

Rank-one group actions with simple mixing \mathbb{Z} -subactions

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ABSTRACT. Let G be a countable Abelian group with \mathbb{Z}^d as a subgroup so that G/\mathbb{Z}^d is a locally finite group. (An Abelian group is locally finite if every element has finite order.) We can construct a rank one action of G so that the \mathbb{Z} -subaction is 2-simple, 2-mixing and only commutes with the other transformations in the action of G .

Applications of this construction include a transformation with square roots of all orders but no infinite square root chain, a transformation with countably many nonisomorphic square roots, a new proof of an old theorem of Baxter and Akcoglu on roots of transformations, and a simple map with no prime factors. The last example, originally constructed by del Junco, was the inspiration for this work.

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1. Introduction and results

Ornstein’s rank-one mixing argument [Orn67] has been refined over the years, and its ideas are used often in the literature to construct interesting examples. Notably, del Junco [Jun98] constructed a measure preserving action of $\mathbb{Z} \oplus \bigoplus_{i=1}^{\infty} \mathbb{Z}_2$. Using an Ornstein style argument along with a joinings argument, he showed the transformation T corresponding to the \mathbb{Z} -subaction was weak mixing, simple and commuted only with the other transformations in the action. del Junco was then able to argue that T was a simple map with no prime factors. This paper provides an extension of del Junco’s construction to a certain class of abelian groups.

Recall that if every element of an abelian group G has finite order then G is a locally finite group.

Theorem 1 (Main Result). *Let G be a countable abelian group with subgroup \mathbb{Z}^d ($d \geq 1$), such that G/\mathbb{Z}^d is a locally finite group. Then there exists a rank-one action of G so that the transformation T corresponding to $(1, 0, 0, \dots, 0)$ in \mathbb{Z}^d is mixing, simple, and only commutes with the other transformations in the group, i.e., $C(T) = G$.*

We note that, in particular, the theorem is valid for groups $G = \mathbb{Z}^d \oplus H$ where H is a locally finite group, possibly finite, or even the trivial group. The theorem is proved in Section 3.

The main theorem allows the construction of simple transformations T with centralizer $C(T)$ prescribed in advance. Since T is simple, this gives us some control over the roots and factors of T . We’ll detail some examples, new and previously known, that can be constructed in this way. First, and most significantly we answer a question posed by King in [Kin00].

1.1. A transformation with square roots of all orders but no infinite square root chain. Let S and T be measure preserving transformations on the same space. If $S^2 = T$ we say S is a square root of T and write $T \rightarrow S$. If $S^{2^n} = T$ (T has a 2^n th root, S) we can find a square root chain for T of length n :

$$T \rightarrow S^{2^{n-1}} \rightarrow S^{2^{n-2}} \rightarrow \dots \rightarrow S.$$

J. King has been investigating the problem of embedding the generic transformation into actions of the rationals [Kin00]. A significant obstruction to embedding the generic transformation in an action of the dyadic rationals is the necessity of existence of an infinite square root chain,

$$T \rightarrow T^{\frac{1}{2}} \rightarrow T^{\frac{1}{4}} \rightarrow T^{\frac{1}{8}} \rightarrow \dots.$$

In [Kin00] King asked: “Is there a transformation with square roots of all orders but no infinite square root chain?” We answer this question affirmatively using an appropriate group.

Definition 1 (Carry group of r). Let $r = \{r_i\}_{i=1}^\infty$ be any countable sequence of natural numbers. Define G , the carry group of r , to consist of the elements of the Cartesian product

$$\mathbb{Z} \times \mathbb{Z}_{r_1} \times \mathbb{Z}_{r_2} \times \cdots,$$

where all but finitely many entries are zero, together with an operation $+_G$ defined by

$$(a_{\mathbb{Z}}, a_1, a_2, \dots) +_G (b_{\mathbb{Z}}, b_1, b_2, \dots) := \left(a_{\mathbb{Z}} + b_{\mathbb{Z}} + \sum_{i=1}^\infty \left\lfloor \frac{a_i + b_i}{r_i} \right\rfloor, a_1 +_{r_1} b_1, a_2 +_{r_2} b_2, \dots \right).$$

That is, addition in the \mathbb{Z} coordinate and addition modulo r_i in the \mathbb{Z}_{r_i} coordinate with a possible carry of 1 into the \mathbb{Z} coordinate.

Let G be the carry group of $r = (2, 4, 8, 16, \dots, 2^n, \dots)$, a group discussed by King in [Kin00]. Note that G/\mathbb{Z} is a locally finite group. Apply Theorem 1 to this group to obtain a transformation T , with G as its centralizer. T is the transformation in the action corresponding to the element $(1, 0, 0, \dots)$. The transformation T has square roots of all orders, since the group element $(0, 0, \dots, 0, 1, 0, \dots)$, where the 1 is in the \mathbb{Z}_{2^n} coordinate, is a 2^n th root of $(1, 0, 0, \dots)$. There are, however, no infinite square root chains in G . An infinite square root chain in G must have some nonzero values in a coordinate other than \mathbb{Z} , say the \mathbb{Z}_{2^n} coordinate. The values in that coordinate would form a nontrivial infinite square root chain in \mathbb{Z}_{2^n} which does not exist. (One easy way to prove this is by induction). Thus, T answers King’s question.

Analogously we can produce transformations with q th roots of all orders but no infinite q th root chains, where q is any positive integer except 1.

1.2. A simple map with no prime factors. Applying Theorem 1 to $G = \mathbb{Z} \oplus \bigoplus_{i=1}^\infty \mathbb{Z}_2$ we obtain a simple map with no prime factors as originally constructed by del Junco. See [Jun98] for details.

1.3. A transformation with $C(T) = \mathbb{Q}$ or \mathbb{Z}^d . Each application of our theorem produces a transformation with countable but (usually) nontrivial centralizer. When $G = \mathbb{Z}$ we have constructed Ornstein’s rank-one mixing transformation that only commutes with its powers [Orn67]. When $G = \mathbb{Z}^d$ we obtain a transformation with centralizer \mathbb{Z}^d , and when $G = \mathbb{Q}$, the transformation has the rationals as centralizer. The latter is possible because \mathbb{Q}/\mathbb{Z} is a locally finite group.

1.4. Transformations with a fixed set of roots.

Question 1. Can you construct a transformation with only a specified set of roots?

Akcoglu and Baxter [AB69] published an interesting theorem on this topic in 1969. We offer an alternate proof, using Theorem 1.

Theorem 2 (Akcoglu and Baxter). *Let P be any set of primes. There exists a weak mixing transformation T , so that T has a p th root if and only if all prime factors of p are in P .*

Proof. Consider the set H of rational numbers (in lowest form) whose denominator is 1 or has all its prime factors in P . H is a subgroup of \mathbb{Q} that includes \mathbb{Z} . Apply Theorem 1 to obtain a transformation T which has a p th root if and only if $1/p$ is in H , that is, if all prime divisors of p are in P . \square

The transformation T we constructed in this proof is actually mixing. Also notice our T will have an infinite number of roots. Define $T \xrightarrow{p} S$ to mean S is a p th root of T . Then T also has infinite root chains. For example, if $p \in P$, the following is a p th root chain:

$$T \xrightarrow{p} T^{1/p} \xrightarrow{p} T^{1/p^2} \xrightarrow{p} T^{1/p^3} \dots$$

One could also use Theorem 1 and a carry group to construct a transformation that satisfies Theorem 2 yet has no infinite root chains.

1.5. A transformation with countably many nonisomorphic square roots.

Let G be the root group of $r = \{2, 2, 2, 2, \dots\}$. Apply the main theorem to G to obtain a simple mixing transformation T corresponding to $(1, 0, \dots)$ in G . T has countably many square roots, each corresponding to an element $(0, \dots, 1, 0, \dots)$ in G . Let S and Q be two distinct square roots of T . Assume ϕ is an isomorphism between S and Q . Then $\phi S = Q\phi$ which implies that $\phi S^2 = Q^2\phi$, and thus $\phi T = T\phi$. Since ϕ commutes with T it is in $C(T)$, which is isomorphic to the commutative group G . S and Q are also in this commutative group $C(T)$ so $\phi S = Q\phi$ implies $S = Q$. This contradicts the assumption that S and Q were distinct. Thus any two square roots of T are nonisomorphic.

