# THE BOLTZMANN-GRAD LIMIT OF THE LORENTZ GAS IN A UNION OF LATTICES 

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#### Abstract

The Lorentz gas describes an ensemble of noninteracting point particles in an infinite array of spherical scatterers. In the present paper we consider the case when the scatterer configuration $\mathcal{P}$ is a fixed union of (translated) lattices in $\mathbb{R}^{d}$, and prove that in the limit of low scatterer density, the particle dynamics converges to a random flight process. In the special case when the lattices in $\mathcal{P}$ are pairwise incommensurable, this settles a conjecture from [20. The proof is carried out by applying a framework developed in recent work by Marklof and Strömbergsson 21, and central parts of our proof are the construction of an admissible marking of the point set $\mathcal{P}$, and the verification of the uniform spherical equidistribution condition required in [21. Regarding the random flight process obtained in the low density limit of the Lorentz gas, we prove that it can be reconstructed from the corresponding limiting flight processes arising from the individual commensurability classes of lattices in $\mathcal{P}$. We furthermore prove that the free path lengths of the limit flight process have a distribution with a power law tail, whose exponent depends on the number of commensurability classes in $\mathcal{P}$.


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## 1. Introduction

The Lorentz gas [14] describes the dynamics of a cloud of non-interacting point particles in an array of fixed spherical scatterers of radius $\rho>0$, centered at the elements of a given locally finite point set $\mathcal{P} \subset \mathbb{R}^{d}$. Each particle travels with constant velocity along straight lines, and each time it hits a scatterer it is deflected by elastic reflection or by a more general (fixed) scattering process. We denote the position and velocity of a point particle at time $t$ by $\boldsymbol{q}(t)$ and $\boldsymbol{v}(t)$. Since the particle speed outside the scatterers is a constant of motion we may without loss of generality consider only point particles having unit speed, viz., $\|\boldsymbol{v}(t)\|=1$. This means that the particle dynamics takes place in the unit tangent bundle $\mathrm{T}^{1}\left(\mathcal{K}_{\rho}\right):=\mathcal{K}_{\rho} \times \mathrm{S}_{1}^{d-1}$ of the domain

$$
\mathcal{K}_{\rho}:=\mathbb{R}^{d} \backslash\left(\mathcal{P}+\mathcal{B}_{\rho}^{d}\right),
$$

where $\mathcal{B}_{\rho}^{d}$ denotes the open ball of radius $\rho$, centered at the origin. The Liouville measure on $\mathrm{T}^{1}\left(\mathcal{K}_{\rho}\right)$ is vol $\times \sigma$, where vol denotes the Lebesgue measure on $\mathbb{R}^{d}$ and $\sigma:=\operatorname{vol}_{\mathrm{S}_{1}^{d-1}}$ is the Lebesgue measure on $S_{1}^{d-1}$.

Since the gas particles are assumed to be non-interacting, to study the evolution of a particle cloud, we may just as well consider the orbit $t \mapsto(\boldsymbol{q}(t), \boldsymbol{v}(t))$ of a single point particle starting from a random point $\left(\boldsymbol{q}_{0}, \boldsymbol{v}_{0}\right)$, chosen according to a given probability measure on the phase space $\mathrm{T}^{1}\left(\mathcal{K}_{\rho}\right)$. Then $t \mapsto(\boldsymbol{q}(t), \boldsymbol{v}(t))$ becomes a random flight process, which we call the Lorentz process. A central challenge is to determine whether, in the limit of small scatterer density (that is as $\rho \rightarrow 0$ ), the Lorentz process converges to a limiting stochastic process. In order to give a precise formulation of this question, we assume from now on that $\mathcal{P}$ has an asymptotic density, meaning that there exists a constant $\bar{n}_{\mathcal{P}}>0$ such that for any bounded set $\mathcal{D} \subset \mathbb{R}^{d}$ with boundary of Lebesgue measure zero, we have

$$
\begin{equation*}
\lim _{R \rightarrow \infty} \frac{\#(\mathcal{P} \cap R \mathcal{D})}{R^{d}}=\bar{n}_{\mathcal{P}} \operatorname{vol}(\mathcal{D}) \tag{1.1}
\end{equation*}
$$

Then a simple heuristic argument shows that the mean free path length, i.e. the mean time between consecutive collisions, should be expected to scale as $\rho^{1-d}$ as $\rho \rightarrow 0$. It is therefore natural to consider the so-called Boltzmann-Grad scaling, in which length and time units are rescaled by a factor of $\rho^{1-d}$. That is, we consider the macroscopic coordinates

$$
(\boldsymbol{Q}(t), \boldsymbol{V}(t))=\left(\rho^{d-1} \boldsymbol{q}\left(\rho^{-(d-1)} t\right), \boldsymbol{v}\left(\rho^{-(d-1)} t\right)\right)
$$

The challenge now is to prove that the rescaled random flight process $t \mapsto(\boldsymbol{Q}(t), \boldsymbol{V}(t))$ converges to a limiting random flight process as $\rho \rightarrow 0$.

There are two instances where this problem has been fully understood for some years. The first is the case when $\mathcal{P}$ is a fixed realisation of a Poisson point process. Here Boldrighini, Bunimovich and Sinai [2] proved that the Lorentz process converges to a limit that is consistent
with the linear Boltzmann equation (cf. also [9] and [29]). In fact the paper [2] is restricted to dimension $d=2$ and hard sphere scatterers, but it was proved in [21 that the results generalise to general dimensions and soft scattering potentials.

The second instance is when the scatterer configuration $\mathcal{P}$ equals a Euclidean lattice $\mathcal{L}$ of full rank in $\mathbb{R}^{d}$. For this case, Marklof and Strömbergsson [16, 17, 15, 18] proved convergence of the Lorentz process to a limiting random flight process which in fact only depends on the asymptotic density of $\mathcal{L}$. The limit process is Markovian only on an extended phase space which, in addition to position and momentum, also includes the impact parameter and distance to the next collision. The corresponding transport equation is in particular not consistent with the linear Boltzmann equation. This new transport equation was obtained independently in dimension $d=2$ for $\mathcal{P}=\mathbb{Z}^{2}$ by Caglioti and Golse [4, 5, subject to a heuristic assumption that was proved (in any dimension) in [17]. In the lattice setting, the limit transport process in fact satisfies a superdiffusive central limit theorem [22], with a mean-square displacement proportional to $t \log t$ (where $t$ is time measured in units of the mean collision time), rather than the standard linear scaling which appears in the case of random scatterer configurations.

In a recent paper by Marklof and Strömbergsson [21, a general framework was developed which, under a certain set of hypotheses on the scatterer configuration $\mathcal{P}$, allows the proof of convergence of the rescaled Lorentz process $t \mapsto(\boldsymbol{Q}(t), \boldsymbol{V}(t))$ to a limiting random flight process. This framework was proved to apply when $\mathcal{P}$ belongs to a certain class of quasicrystals (this includes in particular the case when $\mathcal{P}$ is a general periodic point set), and also in the case when $\mathcal{P}$ is a fixed realisation of a Poisson point process of constant intensity.

Our main goal in the present paper is to prove that the framework from [21] also applies in the case when the scatterer configuration $\mathcal{P}$ is an arbitrary finite union of grids. (By definition, a 'grid' is a translate of a full rank lattice in $\mathbb{R}^{d}$.) That is, we will assume that

$$
\begin{equation*}
\mathcal{P}=\bigcup_{i=1}^{N} \mathcal{L}_{i} \tag{1.2}
\end{equation*}
$$

where each $\mathcal{L}_{i}$ is a grid. As we will see, this case serves as a nice testing ground for the framework developed in [21, and exhibits new features compared with the previous cases where that framework has been proved to apply.

In the special case when the $\mathcal{L}_{i}$ are pairwise incommensurabl $\rrbracket$, the scatterer configurations in (1.2) have previously been considered in the paper [20], where among other things a conjectural description of the Boltzmann-Grad limit of the Lorentz process was given. The main result of the present paper settles that conjecture as a special case. However, the general case of $\mathcal{P}$ as in (1.2), without the incommensurability assumption, is considerably more difficult and involves interesting new phenomena (see Section 6.3).
Remark 1.1. In the case when some scatterers in the family $\left\{\boldsymbol{p}+\mathcal{B}_{\rho}^{d}: \boldsymbol{p} \in \mathcal{P}\right\}$ overlap, certain technical annoyances appear in the definition of the Lorentz process $t \mapsto(\boldsymbol{q}(t), \boldsymbol{v}(t))$. Let us note that overlapping scatterers in general do exist in the situation studied in the present paper. Indeed, if $\mathcal{P}$ is any finite union of grids which is not periodic, then for every $\rho>0$ there exist points $\boldsymbol{p} \neq \boldsymbol{p}^{\prime}$ in $\mathcal{P}$ with $\left\|\boldsymbol{p}-\boldsymbol{p}^{\prime}\right\|<2 \rho$. In the classical case when the particles interact with the scatterers through specular reflection, this issue is handled in a standard manner [21, Ch. 1.2]: in this case the Lorentz flow equals the standard billiard flow in the region $\mathcal{K}_{\rho}$, and this flow is technically defined only on a subset of $\mathrm{T}^{1}\left(\mathcal{K}_{\rho}\right)$ of full measure with respect to the Liouville measure $\operatorname{vol} \times \sigma$; the exceptional points include all points $(\boldsymbol{q}, \boldsymbol{v}) \in \partial \mathcal{K}_{\rho} \times \mathrm{S}_{1}^{d-1}$ for which $\boldsymbol{v}$ points into a scatterer, and also, in the case of scatterer overlaps, any initial condition for which the particle at some time point either in the past or in the future collides with an intersection point of two or more scatterer boundaries. In the case of scatterer interaction through a more general scattering map, for example a scattering

[^1]process generated by a spherically symmetric potential, a simple way to avoid intricacies in the definition of the Lorentz flow is to remove, in an ad hoc manner, scatterer centers in $\mathcal{P}$ causing overlap [21, Ch. 1.3]. Since the probability of the particle hitting a scatterer which overlaps with another scatterer tends to zero in the Boltzmann-Grad limit, the limiting random flight process becomes the same independently of the precise choice of convention.
1.1. The limiting flight process. We will assume throughout the paper that the fixed scattering process of the Lorentz gas satisfies the conditions in [21, Sec. 3.4]. Special cases include the case of hard sphere scatterers as well as a general class of soft scattering potentials.

The main result of the present paper is Theorem 2.1 in Section 2, which states that all the assumptions required in the framework in [21] are satisfied in the case when $\mathcal{P}$ is an arbitrary finite union of grids. By the results of [21], this implies the existence of the Boltzmann-Grad limit of the particle dynamics in the Lorentz gas. To make a precise statement, let us write $F_{t, \rho}$ for the rescaled Lorentz flow;

$$
(\boldsymbol{Q}(t), \boldsymbol{V}(t))=F_{t, \rho}(\boldsymbol{Q}(0), \boldsymbol{V}(0))
$$

For notational reasons we extend the dynamics to the inside of each scatterer trivially, that is, set $F_{t, \rho}=$ id whenever $\boldsymbol{Q}$ is inside the scatterer. Thus $F_{t, \rho}$ is now a flow defined on all of $\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)$, the unit tangent bundle of $\mathbb{R}^{d}$. Let $P_{\mathrm{ac}}\left(\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)\right)$ be the set of Borel probability measures on $\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)$ which are absolutely continuous with respect to the Liouville measure $\operatorname{vol} \times \sigma$.

The following theorem is a consequence of Theorem 2.1 below, in combination with the main results of [21] and in particular [21, Sec. 4.5].

Theorem 1.1. Let the scatterer configuration $\mathcal{P}$ be a finite union of grids. Then, for any $\Lambda \in P_{\mathrm{ac}}\left(\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)\right)$, there exists a random flight process $\Xi=\{\Xi(t): t \geq 0\}$ such that the random process

$$
\Xi^{(\rho)}: t \mapsto \Xi^{(\rho)}(t)=F_{t, \rho}\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right)
$$

obtained by taking the initial data $\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right)$ random with respect to $\Lambda$, converges to $\Xi$ in distribution, as $\rho \rightarrow 0$.

The next result, Theorem 1.2, gives a description of the limiting process $\Xi$ in Theorem 1.1, The key point is that, after introducing a certain marking of the points in $\mathcal{P}$, the process $\Xi$ can be described as the flow with unit speed along a random piecewise linear curve, whose path segments, when considered in combination with the marks of the scatterers involved in the collisions, are generated by a Markov process with memory two. This is a special case of a corresponding result in [21]; in the setting of the present paper with $\mathcal{P}$ being a finite union of grids, the set of marks can be taken to be a certain concrete finite set, which we now turn to describe.

From now on, we will express $\mathcal{P}$ as

$$
\begin{equation*}
\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi}, \tag{1.3}
\end{equation*}
$$

where $\Psi$ is a finite set of indices, and each $\mathcal{L}_{\psi}$ is a grid in $\mathbb{R}^{d}$. In the statement of Theorem 1.2 we must require that the above presentation of $\mathcal{P}$ is admissible. This is a somewhat technical notion which we will define in Section 5.2, we prove in Section 5.3 that every finite union of grids possesses an admissible presentation. We give a couple of simple examples: If the grids $\mathcal{L}_{\psi}$ are pairwise incommensurable, then the presentation in (1.3) is always admissible. On the other hand, for any $q \in \mathbb{Z}_{\geq 2}$ and $\boldsymbol{v} \in \mathbb{R}^{d} \backslash \mathbb{Z}^{d}$, the presentation $\mathcal{P}=\mathbb{Z}^{d} \cup\left(q \mathbb{Z}^{d}+\boldsymbol{v}\right)$ is admissible if and only if $\boldsymbol{v} \notin \mathbb{Q}^{d}$; in the remaining case $\boldsymbol{v} \in \mathbb{Q}^{d} \backslash \mathbb{Z}^{d}$, an admissible presentation of $\mathcal{P}=\mathbb{Z}^{d} \cup\left(q \mathbb{Z}^{d}+\boldsymbol{v}\right)$ can be given as the union of the grids $q \mathbb{Z}^{d}+\boldsymbol{\alpha}$ with $\boldsymbol{\alpha}$ running through $\{1, \ldots, q\}^{d} \sqcup\{\boldsymbol{v}\}$.

Assuming that (1.3) is an admissible presentation of $\mathcal{P}$, let us now fix a choice of a function

$$
\begin{equation*}
\psi: \mathcal{P} \rightarrow \Psi, \quad \text { subject to } \boldsymbol{p} \in \mathcal{L}_{\psi(\boldsymbol{p})} \quad(\forall \boldsymbol{p} \in \mathcal{P}) \tag{1.4}
\end{equation*}
$$

We will call this function a crude marking ${ }^{2}$ of $\mathcal{P}$. Note that the condition in (1.4) forces the choice of $\psi(\boldsymbol{p})$ for every $\boldsymbol{p} \in \mathcal{P}$ which does not lie in more than one of the sets $\mathcal{L}_{\psi}$. That is, the flexibility involved in the choice of the crude marking $\psi: \mathcal{P} \rightarrow \Psi$ concerns only the points in the union $\bigcup_{\psi \neq \psi^{\prime} \in \Psi}\left(\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}\right)$. It will turn out that the exact choice of the mark $\psi(\boldsymbol{p})$ at these points has no influence on the description of the limiting flight process.

For each $\psi \in \Psi$, let $\bar{n}_{\psi}$ be the asymptotic density of $\mathcal{L}_{\psi}$. It follows that the asymptotic density of $\mathcal{P}$ is given by $\bar{n}_{\mathcal{P}}:=\sum_{\psi \in \Psi} \bar{n}_{\psi}$. We equip $\Psi$ with the probability measure m defined by

$$
\begin{equation*}
\mathrm{m}(\psi)=\frac{\bar{n}_{\psi}}{\bar{n}_{\mathcal{P}}} \quad(\psi \in \Psi) \tag{1.5}
\end{equation*}
$$

(We are abusing notation for increased readability; naturally, "m $(\psi)$ " denotes $\mathrm{m}(\{\psi\})$.)
The following result follows from [21, Theorem 4.6], in view of Theorem 2.1 below.
Theorem 1.2. Let $\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi}$ be an admissible presentation of $\mathcal{P}$, and let $\psi: \mathcal{P} \rightarrow \Psi$ be a corresponding crude marking. Let $\Lambda \in P_{\mathrm{ac}}\left(\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)\right)$. For $\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right)$ random with respect to $\Lambda$, let the corresponding random trajectory $t \mapsto F_{t, \rho}\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right)$ be described by the random variables $\xi_{j}^{(\rho)} \in \mathbb{R}_{>0}, \psi_{j}^{(\rho)} \in \Psi$ and $\boldsymbol{V}_{j}^{(\rho)} \in S_{1}^{d-1}$, where $\xi_{j}^{(\rho)}$ is the length of the $j$ th path segment, $\psi_{j}^{(\rho)}$ is the mark of the scatterer involved in the jth collision and $\boldsymbol{V}_{j}^{(\rho)} \in \mathrm{S}_{1}^{d-1}$ is the velocity after the $j$ th collision. Then as $\rho \rightarrow 0$, the random process

$$
\left(\left\langle\xi_{j}^{(\rho)}, \psi_{j}^{(\rho)}, \boldsymbol{V}_{j}^{(\rho)}\right\rangle\right)_{j=1,2, \ldots}
$$

converges in distribution to the second-order Markov process

$$
\begin{equation*}
\left(\left\langle\xi_{j}, \psi_{j}, \boldsymbol{V}_{j}\right\rangle\right)_{j=1,2, \ldots} \tag{1.6}
\end{equation*}
$$

where for any Borel set $A \subset \mathbb{R}_{>0} \times \Psi \times \mathrm{S}_{1}^{d-1}$,

$$
\begin{equation*}
\mathbb{P}\left(\left\langle\xi_{1}, \psi_{1}, \boldsymbol{V}_{1}\right\rangle \in A \mid\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right)\right)=\int_{A} p^{(\psi)}\left(\boldsymbol{V}_{0} ; \xi, \boldsymbol{V}\right) d \xi d \mathbf{m}(\psi) d \sigma(\boldsymbol{V}) \tag{1.7}
\end{equation*}
$$

and for $j \geq 2$,

$$
\begin{align*}
\mathbb{P}\left(\left\langle\xi_{j}, \psi_{j},\right.\right. & \left.\left.\boldsymbol{V}_{j}\right\rangle \in A \mid\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right),\left\langle\left(\xi_{i}, \psi_{i}, \boldsymbol{V}_{i}\right)\right\rangle_{i=1}^{j-1}\right)  \tag{1.8}\\
& =\int_{A} p^{\left(\psi_{j-1} \rightarrow \psi\right)}\left(\boldsymbol{V}_{j-2}, \boldsymbol{V}_{j-1} ; \xi, \boldsymbol{V}\right) d \xi d \mathrm{~m}(\psi) d \sigma(\boldsymbol{V})
\end{align*}
$$

The functions $p^{(\psi)}$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ depend on $\mathcal{P}$ but are independent of $\Lambda$, and for any fixed $\boldsymbol{V}^{\prime}, \boldsymbol{V}^{\prime \prime}, \psi^{\prime}$, both $p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ are probability densities on $\mathbb{R}_{>0} \times$ $\Psi \times \mathrm{S}_{1}^{d-1}$ with respect to the measure $d \xi \times d \mathrm{~m}(\psi) \times d \sigma(\boldsymbol{V})$.

We call the functions $p^{(\psi)}$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ collision kernels. Explicit formulas for these will be given in Section 10.1, as pushforwards of invariant measures on certain homogeneous spaces associated to the given point set $\mathcal{P}$. Both $p^{(\psi)}$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ are rotationally invariant in the sense that $p^{(\psi)}\left(\boldsymbol{V}^{\prime} K ; \xi, \boldsymbol{V} K\right)=p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime} K, \boldsymbol{V}^{\prime} K ; \xi, \boldsymbol{V} K\right)=$ $p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ for any $K \in \mathrm{SO}(d)$. We also mention that by definition we have $p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=0$ unless $\boldsymbol{V} \in \mathcal{V}_{\boldsymbol{V}^{\prime}}$, and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=0$ unless $\boldsymbol{V}^{\prime} \in \mathcal{V}_{\boldsymbol{V}^{\prime \prime}}$ and $\boldsymbol{V} \in \mathcal{V}_{\boldsymbol{V}^{\prime}}$. Here, for any $\boldsymbol{V} \in \mathrm{S}_{1}^{d-1}, \mathcal{V}_{\boldsymbol{V}}$ is the set of possible exit velocities in a scatterer collision with entry velocity $\boldsymbol{V}$ [21, Secs. 3.4-5].

The limiting random flight process $\Xi$ discussed above is closely related to the dynamics of a particle cloud in the Boltzmann-Grad limit of the Lorentz gas. Indeed, the time evolution

[^2]of an initial particle density $f_{0} \in \mathrm{~L}^{1}\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right)$ in the Lorentz gas with fixed scatterer radius $\rho$ is given by $f_{t}=L_{\rho}^{t} f_{0}$ where $L_{\rho}^{t}$ is the Liouville operator defined by
$$
\left[L_{\rho}^{t} f_{0}\right](\boldsymbol{Q}, \boldsymbol{V}):=f_{0}\left(F_{-t, \rho}(\boldsymbol{Q}, \boldsymbol{V})\right) .
$$

The existence of the limiting stochastic process $\Xi(t)$ implies that for every $t>0$ there exists a linear operator $L_{t}: \mathrm{L}^{1}\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right) \rightarrow \mathrm{L}^{1}\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right)$ such that for every $f_{0} \in \mathrm{~L}^{1}\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right)$ and every set $A \subset \mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)$ with boundary of Liouville measure zero,

$$
\lim _{\rho \rightarrow 0} \int_{A}\left[L_{t}^{(\rho)} f_{0}\right](\boldsymbol{Q}, \boldsymbol{V}) d \boldsymbol{Q} d \sigma(\boldsymbol{V})=\int_{A} L_{t} f_{0}(\boldsymbol{Q}, \boldsymbol{V}) d \boldsymbol{Q} d \sigma(\boldsymbol{V})
$$

[21, Cor. 1.4]. Under suitable continuity assumptions we can in fact express $L_{t} f_{0}$ as

$$
\left[L_{t} f_{0}\right](\boldsymbol{Q}, \boldsymbol{V})=\int_{\mathbb{R}_{>0} \times \Psi \times \mathrm{S}_{1}^{d-1}} f^{(\psi)}\left(t, \boldsymbol{Q}, \boldsymbol{V}, \xi, \boldsymbol{V}_{+}\right) d \xi d \mathrm{~m}(\psi) d \sigma\left(\boldsymbol{V}_{+}\right)
$$

where the functions $f^{(\psi)}\left(t, \boldsymbol{Q}, \boldsymbol{V}, \xi, \boldsymbol{V}_{+}\right)$form the unique solution of the system of differential equations

$$
\begin{align*}
\left(\partial_{t}+\boldsymbol{V}\right. & \left.\cdot \nabla_{\boldsymbol{Q}}-\partial_{\xi}\right) f^{\left(\psi_{+}\right)}\left(t, \boldsymbol{Q}, \boldsymbol{V}, \xi, \boldsymbol{V}_{+}\right) \\
& =\int_{\Psi \times \mathrm{S}_{1}^{d-1}} f^{(\psi)}\left(t, \boldsymbol{Q}, \boldsymbol{V}_{0}, 0, \boldsymbol{V}\right) p^{\left(\psi \rightarrow \psi_{+}\right)}\left(\boldsymbol{V}_{0}, \boldsymbol{V} ; \xi, \boldsymbol{V}_{+}\right) d \mathrm{~m}(\psi) d \sigma\left(\boldsymbol{V}_{0}\right) \tag{1.9}
\end{align*}
$$

subject to the initial condition $f^{\left(\psi_{+}\right)}\left(0, \boldsymbol{Q}, \boldsymbol{V}, \xi, \boldsymbol{V}_{+}\right)=f_{0}(\boldsymbol{Q}, \boldsymbol{V}) p^{\left(\psi_{+}\right)}\left(\boldsymbol{V} ; \xi, \boldsymbol{V}_{+}\right)$. See [21, Sections 1.4 and 4.6].

Equation (1.9) may be viewed as a generalization of the linear Boltzmann equation; it is the forward Kolmogorov equation (or Fokker-Planck-Kolmogorov equation) of a natural extension of $\Xi$ to a Markov flight process on the space $\mathrm{T}^{1}\left(\mathbb{R}^{d}\right) \times \mathbb{R}_{>0} \times \Psi \times \mathrm{S}_{1}^{d-1}$ [21, (1.24)].

In the special case when the grids $\mathcal{L}_{i}$ in (1.2) are pairwise incommensurable, the limiting random flight process described in Theorem 1.2, as well as the generalized linear Boltzmann equation in (1.9), agree with the corresponding limits conjectured in [20]; see Remark 10.5 below.
1.2. Expressing $\Xi$ via the commensurability classes. An important result of the present paper is that the collision kernels can be expressed as products over collision kernels associated to the commensurability classes of the grids appearing in $\mathcal{P}$. In the special case when the grids in $\mathcal{P}$ are pairwise incommensurable, such product formulas were given in [20, Sec. 5].

Recall that two grids $\mathcal{L}$ and $\mathcal{L}^{\prime}$ in $\mathbb{R}^{d}$ are said to be commensurable if there exist $\delta>0$ and $\boldsymbol{v} \in \mathbb{R}^{d}$ such that $\mathcal{L} \cap\left(\delta \mathcal{L}^{\prime}+\boldsymbol{v}\right)$ is a grid. This is an equivalence relation on the family of grids in $\mathbb{R}^{d}$. Let us say that two marks $\psi, \psi^{\prime} \in \Psi$ are equivalent (and write $\psi \sim \psi^{\prime}$ ) when $\mathcal{L}_{\psi}$ and $\mathcal{L}_{\psi^{\prime}}$ are commensurable. For any $\psi \in \Psi$ we denote by $[\psi]:=\left\{\psi^{\prime} \in \Psi: \psi^{\prime} \sim \psi\right\}$ the corresponding equivalence class.

Let $C_{\Psi}$ be the family of equivalence classes in $\Psi$, and set, for any $c \in C_{\Psi}$,

$$
\begin{equation*}
\mathcal{P}_{c}:=\bigcup_{\psi \in c} \mathcal{L}_{\psi} . \tag{1.10}
\end{equation*}
$$

Note that $\mathcal{P}_{c}$ is a subset of $\mathcal{P}$, and $\mathcal{P}_{c}$ is itself a finite union of grids, so that all the results discussed so far also hold for the Lorentz gas with $\mathcal{P}_{c}$ as scatterer configuration. The asymptotic density of $\mathcal{P}_{c}$ is

$$
\begin{equation*}
\bar{n}_{c}:=\sum_{\psi \in c} \bar{n}_{\psi}, \tag{1.11}
\end{equation*}
$$

and it will be immediate from our definition of admissibility that for any $c \in C_{\Psi}$, the presentation $\mathcal{P}_{c}=\cup_{\psi \in c} \mathcal{P}_{\psi}$ is admissible; hence Theorem 1.2 applies to the scatterer configuration $\mathcal{P}_{c}$ with this presentation. For any two $\psi, \psi^{\prime} \in c$, let us denote by ${ }^{\circ}{ }^{\circ}(\psi)$ and $\left.{ }^{\circ}{ }^{\circ}{ }^{( } \psi^{\prime} \rightarrow \psi\right)$ the collision kernels for the corresponding limiting flight process. Thus ${ }^{\circ}{ }^{\circ}{ }^{(\psi)}$ is defined for any
$\psi \in \Psi$, and ${ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ is defined for any $\psi^{\prime}, \psi \in \Psi$ satisfying $\psi^{\prime} \sim \psi$. Also, for any $c \in C_{\Psi}$, letting $\mathrm{m}_{c}$ be the probability measure on $c$ defined by

$$
\begin{equation*}
\mathrm{m}_{c}(\psi)=\frac{\bar{n}_{\psi}}{\bar{n}_{c}} \tag{1.12}
\end{equation*}
$$

we have that for any fixed $\boldsymbol{V}^{\prime}, \boldsymbol{V}^{\prime \prime}$ and $\psi^{\prime} \in c$, both ${ }^{\circ} p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ and ${ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)$ are probability densities on $\mathbb{R}_{>0} \times c \times \mathrm{S}_{1}^{d-1}$ with respect to the measure $d \xi \times d \mathrm{~m}_{c}(\psi) \times d \sigma(\boldsymbol{V})$.

The product formula for collision kernels is stated in Theorem 10.3 below (see also Remark (10.4). Here we give an equivalent formulation saying that the Markov process which generates $\Xi$ can be obtained by merging in a certain way mutually independent versions of the Markov processes which arise in the Boltzmann-Grad limit for Lorentz gases which scatterer configurations $\mathcal{P}_{c}$. To prepare for the statement, let $\boldsymbol{V}^{\prime}, \boldsymbol{V}^{\prime \prime}$ be arbitrary, fixed vectors in $\mathrm{S}_{1}^{d-1}$ subject to $\boldsymbol{V}^{\prime} \in \mathcal{V}_{\boldsymbol{V}^{\prime \prime}}$, that is, $\boldsymbol{V}^{\prime}$ is a possible exit velocity in a scatterer collision with entry velocity $\boldsymbol{V}^{\prime \prime}$. For every $c \in C_{\Psi}$, let $\left\langle\widetilde{\xi}_{c}, \widetilde{\psi}_{c}, \widetilde{\boldsymbol{V}}_{c}\right\rangle$ be a random variable in $\mathbb{R}_{>0} \times c \times \mathrm{S}_{1}^{d-1}$, with distribution

$$
\begin{equation*}
{ }^{\circ} p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) d \xi d \mathbf{m}_{c}(\psi) d \sigma(\boldsymbol{V}) \tag{1.13}
\end{equation*}
$$

Also for every $\psi^{\prime} \in \Psi$, let $\left\langle\widetilde{\xi}_{\left[\psi^{\prime}\right], \psi^{\prime}}, \widetilde{\psi}_{\left[\psi^{\prime}\right], \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{\left[\psi^{\prime}\right], \psi^{\prime}}\right\rangle$ be a random variable in $\mathbb{R}_{>0} \times\left[\psi^{\prime}\right] \times \mathrm{S}_{1}^{d-1}$, with distribution

$$
\begin{equation*}
{ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) d \xi d \mathrm{~m}_{\left[\psi^{\prime}\right]}(\psi) d \sigma(\boldsymbol{V}) \tag{1.14}
\end{equation*}
$$

We take all these random triples to be mutually independent. Furthermore, for any $\psi^{\prime} \in \Psi$ and $c \in C_{\Psi}$ with $\psi^{\prime} \notin c$, we set $\left\langle\widetilde{\xi}_{c, \psi^{\prime}}, \widetilde{\psi}_{c, \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{c, \psi^{\prime}}\right\rangle:=\left\langle\widetilde{\xi}_{c}, \widetilde{\psi}_{c}, \widetilde{\boldsymbol{V}}_{c}\right\rangle$. Thus the random variable $\left\langle\widetilde{\xi}_{c, \psi^{\prime}}, \widetilde{\psi}_{c, \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{c, \psi^{\prime}}\right\rangle$ is now defined for all $\psi^{\prime} \in \Psi$ and all $c \in C_{\Psi}$. Note that (1.13) and (1.14) are exactly the distributions which appear in (1.7) and (1.8) when Theorem 1.2 is applied to $\mathcal{P}_{c}$. Next let the random variable $\widetilde{c}$ be equal to that $c \in C_{\Psi}$ which minimizes $\widetilde{\xi}_{c}$, and let the random variable $\widetilde{c}_{\psi^{\prime}}$ be equal to that $c$ which minimizes $\widetilde{\xi}_{c, \psi^{\prime}}$. (These $\widetilde{c}$ and $\widetilde{c}_{\psi^{\prime}}$ are almost surely uniquely defined.)

Theorem 1.3. For given ${\underset{\sim}{V}}^{\prime \prime} \in \mathrm{S}_{1}^{d-1}$ and $\boldsymbol{V}^{\prime} \in \mathcal{V}_{\boldsymbol{V}^{\prime \prime}}$ and $\psi^{\prime} \in \Psi$, introduce the random variables $\left\langle\widetilde{\xi}_{c}, \widetilde{\psi}_{c}, \widetilde{\boldsymbol{V}}_{c}\right\rangle$ and $\left\langle\widetilde{\xi}_{c, \psi^{\prime}}, \widetilde{\psi}_{c, \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{c, \psi^{\prime}}\right\rangle$ (for all $c \in C_{\Psi}$ ) and $\widetilde{c}$ and $\widetilde{c}_{\psi^{\prime}}$ as above. Then the conditional distributions appearing in (1.7) and (1.8) of Theorem 1.2 are given by the following formulas: For any Borel set $A \subset \mathbb{R}_{>0} \times \Psi \times \mathrm{S}_{1}^{d-1}$,

$$
\begin{equation*}
\mathbb{P}\left(\left\langle\xi_{1}, \psi_{1}, \boldsymbol{V}_{1}\right\rangle \in A \mid\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right), \boldsymbol{V}_{0}=\boldsymbol{V}^{\prime}\right)=\mathbb{P}\left(\left\langle\widetilde{\xi}_{\widetilde{c}}, \widetilde{\psi}_{\widetilde{c}}, \widetilde{\boldsymbol{V}}_{\widetilde{c}}\right\rangle \in A\right) \tag{1.15}
\end{equation*}
$$

and for $j \geq 2$,

$$
\mathbb{P}\left(\left\langle\xi_{j}, \psi_{j}, \boldsymbol{V}_{j}\right\rangle \in A \mid\left(\boldsymbol{Q}_{0}, \boldsymbol{V}_{0}\right),\left\langle\left(\xi_{i}, \psi_{i}, \boldsymbol{V}_{i}\right)\right\rangle_{i=1}^{j-1}, \boldsymbol{V}_{j-2}=\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}_{j-1}=\boldsymbol{V}^{\prime}, \psi_{j-1}=\psi^{\prime}\right)
$$

$$
\begin{equation*}
=\mathbb{P}\left(\left\langle\widetilde{\xi}_{\widetilde{c}_{\psi^{\prime}}, \psi^{\prime}}, \widetilde{\psi}_{\widetilde{c}_{\psi^{\prime}}, \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{\widetilde{c}_{\psi^{\prime}}, \psi^{\prime}}\right\rangle \in A\right) \tag{1.16}
\end{equation*}
$$

Note that Theorem 1.3 is a tautology when $\# C_{\Psi}=1$, i.e. when all the grids $\mathcal{L}_{\psi}(\psi \in \Psi)$ are commensurable; but when $\# C_{\Psi}>1$, Theorem 1.3 in effect expresses the collision kernel of the limiting flight process for the scatterer configuration $\mathcal{P}$ in terms of the collision kernels for the subconfigurations $\mathcal{P}_{c}, c \in C_{\Psi}$ (again, see Remark 10.4 below).

We next point out a consequence of Theorem 1.3 concerning the limiting free path length distribution in the Boltzmann-Grad limit. Set, for any $\boldsymbol{V}^{\prime} \in \mathrm{S}_{1}^{d-1}$,

$$
\begin{equation*}
\Phi_{\mathcal{P}}(\xi)=\int_{\Psi \times \mathrm{S}_{1}^{d-1}} p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) d \mathrm{~m}(\psi) d \sigma(\boldsymbol{V}) \tag{1.17}
\end{equation*}
$$

The right hand side in (1.17) is independent of $\boldsymbol{V}^{\prime}$ in view of the rotational invariance of $p^{(\psi)}$, and $\Phi_{\mathcal{P}}(\xi)$ is a probability density on $\mathbb{R}_{>0}$ with respect to Lebesgue measure. In fact $\Phi_{\mathcal{P}}(\xi)$ is the density function for the free path length in the limiting flight process $\Xi$ when starting from
a generic point inside the billiard domain. We will show in Section 10 that the function $\Phi_{\mathcal{P}}(\xi)$ is continuous, decreasing, and $0 \leq \Phi_{\mathcal{P}}(\xi) \leq \bar{n}_{\mathcal{P}} v_{d-1}$ for all $\xi \in \mathbb{R}_{>0}$, where $v_{d-1}:=\operatorname{vol}\left(\mathcal{B}_{1}^{d-1}\right)$. Note that for any $c \in C_{\Psi}$, the formula (1.17) applied to $\mathcal{P}_{c}$ reads

$$
\begin{equation*}
\Phi_{\mathcal{P}_{c}}(\xi)=\int_{c \times S_{1}^{d-1}}{ }^{\circ}{ }^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) d \mathbf{m}_{c}(\psi) d \sigma(\boldsymbol{V}) \tag{1.18}
\end{equation*}
$$

We now have:

## Corollary 1.4.

$$
\Phi_{\mathcal{P}}(\xi)=-\frac{d}{d \xi} \prod_{c \in C_{\Psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime} .
$$

Proof. Given $\xi>0$, we apply (1.15) in Theorem 1.3 to the set $A:=(\xi, \infty) \times \Psi \times \mathrm{S}_{1}^{d-1}$. Then by (1.7) and (1.17), the left hand side of (1.15) equals $\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime}$. On the other hand, by the definition of the (mutualy independent) random variables $\left\langle\widetilde{\xi}_{c}, \widetilde{\psi}_{c}, \widetilde{\boldsymbol{V}}_{c}\right\rangle$, and the definition of $\widetilde{c}$, the right hand side of (1.15) equals $\prod_{c \in C_{\Psi}} \mathbb{P}\left(\widetilde{\xi}_{c}>\xi\right)$, which by (1.7) (for $\left.\mathcal{P}_{c}\right)$ and (1.18) equals $\prod_{c \in C_{\Psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}$. Hence:

$$
\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime}=\prod_{c \in C_{\Psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}
$$

and the corollary follows by differentiation, using the fact that the functions $\Phi_{\mathcal{P}}$ and $\Phi_{\mathcal{P}_{c}}$ are nonnegative and continuous.

In the special case when the grids in $\mathcal{P}$ are pairwise incommensurable, the formula in Corollary 1.4 was given in [20, (2.8)].
1.3. Asymptotic estimates for the free path length distribution. We will prove the following upper and lower bounds showing that the free path length distribution in the Boltzmann-Grad limit has a power-law tail.

Theorem 1.5. Let $\mathcal{P}$ be a finite union of grids, and let $N=\# C_{\Psi}$ be the number of commensurability classes of grids appearing in $\mathcal{P}$. Then there exist constants $0<c_{1}<c_{2}$ such that

$$
\begin{equation*}
c_{1} \xi^{-(N+1)}<\Phi_{\mathcal{P}}(\xi)<c_{2} \xi^{-(N+1)}, \quad \forall \xi \geq 1 \tag{1.19}
\end{equation*}
$$

In particular it follows that the expected value of the free path length from starting at a generic point, $\int_{0}^{\infty} \xi \Phi_{\mathcal{P}}(\xi) d \xi$, is finite if and only if $N \geq 2$.

In the special case when the grids in $\mathcal{P}$ are pairwise incommensurable, a more precise asymptotic formula for $\Phi_{\mathcal{P}}(\xi)$ as $\xi \rightarrow \infty$ was given in [20, Theorem 2]; the proof uses the product formula and the asymptotics for $\Phi_{\mathcal{P}}(\xi)$ in the special case when $\mathcal{P}$ is a lattice, which was obtained in [18, Theorem 1.13]. It would be interesting to try to generalize this asymptotic formula, and also the precise asymptotic estimates for the collision kernels obtained in [18], to the case when $\mathcal{P}$ is a general union of grids as considered in the present paper.

We remark that in the proof of Theorem [1.5, the critical case to handle is $N=1$, as the case of general $N \geq 2$ then follows immediately via the product formula in Corollary 1.4

Next we point out a consequence of Theorem 1.5 concerning the density function $\bar{\Phi}_{\mathcal{P}}(\xi)$ of the free path length between consecutive collisions in the random flight process $\Xi$ describing the Boltzmann-Grad limit of the Lorentz gas. We give an explicit formula for $\bar{\Phi}_{\mathcal{P}}(\xi)$ in (10.14) in Section 10, here we merely point out the relations

$$
\begin{equation*}
\Phi_{\mathcal{P}}(\xi)=\bar{n}_{\mathcal{P}} v_{d-1} \int_{\xi}^{\infty} \bar{\Phi}_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime} \tag{1.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{\infty} \xi \bar{\Phi}_{\mathcal{P}}(\xi) d \xi=\frac{1}{\bar{n} \mathcal{P} v_{d-1}} \tag{1.21}
\end{equation*}
$$

Note that the value in (1.21) is the expected value of the free path length between consecutive collisions. The following result is an easy consequence of Theorem 1.5 and the relation (1.20).
Corollary 1.6. The second moment of the free path length between consecutive collisions, $\int_{0}^{\infty} \xi^{2} \bar{\Phi}_{\mathcal{P}}(\xi) d \xi$, is finite if and only if $N \geq 2$.

It is interesting to ask about the long-time asymptotics of the limiting random flight process $t \mapsto \Xi(t)$; we expect that this process converges to Brownian motion under an appropriate rescaling. In the special case when $\mathcal{P}$ is a lattice, such a convergence result has been proved in [22]; the correct scaling factor in this case is $t \log t$ (superdiffusion). In view of Theorem [1.5] and Corollary 1.6, it is natural to guess that a similar superdiffusive central limit result holds whenever $\mathcal{P}$ is a finite union of commensurable grids $(N=1)$; however if the grids in $\mathcal{P}$ are not all commensurable ( $N \geq 2$ ), we expect that a standard diffusive central limit result holds, i.e. with the rescaling factor being linear in $t$.

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## 2. The hypotheses from [21] on the scatterer configuration $\mathcal{P}$

We now recall the precise formulation of the hypotheses on the scatterer configuration $\mathcal{P}$ which in the paper [21] were proved to imply the convergence of the Lorentz gas particle dynamics in the Boltzmann-Grad limit. We first need to introduce some technical notation.

For any topological space $S$, we write $P(S)$ for the set of Borel probability measures on $S$, equipped with the weak topology. We will only consider $P(S)$ when $S$ is separable and metrizable; recall that then also $P(S)$ is metrizable [1, pp. 72-73]. Also, given any locally compact second countable Hausdorff space $\mathcal{X}$, we let $N(\mathcal{X})$ be the family of locally finite counting measures on $\mathcal{X}$, equipped with the vague topology (then $N(\mathcal{X})$ is a Polish space), and let $N_{s}(\mathcal{X})$ be the subset of simple counting measures; this is a Borel subset of $N(\mathcal{X})$. The elements of $N_{s}(\mathcal{X})$ may be identified with the family of locally finite subsets of $\mathcal{X}$ through $\nu \mapsto \operatorname{supp}(\nu)$. The inverse map is $\left\{x_{i}\right\} \mapsto \sum_{i} \delta_{x_{i}}$. We will use this identification between point sets and simple counting measures throughout this work, often using the same notation for point set and counting measure. A point process in $\mathcal{X}$ is, by definition, a random element $\xi$ in $N(\mathcal{X})$. It is called simple if $\xi \in N_{s}(\mathcal{X})$ almost surely. We identify $P\left(N_{s}(\mathcal{X})\right)$ with the set of probability measures $\nu \in P(N(\mathcal{X}))$ with $\nu\left(N_{s}(\mathcal{X})\right)=1$. Then a point process $\xi$ is simple if and only if its law is in $P\left(N_{s}(\mathcal{X})\right)$.

In the present paper we will always have

$$
\begin{equation*}
\mathcal{X}:=\mathbb{R}^{d} \times \Sigma \tag{2.1}
\end{equation*}
$$

for some compact metric space $\Sigma$. We will view vectors in $\mathbb{R}^{d}$ as row vectors; thus $\mathrm{GL}_{d}(\mathbb{R})$ acts on $\mathbb{R}^{d}$ from the right. We extend the natural actions on $\mathbb{R}^{d}$ of $\mathbb{R}^{d}$ (by translation), of $\mathbb{R}^{\times}$ (by dilation) and of $\mathrm{GL}_{d}(\mathbb{R})$ (by multiplication from the right) by the trivial action on $\Sigma$; thus for any $(\boldsymbol{w}, \varsigma) \in \mathcal{X}$ and any $\boldsymbol{x} \in \mathbb{R}^{d}, T \in \mathbb{R}^{\times}$and $A \in \mathrm{GL}_{d}(\mathbb{R})$ we set $(\boldsymbol{w}, \varsigma)+\boldsymbol{x}:=(\boldsymbol{w}+\boldsymbol{x}, \varsigma)$, $T(\boldsymbol{w}, \varsigma):=(T \boldsymbol{w}, \varsigma)$ and $(\boldsymbol{w}, \varsigma) A:=(\boldsymbol{w} A, \varsigma)$. These actions also give rise to natural, continuous, actions on $N_{s}(\mathcal{X})$. For example, for any $A \in \mathrm{GL}_{d}(\mathbb{R})$ and $\mathcal{Q} \in N_{s}(\mathcal{X})$ (viewed as a locally finite subset of $\mathcal{X}), \mathcal{Q} A:=\{x A: x \in \mathcal{Q}\}$.

For $\rho>0$ we denote by $D_{\rho}$ the diagonal matrix

$$
\begin{equation*}
D_{\rho}=\operatorname{diag}\left(\rho^{d-1}, \rho^{-1}, \cdots, \rho^{-1}\right) \in \operatorname{SL}_{d}(\mathbb{R}) \tag{2.2}
\end{equation*}
$$

We also fix, once and for all, a map $R: \mathrm{S}_{1}^{d-1} \rightarrow \mathrm{SO}(d)$ with the property $\boldsymbol{v} R(\boldsymbol{v})=\boldsymbol{e}_{1}$ for all $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$, where $\boldsymbol{e}_{1}:=(1,0, \ldots, 0)$. We assume that $R$ is continuous when restricted to $\mathrm{S}_{1}^{d-1}$ minus one point; the choice of $R$ is otherwise arbitrary.

Now let $\mathcal{P}$ be an arbitrary fixed locally finite subset of $\mathbb{R}^{d}$ with constant asymptotic density $\bar{n}_{\mathcal{P}}$. We also assume given a compact metric space $\Sigma$, a map $\varsigma: \mathcal{P} \rightarrow \Sigma$, a Borel probability measure m on $\Sigma$, and a continuous map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P(N(\mathcal{X}))=P\left(N\left(\mathbb{R}^{d} \times \Sigma\right)\right)$. Set

$$
\begin{equation*}
\widetilde{\mathcal{P}}=\{(\boldsymbol{p}, \varsigma(\boldsymbol{p})): \boldsymbol{p} \in \mathcal{P}\} \subset \mathcal{X} \tag{2.3}
\end{equation*}
$$

We refer to $\varsigma$ as a marking of $\mathcal{P}$, and $\Sigma$ as the corresponding space of marks 3
For any $\boldsymbol{q} \in \mathbb{R}^{d}, \boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}$ and $\rho>0$, we set

$$
\widetilde{\mathcal{P}}_{\boldsymbol{q}}= \begin{cases}\widetilde{\mathcal{P}} \backslash\{(\boldsymbol{q}, \varsigma(\boldsymbol{q}))\} & (\boldsymbol{q} \in \mathcal{P})  \tag{2.4}\\ \widetilde{\mathcal{P}} & (\boldsymbol{q} \notin \mathcal{P})\end{cases}
$$

and

$$
\begin{equation*}
\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})=\left(\widetilde{\mathcal{P}}_{\boldsymbol{q}}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho} \tag{2.5}
\end{equation*}
$$

Given any $\boldsymbol{q} \in \mathbb{R}^{d}$ and $\lambda \in P\left(\mathrm{~S}_{1}^{d-1}\right)$, if we take $\boldsymbol{v}$ random in $\left(\mathrm{S}_{1}^{d-1}, \lambda\right)$ then $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ becomes a point process; we write $\mu_{\boldsymbol{q}, \rho}^{(\lambda)} \in P\left(N_{s}(\mathcal{X})\right)$ for its distribution. Finally set $\mu_{\mathcal{X}}=\operatorname{vol} \times \mathrm{m}$.

The assumption on the scatterer configuration $\mathcal{P}$, under which the convergence of the rescaled Lorentz process to a limiting flight process was proved in [21], is that $\mathcal{P}$ can be equipped with data $\varsigma, \Sigma$, m and $\varsigma \mapsto \mu_{\varsigma}$ as above, in such a way that the following six conditions [P1]-[P3] and [Q1]-[Q3] are fulfilled (see [21, Sec. 2.3]):
[P1] Uniform density: For any bounded $B \subset \mathcal{X}$ with $\mu_{\mathcal{X}}(\partial B)=0$, we have

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{\#(\widetilde{\mathcal{P}} \cap T B)}{T^{d}}=\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B) \tag{2.6}
\end{equation*}
$$

[P2] Spherical equidistribution: Let $P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ be the set of $\lambda \in P\left(\mathrm{~S}_{1}^{d-1}\right)$ which are absolutely continuous with respect to $\sigma$. There exists a subset $\mathcal{E} \subset \mathcal{P}$ of density zerd ${ }^{4}$ such that for any fixed $T \geq 1$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\mu_{\boldsymbol{q}, \rho}^{(\lambda)} \xrightarrow{\mathrm{w}} \mu_{\varsigma(\boldsymbol{q})} \quad \text { as } \rho \rightarrow 0, \text { uniformly for } \boldsymbol{q} \in \mathcal{P}_{T}(\rho):=\mathcal{P} \cap \mathcal{B}_{T \rho^{1-d}}^{d} \backslash \mathcal{E} .5 \tag{2.7}
\end{equation*}
$$

[P3] No escape of mass: For every bounded Borel set $B \subset \mathbb{R}^{d}$,

$$
\lim _{\xi \rightarrow \infty} \limsup _{\rho \rightarrow 0}[\operatorname{vol} \times \sigma]\left(\left\{(\boldsymbol{q}, \boldsymbol{v}) \in B \times \mathrm{S}_{1}^{d-1}: \mathcal{Q}_{\rho}\left(\rho^{1-d} \boldsymbol{q}, \boldsymbol{v}\right) \cap\left(\mathfrak{Z}_{\xi} \times \Sigma\right)=\emptyset\right\}\right)=0
$$

with the open cylinder $\mathfrak{Z}_{\xi}=(0, \xi) \times \mathcal{B}_{1}^{d-1} \subset \mathbb{R}^{d}$.
[Q1] $\mathrm{SO}(d-1)$-invariance: For every $\varsigma \in \Sigma$,
$\mu_{\varsigma}$ is invariant under the action of $\mathrm{SO}(d-1):=\left\{k \in \mathrm{SO}(d): \boldsymbol{e}_{1} k=\boldsymbol{e}_{1}\right\}$.
[Q2] Coincidence-free first coordinates: For every $\varsigma \in \Sigma$,

$$
\mu_{\varsigma}\left(\left\{\nu \in N(\mathcal{X}): \exists x_{1} \in \mathbb{R} \text { s.t. } \nu\left(\left\{x_{1}\right\} \times \mathbb{R}^{d-1} \times \Sigma\right)>1\right\}\right)=0
$$

[Q3] Small probability of large voids: For every $\varepsilon>0$ there exists $R>0$ such that for all $\varsigma \in \Sigma$ and $\boldsymbol{x} \in \mathbb{R}^{d}$ we have

$$
\begin{equation*}
\mu_{\varsigma}\left(\left\{\nu \in N(\mathcal{X}): \nu\left(\mathcal{B}_{R}^{d}(\boldsymbol{x}) \times \Sigma\right)=0\right\}\right)<\varepsilon \tag{2.8}
\end{equation*}
$$

Here $\mathcal{B}_{R}^{d}(\boldsymbol{x}):=\boldsymbol{x}+\mathcal{B}_{R}^{d}$, the open ball of radius $R$ centered at $\boldsymbol{x}$.

[^3]From a technical point of view, the following is the main result of the present paper:
Theorem 2.1. Let $\mathcal{P}$ be a union of grids in $\mathbb{R}^{d}$. Then there exists a compact metric space $\Sigma$, a marking $\varsigma: \mathcal{P} \rightarrow \Sigma$, a continuous map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P(N(\mathcal{X}))$, and a Borel probability measure m on $\Sigma$, such that all the assumptions [P1]-[P3] and [Q1]-[Q3] hold.

The explicit description of the data $\left[\Sigma, \varsigma, \varsigma \mapsto \mu_{\varsigma}, \mathrm{m}\right]$ which we use in the proof of Theorem 2.1 is given in Section 6.
Remark 2.1. In the special case when $\mathcal{P}$ is periodic, i.e. when $\mathcal{P}$ may be represented as in (1.2) with all the $\mathcal{L}_{i}$ s being translates of one and the same fixed lattice, Theorem 2.1 was proved in [21, Prop. 5.6]. See Remark 6.3 below for a more detailed comparison.

We end this section by pointing out an immediate consequence of [P2], which partly motivates the construction of $\Sigma$ in the present paper (see the discussion in Section 6.3).
Lemma 2.2. Assume given a locally finite subset $\mathcal{P} \subset \mathbb{R}^{d}$, a map $\varsigma$ from $\mathcal{P}$ to a compact metric space $\Sigma$, and a continuous map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P\left(N\left(\mathbb{R}^{d} \times \Sigma\right)\right)$, such that the condition $[P 2]$ holds. Let $p$ be the natural projection from $N\left(\mathbb{R}^{d} \times \Sigma\right)$ to $N\left(\mathbb{R}^{d}\right)$. Then the map $\varsigma \mapsto \widehat{\mu}_{\varsigma}:=p_{*}\left(\mu_{\varsigma}\right)$ from $\Sigma$ to $P\left(N\left(\mathbb{R}^{d}\right)\right)$ is continuous, and for each fixed $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$ (with $\mathcal{E}$ as in [P2]) and each $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have $p_{*}\left(\mu_{\boldsymbol{q}, \rho}^{(\lambda)}\right) \xrightarrow{\mathrm{w}} \widehat{\mu}_{\varsigma(\boldsymbol{q})}$ as $\rho \rightarrow 0$.

Note here that $p_{*}\left(\mu_{\boldsymbol{q}, \rho}^{(\lambda)}\right)$ is the distribution of the random point set $((\mathcal{P} \backslash\{\boldsymbol{q}\})-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho}$ in $\mathbb{R}^{d}$, for $\boldsymbol{v}$ random in $\left(S_{1}^{d-1}, \lambda\right)$.
Proof. As a map of counting measures, we have $p=\left(p_{1}\right)_{*}$ where $p_{1}$ is the projection map $\mathbb{R}^{d} \times \Sigma \rightarrow \mathbb{R}^{d}$; and one immediately verifies that $p$ is continuous (using the fact that $\Sigma$ is compact). Now the lemma follows by the continuous mapping theorem.
2.1. Outline of the rest of the paper. Sections 3 9 are devoted to the proof of Theorem [2.1, from which also Theorems 1.1 and 1.2 follow. Section 10 contains the proofs of Theorems 1.3 and 1.5, and their corollaries.

In Section 3 we introduce some further basic set-up and notation. In particular we fix the index set $\Psi$ in the presentation $\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi}$ to be

$$
\Psi=\left\{\psi=(j, i): j \in\{1, \ldots, N\}, i \in\left\{1, \ldots, r_{j}\right\}\right\},
$$

for some positive integers $N$ and $r_{1}, \ldots, r_{N}$, with the indices chosen in such a way that the grids $\mathcal{L}_{(j, i)}$ and $\mathcal{L}_{\left(j^{\prime}, i^{\prime}\right)}$ are commensurable if and only if $j=j^{\prime}$. We also introduce a certain product of homogeneous spaces $\mathbb{X}=\mathbb{X}_{1} \times \cdots \times \mathbb{X}_{N}$ which naturally parametrizes the families of unions of grids which appear when studying the spherical equidistribution condition [P2]. Here each space $\mathbb{X}_{j}$ is a torus fiber bundle over the space $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ of $d$-dimensional lattices of covolume one in $\mathbb{R}^{d}$; the fiber space is an $r_{j} d$-dimensional torus which we will denote by $\mathbb{T}_{j}^{d}$; it is equipped with a natural action by $\mathrm{SL}_{d}(\mathbb{Z})$ from the right.

In Section 4 we state an equidistribution result in the space $\mathbb{X}$, Theorem 4.2, which can be viewed as a precise description of the limit considered in the spherical equidistribution condition [P2]. In fact Theorem 4.2 applies for an arbitrary translation vector $\boldsymbol{q}$ in $\mathbb{R}^{d}$; however the theorem includes no statement about uniformity with respect to $\boldsymbol{q}$, as is required in [P2]. The corresponding limit probability measure on $\mathbb{X}$ is called $\mu^{(q)}$. In the remainder of Section 4 we give an alternative way to represent these measures $\mu^{(\boldsymbol{q})}$ via certain measures $\omega_{1}^{(\boldsymbol{q})}, \ldots, \omega_{N}^{(\boldsymbol{q})}$, where for each $j, \omega_{j}^{(\boldsymbol{q})}$ is an $\mathrm{SL}_{d}(\mathbb{Z})$-invariant probability measure on $\mathbb{T}_{j}^{d}$. We will denote by $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ the set of all $\mathrm{SL}_{d}(\mathbb{Z})$-invariant probability measures on $\mathbb{T}_{j}^{d}$; this is a compact subset of $P\left(\mathbb{T}_{j}^{d}\right)$.

In Section 5 we study the behaviour of each measure $\omega_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ tends to infinity within a given subgrid $\mathcal{L}_{\psi}$ of $\mathcal{P}$, and we explicitly identify the limit measure(s) obtained. In order for this limit to be unique in a certain sense, we need to introduce the admissibility property
of the presentation $\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi}$ mentioned in Section 1.1. The notion of admissibility is defined in Section 5.2, and we prove in Section 5.3 that an admissible presentation of $\mathcal{P}$ always exists.

In Section 6, building on the results from the previous sections, we are finally able to give the precise definition of the space of marks $\Sigma$, the marking $\varsigma: \mathcal{P} \rightarrow \Sigma$, and the map $\varsigma \mapsto \mu_{\varsigma}$, which we will use in the proof of Theorem 2.1. The construction of the space of marks $\Sigma$ is one of the key steps in the present paper; this construction is more complicated and non-intuitive than for any of the previous types of scatterer configurations for which the framework from [21] is known to apply (see [21, Ch. 5]). The crude marking $\psi: \mathcal{P} \rightarrow \Psi$ which we introduced in Section 1.1 will appear as a factor of $\varsigma: \mathcal{P} \rightarrow \Sigma$; in fact we will define $\Sigma$ to be a subset of the Cartesian product of $\Psi$ and the compact set $\Omega:=\prod_{j=1}^{N} P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$. In Section 6.3 we give a discussion motivating our construction of $\Sigma$.

In Section 7 we prove that the conditions [Q1]-[Q3] and [P1] and [P3] hold, and we reduce the remaining condition, $[\mathrm{P} 2]$, to a certain uniform equidistribution result taking place in the homogeneous space $\mathbb{X}$, Theorem 7.7. In Section 8 we prove a key result, Theorem 8.1, on equidistribution of certain expanding unipotent orbits in a slightly generalized version of the homogeneous space $\mathbb{X}=\Gamma \backslash G$; the proof builds on Ratner's classification of ergodic measures invariant under unipotent flows [25]. Theorem8.1] is then used in Section 9 in order to complete the proof of Theorem [7.7, and thus also of Theorem 2.1.

Finally, in Section 10 we start by recalling the formulas for the collision kernels $p^{(\psi)}$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ (see Theorem [1.2) in the general setting of the paper [21], and we then give the proofs of Theorem 10.3 (which we discussed above in Section 1.2), Theorem 1.3, Theorem 1.5 and Corollary 1.6 .

## 3. Further setup and notation

3.1. Representing the point set $\mathcal{P}$. As mentioned in the introduction, we will always view vectors in $\mathbb{R}^{d}$ as row vectors $6^{6}$

Recall that we are assuming that the scatterer configuration $\mathcal{P}$ is of the form $\mathcal{P}=\bigcup_{i=1}^{M} \mathcal{L}_{i}$, where each $\mathcal{L}_{i}$ is a grid (viz., a translate of a lattice) in $\mathbb{R}^{d}$. Recall also that two grids $\mathcal{L}$ and $\mathcal{L}^{\prime}$ in $\mathbb{R}^{d}$ are said to be commensurable if there exist $c>0$ and $\boldsymbol{v} \in \mathbb{R}^{d}$ such that $\mathcal{L} \cap\left(c \mathcal{L}^{\prime}+\boldsymbol{v}\right)$ is a grid. This is an equivalence relation on the family of grids in $\mathbb{R}^{d}$. By collecting the given grids $\mathcal{L}_{1}, \ldots, \mathcal{L}_{M}$ into commensurability classes, and then applying Lemma 3.1 below to each class, one sees that the set $\mathcal{P}$ can be represented as follows:

$$
\begin{equation*}
\mathcal{P}=\bigcup_{j=1}^{N} \bigcup_{i=1}^{r_{j}} c_{j, i}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i}\right) \tag{3.1}
\end{equation*}
$$

where $\mathcal{L}_{1}, \ldots, \mathcal{L}_{N}$ are pairwise incommensurable lattices in $\mathbb{R}^{d}$ of covolume one, $r_{1}, \ldots, r_{N}$ are positive integers, $c_{j, i}$ are positive real numbers and $\boldsymbol{v}_{j, i}$ are vectors in $\mathbb{R}^{d}$, and furthermore:
$\forall j \in\{1, \ldots, N\}: \forall i \neq i^{\prime} \in\left\{1, \ldots, r_{j}\right\}: c_{j, i} / c_{j, i^{\prime}} \in \mathbb{Q} \Rightarrow\left[c_{j, i}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i}\right) \cap c_{j, i^{\prime}}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i^{\prime}}\right)=\emptyset\right]$.
Our reason for requiring (3.2) is that if $c_{j, i} / c_{j, i^{\prime}} \in \mathbb{Q}$ and the two grids $c_{j, i}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i}\right)$ and $c_{j, i^{\prime}}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i^{\prime}}\right)$ were not disjoint, then their intersection would be a grid and in particular would have positive aymptotic density; this is a situation that we wish to avoid (see Lemma 3.2 and Remark 3.2 below).

In order to arrive at the statement above, we made use of the following lemma.
Lemma 3.1. Let $\mathcal{L}_{1}, \ldots, \mathcal{L}_{m}$ be commensurable grids in $\mathbb{R}^{d}$. Then there exists a lattice $\mathcal{L}$ in $\mathbb{R}^{d}$ of covolume one, a positive integer $r$, positive real numbers $c_{1}, \ldots, c_{r}$ and vectors $\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{r} \in$

[^4]$\mathbb{R}^{d}$ such that $\mathcal{L}_{1} \cup \cdots \cup \mathcal{L}_{m}=\bigcup_{i=1}^{r} c_{i}\left(\mathcal{L}+\boldsymbol{v}_{i}\right)$. In this representation, we may furthermore assume that for any $i \neq i^{\prime}$ in $\{1, \ldots, r\}$, if $c_{i} / c_{i^{\prime}} \in \mathbb{Q}$ then $c_{i}\left(\mathcal{L}+\boldsymbol{v}_{i}\right) \cap c_{i^{\prime}}\left(\mathcal{L}+\boldsymbol{v}_{i^{\prime}}\right)=\emptyset$.
Proof. Choose $\boldsymbol{w}_{1} \in \mathbb{R}^{d}$ such that $\mathcal{L}^{\prime}:=\mathcal{L}_{1}+\boldsymbol{w}_{1}$ is a lattice. Now for each $j \in\{2, \ldots, m\}$, the fact that $\mathcal{L}_{1}$ and $\mathcal{L}_{j}$ are commensurable implies that we can choose $\delta_{j}>0$ and $\boldsymbol{w}_{j} \in \mathbb{R}^{d}$ such that $\mathcal{L}^{\prime} \cap\left(\delta_{j} \mathcal{L}_{j}+\boldsymbol{w}_{j}\right)$ is a lattice. (In particular then $\mathbf{0} \in \delta_{j} \mathcal{L}_{j}+\boldsymbol{w}_{j}$, meaning that $\delta_{j} \mathcal{L}_{j}+\boldsymbol{w}_{j}$ is a lattice.) It follows that also $\cap_{j=2}^{m}\left(\delta_{j} \mathcal{L}_{j}+\boldsymbol{w}_{j}\right)$ intersects $\mathcal{L}^{\prime}$ in a lattice, and we can express this intersection as $\delta \mathcal{L}$, where $\mathcal{L}$ is a lattice of covolume one and $\delta$ is a positive real number. The fact that $\delta \mathcal{L}$ is a finite index subgroup of $\mathcal{L}^{\prime}$ implies that $\mathcal{L}^{\prime}$ is a union of a finite number of translates of $\delta \mathcal{L}$; hence so is $\mathcal{L}_{1}$. Similarly for each $j \in\{2, \ldots, m\}$, since $\delta \mathcal{L}$ is a finite index subgroup of $\delta_{j} \mathcal{L}_{j}+\boldsymbol{w}_{j}$, the grid $\mathcal{L}_{j}$ is a union of a finite number of translates of $\left(\delta / \delta_{j}\right) \mathcal{L}$. The first statement of the lemma follows by collecting all these unions into a single union.

To prove the second statement, we may for example proceed as follows. We start from a union $\bigcup_{i=1}^{r} c_{i}\left(\mathcal{L}+\boldsymbol{v}_{i}\right)$ with arbitrary numbers $c_{1}, \ldots, c_{r} \in \mathbb{R}_{>0}$ and vectors $\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{r} \in \mathbb{R}^{d}$. Partition the index set $\{1, \ldots, r\}$ by the equivalence relation $i \sim i^{\prime} \stackrel{\text { def }}{\Leftrightarrow} c_{i} / c_{i^{\prime}} \in \mathbb{Q}$, and for each equivalence class $J \subset\{1, \ldots, r\}$, let $c_{J}$ be the least common multiple of the numbers $\left\{c_{i}: i \in J\right\}$ (this exists since $c_{i} / c_{i^{\prime}} \in \mathbb{Q}$ for all $\left.i, i^{\prime} \in J\right)$. Then for each $i \in J$ we have $n_{i}:=c_{J} / c_{i} \in \mathbb{Z}^{+}$, and now $\bigcup_{i \in J} c_{i}\left(\mathcal{L}+\boldsymbol{v}_{i}\right)=\bigcup_{i \in J} \bigcup_{\boldsymbol{w} \in R_{i}} c_{J}\left(\mathcal{L}+\boldsymbol{w}+n_{i}^{-1} \boldsymbol{v}_{i}\right)$, where $R_{i} \subset n_{i}^{-1} \mathcal{L}$ is any fixed set of representatives for the quotient $n_{i}^{-1} \mathcal{L} / \mathcal{L}$. Removing any duplicate translate of $c_{J} \mathcal{L}$ from the last union, we arrive at an expression of the form $\bigcup_{i \in J} c_{i}\left(\mathcal{L}+\boldsymbol{v}_{i}\right)=\bigcup_{i=1}^{r_{J}} c_{J}\left(\mathcal{L}+\boldsymbol{v}_{i}^{\prime}\right)$ for some vectors $\boldsymbol{v}_{1}^{\prime}, \ldots, \boldsymbol{v}_{r_{J}}^{\prime} \in \mathbb{R}^{d}$ which are pairwise incongruent modulo $\mathcal{L}$. Applying this rewriting procedure to each equivalence class $J \subset\{1, \ldots, r\}$, we arrive at a representation of $\mathcal{L}_{1} \cup \cdots \cup \mathcal{L}_{m}$ which has the required property.
Remark 3.1. Note that for a given $\mathcal{P}$, the representation (3.1) is far from unique, even when we require (3.2).

Let us rewrite (3.1) slightly: Choose $M_{1}, \ldots, M_{N} \in \mathrm{SL}_{d}(\mathbb{R})$ so that $\mathcal{L}_{j}=\mathbb{Z}^{d} M_{j}$ (this is possible since $\mathcal{L}_{j}$ has covolume one). Set

$$
\begin{equation*}
\Psi=\left\{(j, i): j \in\{1, \ldots, N\}, i \in\left\{1, \ldots, r_{j}\right\}\right\} \tag{3.3}
\end{equation*}
$$

Also for each $\psi=(j, i) \in \Psi$, we set $c_{\psi}=c_{j, i}, \boldsymbol{w}_{\psi}=\boldsymbol{w}_{j, i}=\boldsymbol{v}_{j, i} M_{j}^{-1} \in \mathbb{R}^{d}$, and

$$
\begin{equation*}
\mathcal{L}_{\psi}:=c_{j, i}\left(\mathcal{L}_{j}+\boldsymbol{v}_{j, i}\right)=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j} \tag{3.4}
\end{equation*}
$$

With this notation, the asymptotic density of $\mathcal{L}_{\psi}$ is $\bar{n}_{\psi}=c_{\psi}^{-d}$, and we have:

$$
\begin{equation*}
\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi} \tag{3.5}
\end{equation*}
$$

as in (1.3).
Given $\psi \in \Psi$, we denote by $j_{\psi}$ and $i_{\psi}$ the indices such that $\psi=\left(j_{\psi}, i_{\psi}\right)$. The condition (3.2) now takes the following form:

$$
\begin{equation*}
\forall \psi \neq \psi^{\prime} \in \Psi: \quad\left[j_{\psi}=j_{\psi^{\prime}} \text { and } c_{\psi} / c_{\psi^{\prime}} \in \mathbb{Q}\right] \Rightarrow \mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}=\emptyset \tag{3.6}
\end{equation*}
$$

Also, the fact that $\mathcal{L}_{1}, \ldots, \mathcal{L}_{N}$ are pairwise incommensurable implies that

$$
\begin{equation*}
M_{j} M_{j^{\prime}}^{-1} \notin \mathcal{S}, \quad \forall j \neq j^{\prime} \in\{1, \ldots, N\} \tag{3.7}
\end{equation*}
$$

where $\mathcal{S}$ is the commensurator of $\mathrm{SL}_{d}(\mathbb{Z})$ in $\mathrm{SL}_{d}(\mathbb{R})$, i.e.

$$
\begin{equation*}
\mathcal{S}=\left\{(\operatorname{det} T)^{-1 / d} T: T \in \mathrm{GL}(d, \mathbb{Q}), \operatorname{det} T>0\right\} \tag{3.8}
\end{equation*}
$$

The presentation of $\mathcal{P}$ in (3.5) is the one which we will work with throughout the paper, and the conditions (3.6) and (3.7) will always be assumed to hold.
Lemma 3.2. In the above representation of $\mathcal{P}$, for any two $\psi \neq \psi^{\prime} \in \Psi$, the intersection $\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}$ is contained in an affine hyperplane of $\mathbb{R}^{d}$.

