Invariant measures and the set of exceptions to Littlewood’s conjecture

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Abstract

We classify the measures on $\text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$ which are invariant and ergodic under the action of the group $A$ of positive diagonal matrices with positive entropy. We apply this to prove that the set of exceptions to Littlewood’s conjecture has Hausdorff dimension zero.

1. Introduction

1.1. Number theory and dynamics. There is a long and rich tradition of applying dynamical methods to number theory. In many of these applications, a key role is played by the space $\text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$ which can be identified as the space of unimodular lattices in $\mathbb{R}^k$. Any subgroup $H < \text{SL}(k, \mathbb{R})$ acts on this space in a natural way, and the dynamical properties of such actions often have deep number theoretical implications.

A significant landmark in this direction is the solution by G. A. Margulis [23] of the long-standing Oppenheim Conjecture through the study of the action of a certain subgroup $H$ on the space of unimodular lattices in three space. This conjecture, posed by A. Oppenheim in 1929, deals with density properties of the values of indefinite quadratic forms in three or more variables. So far there is no proof known of this result in its entirety which avoids the use of dynamics of homogeneous actions.

An important property of the acting group $H$ in the case of the Oppenheim Conjecture is that it is generated by unipotents: i.e. by elements of $\text{SL}(k, \mathbb{R})$ all of whose eigenvalues are 1. The dynamical result proved by Margulis was a special case of a conjecture of M. S. Raghunathan regarding the actions...
of general unipotents groups. This conjecture (and related conjectures made shortly thereafter) state that for the action of $H$ generated by unipotents by left translations on the homogeneous space $G/\Gamma$ of an arbitrary connected Lie group $G$ by a lattice $\Gamma$, the only possible $H$-orbit closures and $H$-ergodic probability measures are of an algebraic type. Raghunathan’s conjecture was proved in full generality by M. Ratner in a landmark series of papers ([41], [42] and others; see also the expository papers [40], [43], and the book [28]) which led to numerous applications; in particular, we use Ratner’s work heavily in this paper. Ratner’s theorems provide the model for the global orbit structure for systems with parabolic behavior. See [8] for a general discussion of principal types of orbit behavior in dynamics.

1.2. Weyl chamber flow and Diophantine approximation. In this paper we deal with a different homogeneous action, which is not so well understood, namely the action by left multiplication of the group $A$ of positive diagonal $k \times k$ matrices on $\text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$; $A$ is a split Cartan subgroup of $\text{SL}(k, \mathbb{R})$ and the action of $A$ is also known as a particular case of a Weyl chamber flow [16].

For $k = 2$ the acting group is isomorphic to $\mathbb{R}$ and the Weyl chamber flow reduces to the geodesic flow on a surface of constant negative curvature, namely the modular surface. This flow has hyperbolic structure; it is Anosov if one makes minor allowances for noncompactness and elliptic points. The orbit structure of such flows is well understood; in particular there is a great variety of invariant ergodic measures and orbit closures. For $k > 2$, the Weyl chamber flow is hyperbolic as an $\mathbb{R}^{k-1}$-action, i.e. transversally to the orbits. Such actions are very different from Anosov flows and display many rigidity properties; see e.g. [16], [15]. One of the manifestations of rigidity concerns invariant measures. Notice that one–parameter subgroups of the Weyl chamber flow are partially hyperbolic and each such subgroup still has many invariant measures. However, it is conjectured that $A$-ergodic measures are rare:

**Conjecture 1.1** (Margulis). Let $\mu$ be an $A$-invariant and ergodic probability measure on $X = \text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$ for $k \geq 3$. Then $\mu$ is algebraic; i.e. there is a closed, connected group $L > A$ so that $\mu$ is the $L$-invariant measure on a single, closed $L$-orbit.

This conjecture is a special case of much more general conjectures in this direction by Margulis [25], and by A. Katok and R. Spatzier [17]. This type of behavior was first observed by Furstenberg [6] for the action of the multiplicative semigroup $\Sigma_{m,n} = \{m^n l^k\}_{k,l \geq 1}$ on $\mathbb{R}/\mathbb{Z}$, where $n, m$ are two multiplicatively independent integers (i.e. not powers of the same integer), and the action is given by $k.x = kx \text{ mod } 1$ for any $k \in \Sigma_{m,n}$ and $x \in \mathbb{R}/\mathbb{Z}$. Under these assumptions Furstenberg proved that the only infinite closed invariant
set under the action of this semigroup is the space $\mathbb{R}/\mathbb{Z}$ itself. He also raised the question of extensions, in particular to the measure theoretic analog as well as to the locally homogeneous context.

There is an intrinsic difference regarding the classification of invariant measures between Weyl chamber flows (e.g. higher rank Cartan actions) and unipotent actions. For unipotent actions, every element of the action already acts in a rigid manner. For Cartan actions, there is no rigidity for the action of individual elements, but only for the full action. In stark contrast to unipotent actions, M. Rees [44], [3, §9] has shown there are lattices $\Gamma < \text{SL}(k, \mathbb{R})$ for which there are nonalgebraic $A$-invariant and ergodic probability measures on $X = \text{SL}(k, \mathbb{R})/\Gamma$ (fortunately, this does not happen for $\Gamma = \text{SL}(k, \mathbb{Z})$, see [21], [25] and more generally [48] for related results). These nonalgebraic measures arise precisely because one-parameter subactions are not rigid, and come from $A$ invariant homogeneous subspaces which have algebraic factors on which the action degenerates to a one-parameter action.

While Conjecture 1.1 is a special case of the general question about the structure of invariant measures for higher rank hyperbolic homogeneous actions, it is of particular interest in view of number theoretic consequences. In particular, it implies the following well-known and long-standing conjecture of Littlewood [24, §2]:

**Conjecture 1.2** (Littlewood (c. 1930)). *For every $u, v \in \mathbb{R}$,*

\[
\liminf_{n \to \infty} n \langle nu \rangle \langle nv \rangle = 0,
\]

where $\langle w \rangle = \min_{n \in \mathbb{Z}} |w - n|$ is the distance of $w \in \mathbb{R}$ to the nearest integer.

In this paper we prove the following partial result towards Conjecture 1.1 which has implications toward Littlewood's conjecture:

**Theorem 1.3.** Let $\mu$ be an $A$-invariant and ergodic measure on $X = \text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$ for $k \geq 3$. Assume that there is some one-parameter subgroup of $A$ which acts on $X$ with positive entropy. Then $\mu$ is algebraic.

In [21] a complete classification of the possible algebraic $\mu$ is given. In particular, we have the following:

**Corollary 1.4.** Let $\mu$ be as in Theorem 1.3. Then $\mu$ is not compactly supported. Furthermore, if $k$ is prime, $\mu$ is the unique $\text{SL}(k, \mathbb{R})$-invariant measure on $X$.

Theorem 1.3 and its corollary have the following implication toward Littlewood's conjecture:
Theorem 1.5. Let
\[ \Xi = \left\{ (u,v) \in \mathbb{R}^2 : \liminf_{n \to \infty} n\langle nu \rangle \langle nv \rangle > 0 \right\}. \]
Then the Hausdorff dimension \( \dim_H \Xi = 0 \). In fact, \( \Xi \) is a countable union of compact sets with box dimension zero.

J. W. S. Cassels and H. P. F. Swinnerton-Dyer [1] showed that (1.1) holds for any \( u, v \) which are from the same cubic number field (i.e. any field \( K \) with degree \( [K : \mathbb{Q}] = 3 \)).

It is easy to see that for a.e. \( (u, v) \) equation (1.1) holds — indeed, for almost every \( u \) it is already true that \( \liminf_{n \to \infty} n\langle nu \rangle = 0 \). However, there is a set of \( u \) of Hausdorff dimension 1 for which \( \liminf_{n \to \infty} n\langle nu \rangle > 0 \); such \( u \) are said to be badly approximable. Pollington and Velani [35] showed that for every \( u \in \mathbb{R} \), the intersection of the set
\[ \{ v \in \mathbb{R} : (u, v) \text{ satisfies (1.1)} \} \]
with the set of badly approximable numbers has Hausdorff dimension one. Note that this fact is an immediate corollary of our Theorem 1.5 — indeed, Theorem 1.5 implies in particular that the complement of this set (1.2) has Hausdorff dimension zero for all \( u \). We remark that the proof of Pollington and Velani is effective.

Littlewood’s conjecture is a special case of a more general question. More generally, for any \( k \) linear forms \( m_i(x_1, x_2, \ldots, x_k) = \sum_{j=1}^{k} m_{ij} x_j \), one may consider the product
\[ f_m(x_1, x_2, \ldots, x_k) = \prod_{i=1}^{k} m_i(x_1, \ldots, x_k), \]
where \( m = (m_{ij}) \) denotes the \( k \times k \) matrix whose rows are the linear forms above. Using Theorem 1.3 we prove the following:

Theorem 1.6. There is a set \( \Xi_k \subset \text{SL}(k, \mathbb{R}) \) of Hausdorff dimension \( k-1 \) so that for every \( m \in \text{SL}(k, \mathbb{R}) \setminus \Xi_k \),
\[ \inf_{x \in \mathbb{Z}^k \setminus \{0\}} |f_m(x)| = 0. \]
Indeed, this set \( \Xi_k \) is \( A \)-invariant, and has zero Hausdorff dimension transversally to the \( A \)-orbits.

For more details, see Section 10 and Section 11. Note that (1.3) is automatically satisfied if zero is attained by \( f_m \) evaluated on \( \mathbb{Z}^k \setminus \{0\} \).

We also want to mention another application of our results due to Hee Oh [32], which is related to the following conjecture of Margulis:
Conjecture 1.7 (Margulis, 1993). Let $G$ be the product of $n \geq 2$ copies of $\text{SL}(2, \mathbb{R})$,

$$U_1 = \left\{ \begin{pmatrix} 1 & * \\ 0 & 1 \end{pmatrix} \times \cdots \times \begin{pmatrix} 1 & * \\ 0 & 1 \end{pmatrix} \right\}$$

and

$$U_2 = \left\{ \begin{pmatrix} 1 & 0 \\ * & 1 \end{pmatrix} \times \cdots \times \begin{pmatrix} 1 & 0 \\ * & 1 \end{pmatrix} \right\}.$$  

Let $\Gamma < G$ be a discrete subgroup so that for both $i = 1$ and 2, $\Gamma \cap U_i$ is a lattice in $U_i$ and for any proper connected normal subgroup $N < G$ the intersection $\Gamma \cap N \cap U_i$ is trivial. Then $\Gamma$ is commensurable with a Hilbert modular lattice up to conjunction in $\text{GL}(2, \mathbb{R}) \times \cdots \times \text{GL}(2, \mathbb{R})$.

Hee Oh [33] has shown that assuming a topological analog to Conjecture 1.1 (which is implied by Conjecture 1.1), Conjecture 1.7 is true for $n \geq 3$. As explained in [32] (and following directly from [33, Thm. 1.5]), our result, Theorem 1.3, implies the following weaker result (also for $n \geq 3$): consider the set $\mathcal{D}$ of possible intersections $\Gamma \cap U_1$ for $\Gamma$ as in Conjecture 1.7, which is a subset of the space of lattices in $U_1$. This set $\mathcal{D}$ is clearly invariant under conjugation by the diagonal group in $\text{GL}(2, \mathbb{R}) \times \cdots \times \text{GL}(2, \mathbb{R})$; Theorem 1.3 (or more precisely Theorem 10.2 which we prove using Theorem 1.3 in §10) implies that the set $\mathcal{D}$ has zero Hausdorff dimension transversally to the orbit of this $n$-dimensional group (in particular, this set $\mathcal{D}$ has Hausdorff dimension $n$; see Section 7 and Section 10 for more details regarding Hausdorff dimension and transversals, and [33], [32] for more details regarding this application).

1.3. Measure rigidity. The earliest results for measure rigidity for higher rank hyperbolic actions deal with the Furstenberg problem: [22], [45], [12]. Specifically, Rudolph [45] and Johnson [12] proved that if $\mu$ is a probability measure invariant and ergodic under the action of the semigroup generated by $\times m, \times n$ (again with $m, n$ not powers of the same integer), and if some element of this semigroup acts with positive entropy, then $\mu$ is Lebesgue.

When Rudolph’s result appeared, the second author suggested another test model for the measure rigidity: two commuting hyperbolic automorphisms of the three-dimensional torus. Since Rudolph’s proof seemed, at least superficially, too closely related to symbolic dynamics, jointly with R. Spatzier, a more geometric technique was developed. This allowed a unified treatment of essentially all the classical examples of higher rank actions for which rigidity of measures is expected [17], [13], and in retrospect, Rudolph’s proof can also be interpreted in this framework.

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1 For a definition of Hilbert modular lattices, see [33].
This method (as well as most later work on measure rigidity for these higher rank abelian actions) is based on the study of conditional measures induced by a given invariant measure $\mu$ on certain invariant foliations. The foliations considered include stable and unstable foliations of various elements of the actions, as well as intersections of such foliations, and are related to the Lyapunov exponents of the action. For Weyl chamber flows these foliations are given by orbits of unipotent subgroups normalized by the action.

Unless there is an element of the action which acts with positive entropy with respect to $\mu$, these conditional measures are well-known to be $\delta$-measure supported on a single point, and do not reveal any additional meaningful information about $\mu$. Hence this and later techniques are limited to study actions where at least one element has positive entropy. Under ideal situations, such as the original motivating case of two commuting hyperbolic automorphisms of the three torus, no further assumptions are needed, and a result entirely analogous to Rudolph’s theorem can be proved using the method of [17].

However, for Weyl chamber flows, an additional assumption is needed for the [17] proof to work. This assumption is satisfied, for example, if the flow along every singular direction in the Weyl chamber is ergodic (though a weaker hypothesis is sufficient). This additional assumption, which unlike the entropy assumption is not stable under weak* limits, precludes applying the results from [17] in many cases.

Recently, two new methods of proofs were developed, which overcome this difficulty.

The first method was developed by the first and second authors [3], following an idea mentioned at the end of [17]. This idea uses the noncommutativity of the above-mentioned foliations (or more precisely, of the corresponding unipotent groups). This paper deals with general $\mathbb{R}$-split semisimple Lie groups; in particular it is shown there that if $\mu$ is an $A$-invariant measure on $X = \text{SL}(k, \mathbb{R})/\Gamma$, and if the entropies of $\mu$ with respect to all one-parameter groups are positive, then $\mu$ is the Haar measure. It should be noted that for this method the properties of the lattice do not play any role, and indeed this is true not only for $\Gamma = \text{SL}(k, \mathbb{Z})$ but for every discrete subgroup $\Gamma$. An extension to the nonsplit case appeared in [4]. Using the methods we present in the second part of the present paper, the results of [3] can be used to show that the set of exceptions to Littlewood’s conjecture has Hausdorff dimension at most 1.

A different approach was developed by the third author, and was used to prove a special case of the quantum unique ergodicity conjecture [20]. In its basic form, this conjecture is related to the geodesic flow, which is not rigid, so in order to be able to prove quantum unique ergodicity in certain situations a more general setup for measure rigidity, following Host [9], was needed. A special case of the main theorem of [20] is the following: Let $A$ be an $\mathbb{R}$-split
Cartan subgroup of $\text{SL}(2, \mathbb{R}) \times \text{SL}(2, \mathbb{R})$. Any $A$-ergodic measure on $\text{SL}(2, \mathbb{R}) \times \text{SL}(2, \mathbb{R})/\Gamma$ for which some one-parameter subgroup of $A$ acts with positive entropy is algebraic. Here $\Gamma$ is e.g. an irreducible lattice in $\text{SL}(2, \mathbb{R}) \times \text{SL}(2, \mathbb{R})$. Since the foliations under consideration in this case do commute, the methods of [3] are not applicable.

The method of [20] can be adapted to quotients of more general groups, and in particular to $\text{SL}(k, \mathbb{R})$. It is noteworthy (and gratifying) that for the space of lattices (and more general quotients of $\text{SL}(k, \mathbb{R})$) these two unrelated methods are completely complementary: measures with “high” entropy (e.g. measures for which many one-parameter subgroup have positive entropy) can be handled with the methods of [3], and measures with “low” (but positive) entropy can be handled using the methods of [20]. Together, these methods give Theorem 1.3 (as well as the more general Theorem 2.1 below for more general quotients).

The method of proof in [20], an adaptation of which we use here, is based on study of the behavior of $\mu$ along certain unipotent trajectories, using techniques introduced by Ratner in [39], [38] to study unipotent flows, in particular the H-property (these techniques are nicely exposed in Section 1.5 of [28]). This is surprising because the techniques are applied on a measure $\mu$ which is a priori not even quasi-invariant under these (or any other) unipotent flows.

In showing that the high entropy and low entropy cases are complementary we use a variant on the Ledrappier-Young entropy formula [19]. Such use is one of the simplifying ideas in G. Tomanov and Margulis’ alternative proof of Ratner’s theorem [26].

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Part I. Measure rigidity

Throughout this paper, let $G = \text{SL}(k, \mathbb{R})$ for some $k \geq 3$, let $\Gamma$ be a discrete subgroup of $G$, and let $X = G/\Gamma$. As in the previous section, we let $A < G$ denote the group of $k \times k$ positive diagonal matrices. We shall implicitly
identify
\[ \Sigma = \{ t \in \mathbb{R}^k : t_1 + \cdots + t_k = 0 \} \]
and the Lie algebra of \( A \) via the map \((t_1, \ldots, t_k) \mapsto \text{diag}(t_1, \ldots, t_k)\). We write \( \alpha^t = \text{diag}(e^{t_1}, \ldots, e^{t_k}) \in A \) and also \( \alpha^t \) for the left multiplication by this element on \( X \). This defines an \( \mathbb{R}^{k-1} \) flow \( \alpha \) on \( X \).

A subgroup \( U < G \) is \textit{unipotent} if for every \( g \in U \), \( g - I_k \) is nilpotent; i.e., for some \( n \), \( (g - I_k)^n = 0 \). A group \( H \) is said to be \textit{normalized} by \( g \in G \) if \( gHg^{-1} = H \); \( H \) is normalized by \( L < G \) if it is normalized by every \( g \in L \); and the \textit{normalizer} \( N(H) \) of \( H \) is the group of all \( g \in G \) normalizing it. Similarly, \( g \) \textit{centralizes} \( H \) if \( gh = hg \) for every \( h \in H \), and we set \( C(H) \), the \textit{centralizer} of \( H \) in \( G \), to be the group of all \( g \in G \) centralizing \( H \).

If \( U < G \) is normalized by \( A \) then for every \( x \in X \) and \( a \in A \), \( a(Ux) = U(ax) \), so that the \textit{foliation} of \( X \) into \( U \) orbits is invariant under the action of \( A \). We will say that \( a \in A \) expands \( U \) if all eigenvalues of \( \text{Ad}(a) \) restricted to the Lie algebra of \( U \) are greater than one.

