\gamma_{j} \equiv \omega(q)\right\}}-\frac{\left|\mathrm{SL}_{2}\left(q_{0}\right)\right|}{\left|\mathrm{SL}_{2}(q)\right|}\right| \mu^{(T)}(\gamma),
$$

where $\mu^{(T)}$ is the product measure on $\mathrm{SL}_{2}(\mathbb{Z})^{T}$ and $\gamma=\left(\gamma_{1}, \ldots, \gamma_{T}\right)$. Using (8.15), we have

$$
\mathcal{D}<\eta+\sum_{\gamma \in \mathrm{SL}_{2}(\mathbb{Z})^{T}} \sum_{\xi \in \mathrm{SL}_{2}(\mathbb{Z})^{T}} \max _{q \in \mathfrak{Q}} \max _{\omega \in \mathrm{SL}_{2}(q)}\left|\frac{1}{T} \sum_{j=1}^{T} f_{\omega}\left(\gamma_{j}, \xi_{j}\right)\right| \mu^{(T)}(\gamma) \mu^{(T)}(\xi)
$$

where

$$
f_{\omega}\left(\gamma_{j}, \xi_{j}\right):=\mathbf{1}_{\left\{\gamma_{j} \equiv \omega(q)\right\}}-\mathbf{1}_{\left\{\xi_{j} \equiv \omega(q)\right\}},
$$

and we extended the max over $\omega$ to all of $\mathrm{SL}_{2}(q)$.
Note that for fixed $\omega, f_{\omega}\left(\gamma_{j}, \xi_{j}\right)$ are independent, mean zero random variables and bounded by 1 . Hence the contraction principle gives

$$
\begin{equation*}
\mathcal{D}<\eta+\sum_{\gamma} \sum_{\xi} \mathcal{D}(\gamma, \xi) \mu^{(T)}(\gamma) \mu^{(T)}(\xi), \tag{8.19}
\end{equation*}
$$

where

$$
\mathcal{D}(\gamma, \xi):=\frac{1}{2^{T}} \sum_{\varepsilon \in\{ \pm 1\}^{T}} \max _{q \in \mathfrak{\Omega}} \max _{\omega \in \mathrm{SL}_{2}(q)}\left|\frac{1}{T} \sum_{j=1}^{T} \varepsilon_{j} f_{\omega}\left(\gamma_{j}, \xi_{j}\right)\right| .
$$

Replace the max by an $L^{p}$ norm with $p$ to be chosen later:

$$
\begin{aligned}
\mathcal{D}(\gamma, \xi) & \leq \frac{1}{2^{T}} \sum_{\varepsilon \in\{ \pm 1\}^{T}}\left(\sum_{q \in \mathfrak{Q}} \sum_{\omega \in \operatorname{SL}_{2}(q)}\left|\frac{1}{T} \sum_{j=1}^{T} \varepsilon_{j} f_{\omega}\left(\gamma_{j}, \xi_{j}\right)\right|^{p}\right)^{1 / p} \\
& \ll\left(\sum_{q \in \mathfrak{Q}} \sum_{\omega \in \operatorname{SL}_{2}(q)} \frac{1}{2^{T}} \sum_{\varepsilon \in\{ \pm 1\}^{T}}\left|\frac{1}{T} \sum_{j=1}^{T} \varepsilon_{j} f_{\omega}\left(\gamma_{j}, \xi_{j}\right)\right|^{p}\right)^{1 / p} \\
& \ll\left(\sum_{q \in \mathfrak{Q}} \sum_{\omega \in \operatorname{SL}_{2}(q)} p^{p / 2}\left(\sum_{j=1}^{T}\left|\frac{f_{\omega}\left(\gamma_{j}, \xi_{j}\right)}{T}\right|^{2}\right)^{p / 2}\right)^{1 / p} \\
& \ll Q^{4 / p} p^{1 / 2} T^{-1 / 2}
\end{aligned}
$$

where we applied Khintchine's inequality [Haa81]. (The implied constant is absolute.) Now we choose $p=\log Q$, so that

$$
\mathcal{D}(\gamma, \xi) \ll(\log Q)^{1 / 2} T^{-1 / 2} .
$$

Inserting this into (8.19) and setting $T>\eta^{-2} \log Q$ gives

$$
\mathcal{D} \ll \eta .
$$

The number of $T$-tuples $\gamma$ with distinct $\gamma_{j}$ is $|S|!/(|S|-T)$ !, which is asymptotic to $|S|^{T}$ for $T=o\left(|S|^{1 / 2}\right)$. This completes the proof.

Equipped with this randomness extraction argument, we proceed with the
8.3. Proof of Proposition 3.9. Recalling the parameters $\mathfrak{b}$ in (3.3), $\mathfrak{c}$ from (8.3), and $\mathfrak{B}$ from (8.2), we set $R:=\left|\mathrm{SL}_{2}(\mathfrak{B})\right|$ and define $\alpha_{0}$ in (3.5) by

$$
\begin{equation*}
\alpha_{0}:=\frac{\beta \mathfrak{c}}{40 R} . \tag{8.20}
\end{equation*}
$$

For a parameter

$$
\begin{equation*}
T=N^{c_{1}} \tag{8.21}
\end{equation*}
$$

with small $c_{1}$ to be determined in (8.30), let $H=\mathcal{Q}^{12}, H_{1}=\mathcal{Q}^{6}$, and set

$$
\begin{equation*}
\mathcal{S}(T):=\left\{\gamma \in \Gamma:\left|v_{+}(\gamma)-v\right|<\frac{1}{H},|\lambda(\gamma)-T|<\frac{T}{H_{1}}\right\} . \tag{8.22}
\end{equation*}
$$

By (8.12) with $q=1$, crudely using $\delta<1$ we have that

$$
\begin{equation*}
\# \mathcal{S}(T) \gg T^{2 \delta} / \mathcal{Q}^{18}+O\left(T^{2 \delta} e^{-\mathfrak{c} \sqrt{\log T}}\right) \tag{8.23}
\end{equation*}
$$

We have from (8.21), (8.20), and (3.5) that

$$
e^{-\mathfrak{c} \sqrt{\log T}}=e^{-\mathfrak{c} \sqrt{c_{1}} \sqrt{\log N}}=\mathcal{Q}^{-\sqrt{c_{1}} 40 R / \beta}
$$

So as long as

$$
\begin{equation*}
c_{1}>\left(\frac{3 \mathfrak{b}}{4 R}\right)^{2} \tag{8.24}
\end{equation*}
$$

the estimate (8.23) is significant, with an error of size $\ll T^{2 \delta} / \mathcal{Q}^{30}$.
By the pigeonhole principle, there is some element $\mathfrak{s}_{T} \in \mathcal{S}(T)$ so that

$$
\mathcal{S}^{\prime}(T):=\left\{s \in \mathcal{S}(T): s \equiv \mathfrak{s}_{T}(\bmod \mathfrak{B})\right\}
$$

satisfies

$$
\begin{equation*}
\# \mathcal{S}^{\prime}(T) \geq \frac{1}{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|} \# \mathcal{S}(T) \gg T^{2 \delta} / \mathcal{Q}^{18} \tag{8.25}
\end{equation*}
$$

(Recall our implied constants may depend implicitly on $\mathcal{A}$, and $\mathfrak{B}$ depends only on $\mathcal{A}$.)