Proof. Assume that $\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi}$ is not contained in an affine hyperplane. Then there exist points $\boldsymbol{q}_{0}, \ldots, \boldsymbol{q}_{d} \in \mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}$ such that the vectors $\boldsymbol{q}_{j}-\boldsymbol{q}_{0}$ for $j=1, \ldots, d$ form a linear basis of $\mathbb{R}^{d}$. Then $\boldsymbol{q}_{0}+\mathbb{Z}\left(\boldsymbol{q}_{1}-\boldsymbol{q}_{0}\right)+\cdots+\mathbb{Z}\left(\boldsymbol{q}_{d}-\boldsymbol{q}_{0}\right)$ is a grid contained in $\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}$, and so $\mathcal{L}_{\psi}$ and $\mathcal{L}_{\psi^{\prime}}$ are commensurable, i.e. $j_{\psi}=j_{\psi^{\prime}}$. Hence by (3.4), both $c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right)$ and $c_{\psi^{\prime}}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi^{\prime}}\right)$ contain the points $\boldsymbol{q}_{i} M_{j_{\psi}}^{-1}(i=0,1, \ldots, d)$. Therefore $\left(c_{\psi} / c_{\psi^{\prime}}\right)\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}-\left(c_{\psi^{\prime}} / c_{\psi}\right) \boldsymbol{w}_{\psi^{\prime}}\right)$ intersects $\mathbb{Z}^{d}$ in more than one point, and hence $c_{\psi} / c_{\psi^{\prime}} \in \mathbb{Q}$. When combined with (3.6), this leads to a contradiction.

Given $\mathcal{P}$ as in (3.5) we now also fix a choice of a crude marking $\psi: \mathcal{P} \rightarrow \Psi$, as in (1.4).
Remark 3.2. As we discussed below (1.4), $\psi(\boldsymbol{p})$ is uniquely defined for every point $\boldsymbol{p} \in \mathcal{P}$ outside the union $\bigcup_{\psi \neq \psi^{\prime} \in \Psi}\left(\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}\right)$. By Lemma 3.2, this union is contained in a finite union of affine hyperplanes. In particular $\psi(\boldsymbol{p})$ is uniquely defined for every $\boldsymbol{p} \in \mathcal{P}$ away from a subset of density zero.
3.2. Lie groups and homogeneous spaces. We now introduce appropriate homogeneous spaces to keep track of the unions of grids which will appear in our discussion. (This is very standard; cf. [21, Sections $5.2 \& 5.3 .3]$, [6, Sec. 7] as well as many other references.)

For any positive integer $r$, we let $\mathrm{M}_{r \times d}(\mathbb{R})$ be the space of real $r \times d$ matrices, and let $\mathrm{S}_{r}(\mathbb{R})$ be the semidirect product

$$
\begin{equation*}
\mathrm{S}_{r}(\mathbb{R}):=\mathrm{SL}_{d}(\mathbb{R}) \ltimes \mathrm{M}_{r \times d}(\mathbb{R}) \tag{3.9}
\end{equation*}
$$

with multiplication law

$$
\begin{equation*}
(M, U)\left(M^{\prime}, U^{\prime}\right)=\left(M M^{\prime}, U M^{\prime}+U^{\prime}\right) . \tag{3.10}
\end{equation*}
$$

We also set

$$
\mathrm{S}_{r}(\mathbb{Z}):=\mathrm{SL}_{d}(\mathbb{Z}) \ltimes \mathrm{M}_{r \times d}(\mathbb{Z}) .
$$

In the special case $r=1$, we have $\mathrm{S}_{1}(\mathbb{R})=\mathrm{ASL}_{d}(\mathbb{R})$, the affine special linear group. This group acts on $\mathbb{R}^{d}$ from the right through

$$
\boldsymbol{w}(M, \boldsymbol{u}):=\boldsymbol{w} M+\boldsymbol{u} \quad\left(\boldsymbol{w} \in \mathbb{R}^{d},(M, \boldsymbol{u}) \in \operatorname{ASL}_{d}(\mathbb{R})\right) .
$$

It is natural to identify $\mathrm{S}_{r}(\mathbb{R})$ with a subgroup of $\mathrm{SL}_{d+r}(\mathbb{R})$ through $(M, U) \leftrightarrow\left(\begin{array}{cc}M & 0 \\ U & I\end{array}\right)$, since the group law (3.10) then corresponds to multiplication of matrices; however we will not make use of this identification in the present paper.

We will always consider $\mathrm{SL}_{d}(\mathbb{R})$ to be embedded in $\mathrm{S}_{r}(\mathbb{R})$ through $M \mapsto(M, 0)$; in other words any $M \in \mathrm{SL}_{d}(\mathbb{R})$ is understood to also denote the element $(M, 0)$ in $\mathrm{S}_{r}(\mathbb{R})$. In the opposite direction, we denote by $\iota$ the projection homomorphism

$$
\iota: \mathrm{S}_{r}(\mathbb{R}) \rightarrow \mathrm{SL}_{d}(\mathbb{R}) ; \quad \iota(M, U):=M
$$

For any $U \in \mathrm{M}_{r \times d}(\mathbb{R})$ we let $\mathrm{I}_{U}$ be the element

$$
\begin{equation*}
\mathrm{I}_{U}:=(\mathrm{I}, U) . \tag{3.11}
\end{equation*}
$$

in $\mathrm{S}_{r}(\mathbb{R})$. Note that the map $U \mapsto \mathrm{I}_{U}$ is a homomorphism from $\mathrm{M}_{r \times d}(\mathbb{R})$ to $\mathrm{S}_{r}(\mathbb{R})$. Furthermore, for each $i \in\{1, \ldots, r\}$, we let $\mathrm{r}_{i}: \mathrm{M}_{r \times d}(\mathbb{R}) \rightarrow \mathbb{R}^{d}$ be the projection map which takes any matrix to its $i$ th row, and we also define the following Lie group homomorphism:

$$
\mathrm{a}_{i}: \mathrm{S}_{r}(\mathbb{R}) \rightarrow \mathrm{ASL}_{d}(\mathbb{R}) ; \quad \mathrm{a}_{i}(M, U):=\left(M, \mathrm{r}_{i}(U)\right) .
$$

From now on we assume a point set $\mathcal{P}$ with presentation as in (3.3)-(3.5) to be fixed. In particular this means that the numbers $N, r_{1}, \ldots, r_{N}$ are now fixed; and so are the numbers $c_{\psi}$ and vectors $\boldsymbol{w}_{\psi}(\psi \in \Psi)$. We then set:

$$
G=\mathrm{S}_{r_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{R}) ; \quad \Gamma=\mathrm{S}_{r_{1}}(\mathbb{Z}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{Z}) ;
$$

and

$$
\mathbb{X}=\Gamma \backslash G
$$

We also write $G_{j}=\mathrm{S}_{r_{j}}(\mathbb{R}), \Gamma_{j}=\mathrm{S}_{r_{j}}(\mathbb{Z})$ and $\mathbb{X}_{j}=\Gamma_{j} \backslash G_{j}$, so that $G=G_{1} \times \cdots \times G_{N}$ and $\mathbb{X}=\mathbb{X}_{1} \times \cdots \times \mathbb{X}_{N}$; and we denote by $\pi$ the standard projection maps $G_{j} \rightarrow \mathbb{X}_{j}$ and $G \rightarrow \mathbb{X}$. We also write $p_{j}$ and $\widetilde{p}_{j}$ for the projection maps onto the $j$ th factor:

$$
\begin{equation*}
p_{j}: G \rightarrow \mathrm{~S}_{r_{j}}(\mathbb{R}) \quad \text { and } \quad \widetilde{p}_{j}: \mathbb{X} \rightarrow \mathbb{X}_{j} \quad(j \in\{1, \ldots, N\}), \tag{3.12}
\end{equation*}
$$

and for each $\psi=(j, i) \in \Psi$, we introduce the Lie group homomorphism $p_{\psi}$ through

$$
\begin{equation*}
p_{\psi}:=\mathrm{a}_{i} \circ p_{j}: G \rightarrow \operatorname{ASL}_{d}(\mathbb{R}) . \tag{3.13}
\end{equation*}
$$

We will use the space $\mathbb{X}$ to parametrize the set of all point sets that can be obtained from $\mathcal{P}=$ $\cup_{\psi \in \Psi} \mathcal{L}_{\psi}$ by applying arbitrary translations to each individual grid $\mathcal{L}_{\psi}$, and arbitrary linear maps of determinant one to each individual commensurability class, $\cup_{i=1}^{r_{j}} \mathcal{L}_{(j, i)}$. Specifically, we take $\Gamma g$ in $\mathbb{X}$ to parametrize the point set $J_{0}(\Gamma g)$, where

$$
\begin{equation*}
J_{0}: \mathbb{X} \rightarrow N_{s}\left(\mathbb{R}^{d}\right) ; \quad J_{0}(\Gamma g):=\bigcup_{\psi \in \Psi} c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right) \quad(g \in G) . \tag{3.14}
\end{equation*}
$$

To see that this map is well-defined, note that for each $\psi$ we have $p_{\psi}(\Gamma)=\operatorname{ASL}_{d}(\mathbb{Z})$; hence if $\Gamma g=\Gamma g^{\prime}$ then $p_{\psi}\left(g^{\prime}\right)=\gamma p_{\psi}(g)$ for some $\gamma \in \operatorname{ASL}_{d}(\mathbb{Z})$, implying that $\mathbb{Z}^{d} p_{\psi}\left(g^{\prime}\right)=\mathbb{Z}^{d} p_{\psi}(g)$. However it should be noted that the map $J_{0}$ is discontinuous at any point $\Gamma g \in \mathbb{X}$ for which the grids $c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right)(\psi \in \Psi)$ are not pairwise disjoint.

We next introduce notation for expressing the parametrization in $\mathbb{X}$ of an arbitrary translate of the fixed point set $\mathcal{P}$. For any $\boldsymbol{q} \in \mathbb{R}^{d}$ and $j \in\{1, \ldots, N\}$ we define the matrix $U_{j}^{(\boldsymbol{q})}$ by specifying its row vectors:

$$
\begin{equation*}
U_{j}^{(\boldsymbol{q})} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}) ; \quad \mathrm{r}_{i}\left(U_{j}^{(\boldsymbol{q})}\right)=\boldsymbol{w}_{j, i}-c_{j, i}^{-1} \boldsymbol{q} M_{j}^{-1} \quad\left(i=1, \ldots, r_{j}\right) . \tag{3.15}
\end{equation*}
$$

We also write

$$
\begin{equation*}
U^{(\boldsymbol{q})}:=\left(U_{1}^{(\boldsymbol{q})}, \cdots, U_{N}^{(\boldsymbol{q})}\right) \in \prod_{j=1}^{N} \mathrm{M}_{r_{j} \times d}(\mathbb{R}), \tag{3.16}
\end{equation*}
$$

and we introduce the notation (cf. (3.11))

$$
\begin{equation*}
\mathrm{I}_{V}:=\left(\mathrm{I}_{V_{1}}, \ldots, \mathrm{I}_{V_{N}}\right) \in G \quad \text { for any } V=\left(V_{1}, \ldots, V_{N}\right) \in \prod_{j=1}^{N} \mathrm{M}_{r_{j} \times d}(\mathbb{R}) \tag{3.17}
\end{equation*}
$$

Furthermore, we set

$$
\begin{equation*}
\widetilde{M}:=\left(M_{1}, \ldots, M_{N}\right) \in G, \tag{3.18}
\end{equation*}
$$

and define:

$$
\begin{equation*}
g_{0}^{(\boldsymbol{q})}:=\mathrm{I}_{U^{(q)}} \widetilde{M} \in G \quad\left(\boldsymbol{q} \in \mathbb{R}^{d}\right) . \tag{3.19}
\end{equation*}
$$

The point of these definitions is that we now have, for any $\psi=(j, i) \in \Psi$ :

$$
\begin{equation*}
c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}\left(g_{0}^{(\boldsymbol{q})}\right)\right)=c_{\psi}\left(\mathbb{Z}^{d} \mathrm{a}_{i}\left(\mathrm{I}_{U_{j}^{(\boldsymbol{q})}} M_{j}\right)\right)=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}-c_{\psi}^{-1} \boldsymbol{q} M_{j}^{-1}\right) M_{j}=\mathcal{L}_{\psi}-\boldsymbol{q} \tag{3.20}
\end{equation*}
$$

(cf. (3.4)), and hence

$$
\begin{equation*}
J_{0}\left(\Gamma g_{0}^{(\boldsymbol{q})}\right)=\mathcal{P}-\boldsymbol{q} \quad\left(\forall \boldsymbol{q} \in \mathbb{R}^{d}\right) \tag{3.21}
\end{equation*}
$$

Finally we introduce some further notation relating to the structure of the homogeneous spaces $\mathbb{X}_{j}$ and $\mathbb{X}$. Recall that for each $j \in\{1, \ldots, N\}$ we have a projection map $\iota: \mathrm{S}_{r_{j}}(\mathbb{R}) \rightarrow$ $\mathrm{SL}_{d}(\mathbb{R})$; this map takes $\Gamma_{j}$ to $\mathrm{SL}_{d}(\mathbb{Z})$, and hence induces a projection map from $\mathbb{X}_{j}$ to $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$, which we will denote by $\tilde{\iota}$ :

$$
\begin{equation*}
\tilde{\iota}: \mathbb{X}_{j} \rightarrow \mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R}), \quad \widetilde{\iota}\left(\Gamma_{j} g\right)=\mathrm{SL}_{d}(\mathbb{Z}) \iota(g) \quad\left(g \in G_{j}\right) \tag{3.22}
\end{equation*}
$$

We also set

$$
\mathbb{T}_{j}:=\mathbb{R}^{r_{j}} / \mathbb{Z}^{r_{j}} \quad \text { and } \quad \mathbb{T}_{j}^{d}:=\underbrace{\mathbb{T}_{j} \times \cdots \times \mathbb{T}_{j}}_{d \text { copies }}
$$

Note that $\mathbb{T}_{j}^{d}$ is an $r_{j} d$-dimensional torus. We will use " $\pi$ " to denote both the corresponding projection maps;

$$
\begin{equation*}
\pi: \mathbb{R}^{r_{j}} \rightarrow \mathbb{T}_{j} \quad \text { and } \quad \pi:\left(\mathbb{R}^{r_{j}}\right)^{d} \rightarrow \mathbb{T}_{j}^{d} \tag{3.23}
\end{equation*}
$$

We will identify $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with $\left(\mathbb{R}^{r_{j}}\right)^{d}$, by identifying any matrix in $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with the sequence of its $d$ column vectors (in order). This also induces an identification

$$
\mathbb{T}_{j}^{d}=\mathrm{M}_{r_{j} \times d}(\mathbb{R} / \mathbb{Z})
$$

We write $\widetilde{\mathrm{r}}_{i}: \mathbb{T}_{j}^{d} \rightarrow(\mathbb{R} / \mathbb{Z})^{d}$ for the projection induced by the map $\mathrm{r}_{i}$. For each $j$ we also introduce the embedding

$$
\begin{equation*}
x: \mathbb{T}_{j}^{d} \rightarrow \mathbb{X}_{j} ; \quad x(\pi(U))=\Gamma_{j} \mathrm{I}_{U}, \quad \forall U \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}) \tag{3.24}
\end{equation*}
$$

## 4. Spherical equidistribution (without uniformity)

4.1. The limit measure $\mu^{(\boldsymbol{q})}$. In this section we state a result in homogeneous dynamics, Theorem 4.2 below, which gives a precise description of the limit considered in the spherical equidistribution condition [P2] in Section 2. Theorem 4.2 expresses the answer in terms of the parametrizing space $\mathbb{X}$, and in fact applies for an arbitrary point $\boldsymbol{q}$ in $\mathbb{R}^{d}$; however the theorem includes no statement about uniformity with respect to the translation vector $\boldsymbol{q}$, as is required in [P2].

This Theorem 4.2 is a special case of more general equidistribution results which we will prove in Sections 8 and 9 (see in particular Theorem 9.1), and the main reason for stating this special case already here is to help motivating the introduction of the measure $\mu^{(\boldsymbol{q})}$ (see (4.2) below), which will be a crucial object of study in the remainder of Section 4 as well as in Section 5. Note that we will not make use of the statement of Theorem 4.2 before Section 9 ,

For a given $\boldsymbol{q} \in \mathbb{R}^{d}$ and each $j \in\{1, \ldots, N\}$, we let $U_{j, 1}^{(\boldsymbol{q})}, \ldots, U_{j, d}^{(\boldsymbol{q})} \in \mathbb{R}^{r_{j}}$ be the column vectors of $U_{j}^{(\boldsymbol{q})}$, and then define $L_{j}^{(\boldsymbol{q})}$ to be the identity component of the smallest closed subgroup of $\mathbb{R}^{r_{j}}$ containing both $\mathbb{Z}^{r_{j}}$ and $U_{j, 1}^{(\boldsymbol{q})}, \ldots, U_{j, d}^{(\boldsymbol{q})}$. This $L_{j}^{(\boldsymbol{q})}$ is a rational (linear) subspace of $\mathbb{R}^{r_{j}}$, i.e., $L_{j}^{(\boldsymbol{q})} \cap \mathbb{Z}^{r_{j}}$ is a lattice in $L_{j}^{(\boldsymbol{q})}$. Next, given any linear subspace $L$ of $\mathbb{R}^{r_{j}}$, we let $\mathrm{S}_{L}(\mathbb{R})$ be the closed connected subgroup of $\mathrm{S}_{r_{j}}(\mathbb{R})$ given by

$$
\begin{equation*}
\mathrm{S}_{L}(\mathbb{R}):=\mathrm{SL}_{d}(\mathbb{R}) \ltimes L^{d}=\left\{(M, U) \in \mathrm{S}_{r_{j}}(\mathbb{R}): U \in L^{d}\right\} \tag{4.1}
\end{equation*}
$$

Here " $L$ " should be understood via our identification $\mathrm{M}_{r_{j} \times d}(\mathbb{R})=\left(\mathbb{R}^{r_{j}}\right)^{d}$, viz., $L^{d}$ is the set of matrices in $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ all of whose column vectors lie in $L$.

Lemma 4.1. For any $\boldsymbol{q} \in \mathbb{R}^{d}$, the $\mathrm{S}_{L_{j}^{(\boldsymbol{q})}}(\mathbb{R})$-orbit of the point $\Gamma_{j} \mathrm{I}_{U_{j}^{(\boldsymbol{q})}}$ in $\mathbb{X}_{j}=\Gamma_{j} \backslash G_{j}$ is a closed embedded submanifold of $\mathbb{X}_{j}$ which carries a unique $\mathrm{S}_{L_{j}^{(q)}}(\mathbb{R})$-invariant probability measure.
Proof. Set $U:=U_{j}^{(\boldsymbol{q})}$ and $L:=L_{j}^{(\boldsymbol{q})}$. It follows from the construction of $L$ that there exists a matrix $X \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ such that $U-X \in L^{d}$ (indeed, see Lemma 4.3 below). This implies $\mathrm{I}_{U} \mathrm{I}_{X}^{-1} \in \mathrm{~S}_{L}(\mathbb{R})$, and hence the $\mathrm{S}_{L}(\mathbb{R})$-orbit in the statement of the lemma can be expressed as $\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{U} \mathrm{~S}_{L}(\mathbb{R})=\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R})$. Hence by [24, Theorem 1.13], it suffices to verify that $\Gamma_{j}$ intersects $\mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R}) \mathrm{I}_{X}^{-1}$ in a lattice. To show this, let $n$ be a denominator of $X$, i.e. a positive integer such that $X \in n^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$, and then set

$$
\Lambda:=\left\{(M, V): M \in \Gamma(n), V \in X M-X+\left(L \cap \mathbb{Z}^{r_{j}}\right)^{d}\right\}
$$

where $\Gamma(n)$ is the principal congruence subgroup of $\mathrm{SL}_{d}(\mathbb{Z})$ of level $n$. One verifies that $\Lambda$ is a subgroup of $\Gamma_{j} \cap \mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R}) \mathrm{I}_{X}^{-1}$, and a lattice in $\mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R}) \mathrm{I}_{X}^{-1}$. Hence the lemma is proved.

In view of Lemma 4.1, it makes sense to define $\mu_{j}^{(\boldsymbol{q})} \in P\left(\mathbb{X}_{j}\right)$ to be the unique $S_{L_{j}^{(\boldsymbol{q})}}(\mathbb{R})$ invariant probability measure on the orbit $\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{U_{j}^{(\boldsymbol{q})}} \mathrm{S}_{L_{j}^{(\boldsymbol{q})}}(\mathbb{R})$ in $\mathbb{X}_{j}$. Finally we set

$$
\begin{equation*}
\mu^{(\boldsymbol{q})}:=\mu_{1}^{(\boldsymbol{q})} \otimes \cdots \otimes \mu_{N}^{(\boldsymbol{q})} \in P(\mathbb{X}) \tag{4.2}
\end{equation*}
$$

Let $\varphi: \mathrm{SL}_{d}(\mathbb{R}) \rightarrow G$ be the diagonal embedding. For any topological space $S$, we denote by $\mathrm{C}_{b}(S)$ the space of bounded continuous real-valued functions on $S$.

Theorem 4.2. Given any $\boldsymbol{q} \in \mathbb{R}^{d}, f \in \mathrm{C}_{b}(\mathbb{X})$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\int_{S_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\mathbb{X}} f d \mu^{(\boldsymbol{q})} \tag{4.3}
\end{equation*}
$$

as $\rho \rightarrow 0$.
We will prove Theorem 4.2 at the end of Section 9.1, as a corollary of more general and difficult results which we prove in Section 8. We remark that it is possible to give a considerably simpler deduction of Theorem 4.2 by combining the ideas in the proof of [20, Theorem 10] and the proof of [21, Lemma 5.22]. In any case, the crucial ingredient in either of these two proofs is the deep measure classification theorem of Ratner [25], and strong use is also made of the work of Shah [27].
4.2. Re-expressing the limit measures $\mu_{j}^{(\boldsymbol{q})}$ as $\overline{\omega_{j}^{(\boldsymbol{q})}}$. The task of finding the appropriate space of marks $\Sigma$ and marking $\varsigma: \mathcal{P} \rightarrow \Sigma$ to satisfy the conditions in Section 2 is closely related to the task of understanding the limit measures $\mu^{(\boldsymbol{q})}=\mu_{1}^{(\boldsymbol{q})} \otimes \cdots \otimes \mu_{N}^{(\boldsymbol{q})}$ appearing in Theorem 4.2, and in particular how these vary as $\boldsymbol{q}$ varies within the set $\mathcal{P}$. As a first step towards this goal, in this section we will introduce an alternative way to define the measures $\mu_{j}^{(\boldsymbol{q})}$. Recall that each space $\mathbb{X}_{j}$ is a torus bundle over the space $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})(\mathrm{cf}$. (3.22) $)$, and the point here is that each each measure $\mu_{j}^{(\boldsymbol{q})}$ can be expressed as a product measure, of the normalized Haar measure on $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ and some fixed measure on the fiber $\mathbb{T}_{j}^{d}$. We make this precise in Lemma 4.4 and Proposition 4.7 below.

We keep $j \in\{1, \ldots, N\}$ throughout this section. Our first step will be to introduce the relevant measures which can appear on the fiber $\mathbb{T}_{j}^{d}$.

Recall that $\pi$ denotes (among other things) the projection map $\mathbb{R}^{r_{j}} \rightarrow \mathbb{T}^{r_{j}}:=\mathbb{R}^{r_{j}} / \mathbb{Z}^{r_{j}}$. For any subset $S$ of $\mathbb{R}^{r_{j}}$ (resp. $\mathbb{T}^{r_{j}}$ ), we write $\langle S\rangle$ for the subgroup of $\mathbb{R}^{r_{j}}$ (resp. $\mathbb{T}^{r_{j}}$ ) generated by $S$, and $\overline{\langle S\rangle}$ for its closure. For any topological group $H$ we write $H^{\circ}$ for the connected component of the identity. Now for any non-empty subset $S$ of $\mathbb{T}^{r_{j}}$, we introduce the notation:

$$
\begin{equation*}
\mathfrak{L}(S):={\overline{\left\langle\pi^{-1}(S)\right\rangle}}^{\circ} \tag{4.4}
\end{equation*}
$$

This is a rational subspace of $\mathbb{R}^{r_{j}}$, i.e. a linear subspace of $\mathbb{R}^{r_{j}}$ which is spanned by its vectors in $\mathbb{Z}^{r_{j}}$. If $S$ is a non-empty subset of $\mathbb{R}^{r_{j}}$ then we also write $\mathfrak{L}(S):=\mathfrak{L}(\pi(S))=\overline{\left\langle S+\mathbb{Z}^{r_{j}}\right\rangle}{ }^{\circ}$, and if $V_{1}, \ldots, V_{m}$ is any finite sequence of points in $\mathbb{T}^{r_{j}}$ or $\mathbb{R}^{r_{j}}$, then we also write $\mathfrak{L}\left(V_{1}, \ldots, V_{m}\right):=$ $\mathfrak{L}\left(\left\{V_{1}, \ldots, V_{m}\right\}\right)$.

Lemma 4.3. For any non-empty subset $S \subset \mathbb{R}^{r_{j}}, \mathfrak{L}(S)$ is the unique smallest rational subspace $L \subset \mathbb{R}^{r_{j}}$ with the property that there is some $n \in \mathbb{Z}^{+}$such that $\langle S\rangle \subset n^{-1} \mathbb{Z}^{r_{j}}+L$.

Proof. Set $G=\overline{\langle\pi(S)\rangle}$; this is a closed subgroup of $\mathbb{T}^{r_{j}}$; hence its group of components, $G / G^{\circ}$, is finite, i.e. there exists a positive integer $n$ such that $n \boldsymbol{y} \in G^{\circ}$ for all $\boldsymbol{y} \in G$. We also have $G^{\circ}=\pi(\mathfrak{L}(S))$. It follows that $n \cdot \pi(S) \subset \pi(\mathfrak{L}(S))$, i.e. $n \cdot S \subset \mathbb{Z}^{r_{j}}+\mathfrak{L}(S)$. This is equivalent with $S \subset n^{-1} \mathbb{Z}^{r_{j}}+\mathfrak{L}(S)$, and also with $\langle S\rangle \subset n^{-1} \mathbb{Z}^{r_{j}}+\mathfrak{L}(S)$.

Next assume that $L$ is any rational subspace of $\mathbb{R}^{r_{j}}$ satisfying $\langle S\rangle \subset n^{-1} \mathbb{Z}^{r_{j}}+L$ for some $n \in \mathbb{Z}^{+}$. Now $n^{-1} \mathbb{Z}^{r_{j}}+L$ is a closed subgroup of $\mathbb{R}^{r_{j}}$ which contains $S+\mathbb{Z}^{r_{j}}$; hence $\overline{\left\langle S+\mathbb{Z}^{r_{j}}\right\rangle} \subset$ $n^{-1} \mathbb{Z}^{r_{j}}+L$, and so $\mathfrak{L}(S):=\overline{\left\langle S+\mathbb{Z}^{r_{j}}\right\rangle}{ }^{\circ} \subset\left(n^{-1} \mathbb{Z}^{r_{j}}+L\right)^{\circ}=L$.

Next, for any $V=\left(V_{1}, \ldots, V_{d}\right) \in \mathbb{T}_{j}^{d}$, we introduce the notation

$$
\begin{equation*}
\mathbb{S}_{j}^{(V)}:=\overline{\left\langle V_{1}, \ldots, V_{d}\right\rangle} \tag{4.5}
\end{equation*}
$$

and

$$
\begin{equation*}
L_{j}^{(V)}:=\mathfrak{L}\left(V_{1}, \ldots, V_{d}\right)=\left(\pi^{-1}\left(\mathbb{S}_{j}^{(V)}\right)\right)^{\circ} \tag{4.6}
\end{equation*}
$$

Thus $\mathbb{S}_{j}^{(V)}$ is a closed subgroup of $\mathbb{T}_{j}$ and $L_{j}^{(V)}$ is a rational subspace of $\mathbb{R}^{r_{j}}$.
Recall that we have identified $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with $\left(\mathbb{R}^{r_{j}}\right)^{d}$; this means that for any linear subspace $L \subset \mathbb{R}^{r_{j}}, L^{d}=L \times \cdots \times L$ is a linear subspace of $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$. Similarly for any subgroup $\mathbb{S} \subset \mathbb{T}_{j}, \mathbb{S}^{d}$ is a subgroup of $\mathbb{T}_{j}^{d}=\mathrm{M}_{r_{j} \times d}(\mathbb{R}) / \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$. Note that $\mathrm{SL}_{d}(\mathbb{R})$ acts from the right on $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ by matrix multiplication, and this action preserves the subspace $L^{d}$ for any $L \subset \mathbb{R}^{r_{j}}$. Furthermore, the subgroup $\mathrm{SL}_{d}(\mathbb{Z})$ preserves $\mathrm{M}_{r_{j} \times d}(\mathbb{Z})$, and hence we obtain an induced right action of $\mathrm{SL}_{d}(\mathbb{Z})$ on $\mathbb{T}_{j}^{d}=\mathrm{M}_{r_{j} \times d}(\mathbb{R} / \mathbb{Z})$. This action preserves the subgroup $\mathbb{S}^{d}$, for any subgroup $\mathbb{S} \subset \mathbb{T}_{j}$.

Now for any $V \in \mathbb{T}_{j}^{d}$, we define:

$$
\begin{equation*}
\mathcal{O}_{j}^{(V)}=\bigcup_{\gamma \in \mathrm{SL}_{d}(\mathbb{Z})}\left(V \gamma+\left(\mathbb{S}_{j}^{(V)}\right)^{\circ d}\right) \tag{4.7}
\end{equation*}
$$

Here it should be noted that $V \in\left(\mathbb{S}_{j}^{(V)}\right)^{d}$; hence $\mathcal{O}_{j}^{(V)}$ is a union of some of the connected components of the Lie group $\left(\mathbb{S}_{j}^{(V)}\right)^{d}$. In particular $\mathcal{O}_{j}^{(V)}$ is open and closed in $\left(\mathbb{S}_{j}^{(V)}\right)^{d}$. Note also that since $\left(\mathbb{S}_{j}^{(V)}\right)^{d}$ is compact, its total number of components is finite.

We also define $\omega_{j}^{(V)} \in P\left(\mathbb{T}_{j}^{d}\right)$ to be the restriction to $\mathcal{O}_{j}^{(V)}$ of the Haar measure on $\left(\mathbb{S}_{j}^{(V)}\right)^{d}$. normalized so that $\omega_{j}^{(V)}\left(\mathcal{O}_{j}^{(V)}\right)=1$.

Note that for each $V \in \mathbb{T}_{j}^{d}$, the measure $\omega_{j}^{(V)}$ is $\mathrm{SL}_{d}(\mathbb{Z})$-invariant by construction. We denote by $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ the subset of all $\mathrm{SL}_{d}(\mathbb{Z})$-invariant measures $\omega \in P\left(\mathbb{T}_{j}^{d}\right)$. This is a closed, and hence compact, subset of $P\left(\mathbb{T}_{j}^{d}\right)$.

Next, recall that for any $\boldsymbol{q} \in \mathbb{R}^{d}$ we have defined a matrix $U_{j}^{(\boldsymbol{q})} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$, cf. (3.15); and via our identification $\mathbb{T}_{j}^{d}=\mathrm{M}_{r_{j} \times d}(\mathbb{R}) / \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ we have $\pi\left(U_{j}^{(\boldsymbol{q})}\right) \in \mathbb{T}_{j}^{d}$. Note that the rational subspace $L_{j}^{(\boldsymbol{q})} \subset \mathbb{R}^{r_{j}}$ defined in Section 4.1 equals $L_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$ in our present notation. From now on we will also write $\mathbb{S}_{j}^{(\boldsymbol{q})}, \mathcal{O}_{j}^{(\boldsymbol{q})}$ and $\omega_{j}^{(\boldsymbol{q})}$ to denote $\mathbb{S}_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}, \mathcal{O}_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$ and $\omega_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$, respectively. Thus in particular, $\mathcal{O}_{j}^{(\boldsymbol{q})}$ is a union of some of the components of the closed subgroup $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{d} \subset \mathbb{T}_{j}^{d}$, and $\omega_{j}^{(\boldsymbol{q})}$ is a measure in $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ supported on $\mathcal{O}_{j}^{(\boldsymbol{q})}$. We will see that $\omega_{j}^{(\boldsymbol{q})}$ is the measure on the fiber $\mathbb{T}_{j}^{d}$ which gives back the measure $\mu_{j}^{(\boldsymbol{q})}$ on $\mathbb{X}_{j}$ via the product construction which we will presently describe.

We now introduce a map

$$
\begin{equation*}
P\left(\mathbb{T}_{j}^{d}\right)^{\prime} \rightarrow P\left(\mathbb{X}_{j}\right), \quad \omega \mapsto \bar{\omega} \tag{4.8}
\end{equation*}
$$

as follows. For any $\omega \in P\left(\mathbb{T}_{j}^{d}\right)$, let $\widetilde{\omega}$ be corresponding $\mathrm{M}_{r_{j} \times d}(\mathbb{Z})$-invariant Borel measure on $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$, and define the Borel measure $\widetilde{\bar{\omega}}$ on $G_{j}$ through

$$
\begin{equation*}
d \widetilde{\bar{\omega}}(g)=d \widetilde{\omega}(U) d \nu(A) \quad \text { when } g=\mathrm{I}_{U} A \in G_{j} \quad\left(U \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}), A \in \mathrm{SL}_{d}(\mathbb{R})\right) \tag{4.9}
\end{equation*}
$$

where $\nu$ denotes the Haar measure on $\mathrm{SL}_{d}(\mathbb{R})$, normalized by $\nu\left(\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})\right)=1$. By the following lemma, if $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ then the measure $\widetilde{\bar{\omega}}$ is left $\Gamma_{j}$-invariant, and we finally define $\bar{\omega}$ to be the corresponding Borel measure on $\mathbb{X}_{j}$.

Lemma 4.4. For any $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$, the measure $\widetilde{\bar{\omega}}$ is left $\Gamma_{j}$-invariant, and the corresponding measure $\bar{\omega}$ on $\mathbb{X}_{j}$ satisfies, for any Borel set $E \subset \mathbb{X}_{j}$,

$$
\begin{equation*}
\bar{\omega}(E)=\int_{F_{d}} \int_{\mathbb{T}_{j}^{d}} I(x(U) A \in E) d \omega(U) d \nu(A) \tag{4.10}
\end{equation*}
$$

where $F_{d}$ is any fixed Borel set in $\mathrm{SL}_{d}(\mathbb{R})$ which is a fundamental domain for $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$, and where $I(\cdot)$ is the indicator function.

Note that (4.10) in particular shows that $\bar{\omega}$ is a probability measure, i.e. $\bar{\omega} \in P\left(\mathbb{X}_{j}\right)$, and so we indeed have a map as in (4.8).

Proof. Let $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$. In order to verify that $\widetilde{\bar{\omega}}$ is left $\Gamma_{j}$-invariant, it suffices to verify that for any Borel set $E \subset G_{j}$ we have $\widetilde{\bar{\omega}}\left(\mathrm{I}_{M} E\right)=\widetilde{\bar{\omega}}(E)$ for all $M \in \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ and $\widetilde{\omega}(\gamma E)=\widetilde{\bar{\omega}}(E)$ for all $\gamma \in \mathrm{SL}_{d}(\mathbb{Z})$. The first of these two relations is immediate from (4.9) and the fact that $\widetilde{\omega}$ is $\mathrm{M}_{r_{j} \times d}(\mathbb{Z})$-invariant. The second relation follows by noticing that, for any $U \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ and $A \in \mathrm{SL}_{d}(\mathbb{R}), \mathrm{I}_{U} A \in \gamma E$ holds if and only if $\gamma^{-1} \mathrm{I}_{U} A \in E$, viz., $\mathrm{I}_{U \gamma} \gamma^{-1} A \in E$, and then using the invariance of the Haar measure $\nu$, and the fact that $\widetilde{\omega}$ is $\mathrm{SL}_{d}(\mathbb{Z})$-invariant, since $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$.

Next, in order to verify (4.10), note that if $C_{d}$ is the set of matrices in $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ all of whose entries lie in $[0,1)$, then the set

$$
F_{d}^{\prime}:=\left\{\mathrm{I}_{U} A: U \in C_{d}, A \in F_{d}\right\} \quad \subset G_{j}
$$

is a fundamental domain for $\Gamma_{j} \backslash G_{j}$. Hence for any Borel set $E \subset \mathbb{X}_{j}$ we have

$$
\bar{\omega}(E)=\widetilde{\omega}\left(\pi^{-1}(E) \cap F_{d}^{\prime}\right)=\int_{F_{d}} \int_{C_{d}} I\left(\mathrm{I}_{U} A \in \pi^{-1}(E)\right) d \widetilde{\omega}(U) d \nu(A)
$$

which is the same as (4.10).
For later use we record a few simple properties of the map $\omega \mapsto \bar{\omega}$ just defined.
Lemma 4.5. For any $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$, the measure $\bar{\omega}$ is $\mathrm{SL}_{d}(\mathbb{R})$ invariant, and $\widetilde{\iota}_{*} \bar{\omega}=\nu$.
(Here, by abuse of notation, we write $\nu$ also for the measure on $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ induced by the Haar measure $\nu$.)

Proof. It is obvious from (4.9) that $\widetilde{\bar{\omega}}$ is right $\mathrm{SL}_{d}(\mathbb{R})$ invariant; hence also $\bar{\omega}$ is $\mathrm{SL}_{d}(\mathbb{R})$ invariant. The fact that $\widetilde{\iota}_{*} \bar{\omega}=\nu$ is immediate from (4.10).

Lemma 4.6. The map in (4.8) is continuous. (Here $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ is equipped with the subspace topology from $P\left(\mathbb{T}_{j}^{d}\right)$.)
Proof. Recall that $P(S)$ is metrizable for any separable and metrizable topological space $S$ [1, pp. 72-73]; in particular $P\left(\mathbb{T}_{j}^{d}\right)$ and $P\left(\mathbb{X}_{j}\right)$ are metrizable. Hence it suffices to prove that if $\left(\omega_{k}\right)$ is any sequence in $P\left(\mathbb{T}_{j}^{d}\right)$ converging to $\omega \in P\left(\mathbb{T}_{j}^{d}\right)$, we have $\bar{\omega}_{k} \rightarrow \bar{\omega}$ in $P\left(\mathbb{X}_{j}\right)$. To prove this, we have to prove that $\bar{\omega}_{k}(\varphi) \rightarrow \bar{\omega}(\varphi)$ for any $\varphi \in \mathrm{C}_{b}\left(\mathbb{X}_{j}\right)$; but it is a well-known fact that it suffices to prove this for $\varphi \in \mathrm{C}_{c}\left(\mathbb{X}_{j}\right)$. Thus for a fixed $\varphi \in \mathrm{C}_{c}\left(\mathbb{X}_{j}\right)$, our task is to prove that

$$
\int_{\mathbb{T}_{j}^{d}} \int_{F_{d}} \varphi(x(U) A) d \nu(A) d \omega_{k}(U)
$$

converges to the corresponding integral with $\omega$, as $k \rightarrow \infty$. The fact that $\varphi$ has compact support in $\mathbb{X}_{j}$ implies that $\alpha(U):=\int_{F_{d}} \varphi(x(U) A) d \nu(A)$ is a continuous function on $\mathbb{T}_{j}^{d}$.

Now the task is to prove that $\int_{\mathbb{T}_{j}^{d}} \alpha d \omega_{k} \rightarrow \int_{\mathbb{T}_{j}^{d}} \alpha d \omega$, and this holds by definition since $\omega_{k} \rightarrow \omega$ in $P\left(\mathbb{T}_{j}^{d}\right)$.

Finally, for any $V \in \mathbb{T}_{j}^{d}$ we will now identify the measure $\overline{\omega^{(V)}} \in P\left(\mathbb{X}_{j}\right)$ as an invariant measure on a certain homogeneous submanifold of $\mathbb{X}_{j}$. Given $V \in \mathbb{T}_{j}^{d}$, it follows from Lemma 4.3 that there exists some $X \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ such that $V-\pi(X) \in\left(\mathbb{S}_{j}^{(V)}\right)^{\circ d}$, and this means that we can choose some $\widetilde{V} \in X+\left(L_{j}^{(V)}\right)^{d} \subset \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with $\pi(\widetilde{V})=V$. Let $L=L_{j}^{(V)}$ and recall the definition of $\mathrm{S}_{L}(\mathbb{R})$ in (4.1). It follows that the $\mathrm{S}_{L}(\mathbb{R})$-orbit of the point $x(V)$ in $\mathbb{X}_{j}$ can be expressed as:

$$
\begin{equation*}
x(V) \cdot \mathrm{S}_{L}(\mathbb{R})=\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{\widetilde{V}} \mathrm{~S}_{L}(\mathbb{R})=\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R}) \tag{4.11}
\end{equation*}
$$

By the proof of Lemma 4.1 (applied to $\widetilde{V}$ in place of $U_{j}^{(\boldsymbol{q})}$ ), the orbit in (4.11) is a closed embedded submanifold of $\mathbb{X}_{j}$ which carries a unique $S_{L}(\mathbb{R})$-invariant probability measure.
Proposition 4.7. For any $V \in \mathbb{T}_{j}^{d}$, the unique $\mathrm{S}_{L_{j}^{(V)}}(\mathbb{R})$-invariant probability measure on the orbit $x(V) \cdot \mathrm{S}_{L_{j}^{(V)}}(\mathbb{R})$ equals $\overline{\omega_{j}^{(V)}}$. In particular, for any $\boldsymbol{q} \in \mathbb{R}^{d}$ we have $\mu_{j}^{(\boldsymbol{q})}=\overline{\omega_{j}^{(\boldsymbol{q})}}$.

Proof. Note that the second part of the proposition is an immediate consequence of the first part, since, by the definition in Section 4.1, $\mu_{j}^{(\boldsymbol{q})}$ is the $\mathrm{S}_{L_{j}^{(\boldsymbol{q})}}(\mathbb{R})$-invariant probability measure on the orbit $x\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right) \cdot \mathrm{S}_{L_{j}^{(\boldsymbol{q})}}(\mathbb{R})$ in $\mathbb{X}_{j}$, and furthermore we have $L_{j}^{(\boldsymbol{q})}=L_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$ and $\omega_{j}^{(\boldsymbol{q})}=$ $\omega_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$.

To prove the first part of the proposition, set $\omega:=\omega_{j}^{(V)}$, and let $X, L$ and $\widetilde{V}$ be as in (4.11). Since $\bar{\omega}=\overline{\omega_{j}^{(V)}}$ by definition is the probability measure on $\mathbb{X}_{j}$ which corresponds to the $\Gamma_{j}$-invariant measure $\widetilde{\bar{\omega}}$ on $G_{j}$ given by (4.9), it suffices to verify that $\widetilde{\bar{\omega}}$ is supported on $\Gamma_{j} \mathrm{I}_{\widetilde{V}} \mathrm{~S}_{L}(\mathbb{R})$ and that $\widetilde{\bar{\omega}}$ is right $\mathrm{S}_{L}(\mathbb{R})$ invariant. However, it follows from (4.9) that

$$
\begin{equation*}
\operatorname{supp}(\widetilde{\bar{\omega}})=\left\{\mathrm{I}_{U} A: U \in \pi^{-1}\left(\mathcal{O}_{j}(V)\right), A \in \mathrm{SL}_{d}(\mathbb{R})\right\} \tag{4.12}
\end{equation*}
$$

and here we have

$$
\begin{equation*}
\pi^{-1}\left(\mathcal{O}_{j}(V)\right)=\left\{(\widetilde{V}+W) \gamma: \gamma \in \mathrm{SL}_{d}(\mathbb{Z}), W \in L^{d}+\mathrm{M}_{r_{j} \times d}(\mathbb{Z})\right\} \tag{4.13}
\end{equation*}
$$

(cf. (4.7)). Note here that both $L^{d}$ and $\mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ are $\mathrm{SL}_{d}(\mathbb{Z})$-invariant subsets of $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$. Using (4.12), (4.13) and $\mathrm{I}_{(\widetilde{V}+W) \gamma} A=\gamma^{-1} \mathrm{I}_{\widetilde{V}+W} \gamma A$ it follows that $\operatorname{supp}(\widetilde{\bar{\omega}})=\Gamma_{j} \mathrm{I}_{\widetilde{V}} \mathrm{~S}_{L}(\mathbb{R})$, as desired.

Finally we verify that $\widetilde{\bar{\omega}}$ is right $\mathrm{S}_{L}(\mathbb{R})$-invariant. We have noted that $\widetilde{\bar{\omega}}$ is right $\mathrm{SL}_{d}(\mathbb{R})$ invariant (cf. Lemma 4.5); hence it suffices to verify that $\widetilde{\bar{\omega}}$ is right $\mathrm{I}_{W}$-invariant for every $W \in L^{d}$. However this also follows from (4.9), by noticing that $\mathrm{I}_{U} A \mathrm{I}_{W}=\mathrm{I}_{U+W A^{-1}} A$ for all $U \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ and $A \in \mathrm{SL}_{d}(\mathbb{R})$, and using $W A^{-1} \in L^{d}$ and the fact that $\widetilde{\omega}$ is invariant under $L^{d}$-translations (since $\omega=\omega_{j}^{(V)}$ is invariant under $\left(\mathbb{S}_{j}(V)^{\circ}\right)^{d}$-translations).
4.3. A Siegel integration formula for the measure $\bar{\omega}$. The classical Siegel integration formula [28] states that for any measurable function $f: \mathbb{R}^{d} \rightarrow \mathbb{R}_{\geq 0}$ (or for any $f \in \mathrm{C}_{c}\left(\mathbb{R}^{d}\right)$ ), defining the Siegel transform $\widehat{f}$ on $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ by

$$
\widehat{f}\left(\mathrm{SL}_{d}(\mathbb{Z}) g\right)=\sum_{\boldsymbol{m} \in \mathbb{Z}^{d}} f(\boldsymbol{m} g) \quad\left(g \in \mathrm{SL}_{d}(\mathbb{R})\right)
$$

then

$$
\int_{\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})} \widehat{f} d \nu=\int_{\mathbb{R}^{d}} f(\boldsymbol{x}) d \boldsymbol{x}+f(\mathbf{0}) .
$$

The following proposition gives an analogous formula involving the measure $\bar{\omega}$, for any given $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$. This formula will be used later in our proof of $[\mathrm{P} 2]$ (uniform spherical equidistribution); see the proof of Lemma 7.8 below.

Proposition 4.8. Let $\psi=(j, i) \in \Psi$. For any Borel measurable function $f \in \mathbb{R}^{d} \rightarrow \mathbb{R}_{\geq 0}$, we define its " $\psi$-Siegel transform" $\widehat{f}_{\psi}: \mathbb{X}_{j} \rightarrow \mathbb{R}_{\geq 0} \cup\{+\infty\}$ by

$$
\begin{equation*}
\widehat{f}_{\psi}\left(\Gamma_{j} g\right)=\sum_{\boldsymbol{m} \in \mathbb{Z}^{d}} f\left(c_{\psi} \cdot\left(\boldsymbol{m} \mathrm{a}_{i}(g)\right)\right) \quad\left(g \in G_{j}\right) \tag{4.14}
\end{equation*}
$$

Then for any $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ we have

$$
\begin{equation*}
\int_{\mathbb{X}_{j}} \widehat{f}_{\psi} d \bar{\omega}=\bar{n}_{\psi} \int_{\mathbb{R}^{d}} f(\boldsymbol{x}) d \boldsymbol{x}+\omega\left(\widetilde{r}_{i}^{-1}(\{\mathbf{0}\})\right) f(\mathbf{0}) \tag{4.15}
\end{equation*}
$$

as an equality of extended real numbers in $\mathbb{R}_{\geq 0} \cup\{+\infty\}$.
Remark 4.1. The reason for the formula (4.14) is that this corresponds to the " $\psi$-part" of the union in (3.14). Note that $\widehat{f_{\psi}}$ is well-defined, in the sense that the right hand side of (4.14) remains the same if $g$ is replaced by $\gamma g$ for any $\gamma \in \Gamma_{j}$.

Proof. Let $\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ be given. Define the Borel measure $\mu$ on $\mathbb{R}^{d}$ by $\mu(B):=\int_{\mathbb{X}_{j}} \widehat{\left(\chi_{B}\right)_{\psi}} d \bar{\omega}$ for any Borel set $B \subset \mathbb{R}^{d}$. Let us prove that $\mu$ is finite on compact sets, so that $\mu$ is a Radon measure. For this, it suffices to prove that $\mu(B)<\infty$ for any ball $B=\mathcal{B}_{R}^{d}$. In this case we have, for all $g \in G_{j}$ :

$$
\begin{equation*}
\widehat{\left(\chi_{B}\right)_{\psi}}\left(\Gamma_{j} g\right)=\#\left(B \cap c_{\psi}\left(\mathbb{Z}^{d} \mathrm{a}_{i}(g)\right)\right) \ll \#\left(B \cap c_{\psi} \mathbb{Z}^{d} \iota(g)\right) \tag{4.16}
\end{equation*}
$$

by [10, Proposition 5] (the implied constant in the last bound depends only on $d$ ). Furthermore,

$$
\int_{\mathbb{X}_{j}} \#\left(B \cap c_{\psi} \mathbb{Z}^{d} \iota(g)\right) d \bar{\omega}(g)=\int_{\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})} \#\left(B \cap c_{\psi} \mathbb{Z}^{d} g\right) d \nu(g)=1+\operatorname{vol}\left(c_{\psi}^{-1} B\right)<\infty
$$

where we applied Siegel's original integration formula [28].
Next, one verifies that for any Borel set $B \subset \mathbb{R}^{d}$ and any $T \in \mathrm{SL}_{d}(\mathbb{R})$ and $g \in G_{j}$, we have $\widehat{\left(\chi_{B T}\right)_{\psi}}\left(\Gamma_{j} g\right)=\widehat{\left(\chi_{B}\right)_{\psi}}\left(\Gamma_{j} g T^{-1}\right)$. Using this identity, and the fact that $\bar{\omega}$ is $\mathrm{SL}_{d}(\mathbb{R})$ invariant (cf. Lemma 4.5), it follows that $\mu(B T)=\mu(B)$. Hence $\mu$ is $\mathrm{SL}_{d}(\mathbb{R})$ invariant. By the well-known characterization of $\mathrm{SL}_{d}(\mathbb{R})$ invariant Radon measures on $\mathbb{R}^{d}$, it follows that

$$
\begin{equation*}
\mu=c_{1} \delta_{\mathbf{0}}+c_{2} \mathrm{vol} \tag{4.17}
\end{equation*}
$$

for some constants $c_{1}, c_{2} \in \mathbb{R}_{\geq 0}$, where vol is the Lebesgue measure on $\mathbb{R}^{d}$. Here

$$
c_{1}=\mu(\{\mathbf{0}\})=\int_{F_{d}} \int_{\mathbb{T}_{j}^{d}}\left(\widehat{\left.\chi_{\{\mathbf{0}\}}\right)^{\prime}} \psi(x(U) A) d \omega(U) d \nu(A)=\omega\left(\widetilde{\mathfrak{r}}_{i}^{-1}(\{\mathbf{0}\})\right),\right.
$$

where the second equality holds by (4.10), and the last equality holds since $\widehat{\left(\chi_{\{0\}}\right)_{\psi}}(x(U) A)=$ $I\left(\widetilde{\mathrm{r}}_{i}(U)=\mathbf{0}\right)$ for all $U \in \mathbb{T}_{j}^{d}$ and all $A \in \mathrm{SL}_{d}(\mathbb{R})$. Furthermore, it follows from (4.17) that

$$
c_{2}=\lim _{R \rightarrow \infty} \frac{\mu\left(\mathcal{B}_{R}^{d}\right)}{\operatorname{vol}\left(\mathcal{B}_{R}^{d}\right)}=\lim _{R \rightarrow \infty} \int_{\mathbb{X}_{j}} \frac{\widehat{\left(\chi_{\left.\mathcal{B}_{R}^{d}\right)}\right.} \boldsymbol{\psi}}{\operatorname{vol}\left(\mathcal{B}_{R}^{d}\right)} d \bar{\omega}
$$

and here for each fixed point in $\Gamma_{j} g \in \mathbb{X}_{j}$, the value of the integrand tends to $c_{\psi}^{-d}=\bar{n}_{\psi}$ as $R \rightarrow \infty$, by (4.14) and since $c_{\psi} \cdot\left(\mathbb{Z}^{d} \mathrm{a}_{i}(g)\right)$ is an affine lattice of covolume $c_{\psi}^{d}$ in $\mathbb{R}^{d}$. Hence by Lebesgue's dominated convergence theorem, the application of which is justified by the bound

$$
\frac{\widehat{\left(\mathcal{B}_{R}^{d}\right)_{\psi}}\left(\Gamma_{j} g\right)}{\operatorname{vol}\left(\mathcal{B}_{R}^{d}\right)} \ll \frac{\#\left(\mathcal{B}_{R}^{d} \cap c_{\psi} \mathbb{Z}^{d} \iota(g)\right)}{\operatorname{vol}\left(\mathcal{B}_{R}^{d}\right)} \ll \frac{\#\left(\mathcal{B}_{1}^{d} \cap c_{\psi} \mathbb{Z}^{d} \iota(g)\right)}{\operatorname{vol}\left(\mathcal{B}_{1}^{d}\right)}
$$

(where we first used (4.16) and then [10, Proposition 4]) and Siegel's integration formula, [28], we conclude that $c_{2}=\bar{n}_{\psi}$.

We have thus proved that (4.15) holds whenever $f$ is the characteristic function of a Borel set $B \subset \mathbb{R}^{d}$. By taking finite linear combinations, it follows that (4.15) holds whenever $f$
is a simple function, and finally by a standard approximation argument using the monotone convergence theorem one proves (4.15) in the general case.

## 5. Limit behaviour of $\omega_{j}^{(q)}$

As we have mentioned, it will be important for us to understand how $\omega_{j}^{(\boldsymbol{q})}$ varies as $\boldsymbol{q}$ varies in $\mathcal{P}$; in particular we are interested in the behaviour of $\omega_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ tends to infinity within $\mathcal{P}$. The main result of the present section states that, under a certain admissibility assumption on the presentation of the given point set $\mathcal{P}$ as a union of grids, it holds for each $\psi \in \Psi$ that the measure $\omega_{j}^{(\boldsymbol{q})}$ tends to a unique limit measure $\omega_{j}^{\psi} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ as $\boldsymbol{q}$ tends to infinity within a full density subset of $\mathcal{L}_{\psi}$.
5.1. The limit spaces $L_{j}^{\psi}$ and $L_{j}$. As a first step, we will study the behaviour of the linear spaces $L_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ varies through a fixed grid $\mathcal{L}_{\psi}$. For each $\psi$ and $j$, we will define a certain rational space $L_{j}^{\psi} \subset \mathbb{R}^{r_{j}}$, and will prove that this space is, in a certain sense, the limit of the spaces $L_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ tends to infinity within a full density subset of $\mathcal{L}_{\psi}$.

Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$. We start by introducing some auxiliary notation for keeping track of the matrix $U_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ varies within $\mathbb{R}^{d}$ and $\mathcal{L}_{\psi}$, respectively. Recall from (3.4) that we have introduced $M_{j} \in \mathrm{SL}_{d}(\mathbb{R}), c_{\psi} \in \mathbb{R}_{>0}$ and $\boldsymbol{w}_{\psi} \in \mathbb{R}^{d}$ so that $\mathcal{L}_{\psi}=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$. Now set, for any $j \in\{1, \ldots, N\}$ and $\psi \in \Psi$,

$$
\widetilde{\boldsymbol{c}}_{j}=\left(\begin{array}{c}
c_{j, 1}^{-1}  \tag{5.1}\\
\vdots \\
c_{j, r_{j}}^{-1}
\end{array}\right) \in \mathbb{R}^{r_{j}} \quad \text { and } \quad \boldsymbol{c}_{j}^{\psi}:=c_{\psi} \widetilde{\boldsymbol{c}}_{j} \in \mathbb{R}^{r_{j}}
$$

and also

$$
W_{j}:=\left(\begin{array}{c}
\boldsymbol{w}_{j, 1}  \tag{5.2}\\
\vdots \\
\boldsymbol{w}_{j, r_{j}}
\end{array}\right) \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}) \quad \text { and } \quad W_{j}^{\psi}:=W_{j}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{w}_{\psi} T_{j}^{\psi} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}),
$$

where

$$
T_{j}^{\psi}:=M_{j_{\psi}} M_{j}^{-1} \in \mathrm{SL}_{d}(\mathbb{R}) .
$$

The point of this notation is that we now have (cf. (3.15))

$$
\begin{equation*}
U_{j}^{(\boldsymbol{q})}=W_{j}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j}^{-1}, \quad \forall \boldsymbol{q} \in \mathbb{R}^{d} \tag{5.3}
\end{equation*}
$$

It follows that for any point $\boldsymbol{q}$ in $\mathcal{L}_{\psi}=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$, we have:

$$
\begin{equation*}
U_{j}^{(\boldsymbol{q})}=W_{j}^{\psi}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j}^{\psi} \quad \text { when } \boldsymbol{q}=c_{\psi}\left(\boldsymbol{m}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}\left(\forall \boldsymbol{m} \in \mathbb{Z}^{d}\right) \tag{5.4}
\end{equation*}
$$

Now let $L_{j}^{\psi}$ be the rational subspace of $\mathbb{R}^{r_{j}}$ given by

$$
L_{j}^{\psi}:= \begin{cases}\mathfrak{L}\left(\left\{\boldsymbol{c}_{j}^{\psi}\right\} \cup\left\{W_{j, \ell}^{\psi}: \ell \in\{1, \ldots, d\}\right\}\right) & \text { if } j=j_{\psi}  \tag{5.5}\\ \mathfrak{L}\left(\mathbb{R} c_{j}^{\psi} \cup\left\{W_{j, \ell}^{\psi}: \ell \in\{1, \ldots, d\}\right\}\right) & \text { if } j \neq j_{\psi},\end{cases}
$$

where $W_{j, 1}^{\psi}, \ldots, W_{j, d}^{\psi} \in \mathbb{R}^{r_{j}}$ are the column vectors of $W_{j}^{\psi}$.
Lemma 5.1. Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$. Let $\mathcal{L}^{\prime}$ be any subgrid of $\mathcal{L}_{\psi}$. Then

$$
L_{j}^{\psi}=\mathfrak{L}\left(\left\{U_{j, \ell}^{(\boldsymbol{q})}: \boldsymbol{q} \in \mathcal{L}^{\prime}, \ell \in\{1, \ldots, d\}\right\}\right),
$$

where $U_{j, 1}^{(\boldsymbol{q})}, \ldots, U_{j, d}^{(\boldsymbol{q})} \in \mathbb{R}^{r_{j}}$ are the column vectors of $U_{j}^{(\boldsymbol{q})}$.

Proof. Set $S_{0}=\left\{W_{j, \ell}^{\psi}: \ell \in\{1, \ldots, d\}\right\}$ and $S=\left\{\boldsymbol{c}_{j}^{\psi}\right\} \cup S_{0}$ if $j=j_{\psi}$; otherwise $S=\mathbb{R} \boldsymbol{c}_{j}^{\psi} \cup S_{0}$. Also set $S^{\prime}:=\left\{U_{j, \ell}^{(\boldsymbol{q})}: \boldsymbol{q} \in \mathcal{L}^{\prime}, \ell \in\{1, \ldots, d\}\right\}$. Then the task is to prove that $\mathfrak{L}(S)=\mathfrak{L}\left(S^{\prime}\right)$. Let $T_{j, 1}^{\psi}, \ldots, T_{j, d}^{\psi} \in \mathbb{R}^{r_{j}}$ be the column vectors of $T_{j}^{\psi}$. Then by (5.4) we have

$$
\begin{equation*}
U_{j, \ell}^{(\boldsymbol{q})}=W_{j, \ell}^{\psi}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j, \ell}^{\psi} \quad \text { when } \boldsymbol{q}=c_{\psi}\left(\boldsymbol{m}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}\left(\forall \boldsymbol{m} \in \mathbb{Z}^{d}, \ell \in\{1, \ldots, d\}\right) \tag{5.6}
\end{equation*}
$$

Using (5.6) and the fact that $T_{j}^{\psi}=I$ if $j=j_{\psi}$, it follows that $S^{\prime} \subset\langle S\rangle$, and so $\mathfrak{L}\left(S^{\prime}\right) \subset \mathfrak{L}(S)$. It remains to prove that $\mathfrak{L}(S) \subset \mathfrak{L}\left(S^{\prime}\right)$.

Let $\mathcal{L}^{\prime \prime}$ be the inverse image of $\mathcal{L}^{\prime}$ under the bijection $\boldsymbol{m} \mapsto c_{\psi}\left(\boldsymbol{m}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$ from $\mathbb{Z}^{d}$ onto $\mathcal{L}_{\psi}$. Then $\mathcal{L}^{\prime \prime}$ is a coset of a full rank sublattice of $\mathbb{Z}^{d}$; hence there is some $n \in \mathbb{Z}^{+}$and some $\boldsymbol{m}_{0} \in \mathbb{Z}^{d}$ such that $\mathcal{L}^{\prime \prime}$ contains $\boldsymbol{m}_{0}+n \mathbb{Z}^{d}$, and so, by (5.6):

$$
\begin{equation*}
W_{j, \ell}^{\psi}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j, \ell}^{\psi} \in S^{\prime}, \quad \forall \boldsymbol{m} \in \boldsymbol{m}_{0}+n \mathbb{Z}^{d}, \ell \in\{1, \ldots, d\} . \tag{5.7}
\end{equation*}
$$

By Lemma 4.3 there is some $m \in \mathbb{Z}^{+}$such that $S^{\prime} \subset m^{-1} \mathbb{Z}^{r_{j}}+\mathfrak{L}\left(S^{\prime}\right)$. Hence, since every point in $\mathbb{Z}^{d}$ can be written as an affine linear combination of points in $\boldsymbol{m}_{0}+n \mathbb{Z}^{d}$, with all weights in $n^{-1} \mathbb{Z}$, it follows that

$$
\begin{equation*}
W_{j, \ell}^{\psi}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j, \ell}^{\psi} \in \mathcal{A}, \quad \forall \boldsymbol{m} \in \mathbb{Z}^{d}, \ell \in\{1, \ldots, d\} \tag{5.8}
\end{equation*}
$$

where $\mathcal{A}:=(n m)^{-1} \mathbb{Z}^{r_{j}}+\mathfrak{L}\left(S^{\prime}\right)$ (this is a closed subgroup of $\mathbb{R}^{r_{j}}$ ). We will prove that $S \subset \mathcal{A}$; by Lemma 4.3 this implies $\mathfrak{L}(S) \subset \mathfrak{L}\left(S^{\prime}\right)$, and so the proof of Lemma 5.1 will be complete. By taking $\boldsymbol{m}=\mathbf{0}$ in (5.8) it follows that $S_{0} \subset \mathcal{A}$; hence it suffices to prove that $\boldsymbol{c}_{j}^{\psi} \in \mathcal{A}$ if $j=j_{\psi}$, and $\mathbb{R} \boldsymbol{c}_{j}^{\psi} \subset \mathcal{A}$ if $j \neq j_{\psi}$.

We denote by $\boldsymbol{e}_{k}$ the $k$ th standard unit vector in $\mathbb{R}^{d}$. If $j=j_{\psi}$ then $T_{j}^{\psi}$ is the identity matrix, and applying (5.8) with $\ell=1$ and $\boldsymbol{m} \in\left\{\boldsymbol{e}_{1}, \mathbf{0}\right\}$ it follows that $\mathcal{A}$ contains both $W_{j, \ell}^{\psi}-\boldsymbol{c}_{j}^{\psi}$ and $W_{j, \ell}^{\psi}$; hence also $\boldsymbol{c}_{j}^{\psi} \in \mathcal{A}$.

Next assume $j \neq j_{\psi}$. Then $T_{j}^{\psi} \notin \mathcal{S}($ cf. (3.7) $)$, and so there exist two (non-zero) entries of the matrix $T_{j}^{\psi}$ which have an irrational ratio; say entries $k, \ell$ and $k^{\prime}, \ell^{\prime}$, respectively. Writing $T_{j}^{\psi}=\left(t_{r, s}\right)$ we thus have $t_{k, \ell} / t_{k^{\prime}, \ell^{\prime}} \notin \mathbb{Q}$, which implies that the set $\left\{a t_{k, \ell}+b t_{k^{\prime}, \ell^{\prime}}: a, b \in \mathbb{Z}\right\}$ is dense in $\mathbb{R}$. By considering the difference of two arbitrary vectors as in (5.8) we have $\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j, \ell}^{\psi} \in \mathcal{A}$ for all $\boldsymbol{m} \in \mathbb{Z}^{d}$. Taking here $\boldsymbol{m}=\boldsymbol{e}_{k}$ gives $t_{k, \ell} \boldsymbol{c}_{j}^{\psi} \in \mathcal{A}$. Similarly we also have $t_{k^{\prime}, \ell^{\prime}} \boldsymbol{c}_{j}^{\psi} \in \mathcal{A}$, and hence $\left(a t_{k, \ell}+b t_{k^{\prime}, \ell^{\prime}}\right) \boldsymbol{c}_{j}^{\psi} \in \mathcal{A}$ for all $a, b \in \mathbb{Z}$. Hence since $\mathcal{A}$ is closed, we have $\mathbb{R} \boldsymbol{c}_{j}^{\psi} \subset \mathcal{A}$.
Lemma 5.2. Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$. For every $\boldsymbol{q} \in \mathcal{L}_{\psi}$ we have $L_{j}^{(\boldsymbol{q})} \subset L_{j}^{\psi}$.
Proof. By definition, $L_{j}^{(\boldsymbol{q})}=\mathfrak{L}\left(\left\{U_{j, \ell}^{(\boldsymbol{q})}: \ell \in\{1, \ldots, d\}\right\}\right)$; hence the statement follows from Lemma 5.1, and the fact that $\mathfrak{L}(S)$ is increasing in $S$.

Our goal is now to prove Lemma 5.4 below, which says that, in an appropriate sense, the space $L_{j}^{(\boldsymbol{q})}$ approaches $L_{j}^{\psi}$ as $\boldsymbol{q}$ tends to infinity within a full density subset of $\mathcal{L}_{\psi}$. (Cf. also Remark 5.3 below.) We will need the following auxiliary lemma. We denote by ". "the standard scalar product in the space $\mathbb{R}^{r_{j}}$.

Lemma 5.3. Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$. For any non-empty subset $S \subset \mathbb{R}^{r_{j}}$ and any vector $\boldsymbol{a} \in \mathbb{Q}^{r_{j}}$, we have $\boldsymbol{a} \perp \mathfrak{L}(S)$ if and only if $\{\boldsymbol{a} \cdot \boldsymbol{v}: \boldsymbol{v} \in S\} \subset n^{-1} \mathbb{Z}$ for some $n \in \mathbb{Z}^{+}$.

Proof. The case $\boldsymbol{a}=\mathbf{0}$ is trivial; hence from now on we assume $\boldsymbol{a} \neq \mathbf{0}$. The statement of the lemma is invariant under rescaling of $\boldsymbol{a}$ by a non-zero rational number; hence we may assume that $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$ and $\operatorname{gcd}\left(a_{1}, \ldots, a_{r_{j}}\right)=1$. It follows that $\left\{\boldsymbol{a} \cdot \boldsymbol{m}: \boldsymbol{m} \in \mathbb{Z}^{r_{j}}\right\}=\mathbb{Z}$, and hence $\left\{\boldsymbol{a} \cdot \boldsymbol{m}: \boldsymbol{m} \in n^{-1} \mathbb{Z}^{r_{j}}\right\}=n^{-1} \mathbb{Z}$ for any $n \in \mathbb{Z}^{+}$. Therefore, if $\{\boldsymbol{a} \cdot \boldsymbol{v}: \boldsymbol{v} \in S\} \subset n^{-1} \mathbb{Z}$ then $S \subset n^{-1} \mathbb{Z}^{r_{j}}+\boldsymbol{a}^{\perp}$, and so by Lemma 4.3, $\mathfrak{L}(S) \subset \boldsymbol{a}^{\perp}$, i.e. $\boldsymbol{a} \perp \mathfrak{L}(S)$. Conversely, assume
$\boldsymbol{a} \perp \mathfrak{L}(S)$. By Lemma 4.3 there is some $n \in \mathbb{Z}^{+}$such that $S \subset n^{-1} \mathbb{Z}^{r_{j}}+\mathfrak{L}(S) \subset n^{-1} \mathbb{Z}^{r_{j}}+\boldsymbol{a}^{\perp}$, and this implies $\{\boldsymbol{a} \cdot \boldsymbol{v}: \boldsymbol{v} \in S\} \subset n^{-1} \mathbb{Z}$.