1.6. Future directions. Here are a few of the many natural questions arising from this work.

Question 2. Is the full rank-one group action, constructed in Theorem 1, a mixing action?

In separate work [Mad] we constructed rank-one mixing \mathbb{Z}^d actions so that all times are simple. It is possible to extend our main theorem to ensure the \mathbb{Z}^d -subaction is mixing with all times simple.

Question 3. Given a group G and subgroup H can you construct an action G so that the H -subaction is simple and only commutes with the entire group action?

Our theorem gives an answer for the case $H = \mathbb{Z}$ and a certain class of countable abelian groups G . A more general construction, especially when G is nonabelian, could produce very interesting examples. del Junco outlines several in [Jun98].

Question 4. Which groups have rank-one mixing actions?

Our techniques could lead us to a class of countable abelian groups which have such actions. Can this be extended to any nonabelian groups? To solvable groups?

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2. Definitions and preliminaries

We use $:=$ to indicate definition (or assignment). All our transformations are on $\mathbb{X} := ([0, 1), \beta, \mu)$ which is isomorphic to the unit interval with the Borel σ -algebra and Lebesgue measure. Use $|A|$ to denote absolute value when A is a number and cardinality when A is a finite set. Denote the set of integers $\{0, 1, 2, \dots, n\}$ by $[0, n]$. Similarly $[0, n)$ denotes $\{0, 1, 2, \dots, n - 1\}$. Any other notations are introduced as needed. Unnumbered definitions are in italics.

2.1. Joinings and simplicity. Let T be a measure preserving transformation on the measure space \mathbb{X} . If λ is a $T \times T$ invariant measure on (X^2, β^2) such that for all $A \in \beta$,

$$\lambda(X \times A) = \lambda(A \times X) = \mu(A),$$

then λ is a *self-joining* of T . Every transformation has self-joinings such as product measure, $\mu \times \mu$, or μ lifted onto the diagonal, denoted Δ . More precisely, $\Delta := \mu \circ J^{-1}$ where $J: X \rightarrow X \times X$ and $J(x) = (x, x)$. If S is any measure preserving transformation that commutes with T then $(I \times S)\Delta$ is also a self-joining. These are called *graph joinings*, because $(I \times S)\Delta$ is μ lifted onto the graph of S .

We say T is *simple* if the only ergodic self-joinings of T are product measure or graph joinings.

These definitions in the literature are properly named 2-fold self-joinings, and 2-simple. As we will not discuss higher-order joinings we have opted for the simpler names.

2.2. Rank-one group actions. We follow the definitions and notations for group actions as in [Jun98], [PR91] and [YJ00]. All our groups are amenable, countable and have the discrete topology. In our main theorem G is also abelian, though definitions in this subsection do not assume the group is commutative. Let \mathcal{L} be a homomorphism from the group G into the set of invertible measure preserving transformations on X .

$$\begin{aligned} \mathcal{L}: G &\rightarrow M(X) \\ g &\rightarrow \mathcal{L}^g. \end{aligned}$$

We call \mathcal{L} a measure preserving action of the group G . The range of \mathcal{L} is denoted $\text{times}(\mathcal{L})$. An individual transformation \mathcal{L}^g is called a *time* of the action. Let $C(\mathcal{L})$ denote the centralizer, that is, all invertible measure-preserving transformations that commute with all of the times of the action.

Originally, towers of rank-one transformations were indexed by intervals in \mathbb{Z} . Ferenczi, in [Fer85], credits Thouvenot with the idea that towers could be indexed by other sets. Ferenczi used a special Følner sequence in \mathbb{Z} to define a \mathbb{Z} action and called it funny rank-one. Generalizing this idea to other groups we have rank-one group actions.

Definition 2 (Rank-one group action). \mathcal{L} is called rank-one (with respect to a Følner sequence $\{A_n\}$ in G) if:

1. For all $n \in \mathbb{N}$ there is a partition P_n of X , $P_n = \{E_g^n \mid g \in A_n\} \cup \{X \setminus X_n\}$, where $X_n = \cup_{g \in A_n} E_g^n$.
2. $\mu(X \setminus X_n) \xrightarrow{n} 0$.
3. $P_n \xrightarrow{n} \epsilon$.
4. $\mathcal{L}_g E_h^n = E_{gh}^n$ when $h \in A_n \cap g^{-1}A_n$.

$P_n \xrightarrow{n} \epsilon$ means that given measurable B in X , for all n there is a set B_n made up of a union of sets in P_n and $\mu(B \Delta B_n) \xrightarrow{n} 0$. For each $g \in A_n$, E_g^n is an interval of length a_n where $a_n |A_n| \xrightarrow{n} 1$.

2.3. Cutting and stacking rank-one actions. We will describe a cutting and stacking style of construction that produces a rank-one action \mathcal{L} of G , given a special kind of Følner sequence in G called almost-tiling. It is not known if all rank-one actions can be obtained in this manner.

Definition 3 (Almost-tiling Følner sequences). An almost-tiling Følner sequence consists of two sequences of sets $(\{A_n\}, \{C_n\})$ in G such that:

1. A_n is a Følner sequence.
2. $A_n C_n = \cup_{c \in C_n} A_n c$ is a disjoint union.
3. $A_n C_n \subset A_{n+1}$.
4. $\prod_{n=1}^{\infty} \frac{|A_{n+1}|}{|A_n C_n|} < \infty$.

This is an inductive construction. At stage n (also called time- n) of the construction, the n -tower consists of $|A_n|$ levels, each a distinct interval in the real line that is left-closed and right-open. Levels are indexed by A_n and all have equal length. An individual level is denoted E_a^n , the a th level ($a \in A_n$) of the n -tower.

The action \mathcal{L} is defined partially at time- n by the n -tower. For $g \in A_n A_n^{-1}$, \mathcal{L}^g translates E_a^n onto E_b^n if $g = ba^{-1}$. Thus \mathcal{L}^g is defined on the levels of the n -tower indexed by $(g^{-1} A_n) \cap A_n$. Since A_n is a Følner sequence, as $n \rightarrow \infty$, \mathcal{L}^g becomes defined a.e. To build the $(n+1)$ -tower from the n -tower, each level E_a^n is divided into $|C_n|$ equal intervals, each closed on the left and open on the right. We assume C_n is ordered $\{c_1, c_2, c_3, \dots\}$ and label our new intervals from left to right as

$$E_{ac_1}^{n+1}, E_{ac_2}^{n+1}, E_{ac_3}^{n+1} \dots$$

The $(n+1)$ -tower is formed by these intervals, together with some additional intervals, E_b^{n+1} for all b in $A_{n+1} \setminus A_n C_n$. The additional intervals were called spacers in the classical rank-one construction. Because of property three of the almost-tiling Følner sequence, we only introduce a finite amount of measure in the whole construction.

It's not hard to see that the partially defined action on the $(n+1)$ -tower is consistent with the partially defined action on the n -tower. By normalizing we can assume the action is defined on $[0, 1)$ with probability Lebesgue measure μ .

2.4. Some notations for rank-one actions. The entire space on which the action takes place is called X . The subset of X that is the n -tower is called X_n . The n -tower is composed of left-closed, right-open intervals labeled E_a^n , where a is in A_n . Let $E_s^n := X \setminus X_n$. The s stands for spacer.

The partition P_n divides X into levels E_i^n of the n -tower and the complement E_s^n . View P_n as a function from X to $A_n^* := A_n \cup \{s\}$, defined by $P_n(x) = j$ if $x \in E_j^n$ for $j \in A_n^*$. Define $P_k^n : A_n^* \rightarrow A_k^*$ by $P_k^n(u) = v$ if $E_u^n \subset E_v^k$. This corresponds to the natural map, based on inclusion, from the n -tower to the k -tower ($n > k$). The value $P_k^n(u)$ is the P_k name of a level u in the n -tower. Thus, $P_k^n \circ P_n = P_k$.

Definition 4 (P -name of x). Let be P be a partition on X , so that $P : X \rightarrow A$, and $x \in X$. Then the function $G \rightarrow A$ defined by $g \rightarrow P(\mathcal{L}^g x)$ is the P -name of x .