For this set, the counting statement (8.23) remains significant even with a modular restriction: for any $q<\mathcal{Q}$ with $\mathfrak{B} \mid q$ and any $\omega \in \mathrm{SL}_{2}(q)$ with $\omega \equiv \mathfrak{s}_{T}(\bmod \mathfrak{B})$, applying (8.13) gives

$$
\begin{align*}
& \#\left\{s \in \mathcal{S}^{\prime}(T): s \equiv \omega(\bmod q)\right\}=\#\{s \in \mathcal{S}(T): s \equiv \omega(\bmod q)\}  \tag{8.26}\\
& =\frac{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|}{\left|\mathrm{SL}_{2}(q)\right|} \#\left\{s \in \mathcal{S}(T): s \equiv \omega \equiv \mathfrak{s}_{T}(\bmod \mathfrak{B})\right\}\left(1+O\left(\mathcal{Q}^{-6}\right)\right)+O\left(T^{2 \delta} \mathcal{Q}^{-30}\right) \\
& =\frac{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|}{\left|\mathrm{SL}_{2}(q)\right|} \# \mathcal{S}^{\prime}(T)\left(1+O\left(\mathcal{Q}^{-6}\right)\right)+O\left(T^{2 \delta} \mathcal{Q}^{-30}\right)
\end{align*}
$$

From (8.25), the main term is $\gg T^{2 \delta} / \mathcal{Q}^{21}$, dominating the error.
We just need to play with $\mathcal{S}^{\prime}(T)$ to get good distribution modulo $\mathfrak{B}$. Recall $R=\left|\mathrm{SL}_{2}(\mathfrak{B})\right|$. Then every element of the "coset"

$$
\gamma \in \mathcal{S}^{\prime}(T) \cdot \mathfrak{s}_{T}^{R-1}
$$

satisfies $\gamma \equiv I(\bmod \mathfrak{B})$. Next write $\operatorname{SL}_{2}(\mathfrak{B})=\left\{\gamma_{1}, \gamma_{2}, \ldots, \gamma_{R}\right\}$, and take $x_{1}, \ldots, x_{R} \in \Gamma$ so that

$$
\begin{equation*}
x_{r} \equiv \gamma_{r}(\bmod \mathfrak{B}), \quad r=1, \ldots, R \tag{8.27}
\end{equation*}
$$

(Recall we had assumed in Section 3.1 that $\Gamma(\bmod q)$ is all of $\mathrm{SL}_{2}(q)$ for all $q$, so such $x_{r}$ exist.) Such $x_{r}$ can be found of size $\asymp_{\mathcal{A}} 1$.

Note that any element

$$
\gamma \in \mathcal{S}^{\prime}(T) \cdot \mathfrak{s}_{T}^{R-1} \cdot x_{r}
$$

has $\gamma \equiv \gamma_{r}(\bmod \mathfrak{B})$. Unfortunately, this triple-product does not work since we do not have control on the expanding vector of $x_{r}$. To remedy this, we take a single fixed element $\mathfrak{f}_{0} \in \Gamma$ of size

$$
\begin{equation*}
\lambda\left(\mathfrak{f}_{0}\right) \asymp B^{1 / 100} \tag{8.28}
\end{equation*}
$$

say, with

$$
\begin{equation*}
\left|v_{+}\left(\mathfrak{f}_{0}\right)-\mathfrak{v}\right|<\mathcal{Q}^{-6} \tag{8.29}
\end{equation*}
$$

Then from $(2.7), v_{+}\left(\mathfrak{f}_{0} x_{r}\right)=\mathfrak{v}\left(1+O\left(\mathcal{Q}^{-6}\right)\right)$, and for any $s \in \mathcal{S}^{\prime}(T)$,

$$
v_{+}\left(s \cdot \mathfrak{s}_{T}^{R-1} \cdot \mathfrak{f}_{0} x_{r}\right)=\mathfrak{v}\left(1+O\left(\mathcal{Q}^{-6}\right)\right)
$$

Moreover, from (8.22) and (2.6), we have

$$
\begin{aligned}
\lambda\left(s \cdot \mathfrak{s}_{T}^{R-1} \cdot \mathfrak{f}_{0} x_{r}\right) & =\lambda(s) \lambda\left(\mathfrak{s}_{T}\right)^{R-1} \lambda\left(\mathfrak{f}_{0} x_{r}\right)\left(1+O\left(\mathcal{Q}^{-6}\right)\right) \\
& =T^{R} \lambda\left(\mathfrak{f}_{0} x_{r}\right)\left(1+O\left(\mathcal{Q}^{-6}\right)\right)
\end{aligned}
$$

Now for each $u \in \mathcal{U}, u \asymp B$, and each $r=1, \ldots, R$, take $T=T_{u, r}$ so that

$$
T^{R} \lambda\left(\mathfrak{f}_{0} x_{r}\right)=u
$$

that is, let

$$
\begin{equation*}
T_{u, r}:=\left(\frac{u}{\lambda\left(\mathfrak{f}_{0} x_{r}\right)}\right)^{1 / R} \asymp B^{99 /(100 R)}=N^{99 \mathfrak{b} /(100 R)} \tag{8.30}
\end{equation*}
$$

which, by (3.4), determines $c_{1}$ in (8.21). Note that (8.24) is easily satisfied.
Thus for each $u$ and $r$, we have sets

$$
\mathcal{B}_{u, r}:=\mathcal{S}^{\prime}\left(T_{u, r}\right) \cdot\left(\mathfrak{s}_{T_{u, r}}\right)^{R-1} \cdot \mathfrak{f}_{0} \cdot x_{r} \subset \Gamma
$$

so that for all $\mathfrak{a} \in \mathcal{B}_{u, r}$, the expanding vector is controlled,

$$
\left|v_{+}(\mathfrak{a})-v\right| \ll \mathcal{Q}^{-6}
$$

and the eigenvalue is controlled,

$$
\lambda(\mathfrak{a})=u\left(1+O\left(\mathcal{Q}^{-6}\right)\right)
$$

Since we have saved an extra $\mathcal{Q}$, we can use it to set the implied constant to 1 , getting (3.11) and (3.12).