Remark 5.1. If $S$ is a finite subset of $\mathbb{R}^{r_{j}}$, then Lemma 5.3 implies that a vector $\boldsymbol{a} \in \mathbb{Q}^{r}$ is orthogonal to $\mathfrak{L}(S)$ if and only if $\{\boldsymbol{a} \cdot \boldsymbol{v}: \boldsymbol{v} \in S\} \subset \mathbb{Q}$. Furthermore, if $S=\mathbb{R} \boldsymbol{c} \cup S^{\prime}$ for some vector $\boldsymbol{c} \in \mathbb{R}^{r}$ and a finite subset $S^{\prime} \subset \mathbb{R}^{r}$, then Lemma 5.3 implies that a vector $\boldsymbol{a} \in \mathbb{Q}^{r}$ is orthogonal to $\mathfrak{L}(S)$ if and only if $\boldsymbol{a} \perp \boldsymbol{c}$ and $\left\{\boldsymbol{a} \cdot \boldsymbol{v}: \boldsymbol{v} \in S^{\prime}\right\} \subset \mathbb{Q}$.
Lemma 5.4. Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$. For every $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$ with $\boldsymbol{a} \not \perp L_{j}^{\psi}$, we have

$$
\begin{equation*}
\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{R}^{d}: \boldsymbol{a} \perp L_{j}^{(\boldsymbol{q})}\right\} \ll R^{d-1} \quad \text { as } R \rightarrow \infty \tag{5.9}
\end{equation*}
$$

Proof. As before, let us parametrize the points in $\mathcal{L}_{\psi}=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$ as $\boldsymbol{q}=\boldsymbol{q}(\boldsymbol{m})=$ $c_{\psi}\left(\boldsymbol{m}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$, with $\boldsymbol{m}$ running through $\mathbb{Z}^{d}$. We will prove the following bound, which clearly implies (5.9):

$$
\begin{equation*}
\#\left\{\boldsymbol{m} \in \mathbb{Z}^{d} \cap[-R, R]^{d}: \boldsymbol{a} \perp L_{j}^{(\boldsymbol{q}(\boldsymbol{m}))}\right\} \leq(2 R+1)^{d-1} \quad(\forall R>0) \tag{5.10}
\end{equation*}
$$

$\operatorname{Using} L_{j}^{(\boldsymbol{q})}=\mathfrak{L}\left(\left\{U_{j, \ell}^{(\boldsymbol{q})}: \ell \in\{1, \ldots, d\}\right\}\right)$ together with (5.6) and Lemma 5.3, the bound (5.10) can equivalently be stated as:

$$
\begin{equation*}
\#\left\{\boldsymbol{m} \in \mathbb{Z}^{d} \cap[-R, R]^{d}: \boldsymbol{a} \cdot\left(W_{j, \ell}^{\psi}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{m} T_{j, \ell}^{\psi}\right) \in \mathbb{Q} \forall \ell \in\{1, \ldots, d\}\right\} \leq(2 R+1)^{d-1} \tag{5.11}
\end{equation*}
$$

We are assuming that $\boldsymbol{a}$ is not orthogonal to $L_{j}^{\psi}$. By (5.5) and Lemma5.3 (cf. also Remark 5.1), this implies that

$$
\begin{cases}\text { If } j=j_{\psi}: & \boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi} \notin \mathbb{Q} \text { or } \boldsymbol{a} \cdot W_{j, \ell}^{\psi} \notin \mathbb{Q} \text { for some } \ell \in\{1, \ldots, d\} \\ \text { If } j \neq j_{\psi}: & \boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi} \neq 0 \text { or } \boldsymbol{a} \cdot W_{j, \ell}^{\psi} \notin \mathbb{Q} \text { for some } \ell \in\{1, \ldots, d\}\end{cases}
$$

Let us first assume $j \neq j_{\psi}$. If $\boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi}=0$ then $\boldsymbol{a} \cdot W_{j, \ell}^{\psi} \notin \mathbb{Q}$ for some $\ell$, and it follows that the set in (5.11) is empty, for all $R$. Hence we may assume that $\boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi} \neq 0$. Now $j \neq j_{\psi}$ implies that $T_{j}^{\psi} \notin \mathcal{S}$, which means that there exist two (non-zero) entries of the matrix $T_{j}^{\psi}$ which have an irrational ratio. Hence there exist $k, \ell \in\{1, \ldots, d\}$ such that $\left(\boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi}\right)\left(\boldsymbol{e}_{k} T_{j, \ell}^{\psi}\right) \notin \mathbb{Q}$. This implies that, for any $R$, the set in (5.11) contains at most one point $\boldsymbol{m}$ along any line parallel with $\boldsymbol{e}_{k}$. Hence the bound in (5.11) holds.

Next assume $j=j_{\psi}$. Then $T_{j}^{\psi}=$ I. If $\boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi} \in \mathbb{Q}$ then $\boldsymbol{a} \cdot W_{j, \ell}^{\psi} \notin \mathbb{Q}$ for some $\ell \in\{1, \ldots, d\}$, and it follows that the set in (5.11) is empty, for all $R$. Hence we may assume that $\boldsymbol{a} \cdot \boldsymbol{c}_{j}^{\psi} \notin \mathbb{Q}$. Then for each $\ell \in\{1, \ldots, d\}$ there is at most one integer $m$ such that $\boldsymbol{a} \cdot\left(W_{j, \ell}^{\psi}-m \boldsymbol{c}_{j}^{\psi}\right) \in \mathbb{Q}$, and so the set in (5.11) contains at most one point, and the bound in (5.11) holds.

We end this section by introducing a certain rational space $L_{j} \subset \mathbb{R}^{r_{j}}$ which equals $L_{j}^{(\boldsymbol{q})}$ for $\boldsymbol{q}$ generic within $\mathbb{R}^{d}$ (cf. Lemma 5.5 below). These spaces are closely related to the spaces $L_{j}^{\psi}$ introduced above, and they appear in the explicit description of the macroscopic limit measure $\mu^{\mathrm{g}}$ which we discuss in Section 7.4 below. For each $j \in\{1, \ldots, N\}$, we set

$$
\begin{equation*}
L_{j}:=\mathfrak{L}\left(\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \cup\left\{W_{j, \ell}: \ell \in\{1, \ldots, d\}\right\}\right) \tag{5.12}
\end{equation*}
$$

where $W_{j, 1}, \ldots, W_{j, d} \in \mathbb{R}^{r_{j}}$ are the column vectors of $W_{j}$ (see (5.2)).
Lemma 5.5. Let $j \in\{1, \ldots, N\}$. For every $\boldsymbol{q} \in \mathbb{R}^{d}$ we have $L_{j}^{(\boldsymbol{q})} \subset L_{j}$, and for all except countably many $\boldsymbol{q} \in \mathbb{R}^{d}$ we even have $L_{j}^{(\boldsymbol{q})}=L_{j}$.

Proof. Let $S_{0}=\left\{W_{j, \ell}: \ell \in\{1, \ldots, d\}\right.$ and $S=\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \cup S_{0}$, so that $L_{j}=\mathfrak{L}(S)$. Recall that for every $\boldsymbol{q} \in \mathbb{R}^{d}$ we have $L_{j}^{(\boldsymbol{q})}=\mathfrak{L}\left(\left\{U_{j, \ell}^{(\boldsymbol{q})}: \ell \in\{1, \ldots, d\}\right\}\right)$; and by (5.3),$U_{j, \ell}^{(\boldsymbol{q})}=W_{j, \ell}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j, \ell}^{\prime}$,
where $M_{j, \ell}^{\prime} \in \mathbb{R}^{r_{j}}$ is the $\ell$ th column vector of $M_{j}^{-1}$. Here $\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j, \ell}^{\prime} \in \mathbb{R} \widetilde{\boldsymbol{c}}_{j}$, and thus $U_{j, \ell}^{(\boldsymbol{q})} \in\langle S\rangle$ for every $\ell$. Therefore $L_{j}^{(\boldsymbol{q})} \subset L_{j}$.

Next, since both $L_{j}^{(\boldsymbol{q})}$ and $L_{j}$ are rational subspaces of $\mathbb{R}^{r_{j}}$, if $L_{j}^{(\boldsymbol{q})} \subsetneq L_{j}$ then there exists some $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$ satisfying $\boldsymbol{a} \perp L_{j}^{(\boldsymbol{q})}$ but $\boldsymbol{a} \not \perp L_{j}$. By Lemma 5.3 (cf. also Remark 5.1) this means that

$$
\begin{equation*}
\boldsymbol{a} \cdot\left(W_{j, \ell}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j, \ell}^{\prime}\right) \in \mathbb{Q}, \quad \forall \ell \in\{1, \ldots, d\} \tag{5.13}
\end{equation*}
$$

but either $\boldsymbol{a} \not \perp \widetilde{\boldsymbol{c}}_{j}$ or $\boldsymbol{a} \cdot W_{j, \ell} \notin \mathbb{Q}$ for some $\ell$. Clearly this is not possible if $\boldsymbol{a} \perp \widetilde{\boldsymbol{c}}_{j}$; hence we must have $\boldsymbol{a} \not \perp \widetilde{\boldsymbol{c}}_{j}$. Noticing also that the condition (5.13) is equivalent $\operatorname{td}{ }^{7} \boldsymbol{a}^{\top}\left(W_{j}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j}^{-1}\right) \in \mathbb{Q}^{d}$, we conclude that:

$$
\begin{equation*}
\left\{\boldsymbol{q} \in \mathbb{R}^{d}: L_{j}^{(\boldsymbol{q})} \subsetneq L_{j}\right\}=\bigcup_{\substack{\boldsymbol{a} \in \mathbb{Z}_{j}^{r} \\\left(\boldsymbol{a} \not \subset \widetilde{\boldsymbol{c}}_{j}\right)}} \bigcup_{\boldsymbol{b} \in \mathbb{Q}^{d}}\left\{\boldsymbol{q} \in \mathbb{R}^{d}: \boldsymbol{a}^{\top}\left(W_{j}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j}^{-1}\right)=\boldsymbol{b}\right\} \tag{5.14}
\end{equation*}
$$

But for every $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$ with $\boldsymbol{a} \not \perp \widetilde{\boldsymbol{c}}_{j}$ and every $\boldsymbol{b} \in \mathbb{Q}^{d}$, the set $\left\{\boldsymbol{q} \in \mathbb{R}^{d}: \boldsymbol{a}^{\top}\left(W_{j}-\widetilde{\boldsymbol{c}}_{j} \boldsymbol{q} M_{j}^{-1}\right)=\boldsymbol{b}\right\}$ consists of exactly one point, namely $\boldsymbol{q}=\left(\boldsymbol{a}^{\boldsymbol{\top}} \widetilde{\boldsymbol{c}}_{j}\right)^{-1}\left(\boldsymbol{a}^{\top} W_{j}-\boldsymbol{b}\right) M_{j}$. Hence the set in (5.14) is countable, and the lemma is proved.
Lemma 5.6. For any $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$ we have $L_{j}^{\psi} \subset L_{j}$, and if $j \neq j_{\psi}$ then even $L_{j}^{\psi}=L_{j}$.
Proof. Recall that $\boldsymbol{c}_{j}^{\psi}=c_{\psi} \widetilde{\boldsymbol{c}}_{j}$; hence $\mathbb{R} \boldsymbol{c}_{j}^{\psi}=\mathbb{R} \widetilde{\boldsymbol{c}}_{j}$. Furthermore, by comparing the definitions (5.2) and (5.2) we note that $W_{j, \ell}^{\psi}=W_{j, \ell}-\boldsymbol{c}_{j}^{\psi} \boldsymbol{w}_{\psi} T_{j, \ell}^{\psi} \in W_{j, \ell}+\mathbb{R} \boldsymbol{c}_{j}^{\psi}=W_{j, \ell}+\mathbb{R} \tilde{\boldsymbol{c}}_{j}$. Hence the two sets $\mathbb{R} \boldsymbol{c}_{j}^{\psi} \cup\left\{W_{j, \ell}^{\psi}: \ell \in\{1, \ldots, d\}\right\}$ and $\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \cup\left\{W_{j, \ell}: \ell \in\{1, \ldots, d\}\right\}$ generate the same subgroups of $\mathbb{R}^{r_{j}}$. In view of this fact, the lemma now follows by inspecting the definitions of $L_{j}^{\psi}$ and $L_{j}$, (5.5) and (5.12).
Lemma 5.7. For every $\psi=(j, i) \in \Psi, L_{j}^{\psi}=L_{j} \cap \boldsymbol{e}_{i}^{\perp}$ and $L_{j}=L_{j}^{\psi}+\mathbb{R} \widetilde{\boldsymbol{c}}_{j}$.
Proof. Since $\psi=(j, i)$, the $i$ th coordinate of $\boldsymbol{c}_{j}^{\psi}$ is 1 by (5.1) and $r_{i}\left(W_{j}^{\psi}\right)=\mathbf{0}$ by (5.2); hence it follows from (5.5) that $L_{j}^{\psi} \perp \boldsymbol{e}_{i}$, and so by Lemma 5.6 we have $L_{j}^{\psi} \subset L_{j} \cap \boldsymbol{e}_{i}^{\perp}$. On the other hand, $\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \subset L_{j}$ by (5.12) and $\widetilde{\boldsymbol{c}}_{j} \cdot \boldsymbol{e}_{i}=c_{\psi}^{-1}>0$; in particular $L_{j} \not \perp \boldsymbol{e}_{i}$. Now to complete the proof of the lemma it suffices to prove that

$$
\begin{equation*}
L_{j} \subset \widetilde{L}, \quad \text { with } \widetilde{L}:=L_{j}^{\psi}+\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \tag{5.15}
\end{equation*}
$$

But it follows from (5.5) and Lemma 4.3 that $\boldsymbol{c}_{j}^{\psi} \in \mathbb{Q}^{r_{j}}+L_{j}^{\psi}$; hence since $\widetilde{L}$ equals the linear span of $L_{j}^{\psi}$ and the vector $\boldsymbol{c}_{j}^{\psi}, \widetilde{L}$ is a rational subspace of $\mathbb{R}^{r_{j}}$. It also follows from (5.5) and Lemma 4.3 that each column vector of $W_{j}^{\psi}$ lies in $\mathbb{Q}^{r_{j}}+L_{j}^{\psi}$; thus $W_{j}^{\psi} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+\left(L_{j}^{\psi}\right)^{d}$, and hence we have

$$
W_{j}=W_{j}^{\psi}+\boldsymbol{c}_{j}^{\psi} \boldsymbol{w}_{\psi} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+\left(L_{j}^{\psi}\right)^{d}+\boldsymbol{c}_{j}^{\psi} \boldsymbol{w}_{\psi} \subset \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+\widetilde{L}^{d}
$$

This implies that $W_{\ell}^{\psi} \in \mathbb{Q}^{r_{j}}+\widetilde{L}$ for each $\ell \in\{1, \ldots, d\}$, and since also $\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \subset \widetilde{L}$, the inclusion in (5.15) now follows from (5.12) and Lemma 4.3.
5.2. The limit measures $\omega_{j}^{\psi}$ and $\omega_{j}^{\mathrm{g}}$. The following condition will be of crucial importance for us. It involves the vectors $\boldsymbol{c}_{j}^{\psi} \in \mathbb{R}^{r_{j}}$ and the subspaces $L_{j}^{\psi}$ which were defined in (5.1) and (5.5) in the previous section.

Definition 5.1. We say that a presentation of $\mathcal{P}$ as in (3.5), and satisfying (3.6) and (3.7), is admissible if $\boldsymbol{c}_{j}^{\psi} \in L_{j}^{\psi}+\mathbb{Z}^{r_{j}}$ for all $\psi=(j, i) \in \Psi$.

[^5]It should be carefully noted that the condition in Definition 5.1 only involves vectors $\boldsymbol{c}_{j}^{\psi}$ and spaces $L_{j}^{\psi}$ for pairs of $\psi$ and $j$ with $j=j_{\psi}$.

As we will see below, the admissibility of the presentation of $\mathcal{P}$ is a necessary (as well as sufficient) condition to ensure that all the measures $\omega_{j}^{(\boldsymbol{q})}$ have a unique generic limit as $\boldsymbol{q}$ tends to infinity within any of the grids $\mathcal{L}_{\psi}$. Luckily, it turns out that any point set $\mathcal{P}$ which is a finite union of grids possesses an admissible presentation; we will prove this in Section 5.3 below.

From now on we assume that the given presentation (3.5) of $\mathcal{P}$ is admissible.
We will start by defining, for any $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$, a measure $\omega_{j}^{\psi} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$. We will then prove that $\omega_{j}^{\psi}$ is in fact the generic limit of the measures $\omega_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ tends to infinity within the grid $\mathcal{L}_{\psi}$. As we will see, the assumption about admissibility is needed already for the definition of $\omega_{j}^{\psi}$ to make sense. (Cf. the proof of Lemma 5.8 below.)

Given $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$ we set $\mathbb{S}_{j}^{\psi}:=\pi\left(L_{j}^{\psi}\right)$. This is a closed subtorus of $\mathbb{T}_{j}$, since $L_{j}^{\psi}$ is a rational subspace of $\mathbb{R}^{r_{j}}$. Furthermore, we pick an arbitrary point $\boldsymbol{q} \in \mathcal{L}_{\psi}$, and define:

$$
\begin{equation*}
\widetilde{\mathbb{S}}_{j}^{\psi}:=\mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}^{\psi} \subset \mathbb{T}_{j} ; \quad \text { and } \quad \mathcal{O}_{j}^{\psi}:=\bigcup_{\gamma \in \mathrm{SL}_{d}(\mathbb{Z})}\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right) \gamma+\left(\mathbb{S}_{j}^{\psi}\right)^{d}\right) \subset \mathbb{T}_{j}^{d} \tag{5.16}
\end{equation*}
$$

Lemma 5.8. Both $\widetilde{\mathbb{S}}_{j}^{\psi}$ and $\mathcal{O}_{j}^{\psi}$ are well-defined, i.e. the expressions in (5.16) are independent of the choice of $\boldsymbol{q}$. We have $\mathcal{O}_{j}^{(\boldsymbol{q})} \subset \mathcal{O}_{j}^{\psi}$ for all $\boldsymbol{q} \in \mathcal{L}_{\psi}$. Furthermore, $\widetilde{\mathbb{S}}_{j}^{\psi}$ is a closed subgroup of $\mathbb{T}_{j}$ whose connected component subgroup equals $\mathbb{S}_{j}^{\psi}$, and $\mathcal{O}_{j}^{\psi}$ is a union of some of the connected components of $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$.
(It should be noted that since $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$ is compact, its total number of connected components is finite.)

In view of Lemma [5.8, we may now define $\omega_{j}^{\psi} \in P\left(\mathbb{T}_{j}^{d}\right)$ to be the restriction to $\mathcal{O}_{j}^{\psi}$ of the Haar measure on $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$, normalized so that $\omega_{j}^{\psi}\left(\mathcal{O}_{j}^{\psi}\right)=1$. Note that $\omega_{j}^{\psi}$ is $\mathrm{SL}_{d}(\mathbb{Z})$-invariant by construction, i.e. we actually have $\omega_{j}^{\psi} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$.
Proof of Lemma 5.8. Let us first prove that

$$
\begin{equation*}
U_{j, \ell}^{(\boldsymbol{q})}-U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)} \in L_{j}^{\psi}+\mathbb{Z}^{r_{j}}, \quad \forall \boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathcal{L}_{\psi}, \forall \ell \in\{1, \ldots, d\} . \tag{5.17}
\end{equation*}
$$

Indeed, pick $\boldsymbol{m}, \boldsymbol{m}^{\prime} \in \mathbb{Z}^{d}$ so that $\boldsymbol{q}=c_{\psi}\left(\boldsymbol{m}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$ and $\boldsymbol{q}^{\prime}=c_{\psi}\left(\boldsymbol{m}^{\prime}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$; then by (5.6) we have $U_{j, \ell}^{(\boldsymbol{q})}-U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)}=\boldsymbol{c}_{j}^{\psi}\left(\boldsymbol{m}^{\prime}-\boldsymbol{m}\right) T_{j, \ell}^{\psi}$. If $j \neq j_{\psi}$ then $\mathbb{R} \boldsymbol{c}_{j}^{\psi} \subset L_{j}^{\psi}$ by (5.5) and hence $U_{j, \ell}^{(\boldsymbol{q})}-U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)} \in L_{j}^{\psi}$ for each $\ell$. On the other hand if $j=j_{\psi}$ then $T_{j}^{\psi}=I$ and so $U_{j, \ell}^{(\boldsymbol{q})}-U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)} \in \mathbb{Z} \boldsymbol{c}_{j}^{\psi}$, and using the assumption that $\mathcal{P}$ is admissible (cf. Def. [5.1), this implies that (5.17) holds.

In order to prove that $\widetilde{\mathbb{S}}_{j}^{\psi}$ is well-defined it suffices to verify that for any two $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathcal{L}_{\psi}$ we have $\mathbb{S}_{j}^{\left(\boldsymbol{q}^{\prime}\right)} \subset \mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}^{\psi}$. However this follows from the definition of $\mathbb{S}_{j}^{\left(\boldsymbol{q}^{\prime}\right)}$ and the fact that, by (5.17), $\pi\left(U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)}\right) \in \pi\left(U_{j, \ell}^{(\boldsymbol{q})}\right)+\mathbb{S}_{j}^{\psi} \subset \mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}^{\psi}$. Note also that (5.17) implies that $\pi\left(U_{j}^{(\boldsymbol{q})}\right)-\pi\left(U_{j}^{\left(\boldsymbol{q}^{\prime}\right)}\right) \in\left(\mathbb{S}_{j}^{\psi}\right)^{d}$; and this implies that $\mathcal{O}_{j}^{\psi}$ is well-defined.

Now consider an arbitrary point $\boldsymbol{q} \in \mathcal{L}_{\psi}$. Recall that $L_{j}^{(\boldsymbol{q})} \subset L_{j}^{\psi}$, by Lemma 5.2, hence $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ} \subset \mathbb{S}_{j}^{\psi}$, and by inspection in (4.7) and (5.16) (and recalling $\mathcal{O}_{j}^{(q)}:=\mathcal{O}_{j}^{\left(\pi\left(U_{j}^{(q)}\right)\right)}$ ), this implies that $\mathcal{O}_{j}^{(\boldsymbol{q})} \subset \mathcal{O}_{j}^{\psi}$. Recall also that $\mathbb{S}_{j}^{(\boldsymbol{q})}$ is a finite union of $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ}$-cosets in $\mathbb{T}_{j}$; therefore $\widetilde{\mathbb{S}}_{j}^{\psi}=\mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}^{\psi}$ is a finite union of $\mathbb{S}_{j}^{\psi}$-cosets. Hence $\widetilde{\mathbb{S}}_{j}^{\psi}$ is indeed a closed subgroup of $\mathbb{T}_{j}$
whose connected component subgroup equals $\mathbb{S}_{j}^{\psi}$. Finally, from the definition of $\mathcal{O}_{j}^{\psi}$ and the fact that $\pi\left(U_{j}^{(\boldsymbol{q})}\right) \in\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{d}$, it follows that $\mathcal{O}_{j}^{\psi} \subset\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$; and $\mathcal{O}_{j}^{\psi}$ is by definition a union of $\left(\mathbb{S}_{j}^{\psi}\right)^{d}$-cosets, i.e. a union of some of the connected components of $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$.

This is a convenient point to interject the definition of a measure $\omega_{j}^{\mathrm{g}} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$, which is closely analogous to $\omega_{j}^{\psi}$, and which we will need later for the explicit description of the macroscopic limit measure $\mu^{\mathrm{g}}$ which we discuss in Section 7.4 below.

Given $j \in\{1, \ldots, N\}$ we set $\mathbb{S}_{j}:=\pi\left(L_{j}\right)$. This is a closed subtorus of $\mathbb{T}_{j}$. Picking an arbitrary $\boldsymbol{q} \in \mathbb{R}^{d}$, we define:

$$
\begin{equation*}
\widetilde{\mathbb{S}}_{j}:=\mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j} \quad \text { and } \quad \mathcal{O}_{j}:=\bigcup_{\gamma \in \mathrm{SL}_{d}(\mathbb{Z})}\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right) \gamma+\mathbb{S}_{j}^{d}\right) \subset \mathbb{T}_{j}^{d} \tag{5.18}
\end{equation*}
$$

Lemma 5.9. Both $\widetilde{\mathbb{S}}_{j}$ and $\mathcal{O}_{j}$ are well-defined, i.e. the expressions in (5.18) are independent of the choice of $\boldsymbol{q}$. We have $\mathcal{O}_{j}^{(\boldsymbol{q})} \subset \mathcal{O}_{j}$ for all $\boldsymbol{q} \in \mathbb{R}^{d}$. Furthermore, $\widetilde{\mathbb{S}}_{j}$ is a closed subgroup of $\mathbb{T}_{j}$ whose connected component subgroup equals $\mathbb{S}_{j}$, and $\mathcal{O}_{j}$ is a union of some of the connected components of $\widetilde{\mathbb{S}}_{j}^{d}$.

In view of Lemma 5.9 we may now define $\omega_{j}^{\mathrm{g}} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ to be the restriction to $\mathcal{O}_{j}$ of Haar measure on $\widetilde{\mathbb{S}}_{j}^{d}$, normalized so that $\omega_{j}^{\mathrm{g}}\left(\mathcal{O}_{j}\right)=1$.

Proof. It follows from (5.3) that for all $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathbb{R}^{d}$ we have $U_{j, \ell}^{(\boldsymbol{q})}-U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)} \in \mathbb{R} \boldsymbol{c}_{j} \subset L_{j}$, and thus $\pi\left(U_{j, \ell}^{(\boldsymbol{q})}\right)-\pi\left(U_{j, \ell}^{\left(\boldsymbol{q}^{\prime}\right)}\right) \in \mathbb{S}_{j}$, for every $\ell \in\{1, \ldots, d\}$. Recall that $\mathbb{S}_{j}^{(\boldsymbol{q})}=\overline{\left\langle\pi\left(U_{j, 1}^{(\boldsymbol{q})}\right), \ldots, \pi\left(U_{j, d}^{(\boldsymbol{q})}\right)\right\rangle}$. It follows that $\mathbb{S}_{j}^{\left(\boldsymbol{q}^{\prime}\right)} \subset \mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}$ for any $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathbb{R}^{d}$. Hence we get that $\widetilde{\mathbb{S}}_{j}$ is well-defined. It also follows that $\pi\left(U_{j}^{(\boldsymbol{q})}\right)-\pi\left(U_{j}^{\left(\boldsymbol{q}^{\prime}\right)}\right) \in \mathbb{S}_{j}^{d}$ for any $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathbb{R}^{d}$; hence $\mathcal{O}_{j}$ is well-defined. The proof of the remaining assertions is essentially the same as in Lemma 5.8.

Lemma 5.10. For all except at most countably many $\boldsymbol{q} \in \mathbb{R}^{d}$ we have $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ}=\mathbb{S}_{j}, \mathbb{S}_{j}^{(\boldsymbol{q})}=\widetilde{\mathbb{S}}_{j}$, $\mathcal{O}_{j}^{(\boldsymbol{q})}=\mathcal{O}_{j}$, and $\omega_{j}^{(\boldsymbol{q})}=\omega_{j}^{\mathbf{g}}$.

Proof. We will prove that the stated equalities hold whenever $L_{j}^{(\boldsymbol{q})}=L_{j}$; by Lemma 5.5 this gives the statement of the present lemma. Thus assume that $L_{j}^{(\boldsymbol{q})}=L_{j}$. Then $\left(\mathbb{S}_{\tilde{\sim}}^{(\boldsymbol{q})}\right)^{\circ}=$ $\pi\left(L_{j}^{(\boldsymbol{q})}\right)=\pi\left(L_{j}\right)=\mathbb{S}_{j} . \quad$ In particular $\mathbb{S}_{j} \subset \mathbb{S}_{j}^{(\boldsymbol{q})}$, and so and $\mathbb{S}_{j}^{(\boldsymbol{q})}=\mathbb{S}_{j}^{(\boldsymbol{q})}+\mathbb{S}_{j}=\widetilde{\mathbb{S}}_{j} . \quad$ By comparing (5.18) with (4.7) (applied with $V=\pi\left(U_{j}^{(\boldsymbol{q})}\right)$ ), we have $\mathcal{O}_{j}^{(\boldsymbol{q})}=\mathcal{O}_{j}$. Finally now also $\omega_{j}^{(\boldsymbol{q})}=\omega_{j}^{\mathrm{g}}$ is immediate from the definitions.
Lemma 5.11. For any $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$, if $j \neq j_{\psi}$ then $\mathbb{S}_{j}^{\psi}=\mathbb{S}_{j}, \widetilde{\mathbb{S}}_{j}^{\psi}=\widetilde{\mathbb{S}}_{j}, \mathcal{O}_{j}^{\psi}=\mathcal{O}_{j}$ and $\omega_{j}^{\psi}=\omega_{j}^{\mathrm{g}}$.
Proof. If $j \neq j_{\psi}$. then $L_{j}^{\psi}=L_{j}$ by Lemma 5.6, and the stated equalities are immediate from this fact, by inspection in the relevant definitions.

We now return to the discussion on the measures $\omega_{j}^{\psi}$. In Proposition 5.13 and Remark 5.2 below, we will prove that $\omega_{j}^{\psi}$ is the generic limit of the measures $\omega_{j}^{(\boldsymbol{q})}$ as $\boldsymbol{q}$ tends to infinity within the grid $\mathcal{L}_{\psi}$. We will make use of the following auxiliary lemma.

Lemma 5.12. Let $\psi \in \Psi, j \in\{1, \ldots, N\}$ and $\boldsymbol{q} \in \mathcal{L}_{\psi}$. Let $m=m(\boldsymbol{q})$ be the number of $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$-cosets contained in $\mathcal{O}_{j}(\boldsymbol{q})$, and let $n$ be the number of $\left(\mathbb{S}_{j}^{\psi}\right)^{d}$-cosets contained in $\mathcal{O}_{j}^{\psi}$. Then $n$ divides $m$, and every $\left(\mathbb{S}_{j}^{\psi}\right)^{d}$-coset contained in $\mathcal{O}_{j}^{\psi}$ contains exactly $m / n$ distinct $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$-cosets.

Proof. Let $U:=\pi\left(U_{j}^{(\boldsymbol{q})}\right) \in\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{d}$, and let $\Lambda$ be the stabilizer of the coset $U+\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$, for the action of $\mathrm{SL}_{d}(\mathbb{Z})$ on the finite group $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{d} /\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$. It then follows from the definition in (4.2) that $\mathcal{O}_{j}^{(\boldsymbol{q})}$ equals the disjoint union of the cosets $U \gamma+\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$, when $\gamma$ runs through any set of representatives for $\Lambda \backslash \mathrm{SL}_{d}(\mathbb{Z})$. In particular we have $m=\#\left(\Lambda \backslash \mathrm{SL}_{d}(\mathbb{Z})\right)$. Similarly let $\Lambda^{\prime}$ be the stabilizer of the coset $U+\left(\mathbb{S}_{j}^{\psi}\right)^{d}$ in $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d} /\left(\mathbb{S}_{j}^{\psi}\right)^{d}$; then by (5.16), $\mathcal{O}_{j}^{\psi}$ equals the disjoint union of the cosets $U \gamma+\left(\mathbb{S}_{j}^{\psi}\right)^{d}$, when $\gamma$ runs through any set of representatives for $\Lambda^{\prime} \backslash \mathrm{SL}_{d}(\mathbb{Z})$; and in particular we have $n=\#\left(\Lambda^{\prime} \backslash \mathrm{SL}_{d}(\mathbb{Z})\right)$. Note that $\Lambda \subset \Lambda^{\prime}$, since $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ} \subset \mathbb{S}_{j}^{\psi}$. Therefore $n$ divides $m$, and the last statement of the lemma follows from the fact that each right $\Lambda^{\prime}$-coset in $\mathrm{SL}_{d}(\mathbb{Z})$ contains exactly $m / n$ right $\Lambda$-cosets.

Proposition 5.13. Let $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$, and let $U$ be any open neighbourhood of $\omega_{j}^{\psi}$ in $P\left(\mathbb{T}_{j}^{d}\right)$. Then the set $\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \omega_{j}^{(\boldsymbol{q})} \notin U\right\}$ has density zero.
Proof. To each matrix $A \in \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ corresponds a character $\chi_{A}$ on $\mathbb{T}_{j}^{d}=\mathrm{M}_{r_{j} \times d}(\mathbb{R}) / \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ given by $\chi_{A}(X)=e^{2 \pi i \operatorname{Tr}\left(A X^{\top}\right)}$, and every character on $\mathbb{T}_{j}^{d}$ can be so expressed. Using the fact that the set of finite linear combinations of characters on $\mathbb{T}_{j}^{d}$ is dense in $C\left(\mathbb{T}_{j}^{d}\right)$, it follows that there exists a finite subset $S \subset \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ and some $\varepsilon>0$ such that $U$ contains the set

$$
\left\{\mu \in P\left(\mathbb{T}_{j}^{d}\right):\left|\mu\left(\chi_{A}\right)-\omega_{j}^{\psi}\left(\chi_{A}\right)\right|<\varepsilon \quad \forall A \in S\right\}
$$

Hence, using also the fact that a finite union of density zero sets again has density zero, it suffices to prove that for any fixed $A \in \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$, the set

$$
\begin{equation*}
\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}:\left|\omega_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)-\omega_{j}^{\psi}\left(\chi_{A}\right)\right| \geq \varepsilon\right\} \tag{5.19}
\end{equation*}
$$

has density zero.
For each $\boldsymbol{q} \in \mathcal{L}_{\psi}$, let $\nu_{j}^{(\boldsymbol{q})} \in P\left(\mathbb{T}_{j}^{d}\right)$ be the normalized Haar measure on $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$, and let $R(\boldsymbol{q})$ be a set of representatives containing one point in each $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$-coset contained in $\mathcal{O}_{j}^{(\boldsymbol{q})}$ (recall that $\mathcal{O}_{j}^{(\boldsymbol{q})}$ is a finite union of $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}$-cosets). Then by the definition of $\omega_{j}^{(\boldsymbol{q})}$ around (4.2), we have

$$
\omega_{j}^{(\boldsymbol{q})}=\frac{1}{\# R(\boldsymbol{q})} \sum_{X \in R(\boldsymbol{q})} \nu_{j}^{(\boldsymbol{q})} \circ \tau_{X}^{-1}
$$

where $\tau_{X}: \mathbb{T}_{j}^{d} \rightarrow \mathbb{T}_{j}^{d}$ denotes translation by $X$. Also let $\nu_{j}^{\psi} \in P\left(\mathbb{T}_{j}^{d}\right)$ be the normalized Haar measure on $\left(\mathbb{S}_{j}^{\psi}\right)^{d}$. Comparing (4.2) with (5.16) and the definition of $\omega_{j}^{\psi}$ on p. 26, and using Lemmas 5.8 and 5.12, we then have:

$$
\begin{equation*}
\omega_{j}^{\psi}=\frac{1}{\# R(\boldsymbol{q})} \sum_{X \in R(\boldsymbol{q})} \nu_{j}^{\psi_{0}} \circ \tau_{X}^{-1} \tag{5.20}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
\left|\omega_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)-\omega_{j}^{\psi}\left(\chi_{A}\right)\right| \leq \frac{1}{\# R(\boldsymbol{q})} \sum_{X \in R(\boldsymbol{q})}\left|\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A} \circ \tau_{X}\right)-\nu_{j}^{\psi}\left(\chi_{A} \circ \tau_{X}\right)\right|=\left|\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)-\nu_{j}^{\psi}\left(\chi_{A}\right)\right| \tag{5.21}
\end{equation*}
$$

where the last equality holds since $\chi_{A} \circ \tau_{X}=\chi_{A}(X) \cdot \chi_{A}$ for any $X \in \mathbb{T}_{j}^{d}$.
Recall that $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ d}=\pi\left(\left(L_{j}^{(\boldsymbol{q})}\right)^{d}\right)$; hence $\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)=1$ if $A$ is orthogonal to $\left(L_{j}^{(\boldsymbol{q})}\right)^{d}$, otherwise $\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)=0$. Hence, letting $\boldsymbol{a}_{1}, \ldots, \boldsymbol{a}_{d} \in \mathbb{Z}^{r_{j}}$ be the column vectors of $A$, we have $\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)=1$ if each $\boldsymbol{a}_{\ell}$ is orthogonal to $L_{j}^{(\boldsymbol{q})}$, otherwise $\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)=0$. Similarly $\nu_{j}^{\psi}\left(\chi_{A}\right)=1$ if each $\boldsymbol{a}_{\ell}$ is orthogonal to $L_{j}^{\psi}$, otherwise $\nu_{j}^{(\boldsymbol{q})}\left(\chi_{A}\right)=0$. Recall also that $L_{j}^{(\boldsymbol{q})} \subset L_{j}^{\psi}$ for all $\boldsymbol{q} \in \mathcal{L}_{\psi}$
(cf. Lemma 5.2). These facts, together with (5.21), imply that the set in (5.19) is contained in the following finite union:

$$
\bigcup_{\substack{\ell \in\{1, \ldots, d\} \\\left(\boldsymbol{a}_{\ell} \not \perp L_{j}^{\psi}\right)}}\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \boldsymbol{a}_{\ell} \perp L_{j}^{(\boldsymbol{q})}\right\} .
$$

This set has density zero by Lemma 5.4,
Remark 5.2. We have stated Proposition 5.13 in a way which will be convenient for later applications. However let us note that it can be reformulated as follows: We have $\omega_{j}^{(\boldsymbol{q})} \rightarrow$ $\omega_{j}^{\psi}$ in $P\left(\mathbb{T}_{j}^{d}\right)$ as $\boldsymbol{q}$ tends to infinity within a full density subset of $\mathcal{L}_{\psi}$. Indeed, applying Proposition 5.13 to a sequence of shrinking open sets $U$ containing no other element than $\omega_{j}^{\psi}$ in their intersection, and using Lemma 5.14 below, we conclude that there exists a subset $\mathcal{Z} \subset \mathcal{L}_{\psi}$ of density zero such that $\lim _{k \rightarrow \infty} \omega_{j}^{\left(\boldsymbol{q}_{k}\right)}=\omega_{j}^{\psi}$ for any sequence of points $\boldsymbol{q}_{1}, \boldsymbol{q}_{2}, \ldots$ in $\mathcal{L}_{\psi} \backslash \mathcal{Z}$ with $\left\|\boldsymbol{q}_{k}\right\| \rightarrow \infty$.
Lemma 5.14. Let $\mathcal{L}$ be a grid in $\mathbb{R}^{d}$, and let $\mathcal{Z}_{1}, \mathcal{Z}_{2}, \ldots$ be subsets of $\mathcal{L}$ which each have density zero. Then there exists a subset $\mathcal{Z}$ of $\mathcal{L}$ of density zero which has the property that for every $k$ there exists some $R>0$ such that $\mathcal{Z}_{k} \backslash \mathcal{B}_{R}^{d} \subset \mathcal{Z}$.

Proof. Set $z_{j}(R)=R^{-d} \#\left(\mathcal{Z}_{j} \cap \mathcal{B}_{R}^{d}\right)$. Using the fact that $z_{j}(R) \rightarrow 0$ for each $j$, we can pick $0<R_{1}<R_{2}<\cdots$ such that for each $k \in \mathbb{Z}^{+}$and each $R \geq R_{k}$ we have $\sum_{j=1}^{k} z_{j}(R)<k^{-1}$. Now set

$$
\mathcal{Z}=\bigcup_{j=1}^{\infty}\left(\mathcal{Z}_{j} \backslash \mathcal{B}_{R_{j}}^{d}\right)
$$

This set $\mathcal{Z}$ has density zero, since for each $k \in \mathbb{Z}^{+}$and each $R \in\left[R_{k}, R_{k+1}\right)$ we have $R^{-d} \#(\mathcal{Z} \cap$ $\left.\mathcal{B}_{R}^{d}\right)=R^{-d} \# \cup_{j=1}^{k}\left(\mathcal{Z}_{j} \cap \mathcal{B}_{R}^{d} \backslash \mathcal{B}_{R_{j}^{d}}\right) \leq \sum_{j=1}^{k} z_{j}(R)<k^{-1}$. It is also clear that $\mathcal{Z}$ has the last property stated in the lemma.

Remark 5.3. By analogy with Remark 5.2 we also note that Lemma 5.4 implies that there exists a subset $\mathcal{Z} \subset \mathcal{L}_{\psi}$ of density zero such that for every $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$ with $\boldsymbol{a} \not \perp L_{j}^{\psi}$ there is some $R>0$ such that $\boldsymbol{a} \not \perp L_{j}^{(\boldsymbol{q})}$ holds for all $\boldsymbol{q} \in \mathcal{L}_{\psi} \backslash \mathcal{Z}$ with $\|\boldsymbol{q}\|>R$. As in the proof of Proposition 5.13, this implies that the normalized Haar measure on $\left(\mathbb{S}_{j}^{(\boldsymbol{q})}\right)^{\circ}=\pi\left(L_{j}^{(\boldsymbol{q})}\right) \subset \mathbb{T}_{j}$ tends to the normalized Haar measure $\nu_{j}^{\psi}$ on $\mathbb{S}_{j}^{\psi}=\pi\left(L_{j}^{\psi}\right)$ as $\boldsymbol{q}$ tends to infinity within $\mathcal{L}_{\psi} \backslash \mathcal{Z}$. In this sense, we may say that " $L_{j}^{(\boldsymbol{q})}$ approaches $L_{j}^{\psi "}$ as $\boldsymbol{q}$ tends to infinity within a full density subset of $\mathcal{L}_{\psi}$.
5.3. Admissible presentations of $\mathcal{P}$. In this section we will prove that an admissible presentation of $\mathcal{P}$ can always be obtained:

Proposition 5.15. Let $\mathcal{P}$ be a finite union of grids. Then $\mathcal{P}$ possesses an admissible presentation, i.e. there exist $N \in \mathbb{Z}^{+}, r_{1}, \ldots, r_{N} \in \mathbb{Z}^{+}, M_{1}, \ldots, M_{N} \in \mathrm{SL}_{d}(\mathbb{R})$, and numbers $c_{\psi} \in \mathbb{R}$ and vectors $\boldsymbol{w}_{\psi} \in \mathbb{R}^{d}$ for $\psi \in \Psi$ (with $\Psi$ as in (3.3)), such that $\mathcal{P}$ is given by (3.5), (3.4), and this presentation of $\mathcal{P}$ satisfies the admissibility condition in Definition 5.1.

We will prove Proposition 5.15 by showing that there exists a presentation of $\mathcal{P}$ as in (3.5), (3.4), such that (3.6) and (3.7) hold and also

$$
\begin{equation*}
\boldsymbol{c}_{j}^{\psi} \in \mathfrak{L}\left(\left\{\boldsymbol{c}_{j}^{\psi}\right\}\right)+\mathbb{Z}^{r_{j}}, \quad \forall \psi=(j, i) \in \Psi \tag{5.22}
\end{equation*}
$$

This implies that the presentation is admissible, since we always have $\mathfrak{L}\left(\left\{\boldsymbol{c}_{j}^{\psi}\right\}\right) \subset L_{j}^{\psi}$ for all $\psi=(j, i) \in \Psi$, because of (5.5) and the fact that $\mathfrak{L}(S)$ is increasing in $S$.

For any $r \geq 1$ and any vector $\boldsymbol{u}=\left(u_{1}, \ldots, u_{r}\right)^{\boldsymbol{\top}} \in \mathbb{R}_{>0}^{r}$, let us write $\widetilde{\boldsymbol{u}}:=\left(u_{1}^{-1}, \ldots, u_{r}^{-1}\right)^{\boldsymbol{\top}} \in$ $\mathbb{R}_{>0}^{r}$. We call the vector $\boldsymbol{u} \in \mathbb{R}_{>0}^{r}$ admissible if $u_{i} \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)+\mathbb{Z}^{r}$ for every $i \in\{1, \ldots, r\}$. In view of (5.1), we then have that (5.22) is equivalent with the condition that the vector $\left(c_{j, 1}, \ldots, c_{j, r_{j}}\right)^{\top}$ is admissible, for every $j \in\{1, \ldots, N\}$.

To prove Proposition 5.15, we will start from an arbitrary presentation of $\mathcal{P}$ as obtained in Section (3.1, i.e. we assume that $\mathcal{P}$ is expressed as in (3.5), (3.4), and that the conditions (3.6) and (3.7) hold. If this presentation is not already admissible, then we will modify it by making use of the simple fact that for any positive integer $q$, we can express $\mathbb{Z}^{d}$ as the (disjoint) union of the grids $q \mathbb{Z}^{d}+\boldsymbol{\alpha}$, with $\boldsymbol{\alpha}$ running through the set $\{1, \ldots, q\}^{d}$; hence for any $\boldsymbol{w} \in \mathbb{R}^{d}$ and $M \in \mathrm{SL}_{d}(\mathbb{R})$, the $\operatorname{grid}\left(\mathbb{Z}^{d}+\boldsymbol{w}\right) M$ equals the union of the grids $\left(q \mathbb{Z}^{d}+\boldsymbol{\alpha}+\boldsymbol{w}\right) M$. It follows that for any choice of positive integers $q_{\psi}(\psi \in \Psi)$, we have

$$
\begin{equation*}
\mathcal{P}=\bigcup_{\psi \in \Psi} \mathcal{L}_{\psi}=\bigcup_{\psi \in \Psi} c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}=\bigcup_{\psi \in \Psi} \bigcup_{\boldsymbol{\alpha} \in\left\{1, \ldots, q_{\psi}\right\}^{d}} q_{\psi} c_{\psi}\left(\mathbb{Z}^{d}+q_{\psi}^{-1}\left(\boldsymbol{\alpha}+\boldsymbol{w}_{\psi}\right)\right) M_{j_{\psi}} \tag{5.23}
\end{equation*}
$$

In other words, we have obtained a new presentation of $\mathcal{P}$, analogous to the original one:

$$
\begin{equation*}
\mathcal{P}=\bigcup_{\vartheta \in \Theta} c_{\vartheta}^{\prime}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\vartheta}^{\prime}\right) M_{j_{\vartheta}} \tag{5.24}
\end{equation*}
$$

where

$$
\Theta=\left\{(j, i): j \in\{1, \ldots, N\}, i \in\left\{1, \ldots, r_{j}^{\prime}\right\}\right\}
$$

with

$$
r_{j}^{\prime}=\sum_{i=1}^{r_{j}} q_{j, i}^{d} \quad(j=1, \ldots, N)
$$

and where

$$
\boldsymbol{w}_{\vartheta}^{\prime}=q_{\psi(\vartheta)}^{-1}\left(\boldsymbol{\alpha}_{\vartheta}+\boldsymbol{w}_{\psi(\vartheta)}\right) ; \quad c_{\vartheta}^{\prime}=q_{\psi(\vartheta)} c_{\psi(\vartheta)}
$$

with $\psi(\vartheta) \in \Psi$ and $\boldsymbol{\alpha}_{\vartheta} \in \mathbb{Z}^{d}$ chosen in such a way that $j_{\psi(\vartheta)}=j_{\vartheta}$ for all $\vartheta \in \Theta$, and the $\operatorname{map} \vartheta \mapsto\left\langle\psi(\vartheta), \boldsymbol{\alpha}_{\vartheta}\right\rangle$ is a bijection from $\Theta$ onto $\left\{\langle\psi, \boldsymbol{\alpha}\rangle: \psi \in \Psi, \boldsymbol{\alpha} \in\left\{1, \ldots, q_{\psi}\right\}^{d}\right\}$. It is obvious that the new presentation again satisfies the condition (3.7), since the matrices $M_{1}, \ldots, M_{N}$ are unchanged. It is also immediate that the condition (3.6) remains true for the new presentation. Hence it remains to prove that we can choose the positive integers $q_{\psi}$ in such a way that the analogue of $\left(\sqrt{5.22)}\right.$ holds for the new presentation, i.e. so that $\left(c_{j, 1}^{\prime}, \ldots, c_{j, r_{j}^{\prime}}^{\prime}\right)^{\top}$ is admissible, for every $j \in\{1, \ldots, N\}$.

The proof of this fact is essentially completed by lemmas 5.17 and 5.18 below.
Lemma 5.16. Let $1 \leq r \leq r^{\prime}$ and let $T \in \mathrm{M}_{r^{\prime} \times r}(\mathbb{Q})$. Then for any non-empty subset $S \subset \mathbb{R}^{r}$ we have $\mathfrak{L}(T S)=T \mathfrak{L}(S)$.
(Here for any subset $A \subset \mathbb{R}^{r}$ we write $T A:=\{T \boldsymbol{v}: \boldsymbol{v} \in A\}$. Recall also that we view vectors in $\mathbb{R}^{r}$ as column matrices; hence $T \boldsymbol{v} \in \mathbb{R}^{r^{\prime}}$ for every $\boldsymbol{v} \in \mathbb{R}^{r}$.)

Proof. Both $\mathfrak{L}(T S)$ and $T \mathfrak{L}(S)$ are rational subspaces of $\mathbb{R}^{r^{\prime}}$; hence it suffices to prove that the equivalence $[\boldsymbol{a} \perp \mathfrak{L}(T S) \Leftrightarrow \boldsymbol{a} \perp T \mathfrak{L}(S)]$ holds for all $\boldsymbol{a} \in \mathbb{Q}^{r^{\prime}}$. However, $\boldsymbol{a} \perp T \mathfrak{L}(S)$ holds if and only if $\boldsymbol{a} \cdot T \boldsymbol{v}=0$ for all $\boldsymbol{v} \in \mathfrak{L}(S)$, or equivalently $\left(T^{\top} \boldsymbol{a}\right) \cdot \boldsymbol{v}=0$ for all $\boldsymbol{v} \in \mathfrak{L}(S)$. By Lemma 5.3, this holds if and only if there is some $n \in \mathbb{Z}^{+}$such that $T^{\top} \boldsymbol{a} \cdot \boldsymbol{v} \in n^{-1} \mathbb{Z}$ for all $\boldsymbol{v} \in S$. On the other hand, Lemma 5.3 also gives that $\boldsymbol{a} \perp \mathfrak{L}(T S)$ holds if and only if there is some $n \in \mathbb{Z}^{+}$such that $\boldsymbol{a} \cdot T \boldsymbol{v} \in n^{-1} \mathbb{Z}$ for all $\boldsymbol{v} \in S$, or equivalently $T^{\top} \boldsymbol{a} \cdot \boldsymbol{v} \in n^{-1} \mathbb{Z}$ for all $\boldsymbol{v} \in S$. Hence the equivalence is established.
Lemma 5.17. If the vectors $\boldsymbol{u}=\left(u_{1}, \ldots, u_{r}\right)^{\top} \in \mathbb{R}_{>0}^{r}$ and $\boldsymbol{u}^{\prime}=\left(u_{1}^{\prime}, \ldots, u_{r^{\prime}}^{\prime}\right)^{\top} \in \mathbb{R}_{>0}^{r^{\prime}}$ have the same set of coordinates, i.e. $\left\{u_{1}, \ldots, u_{r}\right\}=\left\{u_{1}^{\prime}, \ldots, u_{r^{\prime}}^{\prime}\right\}$, then $\boldsymbol{u}$ is admissible if and only if $\boldsymbol{u}^{\prime}$ is admissible.

Proof. Assume that $\left\{u_{1}, \ldots, u_{r}\right\}=\left\{u_{1}^{\prime}, \ldots, u_{r^{\prime}}^{\prime}\right\}$. This means that there exist uniquely determined matrices $T \in \mathrm{M}_{r, r^{\prime}}(\mathbb{Z})$ and $T^{\prime} \in \mathrm{M}_{r^{\prime}, r}(\mathbb{Z})$ with exactly one entry of 1 in each row and 0 s elsewhere, such that $\boldsymbol{u}^{\prime}=T \boldsymbol{u}$ and $\boldsymbol{u}=T^{\prime} \boldsymbol{u}^{\prime}$; thus also $\widetilde{\boldsymbol{u}^{\prime}}=T \widetilde{\boldsymbol{u}}$ and $\widetilde{\boldsymbol{u}}=T^{\prime} \widetilde{\boldsymbol{u}^{\prime}}$. Assume that $\boldsymbol{u}$ is admissible. Then for every $i \in\{1, \ldots, r\}$ we have $u_{i} \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)+\mathbb{Z}^{r}$. Multiplying this relation by $T$ from the left, and using Lemma 516 and $T \mathbb{Z}^{r} \subset \mathbb{Z}^{r^{\prime}}$, we obtain $u_{i} \widetilde{\boldsymbol{u}^{\prime}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}^{\prime}}\right)+\mathbb{Z}^{r^{\prime}}$. This holds for all $i \in\{1, \ldots, r\}$, and for each $i^{\prime} \in\left\{1, \ldots, r^{\prime}\right\}$ there exists some $i \in\{1, \ldots, r\}$ such that $u_{i^{\prime}}^{\prime}=u_{i}$. Hence $\boldsymbol{u}^{\prime}$ is admissible. The opposite implication is proved analogously, using $T^{\prime}$.

Lemma 5.18. For any $\boldsymbol{u}=\left(u_{1}, \ldots, u_{r}\right)^{\top} \in \mathbb{R}_{>0}^{r}$, there exist positive integers $q_{1}, \ldots, q_{r}$ such that the vector $\left(q_{1} u_{1}, \ldots, q_{r} u_{r}\right)^{\top}$ is admissible.
Proof. Given $\boldsymbol{q}=\left(q_{1}, \ldots, q_{r}\right)^{\top} \in \mathbb{Z}_{>0}^{r}$ we write $D_{\boldsymbol{q}}=\operatorname{diag}\left(q_{1}, \ldots, q_{r}\right)$; then the task is to prove that we there exists some $\boldsymbol{q} \in \mathbb{Z}_{>0}^{r}$ such that $q_{i} u_{i} D_{\boldsymbol{q}}^{-1} \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(q_{i} u_{i} D_{\boldsymbol{q}}^{-1} \widetilde{\boldsymbol{u}}\right)+\mathbb{Z}^{r}$ for all $i \in\{1, \ldots, r\}$. By Lemma 5.16 applied for the matrix $q_{i} D_{q}^{-1} \in \mathrm{M}_{r}(\mathbb{Q})$, this is equivalent to:

$$
\begin{equation*}
q_{i} u_{i} \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)+D_{\boldsymbol{q}} \mathbb{Z}^{r} \quad \text { for all } i \in\{1, \ldots, r\} . \tag{5.25}
\end{equation*}
$$

For each $i, \mathbb{R} \widetilde{\boldsymbol{u}}+\mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$ is a rational subspace of $\mathbb{R}^{r}$, and $u_{i} \widetilde{\boldsymbol{u}} \in \mathbb{Q}^{r}+\mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$ by Lemma 4.3, and $\mathbb{R} \widetilde{\boldsymbol{u}}=\mathbb{R} u_{i} \widetilde{\boldsymbol{u}}$. But $\mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})$ is the unique smallest rational subspace of $\mathbb{R}^{r}$ containing $\mathbb{R} \widetilde{\boldsymbol{u}}$; hence

$$
\begin{equation*}
\mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}}) \subset \mathbb{R} \widetilde{\boldsymbol{u}}+\mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right) . \tag{5.26}
\end{equation*}
$$

We will now describe a choice of $q_{1}, \ldots, q_{r}$ which makes (5.25) hold. Let $p_{i}: \mathbb{R}^{r} \rightarrow \mathbb{R}$ denote projection onto the $i$ th coordinate. Note that $p_{i}\left(\mathbb{Z}^{r} \cap \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})\right)$ is a subgroup of $\mathbb{Z}$. By considering the expansion of $\widetilde{\boldsymbol{u}}$ with respect to a basis of $\mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})$ consisting of vectors in $\mathbb{Z}^{r}$, it follows that $p_{i}\left(\mathbb{Z}^{r} \cap \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})\right) \neq\{0\}$ for each $i$, and hence there exist unique positive integers $q_{i}$ such that

$$
\begin{equation*}
p_{i}\left(\mathbb{Z}^{r} \cap \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})\right)=q_{i} \mathbb{Z} \quad \text { for all } i \in\{1, \ldots, r\} . \tag{5.27}
\end{equation*}
$$

Choose vectors $\boldsymbol{h}^{(i)} \in \mathbb{Z}^{r} \cap \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})$ with $p_{i}\left(\boldsymbol{h}^{(i)}\right)=q_{i}$. We now claim that

$$
\begin{equation*}
q_{i} u_{i} \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)+\boldsymbol{h}^{(i)} \quad \text { for all } i \in\{1, \ldots, r\} . \tag{5.28}
\end{equation*}
$$

This implies that (5.25) holds, since $\boldsymbol{h}^{(i)} \in \mathbb{Z}^{r} \cap \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}}) \subset D_{\boldsymbol{q}} \mathbb{Z}^{r}$, where the last inclusion follows from (5.27).

In order to prove (5.28), let $i$ be given, and set $\boldsymbol{w}:=q_{i} u_{i} \widetilde{\boldsymbol{u}}-\boldsymbol{h}^{(i)}$. We have $\boldsymbol{w} \in \mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})$, since $\widetilde{\boldsymbol{u}}$ and $\boldsymbol{h}^{(i)}$ lie in $\mathfrak{L}(\mathbb{R} \widetilde{\boldsymbol{u}})$; hence by (5.26) there exists some $t \in \mathbb{R}$ such that $\boldsymbol{w}-t \widetilde{\boldsymbol{u}} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$. However, both $\boldsymbol{w}$ and $\mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$ lie in the orthogonal complement of $\boldsymbol{e}_{i}$, the $i$ th standard unit vector in $\mathbb{R}^{r}$; for $\boldsymbol{w}$ this holds since $\boldsymbol{e}_{i} \cdot \widetilde{\boldsymbol{u}}=u_{i}$ and $\boldsymbol{e}_{i} \cdot \boldsymbol{h}^{(i)}=p_{i}\left(\boldsymbol{h}^{(i)}\right)=q_{i}$, and for $\mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$ it holds by Lemma 5.3, since $\boldsymbol{e}_{i} \cdot u_{i} \widetilde{\boldsymbol{u}}=1$. On the other hand, $\widetilde{\boldsymbol{u}}$ is not orthogonal to $\boldsymbol{e}_{i}$. Hence we must have $t=0$, i.e. $\boldsymbol{w} \in \mathfrak{L}\left(u_{i} \widetilde{\boldsymbol{u}}\right)$, and so (5.28) holds.

Now to finally prove Proposition 5.15 we proceed as follows, for each fixed $j \in\{1, \ldots, N\}$ : Choose $0<u_{1}<\cdots<u_{r}$ so that $\left\{c_{j, 1}, \ldots, c_{j, r_{j}}\right\}=\left\{u_{1}, \ldots, u_{r}\right\}$, and let $\tau:\left\{1, \ldots, r_{j}\right\} \rightarrow$ $\{1, \ldots, r\}$ be the map so that $c_{j, i}=u_{\tau(i)}$ for all $i \in\left\{1, \ldots, r_{j}\right\}$; then by Lemma 5.18 there exist positive integers $q_{1}, \ldots, q_{r}$ such that the vector $\left(q_{1} u_{1}, \ldots, q_{r} u_{r}\right)^{\top}$ is admissible. Setting now $q_{j, i}=q_{\tau(i)}$, Lemma 5.17 implies that the vector $\left(q_{j, 1} c_{j, 1}, \ldots, q_{j, r_{j}} c_{j, r_{j}}\right)^{\top}$ is admissible; and since $c_{j, i}^{\prime}=q_{\psi(j, i)} c_{\psi(j, i)}$ for all $i=1, \ldots, r_{j}^{\prime}$, another application of Lemma 5.17 gives that also the vector $\left(c_{j, 1}^{\prime}, \ldots, c_{j, r_{j}^{\prime}}^{\prime}\right)^{\top}$ is admissible, which was exactly the desired condition. This completes the proof of Proposition 5.15.

## 6. Precise statement of the main Result

In this section the statement of the main technical result of the paper, Theorem [2.1, will be made more explicit, by presenting the precise choice of the data $\varsigma, \Sigma, \mathrm{m}$ and $\varsigma \mapsto \mu_{\varsigma}$ for which we will later prove that the six conditions [P1]-[P3] and [Q1]-[Q3] in Section 2 hold.

Let $\mathcal{P}$ be a finite union of grids in $\mathbb{R}^{d}$, and fix an admissible presentation of $\mathcal{P}$ as in (3.5), (3.4). This is possible by Proposition 5.15,
6.1. The subspace $\mathbb{X}^{\psi} \subset \mathbb{X}$. We first need to make a preparation of a technical nature. For any $\psi \in \Psi$ we let $\mathbb{X}^{\psi}$ be the closed set consisting of those $\Gamma g \in \mathbb{X}$ for which $\mathbb{Z}^{d} p_{\psi}(g)$ is a lattice, i.e.

$$
\begin{equation*}
\mathbb{X}^{\psi}:=\left\{\Gamma g \in \mathbb{X}: g \in G, \mathbf{0} \in \mathbb{Z}^{d} p_{\psi}(g)\right\} . \tag{6.1}
\end{equation*}
$$

(Recall that the grid $\mathbb{Z}^{d} p_{\psi}(g)$ is independent of the choice of representative $g$, i.e. if $\Gamma g=\Gamma g^{\prime}$ then $\mathbb{Z}^{d} p_{\psi}(g)=\mathbb{Z}^{d} p_{\psi}\left(g^{\prime}\right)$.) Clearly, for $\psi=(j, i)$,

$$
\begin{equation*}
\mathbb{X}^{\psi}=p_{j}^{-1}\left(\mathbb{X}_{j}^{(i)}\right), \quad \text { with } \mathbb{X}_{j}^{(i)}=\left\{\Gamma_{j} g: g \in G_{j}, \mathbf{0} \in \mathbb{Z}^{d} \mathrm{a}_{i}(g)\right\} . \tag{6.2}
\end{equation*}
$$

Lemma 6.1. For any $\psi \in \Psi$ and $\boldsymbol{q} \in \mathcal{L}_{\psi}, \Gamma g_{0}^{(\boldsymbol{q})} \in \mathbb{X}^{\psi}$.
Proof. This is immediate from (3.20).
We will write $\delta_{0}$ for the Dirac measure on $(\mathbb{R} / \mathbb{Z})^{d}$ at the point $\mathbf{0} \in(\mathbb{R} / \mathbb{Z})^{d}$. Recall that $\widetilde{\mathrm{r}}_{i}: \mathbb{T}_{j}^{d} \rightarrow(\mathbb{R} / \mathbb{Z})^{d}$ is the projection induced by the $i$ th row map $\mathrm{r}_{i}: \mathrm{M}_{r_{j} \times d}(\mathbb{R}) \rightarrow \mathbb{R}^{d}$.
Lemma 6.2. Assume that $\psi=(j, i) \in \Psi, \omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ and $\widetilde{\mathrm{r}}_{i} * \omega=\delta_{\mathbf{0}}$. Then $\bar{\omega}\left(\mathbb{X}_{j}^{(i)}\right)=1$.
Proof. In view of the definition of $\bar{\omega}$ in (4.9), it suffices to verify that $\Gamma_{j} \mathrm{I}_{U} A \in \mathbb{X}_{j}^{(i)}$ for all $U \in \mathrm{r}_{i}^{-1}\left(\mathbb{Z}^{d}\right) \subset \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ and all $A \in \mathrm{SL}_{d}(\mathbb{R})$. The verification of this fact is immediate.
Lemma 6.3. Let $\psi=(j, i) \in \Psi$ and $V \in \mathbb{T}_{j}^{d}$. If $\widetilde{r}_{i}(V)=\mathbf{0}$ in $(\mathbb{R} / \mathbb{Z})^{d}$ then $\widetilde{\mathfrak{r}}_{i *} \omega_{j}^{(V)}=\delta_{\mathbf{0}}$, while if $\widetilde{\mathrm{r}}_{i}(V) \neq \mathbf{0}$ then $\left(\widetilde{\mathrm{r}}_{i} * \omega_{j}^{(V)}\right)(\{\mathbf{0}\})=0$.
Proof. Write $V=\left(V_{1}, \ldots, V_{d}\right)$. If $\widetilde{\mathrm{r}}_{i}(V)=\mathbf{0}$, then the $i$ th coordinate of each $V_{\ell}$ vanishes; hence all points in $\mathbb{S}_{j}^{(V)}$ have vanishing $i$ th coordinate, and so $r_{i}\left(\left(\mathbb{S}_{j}^{(V)}\right)^{d}\right)=\{\mathbf{0}\}$ in $(\mathbb{R} / \mathbb{Z})^{d}$, which forces $\omega_{j}^{(V)}\left(\widetilde{r}_{i}^{-1}(\{\mathbf{0}\})\right)=1$, i.e., $\widetilde{r}_{i} * \omega_{j}^{(V)}=\delta_{\mathbf{0}}$.

Next assume $\widetilde{\mathrm{r}}_{i}(V) \neq \mathbf{0}$. If $L_{j}^{(V)} \nsucceq \boldsymbol{e}_{i}$, then $\widetilde{\mathrm{r}}_{i}\left(\left(\mathbb{S}_{j}^{(V)}\right)^{\circ d}\right)=(\mathbb{R} / \mathbb{Z})^{d}$ and so by (4.7) and the definition of $\omega_{j}^{(V)}$, we have $\omega_{j}^{(V)}\left(\widetilde{\mathbf{r}}_{i}^{-1}(\{\mathbf{0}\})\right)=0$. In the remaining case, when $L_{j}^{(V)} \perp \boldsymbol{e}_{i}$ and therefore $\widetilde{\mathrm{r}}_{i}\left(\left(\mathbb{S}_{j}^{(V)}\right)^{\circ d}\right)=\{\mathbf{0}\}$, we again have $\omega_{j}^{(V)}\left(\widetilde{r}_{i}^{-1}(\{\mathbf{0}\})\right)=0$, since $\widetilde{\mathrm{r}}_{i}(V \gamma)=\widetilde{\mathrm{r}}_{i}(V) \gamma \neq \mathbf{0}$ for all $\gamma \in \mathrm{SL}_{d}(\mathbb{Z})$.
Lemma 6.4. Let $\psi=(j, i) \in \Psi$ and $\boldsymbol{q} \in \mathbb{R}^{d}$. If $\boldsymbol{q} \in \mathcal{L}_{\psi}$ then $\widetilde{\mathfrak{r}}_{i *} \omega_{j}^{(\boldsymbol{q})}=\delta_{\mathbf{0}}$, while if $\boldsymbol{q} \notin \mathcal{L}_{\psi}$ then $\left(\widetilde{\mathrm{r}}_{i *} \omega_{j}^{(\boldsymbol{q})}\right)(\{\mathbf{0}\})=0$.
Proof. Note that $\boldsymbol{q} \in \mathcal{L}_{\psi}$ holds if and only if $r_{i}\left(U_{j}^{(\boldsymbol{q})}\right) \in \mathbb{Z}^{d}$ (cf. (3.4) and (3.15)), that is, if and only if $\widetilde{\mathrm{r}}_{i}\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)=\mathbf{0}$ in $(\mathbb{R} / \mathbb{Z})^{d}$. Hence the lemma follows from Lemma 6.3,
Lemma 6.5. Let $j \in\{1, \ldots, N\}$ and $i \neq i^{\prime} \in\left\{1, \ldots, r_{j}\right\}$. Then $L_{j}^{\left(j, i^{\prime}\right)} \perp \boldsymbol{e}_{i}$ holds if and only if $L_{j}^{(j, i)} \perp \boldsymbol{e}_{i^{\prime}}$, and in this case also $L_{j}^{\left(j, i^{\prime}\right)}=L_{j}^{(j, i)}$ and $c_{j, i}=c_{j, i^{\prime}}$ and $\boldsymbol{w}_{j, i}-\boldsymbol{w}_{j, i^{\prime}} \notin \mathbb{Z}^{d}$.
Proof. Assume $L_{j}^{\left(j, i^{\prime}\right)} \perp \boldsymbol{e}_{i}$. Then by Lemma[5.7, $L_{j}^{\left(j, i^{\prime}\right)} \subset L_{j} \cap \boldsymbol{e}_{i}^{\perp}=L_{j}^{(j, i)}$, and both $L_{j}^{\left(j, i^{\prime}\right)}$ and $L_{j}^{(j, i)}$ are subspaces of $L_{j}$ of codimension one; hence $L_{j}^{\left(j, i^{\prime}\right)}=L_{j}^{(j, i)}$, and since $L_{j}^{\left(j, i^{\prime}\right)} \perp \boldsymbol{e}_{i^{\prime}}$ by Lemma 5.7, we conclude that $L_{j}^{(j, i)} \perp \boldsymbol{e}_{i^{\prime}}$. Of course, the converse implication, $L_{j}^{(j, i)} \perp \boldsymbol{e}_{i^{\prime}} \Rightarrow$
$L_{j}^{\left(j, i^{\prime}\right)} \perp \boldsymbol{e}_{i}$, holds by symmetry. When both these orthogonality relations hold, admissibility (see Definition 5.11) implies that both the $i$ th coordinate of $\boldsymbol{c}_{j}^{\left(j, i^{\prime}\right)}$ and the $i^{\prime}$ th coordinate of $\boldsymbol{c}_{j}^{(j, i)}$ are integers, i.e. both $c_{j, i^{\prime}} / c_{j, i}$ and $c_{j, i} / c_{j, i^{\prime}}$ are integers; thus $c_{j, i^{\prime}}=c_{j, i}$. Finally, by (3.6), $\boldsymbol{w}_{j, i}-\boldsymbol{w}_{j, i^{\prime}} \notin \mathbb{Z}^{d}$.

Lemma 6.6. For every $\psi \in \Psi$ we have $\widetilde{\mathbf{r}}_{i *} \omega_{j}^{\psi}=\delta_{\mathbf{0}}$ when $(j, i)=\psi$, while $\left(\widetilde{\mathbf{r}}_{i *} \omega_{j}^{\psi}\right)(\{\mathbf{0}\})=0$ for every $(j, i) \in \Psi \backslash\{\psi\}$.