For any locally compact metric space \( Y \) let \( \mathcal{M}_\infty(Y) \) denote the space of Radon measures on \( Y \) equipped with the weak\(^*\) topology, i.e. all locally finite Borel measures on \( Y \) with the coarsest topology for which \( \rho \mapsto \int_Y f(y)d\rho(y) \) is continuous for every compactly supported continuous \( f \). For two Radon measures \( \nu_1 \) and \( \nu_2 \) on \( Y \) we write
\[ \nu_1 \propto \nu_2 \text{ if } \nu_1 = C\nu_2 \text{ for some } C > 0, \]
and say that \( \nu_1 \) and \( \nu_2 \) are proportional.

We let \( B^Y_\varepsilon(y) \) (or \( B_\varepsilon(y) \) if \( Y \) is understood) denote the ball of radius \( \varepsilon \) around \( y \in Y \); if \( H \) is a group we set \( B^H_\varepsilon = B^H_\varepsilon(I) \) where \( I \) is identity in \( H \); and if \( H \) acts on \( X \) and \( x \in X \) we let \( B^H_\varepsilon(x) = B^H_\varepsilon \cdot x \).

Let \( d(\cdot, \cdot) \) be the geodesic distance induced by a right-invariant Riemannian metric on \( G \). This metric on \( G \) induces a right-invariant metric on every closed subgroup \( H \subset G \), and furthermore a metric on \( X = G/\Gamma \). These induced metrics we denote by the same letter.

2. Conditional measures on \( A \)-invariant foliations, invariant measures, and shearing

2.1. \textit{Conditional measures}. A basic construction, which was introduced in the context of measure rigidity in [17] (and in a sense is already used implicitly in [45]), is the restriction of probability or even Radon measures on a foliated space to the leaves of this foliation. A discussion can be found in [17, §4], and a fairly general construction is presented in [20, §3]. Below we consider special cases of this general construction, summarizing its main properties.
Let \( \mu \) be an \( A \)-invariant probability measure on \( X \). For any unipotent subgroup \( U < G \) normalized by \( A \), one has a system \( \{ \mu_{x,U} \}_{x \in X} \) of Radon measures on \( U \) and a co-null set \( X' \subset X \) with the following properties\(^2\):

1. The map \( x \mapsto \mu_{x,U} \) is measurable.
2. For every \( \varepsilon > 0 \) and \( x \in X' \), \( \mu_{x,U}(B^U_\varepsilon) > 0 \).
3. For every \( x \in X' \) and \( u \in U \) with \( ux \in X' \), we have that \( \mu_{x,U} \propto (\mu_{ux,U})u \), where \( (\mu_{ux,U})u \) denotes the push forward of the measure \( \mu_{ux,U} \) under the map \( v \mapsto vu \).
4. For every \( t \in \Sigma \), and \( x, \alpha^t x \in X' \), \( \mu_{\alpha^t x,U} \propto \alpha^t(\mu_{x,U})\alpha^{-t} \).

In general, there is no canonical way to normalize the measures \( \mu_{x,U} \); we fix a specific normalization by requiring that \( \mu_{x,U}(B^U_1) = 1 \) for every \( x \in X' \). This implies the next crucial property.

5. If \( U \subset C(\alpha^t) = \{ g \in G : g\alpha^t = \alpha^t g \} \) commutes with \( \alpha^t \), then \( \mu_{\alpha^t x,U} = \mu_{x,U} \) whenever \( x, \alpha^t x \in X' \).

6. \( \mu \) is \( U \)-invariant if, and only if, \( \mu_{x,U} \) is a Haar measure on \( U \) a.e. (see e.g. [17] or the slightly more general [20, Prop. 4.3]).

The other extreme to \( U \)-invariance occurs when \( \mu_{x,U} \) is atomic. If \( \mu \) is \( A \)-invariant then outside some set of measure zero if \( \mu_{x,U} \) is atomic then it is supported on the identity \( I_k \in U \), in which case we say that \( \mu_{x,U} \) is trivial. This follows from Poincaré recurrence for an element \( a \in A \) that uniformly expands the \( U \)-orbits (i.e. for which the \( U \)-orbits are contained in the unstable manifolds). Since the set of \( x \in X \) for which \( \mu_{x,U} \) is trivial is \( A \)-invariant, if \( \mu \) is \( A \)-ergodic then either \( \mu_{x,U} \) is trivial a.s. or \( \mu_{x,U} \) is nonatomic a.s. Fundamental to us is the following characterization of positive entropy (see [26, §9] and [17]):

7. If for every \( x \in X \) the orbit \( Ux \) is the stable manifold through \( x \) with respect to \( \alpha^t \), then the measure theoretic entropy \( h_\mu(\alpha^t) \) is positive if and only if the conditional measures \( \mu_{x,U} \) are nonatomic a.e.

So positive entropy implies that the conditional measures are nontrivial a.e., and the goal is to show that this implies that they are Haar measures. Quite often one shows first that the conditional measures are translation invariant under some element up to proportionality, which makes the following observation useful.

\(^2\)We are following the conventions of [20] in viewing the conditional measures \( \mu_{x,U} \) as measures on \( U \). An alternative approach, which, for example, is the one taken in [17] and [13], is to view the conditional measures as a collection of measures on \( X \) supported on single orbits of \( U \); in this approach, however, the conditional measure is not a Radon measure on \( X \), only on the single orbit of \( U \) in the topology of this submanifold.
Possibly after replacing $X'$ of (1)–(4) by a conull subset, we see that for any $x \in X'$ and any $u \in U$ with $\mu_{x,U} \propto \mu_{x,U}u$, in fact, $\mu_{x,U} = \mu_{x,U}u$ holds.

This was first shown in [17]. The proof of this fact only uses Poincaré recurrence and (4) above; for completeness we provide a proof below.

**Proof of (8).** Let $t$ be such that $\alpha^t$ uniformly contracts the $U$-leaves (i.e.

for every $x$ the $U$-orbit $Ux$ is part of the stable manifold with respect to $\alpha^t$).

Define for $M > 0$

$$D_M = \{ x \in X' : \mu_{x,U}(B_2^U) < M \}.$$  

We claim that for every $x \in X' \cap \bigcup_M \limsup_{n \to \infty} \alpha^{-nt} D_M$ (i.e.

any $x \in X'$ so that $\alpha^{-nt}$ is in $D_M$ for some $M$ for infinitely many $n$) if $\mu_{x,U} = c\mu_{x,U}u$ then $c \leq 1$.

Indeed, suppose $x \in X' \cap \bigcup_M \limsup_{n \to \infty} \alpha^{-nt} D_M$ and $u \in U$ satisfy $\mu_{x,U} = c\mu_{x,U}u$. Then for any $n,k$

$$\mu_{\alpha^{nt}x,U} = c_k \mu_{\alpha^{nt}x,U}(\alpha^{nt}u^k \alpha^{-nt}).$$

Choose $k > 1$ arbitrary. Suppose $n$ is such that $\alpha^{nt}x \in D_M$ and suppose that $n$ is sufficiently large that $\alpha^{nt}u^k \alpha^{-nt} \in B_1^U$, which is possible since $\alpha^t$ uniformly contracts $U$. Then

$$M \geq \mu_{\alpha^{nt}x,U}(B_1^U) \geq \mu_{\alpha^{nt}x,U}(B_1^U \alpha^{nt}u^k \alpha^{-nt})$$

$$= (\mu_{\alpha^{nt}x,U} \alpha^{nt}u^{-k} \alpha^{-nt})(B_1^U)$$

$$= c_k \mu_{\alpha^{nt}x,U}(B_1^U) = c_k.$$

Since $k$ is arbitrary this implies $c \leq 1$.

If $\mu_{x,U} = c\mu_{x,U}u$ then $\mu_{x,U} = c^{-1} \mu_{x,U}u^{-1}$, so the above argument applied to $u^{-1}$ shows that $c \geq 1$, hence $\mu_{x,U} = \mu_{x,U}u$.

Thus we see that if we replace $X'$ by $X' \cap \bigcup_M \limsup_{n \to \infty} \alpha^{-nt} D_M$ — a conull subset of $X'$, then (8) holds for any $x \in X'$.

Of particular importance to us will be the following one-parameter unipotent subgroups of $G$, which are parametrized by pairs $(i,j)$ of distinct integers in the range $\{1, \ldots, k\}$:

$$u_{ij}(s) = \exp(sE_{ij}) = I_k + sE_{ij}, \quad U_{ij} = \{u_{ij}(s) : s \in \mathbb{R}\},$$

where $E_{ij}$ denotes the matrix with 1 at the $i^{\text{th}}$ row and $j^{\text{th}}$ column and zero everywhere else. It is easy to see that these groups are normalized by $A$; indeed, for $t = (t_1, \ldots, t_k) \in \Sigma$

$$\alpha^t u_{ij}(s) \alpha^{-t} = u_{ij}(e^{t_1-s_1} \ldots e^{t_k-s_k}).$$
Since these groups are normalized by $A$, the orbits of $U_{ij}$ form an $A$-invariant foliation of $X = \text{SL}(k, \mathbb{R})/\Gamma$ with one-dimensional leaves. We will use $\mu_x^{ij}$ as a shorthand for $\mu_{x,U_{ij}}$; any integer $i \in \{1, \ldots, k\}$ will be called an index; and unless otherwise stated, any pair $i, j$ of indices is implicitly assumed to be distinct.

Note that for the conditional measures $\mu_x^{ij}$ it is easy to find a nonzero $t \in \Sigma$ such that (5) above holds; for this all we need is $t_i = t_j$. Another helpful feature is the one-dimensionality of $U_{ij}$ which also helps to show that $\mu_x^{ij}$ are a.e. Haar measures. In particular we have the following:

(9) Suppose there exists a set of positive measure $B \subset X$ such that for any $x \in B$ there exists a nonzero $u \in U_{ij}$ with $\mu_x^{ij} \propto \mu_x^{ij} u$. Then for a.e. $x \in B$ in fact $\mu_x^{ij}$ is a Haar measure of $U_{ij}$, and if $\alpha$ is ergodic then $\mu$ is invariant under $U_{ij}$.

Proof of (9). Recall first that by (8) we can assume $\mu_x^{ij} = \mu_x^{ij} u$ for $x \in B$. Let $K \subset B$ be a compact set of measure almost equal to $\mu(B)$ such that $\mu_x^{ij}$ is continuous for $x \in K$. It is possible to find such a $K$ by Luzin’s theorem. Note however, that here the target space is the space of Radon measures $\mathcal{M}_\infty(U_{ij})$ equipped with the weak* topology so that a more general version [5, p. 69] of Luzin’s theorem is needed. Let $t \in \Sigma$ be such that $U_{ij}$ is uniformly contracted by $\alpha^t$. Suppose now $x \in K$ satisfies Poincaré recurrence for every neighborhood of $x$ relative to $K$. Then there is a sequence $x_\ell = \alpha^{n_\ell}t \in K$ that approaches $x$ with $n_\ell \to \infty$. Invariance of $\mu_x^{ij}$ under $u$ implies invariance of $\mu_{x_\ell}^{ij}$ under the much smaller element $\alpha^{n_\ell}t u \alpha^{-n_\ell}t$ and all its powers. However, since $\mu_x^{ij}$ converges to $\mu_x^{ij}$ we conclude that $\mu_x^{ij}$ is a Haar measure of $U_{ij}$. The final statement follows from (4) which implies that the set of $x$ where $\mu_x^{ij}$ is a Haar measure is $\alpha$-invariant. \qed

Even when $\mu$ is not invariant under $U_{ij}$ we still have the following maximal ergodic theorem [20, Thm. A.1] proved by the last named author in joint work with D. Rudolph, which is related to a maximal ergodic theorem of Hurewicz [11].

(10) For any $f \in L^1(X,\mu)$ and $\alpha > 0$, 

$$\mu \left( \left\{ x : \int_{B_r^{U_{ij}}} f(ux) \, d\mu_x^{ij} > \alpha \mu_x^{ij}(B_r^{U_{ij}}) \text{ for some } r > 0 \right\} \right) < \frac{C \|f\|_1}{\alpha} $$

for some universal constant $C > 0$.

2.2. Invariant measures, high and low entropy cases. We are now in a position to state the general measure rigidity result for quotients of $G$:

Theorem 2.1. Let $X = G/\Gamma$ and $A$ be as above. Let $\mu$ be an $A$-invariant and ergodic probability measure on $X$. For any pair of indices $a, b$, one of the following three properties must hold.
The conditional measures \( \mu_{ab}^x \) and \( \mu_{ba}^x \) are trivial a.e.

The conditional measures \( \mu_{ab}^x \) and \( \mu_{ba}^x \) are Haar a.e., and \( \mu \) is invariant under left multiplication with elements of \( H_{ab} = \langle U_{ab}, U_{ba} \rangle \).

Let \( A'_{ab} = \{ \alpha^s : s \in \Sigma \text{ and } s_a = s_b \} \). Then a.e. ergodic component of \( \mu \) with respect to \( A'_{ab} \) is supported on a single \( C(H_{ab}) \)-orbit, where \( C(H_{ab}) = \{ g \in G : gh = hg \text{ for all } h \in H_{ab} \} \) is the centralizer of \( H_{ab} \).

Remark. If \( k = 3 \) then (3) is equivalent to the following:

(3') There exist a nontrivial \( s \in \Sigma \) with \( s_a = s_b \) and a point \( x_0 \in X \) with \( \alpha^s x_0 = x_0 \) such that the measure \( \mu \) is supported by the orbit of \( x_0 \) under \( C(A'_{ab}) \). In particular, a.e. point \( x \) satisfies \( \alpha^s x = x \).

Indeed, in this case \( C(H_{ab}) \) contains only diagonal matrices, and Poincaré recurrence for \( A'_{ab} \) together with (3) imply that a.e. point is periodic under \( A'_{ab} \). However, ergodicity of \( \mu \) under \( A \) implies that the period \( s \) must be the same a.e. Let \( x_0 \in X \) be such that every neighborhood of \( x_0 \) has positive measure. Then \( x \) close to \( x_0 \) is fixed under \( \alpha^s \) only if \( x \in C(A'_{ab})x_0 \), and ergodicity shows (3'). The examples of M. Rees [44], [3, §9] of nonalgebraic \( A \)-ergodic measures in certain quotients of \( SL(3, \mathbb{R}) \) (which certainly can have positive entropy) are precisely of this form, and show that case (3) and (3') above are not superfluous.

When \( \Gamma = SL(k, \mathbb{Z}) \), however, this phenomenon, which we term exceptional returns, does not happen. We will show this in Section 5; similar observations have been made earlier in [25], [21]. We also refer the reader to [48] for a treatment of similar questions for inner lattices in \( SL(k, \mathbb{R}) \) (a certain class of lattices in \( SL(k, \mathbb{R}) \)).

The conditional measures \( \mu_{ij}^x \) are intimately connected with the entropy. More precisely, \( \mu \) has positive entropy with respect to \( \alpha^t \) if and only if for some \( i, j \) with \( t_i > t_j \) the measures \( \mu_{ij}^x \) are not a.s. trivial (see Proposition 3.1 below for more details; this fact was first proved in [17]). Thus (1) in Theorem 2.1 above holds for all pairs of indices \( i, j \) if, and only if, the entropy of \( \mu \) with respect to every one-parameter subgroup of \( A \) is zero.

In order to prove Theorem 2.1, it is enough to show that for every \( a, b \) for which the \( \mu_{ab}^x \) is a.s. nontrivial either Theorem 2.1.(2) or Theorem 2.1.(3) holds. For each pair of indices \( a, b \), our proof is divided into two cases which we loosely refer to as the high entropy and the low entropy case:

High entropy case. There is an additional pair of indices \( i, j \) distinct from \( a, b \) such that \( i = a \) or \( j = b \) for which \( \mu_{ij}^x \) are nontrivial a.s. In this case we prove:
**Theorem 2.2.** If both $\mu_{x}^{ab}$ and $\mu_{x}^{ij}$ are nontrivial a.s., for distinct pairs of indices $i, j$ and $a, b$ with either $i = a$ or $j = b$, then both $\mu_{x}^{ab}$ and $\mu_{x}^{ba}$ are in fact Haar measures a.s. and $\mu$ is invariant under $H_{ab}$.

The proof in this case, presented in Section 3 makes use of the noncommutative structure of certain unipotent subgroups of $G$, and follows [3] closely. However, by careful use of an adaptation of a formula of Ledrappier and Young (Proposition 3.1 below) relating entropy to the conditional measures $\mu_{x}^{ab}$ we are able to extract some additional information. It is interesting to note that Margulis and Tomanov used the Ledrappier-Young theory for a similar purpose in [26], simplifying some of Ratner’s original arguments in the classification of measures invariant under the action of unipotent groups.

**Low entropy case.** For every pair of indices $i, j$ distinct from $a, b$ such that $i = a$ or $j = b$, $\mu_{x}^{ij}$ are trivial a.s. In this case there are two possibilities:

**Theorem 2.3.** Assume $\mu_{x}^{ab}$ are a.e. nontrivial, and $\mu_{x}^{ij}$ are trivial a.e. for every pair $i, j$ distinct from $a, b$ such that $i = a$ or $j = b$. Then one of the following properties holds.

1. $\mu$ is $U_{ab}$ invariant.
2. Almost every $A_{ab}$-ergodic component of $\mu$ is supported on a single $C(H_{ab})$ orbit.

We will see in Corollary 3.4 that in the low entropy case $\mu_{x}^{ba}$ is also nontrivial; so applying Theorem 2.3 for $U_{ba}$ instead of $U_{ab}$ one sees that either $\mu$ is $H_{ab}$-invariant or almost every $A_{ab}$-ergodic component of $\mu$ is supported on a single $C(H_{ab}) = C(H_{ba})$ orbit.

In this case we employ the techniques developed by the third named author in [20]. There, one considers invariant measures on irreducible quotients of products of the type $SL(2, \mathbb{R}) \times L$ for some algebraic group $L$. Essentially, one tries to prove a Ratner type result (using methods quite similar to Ratner’s [38], [39]) for the $U_{ab}$ flow even though $\mu$ is not assumed to be invariant or even quasi invariant under $U_{ab}$. Implicitly in the proof we use a variant of Ratner’s H-property (related, but distinct from the one used by Witte in [29, §6]) together with the maximal ergodic theorem for $U_{ab}$ as in (9) in Section 2.1.

**3. More about entropy and the high entropy case**

A well-known theorem by Ledrappier and Young [19] relates the entropy, the dimension of conditional measures along invariant foliations, and Lyapunov exponents, for a general $C^2$ map on a compact manifold, and in [26, §9] an adaptation of the general results to flows on locally homogeneous spaces is
provided. In the general context, the formula giving the entropy in terms of the dimensions of conditional measures along invariant foliations requires consideration of a sequence of subfoliations, starting from the foliation of the manifold into stable leaves. However, because the measure $\mu$ is invariant under the full $A$-action one can relate the entropy to the conditional measures on the one-dimensional foliations into orbits of $U_{ij}$ for all pairs of indices $i, j$.