Definition 5 (*n*-block). Given a P_k -name of x (for $k < n$) and a fixed $c \in G$, if for all $u \in A_n$ we find that $P_k(\mathcal{L}^{uc}x) = P_k^n(u)$ then $\{P_k(\mathcal{L}^{uc}x) \mid u \in A_n\}$ is an *n*-block indexed by $A_n c$.

The set C_n determines how *n*-blocks are situated in the $(n + 1)$ -tower.

We often compare names of different points x and y . If the function from G to $A \times A$ is given by $g \rightarrow (P(\mathcal{L}^g x), P(\mathcal{L}^g y))$, then this function is the $P \times P$ name of (x, y) . We refer to the symbols in the first coordinate as *upper* and the symbols in the second coordinate as *lower*.

Definition 6 (Overlap of *n*-blocks). If an *n*-block occurs at $A_n c$ in the P_k name of x and an *n*-block occurs at $A_n c'$ in the P_k name of y then $A_n c \cap A_n c'$ indexes an *n*-block overlap in the $P_k \times P_k$ name of (x, y) . The overlap is given by

$$\{(P_k(\mathcal{L}^h x), P_k(\mathcal{L}^h y)) \mid h \in A_n c \cap A_n c'\}.$$

2.5. Comparing measures on finite sets. If f is a map from a probability space (A, μ, α) to a finite set B , then $\text{Dist}_{a \in A} f(a)$ denotes the probability measure on B given by $\text{Dist}_{a \in A} f(a) := \alpha \circ f^{-1}$.

For $C \subset A$, the *conditional measure* is defined by $\alpha_C(U) := \alpha(U \cap C) / \alpha(C)$. If C is a measurable subset of A define $\text{Dist}_{a \in C} f(a) := \alpha_C \circ f^{-1}$. We also denote this by $\text{Dist}(f(a) \mid a \in C)$. For convenience and readability we also use

$$\text{Dist}_{a \in C}(f(a) \mid g(a) = k)$$

in place of

$$\text{Dist}_{\{a \in C \mid g(a) = k\}} f(a).$$

For finite A the normalized counting measure on A is denoted $\text{Unif } A$.

Example. Define a measure λ on A_n by $\text{Dist}_{x \in X_n}(P_n(x))$. So if E_p^n is a level of the *n*-tower then

$$\lambda(p) = \mu_{X_n}(P_n^{-1}(p)) = \frac{\mu(X_n \cap E_p^n)}{\mu(X_n)} = \frac{1}{m}.$$

Thus $\text{Dist}_{x \in X_n}(P_n(x)) = \text{Unif } A_n$.

Example. Let $X_i : \Omega \rightarrow A$ for $i = 1, 2, 3, \dots$ be a countable family of independent random variables each with uniform distribution over a finite alphabet A . Then for fixed $w \in \Omega$, $\mu := \text{Dist}(X_i(w) \mid i \in [1, m])$ is a measure on A . Precisely,

$$\mu(a) = \frac{|\{i \in [1, m] \mid X_i(w) = a\}|}{m}.$$

To compare two probability measures, p, q , on a finite set A , use the norm

$$\|p - q\| := \sum_{a \in A} |p(a) - q(a)|.$$

The following five lemmas about finite measures, which we use repeatedly, are stated without proof.

Lemma 3. *If p and q are probability measures on A and $\rho : A \rightarrow B$ then*

$$\|p \circ \rho^{-1} - q \circ \rho^{-1}\| \leq \|p - q\|.$$

Lemma 4. *If $C \subset A$ and $p(C) \geq 1 - \epsilon$, then $\|p - p_C\| \leq 2\epsilon$.*

Let $\Pi : A \times B \rightarrow A$ be the projection. If p is a probability measure on $A \times B$, then the *marginal measure* $\bar{p} := \Pi(p)$ is the projection of p onto A . So $\bar{p}(a) = \sum_b p(a, b)$. For $a \in A$ define the *fibre measure* p_a as p conditioned on $\Pi^{-1}(a)$. So

$$p_a(b) = \frac{p(a, b)}{\bar{p}(a)}.$$

Lemma 5. *If p and q are probability measures on $A \times B$ and $\bar{p} = \bar{q}$, then*

$$\|p - q\| = \sum_a \bar{p}(a) \|p_a - q_a\|.$$

A more useful version of Lemma 5 when \bar{p} is not exactly equal to \bar{q} is:

Lemma 6. *If p and q are probability measures on $A \times B$ then*

$$\|p - q\| \leq \sum_a \bar{p}(a) \|p_a - q_a\| + \|\bar{p} - \bar{q}\|.$$

Note that Lemma 5 is a special case of this lemma. When the marginals on A are close and the fibre measures over A are close, Lemma 6 implies the measures will be close.

Conversely, the distance between the measures gives a bound on the distance between the fibre measures.

Lemma 7. *If p and q are probability measures on $A \times B$ then for a fixed $a \in A$,*

$$\|p_a - q_a\| \leq \frac{2}{\bar{p}(a)} \|p - q\|.$$

2.6. Special application of the ergodic theorem. The following is a standard consequence of the mean and pointwise Ergodic Theorems:

Lemma 8. *Suppose (X, β, T, μ) is ergodic, P is a measurable finite partition of X , $\epsilon > 0$, and $\alpha > 0$. Then for μ -almost all x in X , there exists N (depending on x) so that if $E = l + \cup_{m=0}^M [mn, mn + L]$ where:*

1. $n \geq N$,
2. $M \in \mathbb{N}$,
3. $\alpha n \leq L \leq n$, and
4. $-n \leq l \leq n$,

then

$$\left| \text{Dist}_{i \in E} P(T^i x) - \text{Dist } P \right| < \epsilon.$$

Lemma 8 shows the Ergodic Theorem is valid on sets other than initial segments of the integers. The conclusion of this lemma holds even if E is

$$l + \bigcup_{m=0}^M [mn + \delta_m, mn + \delta_m + L],$$

where δ_m satisfies $0 \leq \delta_m \leq (1 - \alpha)n$. The case $\delta_m = 0$, however, when the intervals of length L are regularly spaced, is sufficient for our needs.

3. Main proof

Theorem 9 (Main result). *Let G be a countable abelian group with \mathbb{Z}^d ($d \geq 1$) as a subgroup so that G/\mathbb{Z}^d is a locally finite group. There exists a rank-one action of G so that the transformation T corresponding to $(1, 0, 0, \dots, 0)$ in \mathbb{Z}^d is mixing, simple, and only commutes with the other times of the action, that is, $C(T) = G$.*

We prove this theorem in three parts: the construction of the group action in Subsection 3.1, the proof that T is mixing in Subsection 3.2, and the proof that T is simple and $C(T) = G$ in Subsection 3.3.

3.1. Constructing a rank-one action of G . Let $H := G/\mathbb{Z}^d$. H is a locally finite group. The proof when H is finite or even trivial is an easy modification of the proof given, the case when H is infinite. As H is countable there exists a sequence of finite groups H_n , so that

$$H_1 \subset H_2 \subset H_3 \dots$$

and

$$H = \bigcup_{n=1}^{\infty} H_n.$$

Our group G has a special structure that we will identify. Consider the cosets of \mathbb{Z}^d in G and for each coset select a unique element in the coset. Use $\psi(g)$ to denote the unique element in the coset $g + \mathbb{Z}^d$. So ψ maps from G into G but is not a homomorphism since $\psi(g_1 + g_2)$ need not be equal to $\psi(g_1) + \psi(g_2)$.