Proof. The first statement follows from Lemma 6.4, since $\omega_{j}^{\psi}$ can be obtained as a limit of measures $\omega_{j}^{(\boldsymbol{q})}$ for a sequence of points $\boldsymbol{q}$ in $\mathcal{L}_{\psi}$ (cf. Remark 5.2), and since the map $\widetilde{\mathrm{r}}_{i *}$ : $P\left(\mathbb{T}_{j}^{d}\right) \rightarrow P\left((\mathbb{R} / \mathbb{Z})^{d}\right)$ is continuous. Alternatively it is easy to give a direct verification: In (5.1) and (5.2) one notes that the $i$ th coordinate of $\boldsymbol{c}_{j}^{\psi}$ equals 1 ; thus $\mathrm{r}_{i}\left(W_{j}^{\psi}\right)=\mathbf{0}$, and so by (5.5), $L_{j}^{\psi} \perp \boldsymbol{e}_{i}$, meaning that every point in $\mathbb{S}_{j}^{\psi}$ has a vanishing $i$ th coordinate. We know from the proof of Lemma 6.4 that the same fact holds for every point in $\mathbb{S}_{j}^{(\boldsymbol{q})}$, for any $\boldsymbol{q} \in \mathcal{L}_{\psi}$; hence by (5.16) it also holds for every point in $\widetilde{\mathbb{S}}_{j}^{\psi}$, and so $\widetilde{\mathrm{r}}_{i}\left(\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}\right)=\{\mathbf{0}\}$ and $\widetilde{\mathrm{r}}_{i *} \omega_{j}^{\psi}=\delta_{\mathbf{0}}$.

We turn to the proof of the second statement (this is similar to the proof of the second half of Lemma 6.3). Thus let $(j, i) \in \Psi,(j, i) \neq \psi$. If $L_{j}^{\psi} \not \perp \boldsymbol{e}_{i}$ then $\widetilde{\mathfrak{r}}_{i}\left(\left(\mathbb{S}_{j}^{\psi}\right)^{d}\right)=(\mathbb{R} / \mathbb{Z})^{d}$ and hence $\omega_{j}^{\psi}\left(\tilde{r}_{i}^{-1}(\{\mathbf{0}\})\right)=0\left(\right.$ cf. (5.16) ). Now assume $L_{j}^{\psi} \perp \boldsymbol{e}_{i}$. Then $j=j_{\psi}$, by (5.5). Pick an arbitrary point $\boldsymbol{q} \in \mathcal{L}_{\psi}$, so that (5.16) holds. Now by Lemma 6.5, $c_{\psi}=c_{j, i}$ and $\boldsymbol{w}_{\psi}-\boldsymbol{w}_{j, i} \notin \mathbb{Z}^{d}$, and so via (5.2) and (5.4), $\mathrm{r}_{i}\left(U_{j}^{(\boldsymbol{q})}\right) \notin \mathbb{Z}^{d}$, and hence $\widetilde{\mathrm{r}}_{i}\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right) \gamma\right) \neq \mathbf{0}$ for all $\gamma \in \mathrm{SL}_{d}(\mathbb{Z})$. Also $\widetilde{\mathbf{r}}_{i}\left(\mathbb{S}_{j}^{\psi}\right)=\{\mathbf{0}\}$, since $L_{j}^{\psi} \perp \boldsymbol{e}_{i}$. Hence again $\omega_{j}^{\psi}\left(\widetilde{\mathrm{r}}_{i}^{-1}(\{\mathbf{0}\})\right)=0$ (see (5.16)).
6.2. The space of marks $\Sigma$ and the associated maps. We are now finally in a position to introduce our precise choice of space of marks $\Sigma$ and the associated maps and measure.

To prepare for these definitions, set

$$
\begin{equation*}
\Omega=\prod_{j=1}^{N} P\left(\mathbb{T}_{j}^{d}\right)^{\prime} \tag{6.3}
\end{equation*}
$$

For any $\boldsymbol{q} \in \mathbb{R}^{d}$ we define

$$
\begin{equation*}
\omega^{(\boldsymbol{q})}:=\left(\omega_{1}^{(\boldsymbol{q})}, \ldots, \omega_{N}^{(\boldsymbol{q})}\right) \in \Omega \tag{6.4}
\end{equation*}
$$

and for any $\psi \in \Psi$ we define

$$
\begin{equation*}
\omega^{\psi}:=\left(\omega_{1}^{\psi}, \ldots, \omega_{N}^{\psi}\right) \in \Omega . \tag{6.5}
\end{equation*}
$$

Now we define our space of marks $\Sigma$ through:

$$
\begin{equation*}
\Sigma:=\left\{((j, i), \omega) \in \Psi \times \Omega: \widetilde{\mathrm{r}}_{i *} \omega_{j}=\delta_{0}\right\} \tag{6.6}
\end{equation*}
$$

where $\omega_{j} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ is the $j$ th entry of $\omega$, and where $\delta_{0}$ is the Dirac measure at the point $\mathbf{0} \in(\mathbb{R} / \mathbb{Z})^{d}$. It is to be understood that in (6.6), $\Psi$ is equipped with the discrete topology, and $\Psi \times \Omega$ with the product topology; this makes $\Sigma$ a closed and hence compact subset of $\Psi \times \Omega$.

Remark 6.1. The reason why we cannot simply choose $\Sigma$ to be $\Psi \times \Omega$ is that then the map $\varsigma \mapsto \mu_{\varsigma}$ which we define below would in general not be continuous.

Next we define our marking $\varsigma$ through

$$
\begin{equation*}
\varsigma: \mathcal{P} \rightarrow \Sigma, \quad \varsigma(\boldsymbol{q})=\left(\psi(\boldsymbol{q}), \omega^{(\boldsymbol{q})}\right) \tag{6.7}
\end{equation*}
$$

where $\psi: \mathcal{P} \rightarrow \Psi$ is as in (1.4), i.e. a fixed function such that $\boldsymbol{q} \in \mathcal{L}_{\psi(\boldsymbol{q})}$ for all $\boldsymbol{q} \in \mathcal{P}$. It follows from Lemma 6.4 that $\varsigma$ is indeed a map into $\Sigma$.

For each $\psi \in \Psi$ we set

$$
\begin{equation*}
\sigma^{\psi}:=\left(\psi, \omega^{\psi}\right) \in \Sigma . \tag{6.8}
\end{equation*}
$$

It follows from Lemma 6.6 that $\sigma^{\psi}$ indeed lies in $\Sigma$. Next we define the Borel probability measure m on $\Sigma$ to be the atomic measure on $\Sigma$ supported on the points $\sigma^{\psi}$, with

$$
\begin{equation*}
\mathrm{m}\left(\sigma^{\psi}\right):=\frac{\bar{n}_{\psi}}{\bar{n}_{\mathcal{P}}} \quad(\psi \in \Psi) \tag{6.9}
\end{equation*}
$$

Note that $m$ is supported on the finite subset

$$
\begin{equation*}
\widetilde{\Psi}:=\left\{\sigma^{\psi}: \psi \in \Psi\right\} \tag{6.10}
\end{equation*}
$$

of $\Sigma$.
As in Section 2 we set

$$
\mathcal{X}=\mathbb{R}^{d} \times \Sigma \quad \text { and } \quad \mu_{\mathcal{X}}=\operatorname{vol} \times \mathrm{m}
$$

Finally we will make our choice of the map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P(N(\mathcal{X}))$. To prepare for this, note that we have a natural map $\omega \mapsto \bar{\omega}$, from $\Omega$ to $P(\mathbb{X})$, defined by

$$
\begin{equation*}
\bar{\omega}:=\overline{\omega_{1}} \otimes \cdots \otimes \overline{\omega_{N}} \in P(\mathbb{X}) \quad(\omega \in \Omega) \tag{6.11}
\end{equation*}
$$

This makes sense since $\overline{\omega_{j}} \in P\left(\mathbb{X}_{j}\right)$ for each $j$. Next we introduce the following map:

$$
\begin{equation*}
J: \mathbb{X} \rightarrow N_{s}(\mathcal{X}), \quad J(\Gamma g)=\bigcup_{\psi \in \Psi} c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right) \times\left\{\sigma^{\psi}\right\} \quad(g \in G) \tag{6.12}
\end{equation*}
$$

This map extends the map $J_{0}$ defined in (3.14) in an obvious sense.
Lemma 6.7. $J$ is continuous.
Proof. Recall that $N_{s}(\mathcal{X})$ is equipped with the vague topology; hence the task is to prove that if $x_{1}, x_{2}, \ldots$ is a sequence in $\mathbb{X}$ converging to $x \in \mathbb{X}$, then for any $f \in \mathrm{C}_{c}(\mathcal{X})$ we have $\sum_{p \in J\left(x_{k}\right)} f(p) \rightarrow \sum_{p \in J(x)} f(p)$. This is immediate using the formula

$$
\sum_{p \in J(\Gamma g)} f(p)=\sum_{\psi \in \Psi} \sum_{\boldsymbol{m} \in \mathbb{Z}^{d}} f\left(c_{\psi}\left(\boldsymbol{m} p_{\psi}(g)\right), \sigma^{\psi}\right), \quad \forall g \in G
$$

Next, for every $\psi \in \Psi$ we introduce the following modification of the map $J$ :

$$
\begin{equation*}
J_{\psi}: \mathbb{X}^{\psi} \rightarrow N_{s}(\mathcal{X}) ; \quad J_{\psi}(\Gamma g):=J(\Gamma g) \backslash\left\{\left(\mathbf{0}, \sigma^{\psi}\right)\right\} \tag{6.13}
\end{equation*}
$$

This map $J_{\psi}$ is also continuous; this is proved in the same way as Lemma 6.7, using also the fact that $\left(\mathbf{0}, \sigma^{\psi}\right) \in J(x)$ for all $x \in \mathbb{X}^{\psi}$ (see (6.1) and (6.12)). At last, we now define the map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P\left(N_{s}(\mathcal{X})\right)$ by setting

$$
\begin{equation*}
\mu_{\varsigma}=J_{\psi *} \bar{\omega} \in P\left(N_{s}(\mathcal{X})\right) \quad \text { for } \varsigma=(\psi, \omega) \in \Sigma \tag{6.14}
\end{equation*}
$$

To see that $\mu_{\varsigma}$ is indeed a probability measure, note that

$$
\begin{equation*}
\bar{\omega}\left(\mathbb{X}^{\psi}\right)=1, \quad \text { for all } \varsigma=(\psi, \omega) \in \Sigma ; \tag{6.15}
\end{equation*}
$$

cf. (6.2), Lemma 6.2, (6.6) and (6.11).
Lemma 6.8. The map $\varsigma \mapsto \mu_{\varsigma}$ is continuous.
Proof. Let $\varsigma_{1}, \varsigma_{2}, \ldots$ be an arbitrary sequence in $\Sigma$ converging to a point $\varsigma \in \Sigma$. Write $\varsigma_{k}=$ $\left(\psi^{(k)}, \omega^{(k)}\right)$ with $\omega^{(k)}=\left(\omega_{1}^{(k)}, \ldots, \omega_{N}^{(k)}\right) \in \Omega$; also write $\varsigma=(\psi, \omega)$. Throwing away finitely many initial points from the sequence we may assume that $\psi^{(k)}=\psi$ for all $k$. By (6.15) we then have $\bar{\omega}\left(\mathbb{X}^{\psi}\right)=1$ and $\overline{\omega^{(k)}}\left(\mathbb{X}^{\psi}\right)=1$ for all $k$. Hence we may just as well regard $\bar{\omega}$ and all $\overline{\omega^{(k)}}$ as elements in $P\left(\mathbb{X}^{\psi}\right)$. For each fixed $j \in\{1, \ldots, N\}$ we have $\omega_{j}^{(k)} \rightarrow \omega_{j}$ in $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ as $k \rightarrow \infty$; hence by (6.11), Lemma 4.6 and [1, Thm. 2.8(ii)], $\overline{\omega^{(k)}} \rightarrow \bar{\omega}$ in $P(\mathbb{X})$.

Hence also $\overline{\omega^{(k)}} \rightarrow \bar{\omega}$ in $P\left(\mathbb{X}^{\psi}\right)$ [13, Lemma 4.26], and by the continuous mapping theorem, $J_{\psi *} \overline{\omega^{(k)}} \rightarrow J_{\psi *} \bar{\omega}$, in $P\left(N_{s}(\mathcal{X})\right)$, viz., $\mu_{\varsigma_{k}} \rightarrow \mu_{\varsigma}$.
6.3. Some comments on the construction of $\Sigma$. The facts mentioned in the present section will not be used later in the text, and we will omit some proofs.

The construction of the space of marks $\Sigma$ in the present paper is more complicated and nonintuitive than for any of the previous types of scatterer configurations for which the framework from [21] is known to apply; see [21, Ch. 5]. In the special cases when $\mathcal{P}$ is periodic or when the lattices in $\mathcal{P}$ are pairwise incommensurable, we could have chosen the space of marks to be simply equal to the indexing set $\Psi$; however this simple choice of $\Sigma$ is not possible in the general case; see Remarks 6.2, 6.3 and 6.4 below.

For a general locally finite subset $\mathcal{P} \subset \mathbb{R}^{d}$, when trying to verify the hypotheses from [21, Sec. 2.3] (which we recalled in Section 2) without an apriori guess of what the space of marks $\Sigma$ should be, a natural first task is to ignore the marking, thus thinking of each measure $\mu_{\boldsymbol{q}, \boldsymbol{\rho}}^{(\lambda)}$ as lying in $N_{s}\left(\mathbb{R}^{d}\right), 8$ and then verify that for any fixed $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ and $\boldsymbol{q} \in \mathcal{P}$ (possibly excluding a subset $\mathcal{E} \subset \mathcal{P}$ of density zero), the weak limit of $\mu_{\boldsymbol{q}, \rho}^{(\lambda)}$ as $\rho \rightarrow 0$ exists and is independent of $\lambda$; call this limit measure $\mu_{\boldsymbol{q}, 0} \in P\left(N\left(\mathbb{R}^{d}\right)\right)$. Then by Lemma 2.2, the space of marks $\Sigma$ must admit a continuous map $\varsigma \mapsto \widehat{\mu}_{\varsigma}$ from $\Sigma$ to $P\left(N\left(\mathbb{R}^{d}\right)\right)$ such that $\widehat{\mu}_{\varsigma(\boldsymbol{q})}=\mu_{\boldsymbol{q}, 0}$ for all $\mathcal{P} \backslash \mathcal{E}$.

This to a large extent motivates our choice of $\Sigma$ in Section 6.2, We have constructed the space $\Omega=\prod_{j=1}^{N} P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ to encode all the spherical equidistribution limits $\mu_{\boldsymbol{q}, 0}$ in $P\left(N\left(\mathbb{R}^{d}\right)\right)$. Indeed, it follows from Theorem4.2 that ${ }^{9} \mu_{\boldsymbol{q}, 0}=\left(J_{0, \psi(\boldsymbol{q})}\right)_{*}\left(\overline{\omega^{(\boldsymbol{q})}}\right)$ for any $\boldsymbol{q} \in \mathcal{P} \backslash \cup_{\psi \neq \psi^{\prime} \in \Psi}\left(\mathcal{L}_{\psi} \cap\right.$ $\mathcal{L}_{\psi^{\prime}}$, where for any $\psi \in \Psi, J_{0, \psi}$ is the following modified version of the map defined in (3.14):

$$
\begin{equation*}
J_{0, \psi}: \mathbb{X} \rightarrow N\left(\mathbb{R}^{d}\right) ; \quad J_{0, \psi}(\Gamma g):=\sum_{\psi^{\prime} \in \Psi} \sum_{\substack{x \in c_{\psi^{\prime}}\left(\mathbb{Z}^{d} p_{\psi^{\prime}}(g)\right) \\\left(\text { if } \psi^{\prime}=\psi: \boldsymbol{x} \neq \mathbf{0}\right)}} \delta_{\boldsymbol{x}} \quad(g \in G) \tag{6.16}
\end{equation*}
$$

(That is, $J_{0, \psi}$ is as in (3.14) except that we (1) count multiplicity, and (2) if possible remove the origin from $c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right)$.) This suggests that we could make the choice " $\Sigma=\Omega$ ". It is technically convenient to also include the factor $\Psi$ in the definition of $\Sigma$; and we then restrict by the condition $\widetilde{\mathrm{r}}_{i *} \omega_{j}=\delta_{\mathbf{0}}$ (see (6.6)) to ensure that the map $\varsigma \mapsto \mu_{\varsigma}$ be continuous.

However, even though this space $\Sigma$ by construction classifies all the relevant spherical equidistribution limits in $P\left(N\left(\mathbb{R}^{d}\right)\right.$ ), it might still be too crude to classify the corresponding limits in the refined space $P(N(\mathcal{X}))$ ! This is exactly the reason why we require that the presentation of $\mathcal{P}$ should be admissible (Definition 5.1); only under this assumption have we managed to find a proof of the condition [P2]. The crucial property that we need (see the proof of Lemma 7.11 below) is that for each fixed $\psi \in \Psi, \omega^{(\boldsymbol{q})}$ tends to a unique limit $\omega^{\psi}$ as $\boldsymbol{q} \rightarrow \infty$ within $\mathcal{L}_{\psi} \backslash \mathcal{E}$; see Section 5.2 where the admissibility is needed already to define $\omega_{j}^{\psi}$.

We remark that it would in fact be possible to take the space $\Sigma$ to be a significantly smaller set than we have done 10 Indeed, assuming that the conditions [P1]-[P3] and [Q1]-[Q3] hold for the data $\left[\Sigma, \varsigma, \varsigma \mapsto \mu_{\varsigma}, \mathrm{m}\right]$, and choosing $\mathcal{E} \subset \mathcal{P}$ as in condition [P2], set

$$
\begin{equation*}
\Sigma^{\prime}:=\overline{\{\varsigma(\boldsymbol{q}): \boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}\}} . \tag{6.17}
\end{equation*}
$$

This is a compact subset of $\Sigma$, and as explained in [21, Remark 2.6], after modifying the marking of the points in $\mathcal{E}$ in an arbitrary fashion to ensure that $\varsigma(\boldsymbol{q}) \in \Sigma^{\prime}$ for all $\boldsymbol{q} \in \mathcal{P}$, the conditions $[\mathrm{P} 1]-[\mathrm{P} 3]$ and $[\mathrm{Q} 1]-[\mathrm{Q} 3]$ remain valid with $\Sigma^{\prime}$ in place of $\Sigma$. It is also important to

[^6]note that if [P2] holds for a certain set $\mathcal{E} \subset \mathcal{P}$, then [P2] also holds for any larger subset $\mathcal{E} \subset \mathcal{P}$ of density zero. In our setting, with $\Sigma$ chosen as in Section 6.2, it follows from Remark 5.2 that $\widetilde{\Psi} \subset \Sigma^{\prime}$, but also that, by taking $\mathcal{E}$ appropriately large, we can ensure that the set $\{\varsigma(\boldsymbol{q}): \boldsymbol{q} \in \mathcal{P}\}$ has no limit point outside $\widetilde{\Psi}$. This means that
\[

$$
\begin{equation*}
\Sigma^{\prime}=\{\varsigma(\boldsymbol{q}): \boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}\} \cup \widetilde{\Psi}, \tag{6.18}
\end{equation*}
$$

\]

and that $\Sigma^{\prime}$ has quite a simple structure as a (metrizable, compact) topological space: $\Sigma^{\prime}$ is a finite or countable set containing the finite set $\widetilde{\Psi}$, and its set of limit points is contained in $\widetilde{\Psi}$.

Remark 6.2. The description of $\Sigma^{\prime}$ in the last sentence is essentially as simple as it can be: For the scatterer configurations studied in the present paper, it is in general impossible to satisfy the conditions in Section 2 using a finite space of marks. For example, consider the configuration

$$
\mathcal{P}=\mathbb{Z}^{2} \cup\left(\mathbb{Z}^{2}+(0, \sqrt{2})\right) \cup \mathbb{Z}^{2} M_{2} \quad \text { in } \mathbb{R}^{2}, \quad \text { with } M_{2}=\left(\begin{array}{cc}
1 & \sqrt{2}  \tag{6.19}\\
1 & \sqrt{2}+1
\end{array}\right) \in \mathrm{SL}_{2}(\mathbb{R}) .
$$

In this case, one obtains infinitely many distinct spherical equidistribution limits $\mu_{\boldsymbol{q}, 0}$ in $P\left(N\left(\mathbb{R}^{2}\right)\right)$ as $\boldsymbol{q}$ varies through the lattice $\mathbb{Z}^{2} M_{2} \subset \mathcal{P}$. Indeed, for two arbitrary points $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathbb{Z}^{2} M_{2}$, writing $\boldsymbol{q}=\boldsymbol{m} M_{2}$ with $\boldsymbol{m}=\left(m_{1}, m_{2}\right) \in \mathbb{Z}^{2}$ and similarly $\boldsymbol{q}^{\prime}=\boldsymbol{m}^{\prime} M_{2}$, it turns out that

$$
\begin{equation*}
\mu_{\boldsymbol{q}, 0}=\mu_{\boldsymbol{q}^{\prime}, 0} \quad \Leftrightarrow \quad\left[m_{1}+m_{2}=m_{1}^{\prime}+m_{2}^{\prime} \text { or } m_{1}+m_{2}=1-\left(m_{1}^{\prime}+m_{2}^{\prime}\right)\right] \tag{6.20}
\end{equation*}
$$

We give an outline of the proof below. It follows from (6.20) that for any set $\mathcal{E} \subset \mathcal{P}$ of density zero, the set of spherical equidistribution limits $\left\{\mu_{\boldsymbol{q}, 0}: \boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}\right\}$ is infinite. (Indeed, if it were finite, then by (6.20) the point set $\mathbb{Z}^{2} M_{2} \backslash \mathcal{E}$ would be contained in a finite number of lines and thus be of density zero, whereas in fact it has density one.) Hence as in the discussion in the beginning of this subsection, it follows from Lemma 2.2 that for any marking of $\mathcal{P}$ used to satisfy the conditions [P1]-[P3] and [Q1]-[Q3] in Section 2, the set of marks $\left\{\varsigma(\boldsymbol{q}): \boldsymbol{q} \in \mathcal{L}_{(2,1)} \backslash \mathcal{E}\right\}$ must be infinite, and in particular the space of marks must be infinite. Of course, this also implies that the space of marks must have some limit point, since the space of marks is required to be compact.

Outline of proof of (6.20). We express $\mathcal{P}$ in (6.19) as in (3.3)-(3.5) with $N=2, r_{1}=2$, $r_{2}=1, M_{1}=\left(\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right)$, and $M_{2}$ as above; also $c_{1,1}=c_{1,2}=c_{2,1}=1$ and $\boldsymbol{w}_{1,1}=\boldsymbol{w}_{2,1}=(0,0)$ and $\boldsymbol{w}_{1,2}=(0, \sqrt{2})$. For any point $\boldsymbol{q}=\boldsymbol{m} M_{2}$ in $\mathcal{L}_{(2,1)}=\mathbb{Z}^{2} M_{2}, \omega_{2}^{(\boldsymbol{q})}$ is the Dirac measure at the origin of the 2-dimensional torus $\mathbb{T}_{2}^{2}$. Furthermore, by (3.15) and (4.6),

$$
L_{1}^{(\boldsymbol{q})}=L_{1}^{\left(\pi\left(U_{1}^{(\boldsymbol{q})}\right)\right)}=\mathfrak{L}\left(\binom{0}{0},\binom{-\left(m_{1}+m_{2}\right) \sqrt{2}}{\left(1-m_{1}-m_{2}\right) \sqrt{2}}\right)=\mathbb{R}\binom{m_{1}+m_{2}}{m_{1}+m_{2}-1} \subset \mathbb{R}^{2}
$$

and so by (4.7), $\mathcal{O}_{1}^{(q)}=\mathcal{O}_{1}^{\left(\pi\left(U_{1}^{(q)}\right)\right)}=\pi\left(L_{1}^{(\boldsymbol{q})}\right)^{2}$, which is a 2-dimensional subtorus of $\mathbb{T}_{1}^{2}$, and $\omega_{1}^{(\boldsymbol{q})}$ is the normalized Haar measure of this subtorus. Now $\mu_{0, \boldsymbol{q}}=\left(J_{0,(2,1)}\right)_{*}\left(\overline{\omega^{(\boldsymbol{q})}}\right)=$ $\left(J_{0,(2,1)}\right)_{*}\left(\overline{\omega_{1}^{(\boldsymbol{q})}} \otimes \overline{\omega_{2}^{(\boldsymbol{q})}}\right)$. The left implication in (6.20) follows quite easily. Indeed, it is clear from the above that $\mu_{0, \boldsymbol{q}}$ only depends on $m_{1}+m_{2}$, and to prove that $\mu_{\boldsymbol{q}, 0}=\mu_{\boldsymbol{q}^{\prime}, 0}$ also holds in the case $m_{1}+m_{2}=1-\left(m_{1}^{\prime}+m_{2}^{\prime}\right)$, one makes use of a symmetry originating from the diffeomorphism $\binom{\boldsymbol{v}_{1}}{\boldsymbol{v}_{2}} \mapsto\binom{\boldsymbol{v}_{\boldsymbol{v}}}{\boldsymbol{v}_{1}}$ from $\pi\left(L_{1}^{(\boldsymbol{q})}\right)^{2}$ onto $\pi\left(L_{1}^{\left(\boldsymbol{q}^{\prime}\right)}\right)^{2}$. The proof of the converse is more involved; one approach goes via a study of the support of $\mu_{\boldsymbol{q}, 0}$ in $N\left(\mathbb{R}^{2}\right)$. For example, by a careful explicit analysis one verifies that for any $\boldsymbol{q}=\boldsymbol{m} M_{2}$ in $\mathcal{L}_{(2,1)}$, the set $\left\{k \in \mathbb{Z}_{\geq 3}\right.$ : $\left.\sum_{\boldsymbol{x} \in \mathbb{Z}^{2} \backslash\{0\}} \delta_{\boldsymbol{x}}+\sum_{\boldsymbol{x} \in \mathbb{Z}^{2}}\left(\delta_{\boldsymbol{x}}+\delta_{\boldsymbol{x}+k^{-1} \boldsymbol{e}_{1}}\right) \in \operatorname{supp}\left(\mu_{\boldsymbol{q}, 0}\right)\right\}$ contains exactly those $k \in \mathbb{Z}_{\geq 3}$ which divide $m_{1}+m_{2}$ or $m_{1}+m_{2}-1$; this gives the right implication in (6.20).

Remark 6.3. On the other hand, if $\mathcal{P}$ is periodic, then $\{\varsigma(\boldsymbol{q}): \boldsymbol{q} \in \mathcal{P}\}=\widetilde{\Psi}$, i.e. the space of marks $\Sigma^{\prime}$ in (6.17) is simply equal to $\widetilde{\Psi}$, even if we take $\mathcal{E}$ empty! Indeed, any periodic $\mathcal{P}$ may be represented as in (3.3) (3.5) with $N=1$ and all $c_{1, i}\left(i=1, \ldots, r_{1}\right)$ equal to a common value $c>0$. (This representation is admissible; note also that all the grids $\mathcal{L}_{(1,1)}, \ldots, \mathcal{L}_{\left(1, r_{1}\right)}$ are pairwise disjoint, by (3.6).) Then for any $i \in\left\{1, \ldots, r_{1}\right\}$, all points $\boldsymbol{q} \in \mathcal{L}_{(1, i)}$ give the same $\omega^{(\boldsymbol{q})}$. Indeed, for any two $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathcal{L}_{(1, i)}$ we have $U_{1}^{(\boldsymbol{q})}-U_{1}^{\left(\boldsymbol{q}^{\prime}\right)} \in \mathrm{M}_{1 \times d}(\mathbb{Z})$ by (3.15); thus $\pi\left(U_{1}^{(\boldsymbol{q})}\right)=\pi\left(U_{1}^{\left(\boldsymbol{q}^{\prime}\right)}\right)$ in $\mathbb{T}_{1}^{d}$, and so $\omega_{1}^{(\boldsymbol{q})}=\omega_{1}^{\left(\boldsymbol{q}^{\prime}\right)}$. It then follows from Proposition 5.13, or by direct inspection of the definitions, that in fact $\omega_{1}^{(\boldsymbol{q})}=\omega_{1}^{(1, i)}$ for all $\boldsymbol{q} \in \mathcal{L}_{(1, i)}$. Hence by (6.7) and (6.8), $\varsigma(\boldsymbol{q})=\sigma^{\psi(\boldsymbol{q})} \in \widetilde{\Psi}$ for all $\boldsymbol{q} \in \mathcal{P}$.

As we have mentioned, in the special case of periodic $\mathcal{P}$, Theorem 2.1] was proved in [21, Prop. 5.6], and it is easy to give the translation between the two formulations: Our marking data $\left[\Sigma^{\prime}, \varsigma, \varsigma \mapsto \mu_{\varsigma}, \mathrm{m}\right]$ with $\Sigma^{\prime}=\widetilde{\Psi}$ agrees exactly with the marking data in [21, Sec. 5.2, Prop. 5.6], with the only difference that the mark $\sigma^{(1, \ell)}$ is called simply " $\ell$ " in [21, Sec. 5.2], for $\ell=1, \ldots, r_{1}$ (also our $r_{1}$ is called " $m$ ", and our $c$ is " $\delta^{1 / d ") \text {. In particular, thus, for }}$ each $\ell \in\left\{1, \ldots, r_{1}\right\}$ the measure $\mu_{\sigma^{(1, \ell)}} \in P\left(N_{s}\left(\mathbb{R}^{d} \times \Sigma\right)\right)$ defined by (6.14) above, equals ${ }^{11}$ the measure " $\mu_{\ell}$ " in [21, (5.35)], although it requires some work to verify this fact from the definitions.

As noted in [21, Prop. 5.6], for $\mathcal{P}$ periodic, the condition [P2] holds with $\mathcal{E}=\emptyset$. Indeed, because of the periodicity, we have the even much stronger statement that the convergence $\mu_{\boldsymbol{q}, \rho}^{(\lambda)} \xrightarrow{\mathrm{w}} \mu_{\varsigma(\boldsymbol{q})}$ required in (2.7) holds uniformly over all $\boldsymbol{q} \in \mathcal{P}$.
Remark 6.4. Also in the case when $\mathcal{P}$ is a finite union of pairwise incommensurable grids, one can use $\Sigma^{\prime}=\widetilde{\Psi}$ as the space of marks. Indeed, in this case we have $r_{1}=\cdots=r_{N}=1$, and for any $j \neq j^{\prime}$ in $\{1, \ldots, N\}$ we have $L_{j}^{\left(j^{\prime}, 1\right)}=L_{j}=\mathbb{R}^{1}$ by Lemma 5.6 and (5.12); this also implies that $\omega_{j}^{\left(j^{\prime}, 1\right)}$ equals Haar measure on $\mathbb{T}_{j}^{d}$. Now by Lemma 5.4, applied with $\boldsymbol{a}=1 \in \mathbb{Z}^{1}$, for any $j \neq j^{\prime}$ the set $\left\{\boldsymbol{q} \in \mathcal{L}_{\left(j^{\prime}, 1\right)}: L_{j}^{(\boldsymbol{q})}=\{0\}\right\}$ has density zero. Hence we may assume that $\mathcal{E}$ contains the union of these sets for all pairs of $j \neq j^{\prime}$. Then, for any $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$, taking $j^{\prime}$ so that $\boldsymbol{q} \in \mathcal{L}_{\left(j^{\prime}, 1\right)}$ it follows that for each $j \neq j^{\prime}$ we have $L_{j}^{(\boldsymbol{q})} \neq\{0\}$, viz., $L_{j}^{(\boldsymbol{q})}=\mathbb{R}^{1}$; and therefore $\mathcal{O}_{j}^{(\boldsymbol{q})}=\mathbb{T}_{j}^{d}$ and $\omega_{j}^{(\boldsymbol{q})}=\omega_{j}^{\left(j^{\prime}, 1\right)}$ (this also implies that $\boldsymbol{q} \notin \mathcal{L}_{(j, 1)}$, and so $\psi(\boldsymbol{q})$ must equal $\left(j^{\prime}, 1\right)$ ). On the other hand, immediately from the definitions we have that both $\omega_{j^{\prime}}^{(\boldsymbol{q})}$ and $\omega_{j^{\prime}}^{\left(j^{\prime}, 1\right)}$ equals the Dirac measure at the origin of $\mathbb{T}_{j^{\prime}}^{d}$. Hence $\omega^{(\boldsymbol{q})}=\omega^{\left(j^{\prime}, 1\right)}=\omega^{\psi(\boldsymbol{q})}$ and $\varsigma(\boldsymbol{q})=\sigma^{\psi(\boldsymbol{q})} \in \widetilde{\Psi}$. Since this holds for all $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$, we get $\Sigma^{\prime}=\widetilde{\Psi}$ in (6.17).

## 7. VErification of [Q1],[Q2],[Q3],[P1],[P3], AND initial discussion REGARDing [P2]

In order to prove the main result of the paper, Theorem 2.1, we now wish to prove that all the conditions [P1]-[P3] and [Q1]-[Q3] are satisfied for the maps $\varsigma: \mathcal{P} \rightarrow \Sigma$ and $\varsigma \mapsto \mu_{\varsigma}$ and measure $m$ which we have defined in the previous section. In the present section we will prove all of these conditions except [P2]; we will also reduce the verification of [P2] to a certain statement about uniform equidistribution in the homogeneous space $\mathbb{X}$, Theorem 7.7 below.
7.1. Verification of [Q1], [Q2], [Q3]. We start with the conditions [Q1]-[Q3].

The following lemma shows that [Q1] holds, in a much stronger form.
Lemma 7.1. For every $\varsigma \in \Sigma, \mu_{\varsigma}$ is $\mathrm{SL}_{d}(\mathbb{R})$-invariant.
Proof. It is immediate from the definition in (4.9) that for any $\omega_{j} \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$, the measure $\bar{\omega}_{j}$ on $\mathbb{X}_{j}$ is right $\mathrm{SL}_{d}(\mathbb{R})$ invariant (this was also used in the proof of Proposition 4.7); hence for any $\omega \in \Omega$, the measure $\bar{\omega}$ on $\mathbb{X}$ is right $\mathrm{SL}_{d}(\mathbb{R})^{N}$-invariant, and in particular it is right

[^7]$\varphi\left(\mathrm{SL}_{d}(\mathbb{R})\right)$-invariant, where $\varphi: \mathrm{SL}_{d}(\mathbb{R}) \rightarrow G$ is the diagonal embedding. Next we note that, for any $\psi \in \Psi$,
$$
J_{\psi}(\Gamma g \varphi(h))=J_{\psi}(\Gamma g) h \quad \text { for all } \Gamma g \in \mathbb{X}^{\psi}, h \in \mathrm{SL}_{d}(\mathbb{R}) .
$$

It follows that the measure $J_{\psi *} \bar{\omega}$ on $N_{s}(\mathcal{X})$ is $\mathrm{SL}_{d}(\mathbb{R})$-invariant, for any $\omega \in \Omega$. This implies the lemma, via the definition (6.14).

Next, because of the following general fact, also [Q2] is an immediate consequence of Lemma 7.1:

Lemma 7.2. If $\mu \in P\left(N_{s}(\mathcal{X})\right)$ is invariant under the action of $\mathrm{SO}(d)$ then

$$
\begin{equation*}
\mu\left(\left\{\nu \in N(\mathcal{X}): \exists x_{1} \in \mathbb{R} \text { s.t. } \nu\left(\left\{x_{1}\right\} \times \mathbb{R}^{d-1} \times \Sigma\right)>1\right\}\right)=0 . \tag{7.1}
\end{equation*}
$$

Proof. For any $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$, set

$$
\begin{equation*}
A_{\boldsymbol{v}}=\left\{\nu \in N(\mathcal{X}): \exists x_{1} \in \mathbb{R} \text { s.t. } \nu\left(\left(x_{1} \boldsymbol{v}+\boldsymbol{v}^{\perp}\right) \times \Sigma\right)>1\right\} . \tag{7.2}
\end{equation*}
$$

Then our task is to prove $\mu\left(A_{e_{1}}\right)=0$. The fact that $\mu$ is $\mathrm{SO}(d)$-invariant implies that $\mu\left(A_{v}\right)=\mu\left(A_{e_{1}}\right)$ for all $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$. Hence we have, with $\lambda_{1}$ being the uniform probability measure on $\mathrm{S}_{1}^{d-1}$ :

$$
\mu\left(A_{e_{1}}\right)=\int_{\mathrm{S}_{1}^{d-1}} \mu\left(A_{\boldsymbol{v}}\right) d \lambda_{1}(\boldsymbol{v})=\int_{N_{\boldsymbol{s}}(\mathcal{X})} \int_{\mathrm{S}_{1}^{d-1}} I\left(\nu \in A_{\boldsymbol{v}}\right) d \lambda_{1}(\boldsymbol{v}) d \mu(\nu)=0 .
$$

Here the second equality holds by Fubini's Theorem and since $\mu\left(N(\mathcal{X}) \backslash N_{s}(\mathcal{X})\right)=0$ (because of $\mu \in P\left(N_{s}(\mathcal{X})\right)$ ), and the last equality holds since for any $\nu \in N_{s}(\mathcal{X})$ we have $\int_{\mathrm{S}_{1}^{d-1}} I(\nu \in$ $\left.A_{\boldsymbol{v}}\right) d \lambda_{1}(\boldsymbol{v})=0$, since the set $\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \nu \in A_{\boldsymbol{v}}\right\}$ is a countable union of subspheres of $\mathrm{S}_{1}^{d-1}$ of codimension one.

Next we turn to the condition [Q3].
Lemma 7.3. [Q3] holds.
Proof. (Cf. the proof of [21, Lemma 5.3.13].) Set $\Lambda=\mathrm{SL}_{d}(\mathbb{Z})$ and $\mathbb{Y}=\Lambda \backslash \mathrm{SL}_{d}(\mathbb{R})$. Let us also fix a choice of $\psi=(j, i) \in \Psi$. For $R>0$ we set

$$
\mathbb{Y}(R)=\left\{\Lambda h \in \mathbb{Y}: c_{\psi}\left(\mathbb{Z}^{d} h\right)+\mathcal{B}_{R / 2}^{d}=\mathbb{R}^{d}\right\}
$$

Note that the set $\mathbb{Y}(R)$ is increasing with respect to $R$. Also, for every $\Lambda h \in \mathbb{Y}$, the set $c_{\psi}\left(\mathbb{Z}^{d} h\right)$ is a lattice in $\mathbb{R}^{d}$ and hence there exists some $R=R(h)>0$ such that $c_{\psi}\left(\mathbb{Z}^{d} h\right)+\mathcal{B}_{R / 2}^{d}=\mathbb{R}^{d}$. Hence $\cup_{R>0} \mathbb{Y}(R)=\mathbb{Y}$, and it follows that for any given $\varepsilon>0$ we can choose $R>0$ so that $\eta(\mathbb{Y}(R))>1-\varepsilon$, where $\eta$ is the $\mathrm{SL}_{d}(\mathbb{R})$-invariant probability measure on $\mathbb{Y}$.

With this choice of $R$, we now claim that for any $\varsigma=\left(\psi^{\prime}, \omega\right) \in \Sigma$ and $\boldsymbol{x} \in \mathbb{R}^{d}$, (2.8) holds. By (6.14), this is equivalent to the following:

$$
\begin{equation*}
\bar{\omega}\left(\left\{\Gamma g \in \mathbb{X}^{\psi^{\prime}}: J_{\psi^{\prime}}(\Gamma g) \cap\left(\mathcal{B}_{R}^{d}(\boldsymbol{x}) \times \Sigma\right)=\emptyset\right\}\right)<\varepsilon \tag{7.3}
\end{equation*}
$$

Choose $\boldsymbol{y} \in \mathbb{R}^{d}$ so that $\mathcal{B}_{R / 2}^{d}(\boldsymbol{y}) \subset \mathcal{B}_{R}^{d}(\boldsymbol{x})$ and $\mathbf{0} \notin \mathcal{B}_{R / 2}^{d}(\boldsymbol{y})$. For any $g \in G$, letting $h=$ $\iota\left(p_{j}(g)\right) \in \mathrm{SL}_{d}(\mathbb{R})$ we have that the grid $c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right)$ is a translate of the lattice $c_{\psi}\left(\mathbb{Z}^{d} h\right)$; and in particular if $\Lambda h \in \mathbb{Y}(R)$ then $c_{\psi}\left(\mathbb{Z}^{d} p_{\psi}(g)\right)$ must contain a point in $\mathcal{B}_{R / 2}^{d}(\boldsymbol{y})$. This implies that for every $\Gamma g \in \mathbb{X}$ satisfying $\widetilde{\iota}\left(\widetilde{p}_{j}(\Gamma g)\right) \in \mathbb{Y}(R)$, the point set $J(\Gamma g)$ (cf. (6.12)) must contain a point in $\mathcal{B}_{R / 2}^{d}(\boldsymbol{y}) \times \Sigma$. Using also $\mathbf{0} \notin \mathcal{B}_{R / 2}^{d}(\boldsymbol{y})$ and $\mathcal{B}_{R / 2}^{d}(\boldsymbol{y}) \subset \mathcal{B}_{R}^{d}(\boldsymbol{x})$, it follows that the measure in the left hand side of (7.3) is bounded above by

$$
\begin{equation*}
\bar{\omega}\left(\left\{\Gamma g \in \mathbb{X}^{\psi^{\prime}}: \widetilde{\iota}\left(\widetilde{p}_{j}(\Gamma g)\right) \notin \mathbb{Y}(R)\right\}\right) \tag{7.4}
\end{equation*}
$$

However, writing $\omega=\left(\omega_{1}, \ldots, \omega_{j}\right)$ we have $\widetilde{p}_{j *}(\bar{\omega})=\bar{\omega}_{j}$, which is an $\mathrm{SL}_{d}(\mathbb{R})$-invariant probability measure on $\mathbb{X}_{j}$ (this is immediate from the definition (4.8), as we have noted previously). Hence the pushforward of $\bar{\omega}$ by $\tau \circ \widetilde{p}_{j}$ equals $\eta$, and so the measure in (7.4) equals $\eta(\mathbb{Y} \backslash \mathbb{Y}(R)$ ), which by our choice of $R$ is less than $\varepsilon$. Hence (7.3), and thereby the lemma, is proved.

We end this section by giving a closely related fact, namely a formula for the intensity measure of a point process with distribution $\mu_{\varsigma}$, for any $\varsigma \in \Sigma$. This formula will not be used until Section 10 .

Lemma 7.4. Let $\varsigma=\left(\psi^{\prime}, \omega\right) \in \Sigma$. Then for any Borel set $B \subset \mathcal{X}$ we have

$$
\begin{equation*}
\int_{N_{s}(\mathcal{X})} \#(B \cap Y) d \mu_{\varsigma}(Y)=\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B)+\sum_{\psi \in \Psi \backslash\left\{\psi^{\prime}\right\}} \omega_{j_{\psi}}\left(\tilde{\mathrm{r}}_{i_{\psi}}^{-1}(\{\mathbf{0}\})\right) \cdot I\left(\left(\mathbf{0}, \sigma^{\psi}\right) \in B\right) . \tag{7.5}
\end{equation*}
$$

In particular, for every $\psi^{\prime} \in \Psi$, a point process with distribution $\mu_{\sigma^{\prime}}$ has intensity measure $\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}$.

Proof. By (6.14), (6.13) and (6.12), the left hand side of (7.5) equals

$$
-I\left(\left(\mathbf{0}, \sigma^{\psi^{\prime}}\right) \in B\right)+\int_{\mathbb{X} \psi^{\prime}} \sum_{\psi \in \Psi} \sum_{\boldsymbol{m} \in \mathbb{Z}^{d}} I\left(\left(c_{\psi} \cdot\left(\boldsymbol{m} p_{\psi}(g)\right), \sigma^{\psi}\right) \in B\right) d \bar{\omega}(\Gamma g)
$$

By (6.15), the integration may just as well be taken over all $\mathbb{X}$. Moving out the sum over $\psi$ and then using (6.11) and (3.13), we get

$$
-I\left(\left(\mathbf{0}, \sigma^{\psi^{\prime}}\right) \in B\right)+\sum_{\psi \in \Psi} \int_{\mathbb{X}_{j_{\psi}}} \sum_{\boldsymbol{m} \in \mathbb{Z}^{d}} I\left(\left(c_{\psi} \cdot\left(\boldsymbol{m a}_{i_{\psi}}(g)\right), \sigma^{(\psi)}\right) \in B\right) d \bar{\omega}_{j_{\psi}}\left(\Gamma_{j_{\psi}} g\right)
$$

and by Proposition 4.8, this equals

$$
-I\left(\left(\mathbf{0}, \sigma^{\psi^{\prime}}\right) \in B\right)+\sum_{\psi \in \Psi} \bar{n}_{\psi} \int_{\mathbb{R}^{d}} I\left(\left(\boldsymbol{x}, \sigma^{\psi}\right) \in B\right) d \boldsymbol{x}+\sum_{\psi \in \Psi} \omega_{j_{\psi}}\left(\widetilde{\mathfrak{r}}_{i_{\psi}}^{-1}(\{\mathbf{0}\})\right) \cdot I\left(\left(\mathbf{0}, \sigma^{\psi}\right) \in B\right)
$$

Using $\omega_{j_{\psi^{\prime}}}\left(\tilde{\mathrm{r}}_{i_{\psi^{\prime}}}^{-1}(\{\mathbf{0}\})\right)=1$, which holds since $\left(\psi^{\prime}, \omega\right) \in \Sigma$, together with $\mu_{\mathcal{X}}=\operatorname{vol} \times \mathrm{m}$ and (6.9), we obtain the right hand side of (7.5).

To obtain the last statement of the lemma, we apply (7.5) for $\varsigma=\sigma^{\psi^{\prime}}=\left(\psi^{\prime}, \omega^{\psi^{\prime}}\right)$; in this case the sum over $\Psi \backslash\left\{\psi^{\prime}\right\}$ in (7.5) vanishes, by Lemma 6.6; hence the right hand side of (7.5) equals $\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B)$.

### 7.2. Verification of [P1] (uniform density).

Proposition 7.5. [P1] holds, i.e. for any bounded subset $B \subset \mathcal{X}$ with $\mu_{\mathcal{X}}(\partial B)=0$, we have

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{\#(\widetilde{\mathcal{P}} \cap T B)}{T^{d}}=\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B) \tag{7.6}
\end{equation*}
$$

Proof. Note that $\mathcal{X}$ decomposes as the disjoint union $\sqcup_{\psi \in \Psi} \mathcal{X}_{\psi}$, where $\mathcal{X}_{\psi}=\mathbb{R}^{d} \times \Sigma_{\psi}$ with $\Sigma_{\psi}:=(\{\psi\} \times \Omega) \cap \Sigma$; note also that each set $\mathcal{X}_{\psi}$ is both open and closed is $\mathcal{X}$. It follows that it suffices to prove (7.6) under the extra assumption that $B \subset \mathcal{X}_{\psi}$ for some fixed $\psi$. Using also the fact that the set $\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \psi(\boldsymbol{q}) \neq \psi\right\}$ has density zero (cf. Remark (3.2), it follows that our task is to prove the following, for any bounded set $B \subset \mathcal{X}_{\psi}$ with $\mu_{\mathcal{X}}(\partial B)=0$ :

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}:\left(\boldsymbol{q},\left(\psi, \omega^{(\boldsymbol{q})}\right)\right) \in T B\right\}}{T^{d}}=\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}\left(\mathcal{X}_{\psi}\right) \tag{7.7}
\end{equation*}
$$

Let us first verify that (7.7) holds for any set $B$ of the form

$$
\begin{equation*}
B=\left(\prod_{i=1}^{d}\left[\alpha_{i}, \beta_{i}\right)\right) \times U \tag{7.8}
\end{equation*}
$$

for any real numbers $\alpha_{i}<\beta_{i}(i=1, \ldots, d)$ and any open neighbourhood $U$ of $\sigma^{\psi}$ in $\Sigma_{\psi}$. Indeed, given such a $U$, there exist open neighbourhoods $U_{j}$ of $\omega_{j}^{\psi}$ in $P\left(\mathbb{T}_{j}^{d}\right)$ for $j=1, \ldots, N$ such that $\Sigma_{\psi} \cap\left(\{\psi\} \times \prod_{j=1}^{N} U_{j}\right) \subset U$. Applying now Proposition 5.13 to the set $U_{j}$, for each $j=1, \ldots, N$, it follows that there exists a subset $\mathcal{Z}^{\prime} \subset \mathcal{L}_{\psi}$ of density zero such that $\left(\psi, \omega^{(\boldsymbol{q})}\right) \in U$ for all
$\boldsymbol{q} \in \mathcal{L}_{\psi} \backslash \mathcal{Z}^{\prime}$. Furthermore, for $B$ as in (7.8), we have $\bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B)=\bar{n}_{\psi} \prod_{i=1}^{d}\left(\beta_{i}-\alpha_{i}\right)$, and so (7.7) follows from the fact that the grid $\mathcal{L}_{\psi}$ has asymptotic density $\bar{n}_{\psi}$ in $\mathbb{R}^{d}$.

Next, if $B$ is any bounded open subset of $\mathcal{X}_{\psi}$, then there exists a sequence $B_{1} \subset B_{2} \subset \cdots$ of subsets of $B$ such that each $B_{k}$ is a finite disjoint union of sets of the form (7.8), and $B \cap\left(\mathbb{R}^{d} \times\left\{\sigma^{\psi}\right\}\right) \subset \cup_{k=1}^{\infty} B_{k}$. Using $\mu_{\mathcal{X}}(B)=\mu_{\mathcal{X}}\left(B \cap\left(\mathbb{R}^{d} \times\left\{\sigma^{\psi}\right\}\right)\right)$, it follows that $\mu_{\mathcal{X}}\left(B_{k}\right) \rightarrow$ $\mu_{\mathcal{X}}(B)$ as $k \rightarrow \infty$, and hence since (7.7) holds for each of our sets $B_{k}$, it follows that

$$
\begin{equation*}
\liminf _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}:\left(\boldsymbol{q},\left(\psi, \omega^{(\boldsymbol{q})}\right)\right) \in T B\right\}}{T^{d}} \geq \bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B) \tag{7.9}
\end{equation*}
$$

Next if $\widetilde{B}$ is any bounded closed subset of $\mathcal{X}_{\psi}$, then by taking $R>0$ so that $\widetilde{B} \subset B^{\prime}:=\mathcal{B}_{R}^{d} \times \Sigma_{\psi}$, and noticing that (7.7) holds for $B^{\prime}$ (since the grid $\mathcal{L}_{\psi}$ has asymptotic density $\bar{n}_{\psi}$ in $\mathbb{R}^{d}$ ), and also (7.9) holds for the bounded open set $B^{\prime} \backslash \widetilde{B}$, it follows that

$$
\begin{equation*}
\limsup _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}:\left(\boldsymbol{q},\left(\psi, \omega^{(\boldsymbol{q})}\right)\right) \in T \widetilde{B}\right\}}{T^{d}} \leq \bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(\widetilde{B}) \tag{7.10}
\end{equation*}
$$

Finally consider an arbitrary bounded subset $B \subset \mathcal{X} \psi$ with $\mu_{\mathcal{X}}(\partial B)=0$. Let $B^{\circ}$ and $\bar{B}$ be the interior and the closure of $B$, respectively. Then (7.9) holds for $B^{\circ}$ and (7.10) holds for $\bar{B}$, and furthermore $\mu_{\mathcal{X}}\left(B^{\circ}\right)=\mu_{\mathcal{X}}(\bar{B})$, since $\mu_{\mathcal{X}}(\partial B)=0$. Hence (7.7) holds.
7.3. Initial discussion regarding [P2] (uniform spherical equidistribution). We have the following result:

Theorem 7.6. [P2] holds, i.e. there exists a subset $\mathcal{E} \subset \mathcal{P}$ of density zero such that for any fixed $T \geq 1$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have $\mu_{\boldsymbol{q}, \rho}^{(\lambda)} \xrightarrow{\mathrm{w}} \mu_{\varsigma(\boldsymbol{q})}$ as $\rho \rightarrow 0$, uniformly for $\boldsymbol{q} \in$ $\mathcal{P} \cap \mathcal{B}_{T \rho^{1-d}}^{d} \backslash \mathcal{E}$.

In this section we will prove that Theorem[7.6 follows from the following theorem on uniform equidistribution in the homogeneous space $\mathbb{X}$, the proof of which is the main goal of the later sections in this paper.

Recall that $\varphi: \mathrm{SL}_{d}(\mathbb{R}) \rightarrow G$ is the diagonal embedding.
Theorem 7.7. Given any $\psi \in \Psi$ and any decreasing function $\mathcal{T}:(0,1) \rightarrow \mathbb{R}^{+}$, there exists a subset $\mathcal{E} \subset \mathcal{L}_{\psi}$ of density zero such that for any fixed $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\int_{S_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X} \psi} f \overline{d \omega^{(\boldsymbol{q})}} \rightarrow 0 \tag{7.11}
\end{equation*}
$$

as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$.
To see that the statement of Theorem 7.7 makes sense, note that for every $\boldsymbol{q} \in \mathcal{L}_{\psi}$ we have $\Gamma g_{0}^{(\boldsymbol{q})} \in \mathbb{X}^{\psi}$ by Lemma 6.1, and so $\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \in \mathbb{X}^{\psi}$ for all $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$ and $\rho>$ 0 ; furthermore, $\left.\overline{\omega^{(\boldsymbol{q})}} \mathbb{X}^{\psi}\right)=1$ by Lemmas 6.4 and 6.2 , meaning that $\overline{\omega^{(\boldsymbol{q})}} \in P(\mathbb{X})$ gives a probability measure on the subset $\mathbb{X}^{\psi}$. It is also worth noticing that without the uniformity over $\boldsymbol{q}$, the statement of Theorem 7.7 would be an immediate consequence of Theorem 4.2.

We will now give the proof of Theorem 7.6, assuming Theorem 7.7. As the very first step, let us apply Theorem 7.7 with $\mathcal{T}(\rho)=\rho^{-d}$ and for each $\psi \in \Psi$; this gives the existence of subsets $\mathcal{E}_{\psi} \subset \mathcal{L}_{\psi}$ of density zero such that for any fixed $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$,

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X}^{\psi}} f \overline{d \omega^{(\boldsymbol{q})}} \rightarrow 0 \tag{7.12}
\end{equation*}
$$

as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\rho^{-d}}^{d} \backslash \mathcal{E}_{\psi}$. Let us now set

$$
\begin{equation*}
\mathcal{E}:=\bigcup_{\psi \in \Psi} \mathcal{E}_{\psi} \cup \bigcup_{\psi \neq \psi^{\prime} \in \Psi}\left(\mathcal{L}_{\psi} \cap \mathcal{L}_{\psi^{\prime}}\right) \tag{7.13}
\end{equation*}
$$

By Lemma 3.2, this set $\mathcal{E}$ is a subset of $\mathcal{P}$ of density zero. We will keep this set $\mathcal{E}$ fixed in the rest of the section, and we will prove that the statement of Theorem 7.6 holds with this set $\mathcal{E}$. For any $T \geq 1$ and $\rho>0$ we write, as in (2.7),

$$
\mathcal{P}_{T}(\rho):=\mathcal{P} \cap \mathcal{B}_{T \rho^{1-d}}^{d} \backslash \mathcal{E}
$$

Note that for any given $T$ we have $T \rho^{1-d} \leq \rho^{-d}$ for all sufficiently small $\rho$. Recall also that $\mathcal{P}=\cup_{\psi \in \Psi} \mathcal{L}_{\psi}$. Hence the statement around (7.12) now implies that the following holds:

$$
\left\{\begin{array}{l}
\text { for any fixed } T \geq 1, f \in \mathrm{C}_{b}(\mathbb{X}) \text { and } \lambda \in P_{\mathrm{ac}}\left(\mathrm{~S}_{1}^{d-1}\right)  \tag{7.14}\\
\qquad \int_{\mathrm{S}_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X}} f d \overline{\omega^{(\boldsymbol{q})}} \rightarrow 0 \\
\text { as } \rho \rightarrow 0, \text { uniformly over all } \boldsymbol{q} \in \mathcal{P}_{T}(\rho)
\end{array}\right.
$$

Lemma 7.8. Let $k \in \mathbb{Z}_{>0}$ and let $B$ be a bounded subset of $\mathbb{R}^{d}$ with $\operatorname{vol}(\partial B)=0$. Then for any $V>\operatorname{vol}(B), T \geq 1$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, there exists $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \#\left((\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap B \backslash\{\mathbf{0}\}\right) \geq k\right\}\right)<\bar{n}_{\mathcal{P}} V / k \tag{7.15}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$.
Proof. The assumptions imply that $B$ is Jordan measurable, and hence there is a function $f \in \mathrm{C}_{c}\left(\mathbb{R}^{d}\right)$ such that $f=1$ on $B, 0 \leq f \leq 1$ everywhere, and $V_{f}:=\int_{\mathbb{R}^{d}} f d$ vol $<V$. For each $\psi=(j, i) \in \Psi$, let $\widehat{f}_{\psi} \in \mathrm{C}\left(\mathbb{X}_{j}\right)$ be the " $\psi$ th Siegel transform" of $f$, as defined in (4.14). The function $\widehat{f}_{\psi}$ is typically unbounded; therefore we set $\widetilde{f}_{\psi}=\min \left(k+1, \widehat{f}_{\psi}\right)$; this is a nonnegative function in $\mathrm{C}_{b}\left(\mathbb{X}_{j}\right)$, and hence $\tilde{f}_{\psi} \circ p_{j} \in \mathrm{C}_{b}(\mathbb{X})$. Hence by (7.14), and since $p_{j}\left(g_{0}^{(\boldsymbol{q})}\right)=\mathrm{I}_{U_{j}^{(\boldsymbol{q})}} M_{j}$,

$$
\int_{S_{1}^{d-1}} \widetilde{f}_{\psi}\left(\Gamma_{j} \mathrm{I}_{U_{j}^{(\boldsymbol{q})}} M_{j} R(\boldsymbol{v}) D_{\rho}\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X}} \widetilde{f}_{\psi} \circ p_{j} \overline{d \omega^{(\boldsymbol{q})}} \rightarrow 0
$$

as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$. Here $p_{j *}\left(\overline{\omega^{(\boldsymbol{q})}}\right)=\overline{\omega_{j}^{(\boldsymbol{q})}}$ (cf. (6.11) $)$; hence by Proposition 4.8 and Lemma 6.4.

$$
\int_{\mathbb{X}} \widetilde{f}_{\psi} \circ p_{j} d \overline{\omega^{(\boldsymbol{q})}} \leq \int_{\mathbb{X}} \widehat{f}_{\psi} \circ p_{j} d \overline{\omega^{(\boldsymbol{q})}}=\bar{n}_{\psi} V_{f}+\delta_{\boldsymbol{q} \in \mathcal{L}_{\psi}} \cdot f(\mathbf{0})
$$

Adding the above over all $\psi$ and using $\bar{n}_{\mathcal{P}}:=\sum_{\psi \in \Psi} \bar{n}_{\psi}$, it follows that there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} \sum_{\psi \in \Psi}\left(\widetilde{f}_{\psi}\left(\Gamma_{j} \mathrm{I}_{U_{j}^{(\boldsymbol{q})}} M_{j} R(\boldsymbol{v}) D_{\rho}\right)-\delta_{\boldsymbol{q} \in \mathcal{L}_{\psi}} \cdot f(\mathbf{0})\right) d \lambda(\boldsymbol{v})<\bar{n}_{\mathcal{P}} V \tag{7.16}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$. Here, by (4.14) and (3.20), we have for every $\boldsymbol{v}$ and $\psi$ :

$$
\begin{aligned}
\widetilde{f}_{\psi}\left(\Gamma_{j} \mathrm{I}_{U_{j}^{(\boldsymbol{q})}} M_{j} R(\boldsymbol{v}) D_{\rho}\right)-\delta_{\boldsymbol{q} \in \mathcal{L}_{\psi}} \cdot f(\mathbf{0})=\min (k+1, & \left.\sum_{\boldsymbol{p} \in\left(\mathcal{L}_{\psi}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho}} f(\boldsymbol{p})\right)-\delta_{\boldsymbol{q} \in \mathcal{L}_{\psi}} \cdot f(\mathbf{0}) \\
& \geq \min \left(k, \sum_{\boldsymbol{p} \in\left(\mathcal{L}_{\psi}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho} \backslash\{\mathbf{0}\}} f(\boldsymbol{p})\right)
\end{aligned}
$$

Recalling also that $\cup_{\psi \in \Psi} \mathcal{L}_{\psi}=\mathcal{P}$, it follows that

$$
\int_{\mathrm{S}_{1}^{d-1}} \min \left(k, \sum_{\boldsymbol{p} \in(\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \backslash\{\mathbf{0}\}} f(\boldsymbol{p})\right) d \lambda(\boldsymbol{v})<\bar{n}_{\mathcal{P}} V
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$. Here, for every $\boldsymbol{v}$ such that $\#\left((\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap B \backslash\{\mathbf{0}\}\right) \geq k$, the integrand equals $k$. Hence we obtain the statement of the lemma.

Lemma 7.9. Let $\mathcal{Z}$ be any subset of $\mathcal{P}$ of density zero, and let $T \geq 1, \lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ and $S>0$. Then

$$
\begin{equation*}
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \mathcal{Z} \cap\left(\boldsymbol{q}+\mathcal{B}_{S}^{d} D_{\rho}^{-1} R(\boldsymbol{v})^{-1}\right) \backslash\{\boldsymbol{q}\} \neq \emptyset\right\}\right) \rightarrow 0 \tag{7.17}
\end{equation*}
$$

as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$.
Proof. (This is similar to the proof of [21, Lemma 2.17].) Let $\lambda_{1}$ be the normalized Lebesgue measure on $\mathrm{S}_{1}^{d-1}$. By a standard approximation argument, using the fact that $C_{c}\left(\mathrm{~S}_{1}^{d-1}\right)$ is dense in $L^{1}\left(\mathrm{~S}_{1}^{d-1}\right)$, it suffices to prove (7.17) for those $\lambda$ which have a continuous density with respect to $\lambda_{1}$; and thus in fact it suffices to prove (7.17) for the single case $\lambda=\lambda_{1}$.

Let $T \geq 1, S>0$ and $\varepsilon>0$ be given. Take $0<r<S$ so small that $\bar{n}_{\mathcal{P}} \operatorname{vol}\left(\mathcal{B}_{r}^{d}\right)<\varepsilon$. Set $k=2 S / r>2$ and $T^{\prime}=k^{d-1} T$. By Lemma 7.8, there exists $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\lambda_{1}\left(\left\{\boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}:(\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap \mathcal{B}_{r}^{d} \backslash\{\mathbf{0}\} \neq \emptyset\right\}\right)<\varepsilon \tag{7.18}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T^{\prime}}(\rho)$. Set $\widetilde{B}:=\mathcal{B}_{r}^{d} D_{k}^{-1}$. Replacing $\rho$ by $k \rho$ in (7.18), it follows that for all $\rho \in\left(0, \rho_{0} / k\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T^{\prime}}(k \rho)=\mathcal{P}_{T}(\rho)$ we have

$$
\begin{equation*}
\lambda_{1}\left(\left\{\boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}:(\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap \widetilde{B} \backslash\{\mathbf{0}\} \neq \emptyset\right\}\right)<\varepsilon \tag{7.19}
\end{equation*}
$$

One verifies that $\left|x_{1}\right| \geq k_{1}:=(r / 2)^{d} S^{1-d}$ for all $\boldsymbol{x} \in \mathcal{B}_{S}^{d} \backslash \widetilde{B}$, and hence

$$
\begin{equation*}
\left(\mathcal{B}_{S}^{d} \backslash \widetilde{B}\right) D_{\rho}^{-1} \subset A(\rho):=\mathcal{B}_{S \rho^{1-d}}^{d} \backslash \mathcal{B}_{k_{1} \rho^{1-d}}^{d}, \quad \forall \rho>0 \tag{7.20}
\end{equation*}
$$

Now for any $\rho \in\left(0, \rho_{0} / k\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$ we have, using (7.19) and (7.20):

$$
\begin{align*}
\lambda_{1}\left(\left\{\boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}:(\mathcal{Z}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap \mathcal{B}_{S}^{d} \backslash\{\mathbf{0}\} \neq \emptyset\right\}\right)  \tag{7.21}\\
<\varepsilon+\sum_{\boldsymbol{p} \in \mathcal{Z} \cap(\boldsymbol{q}+A(\rho))} \lambda_{1}\left(\left\{\boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}:(\boldsymbol{p}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \in \mathcal{B}_{S}^{d}\right\}\right)
\end{align*}
$$

But if $(\boldsymbol{p}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \in \mathcal{B}_{S}^{d}$, or equivalently $\boldsymbol{p} \in \boldsymbol{q}+\mathcal{B}_{S}^{d} D_{\rho}^{-1} R(\boldsymbol{v})^{-1}$, then $\boldsymbol{p}$ has a distance less than $S \rho$ to the line $\boldsymbol{q}+\mathbb{R} \boldsymbol{v}$; and if also $\boldsymbol{p} \in \boldsymbol{q}+A(\rho)$ then the angle $\varphi(\boldsymbol{v}, \boldsymbol{p}-\boldsymbol{q})$ between the vectors $\boldsymbol{v}$ and $\boldsymbol{p}-\boldsymbol{q}$ satisfies $\sin \varphi(\boldsymbol{v}, \boldsymbol{p}-\boldsymbol{q})<\left(S / k_{1}\right) \rho^{d}$. The measure of the set of such points $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$ with respect to $\lambda_{1}$ is bounded above by $C_{1} \rho^{d(d-1)}$, where $C_{1}$ depends on $d, S, r$ but not on $\rho$ or $\boldsymbol{p}$. Hence (7.21) is

$$
\leq \varepsilon+\#\left(\mathcal{Z} \cap\left(\boldsymbol{q}+\mathcal{B}_{S \rho^{1-d}}^{d}\right)\right) \cdot C_{1} \rho^{d(d-1)} \leq \varepsilon+\#\left(\mathcal{Z} \cap \mathcal{B}_{(T+S) \rho^{1-d}}^{d}\right) \cdot C_{1} \rho^{d(d-1)}
$$

and since $\mathcal{Z}$ has density zero, the last term is less than $\varepsilon$ for $\rho$ sufficiently small.
Recall that $\mu_{\boldsymbol{q}, \rho}^{(\lambda)}$ is the distribution of $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ for $\boldsymbol{v}$ random in $\left(\mathrm{S}_{1}^{d-1}, \lambda\right)$. We now introduce a certain approximation $\mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})$ to $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$, which will be easier to handle. We set

$$
\begin{equation*}
\widetilde{\mathcal{P}}^{\prime}=\bigcup_{\psi \in \Psi}\left\{\left(\boldsymbol{p}, \sigma^{\psi}\right): \boldsymbol{p} \in \mathcal{L}_{\psi}\right\} \tag{7.22}
\end{equation*}
$$

Note that, unlike the projection $\widetilde{\mathcal{P}} \rightarrow \mathcal{P}$, the projection $\widetilde{\mathcal{P}}^{\prime} \rightarrow \mathcal{P}$ is not necessarily injective! (However, by Remark 3.2, it becomes injective after removing a set of density zero from $\widetilde{\mathcal{P}}^{\prime}$.) For any $\boldsymbol{q} \in \mathcal{P}$, we set

$$
\widetilde{\mathcal{P}}_{\boldsymbol{q}}^{\prime}= \begin{cases}\widetilde{\mathcal{P}}^{\prime} \backslash\left\{\left(\boldsymbol{q}, \sigma^{\psi(\boldsymbol{q})}\right)\right\} & (\boldsymbol{q} \in \mathcal{P})  \tag{7.23}\\ \widetilde{\mathcal{P}}^{\prime} & (\boldsymbol{q} \notin \mathcal{P})\end{cases}
$$

and

$$
\begin{equation*}
\mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})=\left(\widetilde{\mathcal{P}}_{\boldsymbol{q}}^{\prime}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho} \tag{7.24}
\end{equation*}
$$

Lemma 7.10. For every $\boldsymbol{q} \in \mathcal{P}$ we have $\mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})=J_{\psi(\boldsymbol{q})}\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right)$.

Proof. By parsing the definitions ( (7.23), (7.24) and (6.13), we see that it suffices to prove

$$
\begin{equation*}
\widetilde{\mathcal{P}}^{\prime}-\boldsymbol{q}=J\left(\Gamma g_{0}^{(\boldsymbol{q})}\right) \tag{7.25}
\end{equation*}
$$

However, comparing (6.12) and (7.22), we see that (7.25) is an immediate consequence of (3.20) (just as is (3.21)).