We quote the following from [3]; in that paper, this proposition is deduced from the fine structure of the conditional measures on full stable leaves for $A$-invariant measure; however, it can also be deduced from a more general result of Hu regarding properties of commuting diffeomorphisms [10]. It should be noted that the constants $s_{ij}(\mu)$ that appear below have explicit interpretation in terms of the pointwise dimension of $\mu_{ij}^{x}$ [19].

**Proposition 3.1 ([3, Lemma 6.2]).** Let $\mu$ be an $A$-invariant and ergodic probability measure on $X = G/\Gamma$ with $G = \text{SL}(k, \mathbb{R})$ and $\Gamma < G$ discrete. Then for any pair of indices $i, j$ there are constants $s_{ij}(\mu) \in [0, 1]$ so that:

1. $s_{ij}(\mu) = 0$ if and only if for a.e. $x$, $\mu_{ij}^{x}$ are atomic and supported on a single point.
2. If a.s. $\mu_{ij}^{x}$ are Haar (i.e. $\mu$ is $U_{ij}$ invariant), then $s_{ij}(\mu) = 1$.
3. For any $t \in \Sigma$,

   \begin{equation}
   h_{\mu}(\alpha^{t}) = \sum_{i,j} s_{ij}(\mu)(t_{i} - t_{j})^{+}.
   \end{equation}

Here $(r)^{+} = \max(0, r)$ denotes the positive part of $r \in \mathbb{R}$.

We note that the converse to (2) is also true. A similar proposition holds for more general semisimple groups $G$. In particular we get the following (which is also proved in a somewhat different way in [17]):

**Corollary 3.2.** For any $t \in \Sigma$, the entropy $h_{\mu}(\alpha^{t})$ is positive if and only if there is a pair of indices $i, j$ with $t_{i} - t_{j} > 0$ for which $\mu_{ij}^{x}$ are nontrivial a.s.

A basic property of the entropy is that for any $t \in \Sigma$,

\begin{equation}
 h_{\mu}(\alpha^{t}) = h_{\mu}(\alpha^{-t}).
\end{equation}

As we will see this gives nontrivial identities between the $s_{ij}(\mu)$.

The following is a key lemma from [3]; see Figure 1.

**Lemma 3.3 ([3, Lemma 6.1]).** Suppose $\mu$ is an $A$-invariant and ergodic probability measure, $i, j, k$ distinct indices such that both $\mu_{ij}^{x}$ and $\mu_{ik}^{x}$ are nonatomic a.e. Then $\mu$ is $U_{ik}$-invariant.
Figure 1: One key ingredient of the proof of Lemma 3.3 in [3] is the translation produced along $U_{ik}$ when going along $U_{ij}$ and $U_{jk}$ and returning to the same leaf $U_{ik}x$.

Proof of Theorem 2.2. For $\ell = a, b$ we define the sets

$$C_\ell = \{i \in \{1, \ldots, k\} \setminus \{a, b\} : s_{i\ell}(\mu) > 0\},$$
$$R_\ell = \{j \in \{1, \ldots, k\} \setminus \{a, b\} : s_{\ell j}(\mu) > 0\},$$
$$C^L_\ell = \{i \in \{1, \ldots, k\} \setminus \{a, b\} : \mu \text{ is } U_{i\ell}\text{-invariant}\},$$
$$R^L_\ell = \{j \in \{1, \ldots, k\} \setminus \{a, b\} : \mu \text{ is } U_{\ell j}\text{-invariant}\}.$$

Suppose $i \in C_a$; then the conditional measures $\mu^{ia}_x$ are nontrivial a.e. by Proposition 3.1. Since by assumption $\mu^{ab}_x$ are nontrivial a.e., Lemma 3.3 shows that $\mu^{ib}_x$ are Lebesgue a.e. This shows that $C_a \subset C^L_b$, and $R_b \subset R^L_a$ follows similarly.

Let $t = (t_1, \ldots, t_k)$ with $t_i = -1/k$ for $i \neq a$ and $t_a = 1 - 1/k$. For the following expression set $s_{aa} = 0$. By Proposition 3.1 the entropy of $\alpha^t$ equals

$$h_\mu(\alpha^t) = s_{a1}(\mu) + \cdots + s_{ak}(\mu)$$
$$= s_{ab}(\mu) + |R^L_a| + \sum_{j \in R_a \setminus R^L_a} s_{aj}(\mu) > |R^L_a|,$$

where we used our assumption that $s_{ab}(\mu) > 0$. Applying Proposition 3.1 for $\alpha^{-t}$ we see similarly that

$$h_\mu(\alpha^{-t}) = s_{1a}(\mu) + \cdots + s_{ka}(\mu) = s_{ba}(\mu) + \sum_{i \in C_a} s_{ia}(\mu) \leq (1 + |C_a|),$$

where we used the fact that $s_{ia}(\mu) \in [0, 1]$ for $a = 2, \ldots, k$. However, since the entropies of $\alpha^t$ and of $\alpha^{-t}$ are equal, we get $|R^L_a| \leq |C_a|$.

Using $t' = (t'_1, \ldots, t'_k)$ with $t'_i = -1/k$ for $i \neq b$ and $t'_b = 1 - 1/k$ instead of $t$ in the above paragraph shows similarly $|C^L_b| \leq |R_b|$. Recall that $C_a \subset C^L_b$ and $R_b \subset R^L_a$. Combining these inequalities we conclude that

$$|R^L_a| \leq |C_a| \leq |C^L_b| \leq |R_b| \leq |R^L_a|.$$
and so all of these sets have the same cardinality. However, from (3.3) and (3.4) we see that \( s_{ab}(\mu) + |R_a^L| \leq h_\mu(\alpha^t) \leq s_{ba}(\mu) + |C_a| \). Together we see that

\[
(3.5) \quad s_{ba}(\mu) \geq s_{ab}(\mu) > 0.
\]

From this we conclude as before that \( C_a \subset C_b^L \subset C_a^L \), and so \( C_a = C_a^L \).

Similarly, one sees \( R_b = R_b^L \).

This shows that if \( s_{ab}(\mu) > 0 \) and \( s_{ij}(\mu) > 0 \) for some other pair \( i, j \) with either \( i = a \) or \( j = b \), then in fact \( \mu \) is \( U_{ij} \)-invariant. If there was at least one such pair of indices \( i, j \) we could apply the previous argument to \( i, j \) instead of \( a, b \) and get that \( \mu \) is \( U_{ab} \)-invariant.

In particular, we have seen in the proof of Theorem 2.2 that \( s_{ab} > 0 \) implies (3.5). We conclude the following symmetry.

**Corollary 3.4.** For any pair of indices \( (a, b) \), \( s_{ab} = s_{ba} \). In particular, \( \mu_{x}^{ab} \) are nontrivial a.s., if and only if, \( \mu_{x}^{ba} \) are nontrivial a.s.

### 4. The low entropy case

We let \( A_{ab}' = \{ \alpha^s \in A : s_a = s_b \} \), and let \( \alpha^s \in A_{ab}' \). Then \( \alpha^s \) commutes with \( U_{ab} \), which implies that \( \mu_{x}^{ab} = \mu_{x}^{\alpha^s} \) a.e.

For a given pair of indices \( a, b \), we define the following subgroups of \( G \):

\[
L_{(ab)} = C(U_{ab}),
\]

\[
U_{(ab)} = \langle U_{ij} : i = a \text{ or } j = b \rangle,
\]

\[
C_{(ab)} = C(H_{ab}) = C(U_{ab}) \cap C(U_{ba}).
\]

Recall that the metric on \( X \) is induced by a right-invariant metric on \( G \). So for every two \( x, y \in X \) there exists a \( g \in G \) with \( y = gx \) and \( d(x, y) = d(I_k, g) \).

#### 4.1. Exceptional returns.

**Definition 4.1.** We say for \( K \subset X \) that the \( A_{ab}' \)-returns to \( K \) are exceptional (strong exceptional) if there exists a \( \delta > 0 \) so that for all \( x, x' \in K \), and \( \alpha^s \in A_{ab}' \) with \( x' = \alpha^s x \in B_\delta(x) \cap K \) every \( g \in B_\delta^G \) with \( x' = gx \) satisfies \( g \in L_{(ab)} \) \((g \in C_{(ab)}\) respectively).  

**Lemma 4.2.** There exists a null set \( N \subset X \) such that for any compact \( K \subset X \setminus N \) with exceptional \( A_{ab}' \)-returns to \( K \) the \( A_{ab}' \)-returns to \( K \) are in fact strong exceptional.

**Proof.** To simplify notation, we may assume without loss of generality that \( a = 1, b = 2 \), and write \( A', U, L, C \) for \( A'_{12}, U_{(12)}, L_{(12)}, C_{(12)} \) respectively. We
Clearly \( g \subset N \) for all such rational \( \alpha \). Note that condition (1) below is exactly Theorem 2.3(2).

For \( \ell \geq 1 \) let \( D_\ell \) be the set of \( x \in X \) with the property that for all \( z \in B_{1/\ell}(x) \) there exists a unique \( g \in B^C_{1/\ell} \) with \( z = gx \). Note that \( \bigcup_{\ell=1}^\infty D_\ell = X \), and that for every compact set, \( K \subset D_\ell \) for some \( \ell > 0 \).

Let first \( \alpha^s \in A' \) be a fixed element, and let \( E_{\ell,s} \subset D_\ell \) be the set of points \( x \) for which \( x' = \alpha^s x \in B_{1/\ell}(x) \) and \( x' = gx \) with \( g \in B^C_{1/\ell} \cap L = B^L_{1/\ell} \). Since \( g \in B^C_{1/\ell} \) is uniquely determined by \( x \) (for a fixed \( s \)), we can define (in the notation of (4.1)) the measurable function

\[
f(x) = \max(|g_{12}|, |g_{1*}|, |g_{2*}|) \text{ for } x \in E_{\ell,s}.
\]

Let \( t = (-1,1,0,\ldots,0) \in \Sigma \). Then conjugation with \( \alpha^t \) contracts \( U \). In fact for \( g \) as in (4.1) the entries of \( \alpha^t g \alpha^{-t} \) corresponding to \( g_{12}, g_{1*} \) and \( g_{2*} \), are \( e^{-2}g_{12}, e^{-1}g_{1*} \) and \( e^{-1}g_{2*} \), and those corresponding to \( g_{21}, g_{s1} \) and \( g_{s2} \) are \( e^2g_{21}, eg_{s1} \) and \( eg_{s2} \). Notice that the latter are assumed to be zero. This shows that for \( x \in E_{\ell,s} \) and \( \alpha^{-nt} x \in D_\ell \), in fact \( \alpha^{-nt} x \in E_{\ell,s} \). Furthermore \( f(\alpha^{-nt} x) \leq e^{-n} f(x) \). Poincaré recurrence shows that \( f(x) = 0 \) for a.e. \( x \in E_{\ell,s} \) or equivalently \( \alpha^s x \in B^C_{1/\ell}(x) \) for a.e. \( x \in D_\ell \) with \( \alpha^s x \in B^L_{1/\ell}(x) \).

Varying \( s \) over all elements of \( \Sigma \) with rational coordinates and \( \alpha^s \in A' \), we arrive at a nullset \( N_\ell \subset D_\ell \) so that \( \alpha^s x \in B^L_{1/\ell}(x) \) implies \( \alpha^s x \in B^C_{1/\ell}(x) \) for all such rational \( s \). Let \( N \) be the union of \( N_\ell \) for \( \ell = 1, 2, \ldots \). We claim that \( N \) satisfies the lemma.

So suppose \( K \subset X \setminus N \) has \( A' \)-exceptional returns. Choose \( \ell \geq 1 \) so that \( K \subset D_\ell \), and furthermore so that \( \delta = 1/\ell \) can be used in the definition of \( A' \)-exceptional returns to \( K \). Let \( x \in K \), \( x' = \alpha^s x \in B_{1/\ell}(x) \) for some \( s \in \Sigma \) with \( \alpha^s \in A' \), and \( g \in B^C_{1/\ell} \) with \( x' = gx \). By assumption on \( K \), we have that \( g \in L \). Choose a rational \( \tilde{s} \in \Sigma \) close to \( s \) with \( \alpha^{\tilde{s}} \in A' \) so that \( \alpha^{\tilde{s}} x \in B_{1/\ell}(x) \). Clearly \( \tilde{g} = \alpha^{\tilde{s}-s} g \) satisfies \( \alpha^{\tilde{s}} x = \tilde{g} x \) and so \( \tilde{g} \in B^L_{1/\ell} \). Since \( x \in K \subset D_{1/\ell} \setminus N_{1/\ell} \), it follows that \( \tilde{g} \in C \). Going back to \( x' = \alpha^s x \) and \( g \) it follows that \( g \in C \).

Our interest in exceptional returns is explained by the following proposition. Note that condition (1) below is exactly Theorem 2.3(2).
Proposition 4.3. For any pair of indices a, b the following two conditions are equivalent.

(1) A.e. ergodic component of \(\mu\) with respect to \(A'_{ab}\) is supported on a single \(C_{(ab)}\)-orbit.

(2) For every \(\varepsilon > 0\) there exists a compact set \(K\) with measure \(\mu(K) > 1 - \varepsilon\) so that the \(A'_{ab}\)-returns to \(K\) are strong exceptional.

The ergodic decomposition of \(\mu\) with respect to \(A'_{ab}\) can be constructed in the following manner: Let \(\mathcal{E}'\) denote the \(\sigma\)-algebra of Borel sets which are \(A'_{ab}\) invariant. For technical purposes, we use the fact that \((X, \mathcal{B}_X, \mu)\) is a Lebesgue space to replace \(\mathcal{E}'\) by an equivalent countably generated sub-\(\sigma\)-algebra \(\mathcal{E}\). Let \(\mu^E_x\) be the family of conditional measures of \(\mu\) with respect to the \(\sigma\)-algebra \(\mathcal{E}\). Since \(\mathcal{E}\) is countably generated the atom \([x]_\mathcal{E}\) is well defined for all \(x\), and it can be arranged that for all \(x\) and \(y\) with \(y \in [x]_\mathcal{E}\) the conditional measures \(\mu^E_x = \mu^E_y\), and that for all \(x\), \(\mu^E_x\) is a probability measure.

Since \(\mathcal{E}\) consists of \(A'_{ab}\)-invariant sets, a.e. conditional measure is \(A'_{ab}\)-invariant, and can be shown to be ergodic. So the decomposition of \(\mu\) into conditionals

\[
\mu = \int_X \mu^E_x \, d\mu
\]
gives the ergodic decomposition of \(\mu\) with respect to \(A_{ab}^\prime\).

Proof. For simplicity, we write \(A' = A_{ab}^\prime\) and \(C = C_{(ab)}\).

(1) \(\Rightarrow\) (2). Suppose a.e. \(A'\) ergodic component is supported on a single \(C\)-orbit. Let \(\varepsilon > 0\). For any fixed \(r > 0\) we define

\[
f_r(x) = \mu^E_x(B_r^C(x)).
\]

By the assumption \(f_r(x) \not\rightarrow 1\) for \(r \rightarrow \infty\) and a.e. \(x\). Therefore, there exists a fixed \(r > 0\) with \(\mu(C_r) > 1 - \varepsilon\), where \(C_r = \{x : f_r(x) > 1/2\}\).

Fix some \(x \in X\). We claim that for every small enough \(\delta > 0\)

\[
B_{2r}^C(x) \cap B_\delta(x) = B_\delta^C(x).
\]

Indeed, by the choice of the metric on \(X\) there exists \(\delta' > 0\) so that the map \(g \mapsto gx\) from \(B_{3\delta'}^C\) to \(X\) is an isometry. Every \(g \in B_2^C\) satisfies that either \(d(B_\delta^C(x), B_\delta^C(gx)) > 0\), or that there exists \(h \in B_\delta^C\) with \(hx \in B_\delta^C(gx)\). In the latter case \(B_\delta^C(gx) \subset B_{3\delta'}^C(x)\). The sets \(B_\delta^C(g)\) for \(g \in B_{2r}^C\) cover the compact set \(B_{2r}^C\). Taking a finite subcover, we find some \(\eta > 0\) so that \(d(gx, x) > \eta\) or \(gx \in B_{3\delta'}^C(x)\) for every \(g \in B_{2r}^C\). It follows that (4.3) holds with \(\delta = \min(\eta, \delta')\).

In other words, \(C_r = \bigcup_{\delta > 0} D_\delta\), where

\[
D_\delta = \{x \in C_r : B_{2r}^C(x) \cap B_\delta(x) \subset B_\delta^C(x)\},
\]

and there exists \(\delta > 0\) with \(\mu(D_\delta) > 1 - \varepsilon\).
Let $K \subset D_\delta$ be compact. We claim that the $A'$-returns to $K$ are strongly exceptional. So suppose $x \in K$ and $x' = \alpha^s x \in K$ for some $\alpha^s \in A'$. Then since $x$ and $x'$ are in the same atom of $\mathcal{E}$, the conditional measures satisfy $\mu^E_x = \mu^E_{x'}$. By definition of $C_r$ we have $\mu^E_x(B^C_r(x)) > 1/2$ and the same for $x'$. Therefore $B^C_r(x)$ and $B^C_r(x')$ cannot be disjoint, and $x' \in B^C_r(\delta x')$ follows. By definition of $D_\delta$ it follows that $x' \in B^C_r(x)$. Thus the $A'$-returns to $K$ are indeed strongly exceptional.

(2) $\implies$ (1). Suppose that for every $\ell \geq 1$ there exists a compact set $K_\ell$ with $\mu(K_\ell) > 1 - 1/\ell$ so that the $A'$-returns to $K$ are strong exceptional. Then $N = X \setminus \bigcup K_\ell$ is a nullset. It suffices to show that (1) holds for every $A'$ ergodic $\mu^E_x$ which satisfies $\mu^E_x(N) = 0$.

For any such $x$ there exists $\ell > 0$ with $\mu^E_x(K_\ell) > 0$. Choose some $z \in K_\ell$ with $\mu^E_x(B_{1/m}(z) \cap K_\ell) > 0$ for all $m \geq 1$. We claim that $\mu^E_x$ is supported on $C_z$, i.e. that $\mu^E_x(C_z) = 1$. Let $\delta$ be as in the definition of strong exceptional returns. By ergodicity there exists for $\mu^E_x$-a.e. $y_0 \in X$ some $\alpha^s \in A'$ with $y_1 = \alpha^s y_0 \in B_\delta(z) \cap K_\ell$. Moreover, there exists a sequence $y_n \in A' y_0 \cap K_\ell$ with $y_n \to z$. Since $y_n \in B_\delta(y_1)$ for large enough $n$ and since the $A'$-returns to $K_\ell$ are strong exceptional, we conclude that $y_n \in B^C_\delta(y_1)$.