Note that by (8.26), for all $q<\mathcal{Q}$ with $\mathfrak{B} \mid q$, and all $\omega \in \mathrm{SL}_{2}(q)$ with $\omega \equiv \mathfrak{f}_{0} x_{r}(\bmod \mathfrak{B})$, we have, crudely, that

$$
\begin{equation*}
\#\left\{\mathfrak{a} \in \mathcal{B}_{u, r}: \mathfrak{a} \equiv \omega(\bmod q)\right\}=\frac{\left|\mathrm{SL}_{2}(\mathfrak{B})\right|}{\left|\mathrm{SL}_{2}(q)\right|} \# \mathcal{B}_{u, r}\left(1+O\left(\mathcal{Q}^{-5}\right)\right) \tag{8.31}
\end{equation*}
$$

Recall also from (8.23) that the cardinality of $\mathcal{B}_{u, r}$ is

$$
\begin{equation*}
\gg\left(T_{u, r}\right)^{2 \delta} / \mathcal{Q}^{18} \gg N^{c} \tag{8.32}
\end{equation*}
$$

using (8.30).

Hence for fixed $u$, we may apply the randomness extraction argument in Lemma 8.14 to $\mathcal{B}_{u, r}$, with $\eta=\mathcal{Q}^{-5}$ and $q_{0}=\mathfrak{B}$. This gives sets $\mathcal{B}_{u, r}^{\prime} \subset \mathcal{B}_{u, r}$ of size $\gg N^{c}$, for which (8.31) continues to hold and, moreover, we can force them all to have exactly the same cardinality independently of $r$,

$$
\left|\mathcal{B}_{u, r}^{\prime}\right|=\left|\mathcal{B}_{u, r^{\prime}}^{\prime}\right| .
$$

Set

$$
\tilde{\aleph}_{u}:=\bigsqcup_{r=1}^{R} \mathcal{B}_{u, r}^{\prime}
$$

and note that for $\mathfrak{B} \mid q<\mathcal{Q}$ and $\omega \in \mathrm{SL}_{2}(q)$,

$$
\#\left\{\mathfrak{a} \in \tilde{\aleph}_{u}: \mathfrak{a} \equiv \omega(\bmod q)\right\}=\frac{\left|\operatorname{SL}_{2}(\mathfrak{B})\right|}{\left|\operatorname{SL}_{2}(q)\right|} \# \mathcal{B}_{u, r}^{\prime}\left(1+O\left(\mathcal{Q}^{-5}\right)\right)
$$

where $r$ is the index for which $\omega \equiv \mathfrak{f}_{0} x_{r}(\bmod \mathfrak{B})$. Since $\# \mathcal{B}_{u, r}^{\prime}=\left|\tilde{\aleph}_{u}\right| / R$, and $R=\left|\mathrm{SL}_{2}(\mathfrak{B})\right|$, we have that for each $u, \tilde{\aleph}_{u}$ satisfies

$$
\begin{equation*}
\#\left\{\mathfrak{a} \in \tilde{\aleph}_{u}: \mathfrak{a} \equiv \omega(\bmod q)\right\}=\frac{\left|\tilde{\aleph}_{u}\right|}{\left|\operatorname{SL}_{2}(q)\right|}\left(1+O\left(\mathcal{Q}^{-5}\right)\right) \tag{8.33}
\end{equation*}
$$

We can now also drop the condition $\mathfrak{B} \mid q$ in (8.33) by summing along certain arithmetic progressions; since $\mathfrak{B} \ll_{A} 1$, the implied constant still depends only on $A$ (cf. Remark 8.7).

Now we apply Lemma 8.14 again to $\tilde{\aleph}_{u}$, with $\eta=\mathcal{Q}^{-5}$ and $q_{0}=1$, giving sets

$$
\aleph_{u} \subset \tilde{\aleph}_{u}
$$

for which (8.33) still holds, that is (3.14) holds, and which all have the same cardinality, giving (3.10).

This completes the proof of Proposition 3.9.

## 9. Proof of Theorem 1.22

Recall that $\mathfrak{R}_{\mathcal{A}}(N)$ is the set of rationals $b / d$ with partial quotients bounded in the alphabet $\mathcal{A}$ with $d<N, \mathfrak{D}_{\mathcal{A}}(N)$ is the set of continuants up to $N$, and $\mathfrak{C}_{\mathcal{A}}$ is the limit set of $\mathfrak{R}_{\mathcal{A}}$ with Hausdorff dimension $\delta=\delta_{\mathcal{A}}$. Recall the sum-set structure (1.47), that if $a \in \mathcal{A}$ and $b / d \in \mathfrak{R}_{\mathcal{A}}$, then $d$ and $b+a d$ are in $\mathfrak{D}_{\mathcal{A}}$. The same holds for another $a^{\prime} \in \mathcal{A}$; that is, all three of $d, b+a d$, and $b+a^{\prime} d$ are in $\mathfrak{D}_{\mathcal{A}}$.

We wish to show (1.24) that for any $\varepsilon>0$,

$$
\#\left(\mathfrak{D}_{\mathcal{A}} \cap[1, N]\right) \gg_{\varepsilon} N^{\delta+(2 \delta-1)(1-\delta) /(5-\delta)-\varepsilon} .
$$

For ease of notation, we lose no generality by specializing from now on to the case $1,2 \in \mathcal{A}$; whence $b / d \in \mathfrak{R}_{\mathcal{A}}(N)$ implies that

$$
\begin{equation*}
d, b+d, b+2 d \in \mathfrak{D}_{\mathcal{A}}(3 N) \tag{9.1}
\end{equation*}
$$

And again, we can drop the subscript $\mathcal{A}$ from $\mathfrak{D}, \mathfrak{R}$, and $\mathfrak{C}$.

Our new ensemble of focus is the collection $\Omega=\Omega_{N}$ of intervals given by

$$
\begin{equation*}
\Omega:=\bigcup_{\substack{d \in \mathcal{O} \\ N / 2 \leq d<3 N}}\left[\frac{d}{N}, \frac{d+1}{N}\right] \subset[1 / 2,3], \tag{9.2}
\end{equation*}
$$

so that

$$
\begin{equation*}
|\Omega| \ll \frac{\#(\mathfrak{D} \cap[1,3 N])}{N} . \tag{9.3}
\end{equation*}
$$

Hensley's conjecture implies, assuming $\delta>1 / 2$, that we should have

$$
\begin{equation*}
|\Omega| \asymp 1 . \tag{9.4}
\end{equation*}
$$

A priori, we do not even know that $|\Omega|>1$, and the bound (1.16) follows from

$$
\begin{equation*}
|\Omega| \gg N^{-(1-\delta)} \tag{9.5}
\end{equation*}
$$

so this is what we must improve upon.
Note from (9.1) that

$$
\begin{equation*}
\mathbf{1}_{\Omega}(x) \mathbf{1}_{\Omega}(y) \mathbf{1}_{\Omega}(x+y) \mathbf{1}_{\Re}\left(\frac{y}{x}-1\right)=1 \tag{9.6}
\end{equation*}
$$

if $x=d / N$ and $y=(b+d) / N$ with $b / d \in \mathfrak{R}$. Just as we thickened $\mathbf{1}_{\mathfrak{D}}$ to some intervals $\mathbf{1}_{\Omega}$ in (9.2), we wish to thicken $\mathbf{1}_{\mathfrak{R}}$ to some intervals.