The following lemma shows that $\mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})$ approximates $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ in a sense that is appropriate for us.
Lemma 7.11. Let $f \in \mathrm{C}_{c}(\mathcal{X}), T \geq 1, \lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ and $\varepsilon>0$. Then

$$
\begin{equation*}
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}:\left|\sum_{x \in \mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})} f(x)-\sum_{x \in \mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})} f(x)\right|>\varepsilon\right\}\right) \rightarrow 0 \tag{7.26}
\end{equation*}
$$

as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$.
Proof. Choose $S>0$ so that $\operatorname{supp}(f) \subset \mathcal{B}_{S}^{d} \times \Sigma$. Note that for every $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$ and every $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$ we have, using the fact that $\boldsymbol{q} \notin \mathcal{L}_{\psi^{\prime}}$ for all $\psi^{\prime} \neq \psi(\boldsymbol{q})$ (which follows from $\boldsymbol{q} \notin \mathcal{E}$ and (7.13) ):

$$
\sum_{x \in \mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})} f(x)-\sum_{x \in \mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})} f(x)
$$

$$
\begin{equation*}
=\sum_{\boldsymbol{p} \in \mathcal{P} \backslash\{\boldsymbol{q}\}}\left(f\left((\boldsymbol{p}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho}, \varsigma(\boldsymbol{p})\right)-\sum_{\substack{\psi \in \Psi \\\left(\boldsymbol{p} \in \mathcal{L}_{\psi}\right)}} f\left((\boldsymbol{p}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho}, \sigma^{\psi}\right)\right) \tag{7.27}
\end{equation*}
$$

Set

$$
A(\boldsymbol{q}, \boldsymbol{v}, \rho):=\mathcal{P} \cap\left(\boldsymbol{q}+\mathcal{B}_{S}^{d} D_{\rho}^{-1} R(\boldsymbol{v})^{-1}\right) \backslash\{\boldsymbol{q}\}
$$

Note that for every $\boldsymbol{p} \in \mathcal{P} \backslash A(\boldsymbol{q}, \boldsymbol{v}, \rho)$ we have $(\boldsymbol{p}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \notin \mathcal{B}_{S}^{d}$, so that the corresponding term in (7.27) vanishes. Also for every $\boldsymbol{p} \notin \mathcal{E}$ we have $\boldsymbol{p} \notin \mathcal{L}_{\psi^{\prime}}$ for all $\psi^{\prime} \neq \psi(\boldsymbol{p})$ (cf. (7.13)), which implies that the corresponding term in (7.27) is bounded in absolute value by $d(\varsigma(\boldsymbol{p}))$, where the function $d: \Sigma \rightarrow \mathbb{R}_{\geq 0}$ is defined by

$$
\begin{equation*}
d(\psi, \omega)=\sup \left\{\left|f(\boldsymbol{x},(\psi, \omega))-f\left(\boldsymbol{x}, \sigma^{\psi}\right)\right|: \boldsymbol{x} \in \mathbb{R}^{d}\right\} \quad(\psi, \omega) \in \Sigma \tag{7.28}
\end{equation*}
$$

Hence for every $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$ and $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$ such that $\mathcal{E} \cap A(\boldsymbol{q}, \boldsymbol{v}, \rho)=\emptyset$, we have

$$
\begin{equation*}
\left|\sum_{x \in \mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})} f(x)-\sum_{x \in \mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})} f(x)\right| \leq \sum_{\boldsymbol{p} \in A(\boldsymbol{q}, \boldsymbol{v}, \rho)} d(\varsigma(\boldsymbol{p})) \tag{7.29}
\end{equation*}
$$

Now let $\varepsilon^{\prime}>0$ be given. Take $K \in \mathbb{Z}^{+}$and $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \#\left((\mathcal{P}-\boldsymbol{q}) R(\boldsymbol{v}) D_{\rho} \cap \mathcal{B}_{S}^{d}\right)>K\right\}\right)<\varepsilon^{\prime} \tag{7.30}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$. This is possible by Lemma 7.8. Next set

$$
\mathcal{Z}:=\{\boldsymbol{p} \in \mathcal{P}: d(\varsigma(\boldsymbol{p})) \geq \varepsilon / K\} .
$$

We claim that the set $\mathcal{Z}$ has density zero. To prove this, set

$$
U_{\psi}:=\{\omega \in \Omega:(\psi, \omega) \in \Sigma \text { and } d(\psi, \omega)<\varepsilon / K\}
$$

so that $\mathcal{Z} \subset \cup_{\psi \in \Psi}\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \omega^{(\boldsymbol{q})} \notin U_{\psi}\right\}$. Note that $\{\psi\} \times U_{\psi}$ is an open neighbourhood of $\sigma^{\psi}$ in $\Sigma$, since the function $d$ is continuous. Hence as in the proof of Proposition 7.5 (making crucial use of Proposition 5.13), the set $\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \omega^{(\boldsymbol{q})} \notin U_{\psi}\right\}$ has density zero. Hence also $\mathcal{Z}$ has density zero, as claimed.

It follows that also $\mathcal{E} \cup \mathcal{Z}$ has density zero, and so by Lemma 7.9, after possibly shrinking $\rho_{0}$, we have

$$
\begin{equation*}
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}:(\mathcal{E} \cup \mathcal{Z}) \cap A(\boldsymbol{q}, \boldsymbol{v}, \rho) \neq \emptyset\right\}\right)<\varepsilon^{\prime} \tag{7.31}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$. Now note that for any $\boldsymbol{q} \in \mathcal{P} \backslash \mathcal{E}$ and for any $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$ which belongs to neither of the two sets in (7.30) and (7.31) , the set $A(\boldsymbol{q}, \boldsymbol{v}, \rho)$ has cardinality at most $K$ and is disjoint from $\mathcal{E} \cup \mathcal{Z}$; therefore the inequality in (7.29) holds, and the right hand side in that inequality is $\leq \varepsilon$. It follows that for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{P}_{T}(\rho)$, the measure in (7.26) is less than $2 \varepsilon^{\prime}$.

Proof of Theorem 7.6. Let $T \geq 1$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ be given. Let $\rho_{1}, \rho_{2} \ldots$ be an arbitrary sequence in $(0,1)$ with $\rho_{n} \rightarrow 0$, and let $\boldsymbol{q}_{n} \in \mathcal{P}_{T}\left(\rho_{n}\right)$ for $n=1,2, \ldots$ be such that the limit $\varsigma=(\psi, \omega):=\lim _{n \rightarrow \infty} \varsigma\left(\boldsymbol{q}_{n}\right) \in \Sigma$ exists. By [21, Lemma 2.1.2], it suffices to prove that in this situation we have

$$
\begin{equation*}
\mu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} \mu_{\varsigma} \quad \text { as } n \rightarrow \infty . \tag{7.32}
\end{equation*}
$$

Since $\varsigma\left(\boldsymbol{q}_{n}\right) \rightarrow(\psi, \omega)$ implies that $\psi\left(\boldsymbol{q}_{n}\right)=\psi$ for all large $n$, we may without loss of generality assume that $\psi\left(\boldsymbol{q}_{n}\right)=\psi$ for all $n$. This means that $\boldsymbol{q}_{n} \in \mathcal{L}_{\psi}$ for all $n$.

For any $\boldsymbol{q} \in \mathcal{P}, \rho>0$ and $\lambda \in P\left(\mathrm{~S}_{1}^{d-1}\right)$, let $\nu_{\boldsymbol{q}, \rho}^{(\lambda)} \in P(\mathbb{X})$ be the distribution of $\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)$ for $\boldsymbol{v}$ random in $\left(S_{1}^{d-1}, \lambda\right)$. As a first step, let us note that (7.14) implies that

$$
\begin{equation*}
\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} \bar{\omega} \quad \text { as } n \rightarrow \infty \tag{7.33}
\end{equation*}
$$

Indeed, let $f \in \mathrm{C}_{b}(\mathbb{X})$ be given. Then by (7.14) we have $\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}(f)-\overline{\omega^{\left(\boldsymbol{q}_{n}\right)}}(f) \rightarrow 0$ as $n \rightarrow \infty$. Also $\varsigma\left(\boldsymbol{q}_{n}\right) \rightarrow(\psi, \omega)$ implies that $\omega_{j}^{\left(\boldsymbol{q}_{n}\right)} \xrightarrow{\mathrm{w}} \omega_{j}$ in $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ for each $j \in\{1, \ldots, N\}$; hence by Lemma $\frac{4.6}{(\lambda)}$ and [1, Thm. 2.8(ii)], we have $\overline{\omega^{\left(\boldsymbol{q}_{n}\right)}} \xrightarrow{\mathrm{w}} \bar{\omega}$ in $P(\mathbb{X})$, and thus $\overline{\omega^{\left(\boldsymbol{q}_{n}\right)}}(f) \rightarrow \bar{\omega}(f)$. Hence $\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}(f) \rightarrow \bar{\omega}(f)$, and (7.33) is proved.

Next, for each $n$ we have $\Gamma g_{0}^{\left(\boldsymbol{q}_{n}\right)} \in \mathbb{X}^{\psi}$ by Lemma 6.1, and hence $\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}\left(\mathbb{X}^{\psi}\right)=1$; also $\bar{\omega}\left(\mathbb{X}^{\psi}\right)=1$ since $(\psi, \omega) \in \Sigma$; cf. (6.15). Hence all $\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}$ as well as $\bar{\omega}$ may be regarded as elements in $P\left(\mathbb{X}^{\psi}\right)$, and (7.33) implies that $\nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} \bar{\omega}$ also in $P\left(\mathbb{X}^{\psi}\right)$ [13, Lemma 4.26]. Hence by the continuous mapping theorem,

$$
\begin{equation*}
J_{\psi *} \nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} J_{\psi *} \bar{\omega} \quad \text { as } n \rightarrow \infty \tag{7.34}
\end{equation*}
$$

Here $J_{\psi *} \bar{\omega}=\mu_{\varsigma}$, by (6.14). Now let $f \in \mathrm{C}_{c}(\mathcal{X})$ be given, and let $\pi_{f}$ be the continuous map from $N_{s}(\mathcal{X})$ to $\mathbb{R}$ given by $\pi_{f}(Q)=\sum_{x \in Q} f(x)$. Then (7.34) implies that

$$
\begin{equation*}
\pi_{f *} J_{\psi *} \nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} \pi_{f *} \mu_{\varsigma} \quad \text { as } n \rightarrow \infty \tag{7.35}
\end{equation*}
$$

But note that for each $\boldsymbol{q} \in \mathcal{P}$, by Lemma 7.10, $J_{\psi(\boldsymbol{q}) *} \nu_{\boldsymbol{q}, \rho}^{(\lambda)}$ is the distribution of $\mathcal{Q}_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})$ in $N_{s}(\mathcal{X})$ for $\boldsymbol{v}$ random in $\left(\mathrm{S}_{1}^{d-1}, \lambda\right)$. Hence $\pi_{f *} J_{\psi *} \nu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}$ is the distribution of the realvalued random variable $F^{\prime}(\boldsymbol{v})=\sum_{x \in \mathcal{Q}_{\rho_{n}}^{\prime}\left(\boldsymbol{q}_{n}, \boldsymbol{v}\right)} f(x)$, for $\boldsymbol{v}$ random in $\left(\mathrm{S}_{1}^{d-1}, \lambda\right)$. Similarly, $\pi_{f *} \mu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)}$ is the distribution of the real-valued random variable $F(\boldsymbol{v})=\sum_{x \in \mathcal{Q}_{\rho_{n}}\left(\boldsymbol{q}_{n}, \boldsymbol{v}\right)} f(x)$. By Lemma 7.11, $\left|F(\boldsymbol{v})-F^{\prime}(\boldsymbol{v})\right|$ converges in probability to 0 . Hence by [1, Thm. 3.1], (7.35) implies that

$$
\begin{equation*}
\pi_{f *} \mu_{\boldsymbol{q}_{n}, \rho_{n}}^{(\lambda)} \xrightarrow{\mathrm{w}} \pi_{f *} \mu_{\varsigma} \quad \text { as } n \rightarrow \infty . \tag{7.36}
\end{equation*}
$$

We have proved that this holds for any given $f \in \mathrm{C}_{c}(\mathcal{X})$. By [13, Thm. 16.16(ii) $\Rightarrow(\mathrm{i})$ ], this implies that (7.32) holds.

### 7.4. Verification of $[\mathrm{P} 3]$, and the macroscopic limit.

Proposition 7.12. [P3] holds, i.e. for every bounded Borel set $B \subset \mathbb{R}^{d}$ we have

$$
\begin{equation*}
\lim _{\xi \rightarrow \infty} \limsup _{\rho \rightarrow 0}[\operatorname{vol} \times \sigma]\left(\left\{(\boldsymbol{q}, \boldsymbol{v}) \in B \times \mathrm{S}_{1}^{d-1}: \mathcal{Q}_{\rho}\left(\rho^{1-d} \boldsymbol{q}, \boldsymbol{v}\right) \cap\left(\mathfrak{Z}_{\xi} \times \Sigma\right)=\emptyset\right\}\right)=0 \tag{7.37}
\end{equation*}
$$

Proof. Using the fact that $\mathcal{P}$ contains at least one grid, the proposition follows from the existence of a limit distribution on $\mathbb{R}_{>0}$ (with zero mass at $+\infty$ ) for the macroscopic free path length in the Boltzmann-Grad limit of the Lorentz gas on a lattice scatterer configuration, [16. Theorem 1.2]. Indeed, fix an arbitrary $\psi \in \Psi$. For any $\rho>0$ and $\boldsymbol{q} \in \mathbb{R}^{d}, \boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}$, set

$$
Q_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v})=\left(\mathcal{L}_{\psi}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho} \quad\left(\subset \mathbb{R}^{d}\right)
$$

Comparing with the definition of $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ in (2.5), (2.4), and using $\mathcal{L}_{\psi} \subset \mathcal{P}$ and $\mathbf{0} \notin \mathfrak{Z} \xi$, one verifies that for any $\xi>0, \mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v}) \cap\left(\mathfrak{Z}_{\xi} \times \Sigma\right)=\emptyset$ forces $Q_{\rho}^{\prime}(\boldsymbol{q}, \boldsymbol{v}) \cap \mathfrak{Z}_{\xi}=\emptyset$. Hence to prove the proposition it suffices to prove that

$$
\begin{equation*}
\lim _{\xi \rightarrow \infty} \limsup _{\rho \rightarrow 0}[\operatorname{vol} \times \sigma]\left(\left\{(\boldsymbol{q}, \boldsymbol{v}) \in B \times \mathrm{S}_{1}^{d-1}: Q_{\rho}^{\prime}\left(\rho^{1-d} \boldsymbol{q}, \boldsymbol{v}\right) \cap \mathfrak{Z}_{\xi}=\emptyset\right\}\right)=0 . \tag{7.38}
\end{equation*}
$$

By a simple translation and rescaling argument we may reduce to the case when $\mathcal{L}_{\psi}$ has covolume one and is a lattice, and then (7.38) is a simple consequence of [16, Theorem 1.2].

In [21, Sec. 2.5], the condition [P3] is used to prove, for an arbitrary fixed point set $\mathcal{P} \subset \mathbb{R}^{d}$ satisfying the hypotheses in Section 2, the existence of a canonical measure $\mu^{\mathrm{g}} \in P\left(N_{s}(\mathcal{X})\right)^{12}$ giving the limit distribution of $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ in the case of a macroscopic initial condition. That is, $\mu^{\mathbf{g}}$ equals the limit distribution of $\mathcal{Q}_{\rho}\left(\rho^{1-d} \boldsymbol{q}, \boldsymbol{v}\right)$ for $(\boldsymbol{q}, \boldsymbol{v})$ random in $\left(\mathrm{T}^{1}\left(\mathbb{R}^{d}\right), \Lambda\right)$, for any fixed probability measure $\Lambda \in P\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right.$ ) absolutely continuous with respect to the Liouville measure vol $\times \sigma$ [21, Theorem 2.19]. This measure $\mu^{\mathrm{g}}$ also appears in the definition of the transition kernel for generic initial data, $k^{\mathrm{g}}$; see Section 10.1 below.

In our case of $\mathcal{P}$ being a finite union of grids as in (3.5), the macroscopic limit measure $\mu^{\mathrm{g}}$ can be explicitly defined as follows: Set

$$
\begin{equation*}
\omega^{\mathrm{g}}:=\left(\omega_{1}^{\mathrm{g}}, \ldots, \omega_{N}^{\mathrm{g}}\right) \in \Omega, \tag{7.39}
\end{equation*}
$$

and then let

$$
\begin{equation*}
\mu^{\mathrm{g}}=J_{*}(\overline{\omega \mathrm{~m}}), \tag{7.40}
\end{equation*}
$$

with $J: \mathbb{X} \rightarrow N_{s}(\mathcal{X})$ being the map in (6.12).
We next state without proof a limit result which significantly strengthens the above mentioned [21, Theorem 2.19] for our special class of $\mathcal{P}$. For any $\Lambda \in P\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right), s>0$ and $\rho \in(0,1)$, let $\mu_{\rho}^{(\Lambda, s)}$ be the distribution of $\mathcal{Q}_{\rho}(s \boldsymbol{q}, \boldsymbol{v})$ for $(\boldsymbol{q}, \boldsymbol{v})$ random in $\left(\mathrm{T}^{1}\left(\mathbb{R}^{d}\right), \Lambda\right)$.

Theorem 7.13. For any $\Lambda \in P\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right)$ which is absolutely continuous with respect to Liouville measure, and any $s_{0}>0$, we have $\mu_{\rho}^{(\Lambda, s)} \xrightarrow{\mathrm{W}} \mu^{\mathrm{g}}$ as $\rho \rightarrow 0$, uniformly over all $s \geq s_{0}$.

Note that [21, Theorem 2.19] corresponds to the particular choice $s=\rho^{1-d}$ in Theorem 7.13, The formulation of Theorem [7.13 is inspired by [19, Theorem 1.1].

As mentioned, we will not give the proof of Theorem 7.13 in the present paper. However we remark that Theorem 7.13 can be deduced, by similar arguments as in Section 7.3, from the following equidistribution result in the homogeneous space $\mathbb{X}$, which is a kind of macroscopic analogue of Theorem 4.2.
Theorem 7.14. For any $\Lambda \in P\left(\mathrm{~T}^{1}\left(\mathbb{R}^{d}\right)\right)$ which is absolutely continuous with respect to Liouville measure, and any $f \in \mathrm{C}_{b}(\mathbb{X})$ and $s_{0}>0$, we have

$$
\begin{equation*}
\int_{\mathrm{T}^{1}\left(\mathbb{R}^{d}\right)} f\left(\Gamma g_{0}^{(s \boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \Lambda(\boldsymbol{q}, \boldsymbol{v}) \rightarrow \int_{\mathbb{X}} f d \overline{\omega^{\mathrm{g}}} \tag{7.41}
\end{equation*}
$$

as $\rho \rightarrow 0$, uniformly over all $s \geq s_{0}$.
We will not give the proof of Theorem 7.14 either; however we note that it is to a large extent similar to the proof of Theorem 7.7 which we give in Section 9 below.

[^8]We end this section by pointing out a couple of invariance properties of $\mu^{\mathrm{g}}$. By the same argument as in the proof of Lemma 7.1, the measure $\overline{\omega^{g}}$ on $\mathbb{X}$ is $\varphi\left(\mathrm{SL}_{d}(\mathbb{R})\right)$-invariant, and thus $\mu^{\mathrm{g}}$ is invariant under the action of $\mathrm{SL}_{d}(\mathbb{R})$ on $N_{s}(\mathcal{X})$. We also have:
Lemma 7.15. The measure $\mu^{\mathrm{g}}$ is invariant under the action of $\mathbb{R}^{d}$ on $N_{s}(\mathcal{X})$ by translations.
Proof. This is part of [21, Prop. 2.24]; however let us note that it also follows directly from the explicit definition in (7.40). Indeed, it follows from (5.12) that $\mathbb{R} \widetilde{\boldsymbol{c}}_{j} \subset L_{j}$; hence $\widetilde{\boldsymbol{c}}_{j} \boldsymbol{v} \in L_{j}^{d}$ for all $\boldsymbol{v} \in \mathbb{R}^{d}$ (where $\widetilde{\boldsymbol{c}}_{j} \boldsymbol{v}$ is the matrix product of $\widetilde{\boldsymbol{c}}_{j} \in \mathrm{M}_{r_{j} \times 1}(\mathbb{R})$ and $\boldsymbol{v} \in \mathrm{M}_{1 \times d}(\mathbb{R})$ ), and so $\omega_{j}^{\mathrm{g}}$ is invariant under the translation $X \mapsto X+\widetilde{\boldsymbol{c}}_{j} \boldsymbol{v}$ on $\mathbb{T}_{j}^{d}$, for every $\boldsymbol{v} \in \mathbb{R}^{d}$. This implies that the measure $\overline{\overline{\omega_{j}^{\mathbf{g}}}}$ on $G_{j}\left(\right.$ see (4.9)) is invariant under $g \mapsto g \mathrm{I}_{\tilde{\boldsymbol{c}}_{j} \boldsymbol{v}}$ for every $\boldsymbol{v} \in \mathbb{R}^{d}$, and hence $\overline{\omega_{j}^{g}}$ is invariant under the translation $x \mapsto x \mathrm{I}_{\tilde{c}_{j} v}$ on $\mathbb{X}_{j}$. Now the lemma follows from the definition of $\mu^{\mathrm{g}}$ in (7.40) by using the formula

$$
J\left(x_{1} \mathrm{I}_{\widetilde{\boldsymbol{c}}_{1} \boldsymbol{v}}, \ldots, x_{N} \mathrm{I}_{\widetilde{\boldsymbol{c}}_{N} \boldsymbol{v}}\right)=J\left(x_{1}, \ldots, x_{N}\right)+\boldsymbol{v} \quad\left(\forall\left(x_{1}, \ldots, x_{N}\right) \in \mathbb{X}, \boldsymbol{v} \in \mathbb{R}^{d}\right)
$$

which is immediate from (5.1) and (6.12).

## 8. Application of the classification of invariant measures of unipotent flows

In this section we will state and prove a result, Theorem 8.1, on the equidistribution of certain expanding unipotent orbits in a slightly generalized version of the homogeneous space $\mathbb{X}$ introduced in Section 3.2. This theorem is tailor-made to serve as the main ingredient in the proof of Theorem 7.7 which we give in Section 9 below; in particular, it will be crucial for us to have a certain uniformity with respect to the position of the initial point in the torus fiber (that is, uniformity with respect to the variable " $V$ " in Theorem 8.1 below). The proof of Theorem 8.1 builds on Ratner's classification of ergodic measures invariant under unipotent flows [25] and further characterization results by Mozes and Shah [23]. We also remark that if it were not for the uniformity requirement, Theorem 8.1 could be deduced as a consequence of Shah [27, Theorem 1.4].

We start by introducing some notation. We stress that in this Section 8, some of our notation (for example, "X", " $G$ ", " $\Gamma$ ", " $\Gamma_{j} "$ and " $\widetilde{M} "$ ) will be used in a slightly different and more general way than in all the other sections of the paper. The reason is that the results of the present section will be applied, in Section 9, to certain homogeneous submanifolds of our original space " $\mathbb{X}$ " (see the proofs of Theorems 9.1 and 9.2 ). To start with, similarly as before, we set

$$
G=G_{1} \times \cdots \times G_{N}=\mathrm{S}_{r_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{R})
$$

however now we allow $r_{1}, \ldots, r_{N}$ to be arbitrary (fixed) non-negative integers. That is, unlike all the other sections, we allow one or several of the $r_{j}$ s to be zero, with the natural convention that $\mathrm{S}_{0}(\mathbb{R}):=\mathrm{SL}_{d}(\mathbb{R})$. Next, we fix $\Gamma_{1}^{\prime}, \ldots, \Gamma_{N}^{\prime}$ to be arbitrary, fixed, finite index subgroups of $\mathrm{SL}_{d}(\mathbb{Z})$, and set

$$
\begin{equation*}
\Gamma_{j}=\Gamma_{j}^{\prime} \ltimes \mathrm{M}_{r_{j} \times d}(\mathbb{Z})=\left\{(M, U) \in \mathrm{S}_{r_{j}}(\mathbb{Z}): M \in \Gamma_{j}^{\prime}\right\} \quad(j=1, \ldots, N) \tag{8.1}
\end{equation*}
$$

(if $r_{j}=0$, this should be understood as $\Gamma_{j}=\Gamma_{j}^{\prime}$ ), and

$$
\Gamma=\Gamma_{1} \times \cdots \times \Gamma_{N}
$$

Then, as before, we set $\mathbb{X}_{j}:=\Gamma_{j} \backslash G_{j}$ and

$$
\mathbb{X}:=\Gamma \backslash G=\mathbb{X}_{1} \times \cdots \times \mathbb{X}_{N}
$$

and write $p_{j}: G \rightarrow \mathrm{~S}_{r_{j}}(\mathbb{R})$ and $\widetilde{p}_{j}: \mathbb{X} \rightarrow \mathbb{X}_{j}(j=1, \ldots, N)$ for the projection maps.
Recall that we consider $\mathrm{SL}_{d}(\mathbb{R})$ to be an embedded subgroup of each group $G_{j}$, through $M \mapsto(M, 0)$. Now we also set $G^{\prime}:=\mathrm{SL}_{d}(\mathbb{R})^{N}$; this is an embedded subgroup of $G$. We also set

$$
\Gamma^{\prime}:=\Gamma_{1}^{\prime} \times \cdots \times \Gamma_{N}^{\prime} \subset G^{\prime}
$$

and

$$
\mathbb{X}^{\prime}=\Gamma^{\prime} \backslash G^{\prime}
$$

Recall that we have a projection morphism $\iota: G_{j} \rightarrow \mathrm{SL}_{d}(\mathbb{R})$ for each $j$, and $\iota$ induces a projection map

$$
\begin{equation*}
\widetilde{\iota}: \mathbb{X}_{j} \rightarrow \Gamma_{j}^{\prime} \backslash \mathrm{SL}_{d}(\mathbb{R}), \quad \widetilde{\iota}\left(\Gamma_{j} g\right)=\Gamma_{j}^{\prime} \iota(g) \quad\left(g \in G_{j}\right), \tag{8.2}
\end{equation*}
$$

generalizing (3.22). In the present section we will also write $\iota$ for the product morphism from $G$ to $G^{\prime}$, and write $\tau$ for the induced projection map from $\mathbb{X}$ to $\mathbb{X}^{\prime}$.

As before we set

$$
\mathbb{T}_{j}:=\mathbb{R}^{r_{j}} / \mathbb{Z}^{r_{j}} \quad \text { and } \quad \mathbb{T}_{j}^{d}:=\underbrace{\mathbb{T}_{j} \times \cdots \times \mathbb{T}_{j}}_{d \text { copies }} ;
$$

if $r_{j}=0$ this should be understood as $\mathbb{T}_{j}=\mathbb{T}_{j}^{d}=\{\mathbf{0}\}$, the trivial group. The definition in (3.24) of the embedding $x: \mathbb{T}_{j}^{d} \rightarrow \mathbb{X}_{j}$ carries over unchanged to our present setting, although our " $\mathbb{X}_{j}$ " is now more general. (If $r_{j}=0$ then we set $x(\{\mathbf{0}\}):=\Gamma_{j} \in \mathbb{X}_{j}$.) We now also set

$$
\begin{equation*}
\widetilde{\mathbb{T}}:=\mathbb{T}_{1}^{d} \times \mathbb{T}_{2}^{d} \times \cdots \times \mathbb{T}_{N}^{d}, \tag{8.3}
\end{equation*}
$$

and let $\widetilde{p}_{j}: \widetilde{\mathbb{T}} \rightarrow \mathbb{T}_{j}^{d}(j=1, \ldots, N)$ be the projection maps; and we will write " $x$ " also for the map $\widetilde{\mathbb{T}} \rightarrow \mathbb{X}$ which is the product of the maps $x: \mathbb{T}_{j}^{d} \rightarrow \mathbb{X}_{j}$. The fact that both " $x$ " and " $\widetilde{p}_{j}$ " now denote more than one map should not cause any confusion; in particular note that with this abuse of notation we have $x \circ \widetilde{p}_{j}=\widetilde{p}_{j} \circ x: \widetilde{\mathbb{T}} \rightarrow \mathbb{X}_{j}$ for each $j$.

We now come to the statement of the main result of the present section, Theorem 8.1. It concerns the equidistribution of pieces of expanding unipotent orbits in $\mathbb{X}$ of the form

$$
\begin{equation*}
\left\{x(V) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho}\right): \boldsymbol{u} \in \mathbb{R}^{d-1}\right\} \tag{8.4}
\end{equation*}
$$

where $V \in \widetilde{\mathbb{T}} ; \widetilde{M}$ is an arbitrary element in $G^{\prime}$ not belonging to the subset

$$
\begin{equation*}
\mathfrak{D}_{\mathcal{S}}:=\bigcup_{i<j}\left\{\left(M_{1}, \ldots, M_{N}\right) \in G^{\prime}: M_{i} M_{j}^{-1} \in \mathcal{S}\right\} \tag{8.5}
\end{equation*}
$$

with $\mathcal{S}$ as in (3.8); $\varphi$ is the diagonal embedding of $\mathrm{SL}_{d}(\mathbb{R})$ in $G$; and finally

$$
n_{-}(\boldsymbol{u}):=\left(\begin{array}{cc}
1 & \boldsymbol{u}  \tag{8.6}\\
0 & I_{d-1}
\end{array}\right) \in \mathrm{SL}_{d}(\mathbb{R})
$$

(block diagonal notation). The equidistribution is with respect to the $G$-invariant probability measure on $\mathbb{X}$, which we call $\mu$.

In order for orbits of the form (8.4) to equidistribute in $(\mathbb{X}, \mu)$ as $\rho \rightarrow 0$, we have to assume that $V$ avoids a certain 'singular' subset $\Delta_{k}^{(\eta)}$ of $\widetilde{\mathbb{T}}$, which we now introduce. For each $j \in\{1, \ldots, N\}$ with $r_{j} \neq 0$, let us write $\pi_{j}$ for the projection from $\left(\mathbb{R}^{r_{j}}\right)^{d}$ to $\mathbb{T}_{j}^{d}$ (it was called " $\pi$ " in (3.23)). Then for any $q \in \mathbb{Z}^{+}$and $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \backslash\{\mathbf{0}\}$, we set

$$
\begin{equation*}
\Delta_{j, q, \boldsymbol{m}}:=\pi_{j}\left(\left(q^{-1} \mathbb{Z}^{r_{j}}+\boldsymbol{m}^{\perp}\right)^{d}\right) \subset \mathbb{T}_{j}^{d} \tag{8.7}
\end{equation*}
$$

where $\boldsymbol{m}^{\perp}$ is the orthogonal complement of $\boldsymbol{m}$ in $\mathbb{R}^{r_{j}}$. Also, for any $k \in \mathbb{Z}^{+}$, we set

$$
\begin{equation*}
\Delta_{j, k}:=\bigcup_{q=1}^{k} \bigcup_{\substack{\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \\ 0<\|\boldsymbol{m}\| \leq k}} \Delta_{j, q, \boldsymbol{m}} \tag{8.8}
\end{equation*}
$$

Note that $\Delta_{j, q, m}$ and $\Delta_{j, k}$ are only defined when $r_{j} \neq 0$, in which case they are both closed regular submanifolds of $\mathbb{T}_{j}^{d}$ of codimension $d$. Next, for any $\eta>0$ we define $\Delta_{j, k}^{(\eta)}$ to be the open
$\eta$-neighbourhood of $\Delta_{j, k}$ in $\mathbb{T}_{j}^{d}$, with respect to the metric induced by the standard Euclidean metric on $\mathrm{M}_{r_{j} \times d}(\mathbb{R})=\left(\mathbb{R}^{r_{j}}\right)^{d}$. Finally, we set

$$
\begin{equation*}
\Delta_{k}^{(\eta)}=\bigcup_{\substack{j=1 \\\left(r_{j} \neq 0\right)}}^{N} \widetilde{p}_{j}^{-1}\left(\Delta_{j, k}^{(\eta)}\right) \subset \widetilde{\mathbb{T}} \tag{8.9}
\end{equation*}
$$

Let $P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right)$ be the set of Borel probability measures on $\mathbb{R}^{d-1}$ which are absolutely continuous with respect to Lebesgue measure.

Theorem 8.1. Let $f \in \mathrm{C}_{b}(\mathbb{X})$ and $\varepsilon>0$ be given. Then there exists some $k \in \mathbb{Z}^{+}$such that for every $\lambda \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\left|\int_{\mathbb{R}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho}\right)\right) d \lambda(\boldsymbol{u})-\int_{\mathbb{X}} f d \mu\right|<\varepsilon \tag{8.10}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$.
The rest of this section is devoted to the proof of Theorem 8.1.
Definition 8.1. For each $k \in \mathbb{Z}^{+}$, we let $P_{k}$ be the set of all measures $\nu \in P(\mathbb{X})$ which can be obtained as a weak limit of a sequence of probability measures $\nu_{1}, \nu_{2}, \ldots$ given by

$$
\begin{equation*}
\nu_{m}: \quad f \mapsto \int_{\mathbb{R}^{d-1}} f\left(x\left(V_{m}\right) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right) d \lambda(\boldsymbol{u}) \quad\left(f \in \mathrm{C}_{b}(\mathbb{X})\right) \tag{8.11}
\end{equation*}
$$

for some $\lambda \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right), \widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, real numbers $\rho_{1}>\rho_{2}>\cdots \rightarrow 0$, and points $V_{1}, V_{2}, \ldots$ in $\widetilde{\mathbb{T}}$ such that $\left[\exists \eta>0: \forall m \in \mathbb{Z}^{+}: V_{m} \notin \Delta_{k}^{(\eta)}\right]$.

Note that $P_{1} \supset P_{2} \supset \cdots$, since $\Delta_{k}^{(\eta)} \subset \Delta_{k^{\prime}}^{(\eta)}$ whenever $k<k^{\prime}$.
Throughout the rest of this section, we will let $W$ denote the following subgroup of $G$ :

$$
W:=\left\{\varphi\left(n_{-}(\boldsymbol{w})\right): \boldsymbol{w} \in \mathbb{R}^{d-1}\right\} .
$$

Lemma 8.2. Every $\nu \in P_{k}$ is $W$-invariant.
Proof. This is a (very) standard consequence of the fact that for $\rho<1$, the action (from the right) of $\varphi\left(D_{\rho}\right)$ on $\mathbb{X}$ expands any $W$-orbit. The details are as follows. Let $\nu \in P_{k}$ be given. Then the task is to prove that for any given $\boldsymbol{w} \in \mathbb{R}^{d-1}$ and $f \in \mathrm{C}_{b}(\mathbb{X})$ we have $\nu\left(f \circ R_{\boldsymbol{w}}\right)=\nu(f)$, where $R_{\boldsymbol{w}}: \mathbb{X} \rightarrow \mathbb{X}$ denotes right multiplication by $\varphi\left(n_{-}(\boldsymbol{w})\right)$. Choose $\lambda, \widetilde{M},\left(\rho_{m}\right),\left(V_{m}\right)$ as in Definition 8.1, so that $\nu$ is the weak limit of the measures $\nu_{m}$ given by (8.11). Using the relation $D_{\rho} n_{-}(\boldsymbol{w})=n_{-}\left(\rho^{d} \boldsymbol{w}\right) D_{\rho}$, and writing $\lambda^{\prime} \in \mathrm{L}^{1}\left(\mathbb{R}^{d-1}\right)$ for the density of $\lambda$ with respect to Lebesgue measure, we now have:

$$
\nu_{m}\left(f \circ R_{\boldsymbol{w}}\right)=\int_{\mathbb{R}^{d-1}} f\left(x\left(V_{m}\right) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right) \lambda^{\prime}\left(\boldsymbol{u}-\rho_{m}^{d} \boldsymbol{w}\right) d \boldsymbol{u}
$$

Hence $\left|\nu_{m}\left(f \circ R_{\boldsymbol{w}}\right)-\nu_{m}(f)\right| \leq\|f\|_{\mathrm{L}^{\infty}} \cdot\left\|\tau_{\rho_{m}^{d} \boldsymbol{w}} \lambda^{\prime}-\lambda^{\prime}\right\|_{\mathrm{L}^{1}\left(\mathbb{R}^{d-1}\right)}$, and so by [8, Prop. 8.5] we have $\lim _{m \rightarrow \infty} \nu_{m}\left(f \circ R_{\boldsymbol{w}}\right)=\lim _{m \rightarrow \infty} \nu_{m}(f)$, that is, $\nu\left(f \circ R_{\boldsymbol{w}}\right)=\nu(f)$.

Recall that we write $\widetilde{\iota}$ for the natural projection map from $\mathbb{X}$ to $\mathbb{X}^{\prime}$; in particular $\widetilde{\mu}:=\widetilde{\iota}_{*} \mu$ is the unique $G^{\prime}$-invariant probability measure on $\mathbb{X}^{\prime}$.

Lemma 8.3. Every $\nu \in P_{k}$ satisfies $\widetilde{\iota}_{*} \nu=\widetilde{\mu}$.
Proof. Let $\nu \in P_{k}$ be given, and let $\left(\nu_{m}\right)$ be a sequence as in Definition 8.1, tending weakly to $\nu$. For any $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\prime}\right)$ we have

$$
\nu_{m}(f \circ \widetilde{\iota})=\int_{\mathbb{R}^{d-1}} f\left(\widetilde{\iota}\left(x\left(V_{m}\right) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right)\right) d \lambda(\boldsymbol{u})=\int_{\mathbb{R}^{d-1}} f\left(\Gamma^{\prime} \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right) d \lambda(\boldsymbol{u})
$$

This integral tends to $\widetilde{\mu}(f)$ as $m \rightarrow \infty$, by [20, Thm. 5] applied to the function $g \mapsto f\left(\Gamma^{\prime} \widetilde{M} g\right)$ (which is left $\prod_{j=1}^{N}\left(M_{j}^{-1} \Gamma_{j}^{\prime} M_{j}\right)$-invariant). On the other hand, by the definition of $\nu$ we have $\nu_{m}(f \circ \widetilde{\iota}) \rightarrow \nu(f \circ \widetilde{\iota})=\left(\widetilde{\iota}_{*} \nu\right)(f)$. Hence $\left(\widetilde{\iota}_{*} \nu\right)(f)=\widetilde{\mu}(f)$.

Recall that a subgroup $U$ of $G$ is said to be unipotent, if the linear automorphism $\operatorname{Ad}(u)$ of the Lie algebra of $G$ is unipotent for all $u \in U$. For any $h \in G$ let us write $R_{h}: \mathbb{X} \rightarrow \mathbb{X}$ for the map $\Gamma g \mapsto \Gamma g h$. For any $\alpha \in P(\mathbb{X})$, let us define $H_{\alpha}$ to be the identity component of the subgroup of $G$ consisting of all $g$ which preserve $\alpha$;

$$
\begin{equation*}
H_{\alpha}:=\left\{g \in G: R_{g *} \alpha=\alpha\right\}^{\circ} . \tag{8.12}
\end{equation*}
$$

This is a closed connected Lie subgroup of $G$. We let $\mathcal{Q}(\mathbb{X})$ be the set of all $\alpha \in P(\mathbb{X})$ such that the group generated by all unipotent one-parameter subgroups of $G$ contained in $H_{\alpha}$ acts ergodically on $\mathbb{X}$ with respect to $\alpha$. (Note that this definition of $\mathcal{Q}(\mathbb{X})$ is equivalent to the one in [23, p. 150], although our $H_{\alpha}$ equals " $\Lambda(\alpha)^{\circ}$ " in the notation of [23].)

A key ingredient in our proof of Theorem 8.1 will be Ratner's classification of invariant measures of unipotent flows, [25, Thm. 1]. Applied in our setting, this result says that for every $\alpha \in \mathcal{Q}(\mathbb{X})$, there is some $g_{\alpha} \in G$ such that $\alpha\left(\Gamma \backslash \Gamma g_{\alpha} H_{\alpha}\right)=1$. Note that in this situation, $\Gamma \cap g_{\alpha} H_{\alpha} g_{\alpha}^{-1}$ is a lattice in $g_{\alpha} H_{\alpha} g_{\alpha}^{-1}$, and the support of $\alpha$ equals $\Gamma \backslash \Gamma g_{\alpha} H_{\alpha}$, which is a smooth embedded submanifold of $\mathbb{X}$.

Lemma 8.4. For any $\alpha \in \mathcal{Q}(\mathbb{X})$ such that $\widetilde{\iota_{*}} \alpha=\widetilde{\mu}$, we have $\iota\left(H_{\alpha}\right)=G^{\prime}$.
Proof. (This generalizes [6, Lemma 6], and the proof is the same.) Using the fact that the $\operatorname{map} \iota: \mathbb{X} \rightarrow \mathbb{X}^{\prime}$ has compact fibers, we have $\widetilde{\iota}(\operatorname{supp} \alpha)=\operatorname{supp} \widetilde{\iota}_{*} \alpha=\operatorname{supp} \widetilde{\mu}=\mathbb{X}^{\prime}$. But $\operatorname{supp} \alpha=\Gamma \backslash \Gamma g_{\alpha} H_{\alpha}$. Hence $\Gamma^{\prime} \iota\left(g_{\alpha}\right) \iota\left(H_{\alpha}\right)=G^{\prime}$, and thus $\iota\left(H_{\alpha}\right)=G^{\prime}$.

Next, using basic Lie group and Lie algebra theory, we will derive a completely explicit description of any Lie subgroup $H_{\alpha}$ as in Lemma 8.4, cf. Lemma 8.6 below.

For any $r \in \mathbb{Z}_{\geq 0}$, let $\mathfrak{s}_{r}(\mathbb{R})$ be the Lie algebra of $\mathrm{S}_{r}(\mathbb{R})$, which we represent as the set of pairs $(A, X) \in \mathfrak{s l}_{d}(\mathbb{R}) \times \mathrm{M}_{r \times d}(\mathbb{R})$, with the Lie bracket given by

$$
\begin{equation*}
\left[\left(A_{1}, X_{1}\right),\left(A_{2}, X_{2}\right)\right]=\left(\left[A_{1}, A_{2}\right], X_{1} A_{2}-X_{2} A_{1}\right) . \tag{8.13}
\end{equation*}
$$

(For $r=0$ we have $\mathfrak{s}_{0}(\mathbb{R})=\mathfrak{s l}_{d}(\mathbb{R})$, and in (8.13) we should view $\mathrm{M}_{0 \times d}(\mathbb{R})$ as a singleton set containing only the "empty matrix".) Just as for the Lie groups, we always consider $\mathfrak{s l}_{d}(\mathbb{R})$ to be embedded in $\mathfrak{s}_{r}(\mathbb{R})$ through $A \mapsto(A, 0)$. We also set $\mathfrak{g}=\mathfrak{s}_{r_{1}}(\mathbb{R}) \oplus \cdots \oplus \mathfrak{s}_{r_{N}}(\mathbb{R})$ and $\mathfrak{g}^{\prime}=\mathfrak{s l}_{d}(\mathbb{R})^{N}$; these are the Lie algebras of $G$ and of $G^{\prime}$, respectively. Next, as in (4.1), given any linear subspace $L$ of $\mathbb{R}^{r}$, we let $S_{L}(\mathbb{R})$ be the closed connected subgroup of $S_{r}(\mathbb{R})$ given by

$$
\begin{equation*}
\mathrm{S}_{L}(\mathbb{R}):=\mathrm{SL}_{d}(\mathbb{R}) \ltimes L^{d}=\left\{(M, U) \in \mathrm{S}_{r}(\mathbb{R}): U \in L^{d}\right\} \tag{8.14}
\end{equation*}
$$

Recall here that via our identification $\mathrm{M}_{r \times d}(\mathbb{R})=\left(\mathbb{R}^{r}\right)^{d}$, $L^{d}$ is the set of matrices in $\mathrm{M}_{r \times d}(\mathbb{R})$ all of whose column vectors lie in $L$. For any matrix $X \in \mathrm{M}_{r \times d}(\mathbb{R})$, we also write:

$$
\begin{equation*}
\mathrm{S}_{L}^{X}(\mathbb{R}):=\mathrm{I}_{X} \mathrm{~S}_{L}(\mathbb{R}) \mathrm{I}_{X}^{-1} \tag{8.15}
\end{equation*}
$$

(For $r=0$ we have $\mathbb{R}^{r}=\{\mathbf{0}\}$, the only linear subspace $L \subset \mathbb{R}^{0}$ is $L=\mathbb{R}^{0}$, and the only matrix in $\mathrm{M}_{0 \times d}(\mathbb{R})$ is $X=$ the empty matrix, and for these $L, X$ we have $\mathrm{S}_{L}(\mathbb{R})=\mathrm{S}_{L}^{X}(\mathbb{R})=\mathrm{SL}_{d}(\mathbb{R})$.) We write $\mathfrak{s}_{L}(\mathbb{R})$ and $\mathfrak{s}_{L}^{X}(\mathbb{R})$ for the Lie subalgebras of $\mathfrak{s}_{r}(\mathbb{R})$ corresponding to $S_{L}(\mathbb{R})$ and $S_{L}^{X}(\mathbb{R})$, respectively. Thus in particular,

$$
\begin{equation*}
\mathfrak{s}_{L}(\mathbb{R})=\left\{(A, Y) \in \mathfrak{s}_{r}(\mathbb{R}): Y \in L^{d}\right\} . \tag{8.16}
\end{equation*}
$$

Recall that we write $\iota$ for the natural projection $\mathrm{S}_{r}(\mathbb{R}) \rightarrow \mathrm{SL}_{d}(\mathbb{R})$; hence $d \iota$ is the natural projection $\mathfrak{s}_{r}(\mathbb{R}) \rightarrow \mathfrak{s l}_{d}(\mathbb{R})$.
Lemma 8.5. If $\mathfrak{h}$ is a Lie subalgebra of $\mathfrak{s}_{r}(\mathbb{R})$ satisfying $d \iota(\mathfrak{h})=\mathfrak{s l}_{d}(\mathbb{R})$, then there exist a linear subspace $L \subset \mathbb{R}^{r}$ and a matrix $X \in \mathrm{M}_{r \times d}(\mathbb{R})$ such that $\mathfrak{h}=\mathfrak{s}_{L}^{X}(\mathbb{R})$.

Proof. This is a Lie algebra version of [6, Lemma 7], and essentially the same proof works. Therefore we here give a rather terse presentation, referring to the proof in [6] for further details. Of course if $r=0$ then $d^{\prime}(\mathfrak{h})=\mathfrak{s l}_{d}(\mathbb{R})$ implies $\mathfrak{h}=\mathfrak{s l}_{d}(\mathbb{R})$ and the lemma is trivial; hence in the following we may assume $r>0$.

Set $L^{\prime}=\left\{X \in \mathrm{M}_{r \times d}(\mathbb{R}):(0, X) \in \mathfrak{h}\right\}$; this is a linear subspace of $\mathrm{M}_{r \times d}(\mathbb{R})$. Using $d \iota(\mathfrak{h})=\mathfrak{s l}_{d}(\mathbb{R})$ it follows that $X A \in L^{\prime}$ for all $X \in L^{\prime}$ and all $A \in \mathfrak{s l}_{d}(\mathbb{R})$; and this in turn implies that $L^{\prime}$ must be of the form $L^{\prime}=L^{d}$ for some linear subspace $L \subset \mathbb{R}^{r}$. Let $L^{\perp}$ be the orthogonal complement of $L$ in $\mathbb{R}^{r} ;$ then $\mathrm{M}_{r \times d}(\mathbb{R})=L^{d} \oplus\left(L^{\perp}\right)^{d}$, and for each $A \in \mathfrak{s l}_{d}(\mathbb{R})$ there exists a unique $Y \in\left(L^{\perp}\right)^{d}$ such that $(A, Y) \in \mathfrak{h}$. This implies that the Lie subalgebra $\mathfrak{s}_{L^{\perp}}(\mathbb{R}) \cap \mathfrak{h}$ is a Levi subalgebra of $\mathfrak{s}_{L^{\perp}}(\mathbb{R})$, and so by Malcev's Theorem [11, Ch. III.9], there exists some $X \in\left(L^{\perp}\right)^{d}$ such that $\mathfrak{s}_{L^{\perp}}(\mathbb{R}) \cap \mathfrak{h}=\left(\operatorname{AdI}_{X}\right)\left(\mathfrak{s l}_{d}(\mathbb{R})\right)$. But $\mathfrak{h}$ is the vector space direct sum of $\mathfrak{s}_{L^{\perp}}(\mathbb{R}) \cap \mathfrak{h}$ and $\left\{(0, U): U \in L^{d}\right\}$; hence in fact $\mathfrak{h}=\mathfrak{s}_{L}^{X}(\mathbb{R})$.
Lemma 8.6. Assume that $H$ is a connected Lie subgroup of $G$ satisfying $\iota(H)=G^{\prime}$. Then there exist linear subspaces $L_{j} \subset \mathbb{R}^{r_{j}}$ and matrices $X_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ such that

$$
\begin{equation*}
H=\mathrm{S}_{L_{1}}^{X_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{N}}^{X_{N}}(\mathbb{R}) \tag{8.17}
\end{equation*}
$$

Proof. Let $\mathfrak{h}$ be the Lie subalgebra of $\mathfrak{g}=\mathfrak{s}_{r_{1}}(\mathbb{R}) \times \cdots \times \mathfrak{s}_{r_{N}}(\mathbb{R})$ corresponding to $H$. It follows from $\iota(H)=G^{\prime}$ that $d \iota(\mathfrak{h})=\mathfrak{g}^{\prime}$. Recall that $p_{j}: G \rightarrow G_{j}=\mathrm{S}_{r_{j}}(\mathbb{R})$ denotes the projection onto the $j$ th factor. It follows from $d \iota(\mathfrak{h})=\mathfrak{g}^{\prime}$ that, for each $j$, the Lie subalgebra $d p_{j}(\mathfrak{h})$ of $\mathfrak{s}_{r_{j}}(\mathbb{R})$ satisfies $d \iota\left(d p_{j}(\mathfrak{h})\right)=\mathfrak{s l}_{d}(\mathbb{R})$, and so by Lemma 8.5 there exist a linear subspace $L_{j} \subset \mathbb{R}^{r_{j}}$ and a matrix $X_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ such that

$$
\begin{equation*}
d p_{j}(\mathfrak{h})=\mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R}) \quad(\forall j \in\{1, \ldots, N\}) . \tag{8.18}
\end{equation*}
$$

This implies:

$$
\begin{equation*}
\mathfrak{h} \subset \mathfrak{s}_{L_{1}}^{X_{1}}(\mathbb{R}) \times \cdots \times \mathfrak{s}_{L_{N}}^{X_{N}}(\mathbb{R}) \tag{8.19}
\end{equation*}
$$

We claim that the two sides of (8.19) are in fact equal. In order to prove this equality, it suffices to prove that $\varphi_{j}\left(\mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R})\right) \subset \mathfrak{h}$ for each $j$, where

$$
\varphi_{j}: \mathfrak{s}_{r_{j}}(\mathbb{R}) \rightarrow \mathfrak{g}
$$

is the Lie group homomorphism mapping $X$ to $(0, \cdots, X, \cdots, 0)$ ( 0 s in all positions except the $j$ th). Let $j$ be fixed, and set

$$
\mathfrak{l}=\left\{Z \in \mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R}): \varphi_{j}(Z) \in \mathfrak{h}\right\} .
$$

Using (8.18) it follows that $\mathfrak{l}$ is an ideal of $\mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R})$. Hence also $d \iota(\mathfrak{l})$ is an ideal of $\mathfrak{s l}_{d}(\mathbb{R})$. But given any two elements $Y, Y^{\prime} \in \mathfrak{s l}_{d}(\mathbb{R})$, it follows from $d \iota(\mathfrak{h})=\mathfrak{g}^{\prime}$ that there exist $Z, Z^{\prime} \in \mathfrak{h}$ such that $d \iota\left(d p_{j}(Z)\right)=Y, d \iota\left(d p_{j}\left(Z^{\prime}\right)\right)=Y^{\prime}$, and $d \iota\left(d p_{i}(Z)\right)=d \iota\left(d p_{i}\left(Z^{\prime}\right)\right)=0$ for all $i \neq j$. Then also $\left[Z, Z^{\prime}\right] \in \mathfrak{h}$, and one computes that $\left[Z, Z^{\prime}\right]=\varphi_{j}\left(\left(\left[Y, Y^{\prime}\right], C\right)\right)$ for some $C \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ (and then in fact $\left(\left[Y, Y^{\prime}\right], C\right) \in \mathfrak{s}_{L_{j}}^{X_{j}}$, because of (8.19)). Hence we conclude that $\left[Y, Y^{\prime}\right] \in d \iota(\mathfrak{l})$, for all $Y, Y^{\prime} \in \mathfrak{s l}_{d}(\mathbb{R})$. Since $\mathfrak{s l}_{d}(\mathbb{R})$ is a simple Lie algebra, it follows that $d \iota(\mathfrak{l})=\mathfrak{s l}_{d}(\mathbb{R})$.

Next fix some $Y \in \mathfrak{s l}_{d}(\mathbb{R})$ which is invertible as a $d \times d$ matrix. Because of $d \iota(\mathfrak{l})=\mathfrak{s l}_{d}(\mathbb{R})$ there is some $C \in L_{j}^{d}$ such that $\varphi_{j}((Y, C)) \in \mathfrak{h}$. Using also (8.18) it follows that for any $C^{\prime} \in L_{j}^{d}$ there exists some $Z \in \mathfrak{h}$ satisfying $p_{j}(Z)=\left(Y, C^{\prime}\right)$. Then $\mathfrak{h}$ also contains the Lie product $\left[\varphi_{j}((Y, C)), Z\right]=\varphi_{j}\left(\left[(Y, C),\left(Y, C^{\prime}\right)\right]\right)=\varphi_{j}\left(\left(0,\left(C-C^{\prime}\right) Y\right)\right)$. Hence $\left(0,\left(C-C^{\prime}\right) Y\right) \in \mathfrak{l}$. Since $C^{\prime}$ is an arbitrary element in $L_{j}^{d}$ and $Y$ is invertible, it follows that $(0, C) \in \mathfrak{l}$ for all $C \in L_{j}^{d}$. Combining this fact with $d \iota(\mathfrak{l})=\mathfrak{s l}_{d}(\mathbb{R})$, we finally conclude that $\mathfrak{l}=\mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R})$. Hence $\varphi\left(\mathfrak{s}_{L_{j}}^{X_{j}}(\mathbb{R})\right) \subset \mathfrak{h}$. We have proved that this holds for all $j$; hence we finally conclude:

$$
\begin{equation*}
\mathfrak{h}=\mathfrak{s}_{L_{1}}^{X_{1}}(\mathbb{R}) \times \cdots \times \mathfrak{s}_{L_{N}}^{X_{N}}(\mathbb{R}) \tag{8.20}
\end{equation*}
$$

and so (8.17) holds.

Lemma 8.7. Let $H$ be as in (8.17), and assume that $\Gamma \cap H$ is a lattice in $H$. Then for each $j, L_{j}$ is a rational subspace of $\mathbb{R}^{r_{j}}$, and there exist $Y_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ such that

$$
\begin{equation*}
H=\mathrm{S}_{L_{1}}^{Y_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{N}}^{Y_{N}}(\mathbb{R}) \tag{8.21}
\end{equation*}
$$

Proof. The assumption implies that $\Gamma_{j} \cap S_{L_{j}}^{X_{j}}(\mathbb{R})$ is a lattice in $S_{L_{j}}^{X_{j}}(\mathbb{R})$, for each $j$. As in [21, (5.58), (5.59)], this implies that $L_{j} \cap \mathbb{Z}^{r_{j}}$ is a lattice in $L_{j}$, i.e. $L_{j}$ is a rational subspace of $\mathbb{R}^{r_{j}}$, and furthermore $X_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+L_{j}^{d}$. Choosing now any $Y_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ such that $X_{j} \in Y_{j}+L_{j}^{d}$, we have $\mathrm{S}_{L_{j}}^{X_{j}}(\mathbb{R})=\mathrm{S}_{L_{j}}^{Y_{j}}(\mathbb{R})$. Carrying this out for each $j$, we obtain (8.21).

The following very basic observation will also be useful for us:
Lemma 8.8. Let $L$ and $L^{\prime}$ be linear subspaces of $\mathbb{R}^{r}$, and $X, X^{\prime} \in \mathrm{M}_{r \times d}(\mathbb{R})$. Then $\mathrm{S}_{L^{\prime}}^{X^{\prime}}(\mathbb{R}) \subset$ $\mathrm{S}_{L}^{X}(\mathbb{R})$ holds if and only if $L^{\prime} \subset L$ and $X^{\prime}-X \in L^{d}$.

Proof. Let $Y:=X^{\prime}-X$. Then $\mathrm{S}_{L^{\prime}}^{X^{\prime}}(\mathbb{R}) \subset \mathrm{S}_{L}^{X}(\mathbb{R})$ holds if and only if $\mathrm{I}_{Y}(M, U) \mathrm{I}_{Y}^{-1} \in \mathrm{~S}_{L}(\mathbb{R})$ for all $(M, U) \in \mathrm{S}_{L^{\prime}}(\mathbb{R})$, that is, $U+Y(M-I) \in L^{d}$ for all $M \in \mathrm{SL}_{d}(\mathbb{R})$ and all $U \in L^{\prime d}$. Assume that this holds. Then, considering first only $M=I$ it follows that $L^{\prime} \subset L$, thus $L^{\prime d} \subset L^{d}$, and using this fact it follows that we must have $Y(M-I) \in L^{d}$ for all $M \in \mathrm{SL}_{d}(\mathbb{R})$. Considering only the first column of $M-I$ we conclude that $Y \boldsymbol{a} \in L$ for all column vectors $\boldsymbol{a} \in \mathbb{R}^{d} \backslash\left\{-\boldsymbol{e}_{1}\right\}$, and this in turn implies $Y \in L^{d}$, i.e. $X^{\prime}-X \in L^{d}$. The converse direction is immediate.

In the next lemma we derive an important consequence of the condition " $\left[\exists \eta>0: \forall m \in \mathbb{Z}^{+}\right.$: $\left.V_{m} \notin \Delta_{k}^{(\eta)}\right] "$ in Definition 8.1. Given any $j \in\{1, \ldots, N\}$ with $r_{j} \neq 0$ and $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \backslash\{\mathbf{0}\}$, we set:

$$
\begin{equation*}
K_{j, \boldsymbol{m}}:=\left\{\mathrm{I}_{B} A \mathrm{I}_{Y}: B \in\|\boldsymbol{m}\|^{-2} \mathrm{M}_{r_{j} \times d}(\mathbb{Z}), A \in \mathrm{SL}_{d}(\mathbb{R}), Y \in \boldsymbol{m}^{\perp} \times\left(\mathbb{R}^{r_{j}}\right)^{d-1}\right\} \tag{8.22}
\end{equation*}
$$

Note that $K_{j, \boldsymbol{m}}$ is left $\Gamma_{j}$-invariant $F_{\sigma}$ set in $\mathrm{S}_{r_{j}}(\mathbb{R})$; hence $p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)$ is a left $\Gamma$-invariant $F_{\sigma}$ set in $G$, and $\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)$ is an $F_{\sigma}$ set in $\mathbb{X}$.
Lemma 8.9. For any $k \in \mathbb{Z}^{+}, \nu \in P_{k}, j \in\{1, \ldots, N\}$ with $r_{j} \neq 0$ and $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \backslash\{\mathbf{0}\}$, if $\|\boldsymbol{m}\|^{2} \leq k$ then $\nu\left(\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)\right)=0$.
Proof. (This is similar to [7, pp. 114-115] and [6, Lemma 9].) Let $k, \nu, j, \boldsymbol{m}$ be given as in the statement of the lemma. Set $q:=\|\boldsymbol{m}\|^{2}$ (thus $0<q \leq k$ ). Let us fix a vector $\boldsymbol{b} \in \mathbb{Z}^{r_{j}}$ satisfying $\boldsymbol{b} \cdot \boldsymbol{m}=\operatorname{gcd}\left(m_{1}, \ldots, m_{r_{j}}\right)$. Then we have

$$
\begin{equation*}
\mathbb{Z}^{r_{j}}=\left(\boldsymbol{m}^{\perp} \cap \mathbb{Z}^{r_{j}}\right) \oplus \mathbb{Z} \boldsymbol{b} \tag{8.23}
\end{equation*}
$$

Let $p_{\boldsymbol{b}}: \mathbb{R}^{r_{j}} \rightarrow \mathbb{R}$ be the linear map such that $\boldsymbol{v}-p_{\boldsymbol{b}}(\boldsymbol{v}) \boldsymbol{b} \in \boldsymbol{m}^{\perp}$ for all $\boldsymbol{v} \in \mathbb{R}^{r_{j}}$. For any matrix $Z \in \mathrm{M}_{r_{j} \times e}(\mathbb{R})$, we write $Z_{1}, \ldots, Z_{e} \in \mathbb{R}^{r_{j}}$ for its column vectors (in order), and we define $p_{\boldsymbol{b}}(Z)=\left(p_{\boldsymbol{b}}\left(Z_{1}\right), \ldots, p_{\boldsymbol{b}}\left(Z_{e}\right)\right)$, i.e. $p_{\boldsymbol{b}}(Z)$ is the vector in $\mathbb{R}^{e}$ obtained by applying $p_{\boldsymbol{b}}$ individually to each column of $Z$. Furthermore, for any matrix $Z \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ we will write $Z^{\prime}=\left(Z_{2}, \ldots, Z_{d}\right)$ for the matrix in $\mathrm{M}_{r_{j} \times(d-1)}(\mathbb{R})$ formed by removing the first column vector from $Z$. For any $T>0$ and $\delta>0$ we now introduce the following subsets of $\mathrm{M}_{r_{j} \times d}(\mathbb{R})$ :

$$
\begin{aligned}
& \Omega_{T}:=\left\{Z \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}): p_{\boldsymbol{b}}\left(Z_{1}\right)=0,\left\|p_{\boldsymbol{b}}\left(Z^{\prime}\right)\right\|<T\right\} \\
& \Omega_{T, \delta}:=\left\{Z \in \mathrm{M}_{r_{j} \times d}(\mathbb{R}):\left|p_{\boldsymbol{b}}\left(Z_{1}\right)\right|<\delta,\left\|p_{\boldsymbol{b}}\left(Z^{\prime}\right)\right\|<T\right\} .
\end{aligned}
$$

We also introduce the following subsets of $\mathrm{S}_{r_{j}}(\mathbb{R})$ :

$$
K_{T}=\left\{\mathrm{I}_{B} A \mathrm{I}_{Y}: B \in q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z}), A \in \mathrm{SL}_{d}(\mathbb{R}), Y \in \Omega_{T}\right\}
$$

and

$$
K_{T, \delta}=\left\{\mathrm{I}_{B} A \mathrm{I}_{Y}: B \in q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z}), A \in \mathrm{SL}_{d}(\mathbb{R}), Y \in \Omega_{T, \delta}\right\}
$$

Note that both $K_{T}$ and $K_{T, \delta}$ are left $\Gamma_{j}$-invariant; also $K_{T, \delta}$ is open; hence $p_{j}^{-1}\left(K_{T, \delta}\right)$ is an open subset of $G$, and $\pi\left(p_{j}^{-1}\left(K_{T, \delta}\right)\right)$ is an open subset of $\mathbb{X}$.

Now our goal is to prove that for any given $\eta>0, \lambda \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right), \widetilde{M}=\left(M_{1}, \ldots, M_{N}\right) \in$ $G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, real numbers $\rho_{1}>\rho_{2}>\cdots \rightarrow 0$, and points $V_{1}, V_{2}, \ldots$ in $\widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$, if $\nu_{m} \in P(\mathbb{X})$ is defined as in (8.11), then we have:

$$
\begin{equation*}
\forall T, \varepsilon>0: \exists \delta, m_{0}>0: \forall m \geq m_{0}: \quad \nu_{m}\left(\pi\left(p_{j}^{-1}\left(K_{T, \delta}\right)\right)\right)<\varepsilon \tag{8.24}
\end{equation*}
$$

This will complete the proof of Lemma 8.9, Indeed, if $\nu$ is any weak limit of $\nu_{1}, \nu_{2}, \ldots$, then (8.24) together with the Portmanteau Theorem implies that for any $T, \varepsilon>0$ there exists $\delta>0$ such that $\nu\left(\pi\left(p_{j}^{-1}\left(K_{T, \delta}\right)\right)\right) \leq \varepsilon$. This forces $\nu\left(\pi\left(p_{j}^{-1}\left(K_{T}\right)\right)\right)=0$ for all $T>0$, and so $\nu\left(\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)\right)=0$, i.e. the lemma is proved.

Let us note that it suffices to prove (8.24) for special choices of $\lambda$. Indeed, since $\mathrm{C}_{c}\left(\mathbb{R}^{d-1}\right)$ is dense in $L^{1}\left(\mathbb{R}^{d-1}\right)$, it suffices to prove (8.24) for measures $\lambda \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right)$ of the form $d \lambda(\boldsymbol{x})=\lambda^{\prime}(\boldsymbol{x}) d \boldsymbol{x}$ with $\lambda^{\prime} \in \mathrm{C}_{c}\left(\mathbb{R}^{d-1}\right)$. Next, using the fact that any such function $\lambda^{\prime}$ is bounded, it follows that it actually suffices to prove (8.24) when $\lambda=\left.\operatorname{vol}\right|_{\mathcal{B}_{R}^{d-1}}$, i.e. Lebesgue measure restricted to a ball $\mathcal{B}_{R}^{d-1} 13$ with $R>1$ arbitrary and fixed. Hence from now on we assume that $\lambda$ is of this form.

Let us write $V_{m}=\left(V_{m, 1}, \ldots, V_{m, N}\right)$ with $V_{m, j} \in \mathbb{T}_{j}^{d}$. Then for any $T, \delta, m$,

$$
\begin{equation*}
\nu_{m}\left(\pi\left(p_{j}^{-1}\left(K_{T, \delta}\right)\right)\right)=\operatorname{vol}\left(\left\{\boldsymbol{u} \in \mathcal{B}_{R}^{d-1}: x\left(V_{m, j}\right) M_{j} n_{-}(\boldsymbol{u}) D_{\rho_{m}} \in \pi\left(K_{T, \delta}\right)\right\}\right) \tag{8.25}
\end{equation*}
$$

Take $U_{m, j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with $\pi\left(U_{m, j}\right)=V_{m, j}$. Since $K_{T, \delta}$ is $\Gamma_{j}$-invariant, the condition $x\left(V_{m, j}\right) M_{j} n_{-}(\boldsymbol{u}) D_{\rho_{m}} \in \pi\left(K_{T, \delta}\right)$ is equivalent with $\mathrm{I}_{U_{m, j}} M_{j} n_{-}(\boldsymbol{u}) D_{\rho_{m}} \in K_{T, \delta}$, which in turn is equivalent with

$$
\begin{equation*}
\left(U_{m, j}-q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})\right) M_{j} n_{-}(\boldsymbol{u}) D_{\rho_{m}} \bigcap \Omega_{T, \delta} \neq \emptyset \tag{8.26}
\end{equation*}
$$

But for any $Z \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ the condition $Z n_{-}(\boldsymbol{u}) D_{\rho_{m}} \in \Omega_{T, \delta}$ holds if and only if the vector $\boldsymbol{z}:=$ $p_{\boldsymbol{b}}(Z)$ satisfies $\left|z_{1}\right|<\delta \rho_{m}^{1-d}$ and $\left\|z_{1} \boldsymbol{u}+\left(z_{2}, \ldots, z_{d}\right)\right\|<T \rho_{m}$. We also have $p_{\boldsymbol{b}}\left(A M_{j}\right)=p_{\boldsymbol{b}}(A) M_{j}$ for all $A \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$; furthermore $p_{\boldsymbol{b}}\left(\mathbb{Z}^{r_{j}}\right)=\mathbb{Z}$, by (8.23), and thus $p_{\boldsymbol{b}}\left(\mathrm{M}_{r_{j} \times d}(\mathbb{Z})\right)=\mathbb{Z}^{d}$; therefore

$$
\begin{equation*}
\left\{p_{\boldsymbol{b}}(Z): Z \in\left(U_{m, j}-q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})\right) M_{j}\right\}=\left(p_{\boldsymbol{b}}\left(U_{m, j}\right)+q^{-1} \mathbb{Z}^{d}\right) M_{j}=: L_{m} \tag{8.27}
\end{equation*}
$$

Note that this set $L_{m}$ is a grid in $\mathbb{R}^{d}$. It now follows that the measure in (8.25) equals

$$
\begin{align*}
\operatorname{vol}\left(\left\{\boldsymbol{u} \in \mathcal{B}_{R}^{d-1}:\left[\exists \boldsymbol{z} \in L_{m}:\left|z_{1}\right|<\delta \rho_{m}^{1-d} \text { and }\left\|z_{1} \boldsymbol{u}+\left(z_{2}, \ldots, z_{d}\right)\right\|<T \rho_{m}\right]\right\}\right) \\
\leq \sum_{\substack{\boldsymbol{z} \in L_{m} \\
\left(\left|z_{1}\right|<\delta \rho_{m}^{1-d}\right)}} \operatorname{vol}\left(\left\{\boldsymbol{u} \in \mathcal{B}_{R}^{d-1}:\left\|z_{1} \boldsymbol{u}+\left(z_{2}, \ldots, z_{d}\right)\right\|<T \rho_{m}\right\}\right) \tag{8.28}
\end{align*}
$$

Recall that we are assuming $V_{m} \notin \Delta_{k}^{(\eta)}$ for all $m$; this implies that for all $q^{\prime} \in\{1, \ldots, k\}$ and all $\boldsymbol{n} \in \mathbb{Z}^{r_{j}}$ with $0<\|\boldsymbol{n}\| \leq k$, the point $V_{m, j}$ in $\mathbb{T}_{j}^{d}$ has distance $\geq \eta$ from the set $\Delta_{j, q^{\prime}, \boldsymbol{n}}$, and therefore $U_{m, j}$ has distance $\geq \eta$ from the set $\left(q^{\prime-1} \mathbb{Z}^{r_{j}}+\boldsymbol{n}^{\perp}\right)^{d}$ in in $\mathrm{M}_{r_{j} \times d}(\mathbb{R})=\left(\mathbb{R}^{r_{j}}\right)^{d}$. In particular, $U_{m, j}$ has distance $\geq \eta$ from $\left(q^{-1} \mathbb{Z}^{r_{j}}+\boldsymbol{m}^{\perp}\right)^{d}$, and this is seen to be equivalent to the statement that $\left\|p_{\boldsymbol{b}}\left(U_{m, j}\right)-q^{-1} \boldsymbol{a}\right\| \geq \frac{\|\boldsymbol{m}\|}{\operatorname{gcd}(\boldsymbol{m})} \eta$ for all $\boldsymbol{a} \in \mathbb{Z}^{r_{j}}$. Combining this with the definition of $L_{m}$ in (8.27), it follows that there exists a constant $\eta^{\prime}>0$, independent of $m$, such that

$$
\begin{equation*}
\forall m \in \mathbb{Z}^{+}: \forall \boldsymbol{z} \in L_{m}:\|\boldsymbol{z}\| \geq \eta^{\prime} \tag{8.29}
\end{equation*}
$$

Now let $T>0$ be given, and keep $m \in \mathbb{Z}^{+}$so large that $T \rho_{m}<\eta^{\prime} / 6$. Consider any vector $\boldsymbol{z} \in L_{m}$ which gives a non-zero contribution in the sum in (8.28). This means that there exists some $\boldsymbol{u} \in \mathcal{B}_{R}^{d-1}$ such that $\left\|z_{1} \boldsymbol{u}+\left(z_{2}, \ldots, z_{d}\right)\right\|<T \rho_{m}<\eta^{\prime} / 6$, and thus $\left\|\left(z_{2}, \ldots, z_{d}\right)\right\|<$ $\eta^{\prime} / 6+R\left|z_{1}\right|$. If $R\left|z_{1}\right| \leq \eta^{\prime} / 3$ then it would follow that $\|\boldsymbol{z}\|<\eta^{\prime}$, which is impossible by (8.29).