LEMMA 4.4. (1) Under the assumptions of the low entropy case (i.e. $s_{ab}(\mu) > 0$ but $s_{ij}(\mu) = 0$ for all $i, j$ with either $i = a$ or $j = b$), there exists a $\mu$-nullset $N \subset X$ such that for $x \in X \setminus N$,

\[ U_{(ab)} x \cap X \setminus N \subset U_{ab} x. \]

(2) Furthermore, unless $\mu$ is $U_{ab}$-invariant, it can be arranged that

\[ \mu^a_x \neq \mu^a_y \]

for any $x \in X \setminus N$ and any $y \in U_{(ab)} x \setminus N$ which is different from $x$.

Proof. Set $U = U_{(ab)}$ and let $\mu_{x,U}$ be the conditional measures for the foliation into $U$-orbits. By [3, Prop. 8.3] the conditional measure $\mu_{x,U}$ is a.e. – say for $x \notin N$ – a product measure of the conditional measures $\mu^a_{ij}$ over all $i, j$ for which $U_{ij} \subset U$. Clearly, by the assumptions of the low entropy case, $\mu^a_{ab}$ is the only one of these which is nontrivial. Therefore, $\mu_{x,U}$ as a measure on $U$ is supported on the one-dimensional group $U_{ab}$.

By (3) in Section 2.1 the conditional measures satisfy furthermore that there is a null set – enlarge $N$ accordingly – such that for $x, y \notin N$ and $y = u x \in U x$ the conditionals $\mu_{x,U}$ and $\mu_{y,U}$ satisfy that $\mu_{x,U} \propto \mu_{y,U} u$. However, since $\mu_{x,U}$ and $\mu_{y,U}$ are both supported by $U_{ab}$, it follows that $u \in U_{ab}$. This shows Lemma 4.4.(1).
In order to show Lemma 4.4.(2), we note that we already know that \( y \in U_{ab}x \). So if \( \mu_{ab}^y = \mu_{ab}^x \), then \( \mu_{ab}^x \) is again, by (3) in Section 2.1, invariant (up to proportionality) under multiplication by some nontrivial \( u \in U_{ab} \). If this were to happen on a set of positive measure, then by (9) in Section 2.1, \( \mu_{ab}^x \) are in fact Haar a.e. – a contradiction to our assumption.

### 4.2. Sketch of proof of Theorem 2.3.

We assume that the two equivalent conditions in Proposition 4.3 fail (the first of which is precisely the condition of Theorem 2.3 (2)). From this we will deduce that \( \mu \) is \( U_{ab} \)-invariant which is precisely the statement in Theorem 2.3 (1).

For the following we may assume without loss of generality that \( a = 1 \) and \( b = 2 \). Write \( A' \) and \( u(r) = I_k + r E_{12} \in U_{12} \) for \( r \in \mathbb{R} \) instead of \( A_{12}' \) and \( u_{12}(r) \). Also, we shall at times implicitly identify \( \mu_{12}^x \) (which is a measure on \( U_{12} \)) with its push forward under the map \( u(r) \mapsto r \), e.g. write \( \mu_{12}^x([a,b]) \) instead of \( \mu_{12}^x(u([a,b])) \).

By Poincaré recurrence we have for a.e. \( x \in X \) and every \( \delta > 0 \) that

\[
d(\alpha^s x, x) < \delta \text{ for some large } \alpha^s \in A'.
\]

For a small enough \( \delta \) there exists a unique \( g \in B_\delta^{U_{12}} \) such that \( x' = \alpha^s = gx \).

Since \( \alpha^s \) preserves the measure and since \( A' \subset L_{12} = C(U_{12}) \) the conditional measures satisfy

\[
\mu_{12}^x = \mu_{12}^{x'}.
\]

by (5) in Section 2.1. Since \( \mu_{12}^x \) is nontrivial, we can find many \( r \in \mathbb{R} \) so that \( x(r) = u(r)x \) and \( x'(r) = u(r)x' \) are again typical. By (3) in Section 2.1 the conditionals satisfy

\[
\mu_{x(r)} u(r) \propto \mu_{12}^x
\]

and similarly for \( x'(r) \) and \( x' \). Together with (4.4) and the way we have normalized the conditional measures this implies that

\[
\mu_{x(r)}^{12} = \mu_{x'(r)}^{12}.
\]

The key to the low entropy argument, and this is also the key to Ratner’s seminal work on rigidity of unipotent flows, is how the unipotent orbits \( x(r) \) and \( x'(r) \) diverge for \( r \) large (see Figure 2). Ratner’s H-property (which was introduced and used in her earlier works on rigidity of unipotent flows [38], [39] and was generalized by D. Morris-Witte in [29]) says that this divergence occurs only gradually and in prescribed directions. We remark that in addition to our use of the H-property, the general outline of our argument for the low entropy case is also quite similar to [38], [39].

We shall use a variant of this H-property in our paper, which at its heart is the following simple matrix calculation (cf. [38, Lemma 2.1] and [39, Def. 1]).
Let the entries of $g \in B_\delta^G$ be labelled as in (4.1). A simple calculation shows that

$$x'(r) = g(r)x(r)$$

for

$$g(r) = u(r)gu(-r) = \begin{pmatrix} a_1 + g_{21}r & g_{12} + (a_2 - a_1)r - g_{21}r^2 & g_{1*} + g_{2*}r \\ g_{21} & a_2 - g_{21}r & g_{2*} \\ g_{*1} & g_{*2} - g_{*1}r & a_* \end{pmatrix}.$$  

Since the return is not exceptional, $g \notin L_{12} = C(U_{12})$ and one of the following holds; $a_2 - a_1 \neq 0$, $g_{21} \neq 0$, $g_{*1} \neq 0$, or $g_{2*} \neq 0$. From this it is immediate that there exists some $r$ so that $g(r)$ is close to $I_k$ in all entries except at least one entry corresponding to the subgroup $U_{(12)}$. More precisely, there is an absolute constant $C$ so that there exists $r$ with

$$C^{-1} \leq \max(|(a_2 - a_1)r - g_{21}r^2|, \|g_{2*}r\|, \|g_{*1}r\|) \leq C,$$

(4.7)  

$$|g_{21}r| \leq C\delta^{3/8}.\tag{4.8}$$

With some care we will arrange it so that $x(r), x'(r)$ belong to a fixed compact set $X_1 \subset X \setminus N$. Here $N$ is as in Lemma 4.4 and $X_1$ satisfies that $\mu_{z}^{12}$ depends continuously on $z \in X_1$, which is possible by Luzin’s theorem.

If we can indeed find for every $\delta > 0$ two such points $x(r), x'(r)$ with (4.7) and (4.8), we let $\delta$ go to zero and conclude from compactness that there are two different points $y, y' \in X_1$ with $y' \in U_{(12)}y$ which are limits of a sequence of points $x(r), x'(r) \in X_1$. By continuity of $\mu_z^{12}$ on $X_1$ we get that $\mu_y^{12} = \mu_{y'}^{12}$. However, this contradicts Lemma 4.4 unless $\mu$ is invariant under $U_{12}$.

The main difficulty consists in ensuring that $x(r), x'(r)$ belong to the compact set $X_1$ and satisfy (4.7) and (4.8). For this we will need several other compact sets with large measure and various properties.

Our proof follows closely the methods of [20, §8]. The arguments can be simplified if one assumes additional regularity for the conditional measures $\mu_z^{12}$ — see [20, §8.1] for more details.

### 4.3. The construction of a nullset and three compact sets.

As mentioned before we will work with two main assumptions: that $\mu$ satisfies the assump-
tions of the low entropy case and that the equivalent conditions in Proposition 4.3 fail. By the former there exists a nullset $N$ so that all statements of Lemma 4.4 are satisfied for $x \in X \setminus N$. By the latter we can assume that for small enough $\varepsilon$ and for any compact set with $\mu(K) > 1 - \varepsilon$ the $A'$-returns to $K$ are not strong exceptional.

We enlarge $N$ so that $X \setminus N \subset X'$ where $X'$ is as in Section 2.1. Furthermore, we can assume that $N$ also satisfies Lemma 4.2. This shows that for every compact set $K \subset X \setminus N$ with $\mu(K) > 1 - \varepsilon$ the $A'$-returns (which exist due to Poincaré recurrence) are not exceptional, i.e. for every $\delta > 0$ there exists $z \in K$ and $s \in A'$ with $z' = \alpha^s z \in B_\delta(z) \setminus B_\delta^1(z)$.

**Construction of $X_1$.** The map $x \mapsto \mu_x^{12}$ is a measurable map from $X$ to a separable metric space. By Luzin’s theorem [5, p. 76] there exists a compact $X_1 \subset X \setminus N$ of measure $\mu(X_1) > 1 - \varepsilon^4$, and the property that $\mu_x^{12}$ depends continuously on $x \in X_1$.

**Construction of $X_2$.** To construct this set, we use the maximal inequality (10) in Section 2.1 from [20, App. A]. Therefore, there exists a set $X_2 \subset X \setminus N$ of measure $\mu(X_2) > 1 - C_1 \varepsilon^2$ (with $C_1$ some absolute constant) so that for any $R > 0$ and $x \in X_2$

$$\int_{[-R,R]} 1_{X_1}(u(r)x) \, d\mu_x^{12}(r) \geq (1 - \varepsilon^2)\mu_x^{12}([-R, R]).$$

**Construction of $K = X_3$.** Since $\mu_x^{12}$ is assumed to be nontrivial a.e., we have $\mu_x^{12}(\{0\}) = 0$ and $\mu_x^{12}([-1, 1]) = 1$. Therefore, we can find $\rho \in (0, 1/2)$ so that

$$\mathcal{X}(\rho) = \{ x \in X \setminus N : \mu_x^{12}([-\rho, \rho]) < 1/2 \}$$

has measure $\mu(\mathcal{X}(\rho)) > 1 - \varepsilon^2$. Let $t = (1, -1, 0, \ldots, 0) \in \Sigma$ be fixed for the following. By the (standard) maximal inequality we have that there exists a compact set $X_3 \subset X \setminus N$ of measure $\mu(X_3) > 1 - C_2 \varepsilon$ so that for every $x \in X_3$ and $T > 0$ we have

$$\frac{1}{T} \int_0^T 1_{X_3}(\alpha^{-\tau t} x) \, d\tau \geq (1 - \varepsilon),$$

$$\frac{1}{T} \int_0^T 1_{\mathcal{X}(\rho)}(\alpha^{-\tau t} x) \, d\tau \geq (1 - \varepsilon).$$

4.4. **The construction of $z, z' \in X_3, x, x' \in X_2$.** Let $\delta > 0$ be very small (later $\delta$ will approach zero). In particular, the matrix $g \in B_\delta^G$ (with entries as in (4.1)) is uniquely defined by $z' = gz$ whenever $z, z' \in X_3$ and $d(z, z') < \delta$. Since the $A'$-returns to $X_3$ are not exceptional, we can find $z \in X_3$ and $\alpha^s \in A'$
with \( z' = \alpha^s z \in B_3(z) \cap X_3 \) so that
\[
(4.12) \quad \kappa(z, z') = \max\left( |a_2 - a_1|, |g_{21}|^{1/2}, \|g_{12}\|, \|g_{22}\| \right) \in (0, c\delta^{3/2}),
\]
where \( c \) is an absolute constant allowing us to change from the metric \( d(\cdot, \cdot) \) to the norms we used above.

For the moment let \( x = z, \ x' = z' \), and \( r = \kappa(z, z')^{-1} \). Obviously
\[
\max\left( |(a_2 - a_1)r|, |g_{21}|^{1/2}r, \|g_{22}r\|, \|g_{12}r\| \right) = 1.
\]
If the maximum is achieved in one of the last two expressions, then (4.7) and (4.8) are immediate with \( C = 1 \). However, if the maximum is achieved in either of the first two expressions, it is possible that \( (a_2 - a_1)r - g_{21}r^2 \) is very small. In this case we could set \( r = 2\kappa^{-1}(z, z') \), then \( (a_2 - a_1)r \) is about 2 and \( g_{21}r^2 \) is about 4. Now (4.7)–(4.8) hold with \( C = 10 \). The problem with this naive approach is that we do not have any control on the position of \( x(r), x'(r) \). For all we know these points could belong to the null set \( N \) constructed in the last section.

To overcome this problem we want to use the conditional measure \( \mu_2^{12} \) to find a working choice of \( r \) in some interval \( I \) containing \( \kappa(z, z')^{-1} \). Again, this is not immediately possible since \textit{a priori} this interval could have very small \( \mu_2^{12} \)-measure, or even be a nullset. To fix this, we use \( \mathbf{t} = (1, -1, 0, \ldots, 0) \) and the flow along the \( \alpha^{t} \)-direction in Lemma 4.6. However, note that \( x = \alpha^{-\tau t} z \) and \( x' = \alpha^{-\tau t} z' \) differ by \( \alpha^{-\tau t} g\alpha^{rt} \). This results possibly in a difference of \( \kappa(x, x') \) and \( \kappa(z, z') \) as in Figure 3, and so we might have to adjust our interval along the way. The way \( \kappa(x, x') \) changes for various values of \( \tau \) depends on which terms give the maximum.

![Figure 3: The distance function \( \kappa(x, x') \) might be constant for small \( \tau \) and increase exponentially later.](image)

**Lemma 4.5.** For \( z, z' \in X_3 \) as above let \( T = \frac{1}{4} \left| \ln \kappa(z, z') \right|, \eta \in \{0, 1\}, \) and \( \theta \in [4T, 6T] \). There exist subsets \( P, P' \subset [0, T] \) of density at least \( 1 - 9\varepsilon \) such that for any \( \tau \in P, \tau' \in P' \),

(1) \( x = \alpha^{-\tau t} z \in X_2 \ (x' = \alpha^{-\tau t} z' \in X_2) \) and

(2) the conditional measure \( \mu_2^{12} \) satisfies the estimate
\[
(4.13) \quad \mu_2^{12}([-\rho S, \rho S]) < \frac{1}{2} \mu_2^{12}([-S, S])
\]
where \( S = S(\tau) = e^{\theta - \eta t} \) (and similarly for \( \mu_2^{12} \)).
Figure 4: From the way the leaf $U_{12}x$ is contracted along $\alpha^{-t}$ we can ensure (4.13) if $y = \alpha^{-wt}x \in X_\rho$.

Proof. By the first line in (4.11) there exists a set $Q_1 \subset [0, T]$ of density at least $1 - \varepsilon$ (with respect to the Lebesgue measure) such that $x = \alpha^{-rt}z$ belongs to $X_2$ for every $\tau \in Q_1$.

By the second line in (4.11) there exists a set $Q_2 \subset [0, 4T]$ of density at least $1 - \varepsilon$ such that $\alpha^{-vt}z \in X(\rho)$ for $v \in Q_2$. Let

$$Q_3 = \{ \tau \in [0, T] : \frac{1}{2}(\theta + (2 - \eta)\tau) \in Q_2 \}.$$ 

A direct calculations shows that $Q_3$ has density at least $1 - 8\varepsilon$ in $[0, T]$, and for $\tau \in Q_3$ and $v = \frac{1}{2}(\theta + (2 - \eta)\tau)$ we have $y = \alpha^{-vt}z \in X_\rho$.

We claim the set $P = Q_1 \cap Q_3 \subset [0, T]$ satisfies all assertions of the lemma; see Figure 4. First $P$ has at least density $1 - 9\varepsilon$. Now suppose $\tau \in P$; then $x = \alpha^{-rt}z \in X_2$ by definition of $Q_1$. Let $w = \frac{1}{2}(\theta - \eta\tau)$; then

$$y = \alpha^{-wt}x = \alpha^{-vt}z \in X_\rho$$

by the last paragraph. By (4.10)

$$\mu_{y}^{12}([-\rho, \rho]) < \frac{1}{2}\mu_{y}^{12}([-1, 1]) = \frac{1}{2}.$$ 

By property (4) in Section 2 of the conditional measures we get that

$$\frac{\mu_{y}^{12}([-\rho, \rho])}{\mu_{y}^{12}([-1, 1])} = \frac{(\alpha^{-wt}\mu_{x}^{12}\alpha^{wt})([-\rho, \rho])}{(\alpha^{-wt}\mu_{x}^{12}\alpha^{wt})([-1, 1])} = \frac{\mu_{x}^{12}([-\rho e^{2w}, \rho e^{2w}])}{\mu_{x}^{12}([-e^{2w}, e^{2w}])}.$$ 

This implies (4.13) for $S = e^{2w} = e^{\theta - \eta \tau}$. The construction of $P'$ for $z'$ is similar. \hfill \Box

The next lemma uses Lemma 4.5 to construct $x$ and $x'$ with the property that certain intervals containing $\kappa(x, x')^{-1}$ have $\mu_{x}^{12}$-measure which is not too small. This will allow us in Section 4.5 to find $r$ so that both $x(r)$ and $x'(r)$ have all the desired properties.
Lemma 4.6. Let \( z, z' \in X_3 \) and \( T = \frac{1}{4} \| \ln \kappa(z, z') \| \) be as above. If \( \varepsilon < \frac{1}{100} \), then there exists \( \tau \in [0, T] \) such that

1. both \( x = \alpha^{-\tau} z \) and \( x' = \alpha^{-\tau} z' \) are in \( X_2 \),
2. \( \kappa(x, x') < c\delta^{3/8} \), and
3. for \( R = \kappa(x, x')^{-1} \) (as well as \( R' = \rho^{-5} R \))

\[
\mu_x^\alpha([\rho R, \rho R]) < \frac{1}{2} \mu_x^\alpha([-R, R]) \quad \text{and} \quad \mu_x^\alpha([-\rho R, \rho R]) < \frac{1}{2} \mu_x^\alpha([-R, R]).
\]

(4.14)

Proof. Let
\[
\kappa_a(z, z') = |a_2 - a_1|,
\]
\[
\kappa_u(z, z') = \max(\|g_{21}\|^{1/2}, \|g_{1+}\|, \|g_{2+}\|) \in (0, c\delta^{1/2}).
\]

The corresponding quantities for \( x, x' \) are defined similarly. The number \( T \) is chosen so that the two points \( x = \alpha^{-\tau} z \) and \( x' = \alpha^{-\tau} z' \) are still close together for \( \tau \in [0, T] \). In fact,

\[
(4.15) \quad \tilde{g} = \alpha^{-\tau t} \left( \begin{array}{ccc} a_1 & g_{12} & g_{1+} \\ g_{21} & a_2 & g_{2+} \\ g_{1+} & g_{2+} & a_+ \end{array} \right) \alpha^\tau = \left( \begin{array}{ccc} a_1 & e^{-2\tau} g_{12} & e^{-\tau} g_{1+} \\ e^{2\tau} g_{21} & a_2 & e^{\tau} g_{2+} \\ e^{\tau} g_{1+} & e^{-\tau} g_{2+} & a_+ \end{array} \right)
\]

satisfies \( x' = \tilde{g} x \), and so

\[
(4.16) \quad \kappa_a(x, x') = \kappa_a(z, z'), \quad \kappa_u(x, x') = e^\tau \kappa_u(z, z') \leq \kappa(z, z')^{\frac{3}{4}} < c\delta^{\frac{3}{8}} \quad \text{and} \quad \kappa(x, x') = \max(\kappa_a(x, x'), \kappa_u(x, x')) < c\delta^{\frac{3}{8}};
\]

see also Figure 3. Hence the second statement of the lemma holds.