By Frostman's theorem [Fro35], there is a probability measure $\mu$ supported on the Cantor set $\mathfrak{C}$, so that for any interval $\mathcal{I} \subset[0,1]$, we have

$$
\begin{equation*}
\mu(\mathcal{I})<_{\varepsilon}|\mathcal{I}|^{\delta-\varepsilon} . \tag{9.7}
\end{equation*}
$$

The most naive thickening of $\mathfrak{R}$ we could take is at the scale of $1 / N^{2}$; namely, for each $x \in \mathfrak{C}$, there is (by Dirichlet's approximation theorem and properties of continued fractions) some $b / d \in \mathfrak{R}(N)$ with

$$
\left|x-\frac{b}{d}\right|<\frac{1}{N^{2}} .
$$

So we can find a collection of $\asymp N^{2 \delta}$ intervals of length $\asymp 1 / N^{2}$ which cover $\mathfrak{C}$, each of which has a point in $\mathfrak{R}(N)$. Note also that the spacing between consecutive points in $\mathfrak{R}(N)$ is $\geq 1 / N^{2}$. Instead it will be more fruitful to collect points at square-root this scale, $1 / N$, as follows.

By (9.7), there is a collection $\left\{\mathcal{I}_{\ell}\right\}_{\ell \leq L}$ of $L \asymp N^{\delta}$ disjoint intervals $\mathcal{I}_{\ell} \subset$ $[1 / 2,1]$, each of length $1 / N$, so that

$$
\begin{equation*}
\mu\left(\mathcal{I}_{\ell}\right) \gg_{\varepsilon} N^{-\delta-\varepsilon} . \tag{9.8}
\end{equation*}
$$

Denote their union by

$$
\begin{equation*}
\tilde{\mathcal{R}}:=\bigsqcup_{\ell=1}^{L} \mathcal{I}_{\ell} \tag{9.9}
\end{equation*}
$$

Subdividing each $\mathcal{I}_{\ell}$ into intervals $\mathcal{I}_{\ell, n}$ of length $1 / N^{2}$, it follows that there are at least $\gg N^{\delta-\varepsilon}$ of them intersecting $\mathfrak{C}$. For each such $\ell, n$, the intersection $\mathcal{I}_{\ell, n} \cap \mathfrak{R}(N)$ is also nonempty, possibly after replacing $\mathcal{I}_{\ell, n}$ with a doubling. Hence the cardinality of

$$
\tilde{\mathcal{R}}(N):=\left\{\frac{b}{d} \in \mathfrak{R} \cap \tilde{\mathcal{R}}:(b, d)=1, N / 2<b<d<N\right\}
$$

is of the right order

$$
\begin{equation*}
\# \tilde{\mathcal{R}}(N) \gg_{\varepsilon} N^{2 \delta-\varepsilon} . \tag{9.10}
\end{equation*}
$$

We thicken these intervals just a little further, setting

$$
\begin{equation*}
\mathcal{R}:=\bigcup_{\ell=1}^{L}\left(\mathcal{I}_{\ell}+\left[-\frac{2}{N}, \frac{2}{N}\right]\right) \tag{9.11}
\end{equation*}
$$

Note that since $L \ll N^{\delta}$, we have

$$
\begin{equation*}
|\mathcal{R}| \ll N^{-(1-\delta)} . \tag{9.12}
\end{equation*}
$$

Note further that $\mathfrak{C}$, and hence $\mathcal{R}$, is contained strictly inside $[0,1]$; that is, for some $\nu=\nu_{A}>0$, we have

$$
\begin{equation*}
\mathcal{R} \subset[\nu, 1-\nu] . \tag{9.13}
\end{equation*}
$$

Consider now our main integral $\mathcal{J}$, motivated by (9.6), defined by

$$
\begin{equation*}
\mathcal{J}:=\iint_{\mathbb{R}^{2}} \mathbf{1}_{\Omega}(x) \mathbf{1}_{\Omega}(y) \mathbf{1}_{\Omega}(x+y) \mathbf{1}_{\mathcal{R}}\left(\frac{y}{x}-1\right) d y d x . \tag{9.14}
\end{equation*}
$$

By (9.2), the domain of integration above is supported in the box $[1 / 2,3] \times$ [1/2,3].

For each $b / d \in \tilde{\mathcal{R}}(N)$, the intervals

$$
\frac{d}{N} \leq x \leq \frac{d+1 / 2}{N}, \quad \frac{b+d}{N} \leq y \leq \frac{b+d+1 / 2}{N}
$$

belong to $\Omega$, by ( 9.1 ), as does the interval

$$
\frac{b+2 d}{N} \leq x+y \leq \frac{b+2 d+1}{N}
$$

Moreover, the interval

$$
\frac{b}{d}-\frac{2}{N}<\frac{b-1}{d} \leq \frac{y}{x}-1 \leq \frac{b+1}{d}<\frac{b}{d}+\frac{2}{N}
$$

is in $\mathcal{R}$ by the thickening in (9.11). That is, these intervals contribute $1 / N^{2}$ to $\mathcal{J}$ for each $b / d$, and hence by (9.10), we have established the following lower bound.

