A NOTE ON THE KRONECKER-WEYL EQUIDISTRIBUTION THEOREM

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ABSTRACT. We study the relationship between the discrete and the continuous versions of the Kronecker–Weyl equidistribution theorem, as well as their possible extension to manifolds in higher dimensions. We also investigate a way to deduce in some limited way uniformity results in higher dimension from results in lower dimension.

1. Introduction

There are two versions of the classical Kronecker–Weyl equidistribution theorem. Let d be a fixed positive integer. If $v_1, \ldots, v_d, 1$ are linearly independent over \mathbb{Q} , then the vector $\mathbf{v} = (v_1, \ldots, v_d) \in \mathbb{R}^d$ is called a Kronecker vector, and the vector $\mathbf{v}^* = (v_1, \ldots, v_d, 1) \in \mathbb{R}^{d+1}$ is called a Kronecker direction. The continuous version concerns the distribution of half-infinite geodesics with Kronecker directions in the unit torus $[0, 1)^{d+1}$, and much of this monograph concerns extensions of this version from the unit torus $[0, 1)^2$ to arbitrary finite polysquare translation surfaces.

On the other hand, the discrete version of the Kronecker-Weyl equidistribution theorem concerns the distribution of half-infinite Kronecker sequences in the unit torus $[0,1)^d$. A natural first question is then to attempt to extend this version from the unit torus $[0,1)^2$ to arbitrary finite polysquare translation surfaces.

In general, for any fixed positive integer d, we consider the continuous problem concerning the distribution of half-infinite geodesics with Kronecker directions in finite polycube translation manifolds in d+1 dimensions, as well as the discrete problem of the distribution of half-infinite Kronecker sequences in finite polycube translation manifolds in d dimensions. The following are natural questions:

Question 1. Is it true that, for any integer $d \ge 1$, any half-infinite geodesic with a Kronecker direction in a finite polycube translation manifold in d+1 dimensions is uniformly distributed? If not, then under what condition can we guarantee uniform distribution?

Question 2. Is it true that, for any integer $d \ge 2$, any half-infinite Kronecker sequence in a finite polycube translation manifold in d dimensions is uniformly distributed? If not, then under what condition can we quarantee uniform distribution?

Question 3. The classical Kronecker-Weyl equidistribution theorem on the unit torus has some time-quantitative extensions with explicit error terms. Under what conditions can we establish time-quantitative uniformity in these more general settings?

For Question 1, the Gutkin-Veech theorem [4, 5] answers the special case d = 1 in the affirmative. However, for d = 2, we are currently not able to establish uniformity results for half-infinite geodesics in a general finite polycube translation 3-manifold.

In this paper, we study the relationship between the discrete and the continuous versions of possible non-integrable analogues of the Kronecker–Weyl equidistribution

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theorem concerning finite polysquare translation surfaces and some related finite polycube translation 3-manifolds. In Section 2, we investigate a relationship between the discrete version and the continuous version in the special case d=2. Theorem 2, a by-product of this study, gives an affirmative answer to the special case d=2 of Question 2. Then in Section 3, we develop a way to step up the problem by one dimension, and this leads to various infinite classes of polycube translation manifolds where half-infinite Kronecker sequences and half-infinite geodesics with Kronecker directions are uniformly distributed.

2. A SIMPLE EQUIVALENCE PRINCIPLE

Before we go any further, it is appropriate that we understand what we mean by a Kronecker sequence on a finite polysquare translation surface \mathcal{P} . To properly define a half-infinite Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}, \ j = 0, 1, 2, 3, \ldots$, with starting point \mathbf{v}_0 and step vector \mathbf{v} on a polysquare translation surface \mathcal{P} , we need a supporting half-infinite geodesic $\mathcal{L}(t), \ t \geq 0$, with $\mathcal{L}(0) = \mathbf{v}_0$ on \mathcal{P} and a time step $g \in \mathbb{R}$, such that the finite geodesic segment $\mathcal{L}(t), \ 0 \leq t \leq g$, is in the direction of the step vector \mathbf{v} and has length equal to $|\mathbf{v}|$. Then $\mathbf{v}_0 + j\mathbf{v} = \mathcal{L}(jg), \ j = 0, 1, 2, 3, \ldots$. Clearly the Kronecker sequence is half-infinite if and only if \mathbf{v}_0 is a non-pathological starting point of a geodesic with direction \mathbf{v} on \mathcal{P} .

Suppose that \mathcal{P} is a polysquare translation surface with s atomic squares, and that $\mathcal{M} = \mathcal{P} \times [0,1)$ denotes the associated polycube translation 3-manifold which is the cartesian product of \mathcal{P} and the unit torus [0,1).

A vector $\mathbf{v} = (v_1, v_2) \in \mathbb{R}^2$ is said to be a Kronecker vector if $v_1, v_2, 1$ are linearly independent over \mathbb{Q} . We are interested in the distribution of a Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}, \ j = 0, 1, 2, 3, \ldots$, with starting point \mathbf{v}_0 on the polysquare translation surface \mathcal{P} .

A vector $\mathbf{v}^* = (v_1, v_2, 1) \in \mathbb{R}^3$ is said to be a Kronecker direction if $v_1, v_2, 1$ are linearly independent over \mathbb{Q} . We are interested in the distribution of a half-infinite geodesic $\mathcal{L}(t)$, $t \geq 0$, with starting point $\mathcal{L}(0)$ and direction \mathbf{v}^* in the associated polycube translation 3-manifold \mathcal{M} .

We have a simple equivalence principle relating half-infinite Kronecker sequences and half-infinite geodesics with Kronecker direction.

Theorem 1. Suppose that \mathcal{P} is a finite polysquare translation surface, and that $\mathcal{M} = \mathcal{P} \times [0,1)$ denotes the associated polycube translation 3-manifold which is the cartesian product of \mathcal{P} and the unit torus [0,1). Suppose also that $\mathbf{v} = (v_1, v_2) \in \mathbb{R}^2$ is a Kronecker vector, so that $\mathbf{v}^* = (v_1, v_2, 1) \in \mathbb{R}^3$ is a Kronecker direction. Then the following two statements are equivalent:

- (i) Every half-infinite Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, 2, 3, \ldots$, on \mathcal{P} is uniformly distributed.
- (ii) Every half-infinite geodesic $\mathcal{L}(t)$, $t \geqslant 0$, with direction \mathbf{v}^* in \mathcal{M} is uniformly distributed.