Since G/\mathbb{Z}^d is isomorphic to H , the projection $\Pi_H : G \rightarrow H$ is a homomorphism that has \mathbb{Z}^d as its kernel. Define $\phi : H \rightarrow G$ so that $\phi \circ \Pi_H = \psi$. Then we can define an operation, also denoted $+$, on $\mathbb{Z}^d \times H$ as follows:

$$(\mathbf{u}, h_1) + (\mathbf{v}, h_2) := (\mathbf{u} + \mathbf{v} + \phi(h_1) + \phi(h_2) - \phi(h_1 + h_2), h_1 + h_2).$$

We think of this as addition in each coordinate with a carry from the H coordinate into the \mathbb{Z}^d coordinate. The carry, $\phi(h_1) + \phi(h_2) - \phi(h_1 + h_2)$, is an element of \mathbb{Z}^d . Why? If g_1 and g_2 are two elements of G so that $\Pi_H(g_1) = h_1$ and $\Pi_H(g_2) = h_2$ then

$$\begin{aligned} \phi(h_1) + \phi(h_2) - \phi(h_1 + h_2) &= \phi(\Pi_H(g_1)) + \phi(\Pi_H(g_2)) - \phi(\Pi_H(g_1) + \Pi_H(g_2)) \\ &= \psi(g_1) + \psi(g_2) - \phi(\Pi_H(g_1 + g_2)) \\ &= \psi(g_1) + \psi(g_2) - \psi(g_1 + g_2). \end{aligned}$$

Since $\psi(g_1) \in g_1 + \mathbb{Z}^d$ and $\psi(g_2) \in g_2 + \mathbb{Z}^d$, then $\psi(g_1) + \psi(g_2) \in (g_1 + g_2) + \mathbb{Z}^d$. Yet, $\psi(g_1 + g_2)$ is also in $(g_1 + g_2) + \mathbb{Z}^d$, so $\psi(g_1) + \psi(g_2) - \psi(g_1 + g_2)$ is an element of \mathbb{Z}^d .

Claim 10. *With this operation, $\mathbb{Z}^d \times H$ is isomorphic to G .*

Proof. Define an isomorphism $\Phi : G \rightarrow \mathbb{Z}^d \times H$ as $\Phi(g) := (g - \psi(g), \Pi_H(g))$. Then

$$\begin{aligned} \Phi(g_1) + \Phi(g_2) &= (g_1 - \psi(g_1), \Pi_H(g_1)) + (g_2 - \psi(g_2), \Pi_H(g_2)) \\ &= (g_1 - \psi(g_1) + g_2 - \psi(g_2) + \phi(\Pi_H(g_1)) + \phi(\Pi_H(g_2)) \\ &\quad - \phi(\Pi_H(g_1) + \Pi_H(g_2)), \Pi_H(g_1) + \Pi_H(g_2)) \\ &= (g_1 - \psi(g_1) + g_2 - \psi(g_2) + \psi(g_1) + \psi(g_2) \\ &\quad - \phi(\Pi_H(g_1 + g_2)), \Pi_H(g_1 + g_2)) \\ &= (g_1 + g_2 - \psi(g_1 + g_2), \Pi_H(g_1 + g_2)) \\ &= \Phi(g_1 + g_2). \end{aligned}$$

Thus Φ is a homomorphism. Each $g \in G$ can be uniquely written as $(g - \psi(g)) + \psi(g)$. Since $g - \psi(g) \in \mathbb{Z}^d$ and $\psi(g)$ corresponds to $\Pi_H(g) \in H$, it is clear that Φ is injective and surjective. \square

For the remainder of our proof we consider our group G to be presented as $\mathbb{Z}^d \times H$ with the special operation $+$, where we add in each coordinate and carry from the H coordinate into the \mathbb{Z}^d coordinate. An element in G is uniquely represented by (\mathbf{v}, g) , where \mathbf{v} is in \mathbb{Z}^d and g is in H ; sometimes denoted $\mathbf{v} + g$.

For notational convenience:

1. \mathbf{u} denotes (u_1, u_2, \dots, u_d) a vector in \mathbb{Z}^d .
2. $\|\mathbf{u}\|$ is $\max(|u_1|, |u_2|, \dots, |u_d|)$.
3. $a\mathbf{v}$ is the usual scalar multiplication.
4. \mathbf{e}_i represents the i th standard basis vector in \mathbb{Z}^d .
5. For an integer m , \mathbf{m} denotes (m, m, m, \dots, m) .

Thus, $m\mathbf{e}_1$ denotes the \mathbb{Z}^d vector $(m, 0, 0, \dots, 0)$. For $\mathbf{u} = (u_1, u_2, \dots, u_d)$ and $\mathbf{v} = (v_1, v_2, \dots, v_d)$, if $u_i \leq v_i$ for all i , define $[\mathbf{u}, \mathbf{v}]$ to be the set

$$[u_1, v_1] \times [u_2, v_2] \times \dots \times [u_d, v_d].$$

Define the projection from G onto \mathbb{Z}^d by $\Pi_{\mathbb{Z}^d}(\mathbf{v}, g) := \mathbf{v}$. Although Π_H is a homomorphism, $\Pi_{\mathbb{Z}^d}$ is not.

H_{n+1} is composed of $k_n := |H_{n+1}/H_n|$ cosets of H_n . Choose elements of H_{n+1} , one from each coset, to form $\Gamma_n = \{g_1, g_2, \dots, g_{k_n}\}$. Then

$$H_{n+1} = (g_1 + H_n) \cup (g_2 + H_n) \cup \dots \cup (g_{k_n} + H_n).$$

To construct a rank-one action of G we specify an almost-tiling Følner sequence $(\{A_n\}, \{C_n\})$. Our Følner sequence will be defined by $A_n := [0, h_n]^d \times H_n$. More precisely define:

1. Spacer length $s_n := nh_{n-1}$.
2. Window length $w_n := h_n + s_n + 2\theta_{n+1}$.
3. New block length $h_{n+1} := N_n w_n$.
4. The maximum norm of the carry:

$$\theta_n := \sup_{g, g' \in H_n} \|\Pi_{\mathbb{Z}^d}((0, g) + (0, g'))\|.$$

This is the \mathbb{Z}^d norm, as the carry is in \mathbb{Z}^d .

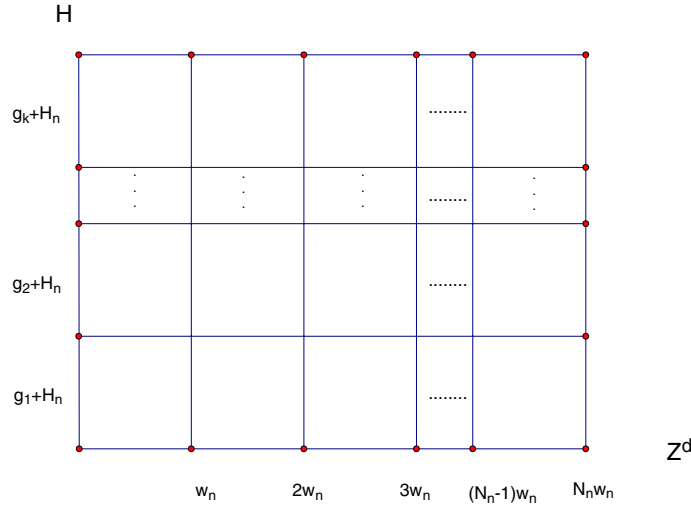


FIGURE 1. Windows in A_{n+1} .

The value of N_n will be specified later. The $(n + 1)$ -tower A_{n+1} has *windows*

$$[w_n \mathbf{i}, w_n(\mathbf{i} + \mathbf{1})] \times (g + H_n),$$

where $\mathbf{i} \in [1, N_n]^d$ and $g \in \Gamma_n$. A window is identified by its *order* (\mathbf{i}, g) . Figure 1 illustrates the case $d=1$, but is sufficient to envision the general case.

To almost-tile A_{n+1} we place a copy of A_n centrally in each window, with a “random perturbation”. More precisely, a spacer function η_n is used to place the copies of A_n in windows of A_{n+1} , that is,

$$\eta_n : [1, N_n]^d \times \Gamma_n \longrightarrow [0, s_n]^d \times H_n,$$

and is defined to be nearly random as detailed later. Now the set of translators C_n is given by

$$C_n := \{(w_n \mathbf{i}, g) + (\theta_{n+1} \mathbf{1}, 0) + \eta_n(\mathbf{i}, g) \mid \mathbf{i} \in [1, N_n]^d, g \in \Gamma_n\}.$$

For $c \in C_n$, we describe exactly how cA_n appears in its window. The n -tower, as seen in Figure 2, is composed of many *rows* of the form $[0, h_n)^d \times r$, for $r \in H_n$. The translation of A_n by c has two components. The \mathbb{Z}^d component called the *shift*,

$$\Pi_{\mathbb{Z}^d}(c) = w_n \mathbf{i} + \theta_{n+1} \mathbf{1} + \Pi_{\mathbb{Z}^d}(\eta_n(\mathbf{i}, g)),$$

and the H component called the *shuffle*,

$$\Pi_H(c) = g + \Pi_H(\eta_n(\mathbf{i}, g)).$$

What is the effect of the shift on A_n ? It shifts the entire n -block to the center of the window $(\mathbf{i}, 0)$ and a further small translation due to the spacer sequence. If we now apply the shuffle what effect does it have? The shuffle “rotates” the rows of A_n as it places them in the window (\mathbf{i}, g) . It also introduces a small \mathbb{Z}^d translation for each row, less than θ_{n+1} in any of the d directions. So the definition of w_n ensures the shuffled n -block is inside the window. Figure 3 illustrates the case $d = 1$.