Proposition 9.15.

$$
\begin{equation*}
\mathcal{J} \gg \varepsilon N^{-2(1-\delta)-\varepsilon .} \tag{9.16}
\end{equation*}
$$

The trivial bound $\mathcal{J} \ll|\Omega|^{2}$ from (9.14) recovers (9.5), which is what we must improve. The rest of the appendix is devoted to establishing the following

Proposition 9.17.

$$
\begin{equation*}
\mathcal{J} \ll_{\varepsilon} N^{-\frac{1-\delta}{2}+\varepsilon}|\Omega|^{2-\frac{3(1-\delta)}{2(2-\delta)}} . \tag{9.18}
\end{equation*}
$$

Then Theorem 1.23 follows immediately from (9.16), (9.18), and (9.3).
9.1. Proof of Proposition 9.17. Let $M>0$ be a parameter to be chosen later (it will be a little less than $N^{-(1-\delta)}$ ), and decompose

$$
\begin{equation*}
\mathcal{J}=\mathcal{J}_{1}+\mathcal{J}_{2}, \tag{9.19}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{J}_{1}:=\iint_{\left(\mathbf{1}_{\Omega * 1}-\Omega\right)(-x)<M} \cdots d x d y \tag{9.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{J}_{2}:=\iint_{\left(\mathbf{1}_{\left.\Omega * \mathbf{1}_{-\Omega}\right)(-x) \geq M}\right.} \cdots d x d y . \tag{9.21}
\end{equation*}
$$

Then writing $\mathbf{1}_{\mathcal{R}} \leq 1$, we have

$$
\begin{align*}
\mathcal{J}_{1} & \leq \iint_{\left(\mathbf{1}_{\Omega * 1} \mathbf{1}_{-\Omega}\right)(-x)<M} \mathbf{1}_{\Omega}(x) \mathbf{1}_{\Omega}(y) \mathbf{1}_{-\Omega}(-x-y) d x d y  \tag{9.22}\\
& \leq \int_{\substack{\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x)<M}} \mathbf{1}_{\Omega}(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}(-x)\right) d x \\
& <M|\Omega| .
\end{align*}
$$

It is clear already that to get a gain on $|\Omega|$, we must take $M$ a power less than $N^{-(1-\delta)}$.

We are left to analyze $\mathcal{J}_{2}$. Note that in the domain of $\mathcal{J}_{2}$, we have

$$
\mathbf{1}_{\Omega}(x+y) \leq 1 \leq \frac{1}{M}\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x) .
$$

Hence we can write

$$
\mathcal{J}_{2} \leq \frac{1}{M} \iint_{\mathbb{R}^{2}} \eta(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x) \mathbf{1}_{\Omega}(y) \mathbf{1}_{\mathcal{R}}\left(\frac{y}{x}-1\right) d y d x
$$

where we have bounded $\mathbf{1}_{\Omega}(x)$ by a smooth bump function $\eta(x)$ with support in [1/4,4], say, and $\eta \geq 1$ on $[1 / 2,3]$ to recall (9.2) that $x \in[1 / 2,3]$.

For a smooth, nonnegative, even function $\Upsilon$ with compact support and $\int \Upsilon=1$, let $\Upsilon_{N}(y):=10 N \Upsilon(10 N y)$, and dominate $\mathbf{1}_{\Omega}(y)$, up to constant, by the smooth function

$$
\mathcal{S}_{\Omega}:=\mathbf{1}_{\Omega} * \Upsilon_{N} .
$$



Figure 1. The Fourier multiplier $\lambda_{0}$.

So we have

$$
\begin{equation*}
\mathcal{J}_{2} \ll \frac{1}{M} \cdot \mathcal{J}_{2}^{\prime}, \tag{9.23}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{J}_{2}^{\prime}:=\iint_{\mathbb{R}^{2}} \eta(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x) \mathcal{S}_{\Omega}(y) \mathbf{1}_{\mathcal{R}}\left(\frac{y}{x}-1\right) d y d x . \tag{9.24}
\end{equation*}
$$

Note that, by the smoothness of $\Upsilon$, the Fourier spectrum of $\mathcal{S}_{\Omega}$ is contained, up to negligible error, in $\left[-N^{1+\varepsilon}, N^{1+\varepsilon}\right]$. So we can decompose $\mathcal{S}_{\Omega}(y)$ by the technique of "slicing." That is, introduce a certain dyadic partition of unity via the Fourier multipliers $\lambda_{k}(\xi)$, defined as follows. Let $\lambda_{0}(\xi)$ be even, $\equiv 1$ on $[1,2]$, and decaying piecewise-linearly to 0 at $\xi=1 / 2$ and $\xi=4$; see Figure 1. For integers $k$ ranging in

$$
\begin{equation*}
0<2^{2 k}<N^{1+\varepsilon}, \tag{9.25}
\end{equation*}
$$

define

$$
\lambda_{k}(\xi):=\lambda_{0}\left(\xi \cdot 2^{-2 k}\right)
$$

Let $\Lambda_{k}$ be the Fourier inverse of $\lambda_{k}$, so

$$
\Lambda_{k}(y)=2^{2 k} \Lambda_{0}\left(2^{2 k} y\right)=\left(\mathcal{D}_{2^{2 k}} \Lambda_{0}\right)(y)
$$

where

$$
\Lambda_{0}(y):=\frac{\sin ^{2}\left(\frac{3}{2} \pi y\right)}{\pi^{2} y^{2}}(2 \cos (2 \pi y)-\cos (\pi y)+\cos (5 \pi y)),
$$

and $\mathcal{D}_{u}$ is the dilation representation:

$$
\left(\mathcal{D}_{u} f\right)(y):=u f(u y) .
$$

We also introduce $\mathcal{T}_{u}$, the translation representation,

$$
\mathcal{T}_{u} f(y):=f(y+u) .
$$

Of course $\widehat{\Lambda_{k} * f}=\lambda_{k} \cdot \widehat{f}$, so we have

$$
\mathcal{S}_{\Omega}=\sum_{2^{2 k}<N^{1+\varepsilon}}\left(\Lambda_{k} * \mathcal{S}_{\Omega}\right)+\text { Err },
$$

where Err is bounded by an arbitrarily large power of $1 / N$, and will henceforth be ignored.

Then we can bound (9.24) as

$$
\begin{equation*}
\left|\mathcal{J}_{2}^{\prime}\right| \ll \sum_{k}\left|\mathcal{J}_{2}^{(k)}\right|, \tag{9.26}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{J}_{2}^{(k)}:=\iint \eta(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x)\left(\Lambda_{k} * \mathcal{S}_{\Omega}\right)(y) \mathbf{1}_{\mathcal{R}}\left(\frac{y}{x}-1\right) d y d x \tag{9.27}
\end{equation*}
$$