Sketch of proof. ((ii) \Rightarrow (i)) The argument here is rather simple. Suppose that \mathcal{P} has s atomic squares. To establish (i), let S be a convex set in an atomic square of \mathcal{P} . Consider the first J terms of the Kronecker sequence, given by $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, \ldots, J - 1$. The number of terms of this finite sequence in S is given by

$$|\{j=0,1,\ldots,J-1:\mathbf{v}_0+j\mathbf{v}\in S\}|,$$

with corresponding expectation given by

$$\frac{J\lambda_2(S)}{s}$$

where λ_2 denotes 2-dimensional Lebesgue measure. To establish uniformity of the Kronecker sequence on \mathcal{P} , we need to show that

$$|\{j=0,1,\ldots,J-1:\mathbf{v}_0+j\mathbf{v}\in S\}|/\frac{J\lambda_2(S)}{s}\to 1 \text{ as } J\to\infty.$$
 (2.1)

Let $\mathcal{L}(t)$, $t \geq 0$, be a half-infinite geodesic with starting point $\mathcal{L}(0) = (\mathbf{v}_0, 0)$ and direction \mathbf{v}^* on \mathcal{M} . Let $S^* \subset \mathcal{M}$ be obtained by sweeping the set $S \times \{0\}$ along by the vector \mathbf{v}^* , so that $\lambda_3(S^*) = \lambda_2(S)$, where λ_3 denotes 3-dimensional Lebesgue measure. Consider a finite geodesic segment $\mathcal{L}(t)$, $0 \leq t \leq T = J(v_1^2 + v_2^2 + 1)^{1/2}$. Then the total length of the parts of this geodesic segment in S^* is given by

$$|\{0 \leqslant t \leqslant T : \mathcal{L}(t) \in S^*\}|$$

$$= (v_1^2 + v_2^2 + 1)^{1/2} |\{j = 0, 1, \dots, J - 1 : \mathbf{v}_0 + j\mathbf{v} \in S\}|, \tag{2.2}$$

with corresponding expectation given by

$$\frac{T\lambda_3(S^*)}{s} = \frac{J(v_1^2 + v_2^2 + 1)^{1/2}\lambda_2(S)}{s}.$$
 (2.3)

The set S^* is a union of finitely many convex sets in atomic cubes of \mathcal{M} . Suppose that (ii) holds. Uniformity of the half-infinite geodesic in \mathcal{M} then implies that

$$|\{0 \leqslant t \leqslant T : \mathcal{L}(t) \in S^*\}| / \frac{T\lambda_3(S^*)}{s} \to 1 \quad \text{as } T \to \infty.$$
 (2.4)

It is clear that (2.1) follows immediately on combining (2.2)–(2.4).

 $((i) \Rightarrow (ii))$ In [3, Section 3.4.1] or [1, Section 3.6], a result concerning uniformity of a half-infinite geodesic is deduced from a result concerning uniformity of a sequence, by an application of the Koksma inequality. Here we apply an analogous Koksma inequality type argument.

The Gutkin–Veech theorem gives uniformity to any half-infinite geodesic with Kronecker direction on a finite polysquare translation surface. As a simple application of Theorem 1, we establish the analogous result for half-infinite Kronecker sequences.

Theorem 2. Any half-infinite Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, 2, 3, \ldots$, on a finite polysquare translation surface \mathcal{P} is uniformly distributed.

Proof. Let $\mathcal{M} = \mathcal{P} \times [0,1)$ denote the associated polycube translation 3-manifold. In view of Theorem 1, it suffices to show that every half-infinite geodesic $\mathcal{L}(t)$, $t \geq 0$, with Kronecker direction in \mathcal{M} is uniformly distributed. This latter condition is the conclusion of Theorem 3.

The following result is [2, Theorem 3].

Theorem 3. Suppose that a polycube translation 3-manifold \mathcal{M} is the cartesian product of a finite polysquare translation surface \mathcal{P} and the unit torus [0,1). Then any half-infinite geodesic in \mathcal{M} with a Kronecker direction $\mathbf{v}^* \in \mathbb{R}^3$ is uniformly distributed unless it hits a singularity.

3. Stepping up principle

Theorem 2, concerning the distribution of half-infinite Kronecker sequences on \mathcal{P} , can be interpreted as a step-up from the Gutkin–Veech theorem, concerning the distribution of half-infinite geodesics with Kronecker direction on \mathcal{P} . In view of the equivalence given by Theorem 1, the stepping-up result is Theorem 3, which establishes the uniformity of a half-infinite geodesic with Kronecker direction on the associated polycube translation 3-manifold \mathcal{M} from the uniformity of a half-infinite

geodesic with Kronceker direction on the finite polysquare translation surface \mathcal{P} . This represents a step-up in dimension in some restricted way.

In this section, we expand on this idea. We establish the following results which, for simplicity, we state only in 2 and 3 dimensions. There are analogues in d and d+1 dimensions for every integer $d \ge 3$.

The discrete versions of these results concern the distribution of half-infinite Kronecker sequences on a finite polysquare translation surface \mathcal{P} and the analogous problem in the associated polycube translation 3-manifold \mathcal{M} which is the cartesian product of \mathcal{P} and the unit torus [0,1).

- **Theorem 4.** Suppose that \mathcal{P} is a finite polysquare translation surface, and that $\mathcal{M} = \mathcal{P} \times [0,1)$ denotes the associated polycube translation 3-manifold which is the cartesian product of \mathcal{P} and the unit torus [0,1). Suppose also that $\mathbf{v} \in \mathbb{R}^2$ is the step vector of a Kronecker sequence on \mathcal{P} . Then the following statements are equivalent:
- (i) Every half-infinite Kronecker sequence with step vector \mathbf{v} on \mathcal{P} is uniformly distributed.
- (ii) For any $w \in \mathbb{R}$ such that $\mathbf{w} = (\mathbf{v}, w) \in \mathbb{R}^3$ is the step vector of a Kronecker sequence on \mathcal{M} , every half-infinite Kronecker sequence with step vector \mathbf{w} in \mathcal{M} is uniformly distributed.