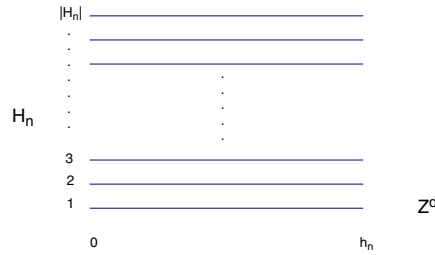


FIGURE 2. A_n is composed of rows.

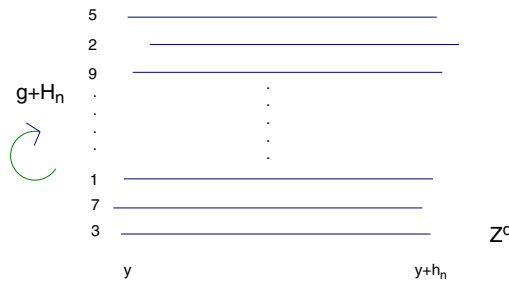


FIGURE 3. *Shuffling* an n -block into a window of A_{n+1} . Here $y = \Pi_{Z^d}(c)$.

When $d = 2$, cA_n could look like a loosely shuffled deck of cards. For any c and c' in C_n , the n -blocks cA_n and $c'A_n$ are inside their respective windows and thus disjoint.

Let m be a positive integer and let (\mathbf{i}, g) and (\mathbf{i}', g') be two starting windows in A_{n+1} .

Definition 7 (Admissibility). The triple $(m, (\mathbf{i}, g), (\mathbf{i}', g'))$ is called *admissible* if:

1. $\mathbf{i}, \mathbf{i}', \mathbf{i} + m\mathbf{e}_1$, and $\mathbf{i}' + m\mathbf{e}_1$ are in $[1, N_n]^d$.
2. (\mathbf{i}, g) is not equal to (\mathbf{i}', g') .
3. $m \geq \frac{N_n}{n^2}$.

Spacer functions η_n must satisfy a uniformity condition: for all admissible triples $(m, (\mathbf{i}, g), (\mathbf{i}', g'))$, starting at windows (\mathbf{i}, g) and (\mathbf{i}', g') in A_{n+1} , and looking in the m consecutive windows in the positive \mathbb{Z} direction, there is a jointly uniform distribution of spacers. Let $\epsilon_n = ((s_n)^d |H_n|)^{-2}$ which clearly $\xrightarrow{n} 0$. Then we require

$$(1) \quad \left\| \text{Dist}_{k \in [0, m]} (\eta_n(\mathbf{i} + k\mathbf{e}_1, g), \eta_n(\mathbf{i}' + k\mathbf{e}_1, g')) - \text{Unif}([0, s_n]^d \times H_n)^2 \right\| \leq \epsilon_n.$$

Why do we want $m \geq N_n/n^2$? The ratio of m to the number of windows in a block required to see uniformity must decrease to zero as n goes to infinity. Yet, m

must increase to infinity to make uniformity possible. Choosing $m > N_n/n^2$ is one way to accomplish this.

The condition (1) on η_n may appear very restrictive. How can we be sure such functions exist? The length of the spacer sequence, N_n , is not yet fixed. As the length of a sequence of random letters from a finite alphabet grows, it is exponentially more likely that its letters are uniformly distributed. From this fact we can construct η_n .

Our spacer sequence η_n will be a realization of a sequence of N_n independent random variables each uniformly distributed over the alphabet $A = ([0, s_n]^d \times H_n)^2$. We show that for large enough N_n there is a high probability that (1) is satisfied.

Lemma 11 (Exponential convergence of random sequences to uniformity). *Given a probability space (Ω, P) , let X_1, X_2, X_3, \dots be independent random variables each with uniform distribution over a finite alphabet A . Specifically, $X_i : \Omega \rightarrow A$ and $P\{X_i = a\} = \frac{1}{|A|}$ for all a in A . Then for all $\epsilon > 0$, there exists M and a constant c so that if $m \geq M$ then*

$$(2) \quad P\{\|\text{Dist}(X_i \mid i \in [0, m]) - \text{Unif } A\| \geq \epsilon\} < e^{-cm}.$$

Proof. This is a consequence of the Central Limit Theorem. See [Orn67] for an elementary proof or [Rud79] for a proof using Stirling’s Formula. \square

Each admissible triple $(m, (\mathbf{i}, g), (\mathbf{i}', g'))$ imposes a constraint on the spacer sequence. Apply Lemma 11 with $A = ([0, s_n]^d \times H_n)^2$ and $\epsilon = \epsilon_n$ to determine c and M . If g is not equal to g' , then clearly the m windows to the right of these starting windows are distinct. So by Lemma 11, if $m > M$ then the probability that (1) is not satisfied is less than e^{-cm} .

If g equals g' , then it may be that the m windows to the right of the starting points do overlap. If so, divide $[0, m]$ into two sets A and B , where $|A|$ and $|B|$ are both greater than $m/3$, and so that the sets $\{\mathbf{i} + k\mathbf{e}_1 \mid k \in A\}$ and $\{\mathbf{i}' + k\mathbf{e}_1 \mid k \in A\}$ are disjoint. Similarly for B , $\{\mathbf{i} + k\mathbf{e}_1 \mid k \in B\}$ and $\{\mathbf{i}' + k\mathbf{e}_1 \mid k \in B\}$ are disjoint. Then the probability that

$$(3) \quad \left\| \text{Dist}_{k \in A} (\eta_n(\mathbf{i} + k\mathbf{e}_1, g), \eta_n(\mathbf{i}' + k\mathbf{e}_1, g')) - \text{Unif}([0, s_n]^d \times H_n)^2 \right\| \leq \epsilon_n.$$

is less than $e^{-cm/3}$, and likewise for B . So the probability that (1) is not true for an admissible triple $(m, (\mathbf{i}, g), (\mathbf{i}', g'))$ is less than the probability that (3) is not true for $k \in A$ or $k \in B$, which is less than $2e^{-cm/3}$.

Since $m > \frac{N_n}{n^2}$ we can conclude that the probability of (1) not being true is less than $2e^{-\frac{cN_n}{3n^2}}$, provided $\frac{N_n}{3n^2} > M$. Remembering $k = |H_{n+1}/H_n|$, there are less than $N_n((N_n)^d k)^2$ admissible triples so the probability that (1) is true for all of them is less than

$$(4) \quad 1 - 2(N_n)^{2d+1} k^2 e^{-cN_n/3n^2}.$$

So, take N_n sufficiently large that $\frac{N_n}{3n^2} > M$ and (4) is greater than zero. This means a realization of η_n will exist to satisfy (1). Additionally we require that N_n be large enough to ensure that $\theta_{n+3} < h_{n+1}$. This means that θ_{n+1} will be much less than s_n , a fact that will be important in later scanning arguments.

The sequences $\{A_n\}$ and $\{C_n\}$ thus defined constitute an almost-tiling Følner sequence. By also requiring that N_n grow exponentially with n , we guarantee that the resulting measure space is finite. Thus, we have \mathcal{L} a rank-one action of G .

3.2. T is mixing. To show T is mixing we must show $\mu(T^i A \cap B) \rightarrow \mu(A)\mu(B)$ for all measurable sets A and B . Since A and B can be approximated by levels of the k -tower, for large k , it suffices to show that for each k , P_k is mixed by T . More precisely, we want to show

$$(5) \quad \left\| \text{Dist}_x (P_k(T^i x), P_k(x)) - \text{Dist}_{(x,y)} (P_k(x), P_k(y)) \right\|$$

tends to 0 as $i \rightarrow \infty$.