For the other two statements of the lemma we will use Lemma 4.5 to define four subsets \( P_a, P_u, P'_a, P'_u \subset [0, T] \), each of density at least \( 1 - 9\varepsilon \), so that for every \( \tau \) in the intersection of these four sets both (1) and (3) hold.

Definition of \( P_a \). If \( \kappa(x, x') > \kappa_a(x, x') \) for all \( \tau \in [0, T] \) (recall that \( x, x' \) depend implicitly on \( \tau \)) we set \( P_a = [0, T] \).

Otherwise, it follows from (4.16) that \( \kappa(z, z') = \kappa_a(z, z') \). We apply Lemma 4.5 for \( \eta = 0 \) and \( \theta = -\log \kappa_a(z, z') = 4T \), and see that (4.13) holds for \( \tau \in P_a \), where \( P_a \subset [0, T] \) has density at least \( 1 - 9\varepsilon \), and \( S_a = \kappa_a(z, z')^{-1} = \kappa_a(x, x')^{-1} \).

Definition of \( P_u \). If \( \kappa(x, x') > \kappa_u(x, x') \) for all \( \tau \in [0, T] \) we set \( P_u = [0, T] \).

Otherwise, it follows from (4.16) that \( \kappa(x, x') = e^\tau \kappa_u(z, z') \geq \kappa(z, z') \) for some \( \tau \in [0, T] \); hence \( \kappa_u(z, z') \in [\kappa(z, z')^{\frac{3}{4}} / \kappa(z, z')] \). This time, we apply Lemma 4.5 with \( \eta = 1 \) and \( \theta = -\log \kappa_u(z, z') \in [4T, 5T] \). We conclude that in
then taking the intersection of particular if $\epsilon$ difference being the use of a slightly different value for $R$ and (4.14) holds for $S = R = \kappa(x, x')^{-1} = \min(\kappa_\alpha(x, x'), \kappa_\nu(x, x'))$.

The sets $P'_a$ and $P'_u$ are defined similarly using $z'$.

The set $P_a \cap P'_a \cap P_u \cap P'_u \subset [0, T]$ has density at least $(1 - 36\epsilon)$, so that in particular if $\epsilon$ is small it is nonempty. For any $\tau$ in this intersection, $x, x' \in X_2$ and (4.14) holds for $R = \kappa(x, x')^{-1}$.

The additional statement in the parenthesis follows similarly, the only difference being the use of a slightly different value for $\theta$ in both cases, and then taking the intersection of $P_a \cap P'_a \cap P_u \cap P'_u$ with four more subsets of $[0, T]$ with similar estimates on their densities.

\[\square\]

### 4.5. Construction of $x(r), x'(r)$ and the conclusion of the proof

Recall that we found $z, z' \in X_3$ using Poincaré recurrence and the assumption that the $A'$-returns to $X_3$ are not exceptional. In the last section we constructed $x = \alpha^{-\tau}z, x' = \alpha^{-\tau}z' = \alpha^\delta \in X_2$ using the properties of $X_3$ to ensure (4.14).

Since $\alpha^\delta$ acts isometrically on the $U_{12}$-leaves, it follows from property (4) of the conditional measures in Section 2 that $\mu_z^{12} = \mu_{z'}^{12}$ and $\mu_x^{12} = \mu_{x'}^{12}$.

Let

\[
P = \{r \in [-R, R] : u(r)x \in X_1\}
\]

\[
P' = \{r \in [-R, R] : u(r)x' \in X_1\}.
\]

By (4.9) we know that $P$ and $P'$ both have density at least $(1 - \epsilon^2)$ with respect to the measure $\mu_z^{12} = \mu_{z'}^{12}$. By (4.14) we know that $[-\rho R, \rho R]$ contains less than one half of the $\mu_z^{12}$-mass of $[-R, R]$. Therefore, if $\epsilon$ is small enough there exists $r \in P \cap P' \setminus [-\rho R, \rho R]$. We define $x(r) = u(r)x$ and $x'(r) = u(r)x'$, and conclude that $x(r), x'(r) \in X_1$ satisfy $\mu_{x(r)}^{12} = \mu_{x'(r)}^{12}$ by property (3) in Section 2.

Let $\tilde{g}$ be defined as in (4.15) and write $\tilde{g}_{12}, \ldots$ for the matrix entries. With $\tilde{g}(r) = u(r)\tilde{g}u(-r)$ we have $x'(r) = \tilde{g}(r)x(r)$ and

\[
\tilde{g}(r) = \begin{pmatrix} a_1 + \tilde{g}_{21r} & \tilde{g}_{12} + (a_2 - a_1)r - \tilde{g}_{21r}^2 & \tilde{g}_{1s} + \tilde{g}_{2s}r \\ \tilde{g}_{21} & a_2 - \tilde{g}_{21r} & \tilde{g}_{2s} \\ \tilde{g}_{s1} & \tilde{g}_{s2} - \tilde{g}_{s1r} & a_s \end{pmatrix}.
\]

We claim it is possible to achieve

\[
C^{-1} \leq \max(|(a_2 - a_1)r - \tilde{g}_{21r}^2|, \|\tilde{g}_{2s}r\|, \|\tilde{g}_{s1r}\|) \leq C,
\]

\[
|\tilde{g}_{21r}| \leq C\delta^{3/8}.
\]

for some constant $C$; see Figure 2.
nonexceptional returns to \(X\) two points \(x_i\), i.e. the problem of almost cancellation cannot happen again. We now turn to prove (4.17). It is immediate from the definition of \(R\) that

\[
(4.19) \quad \rho \leq \max(||(a_2 - a_1)r||, |\bar{g}_{21}|^{1/2}|r||, |\bar{g}_{21}r||, ||\bar{g}_{21}r||) \leq 1.
\]

There are two differences of this estimate to the one in (4.17); first we need to take the square of the second term – this replaces the lower bound \(\rho\) by its square, secondly we looked above at \((a_2 - a_1)r\) and \(\bar{g}_{21}r^2\) separately – taking the difference as in (4.7) might produce a too small a number (almost cancellation). So (4.17) follows with \(C = \rho^{-2}\), unless

\[
(4.20) \quad \max(||(a_2 - a_1)r||, |\bar{g}_{21}|^{1/2}|r||) \geq \rho,
\]

\[
(4.21) \quad ||(a_2 - a_1)r - \bar{g}_{21}r^2|| < \rho^2 < \frac{\rho}{2}.
\]

This is a minor problem, and we can overcome it using the last statement in Lemma 4.6. Assume that for some \(r \in [-R, R] \setminus [-\rho R, \rho R]\) this problem occurs. We deduce a lower estimate on \(|\bar{g}_{21}|\). If the maximum in (4.20) is achieved at \(|\bar{g}_{21}|^{1/2}|r|| \geq \rho\), then \(\bar{g}_{21} \geq \rho^2 R^{-2}\). If the maximum is achieved at \(||(a_2 - a_1)r|| \geq \rho\), (4.21) shows that \(|\bar{g}_{21}r^2|| \geq \rho/2 \geq \rho^2\) (since \(\rho < 1/2\)) and so in both cases

\[
(4.22) \quad |\bar{g}_{21}| \geq \rho^2 R^{-2}.
\]

Now we go through the construction of \(r\) again, only this time using the last statement in Lemma 4.6, and find \(r \in [\rho^{-5} R, \rho^{-5} R] \setminus [\rho^{-4} R, \rho^{-4} R]\). The equivalent to (4.19) is now the estimate

\[
\rho^{-4} \leq \max(||(a_2 - a_1)r||, |\bar{g}_{21}|^{1/2}|r||, |\bar{g}_{21}r||, ||\bar{g}_{21}r||) \leq \rho^{-5}.
\]

This shows that \(||(a_2 - a_1)r|| \leq \rho^{-5}\), and (4.22) shows that

\[
|\bar{g}_{21}r^2|| \geq \rho^2 R^{-2}(\rho^{-4} R)^2 = \rho^{-6}.
\]

Together, we find a lower bound for

\[
||(a_2 - a_1)r - \bar{g}_{21}r^2|| \geq \rho^{-6} - \rho^{-5} > 0,
\]

i.e. the problem of almost cancellation cannot happen again.

Starting with the nonexceptional return \(z' \in X_3\) of \(z \in X_3\) we have found two points \(x_r, x'_r \in X_1\) which satisfy (4.17), (4.18). Since we assume to have nonexceptional returns to \(X_3\) for every \(\delta = \frac{1}{n} > 0\), we get two sequences \(y_n\) and \(y'_n\) of points in \(X_1\) with the same conditional measures

\[
\mu_{y_n}^{12} = \mu_{y'_n}^{12}.
\]
Compactness shows that we can find convergent subsequences with limits \( y, y' \in X_1 \). It follows from (4.17), (4.18) that \( y' \in U_y \), and from (4.17) that \( y' \neq y \). By continuity of \( \mu^y_{12} \) for \( x \in X_1 \) the conditional measures \( \mu^y_{12} = \mu'^y_{12} \) agree. However, this contradicts Lemma 4.4, unless \( \mu \) is invariant under \( U_{12} \).

5. Proof that exceptional returns are not possible for \( \Gamma = \text{SL}(k, \mathbb{Z}) \)

If case (3) in Theorem 2.1 holds, then this gives some restriction on \( \Gamma \). In other words, for some lattices in \( \text{SL}(k, \mathbb{R}) \), exceptional returns cannot occur. As will be shown below, such is the case for \( \Gamma = \text{SL}(k, \mathbb{Z}) \).

We recall that \( H_{ab} \subset \text{SL}(k, \mathbb{R}) \) is an \( A \)-normalized subgroup isomorphic to \( \text{SL}(2, \mathbb{R}) \), and \( A' = A \cap C(H_{ab}) \). If case (3) of Theorem 2.1 holds then any \( A' \)-ergodic component of \( \mu \) is supported on a single \( C(H_{ab}) \)-orbit. In particular, we have an abundance of \( A' \)-invariant probability measures supported on single \( C(H_{ab}) \) orbits. The mere existence of such measures is a restriction on \( \Gamma \).

**Theorem 5.1.** Suppose that \( \nu \) is an \( A' \) invariant probability measure on \( X = \text{SL}(k, \mathbb{R})/\Gamma \), and that \( \text{supp} \nu \subset C(H_{ab})x \) for some \( x \in X \). Then there is a \( \gamma \in \Gamma \) which is

1. diagonalizable over \( \mathbb{R} \),
2. \( \pm 1 \) is not an eigenvalue of \( \gamma \).
3. All eigenvalues of \( \gamma \) are simple except precisely one which has multiplicity two.

Before we prove this theorem, we note the following:

**Proposition 5.2.** There is no \( \gamma \in \text{SL}(k, \mathbb{Z}) \) satisfying the three conditions of Theorem 5.1.

In particular, case (3) of Theorem 2.1 cannot occur for \( \text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z}) \).

**Proof of Proposition 5.2.** Suppose \( \gamma \in \text{SL}(k, \mathbb{Z}) \) is diagonalizable over \( \mathbb{R} \). Then its eigenvalues (with the correct multiplicities) are roots of the characteristic polynomial of \( \gamma \), a polynomial with integer coefficients and both leading term and constant term equal to one. If there is some eigenvalue which is not equal to \( \pm 1 \) and which occurs with multiplicity greater than one then necessarily this eigenvalue is not rational, and its Galois conjugates would also have multiplicity greater than one, contradicting (3).

To prove Theorem 5.1, we need the following standard estimate:

**Lemma 5.3.** There is a neighborhood \( U_0 \) of the identity in \( \text{SL}(m, \mathbb{R}) \) so that for any \( \lambda_1, \lambda_2, \ldots, \lambda_m \) with \( |\lambda_i - \lambda_j| > 1 \) and \( h \in U_0 \) one has that
and more generally $h \operatorname{diag}(e^{\lambda_1}, e^{\lambda_2}, \ldots, e^{\lambda_m})$ is diagonalizable over $\mathbb{R}$ with positive eigenvalues and the eigenvalues $e^{\lambda_1}, e^{\lambda_2}, \ldots, e^{\lambda_m}$ satisfy $|\lambda'_i - \lambda_i| < \frac{1}{2}$.

Proof. Without loss of generality, suppose $\lambda_1 > \lambda_2 > \cdots > \lambda_m$. Let $f = \operatorname{diag}(e_\lambda^1, e_\lambda^2, \ldots, e_\lambda^m)$ and $\eta_1, \eta_2, \ldots$ be the eigenvalues of $f' = hf$ ordered according to descending absolute value. Set for $1 \leq i \leq m$, $\lambda'_i = \log |\eta_i|$.

Clearly, $\lambda'_i = \lim_{n \to \infty} \frac{\log \|f'^n\|}{n}$. Since $f$ is self adjoint, $\|f\| = e^{\lambda_1}$ so

$$\lambda'_i \leq \log \|h\| + \lambda_1.$$ 

Let $\delta > 0$ be small (it will be chosen later and will be independent of $f$).

Consider the cones

$$K = \{(x_1, \ldots, x_m) : |x_l| \leq \delta |x_1| \text{ for every } l \neq 1 \}$$

$$K' = \{(x_1, \ldots, x_m) : |x_l| \leq \delta e^{-1} |x_1| \text{ for every } l \neq 1 \}.$$ 

Then $fK \subset K'$, and for every $x \in K$ 

$$\|fx\| \geq (1 - c\delta) e^{\lambda_1} \|x\|,$$

for some $c$ depending only on $m$. Suppose now that $h$ is close enough to the identity so that $hK' \subset K$. Then $f'K = hfK \subset hK' \subset K$. Again assuming that $h$ is in some fixed neighborhood of the identity, for any $x \in K$, we have

$$\|f'x\| \geq (1 - 2c\delta) e^{\lambda_1} \|x\|$$

so that

$$\|f'^n\| \geq \|f'^n e_1\| \geq \left((1 - 2c\delta) e^{\lambda_1}\right)^n.$$ 

In other words, if $h$ is in some fixed neighborhood of the identity (independently of $f$) then

$$|\lambda'_1 - \lambda_1| < C_1 \delta.$$ 

Similarly, $e^{\lambda_1+\lambda_2}$ is the dominating eigenvalue of $f \wedge f$, i.e. the natural action of $f$ on the space $\mathbb{R}^n \wedge \mathbb{R}^n$. Applying the same logic as before, $\lambda'_1 + \lambda'_2 = \lim_{n \to \infty} \frac{\log \|f^{\wedge n}\|}{n}$, and as long as $h$ is in some fixed neighborhood of the identity, independently of $f$ 

$$|\lambda'_1 + \lambda'_2 - \lambda_1 - \lambda_2| < C_2 \delta$$

and more generally

$$(5.1) \quad \left| \sum_{i=1}^k (\lambda'_i - \lambda_i) \right| < C_3 \delta.$$ 

Clearly, (5.1) implies that there is some $C$ depending only on $m$, and a neighborhood of the identity in $\text{SL}(m, \mathbb{R})$ depending only on $\delta$ so that if $h$ is in that neighborhood 

$$|\lambda'_i - \lambda_i| < C \delta \quad \text{for all } 1 \leq i \leq m.$$
In particular, if $\lambda_i > \lambda_{i-1} + 1$ for every $i$ then if $C\delta < \frac{1}{2}$, all $\lambda'_i$ are distinct. Since this holds for all $h$ in a connected neighborhood of the identity, all the eigenvalues of $f'$ are real and also positive, so that $\eta_i = e^{\lambda'_i}$.

**Proof of Theorem 5.1.** Without loss of generality, we may take $a = 1$, $b = 2$. Let $a(t) = \text{diag}(e^{\lambda_1 t}, \ldots, e^{\lambda_k t}) \in A'$ with $\lambda_1 = \lambda_2$ and for every other pair $i, j$ we have $\lambda_i \neq \lambda_j$.

Take $U_0$ to be a symmetric neighborhood of the identity in $\text{SL}(k-2, \mathbb{R})$ as in Lemma 5.3, and

$$U_1 = \left\{ \begin{pmatrix} e^r & 0 & 0 \\ 0 & e^s & 0 \\ 0 & 0 & e^{-\frac{2r}{k-2}}h \end{pmatrix} : r \in (-1/8, 1/8) \text{ and } h \in U_0 \right\},$$

and let $t_0 = 2 \max_{\lambda_i \neq \lambda_j} |\lambda_i - \lambda_j|^{-1}$. Note that $U_1$ is also symmetric, i.e. $U_1^{-1} = U_1$.

By Poincaré recurrence, for $\nu$-almost every $x = g\Gamma \in \text{SL}(k, \mathbb{R})/\Gamma$ there is a $t > t_0$ so that $a(t)x \in U_1 x$; so in particular $U_1 a(t) \cap g\Gamma g^{-1} \neq \emptyset$. Let

$$\tilde{\gamma} = \begin{pmatrix} e^s & 0 & 0 \\ 0 & e^s & 0 \\ 0 & 0 & f' \end{pmatrix} \in U_1 a(t) \cap g\Gamma g^{-1}$$

be any element from this intersection. By assumption, for every pair $i, j$ except $1, 2$ we have

$$|\lambda_i - \lambda_j| t > 2,$$

and we can apply Lemma 5.3 to deduce that the eigenvalues of

$$f' = he^{-\frac{2r}{k-2}} \text{diag}(e^{\lambda_3 t}, e^{\lambda_4 t}, \ldots, e^{\lambda_k t})$$

for some $h \in U_0$ are of the form $e^{\lambda'_i}, \ldots, e^{\lambda'_k}$ with $|\lambda'_i - \lambda_i t| < 3/4$ for $i = 3, \ldots, k$. Finally $|s - \lambda_1 t| = |s - \lambda_2 t| \leq 1/8$.