[^9]Hence we have proved that every $\boldsymbol{z} \in L_{m}$ which gives a non-zero contribution in the sum in (8.28) satisfies

$$
\begin{equation*}
\left|z_{1}\right|>\frac{\eta^{\prime}}{3 R} \quad \text { and } \quad\left\|\left(z_{2}, \ldots, z_{d}\right)\right\|<2 R\left|z_{1}\right| . \tag{8.30}
\end{equation*}
$$

Let $L_{m, 0}$ be the set of all $\boldsymbol{z} \in L_{m}$ satisfying (8.30) and $\left|z_{1}\right| \leq 2$, and for each $\ell \in \mathbb{Z}^{+}$let $L_{m, \ell}$ be the set of all $\boldsymbol{z} \in L_{m}$ satisfying (8.30) and $2^{\ell}<\left|z_{1}\right| \leq 2^{\ell+1}$. Then the sum in (8.28) is

$$
\leq \sum_{0 \leq \ell<\log _{2}\left(\delta \rho_{m}^{1-d}\right)} \sum_{z \in L_{m, \ell}} \operatorname{vol}\left(\mathcal{B}_{1}^{d-1}\right) \cdot\left(\frac{T \rho_{m}}{\left|z_{1}\right|}\right)^{d-1}
$$

But for every $\ell \geq 0$ and every $\boldsymbol{z} \in L_{m}$ we have

$$
\|\boldsymbol{z}\| \leq\left|z_{1}\right|+\left\|\left(z_{2}, \ldots, z_{d}\right)\right\|<(1+2 R)\left|z_{1}\right|<2^{\ell+3} R .
$$

Recall also that each grid $L_{m}$ is a translate of the fixed lattice $q^{-1} \mathbb{Z}^{d} M_{j}$. It follows that there exists a constant $C>0$ which is independent of $m$ and $\ell$ (but which depends on $R, q$ and $M_{j}$ ) such that $\# L_{m, \ell}<C 2^{d \ell}$ for all $\ell \geq 0$. It follows that our sum is

$$
\leq C \cdot \operatorname{vol}\left(\mathcal{B}_{1}^{d-1}\right) \cdot\left(T \rho_{m}\right)^{d-1} \cdot\left(\left(\frac{\eta^{\prime}}{3 R}\right)^{1-d}+\sum_{1 \leq \ell<\log _{2}\left(\delta \rho_{m}^{1-d}\right)} 2^{\ell}\right)<C^{\prime}\left(\rho_{m}^{d-1}+\delta\right)
$$

where $C^{\prime}$ is a constant which is independent of $m$ or $\delta$ (but which depends on $R, T, \eta^{\prime}$ ).
To sum up, we have proved that for any $T, \delta>0$, and all $m \in \mathbb{Z}^{+}$so large that $T \rho_{m}<\eta^{\prime} / 6$, we have $\nu_{m}\left(\pi\left(p_{j}^{-1}\left(K_{T, \delta}\right)\right)\right) \leq C^{\prime}\left(\rho_{m}^{d-1}+\delta\right)$, with a constant $C^{\prime}$ which may depend on $T$, but is independent of $m$ and $\delta$. This bound implies that (8.24) holds, and hence Lemma 8.9 is proved.

Lemma 8.10. Let $j \in\{1, \ldots, N\}$, let $L$ be a rational subspace of $\mathbb{R}^{r_{j}}, L \neq \mathbb{R}^{r_{j}}$ (thus $r_{j} \neq 0$ ), and let $X \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$. Then for any $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \cap L^{\perp} \backslash\{\mathbf{0}\}$ satisfying $X^{\top} \boldsymbol{m} \in \mathbb{Z}^{d}$, and any $Y \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ satisfying $n_{-}\left(\mathbb{R}^{d-1}\right) \subset \mathrm{S}_{L}^{X-Y}(\mathbb{R})$, we have $\Gamma_{j} \mathrm{~S}_{L}^{X}(\mathbb{R}) \mathrm{I}_{Y} \subset K_{j, m}$.

Proof. Set $\widetilde{X}=\|\boldsymbol{m}\|^{-2} \boldsymbol{m}\left(X^{\top} \boldsymbol{m}\right)^{\top} \in\|\boldsymbol{m}\|^{-2} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$, and note that $\widetilde{X}^{\top} \boldsymbol{m}=X^{\top} \boldsymbol{m}$. Let $Y \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ be such that $n_{-}\left(\mathbb{R}^{d-1}\right) \subset \mathrm{S}_{L}^{X-Y}(\mathbb{R})$. This implies that $(Y-X)\left(n_{-}(\boldsymbol{w})-I\right) \subset L^{d}$ for all $\boldsymbol{w} \in \mathbb{R}^{d-1}$, which forces the first column of $Y-X$ lies in $L$.

Now consider an arbitrary element in $\Gamma_{j} S_{L}^{X}(\mathbb{R}) I_{Y}$. This element can be expressed as follows, for some $(\gamma, B) \in \Gamma_{j}$ and $(A, U) \in \mathrm{S}_{L}(\mathbb{R})$ :

$$
\begin{equation*}
(\gamma, B) \mathrm{I}_{X}(A, U) \mathrm{I}_{-X} \mathrm{I}_{Y}=(\gamma A,(B+X) A+U+Y-X)=\mathrm{I}_{(B+\tilde{X}) \gamma^{-1}} \gamma A \mathrm{I}_{W} \tag{8.31}
\end{equation*}
$$

where $W:=(X-\widetilde{X}) A+U+Y-X$. Here $(B+\widetilde{X}) \gamma^{-1} \in\|\boldsymbol{m}\|^{-2} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$. Furthermore, $\widetilde{X}^{\top} \boldsymbol{m}=X^{\top} \boldsymbol{m}$ implies $X-\widetilde{X} \in\left(\boldsymbol{m}^{\perp}\right)^{d}$ and so $(X-\widetilde{X}) A \in\left(\boldsymbol{m}^{\perp}\right)^{d}$. Also $U \in L^{d} \subset\left(\boldsymbol{m}^{\perp}\right)^{d}$, and finally the first column of $Y-X$ lies in $L$, hence in $\boldsymbol{m}^{\perp}$. It follows that the first column of $W$ lies in $\boldsymbol{m}^{\perp}$. Hence the element in (8.31) lies in $K_{j, \boldsymbol{m}}$, and we have proved that $\Gamma_{j} \mathrm{~S}_{L}^{X}(\mathbb{R}) \mathrm{I}_{Y} \subset$ $K_{j, m}$.
Lemma 8.11. Given any $f \in \mathbb{C}_{b}(\mathbb{X})$ and $\varepsilon>0$, there exists $k \in \mathbb{Z}^{+}$such that $|\nu(f)-\mu(f)|<\varepsilon$ holds for all $\nu \in P_{k}$.
Proof. Assume the opposite; this means (since $\left.P_{1} \supset P_{2} \supset \cdots\right)$ that there exist some $f \in \mathrm{C}_{b}(\mathbb{X})$, $\varepsilon>0$, and measures $\nu_{k} \in P_{k}$ for $k=1,2, \ldots$, such that $\left|\nu_{k}(f)-\mu(f)\right|>\varepsilon$ for all $k$.

For each $k$, since $\nu_{k}$ is $W$-invariant by Lemma 8.2, we can apply ergodic decomposition to $\nu_{k}$ : Let $\mathcal{E}$ be the set of ergodic $W$-invariant probability measures on $\mathbb{X}$, provided with its usual Borel $\sigma$-algebra; then there exists a unique Borel probability measure $\omega_{k}$ on $\mathcal{E}$ such that

$$
\begin{equation*}
\nu_{k}=\int_{\mathcal{E}} \alpha d \omega_{k}(\alpha) . \tag{8.32}
\end{equation*}
$$

Cf., e.g., [31, Theorem 4.4]. Note that (8.32) together with Lemma 8.3 implies $\widetilde{\mu}=\widetilde{\iota}_{*} \nu_{k}=$ $\int_{\mathcal{E}} \widetilde{\iota}_{*} \alpha d \omega_{k}(\alpha)$, and for each $\alpha \in \mathcal{E}, \widetilde{\iota}_{*} \alpha$ is an ergodic $W$-invariant measure on $\mathbb{X}^{\prime}$. Hence in fact $\widetilde{\iota}_{*} \alpha=\widetilde{\mu}$ for $\omega_{k}$-almost all $\alpha \in \mathcal{E}$, by uniqueness of the ergodic decomposition of $\widetilde{\mu}$. Furthermore, for every $j \in\{1, \ldots, N\}$ with $r_{j} \neq 0$ and every $\boldsymbol{m} \in \mathbb{Z}^{r_{j}}$ with $0<\|\boldsymbol{m}\| \leq$ $k$, we have $\int_{\mathcal{E}} \alpha\left(\pi\left(p_{j}^{-1}\left(K_{j, m}\right)\right)\right) d \omega_{k}(\alpha)=\nu_{k}\left(\pi\left(p_{j}^{-1}\left(K_{j, m}\right)\right)\right)=0$, by Lemma 8.9, and hence $\alpha\left(\pi\left(p_{j}^{-1}\left(K_{j, m}\right)\right)\right)=0$ for $\omega_{k}$-almost all $\alpha$. Note also that

$$
\int_{\mathcal{E}}|\alpha(f)-\mu(f)| d \omega_{k}(\alpha) \geq\left|\int_{\mathcal{E}} \alpha(f) d P_{\nu_{k}}(\alpha)-\mu(f)\right|=\left|\nu_{k}(f)-\mu(f)\right|>\varepsilon,
$$

and hence the set $\{\alpha \in \mathcal{E}:|\alpha(f)-\mu(f)|>\varepsilon\}$ must have positive measure with respect to $\omega_{k}$.
It follows from the above discussion that for each $k \in \mathbb{Z}^{+}$, there must exist some $\alpha_{k} \in \mathcal{E}$ satisfying $\widetilde{\iota}_{*} \alpha_{k}=\widetilde{\mu}$ and $\left|\alpha_{k}(f)-\mu(f)\right|>\varepsilon$, and

$$
\begin{equation*}
\forall j \in\{1, \ldots, N\}: r_{j} \neq 0 \Rightarrow \quad \forall \boldsymbol{m} \in \mathbb{Z}^{r_{j}} \backslash\{\mathbf{0}\}:\|\boldsymbol{m}\|^{2} \leq k \Rightarrow \alpha_{k}\left(\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)\right)=0 . \tag{8.33}
\end{equation*}
$$

We will assume that such measures $\alpha_{1}, \alpha_{2}, \ldots$ have now been fixed, and we will derive a contradiction.

Clearly $\mathcal{E} \subset \mathcal{Q}(\mathbb{X})$, and hence for each $k$, we can apply Ratner's theorem, [25, Thm. 1], to $\alpha_{k}$. As discussed below (8.12), this implies that there exists some element $g_{k} \in G$ such that, writing $H_{k}:=H_{\alpha_{k}}, \Gamma \cap g_{k} H_{k} g_{k}^{-1}$ is a lattice in $g_{k} H_{k} g_{k}^{-1}$ and the support of $\alpha_{k}$ equals $\Gamma \backslash \Gamma g_{k} H_{k}$. The validity of the previous statements remain intact when replacing $g_{k}$ by any other element from the double coset $\Gamma g_{k} H_{k}$. Hence, since $\iota\left(H_{k}\right)=G^{\prime}$ by Lemma 8.4, after right multiplying $g_{k}$ by an appropriate element in $H_{k}$ we may assume that $g_{k}$ is of the form $g_{k}=\left(\mathrm{I}_{Y_{k, 1}}, \cdots, \mathrm{I}_{Y_{k, N}}\right)$ for some matrices $Y_{k, j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$. Next, by left multiplying $g_{k}$ by an appropriate element in $\Gamma$ (in fact in $\Gamma \cap \iota^{-1}(\{I\})$ ), we may furthermore assume that every entry of every matrix $Y_{k, j}$ lies in the interval $[0,1)$.

Note that the sequence $\alpha_{1}, \alpha_{2}, \ldots$ in $P(\mathbb{X})$ is tight, since $\widetilde{\iota}_{*} \alpha_{k}=\widetilde{\mu}$ for all $k$ and the map $\widetilde{\iota}$ : $\mathbb{X} \rightarrow \mathbb{X}^{\prime}$ is proper. Hence by Prohorov's Theorem, there exists a subsequence, say $\alpha_{k_{1}}, \alpha_{k_{2}}, \ldots$ where $1 \leq k_{1}<k_{2}<\cdots$, which converges to some limit measure $\nu \in P(\mathbb{X})$. In view of our assumption on the entries of the matrices $Y_{k, j}$, we may also assume that $g_{k_{\ell}}$ converges to some element $\tilde{g}$ in $G$ as $\ell \rightarrow \infty$. We have $\widetilde{\iota_{*}} \nu=\widetilde{\mu}$, by the continuous mapping theorem. Furthermore, by [23, Cor. 1.1] we have $\nu \in \mathcal{Q}(\mathbb{X})$, and hence by Ratner's [25, Thm. 1], there exists some $g_{\nu} \in G$ such that $\Gamma \cap g_{\nu} H_{\nu} g_{\nu}^{-1}$ is a lattice in $g_{\nu} H_{\nu} g_{\nu}^{-1}$ and $\operatorname{supp}(\nu)=\Gamma \backslash \Gamma g_{\nu} H_{\nu}$. Note also that $\iota\left(H_{\nu}\right)=G^{\prime}$, by Lemma 8.4. Next we will apply [23, Thm. 1.1(2)]. As a preparation, note that by [23, Lemma 2.3], for each $\ell$ we can find a one-parameter subgroup $\left\{u_{\ell}(t)\right\}_{t \in \mathbb{R}}$ of $W$ which acts ergodically with respect to $\alpha_{k_{\ell}}$, and we can then find an element $s_{\ell} \in G$ such that the trajectory $\left\{\Gamma g_{\nu} s_{\ell} u_{\ell}(t): t>0\right\}$ is uniformly distributed with respect to $\alpha_{k_{\ell}}$. Note that this last property remains valid if $s_{\ell}$ is replaced by $s_{\ell} u_{\ell}(t)$ for any $t>0$, and in this way we may modify the elements $s_{1}, s_{2}, \ldots$ so that $s_{\ell} \rightarrow e$ in $G$ as $\ell \rightarrow \infty$ (this is possible since $\alpha_{k_{\ell}} \rightarrow \nu$ in $P(\mathbb{X})$ and since the point $\Gamma g_{\nu}$ lies in the support of $\nu$ ). Hence by [23, Thm. 1.1(2)], for all sufficiently large $\ell$ we have $\operatorname{supp}\left(\alpha_{k_{\ell}}\right) \subset \operatorname{supp}(\nu) \cdot s_{\ell}$, or equivalently, $\Gamma \backslash \Gamma g_{k_{\ell}} H_{k_{\ell}} \subset \Gamma \backslash \Gamma g_{\nu} H_{\nu} s_{\ell}$, viz., $\Gamma\left(g_{k_{\ell}} H_{k_{\ell}} g_{k_{\ell}}^{-1}\right) g_{k_{\ell}} s_{\ell}^{-1} g_{\nu}^{-1} \subset \Gamma\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)$. But we know that $\Gamma\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)$ is a closed regular submanifold of $G$, and for each $\gamma \in \Gamma, \gamma\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)$ is a connected component of this submanifold (cf. [24, Theorem 1.13]). Hence for each large $\ell$ there exists some $\gamma_{\ell} \in \Gamma$ such that

$$
\begin{equation*}
\left(g_{k_{\ell}} H_{k_{\ell}} g_{k_{\ell}}^{-1}\right) g_{k_{\ell}} s_{\ell}^{-1} g_{\nu}^{-1} \subset \gamma_{\ell}\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right) \tag{8.34}
\end{equation*}
$$

Recall that we also have $g_{k_{\ell}} \rightarrow \tilde{g}$ as $\ell \rightarrow \infty$; hence $g_{k_{\ell}} s_{\ell}^{-1} g_{\nu}^{-1} \rightarrow \tilde{g} g_{\nu}^{-1}$, and since $g_{k_{\ell}} s_{\ell}^{-1} g_{\nu}^{-1} \in$ $\gamma_{\ell}\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)$ for all large $\ell$ it follows that there is some $\widetilde{\gamma} \in \Gamma$ such that $\gamma_{\ell}\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)=$ $\widetilde{\gamma}\left(g_{\nu} H_{\nu} g_{\nu}^{-1}\right)$ for all sufficiently large $\ell$. For these $\ell$, (8.34) implies

$$
\begin{equation*}
g_{k_{\ell}} s_{\ell}^{-1} g_{\nu}^{-1} \widetilde{\gamma}^{-1} \in \widetilde{\gamma} g_{\nu} H_{\nu} g_{\nu}^{-1} \widetilde{\gamma}^{-1} \quad \text { and } \quad g_{k_{\ell}} H_{k_{\ell}} g_{k_{\ell}}^{-1} \subset \widetilde{\gamma} g_{\nu} H_{\nu} g_{\nu}^{-1} \widetilde{\gamma}^{-1} \tag{8.35}
\end{equation*}
$$

We have proved that (8.35) holds for all sufficiently large $\ell$; however by removing the initial elements from the sequence $k_{1}<k_{2}<\cdots$, we may from now on assume that (8.35) holds for all $\ell \in \mathbb{Z}^{+}$.

Next let us apply Lemmas 8.6 and 8.7 to the group $\widetilde{\gamma} g_{\nu} H_{\nu} g_{\nu}^{-1} \widetilde{\gamma}^{-1}$ and the groups $g_{k_{\ell}} H_{k_{\ell}} g_{k_{\ell}}^{-1}$ for all $\ell$. This gives that there exist rational subspaces $L_{j}$ and $L_{\ell, j}$ of $\mathbb{R}^{r_{j}}$ and matrices $X_{j}$ and $X_{\ell, j}$ in $\mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ (for all $j \in\{1, \ldots, N\}$ and $\ell \in \mathbb{Z}^{+}$), such that

$$
\begin{equation*}
\widetilde{\gamma} g_{\nu} H_{\nu} g_{\nu}^{-1} \widetilde{\gamma}^{-1}=\mathrm{S}_{L_{1}}^{X_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{N}}^{X_{N}}(\mathbb{R}) \tag{8.36}
\end{equation*}
$$

and

$$
\begin{equation*}
g_{k_{\ell}} H_{k_{\ell}} g_{k_{\ell}}^{-1}=\mathrm{S}_{L_{\ell, 1}}^{X_{\ell, 1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{\ell, N}}^{X_{\ell, N}}(\mathbb{R}) \quad\left(\forall \ell \in \mathbb{Z}^{+}\right) \tag{8.37}
\end{equation*}
$$

It now follows from (8.35) that $\mathrm{S}_{L_{\ell, j}}^{X_{\ell, j}}(\mathbb{R}) \subset \mathrm{S}_{L_{j}}^{X_{j}}(\mathbb{R})$, and hence by Lemma 8.8,

$$
\begin{equation*}
L_{\ell, j} \subset L_{j} \quad \text { and } \quad X_{\ell, j}-X_{j} \in L_{j}^{d}, \quad \forall \ell \in \mathbb{Z}^{+}, j \in\{1, \ldots, N\} \tag{8.38}
\end{equation*}
$$

Let us also note that (8.37) and $g_{k}=\left(\mathrm{I}_{Y_{k, 1}}, \cdots, \mathrm{I}_{Y_{k, N}}\right)$ imply that

$$
\begin{equation*}
H_{k_{\ell}}=\mathrm{S}_{L_{\ell, 1}}^{X_{\ell, 1}-Y_{k_{\ell}, 1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{\ell, N}}^{X_{\ell, N}-Y_{k_{\ell}, N}}(\mathbb{R}) \tag{8.39}
\end{equation*}
$$

Now we obtain a contradiction as follows: We have $\nu \neq \mu$, since $\nu(f)=\lim _{\ell \rightarrow \infty} \alpha_{k_{\ell}}(f)$ and $\left|\alpha_{k}(f)-\mu(f)\right|>\varepsilon$ for all $k$. Hence $H_{\nu} \neq G$, and so by (8.36) there is some $j \in\{1, \ldots, N\}$ such that $L_{j} \neq \mathbb{R}^{r_{j}}$ (this implies in particular $r_{j}>0$ ). Hence we can choose some $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \cap L_{j}^{\perp} \backslash\{\mathbf{0}\}$ with $X_{j}^{\top} \boldsymbol{m} \in \mathbb{Z}^{d}$. Now by (8.38) we also have $\boldsymbol{m} \perp L_{\ell, j}$ and $X_{\ell, j}^{\top} \boldsymbol{m}=X_{j}^{\top} \boldsymbol{m} \in \mathbb{Z}^{d}$ for all $\ell$. Let us apply this for some fixed choice of $\ell$ so large that $k_{\ell} \geq\|\boldsymbol{m}\|^{2}$. It is immediate from the definition of $H_{k_{\ell}}=H_{\alpha_{k_{\ell}}}\left(\right.$ cf. (8.12) ) that $W \subset H_{k_{\ell}}$; thus $n_{-}\left(\mathbb{R}^{d-1}\right) \subset \mathrm{S}_{L_{\ell, j}}^{X_{\ell, j}-Y_{k_{\ell}, j}}(\mathbb{R})$ (cf. (8.39)), and so by Lemma 8.10, the set $\Gamma_{j} \mathrm{I}_{Y_{\ell, j}} \mathrm{~S}_{L_{\ell, j}}^{X_{\ell, j}-Y_{k_{\ell}, j}}(\mathbb{R})$ is contained in $K_{j, m}$. By (8.39), this implies that $\Gamma \backslash \Gamma g_{k_{\ell}} H_{k_{\ell}}$, i.e. the support of $\alpha_{k_{\ell}}$, is contained in $\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)$, and so $\alpha_{k_{\ell}}\left(\pi\left(p_{j}^{-1}\left(K_{j, \boldsymbol{m}}\right)\right)\right)=1$. This is a contradiction against (8.33), since $\|\boldsymbol{m}\|^{2} \leq k_{\ell}$.

Hence the lemma is proved.
Proof of Theorem 8.1. Given $f \in \mathrm{C}_{b}(\mathbb{X})$ and $\varepsilon>0$, we choose $k \in \mathbb{Z}^{+}$as in Lemma 8.11. Now also let arbitrary $\lambda \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right), \eta>0, \widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$ be given. Assume that there does not exist any $\rho_{0} \in(0,1)$ such that (8.10) holds for all $\rho \in\left(0, \rho_{0}\right)$ and $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$. This means that there exist sequences $\rho_{1}>\rho_{2}>\cdots \rightarrow 0$ and $V_{1}, V_{2}, \ldots$ in $\widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$ satisfying

$$
\begin{equation*}
\left|\int_{\mathbb{R}^{d-1}} f\left(x\left(V_{m}\right) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right) d \lambda(\boldsymbol{u})-\int_{\mathbb{X}} f d \mu\right| \geq \varepsilon, \quad \forall m \in \mathbb{Z}^{+} \tag{8.40}
\end{equation*}
$$

Define $\nu_{m} \in P(\mathbb{X})$ through $\nu_{m}(g)=\int_{\mathbb{R}^{d-1}} g\left(x\left(V_{m}\right) \widetilde{M} \varphi\left(n_{-}(\boldsymbol{u}) D_{\rho_{m}}\right)\right) d \lambda(\boldsymbol{u})$ for all $g \in \mathrm{C}_{b}(\mathbb{X})$ (just as in (8.11)). By [20, Thm. 5] (applied in the same way as in the proof of Lemma8.3) we have $\widetilde{\iota}_{*} \nu_{m} \rightarrow \widetilde{\mu}$ in $P\left(\mathbb{X}^{\prime}\right)$. Hence the sequence $\nu_{1}, \nu_{2}, \ldots$ in $P(\mathbb{X})$ is tight, and so by Prohorov's Theorem, after passing to a subsequence we may assume that $\nu_{1}, \nu_{2}, \ldots$ converges to some $\nu \in P(\mathbb{X})$. Then $\nu \in P_{k}$ (cf. Def. 8.1), and thus by our choice of $k$ we have $|\nu(f)-\mu(f)|<\varepsilon$. But $\nu_{m}$ converges weakly to $\nu$; in particular $\nu_{m}(f) \rightarrow \nu(f)$ as $m \rightarrow \infty$, and hence we conclude that $\left|\nu_{m}(f)-\mu(f)\right|<\varepsilon$ for all sufficiently large $m$. This is a contradiction against (8.40), and thus Theorem8.1 is proved.

Next we establish a variant of Theorem 8.1, where the unipotent element $n_{-}(\boldsymbol{u})$ is replaced by a rotation:
Theorem 8.12. Let $f \in \mathrm{C}_{b}(\mathbb{X})$ and $\varepsilon>0$ be given. Then there exists some $k \in \mathbb{Z}^{+}$such that for every $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\left|\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X}} f d \mu\right|<\varepsilon \tag{8.41}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$.
Proof. This follows from Theorem 8.1 via a fairly standard approximation argument; cf., e.g., [16, Thm. 5.3 and Cor. 5.4]. Some care is needed to ensure that we can obtain a uniform statement as in the theorem.

To start with, we restrict to functions $f$ of compact support. Thus let $f \in \mathrm{C}_{c}(\mathbb{X})$ and $\varepsilon>0$ be given. Fix a corresponding positive integer $k$ as in Theorem 8.1. Now also let $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$ be given.

For $\delta>0$ we set $\mathcal{U}_{\delta}=\left\{T \in \mathrm{SL}_{d}(\mathbb{R}):\|T-I\|<\delta\right\}$, where $\|\cdot\|$ denotes the entrywise maximum norm on $d \times d$ matrices. Thus $\mathcal{U}_{\delta}$ is an open neighborhood of the identity in $\mathrm{SL}_{d}(\mathbb{R})$. Since $f$ has compact support, we can fix $0<\delta<1$ so small that

$$
\begin{equation*}
|f(x \varphi(T))-f(x)|<\varepsilon, \quad \forall x \in \mathbb{X}, T \in \mathcal{U}_{\delta} \tag{8.42}
\end{equation*}
$$

Recall that $R$ is continuous when restricted to $\mathrm{S}_{1}^{d-1}$ minus one point; it follows that there exists a compact subset $S \subset \mathrm{~S}_{1}^{d-1}$ such that the restriction of $R$ to $S$ is continuous, and

$$
\begin{equation*}
\|f\|_{\infty} \cdot \lambda\left(\mathrm{S}_{1}^{d-1} \backslash S\right) \leq \varepsilon \tag{8.43}
\end{equation*}
$$

For each $\boldsymbol{v}_{0} \in S$ we set $\Omega_{\boldsymbol{v}_{0}}=\left\{\boldsymbol{v} \in S: R\left(\boldsymbol{v}_{0}\right)^{-1} R(\boldsymbol{v}) \in \mathcal{U}_{\delta / 2}\right\}$; this is a relatively open neighborhood of $\boldsymbol{v}_{0}$ in $S$. Since $S$ is compact, we can fix a finite subset $\mathfrak{Q}_{0} \subset \mathrm{~S}_{1}^{d-1}$ such that the sets $\Omega_{\boldsymbol{v}_{0}}$ for $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}$ cover $S$. Let us fix an arbitrary total order $\prec$ on $\mathfrak{Q}_{0}$, and set $\Omega_{v_{0}}^{\prime}:=\Omega_{\boldsymbol{v}_{0}} \backslash\binom{\cup_{v_{0}^{\prime} \in \mathfrak{Q}_{0}}^{v_{0}<\boldsymbol{v}_{0}} ⿺}{\Omega_{v_{0}^{\prime}}}$. Then the sets $\Omega_{\boldsymbol{v}_{0}}^{\prime}$ for $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}$ form a partition of $S$. Set $\mathfrak{Q}_{0}^{\prime}=\left\{\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}: \lambda\left(\Omega_{\boldsymbol{v}_{0}}^{\prime}\right)>0\right\}$, and for each $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}$ let $\lambda_{\boldsymbol{v}_{0}}:=\left.\lambda\left(\Omega_{\boldsymbol{v}_{0}}^{\prime}\right)^{-1} \lambda\right|_{\Omega_{v_{0}}} \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$. Note that we now have

$$
\begin{equation*}
\left.\lambda\right|_{S}=\sum_{\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}} \lambda\left(\Omega_{\boldsymbol{v}_{0}}^{\prime}\right) \lambda_{\boldsymbol{v}_{0}} \tag{8.44}
\end{equation*}
$$

Let us fix $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}$ temporarily, and consider the functions $E: \mathrm{S}_{1}^{d-1} \rightarrow \mathrm{M}_{d}(\mathbb{R})$, $a: \mathrm{S}_{1}^{d-1} \rightarrow \mathbb{R}$, $\boldsymbol{b}, \boldsymbol{c}: \mathrm{S}_{1}^{d-1} \rightarrow \mathbb{R}^{d-1}, D: \mathrm{S}_{1}^{d-1} \rightarrow \mathrm{M}_{d-1}(\mathbb{R})$ defined by

$$
E(\boldsymbol{v})=\left(\begin{array}{cc}
a(\boldsymbol{v}) & \boldsymbol{b}(\boldsymbol{v}) \\
\boldsymbol{c}(\boldsymbol{v})^{\top} & D(\boldsymbol{v})
\end{array}\right):=R\left(\boldsymbol{v}_{0}\right)^{-1} R(\boldsymbol{v}) \quad\left(\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}\right)
$$

Note that $(a(\boldsymbol{v}), \boldsymbol{c}(\boldsymbol{v}))=\boldsymbol{e}_{1} E(\boldsymbol{v})^{\top}=\boldsymbol{v} R\left(\boldsymbol{v}_{0}\right)$ for all $\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}$; in particular $a(\boldsymbol{v})=\boldsymbol{v} \cdot \boldsymbol{v}_{0}$, and, since $\delta<1$, it follows that $\Omega_{\boldsymbol{v}_{0}}$ is contained in the open disc $\mathcal{H}_{\boldsymbol{v}_{0}}=\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \boldsymbol{v} \cdot \boldsymbol{v}_{0}>\frac{1}{2}\right\}$. We introduce the function $\boldsymbol{x}: \mathcal{H}_{\boldsymbol{v}_{0}} \rightarrow \mathbb{R}^{d-1}, \boldsymbol{x}(\boldsymbol{v})=-a(\boldsymbol{v})^{-1} \boldsymbol{c}(\boldsymbol{v})$; this is a diffeomorphism of $\mathcal{H}_{\boldsymbol{v}_{0}}$ onto the open ball $\mathcal{B}_{\sqrt{3}}^{d-1}$. We set $\widetilde{\lambda}_{\boldsymbol{v}_{0}}=\boldsymbol{x}_{*}\left(\lambda_{\boldsymbol{v}_{0}}\right) \in P_{\mathrm{ac}}\left(\mathbb{R}^{d-1}\right)$.

Note that $\widetilde{M} \notin \mathfrak{D}_{\mathcal{S}}$ implies that $\widetilde{M} \varphi\left(R\left(\boldsymbol{v}_{0}\right)\right) \notin \mathfrak{D}_{\mathcal{S}}$ for every $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}$. Hence by Theorem 8.1 and our choice of $k$, there exists $\rho_{0} \in(0,1)$ such that for all $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}, \rho \in\left(0, \rho_{0}\right)$ and $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$, we have:

$$
\begin{equation*}
\left|\int_{\mathbb{R}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R\left(\boldsymbol{v}_{0}\right)\right) \varphi\left(n_{-}(\boldsymbol{x}) D_{\rho}\right)\right) d \widetilde{\lambda}_{\boldsymbol{v}_{0}}(\boldsymbol{x})-\int_{\mathbb{X}} f d \mu\right|<\varepsilon \tag{8.45}
\end{equation*}
$$

Here in the integral over $\mathbb{R}^{d-1}$, we substitute $\boldsymbol{x}=\boldsymbol{x}(\boldsymbol{v})$ and then use (8.42) combined with the fact that for every $\boldsymbol{v} \in \Omega_{\boldsymbol{v}_{0}}^{\prime}$ and $\rho \in\left(0, \rho_{0}\right)$ we have $\left(\begin{array}{c}a(\boldsymbol{v})^{-1} \\ \rho^{d} \boldsymbol{c}(\boldsymbol{v})^{\top} \\ \hline\end{array}\right)$ the fact that $E(\boldsymbol{v}) \in \mathcal{U}_{\delta / 2}$, once we note that $|a-1|<\delta / 2$ implies $\left.\left|a^{-1}-1\right|<2|1-a|<\delta\right)$. This gives:

$$
\begin{align*}
& \mid \int_{\mathbb{R}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R\left(\boldsymbol{v}_{0}\right) n_{-}(\boldsymbol{x}) D_{\rho}\right)\right) d \widetilde{\lambda}_{\boldsymbol{v}_{0}}(\boldsymbol{x}) \\
& \left.\quad-\int_{\Omega_{v_{0}}} f\left(x(V) \widetilde{M} \varphi\left(R\left(\boldsymbol{v}_{0}\right) n_{-}(\boldsymbol{x}(\boldsymbol{v})) D_{\rho}\left(\begin{array}{cc}
a(\boldsymbol{v})^{-1} & \mathbf{0} \\
\rho^{d} \boldsymbol{c}(\boldsymbol{v})^{\top} & D(\boldsymbol{v})
\end{array}\right)\right)\right) d \lambda_{\boldsymbol{v}_{0}}(\boldsymbol{v}) \right\rvert\,<\varepsilon \tag{8.46}
\end{align*}
$$

Here the last integral can be simplified using $R\left(\boldsymbol{v}_{0}\right) n_{-}(\boldsymbol{x}(\boldsymbol{v})) D_{\rho}\left(\begin{array}{c}a(\boldsymbol{v})^{-1} \\ \rho^{d} \boldsymbol{c}(\boldsymbol{v})^{\top} \\ D(\boldsymbol{v})\end{array}\right)=R(\boldsymbol{v}) D_{\rho}$. Hence, by (8.45) and (8.46), we have for any $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}, \rho \in\left(0, \rho_{0}\right)$ and $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$ :

$$
\begin{equation*}
\left|\int_{\Omega_{\boldsymbol{v}_{0}}^{\prime}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda_{\boldsymbol{v}_{0}}(\boldsymbol{v})-\int_{\mathbb{X}} f d \mu\right|<2 \varepsilon . \tag{8.47}
\end{equation*}
$$

Multiplying this inequality by $\lambda\left(\Omega_{\boldsymbol{v}_{0}}^{\prime}\right)$ and then adding over all $\boldsymbol{v}_{0} \in \mathfrak{Q}_{0}^{\prime}$ and using (8.44) and (8.43), we conclude that (8.41) holds with $4 \varepsilon$ in place of $\varepsilon$, for all $\rho \in\left(0, \rho_{0}\right)$ and $V \in \widetilde{\mathbb{T}} \backslash \Delta_{k}^{(\eta)}$.

Thus the theorem is proved under the extra assumption that $f \in \mathrm{C}_{c}(\mathbb{X})$. Finally, the extension to the case of arbitrary functions $f \in \mathrm{C}_{b}(\mathbb{X})$ is achieved by a completely standard approximation argument.

The following is an immediate corollary of Theorem 8.12,
Corollary 8.13. Let $V$ be an arbitrary, fixed point in $\widetilde{\mathbb{T}} \backslash \cup \widetilde{p}_{j}^{-1}\left(\Delta_{j, q, m}\right)$, where the union is taken over all triples $\langle j, q, \boldsymbol{m}\rangle$ with $j \in\{1, \ldots, N\}, r_{j} \neq 0, q \in \mathbb{Z}^{+}$and $\boldsymbol{m} \in \mathbb{Z}^{r_{j}} \backslash\{\mathbf{0}\}$. Then for any $\widetilde{M} \in G^{\prime}, f \in \mathrm{C}_{b}(\mathbb{X})$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\mathbb{X}} f d \mu \quad \text { as } \rho \rightarrow 0 \tag{8.48}
\end{equation*}
$$

Proof. The assumption on $V$ implies that for any $k \in \mathbb{Z}^{+}$there is some $\eta>0$ such that $V \notin \Delta_{k}^{(\eta)}$. Using this fact, the corollary is an immediate consequence of Theorem 8.12,

## 9. PRoof of [P2] (UNIFORM SPHERICAL EQUIDISTRIBUTION)

We now return to using the same notation as in Sections 3-7 in particular $\mathcal{P}$ is a finite union of grids in $\mathbb{R}^{d}$, and we assume that an admissible presentation of $\mathcal{P}$ has been fixed (cf. (3.5), (3.4)), and corresponding to this presentation we let the homogeneous space $\mathbb{X}=\Gamma \backslash G$ be as defined in Section 3.2, In particular we have $r_{1}, \ldots, r_{N}>0$ and

$$
\Gamma=\mathrm{S}_{r_{1}}(\mathbb{Z}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{Z})
$$

which is a stricter requirement than what was imposed in Section 8 ,
Our goal in this section is to complete the proof of [P2]. Recall that by our initial discussion in Section 7.3, the task which remains is to prove Theorem 7.7, A key tool in our proof will be Theorem 8.12, note that this theorem will in general be applied to a certain homogeneous submanifold of our present homogeneous space $\mathbb{X}$.
9.1. Equidistribution without uniformity. We will start by proving the following nonuniform result, which as we will see fairly easily implies Theorem 4.2, and which will also play an important role in our proof of Theorem 7.7.

As in Section 8, we let the subset $\mathfrak{D}_{\mathcal{S}} \subset G^{\prime}$ be given by (8.5). Recall that $\Omega=\prod_{j=1}^{N} P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$; cf. (6.3). As in (8.3) we set $\widetilde{\mathbb{T}}=\mathbb{T}_{1}^{d} \times \mathbb{T}_{2}^{d} \times \cdots \times \mathbb{T}_{N}^{d}$, and we let $x: \widetilde{\mathbb{T}} \rightarrow \mathbb{X}$ be the natural embedding. Finally, for any $V=\left(V_{1}, \ldots, V_{N}\right) \in \widetilde{\mathbb{T}}$ we define:

$$
\begin{equation*}
\omega^{(V)}:=\left(\omega_{1}^{\left(V_{1}\right)}, \ldots, \omega_{N}^{\left(V_{N}\right)}\right) \in \Omega \tag{9.1}
\end{equation*}
$$

Theorem 9.1. For any $V \in \widetilde{\mathbb{T}}, \widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, $f \in \mathrm{C}_{b}(\mathbb{X})$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ we have

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\mathbb{X}} f \overline{d \omega^{(V)}} \tag{9.2}
\end{equation*}
$$

as $\rho \rightarrow 0$.

Proof. Let $V=\left(V_{1}, \ldots, V_{N}\right) \in \widetilde{\mathbb{T}}$ be given. For each $j \in\{1, \ldots, N\}$ we write $L_{j}:=L_{j}^{\left(V_{j}\right)}$, fix some $X_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ such that $V_{j}-\pi\left(X_{j}\right) \in\left(\mathbb{S}_{j}^{\left(V_{j}\right)}\right)^{\circ d}$, and fix some $\widetilde{V}_{j} \in X_{j}+L_{j}^{d} \subset \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ with $\pi\left(\widetilde{V}_{j}\right)=V_{j}$ (cf. the discussion above (4.11)). Using the notation in (8.15), we then set:

$$
H:=S_{L_{1}}^{X_{1}}(\mathbb{R}) \times \cdots \times S_{L_{N}}^{X_{N}}(\mathbb{R})=S_{L_{1}}^{\tilde{V}_{1}}(\mathbb{R}) \times \cdots \times S_{L_{N}}^{\tilde{V}_{N}}(\mathbb{R})
$$

As in (4.11) and the proof of Lemma 4.1, we have for each $j$ that $\Gamma_{j}$ intersects $S_{L_{j}}^{\tilde{V}_{j}}(\mathbb{R})$ in a lattice, and the orbit $x\left(V_{j}\right) \cdot \mathrm{S}_{L_{j}}(\mathbb{R})=\Gamma_{j} \backslash \Gamma_{j} \mathrm{~S}_{L_{j}}(\mathbb{R}) \mathrm{I}_{\tilde{V}_{j}}$ is a closed embedded submanifold of $\mathbb{X}_{j}$ which carries a unique $S_{L_{j}}(\mathbb{R})$-invariant probability measure; and by Proposition 4.7 this measure equals $\overline{\omega_{j}^{\left(V_{j}\right)}}$. Taking the product over all $j$, and writing $\widetilde{V}:=\left(\tilde{V}_{1}, \ldots, \widetilde{V}_{N}\right) \in$ $\prod_{j=1}^{N} \mathrm{M}_{r_{j} \times d}(\mathbb{R})$, it follows that $\Gamma \backslash \Gamma H \mathrm{I}_{\tilde{V}}$ is a closed embedded submanifold of $\mathbb{X}$ and, using also (9.1) and (6.11), that $\overline{\omega^{(V)}}$ is the unique $\prod_{j=1}^{N} \mathrm{~S}_{L_{j}}(\mathbb{R})$-invariant probability measure on $\Gamma \backslash \Gamma H \mathrm{I}_{\tilde{\mathrm{V}}}$.

Let $\Gamma_{H}:=\Gamma \cap H$, and let $\mu$ be the unique $H$-invariant probability measure on the homogeneous submanifold $\Gamma \backslash \Gamma H=\Gamma_{H} \backslash H$ of $\mathbb{X}$. In order to prove the theorem, we will prove that for any $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}, F \in \mathrm{C}_{b}\left(\Gamma_{H} \backslash H\right)$ and $\lambda \in P_{\text {ac }}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} F\left(\Gamma_{H} \mathrm{I}_{\widetilde{V}} \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \mathrm{I}_{\widetilde{V}}^{-1}\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\Gamma_{H} \backslash H} F d \mu \tag{9.3}
\end{equation*}
$$

as $\rho \rightarrow 0$. (To see that the integral to the left in (9.3) is well-defined, note that $\mathrm{I}_{\tilde{V}} g \mathrm{I}_{\widetilde{\mathrm{V}}}^{-1} \in H$ for all $g \in G^{\prime}$.)

To see that the convergence in (9.3) implies the statement of the theorem, let $\tau: \mathbb{X} \rightarrow \mathbb{X}$ be right multiplication by $\mathrm{I}_{\tilde{V}}$; then $\overline{\omega^{(V)}}=\tau_{*}(\mu)$, and the left side in (9.2) can be expressed as

$$
\int_{\mathrm{S}_{1}^{d-1}}(f \circ \tau)\left(\Gamma \mathrm{I}_{\widetilde{V}} \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \mathrm{I}_{\tilde{V}}^{-1}\right) d \lambda(\boldsymbol{v})
$$

Hence (9.2) follows from (9.3) if we let $F$ be the restriction of $f \circ \tau$ to $\Gamma \backslash \Gamma H$.
We now turn to the proof of (9.3). We will start by fixing an isomorphism from $H$ onto the Lie group

$$
\widetilde{G}:=\mathrm{S}_{s_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{s_{N}}(\mathbb{R}),
$$

where $s_{j}:=\operatorname{dim} L_{j}$. Given any linear bijection $\varphi: L \xrightarrow{\sim} \mathbb{R}^{s}$, where $L$ is a linear subspace of $\mathbb{R}^{r}$ (for some $r \in \mathbb{Z}^{+}$) of dimension $s \geq 0$, we write $\mathrm{S}_{\varphi}$ for the following Lie group isomorphism:

$$
\begin{equation*}
\mathrm{S}_{\varphi}: \mathrm{S}_{L}(\mathbb{R}) \xrightarrow[\rightarrow]{\sim} \mathrm{S}_{s}(\mathbb{R}), \quad \mathrm{S}_{\varphi}((M, U))=\left(M, \varphi^{d}(U)\right), \tag{9.4}
\end{equation*}
$$

where $\varphi^{d}$ is the linear bijection from $L^{d}$ onto $\mathrm{M}_{s \times d}(\mathbb{R})$ given by applying $\varphi$ to each column of the matrix. (If $s=0$ so that $L=\{\mathbf{0}\}$ and $\mathrm{S}_{s}(\mathbb{R})=\mathrm{SL}_{d}(\mathbb{R})$, the definition in (9.4) should of course be interpreted to say $S_{\varphi}((M, 0))=M$.) One verifies immediately that $S_{\varphi}$ is indeed a Lie group isomorphism. Also for any $X \in \mathrm{M}_{r \times d}(\mathbb{R})$, we introduce the following Lie group isomorphism:

$$
\begin{equation*}
\mathrm{S}_{\varphi}^{X}: \mathrm{S}_{L}^{X}(\mathbb{R}) \xrightarrow[\rightarrow]{\sim} \mathrm{S}_{s}(\mathbb{R}), \quad \mathrm{S}_{\varphi}^{X}(g)=\mathrm{S}_{\varphi}\left(\mathrm{I}_{X}^{-1} g \mathrm{I}_{X}\right) . \tag{9.5}
\end{equation*}
$$

Next, for each $j$, we fix, once and for all, a linear bijection $\varphi_{j}: L_{j} \xrightarrow{\sim} \mathbb{R}^{s_{j}}$ with the property that $\varphi_{j}\left(L_{j} \cap \mathbb{Z}^{r_{j}}\right)=\mathbb{Z}^{s_{j}}$; this is possible since $L_{j}$ is a rational subspace of $\mathbb{R}^{r_{j}}$. Finally we let $\Phi$ be the Lie group isomorphism

$$
\Phi:=S_{\varphi_{1}}^{X_{1}} \times \cdots \times S_{\varphi_{N}}^{X_{N}}: \quad H \xrightarrow{\sim} \widetilde{G}
$$

For each $j$, fix a positive integer $q_{j}$ such that $X_{j} \in q_{j}^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$, and let $\Gamma_{j}^{\prime}$ be the principal congruence subgroup of $\mathrm{SL}_{d}(\mathbb{Z})$ of level $q_{j}$ :

$$
\begin{equation*}
\Gamma_{j}^{\prime}:=\left\{M \in \mathrm{SL}_{d}(\mathbb{Z}): M \equiv I \bmod q_{j}\right\} \tag{9.6}
\end{equation*}
$$

Then set

$$
\begin{equation*}
\widetilde{\Gamma}_{j}:=\Gamma_{j}^{\prime} \ltimes \mathrm{M}_{s_{j} \times d}(\mathbb{Z}) \quad(j=1, \ldots, N) \quad \text { and } \quad \widetilde{\Gamma}:=\widetilde{\Gamma}_{1} \times \cdots \times \widetilde{\Gamma}_{N} . \tag{9.7}
\end{equation*}
$$

We now claim that

$$
\begin{equation*}
\widetilde{\Gamma} \subset \Phi\left(\Gamma_{H}\right) . \tag{9.8}
\end{equation*}
$$

To verify this, it suffices to verify that $\widetilde{\Gamma}_{j} \subset \mathrm{~S}_{\varphi_{j}}^{X_{j}}\left(\mathrm{~S}_{r_{j}}(\mathbb{Z}) \cap \mathrm{S}_{L_{j}}^{X_{j}}(\mathbb{R})\right)$ for each $j$. To do so, note that given any $(M, U) \in \widetilde{\Gamma}_{j}$, we have $\left(\mathrm{S}_{\varphi_{j}}^{X_{j}}\right)^{-1}(M, U) \in \mathrm{S}_{L_{j}}^{X_{j}}(\mathbb{R})$ and

$$
\left(\mathrm{S}_{\varphi_{j}^{X_{j}}}\right)^{-1}(M, U)=\left(M,\left(\varphi_{j}^{d}\right)^{-1}(U)+X_{j}(M-I)\right) .
$$

Here $\left(\varphi_{j}^{d}\right)^{-1}(U) \in \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ since $U \in \mathrm{M}_{s_{j} \times d}(\mathbb{Z})$ and $\varphi_{j}^{-1}\left(\mathbb{Z}^{s_{j}}\right)=L_{j} \cap \mathbb{Z}^{r_{j}} ;$ also $X_{j}(M-$ $I) \in \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ since $X_{j} \in q_{j}^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$ and $M \in \Gamma_{j}^{\prime}$; hence $\left(\mathrm{S}_{\varphi_{j}}^{X_{j}}\right)^{-1}(M, U) \in \mathrm{S}_{r_{j}}(\mathbb{Z})$. This completes the proof of (9.8). Note that it follows from (9.8) that we have a well-defined covering map

$$
\begin{equation*}
J: \widetilde{\Gamma} \backslash \widetilde{G} \rightarrow \Gamma_{H} \backslash H, \quad J(\widetilde{\Gamma} g)=\Gamma_{H} \Phi^{-1}(g) \tag{9.9}
\end{equation*}
$$

Next, the result of Corollary 8.13 , applied to the homogeneous space $\widetilde{\Gamma} \backslash \widetilde{G}$, can be stated as follows: Let $\widetilde{\mu}$ be the invariant probability measure on $\widetilde{\Gamma} \backslash \widetilde{G}$. Let $W=\left(W_{1}, \ldots, W_{N}\right)$ be an arbitrary element in $\prod_{j=1}^{N} \mathrm{M}_{s_{j} \times d}(\mathbb{R})$ such that for every $j \in\{1, \ldots, N\}$ and every rational subspace $L^{\prime} \subsetneq \mathbb{R}^{s_{j}}$, we have

$$
\begin{equation*}
W_{j} \notin \mathrm{M}_{s_{j} \times d}(\mathbb{Q})+\left(L^{\prime}\right)^{d} . \tag{9.10}
\end{equation*}
$$

Then for any $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}, F_{1} \in \mathrm{C}_{b}(\widetilde{\Gamma} \backslash \widetilde{G})$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, we have

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} F_{1}\left(\widetilde{\Gamma} \mathrm{I}_{W} \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\widetilde{\Gamma} \backslash \widetilde{G}} F_{1} d \widetilde{\mu} \quad \text { as } \rho \rightarrow 0 . \tag{9.11}
\end{equation*}
$$

(Note that $\mathrm{I}_{W} \in \widetilde{G}$ since $W \in \prod_{j=1}^{N} \mathrm{M}_{s_{j} \times d}(\mathbb{R})$; cf. (3.17).) Starting from (9.11) and applying the continuous mapping theorem with the covering map $J$ (cf. (9.9)), we conclude: For any element $W=\left(W_{1}, \ldots, W_{N}\right)$ in $\prod_{j=1}^{N} L_{j}^{d}$ such that for every $j \in\{1, \ldots, N\}$ and every rational subspace $L^{\prime} \subsetneq \mathbb{R}^{s_{j}}$,

$$
\begin{equation*}
\varphi_{j}^{d}\left(W_{j}\right) \notin \mathrm{M}_{s_{j} \times d}(\mathbb{Q})+\left(L^{\prime}\right)^{d}, \tag{9.12}
\end{equation*}
$$

and for any $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}, F_{2} \in \mathrm{C}_{b}\left(\Gamma_{H} \backslash H\right)$ and $\lambda \in P_{\text {ac }}\left(\mathrm{S}_{1}^{d-1}\right)$,

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} F_{2}\left(\Gamma_{H} \mathrm{I}_{W} \mathrm{I}_{X} \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \mathrm{I}_{X}^{-1}\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\Gamma_{H} \backslash H} F_{2} d \mu \quad \text { as } \rho \rightarrow 0 \tag{9.13}
\end{equation*}
$$

where $X:=\left(X_{1}, \ldots, X_{N}\right) \in \prod_{j=1}^{N} \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$. In the above deduction we used the fact that $J_{*}(\widetilde{\mu})=\mu$, the unique $H$-invariant probability measure on $\Gamma_{H} \backslash H$.

We wish to apply the last convergence relation with $W:=\widetilde{V}-X$, i.e. $W_{j}:=\widetilde{V}_{j}-X_{j}$ for each $j$. This $W$ lies in $\prod_{j=1}^{N} L_{j}^{d}$, and we proceed to verify that also the condition (9.12) holds for every $j \in\{1, \ldots, N\}$ and every rational subspace $L^{\prime} \subsetneq \mathbb{R}^{s_{j}}$. Assume the opposite, i.e. assume that there exists $j \in\{1, \ldots, N\}$ and a rational subspace $L^{\prime} \subsetneq \mathbb{R}^{s_{j}}$ such that $\varphi_{j}^{d}\left(W_{j}\right) \in \mathrm{M}_{s_{j} \times d}(\mathbb{Q})+\left(L^{\prime}\right)^{d}$. Using the fact that $\varphi_{j}^{-1}\left(\mathbb{Q}^{s_{j}}\right) \subset \mathbb{Q}^{r_{j}}\left(\right.$ since $\left.\varphi_{j}^{-1}\left(\mathbb{Z}^{s_{j}}\right)=L_{j} \cap \mathbb{Z}^{r_{j}}\right)$, we conclude $W_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+\left(L^{\prime \prime}\right)^{d}$, where $L^{\prime \prime}:=\varphi_{j}^{-1}\left(L^{\prime}\right)$. Since $W_{j}:=\widetilde{V}_{j}-X_{j}$ and $X_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$, it follows that $\widetilde{V}_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})+\left(L^{\prime \prime}\right)^{d}$, or equivalently

$$
\begin{equation*}
\widetilde{V}_{j, 1}, \ldots, \widetilde{V}_{j, d} \in \mathbb{Q}^{r_{j}}+L^{\prime \prime} \tag{9.14}
\end{equation*}
$$

where $\widetilde{V}_{j, 1}, \ldots, \widetilde{V}_{j, d} \in \mathbb{R}^{r_{j}}$ are the column vectors of $\widetilde{V}_{j}$. But we have $L^{\prime \prime} \subsetneq L_{j}$ and $L^{\prime \prime}$ is a rational subspace of $\mathbb{R}^{r_{j}}$; also $L_{j}=L_{j}^{\left(V_{j}\right)}=\mathfrak{J}\left(\left\{\tilde{V}_{j, 1}, \ldots, \widetilde{V}_{j, d}\right\}\right)$ (cf. (4.6) ), and thus by Lemma 4.3, $L_{j}$ is the smallest rational subspace of $\mathbb{R}^{r_{j}}$ with the property that $\tilde{V}_{j, 1}, \ldots, \tilde{V}_{j, d} \in \mathbb{Q}^{r_{j}}+L_{j}$. This is a contradiction against (9.14). This completes the proof that the condition (9.12) is fulfilled for our choice of $W$.

Note also that $\mathrm{I}_{W} \in H$ since $W \in \prod_{j=1}^{N} L_{j}^{d}$, and $\mathrm{I}_{W} \mathrm{I}_{X}=\mathrm{I}_{\widetilde{V}}$. Hence for any given $F \in$ $\mathrm{C}_{b}\left(\Gamma_{H} \backslash H\right)$, also the function $F_{2}$ defined by $F_{2}\left(\Gamma_{H} h\right)=F\left(\Gamma_{H} h \mathrm{I}_{W}^{-1}\right)$ lies in $\mathrm{C}_{b}\left(\Gamma_{H} \backslash H\right)$; and applying (9.13) to this function $F_{2}$ we conclude that (9.3) holds for the given function $F$.

Remark 9.1. If the point $V=\left(V_{1}, \ldots, V_{N}\right) \in \widetilde{\mathbb{T}}$ satisfies $\widetilde{\mathrm{r}}_{i_{\psi}}\left(V_{j_{\psi}}\right)=\mathbf{0}$ for some $\psi \in \Psi$, then the statement of Theorem 9.1 also holds with " $\mathbb{X} \psi$ " in place of "X", i.e. for any $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ we have

$$
\begin{equation*}
\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v}) \rightarrow \int_{\mathbb{X} \psi} f \overline{d \omega^{(V)}} \tag{9.15}
\end{equation*}
$$

as $\rho \rightarrow 0$.
Proof. By [13, Lemma 4.26], this follows from Theorem 9.1]if we can only verify that $\overline{\omega^{(V)}}\left(\mathbb{X}^{\psi}\right)=$ 1 and $x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \in \mathbb{X}^{\psi}$ for all $\rho$ and $\boldsymbol{v}$. The first of these statements is immediate from (6.2), Lemma 6.2 and Lemma 6.3, For the second statement, note that $\mathbb{X}^{\psi}$ is preserved by right multiplication of any $G^{\prime}$-element; hence it suffices to verify that $x(V) \in \mathbb{X}^{\psi}$. But writing $\psi=(j, i)$, and taking $\widetilde{V}_{j} \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ so that $V_{j}=\pi\left(\widetilde{V}_{j}\right)$, we have $\mathrm{r}_{i}\left(\widetilde{V}_{j}\right) \in \mathbb{Z}^{d}$ since $\widetilde{\mathrm{r}}_{i}\left(V_{j}\right)=\mathbf{0}$. Now $x\left(V_{j}\right)=\Gamma_{j} \mathrm{I}_{\widetilde{V}_{j}}$ and $\mathbb{Z}^{d} \mathrm{a}_{i}\left(\mathrm{I}_{\widetilde{V}_{j}}\right)=\mathbb{Z}^{d}\left(\mathrm{I}, \mathrm{r}_{i}\left(\widetilde{V}_{j}\right)\right)=\mathbb{Z}^{d}$, which is a lattice containing $\mathbf{0}$. Hence, by (6.2), $x(V) \in \mathbb{X}^{\psi}$, and the proof is complete.

Let us note that Theorem 4.2 is an immediate consequence of Theorem 9.1;
Proof of Theorem 4.2. By (3.19), $g_{0}^{(\boldsymbol{q})}:=\mathrm{I}_{U^{(\boldsymbol{q})}} \widetilde{M}$, where $\widetilde{M}=\left(M_{1}, \ldots, M_{N}\right)$ with the $M_{j} \mathrm{~s}$ coming from the fixed presentation of $\mathcal{P}$ in (3.4), (3.5). This $\widetilde{M}$ lies outside $\mathfrak{D}_{\mathcal{S}}$, by (3.7). Hence the left hand side of (9.2) equals the left hand side of (4.3) if we choose $V:=\pi\left(U^{(\boldsymbol{q})}\right)$, i.e. $V=\left(V_{1}, \ldots, V_{N}\right)$ with $V_{j}=\pi\left(U_{j}^{(\boldsymbol{q})}\right)$. With this choice, $\omega_{j}^{(\boldsymbol{q})}=\omega_{j}^{\left(V_{j}\right)}$ holds by definition, and hence by Proposition 4.7 and (4.2), (6.11), (9.1), we have $\omega^{(V)}=\mu^{(\boldsymbol{q})}$, meaning that also the right hand sides of (9.2) and (4.3) agree.
9.2. A first uniform result. We will now prove a uniform equidistribution result, Theorem 9.2 below, which, in combination with the non-uniform result of Theorem 9.1, will play a key role in our proof of Theorem 7.7.

For any $\psi \in \Psi$ and $j \in\{1, \ldots, N\}$ we pick an arbitrary point $\boldsymbol{q} \in \mathcal{L}_{\psi}$, and define

$$
\begin{equation*}
\mathbb{Y}_{j}^{\psi}:=\pi\left(U_{j}^{(\boldsymbol{q})}\right)+\left(\mathbb{S}_{j}^{\psi}\right)^{d} \subset \mathbb{T}_{j}^{d} \tag{9.16}
\end{equation*}
$$

This is a connected component of the group $\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d}$; cf. Lemma 5.8. Note that $\mathbb{Y}_{j}^{\psi}$ is independent of the choice of $\boldsymbol{q}$, since $\pi\left(U_{j}^{(\boldsymbol{q})}\right)-\pi\left(U_{j}^{\left(\boldsymbol{q}^{\prime}\right)}\right) \in\left(\mathbb{S}_{j}^{\psi}\right)^{d}$ for any two $\boldsymbol{q}, \boldsymbol{q}^{\prime} \in \mathcal{L}_{\psi}$, as was noted in the proof of Lemma 5.8. Let us also fix a matrix $X_{j}^{\psi} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ with the property that

$$
\begin{equation*}
\mathbb{Y}_{j}^{\psi}=\pi\left(X_{j}^{\psi}+\left(L_{j}^{\psi}\right)^{d}\right) \tag{9.17}
\end{equation*}
$$

(Proof of existence: Choose any $\boldsymbol{q} \in \mathcal{L}_{\psi}$; then by Lemma 4.3 we can choose $X_{j}^{\psi} \in \mathrm{M}_{r_{j} \times d}(\mathbb{Q})$ so that $U_{j}^{(\boldsymbol{q})}-X_{j}^{\psi} \in\left(L_{j}^{(\boldsymbol{q})}\right)^{d}$; using (9.16) and Lemma 5.2 it then follows that (9.17) holds.) Furthermore, we fix a linear bijection $\varphi_{j}^{\psi}: L_{j}^{\psi} \xrightarrow{\sim} \mathbb{R}^{s}\left(\right.$ with $\left.s=s(\psi, j)=\operatorname{dim} L_{j}^{\psi}\right)$ with the
property that $\varphi_{j}^{\psi}\left(L_{j}^{\psi} \cap \mathbb{Z}^{r_{j}}\right)=\mathbb{Z}^{s}$. These matrices $X_{j}^{\psi}$ and bijections $\varphi_{j}^{\psi}$ will be kept fixed throughout the present section. We also introduce the following map:

$$
\begin{equation*}
\widetilde{\varphi_{j}^{\psi}}: \mathbb{Y}_{j}^{\psi} \rightarrow \mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z}) ; \quad \widetilde{\varphi_{j}^{\psi}}\left(\pi\left(X_{j}^{\psi}+W\right)\right)=\pi^{\prime}\left(\left(\varphi_{j}^{\psi}\right)^{d}(W)\right) \quad\left(W \in\left(L_{j}^{\psi}\right)^{d}\right) \tag{9.18}
\end{equation*}
$$

where $\pi^{\prime}$ is the projection map $\mathrm{M}_{s \times d}(\mathbb{R}) \rightarrow \mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z})$. It follows from the defining property of $\varphi_{j}^{\psi}$ that the $\operatorname{map} \widetilde{\varphi_{j}^{\psi}}$ is well-defined, and that $\varphi_{j}^{\psi}$ is a diffeomorphism from $\mathbb{Y}_{j}^{\psi}$ onto the torus $\mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z})$.

Next, for any $\psi \in \Psi, j \in\{1, \ldots, N\}, k \in \mathbb{Z}^{+}$and $\eta>0$, we set

$$
\Delta_{\psi, j, k}^{(\eta)}:= \begin{cases}\left(\widetilde{\varphi_{j}^{\psi}}\right)^{-1}\left(\Delta_{j, k}^{(\eta)}\right) & \text { if } L_{j}^{\psi} \neq\{\mathbf{0}\}  \tag{9.19}\\ \emptyset & \text { if } L_{j}^{\psi}=\{\mathbf{0}\}\end{cases}
$$

where if $L_{j}^{\psi} \neq\{\mathbf{0}\}$, the set $\Delta_{j, k}^{(\eta)} \subset \mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z})$ is defined as on p . 47, but using $s=\operatorname{dim} L_{j}^{\psi}$ $(>0)$ in the place of $r_{j}$, so that " $\mathbb{T}_{j}^{d}$ " on p .47 becomes $\mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z})$.

Next, for each $\psi \in \Psi$, we define

$$
\mathbb{Y}^{\psi}=\mathbb{Y}_{1}^{\psi} \times \cdots \times \mathbb{Y}_{N}^{\psi} \subset \widetilde{\mathbb{T}}
$$

and for any $k \in \mathbb{Z}^{+}$and $\eta>0$ :

$$
\begin{equation*}
\Delta_{\psi, k}^{(\eta)}:=\left\{V=\left(V_{1}, \ldots, V_{N}\right) \in \mathbb{Y}^{\psi}: V_{j} \in \Delta_{\psi, j, k}^{(\eta)} \text { for some } j\right\} \tag{9.20}
\end{equation*}
$$

Theorem 9.2. Let $\psi \in \Psi, f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\varepsilon>0$ be given. Then there exists some $k \in \mathbb{Z}^{+}$ such that for every $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\left|\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X} \psi} f d \overline{\omega^{\psi}}\right|<\varepsilon \tag{9.21}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$.
To see that the statement of Theorem 9.2 makes sense, note that for every $V \in \mathbb{Y}^{\psi}$ we have $\underline{x(V)} \in \mathbb{X}^{\psi}$ by the following Lemma 9.3; thus also $x(V) g \in \mathbb{X}^{\psi}$ for all $g \in G^{\prime}$; and we also have $\overline{\omega^{\psi}}\left(\mathbb{X}^{\psi}\right)=1$, by Lemma 6.2 and Lemma 6.6.

Lemma 9.3. For any $\psi \in \Psi$ and $V \in \mathbb{Y}^{\psi}$, we have $\widetilde{\mathbf{r}}_{i_{\psi}}\left(V_{j_{\psi}}\right)=\mathbf{0}$ and $x(V) \in \mathbb{X}^{\psi}$.
Proof. Assume $V=\left(V_{1}, \ldots, V_{N}\right) \in \mathbb{Y}^{\psi}$. Write $\psi=(j, i)$, and choose a point $\boldsymbol{q} \in \mathcal{L}_{\psi}$. Take $W \in \mathrm{M}_{r_{j} \times d}(\mathbb{R})$ so that $V_{j}=\pi(W)$. It follows from $V_{j} \in \mathbb{Y}_{j}^{\psi}$ that $W \in U_{j}^{(\boldsymbol{q})}+\left(L_{j}^{\psi}\right)^{d}+\mathrm{M}_{r_{j} \times d}(\mathbb{Z})$. It follows from (3.4) and (3.15) that $\mathrm{r}_{i}\left(U_{j}^{(\boldsymbol{q})}\right) \in \mathbb{Z}^{d}$, and we noted in the proof of Lemma 6.6 that $L_{j}^{\psi} \perp \boldsymbol{e}_{i}$. Hence $\mathrm{r}_{i}(W) \in \mathbb{Z}^{d}$. This shows that $\widetilde{\mathbf{r}}_{i}\left(V_{j}\right)=\mathbf{0}$, and it also implies that the grid $\mathbb{Z}^{d} \mathrm{a}_{i}\left(\mathrm{I}_{W}\right)=\mathbb{Z}^{d}+W$ contains $\mathbf{0}$, i.e. the point $x\left(V_{j}\right)=\Gamma_{j} \mathrm{I}_{W}$ lies in $\mathbb{X}_{j}^{(i)}$. Hence $x(V) \in \mathbb{X}^{\psi}$.
Proof of Theorem 9.2. Let $\psi \in \Psi$ be given. Let us set $X:=\left(X_{1}^{\psi}, \ldots, X_{N}^{\psi}\right)$ and

$$
\begin{equation*}
H:=\mathrm{S}_{L_{1}^{\psi}}^{X_{1}^{\psi}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{N}^{\psi}}^{X_{N}^{\psi}}(\mathbb{R})=\mathrm{I}_{X}\left(\mathrm{~S}_{L_{1}^{\psi}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{L_{N}^{\psi}}(\mathbb{R})\right) \mathrm{I}_{X}^{-1} \tag{9.22}
\end{equation*}
$$

Recall that $\overline{\omega^{\psi}}=\overline{\omega_{1}^{\psi}} \otimes \cdots \otimes \overline{\omega_{N}^{\psi}}$. We claim that $\overline{\omega^{\psi}}$ equals the unique $\mathrm{I}_{X}^{-1} H \mathrm{I}_{X}$-invariant probability measure on $\Gamma \backslash \Gamma H \mathrm{I}_{X}$. To prove this, it suffices to prove that for each fixed $j, \overline{\omega_{j}^{\psi}}$ equals the unique $\mathrm{S}_{L_{j}^{\psi}}(\mathbb{R})$-invariant probability measure on $\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{X_{j}^{\psi}} \mathrm{S}_{L_{j}^{\psi}}(\mathbb{R})$. To this end, fix an arbitrary matrix $W \in X_{j}^{\psi}+\left(L_{j}^{\psi}\right)^{d}$ which is generic in the sense that it lies outside $\mathrm{M}_{r_{j} \times d}(\mathbb{Q})+L^{d}$ for every rational subspace $L \subsetneq L_{j}^{\psi}$, and set $V=\pi(W) \in \mathbb{T}_{j}^{d}$. Then $L_{j}^{(V)}=L_{j}^{\psi}$ by Lemma 4.3, and hence by Proposition 4.7, $\overline{\omega_{j}^{(V)}}$ is the unique $\mathrm{S}_{L_{j}^{\psi}}(\mathbb{R})$-invariant probability
measure on $\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{W} \mathrm{~S}_{L_{j}^{\psi}}(\mathbb{R})=\Gamma_{j} \backslash \Gamma_{j} \mathrm{I}_{X_{j}^{\psi}} \mathrm{S}_{L_{j}^{\psi}}(\mathbb{R})$. Furthermore, we have $\left(\mathbb{S}_{j}^{(V)}\right)^{\circ}=\pi\left(L_{j}^{(V)}\right)=$ $\pi\left(L_{j}^{\psi}\right)$ and thus $V+\left(\mathbb{S}_{j}^{(V)}\right)^{\circ d}=\pi\left(X_{j}^{\psi}+\left(L_{j}^{\psi}\right)^{d}\right)=\mathbb{Y}_{j}^{\psi}=\pi\left(U_{j}^{(\boldsymbol{q})}\right)+\left(\mathbb{S}_{j}^{\psi}\right)^{d}$ for any $\boldsymbol{q} \in \mathcal{L}_{\psi}$. This implies $\mathcal{O}_{j}^{(V)}=\mathcal{O}_{j}^{\psi}$ (cf. (5.16) and (4.7)), and hence $\omega_{j}^{(V)}=\omega_{j}^{\psi}$. This completes the proof of the claim.