Theorem 5. Under the hypotheses of Theorem 4, the following statements are equivalent:

- (i) The \mathbf{v} -shift on \mathcal{P} is ergodic.
- (ii) For any $w \in \mathbb{R}$ such that $\mathbf{w} = (\mathbf{v}, w) \in \mathbb{R}^3$ is the step vector of a Kronecker sequence on \mathcal{M} , the \mathbf{w} -shift in \mathcal{M} is ergodic.

The continuous versions of these results concern the distribution of half-infinite geodesics with Kronecker direction on a finite polysquare translation surface \mathcal{P} and the analogous problem in the associated polycube translation 3-manifold \mathcal{M} .

- **Theorem 6.** Suppose that \mathcal{P} is a finite polysquare translation surface, and that $\mathcal{M} = \mathcal{P} \times [0,1)$ denotes the associated polycube translation 3-manifold which is the cartesian product of \mathcal{P} and the unit torus [0,1). Suppose also that $\mathbf{v} \in \mathbb{R}^2$ is a Kronecker direction. Then the following statements are equivalent:
 - (i) Every half-infinite geodesic with direction \mathbf{v} on \mathcal{P} is uniformly distributed.
- (ii) For any $w \in \mathbb{R}$ such that $\mathbf{w} = (\mathbf{v}, w) \in \mathbb{R}^3$ is a Kronecker direction, every half-infinite geodesic with direction \mathbf{w} in \mathcal{M} is uniformly distributed.

Theorem 7. Under the hypotheses of Theorem 4, the following statements are equivalent:

- (i) Geodesic flow with direction \mathbf{v} on \mathcal{P} is ergodic.
- (ii) For any $w \in \mathbb{R}$ such that $\mathbf{w} = (\mathbf{v}, w) \in \mathbb{R}^3$ is a Kronecker direction, geodesic flow with direction \mathbf{w} in \mathcal{M} is ergodic.

Remark. For Theorems 6 and 7, we already know that (i) and (ii) hold, in view of the Gutkin–Veech theorem and Theorem 3, so only the analogues in d and d+1 dimensions for integers $d \ge 3$ are new.

We concentrate our efforts on establishing Theorem 4. That (ii) implies (i) is almost trivial, by projection from \mathcal{M} to \mathcal{P} . The converse is considerably harder.

Let $\mathbf{w} = (v_1, v_2, w) \in \mathbb{R}^3$, where $v_1, v_2, w, 1$ are linearly independent over \mathbb{Q} , and consider the **w**-shift $\mathbf{T}_{\mathbf{w}}^* = \mathbf{T}_{\mathbf{w}}^*(\mathcal{M})$ of geodesic flow in direction \mathbf{w} in $\mathcal{M} = \mathcal{P} \times [0, 1)$. We can consider two projections of $\mathbf{T}_{\mathbf{w}}^*$.

On the one hand, we can project the **w**-shift $\mathbf{T}_{\mathbf{w}}^*$ to the unit torus $[0,1)^3$, simply by taking every coordinate modulo 1, leading to the **w**-shift $\mathbf{T}_{\mathbf{w}} = \mathbf{T}_{\mathbf{w}}([0,1)^3)$ in the unit torus $[0,1)^3$ which is ergodic.

On the other hand, we can project the **w**-shift $\mathbf{T}_{\mathbf{w}}^*$ to the polysquare translation surface \mathcal{P} , simply by dropping reference to the 3-rd coordinates throughout, leading to the **v**-shift $\mathbf{T}_{\mathbf{v}}^* = \mathbf{T}_{\mathbf{v}}^*(\mathcal{P})$ on \mathcal{P} .

Lemma 3.1. Suppose that the condition (i) in Theorem 4 holds. Then the **v**-shift $\mathbf{T}_{\mathbf{v}}^* = \mathbf{T}_{\mathbf{v}}^*(\mathcal{P})$ on \mathcal{P} is ergodic.

Sketch of proof. If a Kronecker sequence on \mathcal{P} becomes undefined after finitely many terms, then the supporting geodesic $\mathcal{L}(t)$, $t \geq 0$, with starting point $\mathcal{L}(0) = \mathbf{v}_0$ and direction \mathbf{v} on \mathcal{P} hits a singular point of \mathcal{P} . The collection of singular points of \mathcal{P} clearly has 2-dimensional Lebesgue measure 0. Thus the collection of starting points \mathbf{v}_0 that lead a Kronecker sequence to be undefined after finitely many terms also has 2-dimensional Lebesgue measure 0. Thus almost every point $\mathbf{v}_0 \in \mathcal{P}$ gives rise to a half-infinite Kronecker sequence. Thus the condition (i) in Theorem 4 guarantees that for almost every starting point \mathbf{v}_0 , the half-infinite Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, 2, 3, \ldots$, on \mathcal{P} is uniformly distributed. Hence the visiting density of the sequence in any Jordan measurable set A is given by $\lambda_2(A)/s$.

Suppose, on the contrary, that the **v**-shift $\mathbf{T}_{\mathbf{v}}^*$ on \mathcal{P} is not ergodic. Then there is a partition $\mathcal{P} = U_1 \cup U_2$, where the subsets $U_1, U_2 \subset \mathcal{P}$ are $\mathbf{T}_{\mathbf{v}}^*$ -invariant, and satisfy $\lambda_2(U_1) > 0$ and $\lambda_2(U_2) > 0$. Moreover, since the projection of $\mathbf{T}_{\mathbf{v}}^*$ to the unit torus $[0,1)^2$, simply by taking every coordinate modulo 1, leads to the **v**-shift $\mathbf{T}_{\mathbf{v}} = \mathbf{T}_{\mathbf{v}}([0,1)^2)$ on the unit torus $[0,1)^2$ which is ergodic, there is a decomposition $\mathcal{P} = M_1 \cup \ldots \cup M_k$ for some integer k satisfying $2 \leq k \leq s$, where s is the number of atomic squares in \mathcal{P} , such that for every $i = 1, \ldots, k$, the set M_i is $\mathbf{T}_{\mathbf{v}}^*$ -invariant and does not contain a proper $\mathbf{T}_{\mathbf{v}}^*$ -invariant subset, and $\lambda_2(M_i)$ is a positive integer. We say that the subsets M_1, \ldots, M_k are minimal.