In view of (5.2) it is clear that $\tilde{\gamma}$ and hence $g^{-1}\tilde{\gamma}g \in \Gamma$ satisfy all the conditions of Theorem 5.1. \qed

**6. Conclusion of the proof of Theorem 1.3**

In this section, we conclude the derivation of Theorem 1.3, and its corollary, Corollary 1.4, from Theorem 2.1. Throughout this section, $X$ will denote the quotient space $\text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z})$, and $\mu$ be an $A$-ergodic and invariant probability measure on $X$. For every pair $a, b$ of distinct indices in $\{1, \ldots, k\}$, one of the three possibilities of Theorem 2.1 holds. However, in view of the results of the previous section, in particular Theorem 5.1 and Proposition 5.2, Theorem 2.1.(3), i.e. the case of exceptional returns, cannot occur for the lattice.
SL(k, ℤ). Therefore, for every pair \(a, b\) of distinct indices one of the following two mutually exclusive possibilities holds:

1. The conditional measures \(\mu_{x}^{ab}\) and \(\mu_{x}^{ba}\) are trivial a.e.

2. The conditional measures \(\mu_{x}^{ab}\) and \(\mu_{x}^{ba}\) are Haar and \(\mu\) is invariant under left multiplication with elements of \(H_{ab} = \langle U_{ab}, U_{ba} \rangle\).

Define a relation \(a \sim b\) if \(\mu\) is \(U_{ab}\)-invariant. By (2) above it follows that \(a \sim b\) if and only if \(b \sim a\). Furthermore, since the group generated by \(U_{ab}\) and \(U_{bc}\) contains \(U_{ac}\), it is clear that \(\sim\) is in fact an equivalence relation on \(\{1, \ldots, k\}\). Let \(H\) be the group generated by all \(U_{ab}\) with \(a \sim b\). Let \(r\) denote the number of equivalence classes for \(\sim\) which contain more than one element, and \(k_1, k_2, \ldots, k_r\) be their sizes; so in particular \(\sum_{i=1}^{r} k_i \leq k\). By permuting the indices if necessary we can assume these equivalence classes are consecutive indices and \(H = \prod_{i=1}^{r} \text{SL}(k_i, ℤ)\). By definition, \(H\) leaves the measure \(\mu\) invariant, normalized by \(A\), and is generated by unipotent one-parameter subgroups of \(\text{SL}(n, ℤ)\) — indeed, \(H\) is precisely the maximal subgroup of \(\text{SL}(n, ℤ)\) satisfying these three conditions.

Measures invariant under groups generated by unipotent one-parameter groups are well understood. In particular, in a seminal series of papers culminating in [41], M. Ratner showed that if \(H\) is such a group the only \(H\)-ergodic and invariant probability measures are the algebraic measures: \(L\)-invariant measures supported on a closed \(L\)-orbit for some \(L > H\) (here and throughout, we use the notation \(L > H\) to denote that \(H\) is a subgroup and \(L \supset H\); specifically \(H\) may be equal to \(L\)) . For the \(A\)-invariant measure \(\mu\) and \(H\) as above we only know that \(\mu\) is \(A\)-ergodic and \(H\)-invariant, but not necessarily \(H\)-ergodic; we shall use the following version of these measure-rigidity results by Margulis and Tomanov [27][3] (similar techniques were used also in [30, proof of Thm. 1]; see also [18, §4.4] and [47]). For any connected real Lie group \(G\), we shall say that \(g \in G\) in an element of class \(A\) if \(\text{Ad} g\) is semisimple, with all eigenvalues integer powers of some \(\lambda \in ℚ \setminus \{±1\}\), and \(g\) is contained in a maximal reductive subgroup of \(G\).

**Theorem 6.1 ([27, Thms. (a) and (b)])**. Let \(G\) be a connected real Lie group, \(\Gamma < G\) a discrete subgroup, and \(\tilde{H}\) generated by unipotent one-parameter groups and elements of class \(A\), with \(H < \tilde{H}\) the subgroup generated by unipotent one-parameter groups. Let \(\mu\) be an \(\tilde{H}\)-invariant and ergodic probability measure on \(G/\Gamma\). Then there is an \(L \geq H\) so that almost every \(H\)-ergodic component of \(\mu\) is the \(L\)-invariant probability measure on a closed \(L\)-orbit.

---

3The main theorem of [27] was substantially more general than what we quote here. In particular, in their theorem \(\Gamma\) can be any closed subgroup of \(G\), and the group \(G\) can be a product of real and \(p\)-adic Lie groups (satisfying some mild additional conditions).
Furthermore, if
\[ SN_G(L) = \{g \in N_G(L) : \text{conjunction by } g \text{ preserves Haar measure on } L \} \]
then \( \hat{H} < SN_G(L) \) and \( \mu \) is supported on a single \( SN_G(L) \)-orbit. In particular, \( L \) is normalized by \( \hat{H} \).

**Lemma 6.2.** Let \( H = \prod_{i=1}^{r} \text{SL}(k_i, \mathbb{R}) \) with \( \sum_{i=1}^{r} k_i < k \). Let
\[ X_H = \{ x \in X : \text{if } Hx \text{ is closed and of finite volume} \} . \]
Then there is a one-parameter subgroup \( a(t) \) of \( A \) so that for every \( x \in X_H \) its trajectory \( a(t)x \to \infty \) as \( t \to \infty \).

**Proof.** Suppose \( Hx \) is closed and of finite volume, with \( x = g \text{SL}(k, \mathbb{Z}) \) and \( g = (g_{ij}) \). Let \( k' = \sum_i k_i < k \). Let \( \Lambda = g^{-1}Hg \cap \Gamma \). Since \( \Lambda \) is Zariski dense in \( g^{-1}Hg \) there is a \( \gamma = g^{-1}h_0g \in \Lambda \) with \( h_0 = \begin{pmatrix} h_0' & 0 \\ 0 & I_{k-k'} \end{pmatrix} \) so that
\[ (6.1) \quad V_g := \left\{ y \in \mathbb{R}^k : y^Tg^{-1}h_0g = y^T \text{ for all } h \in H \right\} = \left\{ y \in \mathbb{R}^k : y^T \gamma = y^T \right\} . \]
Notice that since \( g(g^{-1}h_0g) = h_0g \) (the transpose of) the last \( k - k' \) rows of \( g \) are in \( V_g \).

Clearly \( \dim V_g = k - k' \), and using the right hand side of (6.1) it is clear that \( V_g \) is a rational subspace of \( \mathbb{R}^n \) (i.e. has a basis consisting of rational vectors). Since \( V_g \) is rational, there is an integer vector \( m \in \mathbb{Z}^n \cap (V_g)_{\perp} \). In particular, the last \( k - k' \) entries in the vector \( gm \) (which is a vector in the lattice in \( \mathbb{R}^k \) corresponding to \( g \text{SL}(k, \mathbb{Z}) \)) are zero. For any \( t \in \mathbb{R} \) set \( t = (t_1, \ldots, t_k) \) with \( t_1 = \cdots = t_{k'} = k't \) and \( t_{k'+1} = \cdots = t_k = (k' - k)t \) and \( a(t) = a^t \). Then since the last \( k - k' \) entries in the vector \( gm \) are zero,
\[ a(t)(gm) \to 0 \quad \text{as } t \to \infty, \]
so that by Mahler’s criterion \( a(t)x = a(t)g \text{SL}(k, \mathbb{Z}) \to \infty \).

We are finally in a position to finish the proof of Theorem 1.3:

**Proof of Theorem 1.3.** Let \( H = \prod_{i=1}^{r} \text{SL}(k_i, \mathbb{R}) \) be the maximal group fixing \( \mu \), generated by unipotent one-parameter subgroups, and normalized by \( A \) as above. By Theorem 6.1, applied to \( \hat{H} = AH \), we know that there is some \( L > H \) which is normalized by \( A \) so that almost every \( H \)-ergodic component of \( \mu \) is the \( L \)-invariant measure on a closed \( L \) orbit. In particular \( \mu \) is \( L \)-invariant, which unless \( L < AH \) contradicts the definition of \( H \) as the maximal group with the above properties. Let now \( x = g \text{SL}(k, \mathbb{Z}) \) have a closed \( L \)-orbit \( Lx \) of finite volume. Then \( \Lambda_L = g^{-1}Lg \cap \text{SL}(k, \mathbb{Z}) \) is a lattice in \( g^{-1}Lg \), and so the latter is defined over \( \mathbb{Q} \). Therefore, the same is true for the semi-simple \( g^{-1}Hg = [g^{-1}Lg, g^{-1}Lg] \), \( \Lambda_H = g^{-1}Hg \cap \text{SL}(k, \mathbb{Z}) \) is a lattice in \( g^{-1}Hg \), and \( Hx \) is closed with finite volume. However, this implies \( H = L \).
Thus we conclude that almost every $H$-ergodic component of $\mu$ is supported on a single $H$-orbit; in other words, in the notations of Lemma 6.2, the support of $\mu$ is contained in $X_H$.

By Lemma 6.2, this implies that the sum $\sum_i k_i = k$ since otherwise there is a one-parameter subgroup $a(t)$ of $A$ so that for every $x \in X_H$ its trajectory $a(t)x \to \infty$ as $t \to \infty$, in contradiction to Poincaré recurrence.

But if $\sum_i k_i = k$, the set $SN_G(H)$ of Theorem 6.1 satisfies

$$SN_G(H) = N_G(H) = AH$$

and so by this theorem $\mu$ is supported on a single $AH$-orbit. But $\mu$ is also $AH$-invariant. This show that $\mu$ is algebraic: an $AH$-invariant probability measure on a single $AH$-orbit. Note that this $AH$-orbit has finite volume, hence is closed in $X$.

Proof of Corollary 1.4. Let $\mu$ be an $A$-ergodic probability measure on $X$ with positive entropy. By Theorem 1.3, $\mu$ is algebraic, i.e. there are a subgroup $A < L < G$ and a point $x = g \text{SL}(k, \mathbb{Z}) \in X$ so that $Lx$ is closed and $\mu$ is the $L$-invariant measure on $Lx$.

Since $\mu$ is a probability measure, this implies that $g^{-1}Lg \cap \text{SL}(k, \mathbb{Z})$ is a lattice in $g^{-1}Lg$, which, in turn, implies that $g^{-1}Lg$ is defined over $\mathbb{Q}$. Moreover, the fact that $L$ has any lattice implies it is unimodular, which in view of $A < L < G$ (and since $A$ is the maximal torus in $G$) implies $L$ is reductive (this can also be seen directly from the proof of Theorem 1.3).

We conclude that $g^{-1}Lg$ is a reductive group defined over $\mathbb{Q}$, and $g^{-1}Ag$ is a maximal torus in this group. By [36, Thm. 2.13], there is an $h \in L$ so that $g^{-1}h^{-1}Ahg$ is defined over $\mathbb{Q}$ and is $\mathbb{Q}$-anisotropic. This implies that $Ahx$ is closed and of finite volume (i.e., since $A \cong \mathbb{R}^{k-1}$, compact), so that $Lx$ contains a compact $A$ orbit.

By [21, Thm. 1.3], it follows that (possibly after conjugation by a permutation matrix), $L$ is the subgroup of $g = (g_{ij}) \in \text{SL}(k, \mathbb{R})$ with $g_{ij} = 0$ unless $i$ is congruent to $j$ mod $m$ for some $1 \neq m|k$ (by the Moore ergodicity theorem it is clear that $A$ acts ergodically on $Lx$, hence the condition in that theorem that $Lx$ contains a relatively dense $A$ orbit is satisfied), and that $Lx$ is not compact. Note that if $k$ is prime this implies that $L = \text{SL}(k, \mathbb{R})$.

Part 2. Positive entropy and the set of exceptions to Littlewood’s Conjecture

7. Definitions

We recall the definition of Hausdorff dimension, box dimension, topological and metric entropy. In the following let $Y$ be a metric space with metric $d_Y(\cdot, \cdot)$. 
7.1. Notions of dimension. For $D \geq 0$ the $D$-dimensional Hausdorff measure of a set $B \subset Y$ is defined by

$$\mathcal{H}^D(B) = \lim_{\varepsilon \to 0} \inf \sum_i (\text{diam } C_i)^D,$$

where $C_\varepsilon = \{C_1, C_2, \ldots \}$ is any countable cover of $B$ with sets $C_i$ of diameter $\text{diam}(C_i) < \varepsilon$. Clearly, for $D > m$ any set in the Euclidean space $\mathbb{R}^m$ has Hausdorff measure zero. The Hausdorff dimension $\text{dim}_H(B)$ is defined by

$$\text{dim}_H(B) = \inf \{D : \mathcal{H}^D(B) = 0\} = \sup \{D : \mathcal{H}^D(B) = \infty\}.$$

For every $\varepsilon > 0$ a set $F \subset B$ is $\varepsilon$-separated if $d_Y(x, y) \geq \varepsilon$ for every two different $x, y \in F$. Let $b_\varepsilon(B)$ be the cardinality of the biggest $\varepsilon$-separated subset of $B$; then the (upper) box dimension (upper Minkowski dimension) is defined by

$$(7.2) \quad \text{dim}_\text{box}(B) = \limsup_{\varepsilon \to 0} \frac{\log b_\varepsilon(B)}{\log \varepsilon}.$$

Note that $b_\varepsilon(B) \geq b_{\varepsilon_1}(B)$ if $\varepsilon_2 < \varepsilon_1$. Therefore, it is sufficient to consider a sequence $\varepsilon_n$ in (7.2) if $\log \varepsilon_{n+1}/\log \varepsilon_n \to 1$ for $n \to \infty$.

We recall some elementary properties. First, Hausdorff dimension and box dimension do not change when we use instead of the metric $d_Y(\cdot, \cdot)$ a different but Lipschitz equivalent metric $d'_Y(\cdot, \cdot)$. The Hausdorff dimension of a countable union is given by

$$(7.3) \quad \text{dim}_H(\bigcup_{i=1}^{\infty} B_i) = \sup \text{dim}_H(B_i).$$

(This follows easily from the fact that the measure $\mathcal{H}^D$ is subadditive.) For any $B$ we have

$$(7.4) \quad \text{dim}_H(B) \leq \text{dim}_{\text{box}}(B).$$

If $Y = Y_1 \times Y_2$ is nonempty and

$$d_Y((x_1, x_2), (y_1, y_2)) = \max(d_{Y_1}(x_1, x_2), d_{Y_2}(y_1, y_2)),$$

then $\text{dim}_{\text{box}} Y \leq \text{dim}_{\text{box}} Y_1 + \text{dim}_{\text{box}} Y_2$.

7.2. Entropy and the variational principle. Let $T$ be an endomorphism of a compact metric space $Y$. For $\varepsilon > 0$ and a positive integer $N$ we say that a set $E \subset Y$ is $(N, \varepsilon)$-separated (with respect to $T$) if for any two different $x, y \in E$ there exists an integer $0 \leq n < N$ with $d(T^n x, T^n y) \geq \varepsilon$. Let $s_{N, \varepsilon}(T)$ be the cardinality of the biggest $(N, \varepsilon)$-separated set; then the topological entropy of $T$ is defined by

$$(7.5) \quad h_{\text{top}}(T) = \lim_{\varepsilon \to 0} \limsup_{N \to \infty} \frac{1}{N} \log s_{N, \varepsilon}(T) = \sup_{\varepsilon > 0} \limsup_{N \to \infty} \frac{1}{N} \log s_{N, \varepsilon}(T).$$
Let $\mu$ be a $T$-invariant measure on $Y$, and let $\mathcal{P}$ be a finite partition of $Y$ into measurable sets. Then

$$H_\mu(\mathcal{P}) = -\sum_{P \in \mathcal{P}} \mu(P) \log \mu(P)$$

is the entropy of the finite partition $\mathcal{P}$. (Here $0 \log 0 = 0$.) For two such partitions $\mathcal{P}$ and $\mathcal{Q}$ let $\mathcal{P} \vee \mathcal{Q} = \{ P \cap Q : P \in \mathcal{P}, Q \in \mathcal{Q} \}$ be the common refinement. The metric entropy of $T$ with respect to $\mu$ and $\mathcal{P}$ is defined by

$$h_\mu(T, \mathcal{P}) = \lim_{N \to \infty} \frac{1}{N} H_\mu \left( \bigvee_{i=0}^{N-1} T^{-i} \mathcal{P} \right)$$

and the metric entropy of $T$ with respect to $\mu$ is

$$h_\mu(T) = \sup_{\mathcal{P}} h_\mu(T, \mathcal{P}), \quad (7.6)$$

where the supremum is taken over all finite partitions $\mathcal{P}$ of $Y$ into measurable sets.

Topological and metric entropy are linked: For a compact metric space $Y$, a continuous map $T : Y \to Y$, and a $T$-invariant measure $\mu$ on $Y$ the entropies satisfy

$$h_\mu(T) \leq h_{\text{top}}(T).$$

Furthermore, the variational principle [49, Thm. 8.6] states that

$$h_{\text{top}}(T) = \sup_\mu h_\mu(T), \quad (7.7)$$

where the supremum is taken over all $T$-invariant measures $\mu$ on $Y$.

8. Box dimension and topological entropy

We return to the study of the left action of the positive diagonal subgroup $A$ on $X = \text{SL}(n, \mathbb{R})/\text{SL}(n, \mathbb{Z})$. We fix an element $a \in A$ and study multiplication from the left by $a$ on $X$, in particular we are interested in the dynamical properties of the restriction $a|_K$ of this map to a compact subset $K \subset X$. This will lead to a close connection between topological entropy and box dimension in an unstable manifold.

The following easy lemma shows that the dimensions can be defined using the right invariant metric or a norm on $\text{Mat}(k, \mathbb{R})$.

**Lemma 8.1.** For every $r > 0$ there exists a constant $c_0 \geq 1$ such that

$$c_0^{-1} \|g - h\| \leq d(g, h) \leq c_0 \|g - h\| \quad \text{for all } g, h \in B_r^G,$$

where $\|A\| = \max_{i,j} |a_{ij}|$ for $A = (a_{ij}) \in \text{Mat}(k, \mathbb{R})$. 

$X$ is locally isomorphic to $\text{SL}(k, \mathbb{R})$; more specifically, for every $x \in X$ there exists some $r = r(x) > 0$ such that $B^c_0$ and $B_r(x)$ are isomorphic by sending $g$ to $gx$. For small enough $r$ this is an isometry. For a compact set $K \subset X$ we can choose $r = r(K) > 0$ uniformly with this property for all $x \in K$.

Let $x \in X$, $g \in B^c_0$, $y = gx$, $t \in \Sigma$, and $a = a^t$. Then $ay = (aga^{-1})ax$. In other words, when we use the local description of $X$ as above at $x$ and $ax$, left multiplication by $a$ acts in this local picture like conjugation by $a$ on $B^c_0$.

For this reason we define the subgroups

$$U = \{g \in \text{SL}(k, \mathbb{R}) : a^n ga^{-n} \to 0 \text{ for } n \to -\infty\},$$

$$V = \{g \in \text{SL}(k, \mathbb{R}) : a^n ga^{-n} \to 0 \text{ for } n \to \infty\},$$

$$C = \{g \in \text{SL}(k, \mathbb{R}) : aga^{-1} = g\},$$

which are the unstable, stable, and central subgroup (for conjugation with $a$).