Each level of the k tower is a union of levels of the $(n-1)$ -tower (when $n-1 > k$). We can link i to n by requiring

$$(6) \quad w_n - (s_n + w_{n-1} + \theta_{n+1}) \leq i \leq w_{n+1} - (s_{n+1} + w_n + \theta_{n+2}).$$

So (5) is bounded by

$$(7) \quad \left\| \text{Dist}_x (P_{n-1}(T^i x), P_{n-1}(x)) - \text{Dist}_{(x,y)} (P_{n-1}(x), P_{n-1}(y)) \right\|.$$

Why is this true? The measures in (5) are probability measures on the set $A_k^* \times A_k^*$. Because of Lemma 3 with the map $P_k^{n-1} \times P_k^{n-1}$, (5) is less than or equal to (7). Thus to prove T is mixing it suffices to show (7) tends to zero as n and thus i go to infinity. Why link i to n ? When i is small compared to h_{n+1} , we can investigate the first measure in (7) by examining an $(n+1)$ -block and its overlap with the $(n+1)$ -block shifted by T^i . When i is much larger than w_n , but still smaller than w_{n+1} , we examine the overlap of $(n+2)$ -blocks.

The key idea is that when we shift the $(n+1)$ -tower by an offset i in the \mathbb{Z} coordinate, we see nearly uniform $P_{n-1} \times P_{n-1}$ symbols on the overlap.

Definition 8 (Row overlaps). The $(n+1)$ -block is composed of $|H_{n+1}|$ rows, each of the form $[0, h_n]^d \times g$ for fixed $g \in H_{n+1}$. An overlap of rows is indexed by $[\mathbf{0}, \mathbf{m}]$ and two points \mathbf{a} and \mathbf{a}' . This is a map

$$(8) \quad [\mathbf{0}, \mathbf{m}] \rightarrow ([\mathbf{0}, \mathbf{h}_n], g) \times ([\mathbf{0}, \mathbf{h}_n], g')$$

which is defined by $\mathbf{u} \rightarrow ((\mathbf{a} + \mathbf{u}, g), (\mathbf{a}' + \mathbf{u}, g'))$. Thus the overlap is a pairing of $[\mathbf{a}, \mathbf{a} + \mathbf{m}] \times g$ in the upper block with $[\mathbf{a}', \mathbf{a}' + \mathbf{m}] \times g'$ in the lower block.

Lemma 12 (Row overlap measures). Consider an overlap of two rows in A_{n+1} , as above, where $|m_i| > \frac{N_g}{n^2} w_n$ for all i . If $g + H_n$ equals $g' + H_n$ and $\|\mathbf{a} - \mathbf{a}'\| > w_n - (s_n + w_{n-1} + \theta_{n+1})$, or if $g + H_n$ is not equal to $g' + H_n$ then

$$(9) \quad \left\| \text{Dist}_{\mathbf{u} \in [\mathbf{0}, \mathbf{m}]} (P_{n-1}^{n+1}(\mathbf{a} + \mathbf{u}, g), P_{n-1}^{n+1}(\mathbf{a}' + \mathbf{u}, g')) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\|$$

tends to 0 as $n \rightarrow \infty$.

An overlap of rows that satisfies the conclusion of Lemma 12 is called *good*. Using this lemma we can complete the proof that T is mixing.

Proof that T is mixing. When i is in the range

$$w_n - (s_n + w_{n-1} + \theta_{n+1}) \leq i \leq w_n + s_{n+1} + \theta_{n+2},$$

i is small compared to h_{n+1} , so the overlap between A_{n+1} and $(i\mathbf{e}_1 + A_{n+1})$ is large. Let B equal to $A_{n+1} \cap (i\mathbf{e}_1 + A_{n+1})$ and let X_B represent the levels of the $(n+1)$ -tower indexed by B . Statement (7) is less than

$$\begin{aligned} & \left\| \text{Dist}_x(P_{n-1}(T^i x), P_{n-1}(x)) - \text{Dist}_{x \in X_B}(P_{n-1}(T^i x), P_{n-1}(x)) \right\| \\ & + \left\| \text{Dist}_{x \in X_B}(P_{n-1}(T^i x), P_{n-1}(x)) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\| \\ & + \left\| \text{Unif}(A_{n-1} \times A_{n-1}) - \text{Dist}_{(x,y)}(P_{n-1}(x), P_{n-1}(y)) \right\|. \end{aligned}$$

Lemma 4 shows the first summand is less than $2(1 - \mu(X_B))$ which $\xrightarrow{n} 0$. Why? Notice that X_B is a large fraction of X_{n+1} . More precisely

$$\frac{\mu(X_B)}{\mu(X_{n+1})} \geq \frac{h_{n+1} - w_n - s_{n+1} - \theta_{n+2}}{h_{n+1}}.$$

This fraction $\xrightarrow{n} 1$ so we see that $1 - \mu(X_B) \xrightarrow{n} 0$. Lemma 4 also shows the third summand is less than $2(1 - \mu(X_{n-1}))$ which $\xrightarrow{n} 0$. The second summand is equivalent to

$$(10) \quad \left\| \text{Dist}_{\mathbf{u} \in B}(P_{n-1}^{n+1}(\mathbf{u} + i\mathbf{e}_1), P_{n-1}^{n+1}(\mathbf{u})) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\|,$$

and this distribution over the overlap indexed by B is a weighted average of distributions over overlaps of rows in A_{n+1} . They are all good overlaps so Lemma 12 shows that (10) $\xrightarrow{n} 0$.

When i is in the range

$$(11) \quad w_n + s_{n+1} + \theta_{n+2} \leq i \leq w_{n+1} - (s_{n+1} + \theta_{n+2} + w_n),$$

we consider the overlap of $(n+2)$ -towers. Let B be $A_{n+2} \cap (i\mathbf{e}_1 + A_{n+2})$, and let X_B represent levels of the $(n+2)$ -tower indexed by B . Statement (7) is less than

$$\begin{aligned} & \left\| \text{Dist}_x(P_{n-2}(T^i x), P_{n-2}(x)) - \text{Dist}_{x \in X_B}(P_{n-1}(T^i x), P_{n-1}(x)) \right\| \\ & + \left\| \text{Dist}_{x \in X_B}(P_{n-1}(T^i x), P_{n-1}(x)) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\| \\ & + \left\| \text{Unif}(A_{n-1} \times A_{n-1}) - \text{Dist}_{(x,y)}(P_{n-1}(x), P_{n-1}(y)) \right\|. \end{aligned}$$

Again, Lemma 4 shows the first and third summands $\xrightarrow{n} 0$. The second summand is equal to

$$(12) \quad \left\| \text{Dist}_{\mathbf{u} \in B}(P_{n-1}^{n+1}(\mathbf{u} + i\mathbf{e}_1), P_{n-1}^{n+1}(\mathbf{u})) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\|.$$

Most of the overlap indexed by B is composed of overlaps of $(n+1)$ -blocks. These $(n+1)$ -blocks are not in their standard positions. They have been shifted and shuffled into their respective windows of the $(n+2)$ -tower. This means their overlaps are not indexed by a simple set of the form $(i' + A_{n+1}) \cap A_{n+1}$, rather, they are the union of overlaps of rows of the $(n+1)$ -block. Because of the range for i , most of the overlap indexed by B is composed of good overlaps of rows of $(n+1)$ -blocks. They are good because the range for i guarantees that the row overlaps start in different windows or are too short to be significant. Thus, the portion of B not in good overlaps $\xrightarrow{n} 0$. The distribution of the P_{n-1} symbols in the overlap indexed by B will be a weighted average of distributions over good row overlaps. So (12) is a weighted average of distributions like (9) which $\xrightarrow{n} 0$ by Lemma 12. This concludes the proof that T is mixing. \square

We turn to the proof of Lemma 12. Rather than distributions on overlaps of rows like the map from $[\mathbf{0}, \mathbf{m}]$ to $([\mathbf{0}, \mathbf{h}_n], g) \times ([\mathbf{0}, \mathbf{h}_n], g')$ defined by

$$(13) \quad \mathbf{u} \rightarrow ((\mathbf{a} + \mathbf{u}, g), (\mathbf{a}' + \mathbf{u}, g')),$$

we consider sets formed by fixing a position in the window and looking at that position, in both the upper and lower block, in some number of consecutive windows in the positive \mathbb{Z} direction.