We can find a Jordan measurable subset $A \subset \mathcal{P}$ such that

$$\lambda_2(A \cap M_1) < \frac{1}{10}$$
 and $\lambda_2(A \cap M_2) > \frac{\lambda_2(M_2)}{2}$.

Next we apply the Birkhoff ergodic theorem to both M_1 and M_2 . The restriction of the **v**-shift $\mathbf{T}_{\mathbf{v}}^*$ to M_1 and to M_2 are both ergodic. In each case, the simplest form of the ergodic theorem says that the *time average* is equal to the *space average*. Then for almost every starting point $\mathbf{v}_0 \in M_1$, the visiting density of the Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, 2, 3, \ldots$, on \mathcal{P} in the set A is equal to the relative measure

$$\frac{\lambda_2(A\cap M_1)}{\lambda_2(M_1)}<\frac{1}{10},$$

while for almost every starting point $\mathbf{v}_0 \in M_2$, the visiting density of the Kronecker sequence $\mathbf{v}_0 + j\mathbf{v}$, $j = 0, 1, 2, 3, \ldots$, on \mathcal{P} in the set A is equal to the relative measure

$$\frac{\lambda_2(A \cap M_2)}{\lambda_2(M_2)} > \frac{1}{2}.$$

Thus at least one of these is different from $\lambda_2(A)/s$, leading to a contradiction. \square

Proof of Theorem 4. ((i) \Rightarrow (ii)) We need to prove that the **w**-shift $\mathbf{T}_{\mathbf{w}}^*$ in \mathcal{M} is ergodic. Suppose on the contrary that this is not the case. Then there exists a non-trivial partition $\mathcal{M} = \mathcal{W} \cup \mathcal{S}$ into a union of two $\mathbf{T}_{\mathbf{w}}^*$ -invariant subsets $\mathcal{W}, \mathcal{S} \subset \mathcal{M}$, called White and Silver, say, each with integral measure and such that

$$1 \leqslant \lambda_3(\mathcal{W}), \lambda_3(\mathcal{S}) \leqslant s - 1,$$

where s is the number of atomic cubes in \mathcal{M} and λ_3 denotes 3-dimensional Lebesgue measure.

Let $\mathcal{Y}_1, \ldots, \mathcal{Y}_s$ denote the atomic cubes of \mathcal{M} , and consider the projection of \mathcal{M} to the unit torus $[0,1)^3$. Then for any point $P \in [0,1)^3$, there are precisely s points

$$P_1 \in \mathcal{Y}_1, \ldots, P_s \in \mathcal{Y}_s$$

that have the same image P under this projection. Let

$$f_{\mathcal{W}}(P) = |\{P_1, \dots, P_s\} \cap \mathcal{W}| \quad \text{and} \quad f_{\mathcal{S}}(P) = |\{P_1, \dots, P_s\} \cap \mathcal{S}|.$$

Then $f_{\mathcal{W}}$ and $f_{\mathcal{S}}$ are positive integer valued functions defined on $[0,1)^3$ such that

$$f_{\mathcal{W}}(P) + f_{\mathcal{S}}(P) = s$$
 for almost every $P \in [0, 1)^3$.

The w-shift $\mathbf{T}_{\mathbf{w}}$ in the unit torus $[0,1)^3$ is ergodic. Since the functions $f_{\mathcal{W}}$ and $f_{\mathcal{S}}$ are $\mathbf{T}_{\mathbf{w}}$ -invariant, it follows from the Birkhoff ergodic theorem that they are constant almost everywhere in $[0,1)^3$. Then

$$f_{\mathcal{W}} + f_{\mathcal{S}} = s \quad \text{and} \quad 1 \leqslant f_{\mathcal{W}}, f_{\mathcal{S}} \leqslant s - 1.$$
 (3.1)

Let $\chi_{\mathcal{W}}$ denote the characteristic function of \mathcal{W} in \mathcal{M} . For every $\mathbf{s} \in \mathcal{M}$, write $\mathbf{s} = (\mathbf{x}, z)$, where $\mathbf{x} \in \mathcal{P}$ and $z \in [0, 1)$. The well known Fubini theorem implies that

$$\int_{\mathcal{M}} \chi_{\mathcal{W}}(\mathbf{s}) \, d\mathbf{s} = \int_{\mathcal{P}} \left(\int_{[0,1)} \chi_{\mathcal{W}}(\mathbf{x}; z) \, dz \right) d\lambda_2,$$

where the inner integral

$$\psi(\mathbf{x}) = \int_{[0,1)} \chi_{\mathcal{W}}(\mathbf{x}; z) \, \mathrm{d}z$$

is well defined in the sense of Lebesgue for almost every $\mathbf{x} \in \mathcal{P}$.

Consider now the projection of $\mathbf{T}_{\mathbf{v}}^*$ to the polysquare translation surface \mathcal{P} , resulting in the ergodic \mathbf{v} -shift $\mathbf{T}_{\mathbf{v}}^*$ on \mathcal{P} . Since the function $\psi(\mathbf{x})$ is $\mathbf{T}_{\mathbf{v}}^*$ -invariant, it follows from ergodicity that it has constant value almost everywhere on \mathcal{P} . Thus it follows from (3.1) that

$$\frac{1}{s} \leqslant \psi = \frac{f_{\mathcal{W}}}{s} \leqslant \frac{s-1}{s} \tag{3.2}$$

almost everywhere on \mathcal{P} .

Since the invariant subset $W \subset \mathcal{M}$ is measurable, it follows that for every $\varepsilon_1 > 0$, there exists a finite set of 3-dimensional axis-parallel rectangular boxes such that their union $\mathcal{B} = \mathcal{B}(W; \varepsilon_1)$ has the property that the symmetric difference

$$\mathcal{W} \triangle \mathcal{B} = (\mathcal{W} \setminus \mathcal{B}) \cup (\mathcal{B} \setminus \mathcal{W})$$

has measure

$$\lambda_3(\mathcal{W} \triangle \mathcal{B}) < \varepsilon_1. \tag{3.3}$$

We need the following simple technical result.

Let $\chi_{\mathcal{B}}$ denote the characteristic function of \mathcal{B} in \mathcal{M} .