Let $\delta_{ij} = 0$ for $i \neq j$ and $\delta_{ii} = 1$, so that $I_k = (\delta_{ij})$. It is easy to check that $g \in C$ if $g_{ij} = 0$ for all $i, j$ with $a_{ii} \neq a_{jj}$, $g \in U$ if $g_{ij} = \delta_{ij}$ for all $i, j$ with $a_{ii} \leq a_{jj}$, and similarly $g \in V$ if $g_{ij} = \delta_{ij}$ for all $i, j$ with $a_{ii} \geq a_{jj}$.

Furthermore, there exists a neighborhood $U_0 \subset \text{SL}(k, \mathbb{R})$ of the identity so that every $g \in U_0$ can be written uniquely as $g = g_C g_U g_V$ for some small $g_C \in C$, $g_U \in U$, and $g_V \in V$. If, similarly, $h = h_C h_U h_V$, then

$$c_1^{-1}d(g, h) \leq \max(d(g_C, h_C), d(g_U, h_U), d(g_V, h_V)) \leq c_1 d(g, h)$$

for some constant $c_1 \geq 1$.

Since $A$ is commutative, we have $A \subset C$. The map $T(x) = ax$ on $X$ is partially hyperbolic: $T$ is not hyperbolic (since the identity is not an isolated point of $C$), but part of the local description has hyperbolic structure as follows.

**Lemma 8.2.** Let $K \subset X$ be compact with $aK \subset K$ and let $r = r(K)$ be as above. There exists $\lambda > 1$ and $c_2 > 0$ so that for any small enough $\varepsilon > 0$, any $z \in K$ and $f \in B^c_0$, and any integer $N \geq 1$ with $d(fz, z) \geq \lambda^{-N} \varepsilon$, there exists a nonnegative integer $n < N$ with $d(a^n fz, a^n z) \geq c_2 \varepsilon$.

**Proof.** By continuity there exists $\varepsilon \in (0, r)$ such that $d(afa^{-1}, I_k) < r$ whenever $d(f, I_k) < \varepsilon$. This will be the only requirement on $\varepsilon$. On the other hand, since $U$ is expanded by conjugation with $a$, there exists some $\lambda > 1$ so that $\|afa^{-1} - I_k\| \geq \lambda \|f - I_k\|$ for all $f \in B^c_0$. By Lemma 8.1 $d(a^n fa^{-n}, I_k) \geq c_0^{-2} \lambda^n d(f, I_k)$ for all $f \in B^c_0$ and all $n$ for which

$$\max(d(f, I_k), \ldots, d(a^n fa^{-n}, \text{Id})) < r.$$

By assumption $\lambda^{-N} \varepsilon \leq d(fz, z) = d(f, I_k) < r$. It follows that there exists $n < N$ with $c_0^{-2} \lambda^{-1} \varepsilon < d(a^n fa^{-n}, I_k) < r$. Since $a^n z \in K$ we get

$$d(a^n fz, a^n z) = d((a^n fa^{-n})a^n z, a^n z) = d(a^n fa^{-n}, I_k) > c_0^{-2} \lambda^{-1} \varepsilon.$$
We are ready to give a close connection between box dimension and entropy.

**Proposition 8.3.** Let \( a \in A \) and \( K \subset X \) be compact with \( aK \subset K \). Then one of the following properties holds.

1. The intersection \( Ux \cap K \) of the unstable manifold \( Ux \) with \( K \) is a countable union of compact sets of box dimension zero for every \( x \in X \).

2. The restriction \( a|_K \) of the multiplication operator \( a \) to \( K \) has positive topological entropy.

**Proof.** Note, that the first possibility follows if there exists some \( \varepsilon > 0 \) such that

\[
P_y = K \cap (B_x^U y) \text{ has box dimension zero for every } y \in K.
\]

To see this, suppose \( K \cap Ux \) for \( x \in X \) is nonempty, and cover \( K \cap Ux \) by countably many sets \( P_y \) as in (8.2). Taking the union for every such \( x \) shows the first statement of the proposition.

Now, it suffices to show that if (8.2) fails for \( \varepsilon \) as in Lemma 8.2, then the topological entropy \( h_{\text{top}}(a|_K) > 0 \) is positive. Assume \( 2\varepsilon \leq r \) and that (8.2) fails for \( y \in K \). We use this to construct a sequence of \((N,\varepsilon)\)-separated sets \( F_N \subset K \). Let \( b \in (0, \dim_{\text{box}}(P_y)) \). For every \( N > 0 \) let \( K \subset P_y \) be a maximal (finite) \( \varepsilon \lambda^{-N} \)-separated set. By choice of \( b \) and the definition of box dimension in (7.2) there are infinitely many integers \( N \) with \( |F_N| \geq \lambda^b N \).

We claim that \( F_N \) is an \((N,c_2\varepsilon)\)-separated set for \( a \) restricted to \( K \). Let \( gx,hx \in F_N \) be two different points with \( g,h \in B_x^U \). By construction \( \varepsilon \lambda^{-N} \leq d(gx,hx) < 2\varepsilon \leq r \). By Lemma 8.2 applied to \( z = hx \) and \( f = gh^{-1} \) there exists a nonnegative \( n < N \) with \( d(a^n gx,a^n hx) \geq c_2\varepsilon \). Therefore \( F_N \) is \((N,c_2\varepsilon)\)-separated as claimed, and for infinitely many \( N \) we have \( s_N(a) \geq |F_N| > \lambda^{bN} \).

Finally, the definition of topological entropy in (7.5) implies that \( h_{\text{top}}(\alpha^t) \geq b \log \lambda > 0 \).

The remainder of this section is only needed for Theorem 1.6 and Theorem 10.2. For a compact set which is invariant in both directions we can also look at the stable and unstable subgroup simultaneously. Note however, that the set \( UV \) is not a subgroup of \( \text{SL}(k,\mathbb{R}) \).

**Lemma 8.4.** Let \( K \subset X \) be compact with \( aK = K \). Then \( B_x^U B_y^V \subset B_x^G \) and there exists \( \lambda > 1 \) and \( c_3 \geq 1 \) so that for any small enough \( \varepsilon > 0 \), any \( x \in X \) and \( g,h \in B_x^U B_y^V \) with \( hx \in K \), and any integer \( N \geq 1 \) with \( d(gx,hx) \geq \lambda^{-N} \varepsilon \), there exists an integer \( n \) with \( d(a^n gx,a^n hx) \geq c_3 \varepsilon \) and \( |n| < N \).
Proof. Recall that we use the right invariant metric $d$ to define the balls $B^U_r$, $B^V_r$, and $B^G_r$. Therefore, if $g_U \in B^U_r$ and $g_V \in B^V_r$, then

$$d(g_U g_V, I_k) \leq d(g_U g_V, g_V) + d(g_V, I_k) = d(g_U, I_k) + d(g_V, I_k) < 2r$$

and so $B^U_r B^V_r \subset B^G_{3r}$.

If necessary we reduce the size of $r$ such that (8.1) holds for every $g, h \in B^G_{3r}$. Assume $\varepsilon$ is small enough so that $aB^G_\varepsilon a^{-1} \subset B^G_{\varepsilon^{-1}}$. Let $\lambda > 1$ be such that $\|a x a^{-1} - I_k\| \geq \lambda\|f - I_k\|$ for $f \in U$ and $\|a^{-1} f a - I_k\| \geq \lambda\|f - I_k\|$ for $f \in V$.

Let $g, h \in B^U_r B^V_r$ and $x \in X$ be as in the lemma. Define $f = gh^{-1}$, so that $d(f, I_k) = d(g, h) \geq \lambda^{-N} \varepsilon$. Write $f = f_C f_U f_V$ and $w = \max(d(f_U, I_k), d(f_V, I_k))$. By (8.1)

$$\max(d(f_C, I_k), w) \geq c_1^{-1} d(f, I_k) \geq c_1^{-1} \lambda^{-N} \varepsilon.$$

We need to rule out the case that $d(f_C, I_k)$ is the only big term in this maximum. Clearly $f_C h = f_C h_U h_V$ and $g = g_U g_V$ are the correct decompositions in the sense of (8.1), and so $d(f_C, I_k) \leq c_1 d(f_C, h)$. By right invariance of the metric $d$ we get $d(f_C, h) = d(f_C, f)$ and again by (8.1) we get that $d(f_C, f) \leq c_1 w$. We conclude that $d(f_C, I_k) \leq c_1^2 w$, which allows us to improve (8.3) to $w \geq c_1^{-3} \lambda^{-N} \varepsilon$.

Depending on which term in $w = \max(d(f_U, I_k), d(f_V, I_k))$ achieves the maximum, we find either a positive or a negative $n$ with $|n| < N$ so that $\tilde{f} = a^n f a^{-n} = f_C \tilde{f}_U \tilde{f}_V$ satisfies $d(\tilde{f}, I_k) \in (c_3 \varepsilon, r)$ for some absolute constant $c_3$. Since $hx, a^n hx \in K$, it follows that

$$d(a^n gx, a^n hx) = d(\tilde{f} a^n hx, a^n hx) = d(\tilde{f}, I_k) \geq c_3 \varepsilon. \qed$$

Lemma 8.5. Let $a \in A$ and $K \subset X$ be compact with $aK = K$. Then one of the following properties holds.

1. The intersection $B^U_r B^V_r x \cap K$ has box dimension zero for every $x \in X$.
2. The restriction $a|_K$ of the multiplication operator $a$ to $K$ has positive topological entropy.

Proof. Suppose that $\dim_{box}(B^U_r B^V_r x \cap K) > b > 0$ for some $x \in X$. By (7.2) there exists a $\lambda^{-N} \varepsilon$-separated set $F_N$ for $N \geq 1$, which satisfies $|F_N| \geq \lambda^{b N} \varepsilon^{-b}$ for infinitely many $N$.

Let $g, h \in F_N$ with $g, h \in H$ and $g \neq h$. By Lemma 8.4 there exists an integer $n$ with $|n| < N$ such that $d(a^n gx, a^n hx) \geq c_2 \varepsilon$. This shows that $a^{-N+1} F_N \subset K$ is $(2N - 1, c_2 \varepsilon)$-separated with respect to $a$. It follows that $s_{2N-1}(a) \geq |F_N| \geq \lambda^{b N} \varepsilon^{-b}$ for infinitely many $N$, and so $h_{top}(a|_K) \geq \frac{1}{\lambda} \log \lambda > 0$. \qed
9. Upper semi-continuity of the metric entropy

For the construction of an $A$-ergodic measure $\mu$ as in Theorem 1.3 we need one more property of the metric entropy, namely upper semi-continuity with respect to the measure. More specifically we consider the metric entropy $h_\mu(a)$ as a function of the $a$-invariant measure $\mu$, where we use the weak* topology on the space of probability measures supported on a fixed compact $a$-invariant set $K$. We will show that $\limsup_{\ell \to \infty} h_{\mu_\ell}(a) \leq h_\mu(a)$ whenever $\mu_\ell$ is a sequence of $a$-invariant measures satisfying $\lim_{\ell \to \infty} \mu_\ell = \mu$. This is well known to hold for expansive maps [49, Thm. 8.2] and also for $C^\infty$ automorphisms of compact manifolds [31, Thm. 4.1]. Strictly speaking, neither of the two results applies to our case: the left multiplication by $a$ is not expansive, $X = SL(k, \mathbb{R}) / SL(k, \mathbb{Z})$ is a noncompact manifold, and there is no reason why the compact subsets $K \subset X$ we study should be manifolds at all. However, the proof for the expansive case in [49, §8.1] can be adapted to our purposes – which we provide here for the sake of completeness. We will need a few more facts about entropy and conditional entropy; see [49, Ch. 4] and [34, Ch. 2 and 4].

Let $\mu$ be a probability measure on a compact metric space $Y$. Let $A \subset B_Y$ be a $\sigma$-algebra, which is countably generated by $A_1, \ldots, A_i, \ldots$. Then the atom of $x$ is defined by

$$[x]_A = \bigcap_{i: x \in A_i} A_i \cap \bigcap_{i: x \notin A_i} X \setminus A_i,$$

and the conditional measure $\mu^A_x$ is a probability measure supported on $[x]_A$ a.s. Let $\mathcal{P}$ be a finite partition. We will need the notion of conditional entropy

$$H_\mu(\mathcal{P}|A) = \int H_{\mu^A_x}(\mathcal{P}) \, d\mu$$

and the following basic properties.

For the trivial $\sigma$-algebra $\mathcal{N} = \{\emptyset, Y\}$ the conditional entropy equals the entropy $H_\mu(\mathcal{P}|\mathcal{N}) = H_\mu(\mathcal{P})$. For two partitions $\mathcal{P}$ and $\mathcal{Q}$ we have the addition formula

$$H_\mu(\mathcal{P} \vee \mathcal{Q}|A) = H_\mu(\mathcal{P}|A) + H_\mu(\mathcal{Q}|\mathcal{P} \vee A).$$

If $\mathcal{P}$ is finer than $\mathcal{Q}$ and $\mathcal{C} \subseteq \mathcal{A}$ is another countably generated $\sigma$-algebra, then

$$H_\mu(\mathcal{Q}|A) \leq H_\mu(\mathcal{P}|\mathcal{C}).$$

Finally, the conditional entropy $H_\mu(\mathcal{P}|A)$ vanishes if and only if there exists a nullset $N$ such that $[x]_A \setminus N$ is contained in one of the elements of $\mathcal{P}$ for every $x \in Y \setminus N$.

Suppose $T: Y \to Y$ is measure preserving and invertible. Then the metric entropy (7.6) of $T$ with respect to a finite partition $\mathcal{Q}$ can also be written as

$$h_\mu(T, \mathcal{Q}) = H_\mu\left(\mathcal{Q} \bigvee_{n=1}^\infty T^{-n}\mathcal{Q}\right);$$

see [49, Thm. 4.14].
We will also need the dynamical version of relative entropy. Suppose $\mathcal{A}$ is a countably generated $\sigma$-algebra that satisfies $T\mathcal{A} = \mathcal{A}$. We define

\begin{equation}
(9.1) \quad h_\mu(T, Q|\mathcal{A}) = H_\mu \left( Q \left| \bigvee_{n=1}^\infty T^{-n}Q \lor \mathcal{A} \right. \right);
\end{equation}

then

\begin{equation}
(9.2) \quad h_\mu(T, \mathcal{P} \lor Q) = h_\mu(T, \mathcal{P}) + h_\mu \left( T, Q \left| \bigvee_{i=-\infty}^\infty T^i\mathcal{P} \right. \right).
\end{equation}

The entropy with respect to an invariant measure is defined as a supremum over all finite partitions; see (7.6). For this reason the following general principle will be helpful.

**Lemma 9.1.** Let $a \in A$ and $K \subset X$ be compact with $aK \subset K$. Let $\mu$ be an $a$-invariant measure supported on $K$. There exists a sequence of finite partitions $Q_m$ of $K$ which satisfies for all $m$ that $Q_{m+1}$ is finer than $Q_m$. The boundaries of the elements of $Q_m$ are $\mu$-null sets, and the $\sigma$-algebra $\bigcup_{m=1}^\infty Q_m$ equals the Borel $\sigma$-algebra $\mathcal{B}_K$ of $K$. Furthermore, $h_\mu(a) = \lim_{m \to \infty} h_\mu(a, Q_m)$.

**Proof.** Let $x \in X$ and define $f(y) = d(x, y)$. Then the measure $f_*\mu$ is a probability measure on $\mathbb{R}^+$, there exist arbitrarily small $\varepsilon > 0$ such that $f_*\mu(\{\varepsilon\}) = 0$, and so $\mu(\partial B_\varepsilon(x)) = 0$.

For $m > 0$ we can cover $K$ with finitely many $\varepsilon$-balls with $\varepsilon < 1/m$ whose boundaries are null sets. Let $\mathcal{P}_m$ be the partition generated by these balls. For $P \in \mathcal{P}_m$ the boundary $\partial P$ is contained in the union of the boundaries of the balls; thus it is a null set. To ensure that the sequence of partitions is getting finer we define $Q_m = \bigcup_{i=1}^m \mathcal{P}_m$. It follows that every $Q \in Q_m$ has a null set as boundary, and that $Q_m \not\subset \mathcal{B}_K$ for $m \to \infty$.

The last statement follows from [34, §4, Thm. 3].

**Proposition 9.2.** Let $a \in A$, and $K \subset X$ be compact with $aK \subset K$. For every finite partition $\mathcal{P}$ of $K$ into measurable sets with small enough diameters and for any $a$-invariant measure $\mu$ supported on $K$, $h_\mu(a) = h_\mu(a, \mathcal{P})$.

**Proof.** Every $T$-invariant measure $\mu$ which is supported by $K$ is in fact supported on the set $K' = \bigcap_{a \geq 0} a^iK$. Clearly, $K'$ is compact and satisfies $aK' = K'$. Since a partition of $K$ into small sets induces a partition of $K'$ into small sets, we can assume without loss of generality that $K$ satisfies $aK = K$.

Recall that there exists $r = r(K) > 0$ with $d(x, gx) = d(I_k, g)$ whenever $d(I_k, g) < r$ and $x \in K$. Let $\varepsilon < r$ be as in Lemma 8.4, let $\mathcal{P}$ be a finite partition into measurable sets with diameter less than $\delta$ (to be specified later), and define the $\sigma$-algebra $A = \bigcup_{i=-\infty}^\infty a^{-i}\mathcal{P}$ generated by the orbit of $\mathcal{P}$.

Let $x, y \in K$ belong to the same atom of $\mathcal{A}$; in other words suppose that for all $i \in \mathbb{Z}$ the images $a^ix, a^iy \in P_i$ belong to the same partition element.
of \( P \). We claim that (for small enough \( \delta \)) this implies \( x = f_C y \) for some small \( f_C \in C \).

Let \( x = f y \) with \( f \in B^G_r \) and suppose \( f \notin C \). Let \( f = f_C f_U f_V \) with \( f_C \in C, f_U \in U, \) and \( f_V \in V \). For small enough \( \delta > 0 \) we have \( d(f_C, I_k) < c_3 \varepsilon/2, d(f_U, I_k) < r \) and \( d(f_V, I_k) < r \). Let \( z = f_U f_V y = f_C^{-1} x; \) then \( z \in B^U_r B^V_r y \). Since \( a f_C = f_C a \) we have \( d(a^n x, a^n z) = d(f_C, I_k) < c_3 \varepsilon/2 \) for all \( n \).

By Lemma 8.4 there exists some integer \( n \) with \( d(a^n z, a^n y) \geq c_3 \varepsilon \). We assume \( \delta < c_3 \varepsilon/2; \) then

\[
d(a^n x, a^n y) \geq d(a^n z, a^n y) - d(a^n x, a^n z) > c_3 \varepsilon/2
\]

shows that \( a^n x \) and \( a^n y \) cannot belong to the same partition element of \( P \).

This contradiction shows the claim.