Set $\Gamma_{H}:=\Gamma \cap H$, and let $\mu$ be the $H$-invariant probability measure on $\Gamma_{H} \backslash H$.
Next, for each $j$ let us write $s_{j}:=\operatorname{dim} L_{j}^{\psi}$, and recall that we have fixed a linear bijection $\varphi_{j}^{\psi}: L_{j}^{\psi} \xrightarrow{\sim} \mathbb{R}^{s_{j}}$ with the property that $\varphi_{j}^{\psi}\left(L_{j}^{\psi} \cap \mathbb{Z}^{r_{j}}\right)=\mathbb{Z}^{s_{j}}$. Set

$$
\widetilde{G}:=\mathrm{S}_{s_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{s_{N}}(\mathbb{R})
$$

and let $\Phi$ be the Lie group isomorphism

$$
\Phi:=\mathrm{S}_{\varphi_{1}^{\psi}}^{X_{1}^{\psi}} \times \cdots \times \mathrm{S}_{\varphi_{N}^{\psi}}^{X_{N}^{\psi}}: \quad H \xrightarrow{\sim} \widetilde{G}
$$

(using the notation from (9.5)). For each $j$, choose $q_{j} \in \mathbb{Z}^{+}$so that $X_{j}^{\psi} \in q_{j}^{-1} \mathbb{Z}$, and let $\Gamma_{j}^{\prime}$ be the principal congruence subgroup of $\mathrm{SL}_{d}(\mathbb{Z})$ of order $q_{j}\left(\mathrm{cf}\right.$. (9.6) ) ; then define $\widetilde{\Gamma}_{j}$ and $\widetilde{\Gamma}$ as in (9.7). By an argument entirely similar to the discussion in the proof of Theorem 9.1 (leading up to (9.13)), one verifies that Theorem 8.12 applied to the homogeneous space $\widetilde{\Gamma} \backslash \widetilde{G}$, yields the following result: For any $f \in \mathrm{C}_{b}\left(\Gamma_{H} \backslash H\right)$ and $\varepsilon>0$, there exists some $k \in \mathbb{Z}^{+}$such that for every $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\left|\int_{\mathbb{R}^{d-1}} f\left(\Gamma_{H} \mathrm{I}_{W} \mathrm{I}_{X} \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \mathrm{I}_{X}^{-1}\right) d \lambda(\boldsymbol{v})-\int_{\Gamma_{H} \backslash H} f d \mu\right|<\varepsilon \tag{9.23}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $W=\left(W_{1}, \ldots, W_{N}\right) \in \prod_{j=1}^{N}\left(L_{j}^{\psi}\right)^{d}$ satisfying $\pi_{j}^{\prime}\left(\left(\varphi_{j}^{\psi}\right)^{d}\left(W_{j}\right)\right) \notin \Delta_{j, k}^{(\eta)}$ for every $j \in\{1, \ldots, N\}$ with $s_{j}>0$. In the last condition, the set $\Delta_{j, k}^{(\eta)}$ is defined exactly as on p. 47 but using the dimension $s_{j}$ in place of $r_{j}$ (so that " $\mathbb{T}_{j}^{d}$ " on p .47 becomes $\mathrm{M}_{s_{j} \times d}(\mathbb{R} / \mathbb{Z})$ ), and $\pi_{j}^{\prime}$ is the projection $\mathrm{M}_{s_{j} \times d}(\mathbb{R}) \rightarrow \mathrm{M}_{s_{j} \times d}(\mathbb{R} / \mathbb{Z})$.

Finally, let us write $\widetilde{V}:=W+X$ in the previous result; thus $\widetilde{V}_{j}:=W_{j}+X_{j}^{\psi} \in X_{j}^{\psi}+\left(L_{j}^{\psi}\right)^{d}$ for each $j$, and also $\mathrm{I}_{W} \mathrm{I}_{X}=\mathrm{I}_{\widetilde{V}}$ in (9.23). Then in view of the definitions (9.19) and (9.20), the condition on $W$ is equivalent to $\pi(\widetilde{V}) \notin \Delta_{\psi, k}^{(\eta)}$. Hence, by an argument completely similar to the proof that (9.3) suffices to give Theorem 9.1, the result stated around (9.23) implies the statement of Theorem 9.2,

Next let us note that by combining Theorem 9.2 with Theorem 9.1 , we immediately obtain a variant of Theorem 9.2, where the limit measure $\overline{\omega^{\psi}}$ in (9.21) is replaced by $\overline{\omega^{(V)}}$ :

Theorem 9.4. For any $\psi \in \Psi, f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\varepsilon>0$, there exists some $k \in \mathbb{Z}^{+}$such that for every $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and every $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$, there exists some $\rho_{0} \in(0,1)$ such that

$$
\begin{equation*}
\left|\int_{\mathrm{S}_{1}^{d-1}} f\left(x(V) \widetilde{M} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X} \psi} f \overline{d \omega^{(V)}}\right|<\varepsilon \tag{9.24}
\end{equation*}
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$.
Proof. Given $\psi, f, \varepsilon$, take $k$ as in Theorem 9.2. Now also let $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \eta>0$ and $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$ be given. Take $\rho_{0} \in(0,1)$ as in Theorem 9.2, i.e. so that (9.21) holds for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$. Now for any fixed $V \in \mathbb{Y}^{\psi}$ we have $\widetilde{\mathrm{r}}_{i_{\psi}}\left(V_{j_{\psi}}\right)=\mathbf{0}$ by Lemma 9.3, and so by Theorem 9.1 and Remark 9.1 , the convergence in (9.15) holds as $\rho \rightarrow 0$. Combining this fact with (9.21) gives

$$
\begin{equation*}
\left|\int_{\mathbb{X} \psi} f d \overline{\omega^{(V)}}-\int_{\mathbb{X} \psi} f d \overline{\omega^{\psi}}\right| \leq \varepsilon, \quad \forall V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)} \tag{9.25}
\end{equation*}
$$

Combining (9.21) and (9.25), we conclude that (9.24), with $2 \varepsilon$ in place of $\varepsilon$, holds for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$.
9.3. Proof of the uniformity in [P2]; later steps. Note that for any $\psi \in \Psi$ and $\boldsymbol{q} \in \mathcal{L}_{\psi}$ we have $\pi\left(U_{j}^{(\boldsymbol{q})}\right) \in \mathbb{Y}_{j}^{\psi}$ for all $j \in\{1, \ldots, N\}$, by (9.16), and hence $\pi\left(U^{(\boldsymbol{q})}\right) \in \mathbb{Y}^{\psi}$. The following proposition shows that by taking $\eta$ small, we can ensure that the density of points $\boldsymbol{q} \in \mathcal{L}_{\psi}$ for which $\pi\left(U^{(\boldsymbol{q})}\right)$ falls inside the "singular" set $\Delta_{\psi, k}^{(\eta)}$, is small.

Proposition 9.5. For every $\psi \in \Psi$ and $k \in \mathbb{Z}^{+}$we have

$$
\begin{equation*}
\lim _{\eta \rightarrow 0} \limsup _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{T}^{d}: \pi\left(U^{(\boldsymbol{q})}\right) \in \Delta_{\psi, k}^{(\eta)}\right\}}{T^{d}}=0 . \tag{9.26}
\end{equation*}
$$

Proof. It follows from the definitions in (9.19), (9.20) and on p .47 that $\Delta_{\psi, k}^{(\eta)}$ is a union of sets of the form

$$
\begin{equation*}
\Delta_{j, A, L}^{(\eta)}:=\left\{V=\left(V_{1}, \ldots, V_{N}\right) \in \mathbb{Y}^{\psi}: V_{j}-\pi\left(X_{j}^{\psi}\right) \text { is } \eta \text {-near } \pi\left(A+L^{d}\right)\right\}, \tag{9.27}
\end{equation*}
$$

the union being taken over a finite set of triples $\langle j, A, L\rangle$ with $j \in\{1, \ldots, N\}, L_{j}^{\psi} \neq \mathbf{0}$, $A \in\left(L_{j}^{\psi}\right)^{d} \cap M_{r_{j} \times d}(\mathbb{Q})$ and $L$ being a rational subspace of $L_{j}^{\psi}, L \neq L_{j}^{\psi}$. Note that in (9.27), " $\eta$-near" refers to the Riemannian metric on $\mathbb{Y}_{j}^{\psi}$ induced from the standard Euclidean metric on $\mathrm{M}_{s \times d}(\mathbb{R} / \mathbb{Z})$ (with $s=\operatorname{dim} L_{j}^{\psi}$ ) via the diffeomorphism in (9.18).

It follows that it suffices to prove that for any fixed such triple $\langle j, A, L\rangle$,

$$
\begin{equation*}
\lim _{\eta \rightarrow 0} \limsup _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{T}^{d}: \pi\left(U_{j}^{(\boldsymbol{q})}\right) \in \Delta_{j, A, L}^{(\eta)}\right\}}{T^{d}}=0 \tag{9.28}
\end{equation*}
$$

But it follows from the formula for $U_{j}^{(\boldsymbol{q})}$ in (5.4) and Weyl equidistribution that if we let $Z$ be the closed subgroup of $\mathbb{T}_{j}^{d}$ which is the closure of the set $\left\{\pi\left(U_{j}^{(\boldsymbol{q})}-W_{j}^{\psi}\right): \boldsymbol{q} \in \mathcal{L}_{\psi}\right\}$, and if $\nu$ is the Haar measure on $Z$ normalized so that $\nu(Z)=1$, then for any fixed closed subset $C \subset Z$,

$$
\begin{equation*}
\limsup _{T \rightarrow \infty} \frac{\#\left\{\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{T}^{d}: \pi\left(U_{j}^{(\boldsymbol{q})}-W_{j}^{\psi}\right) \in C\right\}}{\bar{n}_{\psi} \operatorname{vol}\left(\mathcal{B}_{T}^{d}\right)} \leq \nu(C) . \tag{9.29}
\end{equation*}
$$

Note that it follows from (9.16) and (9.17) that $\pi\left(U_{j}^{(\boldsymbol{q})}-W_{j}^{\psi}\right) \in \pi\left(X_{j}^{\psi}-W_{j}^{\psi}\right)+\left(\mathbb{S}_{j}^{\psi}\right)^{d}$ for all $\boldsymbol{q} \in$ $\mathcal{L}_{\psi} ;$ hence also $Z \subset \pi\left(X_{j}^{\psi}-W_{j}^{\psi}\right)+\left(\mathbb{S}_{j}^{\psi}\right)^{d}$, and since $0 \in Z$ it follows that $\pi\left(X_{j}^{\psi}-W_{j}^{\psi}\right) \in\left(\mathbb{S}_{j}^{\psi}\right)^{d}$ and $Z \subset\left(\mathbb{S}_{j}^{\psi}\right)^{d}$. Using (9.29) and (9.27), it follows that in order to prove (9.28), it suffices to prove that $\nu\left(C_{\eta}\right) \rightarrow 0$ as $\eta \rightarrow 0$, where $C_{\eta}$ is the closed $\eta$-neighborhood of $\pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)$. But we have $\cap_{k=1}^{\infty} C_{1 / k}=\pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)$; hence in fact it suffices to prove that

$$
\begin{equation*}
\nu\left(\pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)\right)=0 . \tag{9.30}
\end{equation*}
$$

Fix $q \in \mathbb{Z}^{+}$so that $X_{j}^{\psi}+A \in q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})$. Let $Z^{\circ}$ be the identity component of $Z$; then $Z$ is a union of a finite number of $Z^{\circ}$-cosets. Now if $Z^{\prime}$ is any of these cosets, we may argue as follows. Set

$$
\mathcal{L}^{\prime}:=\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \pi\left(U_{j}^{(\boldsymbol{q})}-W_{j}^{\psi}\right) \in Z^{\prime}\right\}
$$

this is a subgrid of $\mathcal{L}_{\psi}$. Hence by Lemma 5.1, $L_{j}^{\psi}=\mathfrak{L}\left(\left\{U_{j, \ell}^{(\boldsymbol{q})}: \boldsymbol{q} \in \mathcal{L}^{\prime}, \ell \in\{1, \ldots, d\}\right\}\right.$, and so by Lemma 4.3, since $L$ is a rational subspace of $L_{j}^{\psi}$ and $L \neq L_{j}^{\psi}$, there exist some $\boldsymbol{q} \in \mathcal{L}^{\prime}$ and
$\ell \in\{1, \ldots, d\}$ for which $U_{j, \ell}^{(\boldsymbol{q})} \notin q^{-1} \mathbb{Z}^{r_{j}}+L$. This implies that $U_{j}^{(\boldsymbol{q})} \notin q^{-1} \mathrm{M}_{r_{j} \times d}(\mathbb{Z})+L^{d}$, and in particular $\pi\left(U_{j}^{(\boldsymbol{q})}\right) \notin \pi\left(X_{j}^{\psi}+A+L^{d}\right)$. But $\pi\left(U_{j}^{(\boldsymbol{q})}-W_{j}^{\psi}\right) \in Z^{\prime}$. Hence we conclude that

$$
\begin{equation*}
Z^{\prime} \not \subset \pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right) \tag{9.31}
\end{equation*}
$$

Now both $Z^{\prime}$ and $\pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)$ are translates of closed connected subtori of $\left(\mathbb{S}_{j}^{\psi}\right)^{d}$; hence (9.31) implies that $Z^{\prime} \cap \pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)$ is either empty or a submanifold of codimension $\geq 1$ of $Z^{\prime}$. Therefore

$$
\begin{equation*}
\nu\left(Z^{\prime} \cap \pi\left(X_{j}^{\psi}-W_{j}^{\psi}+A+L^{d}\right)\right)=0 \tag{9.32}
\end{equation*}
$$

We have proved that (9.32) holds for every $Z^{\circ}$-coset $Z^{\prime}$ in $Z$; hence (9.30) holds, and the lemma is proved.

Next we prove an auxiliary lemma concerning the type of uniform convergence which we require. We define the upper density of a subset $\mathcal{Z} \subset \mathbb{R}^{d}$ to be the number ${ }^{14}$

$$
\begin{equation*}
\limsup _{T \rightarrow \infty} T^{-d} \#\left(\mathcal{Z} \cap \mathcal{B}_{T}^{d}\right) \tag{9.33}
\end{equation*}
$$

Lemma 9.6. Let $\mathcal{Q}$ be a locally finite subset of $\mathbb{R}^{d}$, let $J$ be a countable set, and let a function $F: J \times \mathcal{Q} \times(0,1) \rightarrow \mathbb{R}$ be given. Assume that
(i) For any fixed $j \in J$ and $\boldsymbol{q} \in \mathcal{Q}, F(j, \boldsymbol{q}, \rho) \rightarrow 0$ as $\rho \rightarrow 0$,
and
(ii) for any $j \in J$ and $\varepsilon, \varepsilon^{\prime}>0$, there exist $\rho_{0} \in(0,1)$ and a subset $\mathcal{Z} \subset \mathcal{Q}$ of upper density $\leq \varepsilon^{\prime}$, such that $|F(j, \boldsymbol{q}, \rho)|<\varepsilon$ for all $\rho \in\left(0, \rho_{0}\right)$ and all $\boldsymbol{q} \in \mathcal{Q} \backslash \mathcal{Z}$.
Then for any decreasing function $\mathcal{T}:(0,1) \rightarrow \mathbb{R}^{+}$, there exists a subset $\mathcal{E} \subset \mathcal{Q}$ of density zero, such that for each fixed $j \in J$, we have $F(j, \boldsymbol{q}, \rho) \rightarrow 0$ as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{Q} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$.
Proof. We assume, without loss of generality, that $J=\mathbb{Z}^{+}$. For any $j, k, m \in \mathbb{Z}^{+}$, by (ii) there exist $\rho_{0}^{\prime} \in(0,1)$ and a subset $\mathcal{Z}^{\prime} \subset \mathcal{Q}$ of upper density $\leq 2^{-m-1}$ such that $|F(j, \boldsymbol{q}, \rho)|<2^{-k}$ for all $\rho \in\left(0, \rho_{0}^{\prime}\right)$ and all $\boldsymbol{q} \in \mathcal{Q} \backslash \mathcal{Z}^{\prime}$. We may then choose $T_{0}>0$ so that $T^{-d} \#\left(\mathcal{Z}^{\prime} \cap \mathcal{B}_{T}^{d}\right)<2^{-m}$ for all $T \geq T_{0}$, and set $\mathcal{Z}=\mathcal{Z}(j, k, m):=\mathcal{Z}^{\prime} \backslash \mathcal{B}_{T_{0}}^{d}$. Now $\mathcal{Z}^{\prime} \backslash \mathcal{Z}$ is finite, and for each $\boldsymbol{q} \in \mathcal{Z}^{\prime} \backslash \mathcal{Z}$ there exists some $\rho_{0}^{(\boldsymbol{q})} \in(0,1)$ such that $|F(j, \boldsymbol{q}, \rho)|<2^{-k}$ for all $\rho \in\left(0, \rho_{0}^{(\boldsymbol{q})}\right)$, by (i). Set

$$
\rho_{0}=\rho_{0}(j, k, m):=\min \left(\left\{\rho_{0}^{\prime}\right\} \cup\left\{\rho_{0}^{(\boldsymbol{q})}: \boldsymbol{q} \in \mathcal{Z}^{\prime} \backslash \mathcal{Z}\right\}\right)
$$

Now for any $j, k, m \in \mathbb{Z}^{+}$, we have constructed a number $\rho_{0}(j, k, m) \in(0,1)$ and a subset $\mathcal{Z}(j, k, m) \subset \mathcal{Q}$, and it is clear from our construction that

$$
\begin{equation*}
\#\left(\mathcal{Z}(j, k, m) \cap \mathcal{B}_{T}^{d}\right)<2^{-m} T^{d}, \quad \forall T>0 \tag{9.34}
\end{equation*}
$$

and that

$$
\begin{equation*}
|F(j, \boldsymbol{q}, \rho)|<2^{-k}, \quad \forall \rho \in\left(0, \rho_{0}(j, k, m)\right), \boldsymbol{q} \in \mathcal{Q} \backslash \mathcal{Z}(j, k, m) \tag{9.35}
\end{equation*}
$$

Let us now also set, for any $m \in \mathbb{Z}^{+}$,

$$
\widetilde{\mathcal{Z}}(m):=\bigcup_{j \in \mathbb{Z}^{+}} \bigcup_{k \in \mathbb{Z}^{+}} \mathcal{Z}(j, k, j+k+m)
$$

Then

$$
\begin{equation*}
\#\left(\widetilde{\mathcal{Z}}(m) \cap \mathcal{B}_{T}^{d}\right) \leq \sum_{j, k \in \mathbb{Z}^{+}} 2^{-j-k-m} T^{d}=2^{-m} T^{d}, \quad \forall T>0 \tag{9.36}
\end{equation*}
$$

[^10]Now let a decreasing function $\mathcal{T}:(0,1) \rightarrow \mathbb{R}^{+}$be given, as in the statement of the lemma. For any $m \in \mathbb{Z}^{+}$we set

$$
\begin{equation*}
\widetilde{\rho}_{0}(m):=\min \left\{\rho_{0}(j, k, j+k+m): j, k \in\{1, \ldots, m\}\right\} . \tag{9.37}
\end{equation*}
$$

Next choose numbers $1 \leq R_{1}<R_{2}<\cdots$ so that $R_{m} \geq \mathcal{T}\left(\widetilde{\rho}_{0}(m+1)\right)$ and $R_{m+1} \geq 2 R_{m}$ for all $m$. Finally set

$$
\mathcal{E}:=\bigcup_{m=1}^{\infty}\left(\widetilde{\mathcal{Z}}(m) \cap \mathcal{B}_{R_{m}}^{d}\right) .
$$

Let us prove that this set $\mathcal{E}$ has density zero. Given any $T \geq R_{1}$, choosing $n$ so that $R_{n} \leq$ $T<R_{n+1}$ we have

$$
\#\left(\mathcal{E} \cap \mathcal{B}_{T}^{d}\right) \leq \sum_{m=1}^{n} \#\left(\widetilde{\mathcal{Z}}(m) \cap \mathcal{B}_{R_{m}}^{d}\right)+\sum_{m=n+1}^{\infty} \#\left(\widetilde{\mathcal{Z}}(m) \cap \mathcal{B}_{T}^{d}\right) \leq \sum_{m=1}^{n} 2^{-m} R_{m}^{d}+\sum_{m=n+1}^{\infty} 2^{-m} T^{d}
$$

by (9.36). Here for all $m \leq n$ we have $R_{m} \leq 2^{m-n} R_{n}$ and hence $R_{m}^{d} \leq 2^{2(m-n)} T^{d}$, since $d \geq 2$. Plugging in these bounds we get $\#\left(\mathcal{E} \cap \mathcal{B}_{T}^{d}\right)<3 \cdot 2^{-n} T^{d}$. Hence since $n \rightarrow \infty$ as $T \rightarrow \infty$, we conclude that $\mathcal{E}$ has density zero.

It remains to prove the uniform convergence stated in the lemma. Thus let $j \in \mathbb{Z}^{+}$and $\varepsilon>0$ be given. By (i), we can take $\rho_{1} \in(0,1)$ so small that $|F(j, \boldsymbol{q}, \rho)|<\varepsilon$ for all $\boldsymbol{q} \in \mathcal{Q} \cap \mathcal{B}_{R_{j}}^{d}$ and all $\rho \in\left(0, \rho_{1}\right)$. Choose $k \in \mathbb{Z}^{+}$so that $2^{-k}<\varepsilon$, and set

$$
\begin{equation*}
\rho_{0}:=\min \left(\left\{\rho_{1}\right\} \cup\left\{\rho_{0}(j, k, j+k+m): m \in\{1, \ldots, k\}\right\}\right) . \tag{9.38}
\end{equation*}
$$

Now for any $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{Q} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$ we can argue as follows: If $\|\boldsymbol{q}\|<R_{j}$ then $|F(j, \boldsymbol{q}, \rho)|<\varepsilon$ since $\rho<\rho_{0} \leq \rho_{1}$. Next assume instead that $\|\boldsymbol{q}\| \geq R_{j}$, and let $m>j$ be the minimal positive integer such that $\|\boldsymbol{q}\|<R_{m}$. Then $\mathcal{T}(\rho)>\|\boldsymbol{q}\| \geq R_{m-1} \geq \mathcal{T}\left(\widetilde{\rho}_{0}(m)\right)$, and so $\rho<\widetilde{\rho}_{0}(m)$. If $m \geq k$ then $\widetilde{\rho}_{0}(m) \leq \rho_{0}(j, k, j+k+m)$ by (9.37); on the other hand if $m<k$ then $\rho_{0} \leq \rho_{0}(j, k, j+k+m)$ by (9.38). Hence we always have $\rho<\rho_{0}(j, k, j+k+m)$. Furthermore, $\boldsymbol{q} \notin \mathcal{E}$ and $\|\boldsymbol{q}\|<R_{m}$ implies $\boldsymbol{q} \notin \widetilde{\mathcal{Z}}(m)$, and in particular $\boldsymbol{q} \notin \mathcal{Z}(j, k, j+k+m)$. Hence by (9.35), $|F(j, \boldsymbol{q}, \rho)|<2^{-k}<\varepsilon$. Summing up, we have proved:

$$
\begin{equation*}
|F(j, \boldsymbol{q}, \rho)|<\varepsilon, \quad \forall \rho \in\left(0, \rho_{0}\right), \boldsymbol{q} \in \mathcal{Q} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E} \tag{9.39}
\end{equation*}
$$

This completes the proof of the uniform convergence stated in the lemma.
We are now ready to prove Theorem 7.7
Proof of Theorem [7.7. Let $\psi$ and $\mathcal{T}$ be given as in the statement of Theorem 7.7 .
Let $J_{1}$ be a fixed countable dense subset of $\mathrm{C}_{c}\left(\mathbb{X}^{\psi}\right)$ with respect to the uniform norm. We equip $P_{\text {ac }}\left(\mathrm{S}_{1}^{d-1}\right)$ the metric $d$ defined by $d\left(\lambda_{1}, \lambda_{2}\right):=\int_{\mathrm{S}_{1}^{d-1}}\left|\lambda_{1}^{\prime}-\lambda_{2}^{\prime}\right| d \boldsymbol{v}$, where $\lambda_{j}^{\prime}$ is the density of $\lambda_{j}$ with respect to $\sigma$, that is, $\lambda_{j}^{\prime} \in \mathrm{L}^{1}\left(\mathrm{~S}_{1}^{d-1}\right)$ and $\lambda_{j}=\lambda_{j}^{\prime} d \sigma(j=1,2)$. Now let $J_{2}$ be a fixed countable dense subset of $P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ with respect to the metric $d$; such a set exists since $\mathrm{C}\left(\mathrm{S}_{1}^{d-1}\right)$ is dense in $\mathrm{L}^{1}\left(\mathrm{~S}_{1}^{d-1}\right)$ [26, Thm. 3.14].

We will apply Lemma 9.6 with $\mathcal{Q}=\mathcal{L}_{\psi}, J:=J_{1} \times J_{2}$, and with the function $F: J \times \mathcal{L}_{\psi} \times$ $(0,1) \rightarrow \mathbb{R}$ defined by

$$
\begin{equation*}
F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho):=\int_{\mathrm{S}_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathbb{X}^{\psi}} f \overline{d \omega^{(\boldsymbol{q})}} . \tag{9.40}
\end{equation*}
$$

Recall that $\left.\overline{\omega^{(\boldsymbol{q})}} \mathbb{X}^{\psi}\right)=1$ for every $\boldsymbol{q} \in \mathcal{L}_{\psi}$, by Lemmas 6.4 and 6.2, hence the integral over $\mathbb{X}^{\psi}$ in (9.40) may just as well be taken over all $\mathbb{X}$. Recall also that for any $\boldsymbol{q} \in \mathbb{R}$ we have $\overline{\omega^{(q)}}=\mu^{(q)}$, by Proposition 4.7 and (4.2), (6.4), (6.11). Hence by Theorem 4.2 we have $F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho) \rightarrow 0$ for any fixed $\langle f, \lambda\rangle \in J$ and $\boldsymbol{q} \in \mathcal{\mathcal { L } _ { \psi }}$, i.e. the assumption (i) in Lemma 9.6 holds.

We next verify that also the assumption (ii) in Lemma 9.6 holds. Thus let $\langle f, \lambda\rangle \in J$ and $\varepsilon, \varepsilon^{\prime}>0$ be given. Choose $k \in \mathbb{Z}^{+}$as in Theorem 9.4 (for our given $\left.\psi, f, \varepsilon\right)$. By Proposition 9.5, we can now fix some $\eta>0$ such that the set $\mathcal{Z}:=\left\{\boldsymbol{q} \in \mathcal{L}_{\psi}: \pi\left(U^{(\boldsymbol{q})}\right) \in \Delta_{\psi, k}^{(\eta)}\right\}$ has upper density $<\varepsilon^{\prime}$. We keep $\widetilde{M} \in G^{\prime} \backslash \mathfrak{D}_{\mathcal{S}}$ fixed as in (3.18), for our given, fixed union of grids $\mathcal{P}$. Because of our choice of $k$, we may now fix $\rho_{0} \in(0,1)$ in such a way that (9.24) holds for all $\rho \in\left(0, \rho_{0}\right)$ and all $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$, for our chosen $\widetilde{M}, \psi, f, \lambda, \varepsilon, \eta$. Now for any $\boldsymbol{q} \in \mathcal{L}_{\psi} \backslash \mathcal{Z}$, we may apply (9.24) with $V:=\pi\left(U^{(q)}\right) \in \widetilde{\mathbb{T}}$; indeed, for this $V$ we have $V \in \mathbb{Y}^{\psi} \backslash \Delta_{\psi, k}^{(\eta)}$ since $\boldsymbol{q} \notin \mathcal{Z}$; also $\Gamma g_{0}^{(\boldsymbol{q})}=x(V) \widetilde{M}$ (by (3.19)) and $\omega^{(V)}=\omega^{(\boldsymbol{q})}$. Hence we conclude that for all $\boldsymbol{q} \in \mathcal{L}_{\psi} \backslash \mathcal{Z}$ and all $\rho \in\left(0, \rho_{0}\right)$ we have $|F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho)|<\varepsilon$. Hence assumption (ii) in Lemma 9.6 is indeed fulfilled.

It now follows from Lemma 9.6 that there exists a subset $\mathcal{E} \subset \mathcal{L}_{\psi}$ of density zero such that for any $f \in J_{1}$ and $\lambda \in J_{2}$ we have $F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho) \rightarrow 0$ as $\rho \rightarrow 0$, uniformly over all $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$, viz., the uniform convergence in (7.11) in Theorem 7.7 holds. To complete the proof of Theorem 7.7, we will give a (standard) approximation argument to show that the uniform convergence in (7.11) in fact holds for arbitrary $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$.

Note that equation (9.40) defines $F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho)$ for arbitrary $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$; and for any $f_{1}, f_{2} \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right), \lambda_{1}, \lambda_{2} \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right), \boldsymbol{q} \in \mathcal{L}_{\psi}$ and $\rho \in(0,1)$ we have, with $\|\cdot\|_{u}$ denoting the uniform norm on $\mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$ :

$$
\begin{aligned}
\mid F\left(\left\langle f_{1}, \lambda_{1}\right\rangle, \boldsymbol{q}, \rho\right) & -F\left(\left\langle f_{2}, \lambda_{2}\right\rangle, \boldsymbol{q}, \rho\right) \mid \\
& \leq 2\left\|f_{1}-f_{2}\right\|_{u}+\int_{\mathrm{S}_{1}^{d-1}}\left|f_{1}\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right)\right| \cdot\left|\lambda_{1}^{\prime}(\boldsymbol{v})-\lambda_{2}^{\prime}(\boldsymbol{v})\right| d \sigma(\boldsymbol{v}) \\
& \leq 2\left\|f_{1}-f_{2}\right\|_{u}+\left\|f_{1}\right\|_{u} \cdot d\left(\lambda_{1}, \lambda_{2}\right) .
\end{aligned}
$$

Using this bound, and the fact that $J_{1}$ and $J_{2}$ are dense in $\mathrm{C}_{c}\left(\mathbb{X}^{\psi}\right)$ and $P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$, respectively, it is immediate to extend the uniform convergence in (7.11) (with the subset $\mathcal{E} \subset \mathcal{L}_{\psi}$ fixed once and for all) from $f \in J_{1}$ and $\lambda \in J_{2}$ to arbitrary $f \in \mathrm{C}_{c}\left(\mathbb{X}^{\psi}\right)$ and $\lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$.

It remains to extend to arbitrary functions $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right)$. Thus let $f \in \mathrm{C}_{b}\left(\mathbb{X}^{\psi}\right), \lambda \in P_{\mathrm{ac}}\left(\mathrm{S}_{1}^{d-1}\right)$ and $\varepsilon>0$ be given. Set $B:=\|f\|_{u}$; we may assume $B>0$ since otherwise $f$ is identically zero. Let $\nu$ be the $\mathrm{SL}_{d}(\mathbb{R})$ invariant probability measure on $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$; fix a compact subset $K^{\prime}$ of $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ with $\nu\left(K^{\prime}\right)>(1-\varepsilon /(4 B))^{1 / N}$, and then set $K:=\mathbb{X}^{\psi} \cap \prod_{j=1}^{N} \widetilde{\iota}_{j}^{-1}\left(K^{\prime}\right)$, where $\widetilde{\iota}_{j}$ is the projection from $\mathbb{X}_{j}$ to $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R}){ }^{15}$. Then $K$ is a compact subset of $\mathbb{X}^{\psi}$, and for every $\boldsymbol{q} \in \mathcal{L}_{\psi}$ we have

$$
\begin{equation*}
\overline{\omega^{(\boldsymbol{q})}}(K)=\overline{\omega^{(\boldsymbol{q})}}\left(\prod_{j=1}^{N} \widetilde{\iota}_{j}^{-1}\left(K^{\prime}\right)\right)=\prod_{j=1}^{N} \overline{\omega_{j}^{(\boldsymbol{q})}}\left(\widetilde{\iota}_{j}^{-1}\left(K^{\prime}\right)\right)=\nu\left(K^{\prime}\right)^{N}>1-\frac{\varepsilon}{4 B}, \tag{9.41}
\end{equation*}
$$

where we first used the fact that $\overline{\omega^{(\boldsymbol{q})}}\left(\mathbb{X}^{\psi}\right)=1$ (by Lemmas 6.4 and 6.2), and then used Lemma 4.5. Next fix a function $h \in \mathrm{C}_{c}\left(\mathbb{X}_{\psi}\right)$ with $0 \leq h \leq 1$ and $\left.h\right|_{K} \equiv 1$. By the uniform convergence which we have already proved, there exists $\rho_{0} \in(0,1)$ such that for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$ we have $|F(\langle h, \lambda\rangle, \boldsymbol{q}, \rho)|<\varepsilon /(4 B)$; but also $\int_{\mathbb{X}^{\psi}} h \overline{\omega^{(\boldsymbol{q})}} \geq \overline{\omega^{(\boldsymbol{q})}}(K)>$ $1-\varepsilon /(4 B)$, and hence $\int_{S_{1}^{d-1}} h\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})>1-\varepsilon /(2 B)$. It follows that

$$
\lambda\left(\left\{\boldsymbol{v} \in \mathrm{S}_{1}^{d-1}: \Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \in \operatorname{supp}(h)\right\}\right)>1-\frac{\varepsilon}{2 B},
$$

for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$.

[^11]Next fix a function $h_{1} \in \mathrm{C}_{c}\left(\mathbb{X}^{\psi}\right)$ with $0 \leq h_{1} \leq 1$ and $\left.h_{1}\right|_{\operatorname{supp}(h)} \equiv 1$, and set $f_{1}:=h_{1} f \in$ $\mathrm{C}_{c}\left(\mathbb{X}^{\psi}\right)$. Then note that for all $x \in \mathbb{X}^{\psi}$ we have

$$
\left|f(x)-f_{1}(x)\right| \leq|f(x)| \cdot\left|1-h_{1}(x)\right| \leq B\left|1-h_{1}(x)\right| \leq B \cdot I(x \notin \operatorname{supp}(h)) .
$$

Hence for all $\rho \in\left(0, \rho_{0}\right)$ and $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$ we have

$$
\begin{aligned}
\left|\int_{\mathrm{S}_{1}^{d-1}} f\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})-\int_{\mathrm{S}_{1}^{d-1}} f_{1}\left(\Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right)\right) d \lambda(\boldsymbol{v})\right| \\
\leq B \cdot \lambda\left(\left\{\boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}: \Gamma g_{0}^{(\boldsymbol{q})} \varphi\left(R(\boldsymbol{v}) D_{\rho}\right) \notin \operatorname{supp}(h)\right\}\right)<\frac{\varepsilon}{2},
\end{aligned}
$$

and also

$$
\left|\int_{\mathbb{X}^{\psi}} f d \overline{\omega^{(\boldsymbol{q})}}-\int_{\mathbb{X}^{\psi}} f_{1} \overline{d \omega^{(\boldsymbol{q})}}\right| \leq B \cdot \overline{\omega^{(\boldsymbol{q})}}\left(\mathbb{X}^{\psi} \backslash \operatorname{supp}(h)\right) \leq B \cdot \overline{\omega^{(\boldsymbol{q})}}\left(\mathbb{X}^{\psi} \backslash K\right)<\frac{\varepsilon}{4} .
$$

Hence for these $\rho$ and $\boldsymbol{q}$ we have $\left|F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho)-F\left(\left\langle f_{1}, \lambda\right\rangle, \boldsymbol{q}, \rho\right)\right|<3 \varepsilon / 4$. Furthermore, by again applying the uniform convergence result which we have already proved it follows that after possibly shrinking $\rho_{0}$, we have $\left|F\left(\left\langle f_{1}, \lambda\right\rangle, \boldsymbol{q}, \rho\right)\right|<\varepsilon / 4$ for all $\rho \in\left(0, \rho_{0}\right)$ and all $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$. Combining the last two inequalities, we conclude that $|F(\langle f, \lambda\rangle, \boldsymbol{q}, \rho)|<\varepsilon$ for all $\rho \in\left(0, \rho_{0}\right)$ and all $\boldsymbol{q} \in \mathcal{L}_{\psi} \cap \mathcal{B}_{\mathcal{T}(\rho)}^{d} \backslash \mathcal{E}$. This completes the proof of the uniform convergence in (7.11), viz., the proof of Theorem 7.7.

Note that in view of the results in Section [7, we have now also completed the proof of the main result of the paper, Theorem 2.1.

## 10. The transition kernels

10.1. Definitions; collision kernels and transition kernels. We will now give the explicit formulas for the collision kernels $p^{(\psi)}$ and $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ appearing in Theorem 1.2, It should be noted that Theorem [1.2 is a reformulation of [21, Theorem 4.6], in our special case of $\mathcal{P}$ being a finite union of grids. Crucially, our main result, Theorem [2.1, is required to ensure that the assumptions in [21, Theorem 4.6] are fulfilled. We also use the fact that any integral over the space $\Sigma$ with respect to the measure m can be reduced to an integral over $\widetilde{\Psi}$, since $\mathrm{m}(\Sigma \backslash \widetilde{\Psi})=0$; and in the statement of Theorem 1.2 we view m as a probability measure on $\Psi$, via the bijection $\psi \leftrightarrow \sigma^{\psi}$ between $\Psi$ and $\widetilde{\Psi}$.

The defining formulas for the collision kernels are [21, (3.41) and (3.44)] ${ }^{16}$

$$
\begin{equation*}
p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=\frac{\sigma\left(\boldsymbol{V}^{\prime}, \boldsymbol{V}\right)}{v_{d-1}} k^{\mathbf{g}}(\xi, \boldsymbol{w}, \psi) \tag{10.1}
\end{equation*}
$$

and

$$
\begin{equation*}
p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=\frac{\sigma\left(\boldsymbol{V}^{\prime}, \boldsymbol{V}\right)}{v_{d-1}} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right), \tag{10.2}
\end{equation*}
$$

where $k$ and $k^{\mathrm{g}}$ are the transition kernels, whose definition we will recall below; $\sigma\left(\boldsymbol{V}^{\prime}, \boldsymbol{V}\right)$ is the differential cross section of the fixed scattering process (see [21, Sec. 3.4] and the start of Section 1.1 above); and finally $\boldsymbol{w} \in \mathcal{B}_{1}^{d-1}$ is the impact parameter which is a function of $\boldsymbol{V}^{\prime}$ and $\boldsymbol{V}$, and $\boldsymbol{w}^{\prime} \in \mathcal{B}_{1}^{d-1}$ is the exit parameter which is a function of $\boldsymbol{V}^{\prime \prime}$ and $\boldsymbol{V}^{\prime}$. We refer to [21, Sections 3.4 and 3.5] for a more detailed description; in particular it should be noted that the formula (10.1) only makes sense for $\boldsymbol{V}^{\prime}, \boldsymbol{V}$ which can form entry and exit velocities in a scatterer collision; for any other $\boldsymbol{V}^{\prime}, \boldsymbol{V} \in \mathrm{S}_{1}^{d-1}$ we have $p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=0$ by definition; a similar statement holds for (10.2).

Next we recall in detail the definition of the transition kernels $k$ and $k^{g}$, from [21, Section 3.1], with some parts of the notation mildly modified.

[^12]Recall that for any $\xi>0$ we denote $\mathfrak{Z}_{\xi}=(0, \xi) \times \mathcal{B}_{1}^{d-1}$, an open cylinder in $\mathbb{R}^{d}$. In the following we will also use this notation with $\xi=\infty$, i.e., $\mathfrak{Z}_{\infty}=(0, \infty) \times \mathcal{B}_{1}^{d-1}$. Furthermore, we define the " $\boldsymbol{e}_{1}$-coordinate" of any point $(\boldsymbol{x}, \varsigma) \in \mathcal{X}$ to be the number $\boldsymbol{x} \cdot \boldsymbol{e}_{1}$. Next we define a map

$$
\begin{equation*}
z: N_{s}(\mathcal{X}) \rightarrow \Delta:=\left(\mathfrak{Z}_{\infty} \times \Sigma\right) \sqcup\{\text { undef }\} \tag{10.3}
\end{equation*}
$$

in the following way: ${ }^{17}$ Given $Y \in N_{s}(\mathcal{X})$, let $\boldsymbol{z}=\boldsymbol{z}(Y)$ be that point in $Y \cap\left(\mathfrak{Z}_{\infty} \times \Sigma\right)$ which has minimal $\boldsymbol{e}_{1}$-coordinate; if there does not exist a unique such point then let $\boldsymbol{z}(Y)=$ undef. Here "undef" is a dummy element not in $\mathfrak{Z}_{\infty} \times \Sigma$, and we provide $\Delta$ with the disjoint union topology. We also introduce the reflection map

$$
\begin{equation*}
\iota: \mathcal{X} \rightarrow \mathcal{X}, \quad \iota\left(\left(w_{1}, \boldsymbol{w}\right), \varsigma\right)=\left(\left(w_{1},-\boldsymbol{w}\right), \varsigma\right) \quad\left(w_{1} \in \mathbb{R}, \boldsymbol{w} \in \mathbb{R}^{d-1}, \varsigma \in \Sigma\right) \tag{10.4}
\end{equation*}
$$

Next, given $\boldsymbol{w} \in \mathbb{R}^{d-1}$ and $\psi \in \Psi$, we let $\kappa_{\boldsymbol{w}, \psi} \in P(\Delta)$ be the pushforward of the measure $\mu_{\sigma^{\psi}}$ by the map $Y \mapsto \iota(\boldsymbol{z}(Y-(0, \boldsymbol{w})))$ from $N_{s}(\mathcal{X})$ to $\Delta$. Equivalently,
(10.5) $\quad \kappa_{\boldsymbol{w}, \psi}$ is the pushforward of $\overline{\omega^{\psi}}$ by the map $x \mapsto \iota(\boldsymbol{z}(J(x)-(0, \boldsymbol{w})))$ from $\mathbb{X}$ to $\Delta$.

Indeed, this holds since $\mu_{\sigma^{\psi}}=J_{\psi_{*}} \overline{\omega^{\psi}}$ by definition, and furthermore $\boldsymbol{z}\left(J_{\psi}(x)\right)=\boldsymbol{z}(J(x))$ for all $x \in \mathbb{X}^{\psi}$, and $\overline{\omega^{\psi}}\left(\mathbb{X} \backslash \mathbb{X}^{\psi}\right)=0$ by (6.15). By [21, Lemma 3.1], $\kappa_{\boldsymbol{w}, \psi}(\{$ undef $\})=0$, i.e. $\kappa_{\boldsymbol{w}, \psi}$ restricts to a probability measure on $\mathfrak{Z}_{\infty} \times \Sigma$. Furthermore, since $J(x) \subset \mathbb{R}^{d} \times \widetilde{\Psi}$ for all $x \in \mathbb{X}^{\psi}$, by (6.12), $\kappa_{\boldsymbol{w}, \psi}$ in fact restricts to a probability measure on $\mathfrak{Z}_{\infty} \times \widetilde{\Psi}$. By Lemma 7.4, we have $\kappa_{\boldsymbol{w}, \psi}(B) \leq \bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B)$ for every Borel set $B \subset \mathfrak{Z}_{\infty} \times \Sigma \subset \mathcal{X}$; in particular $\kappa_{\boldsymbol{w}, \psi}$ is absolutely continuous with respect to $\mu_{\mathcal{X}}$.

Now the transition kernel $k$ is defined as follows: For any given $\boldsymbol{w}^{\prime} \in \mathbb{R}^{d-1}$ and $\psi^{\prime} \in \Psi$, we define $k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \cdot\right)$ to be the probability density of the measure $\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}$ with respect to the measure $v_{d-1}^{-1} \mu_{\mathcal{X}}$, but restricted to the set $\mathfrak{Z}_{\infty} \times \widetilde{\Psi}$, parametrized by $(\xi, \boldsymbol{w}, \psi) \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi$. That is, we define the function

$$
\begin{equation*}
k: \mathbb{R}^{d-1} \times \Psi \times \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi \rightarrow\left[0, \bar{n}_{\mathcal{P}} v_{d-1}\right] \tag{10.6}
\end{equation*}
$$

so that for each $\boldsymbol{w}^{\prime} \in \mathbb{R}^{d-1}$ and $\psi^{\prime} \in \Psi, k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \cdot, \cdot, \cdot\right)$ is uniquely defined as an element in $\mathrm{L}^{1}\left(\mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi, d \xi d \boldsymbol{w} d \mathrm{~m}\right)$, and

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=v_{d-1}^{-1} \int_{\left(\xi, \boldsymbol{w}, \sigma^{\psi}\right) \in B} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right) d \xi d \boldsymbol{w} d \mathrm{~m}(\psi) \tag{10.7}
\end{equation*}
$$

for all Borel sets $B \subset \mathfrak{Z}_{\infty} \times \Sigma$. Note that in (10.7) we view $m$ as a probability measure on $\Psi$ through $\mathrm{m}(\psi):=\mathrm{m}\left(\sigma^{\psi}\right)$, just as we did in the statement of Theorem 1.2.

Similarly, we let $\kappa^{\mathrm{g}}$ be the pushforward of the measure $\mu^{\mathrm{g}}$ by the map $Y \mapsto \iota(\boldsymbol{z}(Y))$ from $N_{s}(\mathcal{X})$ to $\Delta$. Then by (7.40) and Lemma 7.15 we have, for any fixed $\boldsymbol{v} \in \mathbb{R}^{d}$ :

$$
\begin{equation*}
\kappa^{\mathrm{g}} \text { is the pushforward of } \overline{\omega^{\mathrm{g}}} \text { by the map } x \mapsto \iota(\boldsymbol{z}(J(x)+\boldsymbol{v})) \text { from } \mathbb{X} \text { to } \Delta . \tag{10.8}
\end{equation*}
$$

By [21, Lemma 3.5], and since $J(x) \subset \mathbb{R}^{d} \times \widetilde{\Psi}$ for all $x \in \mathbb{X}$ by (6.10) and (6.12), $\kappa^{\mathrm{g}}$ in fact restricts to a probability measure on $\mathfrak{Z}_{\infty} \times \widetilde{\Psi}$. By [21, Prop. 2.27] applied with $A=N_{s}(\mathcal{X})$, we have $\kappa^{\mathrm{g}}(B) \leq \bar{n}_{\mathcal{P}} \mu_{\mathcal{X}}(B)$ for every Borel set $B \subset \mathfrak{Z}_{\infty} \times \Sigma \subset \mathcal{X}$ (in our situation this may also be derived from (7.40) and Proposition 4.8, by an argument similar to that in the proof of Lemma (7.4). In particular $\kappa^{\mathrm{g}}$ is absolutely continuous with respect to $\mu_{\mathcal{X}}$. We now define the transition kernel $k^{\mathrm{g}}$ to be the corresponding probability density with respect to the measure $v_{d-1}^{-1} \mu_{\mathcal{X}}$, but restricted to the set $\mathfrak{Z}_{\infty} \times \widetilde{\Psi}$, which we parametrize by $(\xi, \boldsymbol{w}, \psi) \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi$. That is, we define the function

$$
\begin{equation*}
k^{\mathrm{g}}: \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi \rightarrow\left[0, \bar{n}_{\mathcal{P}} v_{d-1}\right] \tag{10.9}
\end{equation*}
$$

[^13]so that
\[

$$
\begin{equation*}
\kappa^{\mathrm{g}}(B)=v_{d-1}^{-1} \int_{\left(\xi, \boldsymbol{w}, \sigma^{\psi}\right) \in B} k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi) d \xi d \boldsymbol{w} d \mathrm{~m}(\psi) \tag{10.10}
\end{equation*}
$$

\]

for all Borel sets $B \subset \mathfrak{Z}_{\infty} \times \Sigma$. The function $k^{\mathrm{g}}$ is uniquely defined as an element in $\mathrm{L}^{1}\left(\mathbb{R}_{>0} \times\right.$ $\left.\mathcal{B}_{1}^{d-1} \times \Psi, d \xi d \boldsymbol{w} d \mathrm{~m}\right)$.

To facilitate comparison, let us note that in the notation of 21] we have:

$$
\begin{equation*}
k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right)=" k\left(\left(\boldsymbol{w}^{\prime}, \sigma^{\psi^{\prime}}\right), \xi,\left(\boldsymbol{w}, \sigma^{\psi}\right)\right)^{\prime} ; \quad k^{\mathbf{g}}(\xi, \boldsymbol{w}, \psi)=" k^{\mathrm{g}}\left(\xi,\left(\boldsymbol{w}, \sigma^{\psi}\right)\right) " \tag{10.11}
\end{equation*}
$$

Let us also recall from [21] the following formula expressing $k^{\mathrm{g}}$ in terms of $k$.
Lemma 10.1. There is a continuous version of $k^{\mathrm{g}}$, which for all $(\xi, \boldsymbol{w}, \psi) \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi$ satisfies

$$
\begin{equation*}
k^{\mathbf{g}}(\xi, \boldsymbol{w}, \psi)=\bar{n}_{\mathcal{P}} \int_{\xi}^{\infty} \int_{\mathcal{B}_{1}^{d-1}} \int_{\Psi} k\left(-\boldsymbol{w}, \psi, \xi^{\prime}, \boldsymbol{w}^{\prime}, \psi^{\prime}\right) d \mathrm{~m}\left(\psi^{\prime}\right) d \boldsymbol{w}^{\prime} d \xi^{\prime} \tag{10.12}
\end{equation*}
$$

Proof. This follows as in the proof of [21, Cor. 3.24]. (Indeed, for $d=2$ we are forced to take $R=-I_{1} \in O(1)$ in that proof, which directly gives the lemma; for $d \geq 3$ we apply the proof with a fixed $R \in O(d-1)$ with $\operatorname{det} R=-1$, and then also use [21, Lemma 3.18].)

In the rest of Section 10 we will always assume that the continuous version of $k^{5}$ is used.
Let us now note that the formula for the free path length density $\Phi_{\mathcal{P}}$, (1.17), can be rewritten as follows using (10.1):

$$
\begin{equation*}
\Phi_{\mathcal{P}}(\xi)=\frac{1}{v_{d-1}} \int_{\mathcal{B}_{1}^{d-1} \times \Psi} k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi) d \boldsymbol{w} d \mathrm{~m}(\psi) \tag{10.13}
\end{equation*}
$$

We also introduce the density function for the free path length between consecutive collisions,

$$
\begin{equation*}
\Phi_{\mathcal{P}}(\xi)=\frac{1}{v_{d-1}^{2}} \int_{\mathcal{B}_{1}^{d-1} \times \Psi} \int_{\mathcal{B}_{1}^{d-1} \times \Psi} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right) d \boldsymbol{w} d \mathrm{~m}(\psi) d \boldsymbol{w}^{\prime} d \mathrm{~m}\left(\psi^{\prime}\right) . \tag{10.14}
\end{equation*}
$$

It is immediate from the definitions of $k^{\mathrm{g}}$ and $k$ that $\int_{0}^{\infty} \Phi_{\mathcal{P}}(\xi) d \xi=\int_{0}^{\infty} \bar{\Phi}_{\mathcal{P}}(\xi) d \xi=1$ and that $0 \leq \Phi_{\mathcal{P}}(\xi), \bar{\Phi}_{\mathcal{P}}(\xi) \leq \bar{n}_{\mathcal{P}} v_{d-1}$. Furthermore $\Phi_{\mathcal{P}}(\xi)$ is continuous since $k^{\mathrm{g}}$ is continuous. The formula (1.20) expressing $\Phi_{\mathcal{P}}(\xi)$ in terms of $\bar{\Phi}_{\mathcal{P}}(\xi)$ follows from Lemma 10.1, and (1.20) implies that $\Phi_{\mathcal{P}}(\xi)$ is a decreasing function. Also, (1.21) follows from $\int_{0}^{\infty} \Phi_{\mathcal{P}}(\xi) d \xi=1$ and (1.20). It should be noted that the formula (1.21) holds more generally for any scattering configuration $\mathcal{P}$ which fits within the framework of [21; see [21, Cor. 3.23].

Remark 10.1. In the special case when $\mathcal{P}$ is a single grid (presented with $\# \Psi=1$ ), the transition kernels $k$ and $k^{\mathrm{g}}$ have previously been studied in detail in the papers [15, 16, 17, 18. In this case we have:

$$
k(\boldsymbol{z}, \psi, \xi, \boldsymbol{w}, \psi)=\bar{n}_{\mathcal{P}} v_{d-1} \Phi_{\mathbf{0}}\left(\bar{n}_{\mathcal{P}} \xi, \boldsymbol{w}, \boldsymbol{z}\right) \quad \text { and } \quad k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)=\bar{n}_{\mathcal{P}} v_{d-1} \Phi\left(\bar{n}_{\mathcal{P}} \xi, \boldsymbol{w}\right)
$$

where " $\Phi_{0}$ " and " $\Phi$ " is the notation used in [15, 18]. Also the free path length densities are given by

$$
\bar{\Phi}_{\mathcal{P}}(\xi)=\bar{n}_{\mathcal{P}} \bar{\Phi}_{\mathbf{0}}\left(\bar{n}_{\mathcal{P}} \xi\right) \quad \text { and } \quad \Phi_{\mathcal{P}}(\xi)=\bar{n}_{\mathcal{P}} \Phi\left(\bar{n}_{\mathcal{P}} \xi\right),
$$

where, again, " $\Phi_{0}$ " and " $\Phi$ " is the notation used in [15, 18].
Remark 10.2. The more general special case when all the grids appearing in $\mathcal{P}$ are incommensurable was previously studied in [20]. In that paper the normalizing condition $\bar{n}_{\mathcal{P}}=1$ was imposed. Given such a point set $\mathcal{P}$, presented with $\Psi=\{(j, 1): j=1, \ldots, N\}$, we have:

$$
\begin{equation*}
k\left(\boldsymbol{w}, \psi^{\prime}, \xi, \boldsymbol{z}, \psi\right)=\frac{v_{d-1}}{\bar{n}_{j_{\psi}}} \cdot \Phi_{\mathbf{0}}^{\left(j_{\psi}^{\prime} \rightarrow j_{\psi}\right)}(\xi, \boldsymbol{z}, \boldsymbol{w}) \quad \text { and } \quad k^{\mathbf{g}}(\xi, \boldsymbol{w}, \psi)=\frac{v_{d-1}}{\bar{n}_{j_{\psi}}} \cdot \Phi^{\left(j_{\psi}\right)}(\xi, \boldsymbol{w}) \tag{10.15}
\end{equation*}
$$

where " $\Phi_{0}^{(j \rightarrow k)}$ " and " $\Phi^{(j) "}$ is the notation used in [20]. The first translation formula in (10.15) can be obtained e.g. by comparing the statements of [20, Theorem 4 and (5.4)] and [21, Thm. 3.6 ], and the second formula by comparing [20, Theorem 4 and (5.11)] and [21, Thm. 3.14]. See also Remark 10.5 below.

Remark 10.3. For $\mathcal{P}$ a general finite union of grids as in the present paper, it seems that it should be possible to give explicit formulas for the transition kernels $k$ and $k^{\mathrm{g}}$ in terms of Haar measures on appropriate homogeneous submanifolds of $\mathbb{X}$, similar in spirit to the formulas in [16, Sec. 7-8] for the case when $\mathcal{P}$ is a single grid. Obtaining such formulas would be a first step towards the problem of generalizing the precise asymptotic estimates for the transition kernels which were obtained in [18] in the single grid case.
10.2. The product formula. In view of our explicit set-up from Section 3.1, any two $\psi, \psi^{\prime} \in$ $\Psi$ are equivalent ( " $\psi \sim \psi^{\prime \prime}$ ") in the sense defined in Section 1.2 if and only if $j_{\psi}=j_{\psi^{\prime}}$. Hence we have an obvious bijection between $C_{\Psi}$ and $\{1, \ldots, N\}$. For any $j \in\{1, \ldots, N\}$, let us denote the corresponding equivalence class by $\Psi_{j}:=\left\{(j, i): i \in\left\{1, \ldots, r_{j}\right\}\right\} \in C_{\Psi}$. Let us also introduce the short-hand notation $\mathcal{P}_{j}:=\mathcal{P}_{\Psi_{j}}$; that is

$$
\begin{equation*}
\mathcal{P}_{j}=\bigcup_{\psi \in \Psi_{j}} \mathcal{L}_{\psi} \tag{10.16}
\end{equation*}
$$

The asymptotic density of this point set is $\bar{n}_{j}:=\bar{n}_{\mathcal{P}_{j}}=\sum_{\psi \in \Psi_{j}} \bar{n}_{\psi}$.
Lemma 10.2. For each $j \in\{1, \ldots, N\}$, the presentation of $\mathcal{P}_{j}$ in (10.16) is admissible.
Proof. This lemma is essentially clear by inspection; we still write out some details in order to explain certain notation which we will use later. Let $j \in\{1, \ldots, N\}$ be fixed. The setup in Section 3.1 applies to $\mathcal{P}_{j}$ if we set

$$
\left[\Psi \text { for } \mathcal{P}_{j}\right]:=\left\{(1, i): i \in\left\{1, \ldots, r_{j}\right\}\right\}
$$

(thus $\left[N\right.$ for $\left.\mathcal{P}_{j}\right]=1$ and $\left[r_{1}\right.$ for $\left.\mathcal{P}_{j}\right]=r_{j}$ ), and then, for each $i \in\left\{1, \ldots, r_{j}\right\}$, set

$$
\left[c_{1, i} \text { for } \mathcal{P}_{j}\right]:=c_{j, i}, \quad\left[\boldsymbol{w}_{1, i} \text { for } \mathcal{P}_{j}\right]:=\boldsymbol{w}_{j, i}, \quad \text { and }\left[M_{1} \text { for } \mathcal{P}_{j}\right]:=M_{j}
$$

so that $\left[\mathcal{L}_{(1, i)}\right.$ for $\left.\mathcal{P}_{j}\right]=\mathcal{L}_{(j, i)}$ (see (3.4)). Indeed, with these definitions the analogue of the formula in (3.5) holds for $\mathcal{P}_{j}$. Of course, this formula is exactly the same as (10.16), once we identify $\left[\Psi\right.$ for $\mathcal{P}_{j}$ ] with $\Psi_{j}$ through $(1, i) \leftrightarrow(j, i)$.

Now the condition (3.6) trivially holds for $\mathcal{P}_{j}$ since it holds for $\mathcal{P}$, and the condition (3.7) for $\mathcal{P}_{j}$ is void, since $\left[N\right.$ for $\left.\mathcal{P}_{j}\right]=1$. Next, we have $\left[\boldsymbol{c}_{1}^{(1, i)}\right.$ for $\left.\mathcal{P}_{j}\right]=\boldsymbol{c}_{j}^{(j, i)}$ by (5.1), and $\left[W_{1}^{(1, i)}\right.$ for $\left.\mathcal{P}_{j}\right]=W_{j}^{(j, i)}$ by (5.2); therefore $\left[L_{j}^{(1, i)}\right.$ for $\left.\mathcal{P}_{j}\right]=L_{j}^{(j, i)}$ by (5.5). Hence by Definition 5.1, our presentation of $\mathcal{P}_{j}$ is admissible if and only if $\boldsymbol{c}_{j}^{(j, i)} \in L_{j}^{(j, i)}+\mathbb{Z}^{r_{j}}$ for all $i \in\left\{1, \ldots, r_{j}\right\}$; and these conditions are certainly fulfilled, since we are assuming that the fixed presentation of $\mathcal{P}$ in (3.5) is admissible.

Recall from Section 1.2 that for any $j \in\{1, \ldots, N\}$ and $\psi, \psi^{\prime} \in \Psi_{j}$, the collision kernels ${ }^{\circ} p^{(\psi)}$ and ${ }^{\circ} p^{\left(\psi \rightarrow \psi^{\prime}\right)}$ are defined as $\left[p^{(\psi)}\right.$ and $p^{\left(\psi \rightarrow \psi^{\prime}\right)}$ for $\left.\mathcal{P}_{j_{\psi}}\right]$. Similarly let us define the transition kernels ${ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right)$ and ${ }^{\circ} k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)$ as $\left[k(\cdots)\right.$ and $k^{\mathrm{g}}(\cdots)$ for $\left.\mathcal{P}_{j_{\psi}}\right]$. Now the product formulas announced in Section 1.2 are as follows.

Theorem 10.3. Given any $\left\langle\boldsymbol{w}^{\prime}, \psi^{\prime}\right\rangle \in \mathcal{B}_{1}^{d-1} \times \Psi$ and $\psi \in \Psi$, the following formula holds for Lebesgue almost all $\langle\xi, \boldsymbol{w}\rangle \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1}$ :

$$
k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right)=\left\{\begin{array}{lc}
\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{\psi}}} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right) \prod_{j \neq j_{\psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime} & \text { if } j_{\psi^{\prime}}=j_{\psi}  \tag{10.17}\\
\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{\psi}} \bar{n}_{j_{\psi^{\prime}}} v_{d-1}}{ }^{\circ} k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)^{\circ} k^{\mathrm{g}}\left(\xi,-\boldsymbol{w}^{\prime}, \psi^{\prime}\right) \prod_{j \neq j_{\psi}, j_{\psi^{\prime}}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime} \\
& \text { if } j_{\psi^{\prime}} \neq j_{\psi}
\end{array}\right.
$$

Furthermore, for all $\langle\xi, \boldsymbol{w}, \psi\rangle \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi$ :

$$
\begin{equation*}
k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{\psi}}}{ }^{\circ} \mathrm{g}^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi) \prod_{j \neq j_{\psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime} \tag{10.18}
\end{equation*}
$$

Each product over $j$ in (10.17) and (10.18) should be understood to run over all $j \in\{1, \ldots, N\}$ except those numbers which are explicitly excluded.

Remark 10.4. It is not difficult to translate Theorem 10.3 into a product formula involving the collision kernels instead, using (10.1) and (10.2). For example, when $j_{\psi^{\prime}}=j_{\psi}$, the formula in (10.17) is equivalent to

$$
\begin{equation*}
p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{\psi}}}{ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) \prod_{j \neq j_{\psi}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime} \tag{10.19}
\end{equation*}
$$

But when $j_{\psi^{\prime}} \neq j_{\psi}$ (and assuming $\boldsymbol{V}^{\prime} \in \mathcal{V}_{\boldsymbol{V}^{\prime \prime}}$, (10.17) translates into the following somewhat more complicated formula:
$p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{\psi}} \bar{n}_{j_{\psi^{\prime}}} \sigma\left(\boldsymbol{V}^{\prime}, \boldsymbol{W}\right)} p^{(\psi)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) p^{\left(\psi^{\prime}\right)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{W}\right) \prod_{j \neq j_{\psi}, j_{\psi^{\prime}}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}$,
with $\boldsymbol{W}=\Psi_{1}\left(\boldsymbol{V}^{\prime},-\boldsymbol{\beta}_{\boldsymbol{V}^{\prime \prime}}^{+}\left(\boldsymbol{V}^{\prime}\right)\right)$, where the maps " $\Psi_{1}$ " and " $\boldsymbol{\beta}^{+}$" are as in [21, Sec. 3.4] (these are defined in terms of the fixed scattering process of the Lorentz gas). If the scattering process preserves angular momentum, then we have in fact $\boldsymbol{W}=\boldsymbol{V}^{\prime \prime}$, by [21, (3.33)].