Lemma 3.2. Let \mathcal{B} be a finite union of axis-parallel rectangular boxes satisfying (3.3). For every $\varepsilon_2 > 0$, there exists $\varepsilon_3 > 0$ and a finite set of disjoint axis-parallel rectangular boxes such that their union $\mathcal{B}^* = \mathcal{B}^*(\mathcal{B}; \varepsilon_2; \varepsilon_3)$ satisfies the following conditions:

- (i) The measure $\lambda_3(\mathcal{M} \setminus \mathcal{B}^*) < \varepsilon_2$.
- (ii) We have $\chi_{\mathcal{B}}(\mathbf{s}') = \chi_{\mathcal{B}}(\mathbf{s}'')$ if the points $\mathbf{s}', \mathbf{s}'' \in \mathcal{M}$ belong to the same axisparallel rectangular box in the disjoint union \mathcal{B}^* .
- (iii) The side lengths of each axis-parallel rectangular box in the disjoint union \mathcal{B}^* is greater than ε_3 .

(iv) Let s denote the number of atomic squares of \mathcal{P} . Then

$$\lambda_2\left(\left\{\mathbf{x} \in \mathcal{P} : \lambda_1(\{z \in [0,1) : (\mathbf{x},z) \in \mathcal{B}^*\}\right) > 1 - \varepsilon_2^{1/2}\right\}\right) > s - \varepsilon_2^{1/2}.$$

Proof. We assume, without loss of generality, that every axis-parallel rectangular box in the finite union \mathcal{B} lies in the interior of an atomic cube of \mathcal{M} , as illustrated in the picture on the left in Figure 1, which only shows 2 of the 3 dimensions. For each face of each axis-parallel rectangular box in an atomic cube of \mathcal{M} , we extend it to an axis-parallel special square face of area 1 in the same atomic cube, as illustrated by the dashed lines in the picture on the right in Figure 1. We repeat this process for every axis-parallel rectangular box in the finite union \mathcal{B} . Suppose that $\mathcal{N} = \mathcal{N}(\mathcal{B})$ denotes the total number of distinct axis-parallel rectangular boxes in the union \mathcal{B} . Then there are at most $2\mathcal{N}$ such distinct special square faces in each of the 3 axisparallel directions. It follows that there exists a number $\varepsilon_4 = \varepsilon_4(\mathcal{B}) > 0$ such that any special square face in an atomic cube has distance at least ε_4 from any other parallel special square face in the atomic cube or from any parallel boundary square face of the atomic cube. We now remove every point in \mathcal{M} that lies a distance less than $\varepsilon_5 > 0$ from some special square face in the atomic cube that contains that point. Then the set of such points in \mathcal{M} that are removed has measure at most $12 \mathcal{N} \varepsilon_5$, and is represented by the regions shaded in light gray in the picture on the right in Figure 1.

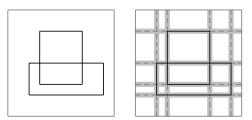


Figure 1: idea behind the construction of the set \mathcal{B}^*

Let \mathcal{B}^* denote the remainder of \mathcal{M} after these points are removed. Then \mathcal{B}^* is clearly a finite union of disjoint axis-parallel rectangular boxes in \mathcal{M} , where each box is either contained in \mathcal{B} or disjoint from \mathcal{B} , so that the condition (ii) is satisfied. Suppose now that $\varepsilon_5 > 0$ is chosen to satisfy

$$12\mathcal{N}\varepsilon_5 < \varepsilon_2 \quad \text{and} \quad 2\varepsilon_5 < \varepsilon_4.$$
 (3.4)

Then the first condition in (3.4) ensures that the condition (i) is satisfied, while the second condition in (3.4) ensures that the side lengths of each axis-parallel rectangular box in the union \mathcal{B}^* is greater than $\varepsilon_3 = \varepsilon_4 - 2\varepsilon_5$, so that the condition (iii) is satisfied.

To establish the condition (iv), note that it follows from the condition (i) and the Fubini theorem that

$$\int_{\mathcal{P}} \lambda_1(\{z \in [0,1) : (\mathbf{x},z) \notin \mathcal{B}^*\}) \, \mathrm{d}\mathbf{x} < \varepsilon_2.$$
(3.5)

Let

$$\mathcal{E} = \left\{ \mathbf{x} \in \mathcal{P} : \lambda_1(\{z \in [0,1) : (\mathbf{x},z) \notin \mathcal{B}^*\}) \geqslant \varepsilon_2^{1/2} \right\}.$$

Then it follows from (3.5) that $\lambda_2(\mathcal{E})\varepsilon_2^{1/2} < \varepsilon_2$, from which the condition (iv) follows immediately, since $\lambda_2(\mathcal{P}) = s$.

For almost every $\mathbf{x} \in \mathcal{P}$, the set

$$U(\mathcal{W}; \mathbf{x}) = \{ z \in [0, 1) : (\mathbf{x}, z) \in \mathcal{W} \}$$

is measurable, and it follows from (3.2) that

$$\frac{1}{s} \leqslant \lambda_1(U(\mathcal{W}; \mathbf{x})) = \frac{f_{\mathcal{W}}}{s} \leqslant \frac{s-1}{s}.$$
(3.6)

We consider Lebesgue measurable subsets $U_{\sigma} \subset [0,1)$, $\sigma = 1, \ldots, s$, of the unit torus [0,1). In particular, we make the assumption that $0 < \lambda_1(U_1) < 1$, where λ_1 denotes 1-dimensional Lebesgue measure. Furthermore, for any real number $u \in \mathbb{R}$ and any $\sigma = 1, \ldots, s$, we consider the (-u)-translated copy of U_{σ} , given by

$$U_{\sigma} - u = \{ \{ z - u \} : z \in U_{\sigma} \}.$$

Let $\mathbf{x}_1, \ldots, \mathbf{x}_s \in \mathcal{P}$ be distinct points such that their images under projection modulo 1 to the unit torus $[0,1)^2$ coincide. We now apply the following variant of [2, Lemma 6.1] to the sets

$$U_1 = U(\mathcal{W}; \mathbf{x}_1), \quad \dots, \quad U_s = U(\mathcal{W}; \mathbf{x}_s),$$
 (3.7)

so that for every $\sigma = 1, \ldots, s$,

$$z \in U_{\sigma}$$
 if and only if $(\mathbf{x}_{\sigma}, z) \in \mathcal{W}$.