Suppose \( Q = \{Q_1, \ldots, Q_m\} \) is one of the partitions of Lemma 9.1. We remove all the boundaries of the elements of the partition and obtain a partition modulo \( \mu \) into open sets of small diameter.

By (9.2) we have

\[
h_{\mu}(T, Q) \leq h_{\mu}(T, P \vee Q) = h_{\mu}(T, P) + h_{\mu}(T, Q|A),
\]

where \( h_{\mu}(T, Q|A) \) is the relative entropy as in (9.1). We will show that this last term vanishes, which together with Lemma 9.1 will conclude the proof of Proposition 9.2.

Let \( B \subset X \) be measurable. By Poincaré recurrence, there exists a null set \( N \) such that for every \( x \notin N \) and \( x \in B \) there exists some \( n \geq 1 \) with \( a^n x \in B \).

We apply this simultaneously to the countable family of sets

\[
B_{i,j,g_C,\ell} = \{x : B_{1/\ell}(x) \subset Q_i \cap g_C^{-1} Q_j\}
\]

for \( g_C \in C \cap \text{SL}(k, \mathbb{Q}), Q_i, Q_j \in Q \) and \( \ell \geq 1 \). To show that the relative entropy

\[
h_{\mu}(T, Q|A) = H_{\mu}(Q|\tilde{A}) \text{ with } \tilde{A} = \bigcup_{n=1}^{\infty} T^{-n} Q \vee A(\text{mod } \mu)
\]

vanishes, we have to show that for \( x, y \notin N \) which are in the same atom with respect to \( \tilde{A} \) and satisfy \( x \in Q_i \in Q \) and \( y \in Q_j \in Q \), in fact, \( i = j \) holds. Since \( x \) and \( y \) belong to the same atom with respect to \( A \), we know from the above claim that \( y = f_C x \) for some small \( f_C \in C \). Therefore, \( x \in Q_i \cap f_C^{-1} Q_j \) and there exists some rational \( g_C \) close to \( f_C \) with \( x \in Q_i \cap g_C^{-1} Q_j \). Furthermore, we can ensure that \( B_{1/\ell}(x) \subset Q_i \cap g_C^{-1} Q_j \), \( d(g_C, f_C) < 1/\ell, \) and \( 1/\ell < r \). It follows that \( x \in B_{i,j,g_C,\ell} \). By construction of \( N \) there exists \( n > 0 \) with \( a^n x \in B_{i,j,g_C,\ell} \). Therefore \( a^n x \in Q_i \) and \( B_{1/\ell}(a^n x) \subset g_C^{-1} Q_j \). From

\[
d(g_C^{-1} f_C a^n x, a^n y) = d(g_C^{-1} f_C, I_k) < 1/\ell < r
\]

we see that \( g_C^{-1} f_C a^n x \in g_C^{-1} Q_j \). Since \( a \) commutes with \( f_C, f_C a^n x = a^n y \in Q_j \).

We have shown that \( x \in a^{-n} Q_i \) and \( y \in a^{-n} Q_j \). Since \( a^{-n} Q_i, a^{-n} Q_j \) belong to \( \tilde{A} \) and \( x, y \) are assumed to belong to the same atom with respect to \( \tilde{A} \), it follows that \( i = j \) as claimed.
The above proposition has the following important consequence.

**Corollary 9.3.** Let \( a \in A \) and \( K \subset X \) be compact with \( aK \subset K \). Then the metric entropy \( h_\mu(a|K) \) is upper semi-continuous with respect to the measure \( \mu \), i.e., for every \( a \)-invariant \( \mu \) and every \( \varepsilon > 0 \) there is a neighborhood \( U \) of \( \mu \) in the weak* topology of probability measures on \( K \) such that \( h_\nu(a) \leq h_\mu(a) + \varepsilon \) for every \( a \)-invariant \( \nu \in U \).

**Proof.** As in the proof of Lemma 9.1 we can find a partition \( P \) of \( K \) whose elements have small enough diameter to satisfy Proposition 9.2 and whose boundaries are null sets with respect to \( \mu \). Therefore \( h_\nu(a) = h_\nu(a, P) \) for every \( a \)-invariant measure \( \nu \) supported on \( K \). Let \( \varepsilon > 0 \). By the definition of entropy there exists \( N \geq 1 \) with

\[
\frac{1}{N} H_\mu \left( \bigvee_{n=0}^{N-1} a^{-n} P \right) < h_\mu(a, P) + \varepsilon/2.
\]

Since the sets in the partition \( Q = \bigvee_{n=0}^{N-1} a^{-n} P \) all have boundaries which are null sets with respect to \( \mu \), there exists a weak* neighborhood \( U \) of \( \mu \) such that \( \nu(Q) \) is very close to \( \mu(Q) \) for every \( Q \in Q \). The entropy of the partition \( Q \) depends only on the measures of the elements of \( Q \); therefore we can make sure that

\[
\frac{1}{N} \left| H_\nu(Q) - H_\mu(Q) \right| < \varepsilon/2.
\]

For any \( a \)-invariant \( \nu \in U \),

\[
h_\nu(a) = h_\nu(a, P) \leq \frac{1}{N} H_\nu(Q) \leq \frac{1}{N} H_\mu(Q) + \varepsilon/2 < h_\mu(a, P) + \varepsilon,
\]

where we used Proposition 9.2 for \( \nu \) and \( \mu \). Furthermore, \( h_\nu(a, P) \) is the infimum over \( \frac{1}{M} H_\nu(\bigvee_{n=0}^{M-1} a^n P) \) by subadditivity [49, Thm. 4.10].

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10. **Transversal Hausdorff dimension for the set of points with bounded orbits**

In this section we apply Theorem 1.3 to prove two theorems about sets with bounded orbits.

For a unimodular lattice \( \Lambda \subset \mathbb{R}^k \) we define

\[
\delta_{\mathbb{R}^k}(\Lambda) = \min_{y \in \Lambda \setminus \{0\}} \|y\|.
\]

Clearly, every point \( x = m \text{SL}(k, \mathbb{Z}) \) with \( m \in \text{SL}(k, \mathbb{R}) \) can be identified with the unimodular lattice generated by the columns of \( m \). By this identification \( \delta_{\mathbb{R}^k} \) becomes a positive continuous function on \( X \) with the property that the preimages \( K_\rho = \delta_{\mathbb{R}^k}^{-1}[\rho, \infty) \) are compact sets for every \( \rho > 0 \) by Mahler’s criterion. In other words \( B \subset X \) is bounded if and only if \( \inf_{x \in B} \delta_{\mathbb{R}^k}(x) > 0 \).
A nonempty subset \( \Sigma' \subset \Sigma \) is a cone if \( \Sigma' \) is convex and satisfies \( rt \in \Sigma' \) whenever \( r > 0 \) and \( t \in \Sigma' \).

**Theorem 10.1.** Let \( X = \text{SL}(k, \mathbb{R})/\text{SL}(k, \mathbb{Z}) \) with \( k \geq 3 \), and let \( \Sigma' \) be an open cone in \( \Sigma \). Define

\[
D = \{ x \in X : \inf_{t \in \Sigma'} \delta_{\mathbb{R}^k}(\alpha^t x) > 0 \}
\]

to be the set of points with bounded \( \Sigma' \)-orbits. Then for every \( t \in \Sigma' \) and \( x \in X \) the \( \alpha^t \)-unstable manifold \( U_x \) through \( x \) intersects \( D \) in a set \( D \cap U_x \) of Hausdorff dimension zero. In fact, \( D \cap U_x \) is a countable union of sets with upper box dimension zero.

**Proof.** For \( \rho > 0 \) we define the compact set

\[
(10.1) \quad D_\rho = \{ x \in X : \inf_{t \in \Sigma'} \delta_{\mathbb{R}^k}(\alpha^t x) \geq \rho \}.
\]

Clearly \( D = \bigcup_{n=1}^{\infty} D_{1/n} \). Let \( t \in \Sigma' \), \( a = \alpha^t \), and \( x \in X \). Then \( aD_\rho \subset D_\rho \). By Proposition 8.3 there are two possibilities; \( D_\rho \cap U_x \) is a countable union of compact sets of box dimension zero, or \( a|_{D_\rho} \) has positive topological entropy. If the first possibility takes place for all \( \rho > 0 \), the theorem follows from (7.3) and (7.4).

We will show that the second possibility cannot happen ever. Suppose \( a|_{D_\rho} \) has positive topological entropy. By the variational principle ([49, Thm. 8.6]) there exists an \( a \)-invariant measure \( \nu \) supported on \( D_\rho \) with positive metric entropy \( h_\nu(a) > 0 \). However, we need to find an \( A \)-ergodic measure with this property in order to get a contradiction to Theorem 1.3.

Since \( \Sigma' \subset \Sigma \) is open we can find a basis \( t_1, \ldots, t_{k-1} \in \Sigma' \) of \( \Sigma \). By construction \( K = D_\rho \) is compact and satisfies \( \alpha^s K \subset K \) for all \( s \in \mathbb{R}^+ t_1 + \cdots + \mathbb{R}^+ t_{k-1} \). For \( N > 0 \) the measure

\[
\nu_N = \frac{1}{N^{k-1}} \int_0^N \cdots \int_0^N (\alpha^{s_1 t_1 + \cdots + s_{k-1} t_{k-1}})_s \nu \, ds_1 \cdots ds_{k-1}
\]

is supported on \( K \) and \( a \)-invariant. Since entropy is affine ([49, Thm. 8.1]) and upper semi-continuous by Corollary 9.3 with respect to the measure, entropy with respect to a generalized convex combination of measures is the integral of the entropies. In particular \( h_{\nu_N}(T) = h_\nu(T) \).

Let \( \mu \) be a weak* limit of a subsequence of \( \nu_N \). From the definition of \( \nu_N \) it follows that \( \mu \) is \( A \)-invariant. It is also clear that \( \mu \) is supported on \( K \). From upper semi-continuity, the entropy \( h_\mu(a) \geq h_\nu(a) > 0 \) is positive. The ergodic decomposition (4.2) of \( \mu \) defines \( \mu \) as a generalized convex combination of \( A \)-ergodic measures \( \mu_\tau \), which have, almost surely, support contained in \( K \). Since \( h_\mu(T) > 0 \), there exists some \( A \)-ergodic measure \( \mu_\tau \) with \( h_{\mu_\tau}(T) > 0 \) and support in \( K \). This contradicts Theorem 1.3 and concludes the proof of Theorem 10.1.
Let $D \subseteq X$ be $A$-invariant. We say $D$ has transversal box dimension zero if \( \{g \in B^U_f : gx \in D\} \) and \( g_x = 1 \) for \( i = 1, \ldots, k \) has box dimension zero for all \( x \in D \). (Note that the particular shape of the set used here does not matter as long as this set is still transversal to the subgroup \( A \).) It is easy to check, that an \( A \)-invariant set \( D \) with transversal box dimension zero has box dimension \( k - 1 \) (unless \( D \) is empty).

**Theorem 10.2.** Let \( X = \text{SL}(k, \mathbb{R}) / \text{SL}(k, \mathbb{Z}) \) with \( k \geq 3 \), let \( A \subset \text{SL}(k, \mathbb{R}) \) be the subgroup of positive diagonal matrices. Define

\[
D = \{ x \in X : \inf_{a \in A} \delta_{\mathbb{R}^+}(ax) > 0 \}
\]

to be the set of points with bounded \( A \)-orbits. Then \( D \) is a countable union of sets with transversal box dimension zero and has Hausdorff dimension \( k - 1 \).

Clearly \( D \) is \( A \)-invariant, and nonempty since it contains every periodic \( A \)-orbit.

**Proof.** As before we define the \( A \)-invariant compact sets \( D_\rho \) as in (10.1) with \( \Sigma' = \Sigma \). Pick an element \( a = \alpha^t \in A \) with \( t \in \Sigma \), \( t_i \neq t_j \) for \( i \neq j \). Then the corresponding central subgroup equals \( C = A \) and \( B^U_r B^V_r \) is transversal to \( A \). Let \( x \in X \) and \( \rho > 0 \).

We give some conditions on \( r > 0 \). Our first restriction is that \( B^u_r(x) \) and \( B^V_{\rho r} \) should be isometric. Let \( O = B^A_r \times B^U_r \times B^V_r \) and use the metric

\[
d_O((f_C, f_U, f_V), (g_C, g_U, g_V)) = \max(d(f_C, g_C), d(f_U, g_U), d(f_V, g_V)).
\]

Furthermore, define \( \psi : O \to \text{SL}(k, \mathbb{R}) \) by \( \psi(f_C, f_U, f_V) = f_C f_U f_V \) and assume \( \psi \) is invertible and Lipschitz in both directions (as in (8.1)).

Let \( P = \{(f_C, f_U, f_V) \in O : \psi(f_C, f_U, f_V)x \in D_\rho \} \). Since \( D_\rho \) is \( A \)-invariant, the set \( P' = P \cap \{(1_k) \times B^U_r \times B^V_r \} \) determines \( P = \{(f_C, f_U, f_V) \in O : (1_k, f_U, f_V) \in P' \} \). Clearly \( \psi(P')x = (B^U_r B^V_r x) \cap D_\rho \). By Lemma 8.5 there are two possibilities; \( P' \) has box dimension zero or \( a \) has positive topological entropy when restricted to \( D_\rho \). As in the proof of Theorem 10.1 the latter contradicts Theorem 1.3. Therefore \( D_\rho \) has transversal box dimension zero. \( \square \)

11. The set of exceptions to Littlewood’s Conjecture

For any \( u, v \in \mathbb{R} \), define \( \tau_{u,v} \) to be the point

\[
\tau_{u,v} = \begin{pmatrix} 1 & 0 & 0 \\ u & 1 & 0 \\ v & 0 & 1 \end{pmatrix} \text{ SL}(3, \mathbb{Z});
\]

in other words, \( \tau_{u,v} \) is the point in \( X \) corresponding to the lattice in \( \mathbb{R}^3 \) generated by \((1, u, v), (0, 1, 0), \) and \((0, 0, 1) \). The following well-known proposition
gives the reduction of Littlewood’s conjecture to the dynamical question which we studied in Section 10; see also [24, §2] and [46, §30.3]. We include the proof for completeness.

**Proposition 11.1.** The tuple \((u, v)\) satisfies

\[
\liminf_{n \to \infty} n \langle nu \rangle \langle nv \rangle = 0,
\]

if and only if the orbit \(A^+ \tau_{u,v}\) is unbounded where \(A^+\) is the semigroup

\[
A^+ = \left\{ \begin{pmatrix} e^{-r-s} & 0 & 0 \\ e^r u & e^r & 0 \\ e^s v & 0 & e^s \end{pmatrix} : \text{for } r, s \in \mathbb{R}^+ \right\}.
\]

**Proof.** By the the properties of \(\delta_{\mathbb{R}^3}\) we have to show for \((u, v) \in \mathbb{R}^2\) that (11.1) holds if and only if \(\inf_{a \in A^+} \delta_{\mathbb{R}^3}(a \tau_{u,v}) = 0\).

Suppose \(\varepsilon > 0\) and there exists \(a \in A^+\) with \(\delta_{\mathbb{R}^3}(a \tau_{u,v}) < \varepsilon\). Then

\[
a \tau_{u,v} = \begin{pmatrix} e^{-r-s} & 0 & 0 \\ e^r u & e^r & 0 \\ e^s v & 0 & e^s \end{pmatrix} \text{SL}(3, \mathbb{Z})
\]

and by definition of \(\delta_{\mathbb{R}^3}\) there exists nonzero \((n, m_1, m_2) \in \mathbb{Z}^3 \) with

\[
\left\| \begin{pmatrix} ne^{-r-s} \\ ne^r u + m_1 e^r \\ ne^s v + m_2 e^s \end{pmatrix} \right\| < \varepsilon.
\]

Taking the product of all three entries of this vector we find that

\[
|ne^{-r-s}(ne^r u + m_1 e^r)(ne^s v + m_2 e^s)| = |n(nu + m_1)(nv + m_2)| < c \varepsilon^3
\]

is small (c depends only on the norm used in \(\mathbb{R}^3\)), and so (11.1) follows. Note that \(n \neq 0\) since otherwise the lower two entries in the vector cannot be small.

Suppose now that (11.1) holds for \((u, v)\). Let \(\varepsilon > 0\) and find \(n > 0 \) and \((m_1, m_2) \in \mathbb{Z}^2 \) with \(|n(nu + m_1)(nv + m_2)| < \varepsilon^5\). We would like to have additionally that

\[
\max(|nu + m_1|, |nv + m_2|) < \varepsilon.
\]

Suppose this is not true, and assume without loss of generality that \(|nv + m_2| \geq \varepsilon\) and \(|n(nu + m_1)| < \varepsilon^4\). Then by Dirichlet’s theorem there exists an integer \(q < 1/\varepsilon\) so that \(\langle qnu \rangle < \varepsilon\). It follows that \(|qn(nu + qm_1)| < \varepsilon^2\), and \(|qnv + m'_2| < \varepsilon\) for some \(m'_2 \in \mathbb{Z}\). In other words when we replace \(n\) by \(nq\) and \(m_1, m_2\) by \(qm_1\) and \(m'_2\) respectively, we see that (11.2) and \(|n(nu + m_1)(nv + m_2)| < \varepsilon^3\) hold simultaneously. Therefore we can find \(r > 0\) and \(s > 0\) with \(e^r|nu + m_1| = \varepsilon\) and \(e^s|nv + m_2| = \varepsilon\). (If one of the expressions vanishes, we use some large \(r\), resp. \(s\), instead.) Then \(e^{-r-s}n < \varepsilon\) and \(\delta_{\mathbb{R}^3}(a \tau_{u,v}) < c \varepsilon\) follows. \(\square\)
Proof of Theorem 1.5. By Proposition 11.1 the set $\Xi$ is embedded by the map $(u,v) \mapsto \tau_{u,v}$ to the set $D$ with $A^+$-bounded orbits. We apply Theorem 10.1 with $\Sigma' = \{(-r - s, r, s) : r, s > 0\}$. Therefore $D$ intersects every unstable manifold of $\alpha^t$ in a set of Hausdorff dimension zero where $t = (-2, 1, 1)$. Note that the unstable manifold of $\alpha^t$ through $I_k \text{SL}(3, \mathbb{Z})$ is the image of $\tau$. It follows that $\Xi$ has Hausdorff dimension zero, and similarly that $\Xi$ is a countable union of sets with box dimension zero.

Proof of Theorem 1.6. We apply Theorem 10.2 and set $\Xi_k = D$. Suppose $m \notin \Xi_k$; then $\delta_{2^k}(am) < \varepsilon$ for some $a \in A$. By definition of $\delta_{2^k}$ there exists some $n \in \mathbb{Z}^k$ such that $\left\| \begin{bmatrix} a_{11}m_1(n) \\ \vdots \\ a_{kk}m_k(n) \end{bmatrix} \right\| < \varepsilon$ and (1.3) follows.

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