Proof of Theorem 10.3. We first need to introduce some further notation. For each $j \in$ $\{1, \ldots, N\}$ let us set $\Sigma_{j}:=\left\{(\psi, \omega) \in \Sigma: j_{\psi}=j\right\}$ and $\mathcal{X}_{j}=\mathbb{R}^{d} \times \Sigma_{j}$. Note that each $\Sigma_{j}$ is a clopen subset of $\Sigma$, and each $\mathcal{X}_{j}$ is a clopen subset of $\mathcal{X}$, and $\Sigma$ and $\mathcal{X}$ decompose as disjoint unions $\Sigma=\sqcup_{j=1}^{N} \Sigma_{j}$ and $\mathcal{X}=\sqcup_{j=1}^{N} \mathcal{X}_{j}$. Also for each $j \in\{1, \ldots, N\}$ we define the map $J^{(j)}: \mathbb{X}_{j} \rightarrow N_{s}\left(\mathcal{X}_{j}\right)$ through

$$
\begin{equation*}
J^{(j)}\left(\Gamma_{j} g\right)=\bigcup_{i=1}^{r_{j}} c_{j, i}\left(\mathbb{Z}^{d} \mathrm{a}_{i}(g)\right) \times\left\{\sigma^{(j, i)}\right\} \quad\left(g \in G_{j}\right) \tag{10.20}
\end{equation*}
$$

Using these maps $J^{(j)}$, the formula for $J: \mathbb{X} \rightarrow N_{s}(\mathcal{X})$ in (6.12) can be written

$$
\begin{equation*}
J(x)=\bigcup_{j=1}^{N} J^{(j)}\left(\widetilde{p}_{j}(x)\right) \quad(x \in \mathbb{X}) \tag{10.21}
\end{equation*}
$$

Throughout the following, we will identify [ $\Psi$ for $\mathcal{P}_{j}$ ] with $\Psi_{j}$ through $(1, i) \leftrightarrow(j, i)$ whenever convenient, as explained in the proof of Lemma 10.2 above. It is now straightforward to express the concepts defined in Section 6.2 when applied to $\mathcal{P}_{j}$ with its presentation in (10.16). For example, $\left[\Omega\right.$ for $\left.\mathcal{P}_{j}\right]$ now equals $P\left(\mathbb{T}_{j}^{d}\right)^{\prime},\left[\Sigma\right.$ for $\left.\mathcal{P}_{j}\right]$ equals $\Sigma_{j}$, and

$$
\begin{equation*}
\mathrm{m}_{j}:=\left[\mathrm{m} \text { for } \mathcal{P}_{j}\right]=\left.\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j}} \mathrm{~m}\right|_{\Sigma_{j}} \quad \in P\left(\Sigma_{j}\right) \tag{10.22}
\end{equation*}
$$

As we did for m in Section 10.1, we also use $\mathrm{m}_{j}$ to denote the probability measure on $\Psi_{j}$ given by $\mathrm{m}_{j}(\psi):=\mathrm{m}_{j}\left(\sigma^{\psi}\right)\left(\psi \in \Psi_{j}\right)$; this gives back the measure defined in (1.12) (but now writing $\Psi_{j}$ in place of $c$ ). Furthermore, we have $\left[\mathcal{X}\right.$ for $\left.\mathcal{P}_{j}\right]=\mathcal{X}_{j}$ and $\left[\mu_{\mathcal{X}}\right.$ for $\left.\mathcal{P}_{j}\right]=\operatorname{vol} \times \mathrm{m}_{j}=\left.\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j}} \mu_{\mathcal{X}}\right|_{\mathcal{X}_{j}}$.
Also, the map $\left[J\right.$ for $\left.\mathcal{P}_{j}\right]$ is given by $J^{(j)}$ defined in (10.20).
Next, for any $\boldsymbol{w}^{\prime} \in \mathbb{R}^{d-1}$ and $\psi^{\prime} \in \Psi_{j}$, let us write

$$
{ }^{\circ} \kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}:=\left[\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}} \text { for } \mathcal{P}_{j}\right], \quad \text { and also } \quad{ }^{\circ} \kappa_{j}^{\mathrm{g}}:=\left[\kappa^{\mathrm{g}} \text { for } \mathcal{P}_{j}\right] \text {. }
$$

Now it is easy to write out the definitions and facts from Section 10.1 when applied for $\mathcal{P}_{j}$ in the place of $\mathcal{P}$. For example, ${ }^{\circ} k$, which we defined above to be the transition kernel $\left[k\right.$ for $\left.\mathcal{P}_{j}\right]$, is a function

$$
\begin{equation*}
{ }^{\circ} k: \mathbb{R}^{d-1} \times \Psi_{j} \times \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi_{j} \rightarrow\left[0, \bar{n}_{j} v_{d-1}\right] \tag{10.23}
\end{equation*}
$$

such that for each $\boldsymbol{w}^{\prime} \in \mathbb{R}^{d-1}$ and $\psi^{\prime} \in \Psi_{j}$, and for any Borel set $B \subset \mathfrak{Z}_{\infty} \times \Sigma_{j}$, we have

$$
{ }^{\circ} \kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=v_{d-1}^{-1} \int_{\left(\xi, \boldsymbol{w}, \sigma^{\psi}\right) \in B}{ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right) d \xi d \boldsymbol{w} d \mathrm{~m}_{j}(\psi) .
$$

We now start with the actual proof of Theorem 10.3, Let $\boldsymbol{w}^{\prime} \in \mathcal{B}_{1}^{d-1}$ and $\psi^{\prime} \in \Psi$; these will be kept fixed throughout the proof. Let $j_{0} \in\{1, \ldots, N\}$, and let $B$ be a Borel subset of $\mathcal{Z}_{\infty} \times \Sigma_{j_{0}}$. Then by (10.5),

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\overline{\omega^{\psi^{\prime}}}(A), \quad \text { with } A:=\left\{x \in \mathbb{X}: \iota\left(\boldsymbol{z}\left(J(x)-\left(0, \boldsymbol{w}^{\prime}\right)\right)\right) \in B\right\} . \tag{10.24}
\end{equation*}
$$

Recall that $\mathbb{X}=\mathbb{X}_{1} \times \cdots \times \mathbb{X}_{N}$ and $\overline{\omega^{\psi^{\prime}}}=\overline{\omega_{1}^{\psi^{\prime}}} \otimes \cdots \otimes \overline{\omega_{N}^{\psi^{\prime}}}$. Let us write $\mathbb{X}^{\left(j_{0}\right)}=\prod_{j \neq j_{0}} \mathbb{X}_{j}$, so that

$$
\mathbb{X}=\mathbb{X}_{j_{0}} \times \mathbb{X}^{\left(j_{0}\right)}=\left\{\left(x, x^{\prime}\right): x \in \mathbb{X}_{j_{0}}, x^{\prime} \in \mathbb{X}^{\left(j_{0}\right)}\right\}
$$

in an obvious identification. Corresponding to this product decomposition we have $\overline{\omega^{\psi^{\prime}}}=$ $\overline{\omega_{j_{0}}^{\psi^{\prime}}} \otimes \nu$ with $\nu:=\otimes_{j \neq j_{0}} \overline{\omega_{j}^{\psi^{\prime}}}$, and hence, by Fubini's Theorem,

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\overline{\omega^{\psi^{\prime}}}(A)=\int_{\mathbb{X}_{j_{0}}} \nu\left(A_{x}\right) \overline{d \omega_{j_{0}}^{\psi^{\prime}}}(x), \tag{10.25}
\end{equation*}
$$

where for each $x \in \mathbb{X}_{j_{0}}$,

$$
A_{x}:=\left\{x^{\prime} \in \mathbb{X}^{\left(j_{0}\right)}: \iota\left(\boldsymbol{z}\left(J\left(x, x^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)\right) \in B\right\} .
$$

For each $j \in\{1, \ldots, N\}$ we set $\widetilde{\mathbb{X}}_{j}:=\left\{x \in \mathbb{X}_{j}: \boldsymbol{z}\left(J^{(j)}(x)-\left(0, \boldsymbol{w}^{\prime}\right)\right) \neq\right.$ undef $\}$. Then

$$
\begin{equation*}
\overline{\omega_{j}^{\psi^{\prime}}}\left(\mathbb{X}_{j} \backslash \widetilde{\mathbb{X}}_{j}\right)=0 \tag{10.26}
\end{equation*}
$$

Indeed, if $j=j_{\psi^{\prime}}$ then (10.26) follows from ${ }^{\circ} \kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(\{$ undef $\})=0$ (which holds by [21, Lemma 3.1]), and if $j \neq j_{\psi^{\prime}}$ then $\omega_{j}^{\psi^{\prime}}=\omega_{j}^{\mathrm{g}}$ by Lemma 5.11, and so (10.26) follows from the fact that ${ }^{\circ} \kappa_{j}^{\mathrm{g}}(\{$ undef $\})=0$ (which holds by [21, Lemma 3.5]) and the formula (10.8) with $\boldsymbol{v}=-\left(0, \boldsymbol{w}^{\prime}\right)$, applied for $\mathcal{P}_{j}$ in the place of $\mathcal{P}$. Hence, writing also $\widetilde{\mathbb{X}}=\prod_{j=1}^{N} \widetilde{\mathbb{X}}_{j}$ and $\widetilde{\mathbb{X}}^{\left(j_{0}\right)}=\prod_{j \neq j_{0}} \widetilde{\mathbb{X}}_{j}$, we have $\overline{\omega^{\psi^{\prime}}}(\mathbb{X} \backslash \widetilde{\mathbb{X}})=0$ and $\nu\left(\mathbb{X}^{\left(j_{0}\right)} \backslash \widetilde{\mathbb{X}}^{\left(j_{0}\right)}\right)=0$, and using these facts together with (10.26), it follows that:

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\overline{\omega^{\psi^{\prime}}}(A \cap \widetilde{\mathbb{X}})=\int_{\widetilde{\mathbb{X}}_{j_{0}}} \nu\left(A_{x} \cap \widetilde{\mathbb{X}}^{\left(j_{0}\right)}\right) \overline{d \omega_{j_{0}}^{\psi^{\prime}}}(x), \tag{10.27}
\end{equation*}
$$

For any $x \in \widetilde{\mathbb{X}}_{j}$, let us define $\xi(x) \in \mathbb{R}_{>0}, \boldsymbol{w}(x) \in \mathcal{B}_{1}^{d-1}$ and $\psi(x) \in \Psi_{j_{0}}$ through the relation

$$
\begin{equation*}
((\xi(x), \boldsymbol{w}(x)), \psi(x))=\iota\left(\boldsymbol{z}\left(J^{\left(j_{0}\right)}(x)-\left(0, \boldsymbol{w}^{\prime}\right)\right)\right) \quad \text { in } \mathfrak{Z}_{\infty} \times \Sigma_{j_{0}} . \tag{10.28}
\end{equation*}
$$

Let us also set

$$
\widetilde{B}=\left\{x \in \widetilde{\mathbb{X}}_{j_{0}}:((\xi(x), \boldsymbol{w}(x)), \psi(x)) \in B\right\}
$$

We will write any element $x^{\prime}$ in $\mathbb{X}^{\left(j_{0}\right)}=\prod_{j \neq j_{0}} \mathbb{X}_{j}$ as $x^{\prime}=\left(x_{j}^{\prime}\right)_{j \neq j_{0}}$ with $x_{j}^{\prime} \in \mathbb{X}_{j}$. Also, for any $Y \in N_{s}(\mathcal{X})$ such that $\boldsymbol{z}(Y) \neq$ undef, let us denote by $z_{1}(Y)$ the $\boldsymbol{e}_{1}$-coordinate of the point $\boldsymbol{z}(Y)$; in particular we then have $z_{1}\left(J^{\left(j_{0}\right)}(x)-\left(0, \boldsymbol{w}^{\prime}\right)\right)=\xi(x)$ for every $x \in \widetilde{\mathbb{X}}_{j_{0}}$. Using the definition of the map $\boldsymbol{z}$ and our assumption that $B \subset \mathfrak{Z}_{\infty} \times \Sigma_{j_{0}}$, it follows that for every $\left(x, x^{\prime}\right) \in \widetilde{\mathbb{X}}_{j_{0}} \times \widetilde{\mathbb{X}}^{\left(j_{0}\right)}$, the condition $\iota\left(\boldsymbol{z}\left(J\left(x, x^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)\right) \in B$ holds if and only if $x \in \widetilde{B}$ and $z_{1}\left(J^{(j)}\left(x_{j}^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)>\xi(x)$ for all $j \neq j_{0}$. Hence (10.27) can be rewritten as

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\int_{\widetilde{B}} \prod_{j \neq j_{0}} \overline{\omega_{j}^{\psi^{\prime}}}\left(\left\{x_{j}^{\prime} \in \widetilde{\mathbb{X}}_{j}: z_{1}\left(J^{(j)}\left(x_{j}^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)>\xi(x)\right\}\right) d \overline{\omega_{j 0}^{\psi^{\prime}}}(x) \tag{10.29}
\end{equation*}
$$

Here for each $j \neq j_{\psi^{\prime}}$ we have

$$
\begin{array}{r}
\overline{\omega_{j}^{\psi^{\prime}}}\left(\left\{x_{j}^{\prime} \in \widetilde{\mathbb{X}}_{j}: z_{1}\left(J^{(j)}\left(x_{j}^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)>\xi(x)\right\}\right)={ }^{\circ} \kappa_{j}^{\mathrm{g}}\left((\xi(x), \infty) \times \mathcal{B}_{1}^{d-1} \times \widetilde{\Psi}_{j}\right) \\
=v_{d-1}^{-1} \int_{\xi(x)}^{\infty} \int_{\mathcal{B}_{1}^{d-1}} \int_{\Psi_{j}}{ }^{\circ} k^{\mathrm{g}}\left(\xi^{\prime}, \boldsymbol{w}, \psi\right) d \mathrm{~m}_{j}(\psi) d \boldsymbol{w} d \xi=\int_{\xi(x)}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}
\end{array}
$$

where the first equality follows from $\omega_{j}^{\psi^{\prime}}=\omega_{j}^{\mathrm{g}}$ and the formula (10.8) with $\boldsymbol{v}=-\left(0, \boldsymbol{w}^{\prime}\right)$, applied for $\mathcal{P}_{j}$ in place of $\mathcal{P}$, and the last equality holds by (10.13) applied for $\mathcal{P}_{j}$. On the other hand, for $j=j_{\psi^{\prime}}$ we have

$$
\begin{aligned}
\overline{\omega_{j}^{\psi^{\prime}}} & \left(\left\{x_{j}^{\prime} \in \widetilde{\mathbb{X}}_{j}: z_{1}\left(J^{(j)}\left(x_{j}^{\prime}\right)-\left(0, \boldsymbol{w}^{\prime}\right)\right)>\xi(x)\right\}\right)={ }^{\circ} \kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}\left((\xi(x), \infty) \times \mathcal{B}_{1}^{d-1} \times \widetilde{\Psi}_{j}\right) \\
& =v_{d-1}^{-1} \int_{\xi(x)}^{\infty} \int_{\mathcal{B}_{1}^{d-1}} \int_{\Psi_{j}}{ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi^{\prime}, \boldsymbol{w}, \psi\right) d \mathrm{~m}_{j}(\psi) d \boldsymbol{w} d \xi=\frac{1}{\bar{n}_{j} v_{d-1}}{ }^{\circ} k_{j}^{\mathrm{g}}\left(\xi(x),-\boldsymbol{w}^{\prime}, \psi^{\prime}\right)
\end{aligned}
$$

where in the last equality we used Lemma 10.1 for $\mathcal{P}_{j}$.
Let us first assume $j_{0}=j_{\psi^{\prime}}$. Then the above formulas imply that

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\int_{\widetilde{B}}\left(\prod_{j \neq j_{0}} \int_{\xi(x)}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) \overline{\omega_{j_{0}}^{\psi^{\prime}}}(x) \tag{10.30}
\end{equation*}
$$

We here change to new integration variables $(\xi, \boldsymbol{w}, \psi)$ through (10.28); then by the definition of ${ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \cdots\right)$ we have $d \overline{\omega_{j_{0}}^{\psi^{\prime}}}(x)=v_{d-1}^{-1}{ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right) d \xi d \boldsymbol{w} d \mathbf{m}_{j_{0}}(\psi)$. Hence, using also (10.22) to express $\mathrm{m}_{j_{0}}$ in terms of m , we get

$$
\begin{equation*}
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{0}} v_{d-1}} \int_{B}{ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}, \psi\right)\left(\prod_{j \neq j_{0}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \boldsymbol{w} d \mathbf{m}(\psi) \tag{10.31}
\end{equation*}
$$

The fact that (10.31) holds for all Borel sets $B \subset \mathfrak{Z}_{\infty} \times \Sigma_{j_{0}}$ implies that the first formula in (10.17) holds for almost all $\langle\xi, \boldsymbol{w}, \psi\rangle \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi_{j_{0}}$.

Next assume $j_{0} \neq j_{\psi^{\prime}}$. Then in a similar way, we get

$$
\kappa_{\boldsymbol{w}^{\prime}, \psi^{\prime}}(B)=\frac{1}{\bar{n}_{j_{\psi^{\prime}}} v_{d-1}} \int_{\widetilde{B}}{ }^{\circ} k_{j_{\psi^{\prime}}}^{\mathrm{g}}\left(\xi(x),-\boldsymbol{w}^{\prime}, \psi^{\prime}\right)\left(\prod_{j \neq j_{\psi^{\prime}}, j_{0}} \int_{\xi(x)}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \overline{\omega_{j_{0}}^{\psi^{\prime}}}(x)
$$

$$
\begin{equation*}
=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{0}} \bar{n}_{j_{\psi^{\prime}}} v_{d-1}^{2}} \int_{B}{ }^{\circ} k_{j_{0}}^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)^{\circ} k_{j_{\psi^{\prime}}}^{\mathrm{g}}\left(\xi,-\boldsymbol{w}^{\prime}, \psi^{\prime}\right)\left(\prod_{j \neq j_{\psi^{\prime}, j_{0}}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \boldsymbol{w} d \mathrm{~m}(\psi) \tag{10.32}
\end{equation*}
$$

Here, to obtain the last equality, we again changed to new integration variables $(\xi, \boldsymbol{w}, \psi)$ through (10.28) and used $\omega_{j_{0}}^{\psi^{\prime}}=\omega_{j_{0}}^{\mathrm{g}}$ and the formula (10.8) with $\boldsymbol{v}=-\left(0, \boldsymbol{w}^{\prime}\right)$ (applied for $\mathcal{P}_{j}$ in place of $\mathcal{P})$, to see that $\overline{d \omega_{j_{0}}^{\psi^{\prime}}}(x)=v_{d-1}^{-1}{ }^{\circ}{ }_{j_{\psi^{\prime}}}^{\mathrm{g}}\left(\xi,-\boldsymbol{w}^{\prime}, \psi^{\prime}\right) d \xi d \boldsymbol{w} d \mathrm{~m}_{j_{0}}(\psi)$. Finally, the fact that (10.32) holds for all Borel sets $B \subset \mathfrak{Z}_{\infty} \times \Sigma_{j_{0}}$ implies that the second formula in (10.17) holds for almost all $\langle\xi, \boldsymbol{w}, \psi\rangle \in \mathbb{R}_{>0} \times \mathcal{B}_{1}^{d-1} \times \Psi_{j_{0}}$.

Since $j_{0}$ was arbitrary, we have now proved all of (10.17). The formula (10.18) follows by a very similar computation, starting from the relation $\kappa^{\mathrm{g}}(B)=\overline{\omega \mathrm{g}}(A)$ instead of (10.24), leading to

$$
\kappa^{\mathrm{g}}(B)=\frac{\bar{n}_{\mathcal{P}}}{\bar{n}_{j_{0}} v_{d-1}} \int_{B}{ }^{\circ} k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)\left(\prod_{j \neq j_{0}} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \boldsymbol{w} d \mathrm{~m}(\psi) .
$$

The fact that (10.18) holds for all (and not just almost all) $\xi, \boldsymbol{w}, \psi$ follows by continuity, see Lemma 10.1 .

This completes the proof of Theorem [10.3, and thus also the proof of Theorem 10.4.
Remark 10.5. In the special case when all the grids appearing in $\mathcal{P}$ are incommensurable and $\bar{n}_{\mathcal{P}}=1$, the product formulas in Theorem 10.3 agree with the formulas [20, (5.5), (5.6), (5.12)]. This is immediately verified using the formulas in Remark 10.1 and Remark 10.2, and the fact that ${ }^{\circ} k^{\mathrm{g}}\left(\xi,-\boldsymbol{w}^{\prime}, \psi^{\prime}\right)={ }^{\circ} k^{\mathrm{g}}\left(\xi, \boldsymbol{w}^{\prime}, \psi^{\prime}\right)$, since $\mathcal{P}_{j_{\psi^{\prime}}}$ is a single grid. (The last mentioned symmetry relation follows from Remark 10.1 and [16, Remark 4.5 18 . In fact, in dimension $d \geq 3$ we have for our general $\mathcal{P}$ that $k^{\mathrm{g}}(\xi, \boldsymbol{w}, \psi)$ only depends on $\xi,\|\boldsymbol{w}\|, \psi$, by [21, Lemma 3.18].)

We next give the proof of Theorem 1.3. We will here stay close to the notation used in the statement of that theorem, thus, e.g., we return to parametrizing the equivalence classes in $\Psi$ by $c \in C_{\Psi}$ instead of $j \in\{1, \ldots, N\}$.
Proof of Theorem 1.3. Since the set $A$ may be decomposed as the disjoint union of the subsets $A \cap\left(\mathbb{R}_{>0} \times\{\psi\} \times \mathrm{S}_{1}^{d-1}\right)$ with $\psi$ running through $\Psi$, it suffices to prove the statements of the corollary when $A$ is a subset of $\mathbb{R}_{>0} \times\left\{\psi_{0}\right\} \times \mathrm{S}_{1}^{d-1}$ for some fixed $\psi_{0} \in \Psi$. In this case, for any given $\psi^{\prime} \in \Psi$, it follows from the definition of $\widetilde{c}_{\psi^{\prime}}$ that the right hand side of (1.16) equals

$$
\begin{equation*}
\mathbb{P}\left(\widetilde{\xi}_{\left[\psi_{0}\right], \psi^{\prime}}<\widetilde{\xi}_{c, \psi^{\prime}} \forall c \neq\left[\psi_{0}\right] \text { and }\left\langle\widetilde{\xi}_{\left[\psi_{0}\right], \psi^{\prime}}, \widetilde{\psi}_{\left[\psi_{0}\right], \psi^{\prime}}, \tilde{\boldsymbol{V}}_{\left[\psi_{0}\right], \psi^{\prime}}\right\rangle \in A\right) . \tag{10.33}
\end{equation*}
$$

Using the definition of the random triples $\left\langle\widetilde{\xi}_{c, \psi^{\prime}}, \widetilde{\psi}_{c, \psi^{\prime}}, \widetilde{\boldsymbol{V}}_{c, \psi^{\prime}}\right\rangle$, the above probability can be expressed as follows, if $\psi^{\prime} \nsim \psi_{0}$ :

$$
\begin{aligned}
& \int_{A}{ }^{\circ}{ }^{\circ}{ }^{\left(\psi_{0}\right)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)\left(\int_{(\xi, \infty) \times\left[\psi^{\prime}\right] \times S_{1}^{d-1}}{ }^{\circ}{ }^{\circ}\left(\psi^{\prime} \rightarrow \psi_{1}\right)\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi_{1}, \boldsymbol{V}_{1}\right) d \xi_{1} d \mathrm{~m}_{\left[\psi^{\prime}\right]}\left(\psi_{1}\right) d \sigma\left(\boldsymbol{V}_{1}\right)\right) \\
& \quad \times\left(\prod_{c \neq\left[\psi_{0}\right],\left[\psi^{\prime}\right]} \int_{(\xi, \infty) \times c \times S_{1}^{d-1}}{ }^{\circ} p^{\left(\psi_{1}\right)}\left(\boldsymbol{V}^{\prime} ; \xi_{1}, \boldsymbol{V}_{1}\right) d \xi_{1} d \mathbf{m}_{c}\left(\psi_{1}\right) d \sigma\left(\boldsymbol{V}_{1}\right)\right) d \xi d \mathrm{~m}_{\left[\psi_{0}\right]}(\psi) d \sigma(\boldsymbol{V}) \\
& =\int_{A}{ }^{\circ} p^{\left(\psi_{0}\right)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)\left(\frac{1}{v_{d-1}} \int_{(\xi, \infty) \times\left[\psi^{\prime}\right] \times \mathcal{B}_{1}^{d-1}}{ }^{\circ} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi_{1}, \boldsymbol{w}_{1}, \psi_{1}\right) d \xi_{1} d \mathrm{~m}_{\left[\psi^{\prime}\right]}\left(\psi_{1}\right) d \boldsymbol{w}_{1}\right) \\
& \quad \times\left(\prod_{c \neq\left[\psi_{0}\right],\left[\psi^{\prime}\right]} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \mathrm{~m}_{\left[\psi_{0}\right]}(\psi) d \sigma(\boldsymbol{V}),
\end{aligned}
$$

where $\boldsymbol{w}^{\prime} \in \mathcal{B}_{1}^{d-1}$ is the exit parameter determined by the fixed vectors $\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime}$. To get the last equality, in the integral over $(\xi, \infty) \times\left[\psi^{\prime}\right] \times \mathrm{S}_{1}^{d-1}$ we introduced the new variable $\boldsymbol{w}_{1}:=$ the impact parameter determined by $\boldsymbol{V}^{\prime}, \boldsymbol{V}_{1}$, and used (10.2) and the defining property of the differential cross section (cf. [21, Sec. 3.4 and Lemma 3.26]); also, inside the product over $c$ we applied (1.18).

Next we apply Lemma 10.1 to identify the integral over $(\xi, \infty) \times\left[\psi^{\prime}\right] \times \mathcal{B}_{1}^{d-1}$, and in the integral over $A$ we substitute $\boldsymbol{w}_{2}:=$ the impact parameter determined by $\boldsymbol{V}^{\prime}, \boldsymbol{V}$, and use (10.1). It follows that the above expression equals

$$
\frac{1}{\bar{n}_{\left[\psi^{\prime}\right]} v_{d-1}^{2}} \int_{\bar{A}_{\boldsymbol{V}^{\prime}}}{ }^{\circ} k\left(\xi, \boldsymbol{w}_{2}, \psi_{0}\right)^{\circ} k\left(\xi,-\boldsymbol{w}^{\prime}, \psi^{\prime}\right)\left(\prod_{c \neq\left[\psi_{0}\right],\left[\psi^{\prime}\right]} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \mathrm{~m}_{\left[\psi_{0}\right]}(\psi) d \boldsymbol{w}_{2},
$$

[^14]where $\bar{A}_{\boldsymbol{V}^{\prime}}$ is the set of triples $\left(\xi, \psi, \boldsymbol{w}_{2}\right)$ which arise as images of the $(\xi, \psi, \boldsymbol{V}) \in A$. Finally, using $\mathrm{m}_{\left[\psi_{0}\right]}=\left.\left(\bar{n}_{\mathcal{P}} / \bar{n}_{\left[\psi_{0}\right]}\right) \mathrm{m}\right|_{\left[\psi_{0}\right]}$ and the product formula (10.17) in Theorem 10.3, and our assumption that $A \subset \mathbb{R}_{>0} \times\left\{\psi_{0}\right\} \times \mathrm{S}_{1}^{d-1}\left(\Rightarrow \bar{A}_{\boldsymbol{V}^{\prime}} \subset \mathbb{R}_{>0} \times\left\{\psi_{0}\right\} \times \mathcal{B}_{1}^{d-1}\right)$, the above expression becomes
$$
\frac{1}{v_{d-1}} \int_{\bar{A}_{\boldsymbol{V}^{\prime}}} k\left(\boldsymbol{w}^{\prime}, \psi^{\prime}, \xi, \boldsymbol{w}_{2}, \psi\right) d \xi d \mathrm{~m}(\psi) d \boldsymbol{w}_{2}=\int_{A} p^{\left(\psi^{\prime} \rightarrow \psi\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right) d \xi d \mathrm{~m}(\psi) d \sigma(\boldsymbol{V})
$$

By (1.8) in Theorem (1.2, this equals the left hand side of (1.16)!
In the case $\psi^{\prime} \sim \psi_{0}$, the probability in (10.33) instead equals

$$
\int_{A}{ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi_{0}\right)}\left(\boldsymbol{V}^{\prime \prime}, \boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)\left(\prod_{c \neq\left[\psi_{0}\right]} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \mathrm{~m}_{\left[\psi_{0}\right]}(\psi) d \sigma(\boldsymbol{V})
$$

and by the same type of computation as above (but easier, and using the first equality in (10.17) instead of the second), this is again seen to be equal to the left hand side of (1.16). This completes the proof of (1.16).

The proof of (1.15) is completely similar, but easier: Again assuming $A \subset \mathbb{R}_{>0} \times\left\{\psi_{0}\right\} \times \mathrm{S}_{1}^{d-1}$, the right hand side of (1.15) equals

$$
\begin{aligned}
\mathbb{P}\left(\widetilde{\xi}_{\left[\psi_{0}\right]}<\right. & \left.\widetilde{\xi}_{c} \forall c \neq\left[\psi_{0}\right] \text { and }\left\langle\widetilde{\xi}_{\left[\psi_{0}\right]}, \widetilde{\psi}_{\left[\psi_{0}\right]}, \tilde{\boldsymbol{V}}_{\left[\psi_{0}\right]}\right\rangle \in A\right) \\
& =\int_{A}{ }^{\circ} p^{\left(\psi_{0}\right)}\left(\boldsymbol{V}^{\prime} ; \xi, \boldsymbol{V}\right)\left(\prod_{c \neq\left[\psi_{0}\right]} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{c}}\left(\xi^{\prime}\right) d \xi^{\prime}\right) d \xi d \mathrm{~m}_{\left[\psi_{0}\right]}(\psi) d \sigma(\boldsymbol{V})
\end{aligned}
$$

and by the same type of computation as above, using now (10.18), this is seen to be equal to the left hand side of (1.15).
10.3. Asymptotic estimates for the free path length distribution. We will now prove Theorem 1.5, This theorem will be a quite easy consequence of the following result.

Proposition 10.4. There exist constants $0<c_{1}<c_{2}$, which only depend on $\mathcal{P}$, such that

$$
\begin{equation*}
c_{1} \xi^{-N}<\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime}<c_{2} \xi^{-N}, \quad \forall \xi \geq 1 \tag{10.34}
\end{equation*}
$$

Proof. By Corollary 1.4, $\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime}=\prod_{j=1}^{N} \int_{\xi}^{\infty} \Phi_{\mathcal{P}_{j}}\left(\xi^{\prime}\right) d \xi^{\prime}$; thus, if we can prove the proposition when $N=1$, it follows for general $N$. Hence from now on we assume $N=1$.

By (10.13) and the definition of $k^{\mathrm{g}}$, we have

$$
\begin{align*}
\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime} & =\frac{1}{v_{d-1}} \int_{[\xi, \infty) \times \mathcal{B}_{1}^{d-1} \times \Psi} k^{\mathrm{g}}\left(\xi^{\prime}, \boldsymbol{w}, \psi\right) d \xi^{\prime} d \boldsymbol{w} d \mathrm{~m}(\psi) \\
& =\kappa^{\mathrm{g}}\left([\xi, \infty) \times \mathcal{B}_{1}^{d-1} \times \Sigma\right)=\overline{\omega^{\mathrm{g}}}\left(\left\{x \in \mathbb{X}: J(x) \cap\left(\mathfrak{Z}_{\xi} \times \Sigma\right)=\emptyset\right\}\right) \tag{10.35}
\end{align*}
$$

where in the last equality we used (10.8), $\kappa^{\mathrm{g}}(\{$ undef $\})=0$, and the definition of the map $\boldsymbol{z}$. Now since $N=1$ we have $\mathbb{X}=\mathbb{X}_{1}$ and $\omega^{\mathrm{g}}=\omega_{1}^{\mathrm{g}}$, and $\overline{\omega^{\mathrm{g}}}$ is $\mathrm{SL}_{d}(\mathbb{R})$-invariant by Lemma 4.5, in particular $\overline{\omega^{g}}$ is invariant under (right) multiplication by $D:=D_{\xi^{1 / d}}=$ $\operatorname{diag}\left(\xi^{1-\frac{1}{d}}, \xi^{-\frac{1}{d}}, \ldots, \xi^{-\frac{1}{d}}\right)$, and so we may replace $J(x)$ by $J(x D)$ in the last expression in (10.35). Using also $J(x D)=J(x) D$, we see that this is equivalent with replacing, in the same expression, $\mathfrak{Z}_{\xi}$ by

$$
\widetilde{\mathfrak{Z}}_{\xi}:=\mathfrak{\mathfrak { Z }}_{\xi} D^{-1}=\left(0, \xi^{1 / d}\right) \times \mathcal{B}_{\xi^{1 / d}}^{d-1}
$$

Using also Lemma 4.4 for $\omega^{\mathrm{g}}=\omega_{1}^{\mathrm{g}}$, and the definition of $J$ in (6.12), we conclude:
(10.36)

$$
\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime}=\int_{F_{d}} \omega^{\mathrm{g}}\left(\left\{U \in \mathbb{T}_{1}^{d}: c_{1, i}\left(\mathbb{Z}^{d}+\widetilde{\mathrm{r}}_{i}(U)\right) A \cap \widetilde{\mathfrak{Z}}_{\xi}=\emptyset \text { for } i=1, \ldots, r_{1}\right\}\right) d \nu(A)
$$

Here $\widetilde{\mathrm{r}}_{i}(U) \in(\mathbb{R} / \mathbb{Z})^{d}$, and in the above expression we use the convention that for any $\boldsymbol{u} \in$ $(\mathbb{R} / \mathbb{Z})^{d}, \mathbb{Z}^{d}+\boldsymbol{u}$ denotes the inverse image of $\boldsymbol{u}$ in $\mathbb{R}^{d}$ (that is, $\mathbb{Z}^{d}+\boldsymbol{u}:=\mathbb{Z}^{d}+\boldsymbol{v}$, where $\boldsymbol{v}$ is any lift of $\boldsymbol{u}$ to $\mathbb{R}^{d}$ ).

To get an upper bound, we fix an arbitrary $i \in\left\{1, \ldots, r_{1}\right\}$, and note that the previous integral is bounded above by

$$
\int_{F_{d}}\left[\widetilde{\mathbf{r}}_{i *}\left(\omega^{\mathrm{g}}\right)\right]\left(\left\{\boldsymbol{u} \in(\mathbb{R} / \mathbb{Z})^{d}: c_{1, i}\left(\mathbb{Z}^{d}+\boldsymbol{u}\right) A \cap \widetilde{\mathfrak{Z}}_{\xi}=\emptyset\right\}\right) d \nu(A)
$$

Here note that the measure $\widetilde{\mathrm{r}}_{i *}\left(\omega^{\mathrm{g}}\right)$ is invariant under translations, since, as we saw in the proof of Lemma [7.15, $\omega^{\mathrm{g}}=\omega_{1}^{\mathrm{g}}$ is invariant under $X \mapsto X+\widetilde{\boldsymbol{c}}_{1} \boldsymbol{v}$ for all $\boldsymbol{v} \in \mathbb{R}^{d}$, and $\mathrm{r}_{i}\left(\widetilde{\boldsymbol{c}}_{1} \boldsymbol{v}\right)=c_{1, i}^{-1} \boldsymbol{v}$ with $c_{1, i}>0$. Hence $\widetilde{\mathbf{r}}_{i *}\left(\omega^{\mathrm{g}}\right)$ is in fact the Lebesgue probability measure on $(\mathbb{R} / \mathbb{Z})^{d}$, and it follows that in the last integral equals

$$
\int_{\mathrm{ASL}_{d}(\mathbb{Z}) \backslash \mathrm{ASL}_{d}(\mathbb{R})} I\left(j(x) \cap c_{1, i}^{-1} \widetilde{\mathfrak{\jmath}}_{\xi}=\emptyset\right) d \widetilde{\nu}(x)
$$

where $\widetilde{\nu}$ is the $\mathrm{ASL}_{d}(\mathbb{R})$ invariant probability measure on $\mathrm{ASL}_{d}(\mathbb{Z}) \backslash \mathrm{ASL}_{d}(\mathbb{R})$ and $j$ is the standard identification between $\mathrm{ASL}_{d}(\mathbb{Z}) \backslash \mathrm{ASL}_{d}(\mathbb{R})$ and the space of grids of density one in $\mathbb{R}^{d}$, that is, $j\left(\mathrm{ASL}_{d}(\mathbb{Z}) g\right):=\mathbb{Z}^{d} g$ for any $g \in \mathrm{ASL}_{d}(\mathbb{R})$. This integral is known to be asymptotically equal to $c \xi^{-1}$ as $\xi \rightarrow \infty$ with an explicit constant $c$ [18, Theorem 3.5 and (3.21), (3.22)]. Hence the upper bound in (10.34) (for $N=1$ ) holds for an appropriate constant $c_{2}$.

It remains to prove the lower bound. We will need the following lemma.
Lemma 10.5. There exists a constant $c=c(d)>0$ such that for any $B \geq 1$, there exists a subset $M_{B} \subset F_{d}$ with $\nu\left(M_{B}\right) \geq c B^{-d}$ such that for every $A \in M_{B}$ there exists a vector $\ell \in \mathbb{R}^{d}$ satisfying $\|\ell\| \geq B$ and $\mathbb{Z}^{d} A \subset \mathbb{Z} \boldsymbol{\ell}+\boldsymbol{\ell}^{\perp}$.
(Here $\boldsymbol{\ell}^{\perp}$ denotes the hyperplane orthogonal to $\boldsymbol{\ell}$.)
Proof. Consider the Siegel set

$$
\mathcal{S}_{d}:=\left\{u \operatorname{diag}\left(a_{1}, \ldots, a_{d}\right) k: u \in F_{U}, 0<a_{j+1} \leq \frac{2}{\sqrt{3}} a_{j}(j=1, \ldots, d-1), k \in \mathrm{SO}(d)\right\}
$$

where $F_{U}$ is the set of all unipotent upper triangular matrices in $\mathrm{SL}_{d}(\mathbb{R})$ all of whose elements above the diagonal belong to $\left[-\frac{1}{2}, \frac{1}{2}\right]$. It is known that $\mathcal{S}_{d}$ contains a fundamental region for $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$, and on the other hand, the number of points in $\mathcal{S}_{d}$ which project to any given point in $\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ is bounded above by a constant $K=K(d)$ (see [3]). Note that for any $A=u \operatorname{diag}\left(a_{1}, \ldots, a_{d}\right) k \in \mathcal{S}_{d}$ we have $\boldsymbol{e}_{1} u \subset \boldsymbol{e}_{1}+\boldsymbol{e}_{1}^{\perp}$ and $\boldsymbol{e}_{j} u \subset \boldsymbol{e}_{1}^{\perp}$ for all $j \in\{2, \ldots, d\}$; hence the vector $\boldsymbol{\ell}:=a_{1} \boldsymbol{e}_{1} k$ satisfies $\mathbb{Z}^{d} A \subset \mathbb{Z} \boldsymbol{\ell}+\ell^{\perp}$. Therefore, given $B \geq 1$, if we set

$$
M_{B}^{\prime}:=\left\{u \operatorname{diag}\left(a_{1}, \ldots, a_{d}\right) k \in \mathcal{S}_{d}: a_{1} \geq B\right\}
$$

then for every $A \in M_{B}^{\prime}$ there exists $\ell \in \mathbb{R}^{d}$ satisfying $\|\ell\| \geq B$ and $\mathbb{Z}^{d} A \subset \mathbb{Z} \ell+\ell^{\perp}$. Note also that this property only depends on the lattice $\mathbb{Z}^{d} A$, that is, it depends only on the coset $\mathrm{SL}_{d}(\mathbb{Z}) A$. Hence if we set $M_{B}:=F_{d} \cap\left(\mathrm{SL}_{d}(\mathbb{Z}) M_{B}^{\prime}\right)$, then for every $A \in M_{B}$ there again exists $\ell \in \mathbb{R}^{d}$ satisfying $\|\ell\| \geq B$ and $\mathbb{Z}^{d} A \subset \mathbb{Z} \boldsymbol{\ell}+\boldsymbol{\ell}^{\perp}$. Furthermore,

$$
\nu\left(M_{B}\right) \geq K^{-1} \nu\left(M_{B}^{\prime}\right) \gg \int_{B}^{\infty} a_{1}^{-d-1} d a_{1} \gg B^{-d}
$$

where the lower bound on $\nu\left(M_{B}^{\prime}\right)$ follows from a standard integration using the expression for $\nu$ in terms of the Iwasawa decomposition of $\mathrm{SL}_{d}(\mathbb{R})$ (see, e.g., [30, (2.11) and (2.12)], and use also the fact that for $B$ sufficiently large, the set of $A^{\prime} \in \mathcal{S}_{d-1}$ with $a_{1} \leq \frac{2}{\sqrt{3}} B^{d /(d-1)}$ has Haar measure bounded away from zero independently of $B$ ).

Now let $M_{B} \subset F_{d}$ be the set provided by Lemma 10.5, applied with

$$
\begin{equation*}
B:=4\left(\sum_{i=1}^{r_{1}} c_{1, i}^{-1}\right) C_{\xi} \quad \text { where } C_{\xi}:=\sup \left\{|\boldsymbol{v} \cdot \boldsymbol{w}|: \boldsymbol{v} \in \mathrm{S}_{1}^{d-1}, \boldsymbol{w} \in \widetilde{\mathfrak{Z}}_{\xi}\right\} \tag{10.37}
\end{equation*}
$$

We claim that for every $A \in M_{B}$,

$$
\begin{equation*}
\omega^{\mathrm{g}}\left(\left\{U \in \mathbb{T}_{1}^{d}: c_{1, i}\left(\mathbb{Z}^{d}+\widetilde{\mathrm{r}}_{i}(U)\right) A \cap \widetilde{\mathfrak{Z}}_{\xi}=\emptyset \text { for } i=1, \ldots, r_{1}\right\}\right) \geq \frac{1}{2} . \tag{10.38}
\end{equation*}
$$

To prove this, we first use that the left hand side of (10.38) is

$$
\begin{equation*}
\geq 1-\sum_{i=1}^{r_{1}} \omega^{\mathrm{g}}\left(\left\{U \in \mathbb{T}_{1}^{d}: c_{1, i}\left(\mathbb{Z}^{d}+\widetilde{\mathfrak{r}}_{i}(U)\right) A \cap \widetilde{\mathfrak{Z}}_{\xi} \neq \emptyset\right\}\right) . \tag{10.39}
\end{equation*}
$$

Now $A \in M_{B}$ implies that there is a vector $\boldsymbol{\ell} \in \mathbb{R}^{d}$ satisfying $\|\ell\| \geq B$ and $\mathbb{Z}^{d} A \subset \mathbb{Z} \ell+\ell^{\perp}$. It follows that for every $U \in \mathbb{T}_{1}^{d}$ we have $\left(\mathbb{Z}^{d}+\widetilde{\mathfrak{r}}_{i}(U)\right) A \subset(\mathbb{Z}+x) \boldsymbol{\ell}+\boldsymbol{\ell}^{\perp}$, where $x:=\|\boldsymbol{\ell}\|^{-2} \boldsymbol{\ell} \cdot(\boldsymbol{v} A)$ for some vector $\boldsymbol{v} \in \mathbb{Z}^{d}+\widetilde{\mathrm{r}}_{i}(U)$. Note that for every $\boldsymbol{m} \in \mathbb{Z}^{d}$ we have $\boldsymbol{m} A \in \mathbb{Z} \boldsymbol{\ell}+\boldsymbol{\ell}^{\perp}$ and so $\|\boldsymbol{\ell}\|^{-2} \boldsymbol{\ell} \cdot(\boldsymbol{m} A) \in \mathbb{Z}$; this implies that the map $\boldsymbol{v} \mapsto\|\boldsymbol{\ell}\|^{-2} \boldsymbol{\ell} \cdot(\boldsymbol{v} A)$ from $\mathbb{R}^{d}$ to $\mathbb{R}$ induces a well-defined map $\delta:(\mathbb{R} / \mathbb{Z})^{d} \rightarrow \mathbb{R} / \mathbb{Z}$. Using this map, the inclusion just mentioned can be written:

$$
\left(\mathbb{Z}^{d}+\widetilde{\mathrm{r}}_{i}(U)\right) A \subset\left(\mathbb{Z}+\delta\left(\widetilde{\mathrm{r}}_{i}(U)\right)\right) \ell+\ell^{\perp}, \quad \forall U \in \mathbb{T}_{1}^{d}
$$

(Here for any $y \in \mathbb{R} / \mathbb{Z}$, we write $\mathbb{Z}+y$ for the inverse image of $y$ in $\mathbb{R}$.) Hence we conclude that the expression in (10.39) is

$$
\begin{equation*}
\geq 1-\sum_{i=1}^{r_{1}}\left[\left(\delta \circ \widetilde{\mathfrak{r}}_{i}\right)_{*}\left(\omega^{\mathrm{g}}\right)\right]\left(\left\{y \in \mathbb{R} / \mathbb{Z}:(\mathbb{Z}+y) \ell+\boldsymbol{\ell}^{\perp} \cap c_{1, i}^{-1} \widetilde{\mathfrak{z}}_{\xi} \neq \emptyset\right\}\right) . \tag{10.40}
\end{equation*}
$$

But we have noted that $\widetilde{\mathfrak{r}}_{i *}\left(\omega^{\mathfrak{g}}\right)$ is Lebesgue measure on $(\mathbb{R} / \mathbb{Z})^{d}$; hence $\left(\delta \circ \widetilde{\mathrm{r}}_{i}\right)_{*}\left(\omega^{\mathfrak{g}}\right)$ is Lebesgue measure on $\mathbb{R} / \mathbb{Z}$. Note also that for any $x \in \mathbb{R}$, the set $c_{1, i}^{-1} \widetilde{\mathcal{Z}}_{\xi}$ intersects $x \boldsymbol{\ell}+\boldsymbol{\ell}^{\perp}$ if and only there is some $\boldsymbol{w} \in c_{1, i}^{-1} \widetilde{\boldsymbol{i}}_{\xi}$ satisfying $\|\boldsymbol{\ell}\|^{-2} \boldsymbol{\ell} \cdot \boldsymbol{w}=x$. Furthermore, because of our choice of $C_{\xi}$ in (10.37), we have $\left|\|\boldsymbol{\ell}\|^{-2} \boldsymbol{\ell} \cdot \boldsymbol{w}\right| \leq C_{\xi} /\left(c_{1, i}\|\ell\|\right)$ for every $\boldsymbol{w} \in c_{1, i}^{-1} \widetilde{\mathfrak{z}} \xi$. Therefore, the expression in (10.40) is

$$
\geq 1-\sum_{i=1}^{r_{1}} \operatorname{Leb}\left(\left\{y \in \mathbb{R} / \mathbb{Z}:(\mathbb{Z}+y) \cap\left[-\frac{C_{\xi}}{c_{1, i}\|\ell\|}, \frac{C_{\xi}}{c_{1, i}\|\ell\|}\right] \neq \emptyset\right\}\right)=1-\sum_{i=1}^{r_{1}} \frac{2 C_{\xi}}{c_{1, i}\|\ell\|} \geq \frac{1}{2}
$$

where the last bound holds since $\|\ell\| \geq B$ and because of our choice of $B$ in (10.37) (and the preceding equality holds since $C_{\xi} /\left(c_{1, i}\|\ell\|\right) \leq \frac{1}{4} \leq \frac{1}{2}$ for every $i$ ). This completes the proof of (10.38).

Using (10.38) and (10.36), we conclude

$$
\int_{\xi}^{\infty} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime} \geq \frac{1}{2} \nu\left(M_{B}\right) \gg B^{-d} \gg \xi^{-1}
$$

where we first used the bound from Lemma 10.5, and then the fact that $B \ll C_{\xi} \ll \xi^{1 / d}$. Hence the lower bound in (10.34) (for $N=1$ ) holds for an appropriate constant $c_{1}$, and Proposition 10.4 is proved.

Proof of Theorem [1.5. This is immediate from Proposition 10.4 and the fact that $\Phi_{\mathcal{P}}(\xi)$ is a nonnegative, decreasing function. Indeed, the upper bound for $\xi \geq 2$ follows using $(\xi / 2) \Phi_{\mathcal{P}}(\xi) \leq \int_{\xi / 2}^{\xi} \Phi_{\mathcal{P}}\left(\xi^{\prime}\right) d \xi^{\prime} ;$ for $1 \leq \xi \leq 2$ one may e.g. use $\Phi_{\mathcal{P}}(\xi) \leq \bar{n}_{\mathcal{P}} v_{d-1}$. For the lower bound, let $c_{1}, c_{2}$ be as in Proposition 10.4, and fix $a>1$ so that $c_{1}-c_{2} a^{-N}>0$. Then for any $\xi \geq 1$,

$$
(a-1) \xi \Phi_{\mathcal{P}}(\xi) \geq \int_{\xi}^{a \xi} \Phi_{\mathcal{P}}(x) d x=\int_{\xi}^{\infty} \Phi_{\mathcal{P}}(x) d x-\int_{a \xi}^{\infty} \Phi_{\mathcal{P}}(x) d x>c_{1} \xi^{-N}-c_{2}(a \xi)^{-N}
$$

which implies that the lower bound in (1.19) holds with $c_{2}^{\text {new }}=(a-1)^{-1}\left(c_{1}-c_{2} a^{-N}\right)$.

Proof of Corollary 1.6. By (1.20), $-\left(\bar{n}_{\mathcal{P}} v_{d-1}\right)^{-1} \Phi_{\mathcal{P}}(\xi)$ is a primitive function of $\bar{\Phi}_{\mathcal{P}}(\xi)$; hence by integration by parts,

$$
\int_{0}^{A} \xi^{2} \bar{\Phi}_{\mathcal{P}}(\xi)=\frac{1}{\bar{n}_{\mathcal{P}} v_{d-1}}\left(-A^{2} \Phi_{\mathcal{P}}(A)+2 \int_{0}^{A} \xi \Phi_{\mathcal{P}}(\xi) d \xi\right)
$$

for any $A>0$. Now the corollary follows by letting $A \rightarrow \infty$ and using Theorem 1.5 ,

## Index of notation

For the convenience of the reader, we include an index of some of the most important notation. Note that in Section 8, some of the notation listed below (for example, " $\Gamma$ ", " $\Gamma_{j}$ " and " $\widetilde{M}$ ") is used in a slightly different and more general way; this is explained in the beginning of that section.

| $\mathcal{B}_{r}^{d}$ | open ball in $\mathbb{R}^{d}$ with center $\mathbf{0}$ and radius $r$ | 2 |
| :---: | :---: | :---: |
| $\mathcal{B}_{r}^{d}(\boldsymbol{x})$ | open ball in $\mathbb{R}^{d}$ with center $\boldsymbol{x}$ and radius $r$ | 10 |
| $\mathrm{C}_{b}(S)$ | the space of bounded continuous functions on $S$ | 17 |
| $C_{\Psi}$ | family of equivalence classes in $\Psi$ | 6 |
| $\widetilde{\boldsymbol{c}}_{j}$ | the vector $\left(c_{j, 1}^{-1} \cdots c_{j, r_{j}}^{-1}\right)^{\top}$ in $\mathbb{R}^{r_{j}}$ | 22 |
| $\boldsymbol{c}_{j}^{\psi}$ | the vector $c_{\psi} \widetilde{\boldsymbol{c}}_{j}$ in $\mathbb{R}^{r_{j}}$ | 22 |
| $c_{\psi}$ | fixed positive real numbers such that (3.4)-(3.5) hold | 13 |
| $D_{\rho}$ | diagonal matrix $\operatorname{diag}\left(\rho^{d-1}, \rho^{-1}, \cdots, \rho^{-1}\right)$ | 9 |
| $\mathfrak{D}_{\mathcal{S}}$ | the set $\bigcup_{i<j}\left\{\left(M_{1}, \ldots, M_{N}\right) \in G^{\prime}: M_{i} M_{j}^{-1} \in \mathcal{S}\right\}$ | 47 |
| $\mathcal{E}$ | subset of $\mathcal{P}$ of density zero ("exceptional points") | 1040 |
| $\boldsymbol{e}_{k}$ | the $k$ th standard unit vector in $\mathbb{R}^{d}$ (or in $\mathbb{R}^{r}$ or $\mathbb{R}^{r_{j}}$ ) | (9) 2331 |
| $G$ | $\mathrm{S}_{r_{1}}(\mathbb{R}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{R})$ | 14 |
| $G_{j}$ | $\mathrm{S}_{\mathrm{r}_{j}}(\mathbb{R})$ | 15 |
| $G^{\prime}$ | $\mathrm{SL}_{d}(\mathbb{R})^{N}$ (a subgroup of $G$ ) | 46 |
| $g_{0}^{(\underline{q})}$ | $\mathrm{I}_{U^{(q)}} \widetilde{M}$ | 15 |
| $\mathrm{I}_{U}$ | for $U \in \mathrm{M}_{r \times d}(\mathbb{R}), \mathrm{I}_{U}:=(\mathrm{I}, U) \in \mathrm{S}_{r}(\mathbb{R})$; for $U \in \prod_{j=1}^{N} \mathrm{M}_{r_{j} \times d}(\mathbb{R}), \mathrm{I}_{U}:=\left(\mathrm{I}_{U_{1}}, \ldots, \mathrm{I}_{U_{N}}\right)$ | 14 15 |
| $I(\cdot)$ | indicator function: $I(P)=1$ if statement $P$ is true, otherwise $I(P)=0$ | 19 |
| $i_{\psi}$ | the second coordinate of $\psi$ (for any $\psi \in \Psi$ ) | 13 |
| $J_{0}$ | a map from $\mathbb{X}$ to $N_{s}\left(\mathbb{R}^{d}\right)$ | 15 |
| $J$ | a map from $\mathbb{X}$ to $N_{s}(\mathcal{X})$ | 34 |
| $J_{\psi}$ | a map from $\mathbb{X}^{\psi}$ to $N_{s}(\mathcal{X})$ | 34 |
| ${ }^{\prime} \psi$ | the first coordinate of $\psi$ (for any $\psi \in \Psi$ ) | 13 |
| $k$ | transition kernel, defined in (10.6)-10.7) | 68 |
| $k^{\text {g }}$ | transition kernel, defined in (10.9)-(10.10) | 68 |
| $\mathfrak{L}(S)$ | for $\emptyset \neq S \subset \mathbb{T}^{r_{j}}, \mathfrak{L}(S):={\overline{\left\langle\pi^{-1}(S)\right\rangle}}^{\circ}$ ( ( rational subspace of $\mathbb{R}^{r_{j}}$ ) | 17 |
| $\mathcal{L}_{\psi}$ | a grid in $\mathbb{R}^{d}$. After (3.4) we have $\mathcal{L}_{\psi}=c_{\psi}\left(\mathbb{Z}^{d}+\boldsymbol{w}_{\psi}\right) M_{j_{\psi}}$ | 44 |
| $L_{j}^{(q)}$ | $L_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$ | 16 18 |
| $L_{j}^{(V)}$ | $\mathfrak{L}\left(V_{1}, \ldots, V_{d}\right)$, for $V=\left(V_{1}, \ldots, V_{d}\right) \in \mathbb{T}_{j}^{d}$ | 18 |
| $L_{j}^{\psi}$ | the rational subspace of $\mathbb{R}^{r_{j}}$ given by (5.5) | 22 |
| $L_{j}$ | the rational subspace of $\mathbb{R}^{r_{j}}$ defined in (5.12) | 24 |
| m | probability measure on $\Psi$ or on $\Sigma$ | 5. 10 34 |
| $\mathrm{m}_{c}$ | probability measure on an equivalence class $c \subset \Psi$ | 7 |
| $M_{j}$ | $M_{1}, \ldots, M_{N}$ are fixed elements in $\mathrm{SL}_{d}(\mathbb{R})$ such that (3.4)-(3.5) hold | 13 |
| $\widetilde{M}$ | $\left(M_{1}, \ldots, M_{N}\right)$ (an element in $G^{\prime}$, thus in $G$ ) | 15 |
| $N(\mathcal{X})$ | the set of locally finite counting measures on $\mathcal{X}$ | 9 |
| $N_{s}(\mathcal{X})$ | the set of simple measures in $N(\mathcal{X})$ | 9 |
| $\bar{n}_{\mathcal{P}}$ | the asymptotic density of the point set $\mathcal{P} ; \bar{n}_{\mathcal{P}}=\sum_{\psi \in \Psi} \bar{n}_{\psi}$ | 2 [5] |
| $\bar{n}_{\psi}$ | the asymptotic density of $\mathcal{L}_{\psi}$. (After (3.4): $\bar{n}_{\psi}=c_{\psi}^{-d}$.) | 5 |
| $\mathcal{O}_{j}^{(\boldsymbol{q})}$ | $\mathcal{O}_{j}^{\left(\pi\left(U_{j}^{(q)}\right)\right)}$ | 18 |
| $\mathcal{O}_{j}{ }^{(V)}$ | the subset of $\left(\mathbb{S}_{j}^{(V)}\right)^{d}$ given by (4.7) | 18 |
| $\mathcal{O}_{j}^{\psi}$ | the subset of $\mathbb{T}_{j}^{d}$ given in (5.16) | 26 |
| $\mathcal{O}_{j}$ | the subset of $\mathbb{T}_{j}^{d}$ given in (5.18) | 27 |
| $\mathcal{P}$ | the scatterer configuration; a fixed union of grids in $\mathbb{R}^{d}$ | 3 |
| $\widetilde{\mathcal{P}}$ | $\{(\boldsymbol{p}, \varsigma(\boldsymbol{p})): \boldsymbol{p} \in \mathcal{P}\}$ (a subset of $\mathcal{X}$ ) | 10 |
| $\widetilde{\mathcal{P}}_{\boldsymbol{q}}$ | $\widetilde{\mathcal{P}} \backslash\{(\boldsymbol{q}, \varsigma(\boldsymbol{q}))\}$ if $\boldsymbol{q} \in \mathcal{P}$, otherwise $\widetilde{\mathcal{P}}$ | 10 |
| $\mathcal{P}_{T}(\rho)$ | the set $\mathcal{P} \cap \mathcal{B}_{T \rho^{1-d}}^{d} \backslash \mathcal{E}$ | 10 |
| $P(S)$ | the set of Borel probability measures on $S$ (for any topological space $S$ ) | 9 |
| $P_{\text {ac }}\left(\mathrm{S}_{1}^{d-1}\right)$ | the set of $\lambda \in P\left(\mathrm{~S}_{1}^{d-1}\right)$ which are absolutely continuous with respect to $\sigma$ | 10 |
| $P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ | the subset of $\mathrm{SL}_{d}(\mathbb{Z})$-invariant measures in $P\left(\mathbb{T}_{j}^{d}\right)$ | 18 |


| $p^{(\psi)}$ | collision kernel | 5 |
| :---: | :---: | :---: |
| $p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ | collision kernel | 5 |
| ${ }^{\circ} p^{(\psi)}$ | collision kernel for the scatterer configuration $\mathcal{P}_{[\psi]}$ | 5 |
| ${ }^{\circ} p^{\left(\psi^{\prime} \rightarrow \psi\right)}$ | collision kernel for the scatterer configuration $\mathcal{P}_{[\psi]}$ (here $\left.\left[\psi^{\prime}\right]=[\psi]\right)$ | 5 |
| $p_{j}$ | the projection map $G \rightarrow \mathrm{~S}_{r_{j}}(\mathbb{R})$ | 15 |
| $\widetilde{p}_{j}$ | either of the projection maps $\mathbb{X} \rightarrow \mathbb{X}_{j}$ or $\widetilde{\mathbb{T}} \rightarrow \mathbb{T}_{j}^{d}$ | [15, 47 |
| $p_{\psi}$ | for $\psi=(i, j) \in \Psi, p_{\psi}$ is the map $\mathrm{a}_{i} \circ p_{j}$ from $G$ to $\mathrm{ASL}_{d}(\mathbb{R})$ | 15 |
| $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ | $\left(\widetilde{\mathcal{P}}_{\boldsymbol{q}}-\boldsymbol{q}\right) R(\boldsymbol{v}) D_{\rho}$ | 10 |
| $R$ | a fixed map $\mathrm{S}_{1}^{d-1} \rightarrow \mathrm{SO}(d)$ such that $\boldsymbol{v} R(\boldsymbol{v})=\boldsymbol{e}_{1}, \forall \boldsymbol{v} \in \mathrm{~S}_{1}^{d-1}$ | 9 |
| $\mathrm{r}_{i}$ | the projection map $\mathrm{M}_{r \times d}(\mathbb{R}) \rightarrow \mathbb{R}^{d}$ which takes any matrix to its $i$ th row | 14 |
| $\widetilde{r}_{i}$ | for $1 \leq i \leq r_{j}, \widetilde{\mathbf{r}}_{i}$ is the projection map $\mathbb{T}_{j}^{d} \rightarrow(\mathbb{R} / \mathbb{Z})^{d}$ induced by $\mathrm{r}_{i}$ | 16 |
| $\mathcal{S}$ | the commenturator of $\mathrm{SL}_{d}(\mathbb{Z})$ in $\mathrm{SL}_{d}(\mathbb{R})$ | 13 |
| $\mathrm{S}_{r}(\mathbb{R})$ | $\mathrm{SL}_{d}(\mathbb{R}) \ltimes \mathrm{M}_{r \times d}(\mathbb{R})$ | 14 |
| $\mathrm{S}_{L}(\mathbb{R})$ | for $L$ a linear subspace of $\mathbb{R}^{r}, \mathrm{~S}_{L}(\mathbb{R}):=\mathrm{SL}_{d}(\mathbb{R}) \ltimes L^{d}$, a subgroup of $\mathrm{S}_{r}(\mathbb{R})$ | 16 |
| $\mathbb{S}_{j}^{(V)}$ | $\overline{\left\langle V_{1}, \ldots, V_{d}\right\rangle}$, a closed subgroup of $\mathbb{T}_{j}$ | 18 |
| $\mathbb{S}_{j}^{(\boldsymbol{q})}$ | $\mathbb{S}_{j}^{\left(\pi\left(U_{j}^{(\boldsymbol{q})}\right)\right)}$ | 18 |
| $\mathbb{S}_{j}^{\psi}$ | $\pi\left(L_{j}^{\psi}\right)$ | 26 |
| $\mathbb{S}_{j}$ | $\pi\left(L_{j}\right)$ | 27 |
| $\mathbb{T}_{j}$ | $(\mathbb{R} / \mathbb{Z})^{r_{j}}$ | 16 |
| $\mathbb{T}^{d}{ }^{d}$ | $\mathbb{T}_{j} \times \cdots \times \mathbb{T}_{j}=\mathrm{M}_{r_{j} \times d}(\mathbb{R} / \mathbb{Z})$ | 16 |
| $\widetilde{\mathbb{T}}$ | $\mathbb{T}_{1}^{d} \times \mathbb{T}_{2}^{d} \times \cdots \times \mathbb{T}_{N}^{d}$ | 47 |
| $U_{j}^{(\boldsymbol{q})}$ | the $r_{j} \times d$ matrix with row vectors $\boldsymbol{w}_{j, i}-c_{j, i}^{-1} \boldsymbol{q} M_{j}^{-1}\left(i=1, \ldots, r_{j}\right)$ | 15 |
| $U^{(q)}$ | $\left(U_{1}^{(\boldsymbol{q})}, \cdots, U_{N}^{(\boldsymbol{q})}\right)$ | 15 |
| $v_{d-1}$ | $\operatorname{vol}\left(\mathcal{B}_{1}^{d-1}\right)$ | 8 |
| $\boldsymbol{w}_{\psi}$ | fixed vectors in $\mathbb{R}^{d}$ such that (3.4)-(3.5) hold | 13 |
| $W_{j}$ | The $r_{j} \times d$ matrix with row vectors $\boldsymbol{w}_{j, i}\left(i=1, \ldots, r_{j}\right)$ | 22 |
| $\mathcal{X}$ | $\mathbb{R}^{d} \times \Sigma$ | 9 |
| $\mathbb{X}$ | $\Gamma \backslash G$ | 14 |
| $\mathbb{X}_{j}$ | $\Gamma_{j} \backslash G_{j}$ | 15 |
| $\mathbb{X}^{\psi}$ | $\left\{\Gamma g \in \mathbb{X}: g \in G, \mathbf{0} \in \mathbb{Z}^{d} p_{\psi}(g)\right\}$ | 32 |
| $x$ | the embedding $\mathbb{T}_{j}^{d} \rightarrow \mathbb{X}_{j}$ defined in (3.24), or the embedding $\widetilde{\mathbb{T}} \rightarrow \mathbb{X}$ | 16) 47 |
| $\mathfrak{Z}_{\xi}$ | The open cylinder $(0, \xi) \times \mathcal{B}_{1}^{d-1}$ | 10 |
| $z$ | map from $N_{s}(\mathcal{X})$ to $\Delta$ | 68 |
| $\Gamma$ | $\mathrm{S}_{r_{1}}(\mathbb{Z}) \times \cdots \times \mathrm{S}_{r_{N}}(\mathbb{Z})$ | 14 |
| $\Gamma_{j}$ | $\mathrm{S}_{r_{j}}(\mathbb{Z})$ | 15 |
| $\Delta$ | the set $\left(\mathfrak{Z}_{\infty} \times \Sigma\right) \sqcup\{$ undef $\}$ | 68 |
| $\iota$ | the projection $\mathrm{S}_{r}(\mathbb{R}) \rightarrow \mathrm{SL}_{d}(\mathbb{R})$, or the reflection map $\mathcal{X} \rightarrow \mathcal{X}$ | (14) 68 |
| $\tau$ | the projection $\mathbb{X}_{j} \rightarrow \mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})$ | 15 |
| $\mu_{\boldsymbol{q}, \rho}^{(\lambda)}$ | distribution of $\mathcal{Q}_{\rho}(\boldsymbol{q}, \boldsymbol{v})$ for $\boldsymbol{v}$ random in $\left(\mathrm{S}_{1}^{d-1}, \lambda\right)$ | 10 |
| $\mu_{\varsigma}$ | we have fixed a continuous map $\varsigma \mapsto \mu_{\varsigma}$ from $\Sigma$ to $P(N(\mathcal{X})$ ) | 10, 34 |
| $\mu_{\mathcal{X}}$ | $\mathrm{vol} \times \mathrm{m}$ | 10 |
| $\mu_{j}^{(\boldsymbol{q})}$ | measure in $P\left(\mathbb{X}_{j}\right)$ | 17 |
| $\mu^{(q)}$ | the measure $\mu_{1}^{(\boldsymbol{q})} \otimes \cdots \otimes \mu_{N}^{(\boldsymbol{q})}$ in $P(\mathbb{X})$ | 17 |
| $\mu^{\mathrm{g}}$ | the measure $J_{*}\left(\overline{\omega^{\mathrm{g}}}\right)$ on $N_{s}(\mathcal{X})$ | 45 |
| $\nu$ | Haar measure on $\mathrm{SL}_{d}(\mathbb{R})$, normalized by $\nu\left(\mathrm{SL}_{d}(\mathbb{Z}) \backslash \mathrm{SL}_{d}(\mathbb{R})\right)=1$ | 18 |
| $\Xi$ | random flight process | 4 |
| $\pi$ | either of the projection maps $G_{j} \rightarrow \mathbb{X}_{j}, G \rightarrow \mathbb{X}, \mathbb{R}^{r_{j}} \rightarrow \mathbb{T}_{j}$ or $\left(\mathbb{R}^{r_{j}}\right)^{d} \rightarrow \mathbb{T}_{j}^{d}$ | 15, 16 |
| $\Sigma$ | fixed compact metric space (the space of marks) | (9) 33 |
| $\sigma$ | $\operatorname{vol}_{\mathrm{S}_{1}^{d-1}}$, Lebesgue measure on $\mathrm{S}_{1}^{d-1}$ | 2 |
| $\sigma^{\psi}$ | $\left(\psi, \omega^{\psi}\right)($ an element in $\Sigma$ ) | 33 |
| $\varsigma$ | function $\mathcal{P} \rightarrow \Sigma$ fixed in (6.7) | 33 |
| $\Phi_{\mathcal{P}}(\xi)$ | free path length density | 7 |
| $\varphi$ | the diagonal embedding $\mathrm{SL}_{d}(\mathbb{R}) \rightarrow G$ | 17 |
| $\Psi$ | a set of indices. After (3.3) : $\Psi=\left\{(j, i): j \in\{1, \ldots, N\}, i \in\left\{1, \ldots, r_{j}\right\}\right\}$ | (4) 13 |
| $\widetilde{\Psi}$ | the set $\left\{\sigma^{\psi}: \psi \in \Psi\right\}$ (subset of $\Sigma$ ) | 34 |
| $\psi$ | function $\mathcal{P} \rightarrow \Psi$ fixed in (1.4). But $\psi$ is also used to denote a variable element in $\Psi$. | 5 |
| $\Omega$ | $\prod_{j=1}^{N} P\left(\mathbb{T}_{j}^{d}\right)^{\prime}$ | 33 |

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\(\omega_{j}^{(\boldsymbol{q})} \quad \omega_{j}^{\left(\pi\left(U_{j}(\boldsymbol{q})\right)\right)}\)
\(\omega_{j}^{(V)} \quad\) normalized restriction to \(\mathcal{O}_{j}(V)\) of the Haar measure on \(\mathbb{S}_{j}(V)^{d}\)
\(\omega_{j}^{\psi} \quad\) normalized restriction to \(\mathcal{O}_{j}^{\psi}\) of the Haar measure on \(\left(\widetilde{\mathbb{S}}_{j}^{\psi}\right)^{d} \quad 26\)
\(\omega_{j}^{\mathrm{g}} \quad\) normalized restriction to \(\mathcal{O}_{j}\) of the Haar measure on \(\left(\widetilde{\mathbb{S}}_{j}\right)^{d}\)
\(\omega^{(\boldsymbol{q})} \quad\left(\omega_{1}^{(\boldsymbol{q})}, \ldots, \omega_{N}^{(\boldsymbol{q})}\right)\) (an element in \(\Omega\) )
\(\omega^{(V)} \quad\left(\omega_{1}^{\left(V_{1}\right)}, \ldots, \omega_{N}^{\left(V_{N}\right)}\right)\) (an element in \(\Omega\) )
\(\omega^{\psi} \quad\left(\omega_{1}^{\psi}, \ldots, \omega_{N}^{\psi}\right)\) (an element in \(\Omega\) )
\(\omega^{\mathrm{g}} \quad\left(\omega_{1}^{\mathrm{g}}, \ldots, \omega_{N}^{\mathrm{g}}\right)\) (an element in \(\Omega\) )
\(\bar{\omega}\)
for \(\omega \in P\left(\mathbb{T}_{j}^{d}\right)^{\prime}, \bar{\omega}\) is the probability measure on \(\mathbb{X}_{j}\) defined below (4.8);
for \(\omega \in \Omega, \bar{\omega}\) is the probability measure on \(\mathbb{X}\) defined in 6.11)
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[^1]:    ${ }^{1}$ Two grids $\mathcal{L}$ and $\mathcal{L}^{\prime}$ in $\mathbb{R}^{d}$ are said to be commensurable if there exist $\delta>0$ and $\boldsymbol{v} \in \mathbb{R}^{d}$ such that $\mathcal{L} \cap\left(\delta \mathcal{L}^{\prime}+\boldsymbol{v}\right)$ is a grid.

[^2]:    ${ }^{2}$ See Section 2.1 regarding why we call it "crude".

[^3]:    ${ }^{3}$ For the point sets $\mathcal{P}$ which we consider in the present paper, the marking $\varsigma: \mathcal{P} \rightarrow \Sigma$ will be constructed to be a certain refinement of the crude marking $\psi: \mathcal{P} \rightarrow \Psi$ which we introduced in Section 1.1
    ${ }^{4}$ We say that a subset $\mathcal{E}$ of $\mathbb{R}^{d}$ has density zero if it has asymptotic density zero in the sense of (1.1); note that this holds if and only if it holds for $\mathcal{D}=\mathcal{B}_{1}^{d}$, viz., if and only if $\lim _{R \rightarrow \infty} R^{-d} \#\left(\mathcal{E} \cap \mathcal{B}_{R}^{d}\right)=0$.
    ${ }^{5}$ The statement in (2.7) means by definition: $\left[\forall \varepsilon>0: \exists \rho_{0}>0: \forall \rho \in\left(0, \rho_{0}\right): \forall \boldsymbol{q} \in \mathcal{P}_{T}(\rho): d\left(\mu_{\boldsymbol{q}, \rho}^{(\lambda)}, \mu_{\varsigma(\boldsymbol{q})}\right)<\varepsilon\right]$, where $d$ is some metric on $P(N(\mathcal{X}))$ realizing the weak topology. This definition is independent of the choice of $d$; indeed see [21, Lemma 2.1] and note that $\left\{\mu_{\varsigma}: \varsigma \in \Sigma\right\}$ is a compact subset of $P(N(\mathcal{X}))$, since it is the continuous image of the compact set $\Sigma$.

[^4]:    ${ }^{6}$ But note that for other spaces $\mathbb{R}^{r}$ appearing below, it will be natural to instead consider the vectors to be column vectors. The first instance of this is on p . 16. for the space $\mathbb{R}^{r_{j}}$ in the definition of the torus $\mathbb{T}_{j}=\mathbb{R}^{r_{j}} / \mathbb{Z}^{r_{j}}$.

[^5]:    ${ }^{7}$ Note: We view $\boldsymbol{a}$ as a column vector, or equivalently as an $r_{j} \times 1$ matrix; hence $\boldsymbol{a}^{\top}$ is a $1 \times r_{j}$ matrix.

[^6]:    ${ }^{8}$ In the notation of Lemma 2.2 this means: consider $p_{*}\left(\mu_{\boldsymbol{q}, \boldsymbol{\rho}}^{(\lambda)}\right)$ in place of $\mu_{\boldsymbol{q}, \rho}^{(\lambda)}$.
    ${ }^{9}$ The proof of $\mu_{\boldsymbol{q}, 0}=\left(J_{0, \psi(\boldsymbol{q})}\right)_{*}\left(\overline{\omega^{(\boldsymbol{q})}}\right)$ from Theorem4.2 is somewhat technical, and we omit it. (The details are similar to some of the arguments in Sec. 7.3 )
    ${ }^{10}$ This is an aposteriori fact; as far as we can see, such a choice would not lead to a simplification of the presentation of the proof of our main Theorem 2.1

[^7]:    ${ }^{11}$ After noticing that $\mu_{\sigma^{(1, \ell)}}$ is in fact supported on $N_{s}\left(\mathbb{R}^{d} \times \widetilde{\Psi}\right)$, and then identifying $\widetilde{\Psi}$ with $\left\{1, \ldots, r_{1}\right\}$ in the way we have just explained.

[^8]:    ${ }^{12}$ In 21 this measure is called " $\mu$ ".

[^9]:    ${ }^{13}$ This measure should really be normalized by a factor $\operatorname{vol}\left(\mathcal{B}_{R}^{d-1}\right)^{-1}$, to make $\lambda$ and $\nu_{m}$ probability measures; however such a normalizing factor clearly does not affect the validity of (8.24), and so we will ignore it.

[^10]:    ${ }^{14}$ To conform with the definition of asymptotic density in (1.1), it would be more natural to divide with $\operatorname{vol}\left(\mathcal{B}_{T}^{d}\right)$ instead of $T^{d}$ in (9.33); however using (9.33) makes some computations in the following slightly cleaner.

[^11]:    ${ }^{15}$ This is the map which we called " $\widetilde{\imath}$ " in (3.22); we now call it $\widetilde{\iota}_{j}$ for clarity.

[^12]:    ${ }^{16}$ See also (10.11) below.

[^13]:    ${ }^{17}$ This map $\boldsymbol{z}$ is the same map as in 21, (3.3)]; in particular, note that in [21, " $\Omega$ " denotes $\mathcal{B}_{1}^{d-1} \times \Sigma$, and hence " $\mathbb{R}_{>0} \times \Omega$ " can in an obvious way be identified with the set $\mathcal{Z}_{\infty} \times \Sigma$. Recall that in the present paper, " $\Omega$ " has a completely different meaning; see (6.3).

[^14]:    ${ }^{18}$ Note that the function " $\Phi(\xi, \boldsymbol{w})$ " in [15, 18] is the same as " $\Phi_{\boldsymbol{\alpha}}(\xi, \boldsymbol{w}, \boldsymbol{z})$ with $\boldsymbol{\alpha} \notin \mathbb{Q}^{d "}$ in 16.