Since the Fourier spectrum of $\mathcal{S}_{\Omega}$ is now controlled, so is that of $\mathbf{1}_{\mathcal{R}}$, as follows. Write $\mathbf{1}_{\mathcal{R}}\left(\frac{y}{x}-1\right)=\mathbf{1}_{x \mathcal{R}}(y-x)=\left(\mathcal{T}_{-x} \mathbf{1}_{x \mathcal{R}}\right)(y)$. Then the $y$ integral can be written as

$$
\begin{align*}
\int_{\mathbb{R}}\left(\Lambda_{k}\right. & \left.* \mathcal{S}_{\Omega}\right)(y)\left(\mathcal{T}_{-x} \mathbf{1}_{x \mathcal{R}}\right)(y) d y  \tag{9.28}\\
& =\int_{\mathbb{R}} \lambda_{k}(\xi) \widehat{\mathcal{S}}_{\Omega}(\xi) \overline{\overline{\left(\mathcal{T}_{-x} \mathbf{1}_{x \mathcal{R}}\right)}(\xi)} d \xi \\
& =\int_{\mathbb{R}} \lambda_{k}(\xi) \widehat{\mathcal{S}}_{\Omega}(\xi) \lambda_{k}^{\prime}(x \xi) \overline{\overline{\left(\mathcal{T}_{-x} \mathbf{1}_{x \mathcal{R}}\right)}(\xi)} d \xi
\end{align*}
$$

where we inserted another bump function $\lambda_{k}^{\prime}$ which is smooth in addition to other properties of $\lambda_{k}$. Namely, let $\lambda_{0}^{\prime}$ be even, $\equiv 1$ on $\pm[1 / 16,16]$, and decay smoothly to 0 outside of $\pm[1 / 32,32]$; then set $\lambda_{k}^{\prime}(\xi):=\lambda_{0}^{\prime}\left(\xi 2^{-2 k}\right)$. The point is that $\lambda_{k}^{\prime}(x \xi) \equiv 1$ on the support of $\lambda_{k}$, since $x \in[1 / 4,4]$ by the support of $\eta$, so the above equality holds.

Then writing $\Lambda_{k}^{\prime}$ for the inverse transform of $\lambda_{k}^{\prime}$, we have

$$
\begin{equation*}
\mathcal{J}_{2}^{(k)}=\iint \eta(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x)\left(\Lambda_{k} * \mathcal{S}_{\Omega}\right)(y)\left[\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right]\left(\frac{y}{x}-1\right) d y d x \tag{9.29}
\end{equation*}
$$

Now we handle two ranges of $k$ separately. We introduce a cutoff parameter $\mathcal{K}$ to be chosen later, see (9.41).
9.1.1. The range $k \leq \mathcal{K}$. We wish to prove

Lemma 9.30. For $k \leq \mathcal{K}$,

$$
\begin{equation*}
\left|\mathcal{J}_{2}^{(k)}\right| \lll \varepsilon|\Omega|^{3} 2^{2 k(1-\delta)} N^{-(1-\delta)+\varepsilon} \tag{9.31}
\end{equation*}
$$

This is a gain of a power of $|\Omega|$. (Recall we are assuming $|\Omega|<1$.)

Proof. Estimate (9.29) by

$$
\begin{equation*}
\left|\mathcal{J}_{2}^{(k)}\right| \ll\left\|\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right\|_{1} \cdot\left\|\Lambda_{k} * \mathcal{S}_{\Omega}\right\|_{1} \cdot\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{\infty} \tag{9.32}
\end{equation*}
$$

The first factor contributes $|\Omega|^{2}$, and the second is $\ll|\Omega|$, since $\Lambda_{0}$ is integrable. For the last term, write

$$
\begin{align*}
\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{\infty} & \ll \sup _{z} \int_{\mathbb{R}} 2^{2 k}\left|\Lambda_{0}^{\prime}\left(u 2^{2 k}\right)\right| \mathbf{1}_{\mathcal{R}}(z-u) d u  \tag{9.33}\\
& \ll \sum_{m \geq 0} 2^{2(k-m)} \sup _{|\mathcal{U}|=2^{m+1-2 k}}|\mathcal{R} \cap \mathcal{U}|
\end{align*}
$$

where the supremum is taken over intervals $\mathcal{U}$. Here we used that $\Lambda_{0}^{\prime}$ has rapid decay, so certainly

$$
\Lambda_{0}^{\prime}(y) \ll \mathbf{1}_{|y|<1}+\frac{1}{2^{2}} \mathbf{1}_{|y|<2}+\frac{1}{4^{2}} \mathbf{1}_{|y|<4}+\frac{1}{8^{2}} \mathbf{1}_{|y|<8}+\cdots
$$

Note by (9.25) that

$$
|\mathcal{U}| \geq 2^{-2 k+1} \geq \frac{2}{N^{1+\varepsilon}}
$$

We have yet to exploit the structure of $\mathcal{R}$ and do so now. This requires the following

Lemma 9.34. For any interval $\mathcal{U}$ of length at least $1 / N^{1+\varepsilon}$, we have

$$
\begin{equation*}
|\mathcal{R} \cap \mathcal{U}| \ll N^{-(1-\delta)+\varepsilon}|\mathcal{U}|^{\delta+\varepsilon} \tag{9.35}
\end{equation*}
$$

Postponing the proof of this lemma, we see that applying (9.35) in (9.33) gives

$$
\begin{equation*}
\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{\infty} \ll_{\varepsilon} 2^{2 k(1-\delta)} N^{-(1-\delta)+\varepsilon} \tag{9.36}
\end{equation*}
$$

Putting (9.36) into (9.32) gives (9.31), as claimed.
It remains to establish (9.35).
Proof of Lemma 9.34. From the structure of $\mathcal{R}$ in (9.11), we have that

$$
\begin{aligned}
|\mathcal{R} \cap \mathcal{U}| & \leq \sum_{\ell}\left|\mathcal{U} \cap\left(\mathcal{I}_{\ell}+[-2 / N, 2 / N]\right)\right| \\
& \ll \frac{1}{N} \#\left\{\ell \leq L: \mathcal{U} \cap\left(\mathcal{I}_{\ell}+[-2 / N, 2 / N]\right) \neq \emptyset\right\} \\
& \ll \varepsilon \frac{1}{N} \frac{\mu(\mathcal{U}+[0,1 / N])}{N^{-\delta-\varepsilon}} \\
& \ll \varepsilon N^{-(1-\delta)+\varepsilon}|\mathcal{U}|^{\delta+\varepsilon}
\end{aligned}
$$

where we used (9.8) and (9.7) in the penultimate and final lines, respectively.
9.1.2. The range $k>\mathcal{K}$. In this range, we will establish

Lemma 9.37. For $k>\mathcal{K}$ and any $\varepsilon>0$,

$$
\begin{equation*}
\left|\mathcal{J}_{2}^{(k)}\right| \lll|\Omega|^{3 / 2} 2^{-k \delta} N^{-(1-\delta)+\varepsilon} . \tag{9.38}
\end{equation*}
$$

Proof. Changing variables $y \mapsto y x$ in (9.29), we have

$$
\left|\mathcal{J}_{2}^{(k)}\right| \ll\left|\iint \mathfrak{f}(x)\left(\Lambda_{k} * \mathcal{S}_{\Omega}\right)(x y) \eta(y)\left(\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right)(y-1) d y d x\right|,
$$

where we set

$$
\begin{equation*}
\mathfrak{f}(x):=x \eta(x)\left(\mathbf{1}_{\Omega} * \mathbf{1}_{-\Omega}\right)(-x) . \tag{9.39}
\end{equation*}
$$