Lemma 3.3. The set of values $u_0 \in [0,1)$ for which the inequalities

$$\lambda_1(U_1 \triangle (U_\sigma - u_0)) \geqslant \frac{1}{32s^2} \lambda_1(U_1)(1 - \lambda_1(U_1)), \quad \sigma = 1, \dots, s,$$
 (3.8)

hold simultaneously has Lebesgue measure at least 1/2.

It then follows on combining (3.6)–(3.8) that the set of values $u_0 \in [0,1)$ for which the inequalities

$$\lambda_1(U_1 \triangle (U_\sigma - u_0)) \geqslant \frac{1}{32s^4}, \quad \sigma = 1, \dots, s, \tag{3.9}$$

hold simultaneously has Lebesgue measure at least 1/2. Let

$$\mathscr{S}(W; \mathbf{x}_1) = \{ u_0 \in [0, 1) : (3.9) \text{ holds for every } \sigma = 1, \dots, s \}$$
 (3.10)

denote this set of such values of u_0 . Note that the condition (3.9) is equivalent to the condition

$$\lambda_1(\{z \in [0,1) : \chi_{\mathcal{W}}(\mathbf{x}_1, z)\chi_{\mathcal{W}}(\mathbf{x}_{\sigma}, \{z + u_0\}) = 0\}) \geqslant \frac{1}{32s^4}, \quad \sigma = 1, \dots, s. \quad (3.11)$$

Let us revisit the disjoint union $\mathcal{B}^* = \mathcal{B}^*(\mathcal{B}; \varepsilon_2; \varepsilon_3)$ of axis-parallel rectangular boxes. For each axis-parallel rectangular box in this union, we push each boundary face inwards by a distance ε_6 , where the parameter $\varepsilon_6 > 0$, to be specified later, satisfies

$$0 < \varepsilon_6 < \frac{\varepsilon_3}{2}$$

where ε_3 is a lower bound of the side lengths of these axis-parallel rectangular boxes. Then the resulting smaller axis-parallel rectangular box has volume which is at least $(1-2\varepsilon_6/\varepsilon_3)^3$ times that of the original axis-parallel rectangular box. This means that if $\mathcal{B}^{**} = \mathcal{B}^{**}(\mathcal{B}^*; \varepsilon_6)$ is the disjoint union of these smaller axis-parallel rectangular boxes, then using Lemma 3.2(i) and writing

$$\varepsilon_7 = \varepsilon_2 + \frac{6s\varepsilon_6}{\varepsilon_3},\tag{3.12}$$

we have

$$\lambda_{3}(\mathcal{B}^{**}) \geqslant \left(1 - \frac{2\varepsilon_{6}}{\varepsilon_{3}}\right)^{3} \lambda_{3}(\mathcal{B}^{*}) \geqslant \left(1 - \frac{2\varepsilon_{6}}{\varepsilon_{3}}\right)^{3} (s - \varepsilon_{2})$$

$$\geqslant \left(1 - \frac{6\varepsilon_{6}}{\varepsilon_{3}}\right) (s - \varepsilon_{2}) \geqslant s - \varepsilon_{2} - \frac{6s\varepsilon_{6}}{\varepsilon_{3}} = s - \varepsilon_{7}. \tag{3.13}$$

Remark. Note that the ε_6 -neighbourhood of every point in \mathcal{B}^{**} is contained in \mathcal{B}^* .

Observe that $\chi_{\mathcal{B}}(\mathbf{s}') = \chi_{\mathcal{B}}(\mathbf{s}'')$ if the points $\mathbf{s}', \mathbf{s}'' \in \mathcal{M}$ belong to the same axis-parallel rectangular box in the disjoint union \mathcal{B}^{**} , and that the side lengths of each axis-parallel rectangular box in the disjoint union \mathcal{B}^{**} is greater than $\varepsilon_3 - 2\varepsilon_6$. Furthermore, analogous to Lemma 3.2(iv), we have the inequality

$$\lambda_2\left(\left\{\mathbf{x} \in \mathcal{P} : \lambda_1(\left\{z \in [0,1) : (\mathbf{x},z) \in \mathcal{B}^{**}\right\}) > 1 - \varepsilon_7^{1/2}\right\}\right) > s - \varepsilon_7^{1/2}.$$
 (3.14)

For any $\mathbf{x} \in \mathcal{P}$, let $\mathbf{x}_1, \dots, \mathbf{x}_s \in \mathcal{P}$ be distinct points such that their images under projection modulo 1 to the unit torus $[0,1)^2$ coincides with the image of \mathbf{x} under the same projection. Then it follows from (3.14) that

$$\lambda_2 \left(\left\{ \mathbf{x} \in \mathcal{P} : \lambda_1(\left\{ z \in [0, 1) : (\mathbf{x}_{\sigma}, z) \in \mathcal{B}^{**} \right\}) > 1 - \varepsilon_7^{1/2} \text{ for every } \sigma = 1, \dots, s \right\} \right)$$

$$> s - s\varepsilon_7^{1/2}.$$
(3.15)

For every $u_0 \in [0, 1)$, let

$$\mu(u_0) = \lambda_2(\{\mathbf{x}_1 \in \mathcal{P} : u_0 \in \mathscr{S}(\mathcal{W}; \mathbf{x}_1)\})$$

denote some relevant multiplicity of u_0 . Then

$$\int_{[0,1)} \mu(u_0) du_0 = \int_{\mathcal{P}} \lambda_1(\mathscr{S}(\mathcal{W}; \mathbf{x}_1)) d\mathbf{x}_1 \geqslant \frac{s}{2},$$

in view of (3.10). This can be interpreted to say that the average multiplicity $\mu(u_0)$ of $u_0 \in [0,1)$ is at least s/2. Thus there exists a shift $u_0^* \in [0,1)$ such that

$$\mu(u_0^*) = \lambda_2(\{\mathbf{x}_1 \in \mathcal{P} : u_0^* \in \mathscr{S}(\mathcal{W}; \mathbf{x}_1)\}) \geqslant \frac{s}{2}.$$

It also follows from (3.15) that with the exception of at most $s\varepsilon_7^{1/2}$ part of $\mathbf{x} \in \mathcal{P}$, at least $1 - 2\varepsilon_7^{1/2}$ part of the real numbers $z \in [0, 1)$ are such that with $\mathbf{x} = \mathbf{x}_1$,

$$(\mathbf{x}_{\sigma}, z) \in \mathcal{B}^{**}, \quad (\mathbf{x}_{\sigma}, \{z + u_0^*\}) \in \mathcal{B}^{**}, \quad \sigma = 1, \dots, s.$$

In order to derive a contradiction, we need a density analogue of [2, Lemma 6.2]. Here $\|\beta\|$ denotes the distance of $\beta \in \mathbb{R}$ to the nearest integer.