Definition 9 (Span). A span is a map from $[0, M]$ to $([\mathbf{0}, \mathbf{h}_n], g) \times ([\mathbf{0}, \mathbf{h}_n], g')$ of the form

$$(14) \quad i \rightarrow ((\mathbf{b} + iw_n \mathbf{e}_1, g), (\mathbf{b}' + iw_n \mathbf{e}_1, g'))$$

where:

1. \mathbf{b} and \mathbf{b}' are in $[\mathbf{0}, \mathbf{h}_{n+1}]$.
2. g and g' are in H_{n+1} .
3. $(\mathbf{b} + Mw_n \mathbf{e}_1)$ and $(\mathbf{b}' + Mw_n \mathbf{e}_1)$ are in $[\mathbf{0}, \mathbf{h}_{n+1}]$.
4. $M > N_n/n^2$.

Let $R(b, m)$ be $b \bmod m$. Call a vector \mathbf{b} *good* if each coordinate b_i satisfies

$$s_n + \theta_{n+1} \leq R(b_i, w_n) \leq w_n - (s_n + \theta_{n+1}).$$

If both \mathbf{b} and \mathbf{b}' are good and in windows of different order in A_{n+1} , we call (14) a *good span*. The constraint on b_i means that if \mathbf{b} is a good vector and $g \in H_{n+1}$ then (\mathbf{b}, g) is contained in $C_n A_n$. That's because \mathbf{b} is far enough away from the boundary of the \mathbb{Z}^d component of the window that (\mathbf{b}, g) is always included in the n -block in the window. So the $P_n \times P_n$ symbols in good spans are never spacers.

Overlaps like (13) are unions of spans of nearly equal size. So (9) is an almost equally weighted average of distributions over spans D ,

$$(15) \quad \left\| \text{Dist}_{d \in D} (P_{n-1}^{n+1} \times P_{n-1}^{n+1})(d) - \text{Unif}(A_{n-1} \times A_{n-1}) \right\|.$$

Most of the spans considered are good but those that are not have weight at most $2d(s_n + \theta_{n+1} + w_{n-1})/w_n$. Since their weight $\xrightarrow{n} 0$, to establish Lemma 12 it suffices to prove the following:

Lemma 13. *If D is a good span then (15) $\xrightarrow{n} 0$.*

Proof. Let D be as in (14), and let the starting points (\mathbf{b}, g) and (\mathbf{b}', g') be in windows (\mathbf{I}, g_I) and $(\mathbf{I}', g_{I'})$ respectively. Here g_I and $g_{I'}$ are in Γ_n so that $g + H_n = g_I + H_n$ and $g' + H_n = g_{I'} + H_n$. Let ν be the measure

$$\text{Dist}_{d \in D} (P_{n-1}^{n+1} \times P_{n-1}^{n+1})(d).$$

The marginal measure $\bar{\nu}$ observes the P_{n-1} symbols at the fixed position in the windows of the upper $(n+1)$ -block. Because the position is fixed, the value of the spacer function in that window determines the symbol seen. We look at this position for M consecutive windows along the positive \mathbb{Z} axis. Over that range the spacers are nearly uniformly distributed,

$$\left\| \text{Dist}_{k \in [0, M]} \eta_n(\mathbf{I} + k\mathbf{e}_1, g_I) - \text{Unif}([0, s_n]^d \times H_n) \right\| \leq \epsilon_n,$$

so we have the effect of “scanning” a region of the n -block nearly uniformly. Why? The measure $\bar{\nu}$ is equal to

$$\text{Dist}_{k \in [0, M]} P_{n-1}^n ((\mathbf{b} + kw_n \mathbf{e}_1, g) - kw_n \mathbf{I} - g_I - \eta_n(\mathbf{I} + i\mathbf{e}_1, g_I)).$$

Since the H component of the spacers are uniform over H_n we scan a portion of each row of A_n . Because the spacers are jointly uniform over $[0, s_n]^d$ for each row of A_n we scan a portion that is $[0, s_n]^d$ in size. The region scanned is like the shuffled block seen in Figure 3, but we scan only a portion of each row that is $[0, s_n]^d$ in size. Also because we have arranged for the maximum carry θ_{n+1} to be small compared to s_n , the scanned region of A_n contains $n^d|H_n/H_{n-1}|$ complete copies of the $(n - 1)$ -block, and the portion scanned that is not in those blocks $\xrightarrow{n} 0$. Thus we observe nearly uniform P_{n-1} symbols at our fixed window position, that is, $\|\bar{\nu} - \text{Unif } A_{n-1}\| \xrightarrow{n} 0$.

Now examine the fibres of ν . The fibre measure ν_a looks only at the windows in the upper $(n + 1)$ -block where we see a fixed symbol a ,

$$\nu_a := \text{Dist}_{i \in [0, M]} (P_{n-1}^{n+1}(\mathbf{b}' + iw_n \mathbf{e}_1, g') \mid P_{n-1}^{n+1}(\mathbf{b} + iw_n \mathbf{e}_1, g) = a).$$

The fixed symbol a only occurs for a finite number of spacer values $\gamma_1, \gamma_2, \dots, \gamma_t$. Because the spacers are jointly uniform,

$$\|\text{Dist}_{k \in [0, M]} (\eta_n(\mathbf{I} + kw_n \mathbf{e}_1, g_I), \eta_n(\mathbf{I}' + kw_n \mathbf{e}_1, g')) - \text{Unif}([0, s_n]^d \times H_n)^2\| \leq \epsilon_n,$$

for a fixed spacer value in the window of the upper block, Lemma 7 shows the distribution of spacers in the windows of the lower block are also nearly uniform,

$$\begin{aligned} & \|\text{Dist}_{k \in [0, M]} (\eta_n(\mathbf{I}' + kw_n \mathbf{e}_1, g') \mid \eta_n(\mathbf{I} + kw_n \mathbf{e}_1, g) = \gamma_j) - \text{Unif}([0, s_n]^d \times H_n)\| \\ & \leq \epsilon_n s_n^d |H_n|. \end{aligned}$$

Since $\epsilon_n = (s_n^d |H_n|)^2$ then $\epsilon_n s_n^d |H_n| = \sqrt{\epsilon_n}$ which also $\xrightarrow{n} 0$. Thus, we get nearly uniform scanning of a region of the n -block by observing the windows of the lower block at our fixed position. By the same nearly uniform scanning argument, we see nearly uniform P_{n-1} symbols. This occurs for each γ_j ; so taking a weighted average, $\|\nu_a - \text{Unif } A_{n-1}\| \xrightarrow{n} 0$.

Since both $\bar{\nu}$ and ν_a are nearly $\text{Unif } A_{n-1}$, by Lemma 6

$$\|\nu - \text{Unif}(A_{n-1} \times A_{n-1})\| \xrightarrow{n} 0. \quad \square$$

Thus Lemma 13 is true which concludes the proof of Lemma 12.

3.3. T is simple and $C(T) = G$. Let λ be an ergodic self-joining of T . We show that λ is either the product measure $\mu \times \mu$, or an off diagonal $(I \times \mathcal{L}^g)\Delta$. Since $(I \times S)\Delta$ is an ergodic joining for all $S \in C(T)$, this implies $C(T)$ is merely times(\mathcal{L}) (all the transformations in the action of G).