By the rapid decay of $\Lambda_{k}^{\prime}$ and (9.13), we may restrict the integral to $y \asymp 1$ with a negligible error. Now reverse orders, apply Parseval in $x$, reverse orders again, use the definition of the Fourier multipliers $\lambda_{k}$, apply Cauchy-Schwarz in $y$, change variables $y \mapsto \xi / y$, and estimate

$$
\begin{aligned}
\left|\mathcal{J}_{2}^{(k)}\right| & \ll\left|\int_{y \asymp 1}\left(\int_{|\xi| / y \asymp 2^{2 k}} \widehat{\mathfrak{f}}(\xi) \overline{\hat{\mathcal{S}}_{\Omega}(\xi / y)} \frac{1}{y} d \xi\right)\left(\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right)(y-1) d y\right| \\
& \ll \int_{|\xi| \asymp 2^{2 k}}|\widehat{\mathfrak{f}}(\xi)|\left(\frac{1}{|\xi|} \int\left|\widehat{\mathcal{S}}_{\Omega}(y)\right|^{2} d y\right)^{1 / 2}\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{2} d \xi \\
& \ll \varepsilon\left(2^{-k} \int_{\xi \in \mathbb{R}}|\widehat{\mathfrak{f}}(\xi)| d \xi\right)|\Omega|^{1 / 2} 2^{k(1-\delta)} N^{-(1-\delta)+\varepsilon},
\end{aligned}
$$

where we estimated the last piece by

$$
\begin{aligned}
\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{2} & \leq\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{\infty}^{1 / 2}\left\|\Lambda_{k}^{\prime} * \mathbf{1}_{\mathcal{R}}\right\|_{1}^{1 / 2} \\
& <_{\varepsilon} 2^{k(1-\delta)} N^{-(1-\delta) / 2+\varepsilon}|\mathcal{R}|^{1 / 2} \\
& \ll_{\varepsilon} 2^{k(1-\delta)} N^{-(1-\delta)+\varepsilon},
\end{aligned}
$$

using the $\mathcal{L}^{\infty}$ bound in (9.36) and (9.12). We easily estimate from (9.39) that $\|\widehat{\mathfrak{f}}\|_{1} \ll|\Omega|$, giving (9.38), as claimed.
9.1.3. Completion of Proof. It is now a simple matter to establish Proposition 9.17. Putting (9.31), (9.38), and (9.26) into (9.23) gives

$$
\begin{align*}
\mathcal{J}_{2} & \ll \varepsilon \frac{1}{M} N^{-(1-\delta)+\varepsilon}\left(|\Omega|^{3} 2^{2 \mathcal{K}(1-\delta)}+|\Omega|^{3 / 2} 2^{-\mathcal{K} \delta}\right)  \tag{9.40}\\
& \ll \varepsilon \frac{1}{M} N^{-(1-\delta)+\varepsilon}|\Omega|^{3 /(2-\delta)},
\end{align*}
$$

on setting

$$
\begin{equation*}
\mathcal{K}:=\frac{-3 \log _{2}|\Omega|}{2(2-\delta)} . \tag{9.41}
\end{equation*}
$$

Combining (9.40) with (9.22) and choosing

$$
M=N^{-(1-\delta) / 2+\varepsilon}|\Omega|^{(1+\delta) /(4-2 \delta)}
$$

gives (9.18), as claimed. This competes the proof of Theorem 1.23.

## References

[Bou03] J. Bourgain, On the Erdős-Volkmann and Katz-Tao ring conjectures, Geom. Funct. Anal. 13 (2003), 334-365. MR 1982147. Zbl 1115.11049. http://dx.doi.org/10.1007/s000390300008.
[Bou05] J. Bourgain, Green's Function Estimates for Lattice Schrödinger Operators and Applications, Ann. of Math. Stud. 158, Princeton Univ. Press, Princeton, NJ, 2005. MR 2100420. Zbl 1137. 35001.
[Bou10] J. Bourgain, The discretized sum-product and projection theorems, J. Anal. Math. 112 (2010), 193-236. MR 2763000. Zbl 1234.11012. http: //dx.doi.org/10.1007/s11854-010-0028-x.
[Bou12] J. Bourgain, Partial quotients and representation of rational numbers, C. R. Math. Acad. Sci. Paris 350 (2012), 727-730. MR 2981341. Zbl 06100614. http://dx.doi.org/10.1016/j.crma.2012.09.002.
[BG08] J. Bourgain and A. Gamburd, Uniform expansion bounds for Cayley graphs of $\mathrm{SL}_{2}\left(\mathbb{F}_{p}\right)$, Ann. of Math. 167 (2008), 625-642. MR 2415383. Zbl 1216.20042. http://dx.doi.org/10.4007/annals.2008.167.625.
[BGS10] J. Bourgain, A. Gamburd, and P. Sarnak, Affine linear sieve, expanders, and sum-product, Invent. Math. 179 (2010), 559-644. MR 2587341. Zbl 1239.11103. http://dx.doi.org/10.1007/ s00222-009-0225-3.
[BGS11] J. Bourgain, A. Gamburd, and P. Sarnak, Generalization of Selberg's $\frac{3}{16}$ theorem and affine sieve, Acta Math. 207 (2011), 255-290. MR 2892611. Zbl 1276.11081. http://dx.doi.org/10.1007/s11511-012-0070-x.
[BK10] J. Bourgain and A. Kontorovich, On representations of integers in thin subgroups of $\mathrm{SL}_{2}(\mathbb{Z})$, Geom. Funct. Anal. 20 (2010), 1144-1174. MR 2746949. Zbl 1230.11050. http://dx.doi.org/10.1007/ s00039-010-0093-4.
[BK11] J. Bourgain and A. Kontorovich, On Zaremba's conjecture, C. R. Math. Acad. Sci. Paris 349 (2011), 493-495. MR 2802911. Zbl 1215. 11005. http://dx.doi.org/10.1016/j.crma.2011.03.023.
[BKS10] J. Bourgain, A. Kontorovich, and P. Sarnak, Sector estimates for hyperbolic isometries, Geom. Funct. Anal. 20 (2010), 1175-1200. MR 2746950. Zbl 1208.22007. http://dx.doi.org/10.1007/ s00039-010-0092-5.
[BV12] J. Bourgain and P. P. Varjú, Expansion in $S L_{d}(\mathbf{Z} / q \mathbf{Z}), q$ arbitrary, Invent. Math. 188 (2012), 151-173. MR 2897695. Zbl 1247.20052. http: //dx.doi.org/10.1007/s00222-011-0345-4.
[Bum85] R. T. Bumby, Hausdorff dimension of sets arising in number theory, in Number Theory (New York, 1983-84), Lecture Notes in Math. 1135,