Lemma 3.4. Let $\mathbf{v} = (v_1, v_2) \in \mathbb{R}^2$ be a Kronecker vector, and let $w \in \mathbb{R}$ be arbitrary such that $\mathbf{w} = (v_1, v_2, w) \in \mathbb{R}^3$ is a Kronecker vector. Let $\varepsilon_6 > 0$ be given. There exists a finite sequence

$$1 \leqslant m_1(\varepsilon_6) < \ldots < m_k(\varepsilon_6)$$

of positive integers such that

$$||m_j(\varepsilon_6)v_1|| < \varepsilon_6, \quad ||m_j(\varepsilon_6)v_2|| < \varepsilon_6, \quad j = 1, \dots, k,$$
 (3.16)

and the finite sequence $\{m_j(\varepsilon_6)w\}$, $j=1,\ldots,k$, visits every subinterval of [0,1) with length ε_6 .

Remark. Note that the number $k = k(\mathbf{w}; \varepsilon_6)$ of terms of the finite sequence may depend on the choice of \mathbf{w} and the value of ε_6 .

Proof of Lemma 3.4. By the Kronecker density theorem, the sequence

$$j\mathbf{w} = j(v_1, v_2, w), \quad j = 1, 2, 3, \dots,$$

modulo 1 is dense in the unit torus $[0,1)^3$. Let

$$m_1(\varepsilon_6), m_2(\varepsilon_6), m_3(\varepsilon_6), \dots$$

be the infinite subsequence of $1, 2, 3, \ldots$ such that

$$\{m_j(\varepsilon_6)v_1\}, \{m_j(\varepsilon_6)v_2\} \in [0, \varepsilon_6) \cup (1 - \varepsilon_6, 1), \quad j = 1, 2, 3, \dots$$

Clearly (3.16) holds. Also, the subsequence $m_j(\varepsilon_6)\mathbf{w}$, $j=1,2,3,\ldots$, modulo 1 is dense in $([0,\varepsilon_6)\cup(1-\varepsilon_6,1))^2\times[0,1)\subset[0,1)^3$. This implies that the sequence $\{m_j(\varepsilon_6)w\}$, $j=1,2,3,\ldots$, is dense in [0,1).

Next, let the integer n satisfy $n > 2/\varepsilon_6$. We now partition the unit torus [0,1) into n short intervals I_1, \ldots, I_n of length 1/n in the standard way. Then for every $i = 1, \ldots, n$, there exists an integer k_i such that the finite sequence

$$\{m_j(\varepsilon_6)w\}, \quad j=1,\ldots,k_i,$$

visits I_i . Let $k = \max\{k_1, \ldots, k_n\}$. Then the finite sequence

$$\{m_j(\varepsilon_6)w\}, \quad j=1,\ldots,k, \tag{3.17}$$

visits every interval I_1, \ldots, I_n . Now, since $1/n < \varepsilon_6/2$, every subinterval of [0, 1) with length ε_6 must contain at least one of the intervals I_1, \ldots, I_n , so is visited by the finite sequence (3.17).

It follows that there exists $j_0 = j_0(u_0^*)$ satisfying $1 \leq j_0 \leq k$ such that

$$||m_{j_0}(\varepsilon_6)v_1|| < \varepsilon_6, \quad ||m_{j_0}(\varepsilon_6)v_2|| < \varepsilon_6, \quad ||m_{j_0}(\varepsilon_6)w - u_0^*|| < \varepsilon_6.$$
 (3.18)

Remark. For a fixed point $\mathbf{x} \in \mathcal{P}$, we can visualize the set $\{(\mathbf{x}, z) : z \in [0, 1)\}$ as a circle over the point \mathbf{x} , since [0, 1) is the unit torus. For any point (\mathbf{x}, z) on this circle, we have $\mathbf{T}^*_{\mathbf{w}}(\mathbf{x}, z) = (\mathbf{x} + \mathbf{v}, \{z + w\})$. It follows that the image of the circle under the transformation $\mathbf{T}^*_{\mathbf{w}}$ is a circle $\{(\mathbf{y}, z') : z' \in [0, 1)\}$ over the point $\mathbf{y} = \mathbf{x} + \mathbf{v}$. Clearly, this new circle is rotated from the original circle by a quantity w and its position on \mathcal{P} is translated from the original position by a vector \mathbf{v} . This action is repeated multiple times when we apply the transformation $\mathbf{T}^*_{\mathbf{w}}$ successively. Our proof of Theorem $4((i) \Rightarrow (ii))$ is based on a combination of this fact with Lemmas 3.2-3.4.

Let us continue the discussion prior to Lemma 3.4. For at least $s/2 - s\varepsilon_7^{1/2}$ part of the points $\mathbf{x} \in \mathcal{P}$, we have $u_0^* \in \mathscr{S}(\mathcal{W}; \mathbf{x})$, and that for at least $1 - 2\varepsilon_7^{1/2}$ part of the real numbers $z \in [0, 1)$, writing $\mathbf{x} = \mathbf{x}_1$, we have

$$(\mathbf{x}, z) \in \mathcal{B}^{**}, \quad (\mathbf{x}_{\sigma}, \{z + u_0^*\}) \in \mathcal{B}^{**}, \quad \sigma = 1, \dots, s.$$
 (3.19)

We say that the points in (3.19) form a good (s+1)-tuple, in the sense that they all belong to \mathcal{B}^{**} .