Take $\epsilon > 0$, $\alpha = 1/4$ and a partition Q of $X \times X$; apply Lemma 8. Then for λ -almost all (x, x') there exists N (depending on (x, x')) so that if E is a set of the form $E = l + \cup_{i=0}^M [iw_n, iw_n + L]$ where:

1. $n \geq N$,
2. $M \in \mathbb{N}$,
3. $\alpha w_n \leq L \leq w_n$, and
4. $-w_n \leq l \leq w_n$,

then

$$\|\text{Dist}_{i \in E} Q(T^i x, T^i x') - \text{Dist}_{(y, y')} Q(y, y')\| < \epsilon.$$

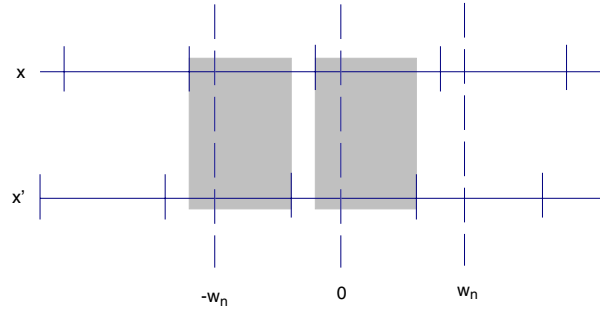


FIGURE 4. Overlap of \mathbb{Z} - n -windows in the names of x and x' .

A point (x, x') that satisfies the above for all rational $\epsilon > 0$, all partitions $P_k \times P_k$ and $\alpha = 1/4$, we call *generic*. Since this is only a countable number of applications of the lemma, λ -almost all points are generic. Fix one such point $(x, x') \in X \times X$.

Definition 10 (Time- n order of x). If $x \in X_{n+1}$ then it is contained in a level of the tower E_u^n for $u \in A_{n+1}$. Let $\Theta_n(x)$ denote the time- n order of x , defined by

$$\Theta_n(x) := (\mathbf{i}, g),$$

where (\mathbf{i}, g) is the window of A_{n+1} that contains u .

Definition 11 (\mathbb{Z} - n -window). For a point $x \in X_{n+1}$ the \mathbb{Z} - n -window is $[a, b]$, a maximal interval in \mathbb{Z} , so that $\Theta_n(T^i x)$ is constant for $i \in [a, b]$.

Definition 12 (Centrally Located). Say x is centrally located at time- n if

$$N_n/n^2 < i_1 < N_n - N_n/n^2,$$

where $\Theta_n(x) = (\mathbf{i}, g)$. Define eventually centrally located to mean there exists an N depending on x so that for all times $n \geq N$, x is centrally located at time n .

By the Borel-Cantelli lemma, since the amount of measure we are outlawing is finite, μ -almost all x are eventually centrally located. We want x and x' to be away from the edge of the n -tower in the \mathbb{Z} coordinate.

Fix a pair (x, x') that is generic and both x and x' are eventually centrally located. Examining the structure of the $P_k \times P_k$ names of (x, x') will reveal the nature of λ . The placement of \mathbb{Z} - n -windows in the name of x' relative to those in the name of x , will be of particular importance. When x and x' are both centrally located, \mathbb{Z} - n -windows for x and x' fill the interval $[-N_n w_n/n^2, N_n w_n/n^2]$.

Definition 13 (Overlap of \mathbb{Z} - n -windows). An overlap of \mathbb{Z} - n -windows is a maximal interval $[a, b]$ such that $\Theta_n(T^i x)$ and $\Theta_n(T^i x')$ are constant for i in $[a, b]$.

Definition 14 (Offset time). Call n an offset time if there is an overlap of length at least $w_n/4$ in $[-w_n, w_n]$ so that the overlapping \mathbb{Z} - n -windows (in x and x') have different orders.

Remarkably, the occurrence of offset times completely determines the structure of λ .

Claim 14 (Finitely many offset times). *If there are only a finite number of offset times then $\lambda = (I \times \mathcal{L}^g)\Delta$, for some $g \in G$.*

Proof. Let N be sufficiently large so that for all $n \geq N$, n is not an offset time, and x and x' are centrally located. We know $[-w_N, w_N]$ is filled with \mathbb{Z} - n -windows for both x and x' , so it must contain at least one overlap of size $> w_n/4$. Since N is not an offset time this overlap is from windows whose order must be the same. If x and x' do not have the same order at time- N their orders differ by no more than $(1, 0, \dots, 0) \in \mathbb{Z}^d$. The \mathbb{Z} -windows of x and x' at time- $(N+1)$ must overlap almost completely, certainly $> w_{N+1}/4$ in length. Since $N+1$ is not an offset time $\Theta_{N+1}(x) = \Theta_{N+1}(x')$. We repeat, inductively, to see that for all $n \geq N+1$, $\Theta_n(x) = \Theta_n(x')$, that is, x and x' have the same order. If x is not equal to $\mathcal{L}^g x'$ for some $g \in A_{N+1} - A_{N+1}$ then at some time $n > N$, x and x' would be in different windows. But that's not possible. This implies $x = \mathcal{L}^g x'$ for some $g \in G$. Thus $\lambda = (I \times \mathcal{L}^g)\Delta$. \square

Claim 15 (Infinitely many offset times). *If there is an infinite sequence of offset times $\{n_j\}_{j=1}^\infty$, then $\lambda = \mu \times \mu$.*

Proof. We can assume that for n_j where $j \geq 1$, x and x' are centrally located. Consider the name of (x, x') with respect to partitions $P_k \times P_k$. In the interval $[-w_{n_j}, w_{n_j}]$ there will be an overlap of $\mathbb{Z} - n_j$ -windows of length greater than $w_{n_j}/4$. Since n_j is an offset time, this is an overlap of windows of different orders and because x and x' are centrally located this overlap pattern is repeated at least N_{n_j}/n_j^2 times to the right. Let L_j be the width of the overlap. Thus we can define a set E_{n_j} of the form $l + \cup_{i=0}^M [i w_{n_j}, i w_{n_j} + L_j]$ where:

1. $M = N_{n_j}/n_j^2$.
2. $w_{n_j}/4 \leq L_j \leq w_{n_j}$.
3. $-w_{n_j} \leq l \leq w_{n_j}$.

Because (x, x') was chosen to be generic for λ over sets like E_{n_j} , for all partitions $P_k \times P_k$,

$$(16) \quad \lim_{j \rightarrow \infty} \text{Dist}_{i \in E_{n_j}}(P_k(T^i x), P_k(T^i x')) = \text{Dist}_\lambda P_k \times P_k.$$

If we also knew that for all k ,

$$(17) \quad \lim_{j \rightarrow \infty} \text{Dist}_{i \in E_{n_j}}(P_k(T^i x), P_k(T^i x')) = \text{Dist}_{\mu \times \mu} P_k \times P_k,$$

then, since the partitions $P_k \times P_k$ generate the σ -algebra, from (16) and (17) we could conclude that $\lambda = \mu \times \mu$.

To establish (17) we will show

$$\left\| \text{Dist}_{i \in E_{n_j}}(P_k(T^i x), P_k(T^i x')) - \text{Dist}_{\mu \times \mu}(P_k \times P_k) \right\| \xrightarrow{n_j} 0.$$

Proceeding similarly to the proof that T is mixing, it suffices to show

$$\left\| \text{Dist}_{i \in E_{n_j}}(P_{n_j-1}(T^i x), P_{n_j-1}(T^i x')) - \text{Dist}_{\mu \times \mu}(P_{n_j-1} \times P_{n_j-1}) \right\| \xrightarrow{n_j} 0.$$

This is because P_{n_j-1} refines P_k when n_j is large enough. Of course

$$\left\| \text{Dist}_{\mu \times \mu}(P_{n_j-1} \times P_{n_j-1}) - \text{Unif}(A_{n_j-1} \times A_{n_j-1}) \right\|$$

is less than $2\mu(X_{n_j-1}^C)$ which $\xrightarrow{n_j} 0$. So it suffices to show that

$$(18) \quad \left\| \text{Dist}_{i \in E_{n_j}}(P_{n_j-1}(T^i x), P_{n_j-1}(T^i x')) - \text{Unif}(A_{n_j-1} \times A_{n_j-1}) \right\| \xrightarrow{n_j} 0.$$

But the measure

$$\text{Dist}_{i \in E_{n_j}}(P_{n_j-1}(T^i x), P_{n_j-1}(T^i x'))$$

is an average of distributions over a fixed position in the windows, a span. A small proportion of the spans comprising E_{n_j} are not good, at most $4(s_{n_j} + \theta_{n_j+1})/w_{n_j}$, which goes to 0. Since n_j is an offset time the rest are good. By Claim 13, statement (18) tends to 0. \square

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