Springer-Verlag, New York, 1985, pp. 1-8. MR 0803348. Zbl 0575.28004. http://dx.doi.org/10.1007/BFb0074599.
[Dol98] D. Dolgopyat, On decay of correlations in Anosov flows, Ann. of Math. 147 (1998), 357-390. MR 1626749. Zbl 0911.58029. http://dx.doi.org/10. 2307/121012.
[FK13] D. A. Frolenkov and I. D. Kan, A reinforcement of the BourgainKontorovich's theorem by elementary methods II, 2013. arXiv 1303.3968.
[Fro35] O. Frostman, Potentiel d'equilibre et capacité des ensembles avec quelques applications à la théorie des fonctions, Meddel. Lunds Univ. Math. Sen. 3 (1935), 1-118. JFM 61.1262. 02.
[GS01] M. Goldstein and W. Schlag, Hölder continuity of the integrated density of states for quasi-periodic Schrödinger equations and averages of shifts of subharmonic functions, Ann. of Math. 154 (2001), 155-203. MR 1847592. Zbl 0990.39014. http://dx.doi.org/10.2307/3062114.
[GV12] A. S. Golsefidy and P. P. Varjú, Expansion in perfect groups, Geom. Funct. Anal. 22 (2012), 1832-1891. MR 3000503. Zbl 06132223. http:// dx.doi.org/10.1007/s00039-012-0190-7.
[Goo83] A. Good, Local Analysis of Selberg's Trace Formula, Lecture Notes in Math. 1040, Springer-Verlag, New York, 1983. MR 0727476. Zbl 0525. 10013. http://dx.doi.org/10.1007/BFb0073074.
[Goo41] I. J. Good, The fractional dimensional theory of continued fractions, Proc. Cambridge Philos. Soc. 37 (1941), 199-228. MR 0004878. Zbl 0061.09408. http://dx.doi.org/10.1017/S030500410002171X.
[Haa81] U. HaAgerup, The best constants in the Khintchine inequality, Studia Math. 70 (1981), 231-283 (1982). MR 0654838. Zbl 0501.46015.
[Hen89] D. Hensley, The distribution of badly approximable numbers and continuants with bounded digits, in Théorie des Nombres (Quebec, PQ, 1987), de Gruyter, Berlin, 1989, pp. 371-385. MR 1024576. Zbl 0689. 10042. http://dx.doi.org/10.1515/9783110852790.371.
[Hen92] D. Hensley, Continued fraction Cantor sets, Hausdorff dimension, and functional analysis, J. Number Theory 40 (1992), 336-358. MR 1154044. Zbl 0745.28005. http://dx.doi.org/10.1016/0022-314X(92)90006-B.
[Hen06] D. Hensley, Continued Fractions, World Scientific Publishing Co. Pte. Ltd., Hackensack, NJ, 2006. MR 2351741. Zbl 1161.11028. http://dx.doi. org/10.1142/9789812774682.
[Hen96] D. Hensley, A polynomial time algorithm for the Hausdorff dimension of continued fraction Cantor sets, J. Number Theory 58 (1996), 9-45. MR 1387719. Zbl 0858.11039. http://dx.doi.org/10.1006/jnth.1996.0058.
[Hua13] S. Huang, An improvement on Zaremba's conjecture, 2013. arXiv 1310. 3772.
[Jen04] O. Jenkinson, On the density of Hausdorff dimensions of bounded type continued fraction sets: the Texan conjecture, Stoch. Dyn. 4 (2004), 63-76. MR 2069367. Zbl 1089.28006. http://dx.doi.org/10.1142/ S0219493704000900.
[JP01] O. Jenkinson and M. Pollicott, Computing the dimension of dynamically defined sets: $E_{2}$ and bounded continued fractions, Ergodic Theory Dynam. Systems 21 (2001), 1429-1445. MR 1855840. Zbl 0991.28009. http://dx.doi.org/10.1017/S0143385701001687.
[Kim03] H. H. Kim, Functoriality for the exterior square of $\mathrm{GL}_{4}$ and the symmetric fourth of $\mathrm{GL}_{2}$, J. Amer. Math. Soc. 16 (2003), 139-183. MR 1937203. Zbl 1018.11024. http://dx.doi.org/10.1090/S0894-0347-02-00410-1.
[Kon13] A. Kontorovich, From Apollonius to Zaremba: local-global phenomena in thin orbits, Bull. Amer. Math. Soc. 50 (2013), 187-228. MR 3020826. Zbl 06157470. http://dx.doi.org/10.1090/S0273-0979-2013-01402-2.
[Lal89] S. P. Lalley, Renewal theorems in symbolic dynamics, with applications to geodesic flows, non-Euclidean tessellations and their fractal limits, Acta Math. 163 (1989), 1-55. MR 1007619. Zbl 0701.58021. http://dx.doi.org/ 10.1007/BF02392732.
[MVW84] C. R. Matthews, L. N. Vaserstein, and B. Weisfeiler, Congruence properties of Zariski-dense subgroups. I, Proc. London Math. Soc. 48 (1984), 514-532. MR 0735226. Zbl 0551.20029. http://dx.doi.org/10. 1112/plms/s3-48.3.514.
[McM09] C. T. McMullen, Uniformly Diophantine numbers in a fixed real quadratic field, Compos. Math. 145 (2009), 827-844. MR 2521246. Zbl 1176. 11032. http://dx.doi.org/10.1112/S0010437X09004102.
[Nau05] F. NAUD, Expanding maps on Cantor sets and analytic continuation of zeta functions, Ann. Sci. École Norm. Sup. 38 (2005), 116-153. MR 2136484. Zbl 1110.37021. http://dx.doi.org/10.1016/j.ansens.2004.11.002.
[Nie78] H. Niederreiter, Quasi-Monte Carlo methods and pseudo-random numbers, Bull. Amer. Math. Soc. 84 (1978), 957-1041. MR 0508447. Zbl 0404. 65003. http://dx.doi.org/10.1090/S0002-9904-1978-14532-7.
[Zar72] S. K. Zaremba, La méthode des "bons treillis" pour le calcul des intégrales multiples, in Applications of Number Theory to Numerical Analysis (Proc. Sympos., Univ. Montreal, Montreal, Que., 1971), Academic Press, New York, 1972, pp. 39-119. MR 0343530. Zbl 0246.65009.
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[^0]:    ${ }^{1}$ The program is available at http://math.yale.edu/~avk23/maths/HausdorffZaremba.nb.

[^1]:    ${ }^{2}$ http://arxiv.org/abs/1107.3776v1
    ${ }^{3}$ Added in press: Huang [Hua13] has combined the methods of Frolenkov-Kan with ours to obtain the full density one statement with $\delta_{0}=5 / 6$.

[^2]:    ${ }^{4}$ We could take $J$ even a bit smaller, but choose not to for the sake of exposition.