Let

$$Q = (\mathbf{x}, z) = (\mathbf{x}_1, z), \quad Q_{\sigma} = (\mathbf{x}_{\sigma}, \{z + u_0^*\}), \quad \sigma = 1, \dots, s,$$

be such a good (s+1)-tuple, and consider the new point $Q^* = (\mathbf{T}_{\mathbf{w}}^*)^{m_{j_0}}(Q)$, obtained from Q by m_{j_0} successive applications of the transformation $\mathbf{T}_{\mathbf{w}}^*$, where j_0 is chosen so that the inequalities (3.18) hold. Since the subset $\mathcal{W} \subset \mathcal{M}$ is invariant under the transformation $\mathbf{T}_{\mathbf{w}}^*$, it follows that $\chi_{\mathcal{W}}(Q) = \chi_{\mathcal{W}}(Q^*)$.

On the other hand, as observed in the Remark above, the image of the circle $\{(\mathbf{x}, z) : z \in [0, 1)\}$ under the transformation $(\mathbf{T}_{\mathbf{w}}^*)^{m_{j_0}}$ is a circle $\{(\mathbf{y}, z') : z' \in [0, 1)\}$ with some particular $\mathbf{y} = \mathbf{x}_1 + m_{j_0}\mathbf{v} \in \mathcal{P}$. It then follows from (3.18) that the coordinates of \mathbf{y} are ε_6 -close to the corresponding coordinates of \mathbf{x}_{σ_0} for a particular $\sigma_0 = 1, \ldots, s$, and the last coordinate of Q^* is ε_6 -close to $\{z + u_0^*\}$. Since $Q_{\sigma_0} \in \mathcal{B}^{**}$, it follows from the Remark after (3.13) that Q_{σ_0} and Q^* are in the same axis-parallel rectangular box in the disjoint union \mathcal{B}^* , and so $\chi_{\mathcal{B}}(Q_{\sigma_0}) = \chi_{\mathcal{B}}(Q^*)$, in view of Lemma 3.2(ii).

Recall that $u_0^* \in \mathcal{S}(\mathcal{W}; \mathbf{x})$ for at least $s/2 - s\varepsilon_7^{1/2}$ part of the points $\mathbf{x} \in \mathcal{P}$. This and the condition (3.11) together imply that

$$\lambda_1(\{z \in [0,1) : \chi_{\mathcal{W}}(\mathbf{x}_1, z)\chi_{\mathcal{W}}(\mathbf{x}_{\sigma_0}, \{z + u_0^*\}) = 0\}) \geqslant \frac{1}{32s^4},$$

and this is equivalent to

$$\lambda_1(\{z \in [0,1) : \chi_{\mathcal{W}}(Q) \neq \chi_{\mathcal{W}}(Q_{\sigma_0})\}) \geqslant \frac{1}{32s^4}.$$
 (3.20)

Consider now the three relations

$$\chi_{\mathcal{W}}(Q) = \chi_{\mathcal{W}}(Q^*), \quad \chi_{\mathcal{B}}(Q_{\sigma_0}) = \chi_{\mathcal{B}}(Q^*), \quad \chi_{\mathcal{W}}(Q) \neq \chi_{\mathcal{W}}(Q_{\sigma_0}),$$

which clearly imply the two relations

$$\chi_{\mathcal{W}}(Q_{\sigma_0}) \neq \chi_{\mathcal{W}}(Q^*), \quad \chi_{\mathcal{B}}(Q_{\sigma_0}) = \chi_{\mathcal{B}}(Q^*).$$

It follows that

$$\chi_{\mathcal{W}}(Q_{\sigma_0}) \neq \chi_{\mathcal{B}}(Q_{\sigma_0}) \quad \text{or} \quad \chi_{\mathcal{W}}(Q^*) \neq \chi_{\mathcal{B}}(Q^*).$$
(3.21)

Intuitively speaking, (3.21) represents two *negligible* cases, with total measure less than ε_1 , in view of (3.3), which contradict the substantial constant lower bound in (3.20) if $\varepsilon_1 > 0$ is sufficiently small. To make this precise, we need to study more closely the various parameters.

We have $u_0^* \in \mathscr{S}(\mathcal{W}; \mathbf{x})$ for at least $s/2 - s\varepsilon_7^{1/2}$ part of the points $\mathbf{x} \in \mathcal{P}$. Using (3.20), we deduce that for at least

$$\frac{1}{32s^4} - 2\varepsilon_7^{1/2} \tag{3.22}$$

part of the real numbers $z \in [0, 1)$, the points

$$Q_{\sigma_0} = (\mathbf{x}_{\sigma_0}, \{z + u_0^*\}) \text{ and } Q^* = (\mathbf{T}_{\mathbf{w}}^*)^{m_{j_0}}(Q) = (\mathbf{T}_{\mathbf{w}}^*)^{m_{j_0}}(\mathbf{x}_1, z)$$
 (3.23)

exhibit the property (3.21). It follows from (3.22) that the 3-dimensional Lebesgue measure of the points $Q = (\mathbf{x}_1, z) \in \mathcal{M}$ such that the points (3.23) exhibit the property (3.21) is at least

$$\left(\frac{s}{2} - s\varepsilon_7^{1/2}\right) \left(\frac{1}{32s^4} - 2\varepsilon_7^{1/2}\right),\tag{3.24}$$

where ε_7 is given by (3.12).

On the other hand, the property (3.21) is *exceptional*, and (3.3) implies that the quantity in (3.24) is less than $2\varepsilon_1$. We emphasize the fact that the choice of the parameter ε_1 is independent of the choices of the other parameters $\varepsilon_2, \varepsilon_3, \varepsilon_6$. Thus we can make ε_7 in (3.24) arbitrarily small independently of the fixed value of ε_1 . It is therefore easy to specify the parameters $\varepsilon_1, \varepsilon_2, \varepsilon_3, \varepsilon_6$ so that the value of the quantity in (3.24) is greater than $2\varepsilon_1$, leading to a contradiction.

The contradiction establishes the ergodicity of the w-shift in \mathcal{M} .

The last step is to extend ergodicity of the **w**-shift in \mathcal{M} to unique ergodicity by using the standard argument in functional analysis. This is possible, since the projection of \mathcal{M} to the unit torus $[0,1)^3$ leads to unique ergodicity there.

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