

An infinite interval version of the α -Kakutani equidistribution problem

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March 17, 2021

Abstract

In this article we extend results of Kakutani, Adler–Flatto, Smilansky and others on the classical α -Kakutani equidistribution result for sequences arising from finite partitions of the interval. In particular, we describe a generalization of the equidistribution result to infinite partitions. In addition, we give discrepancy estimates, extending results of Drmota–Infusino [8].

1 Introduction

Uniform distribution of sequences of numbers $(x_n)_{n=1}^{\infty}$ in the unit interval has been an important area of interest for over a century. For example, it was shown by Weyl [20] that, for any irrational α , the sequence $x_n = \alpha n \pmod{1}$ is uniformly distributed and Hardy and Littlewood showed that, for almost all $\lambda > 1$, the sequence $x_n = \lambda^n \pmod{1}$ is uniformly distributed [13]. In this note we consider another simple family of examples based on subdividing intervals. Before introducing the original motivating example, we first fix our terminology: a *partition* \mathcal{P} is a set of closed, positive-length intervals, which have pairwise disjoint interiors and cover $[0, 1]$ up to a set of Lebesgue measure zero.

1.1 The original α -Kakutani equidistribution result

In 1973 Araki posed a problem which led to Kakutani to prove an elegant equidistribution result. (An interesting historical background is presented in [1]). For clarity, we give a description of his original partition scheme, which we generalise in the next section.

*Partly supported by ERC-Advanced Grant 833802-Resonances and EPSRC grant EP/T001674/1

Definition (α -Kakutani scheme). For a fixed $0 < \alpha < 1$, the α -Kakutani scheme is a sequence of partitions $(\mathcal{P}_n)_{n=0}^\infty$ defined inductively:

- $\mathcal{P}_0 = \{[0, 1]\}$ is the trivial partition; and
- \mathcal{P}_{n+1} is obtained from \mathcal{P}_n by taking each interval of maximal length and subdividing it into two smaller intervals in the ratio $\alpha : 1 - \alpha$.

Example 1. Figure 1 shows the first seven partitions for the choice $\alpha = 1/3$. Notice that \mathcal{P}_5 is obtained by splitting two maximal length intervals in \mathcal{P}_4 simultaneously (each of length $2/9$). By contrast, the choice $\alpha = 1/2$ gives the trivial dyadic splitting.

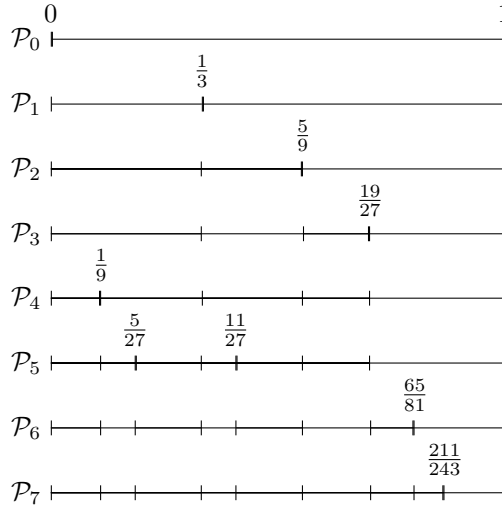


Figure 1: The first seven partitions $(\mathcal{P}_n)_{n=0}^7$ of the $\frac{1}{3}$ -Kakutani scheme.

Consider the set of endpoints at the n th stage of this process, $E_n = \bigcup_{I \in \mathcal{P}_n} \partial I$. Kakutani's result is the following.

Theorem (Kakutani). For all $\alpha \in (0, 1)$, the set E_n becomes uniformly distributed as $n \rightarrow \infty$.

1.2 Interval substitutions using multiple intervals

A natural generalisation of the α -Kakutani scheme, first introduced by Volčič in [18], is to alter the above process by splitting intervals of maximal length according to a fixed finite partition consisting of $N \geq 2$ subintervals, say. That is, at each stage, one splits all intervals of maximal length into N pieces whose lengths have a certain fixed ratio, $\alpha_1 : \alpha_2 : \dots : \alpha_N$, where the α_i sum to 1.

Example 2. In Figure 2, we have the first seven partitions of the interval substitution scheme in which one splits maximal intervals according to the partition $\{[0, \frac{1}{2}], [\frac{1}{2}, \frac{2}{3}], [\frac{2}{3}, 1]\}$, i.e., with ratio $\frac{1}{2} : \frac{1}{6} : \frac{1}{3}$. By contrast, the α -Kakutani scheme corresponds to splitting according to the partition $\{[0, \alpha], [\alpha, 1]\}$.

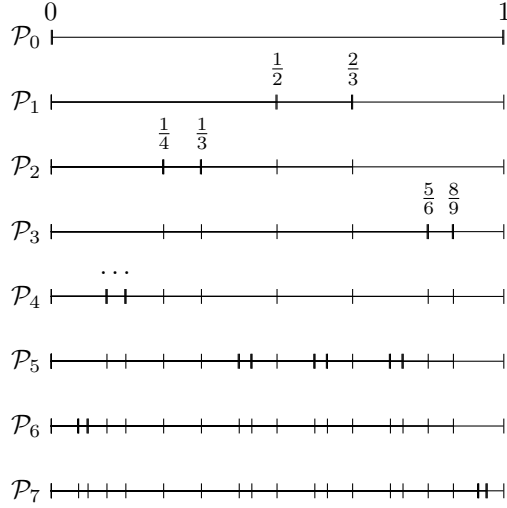


Figure 2: The first seven partitions $(\mathcal{P}_n)_{n=0}^7$ of the interval substitution scheme where one splits maximal-length intervals according to the partition $\mathcal{P}_1 = \{[0, \frac{1}{2}], [\frac{1}{2}, \frac{2}{3}], [\frac{2}{3}, 1]\}$.

One particular family of these schemes, that of so-called LS-sequences, which include the ϕ -Kakutani scheme for $\phi = \frac{1}{2}(\sqrt{5} - 1)$, has received particular attention in the context of low discrepancy sequences: see, e.g., [3, 5, 6, 14, 19].

1.3 Interval substitutions using infinitely many intervals

In this note, we continue the process and ask in what way the result above still holds if at every stage we insert an *infinite* partition \mathcal{P} into each maximal-length subinterval. We denote by \hat{E}_n the finite set of endpoints of those intervals which have been split up to the $(n + 1)$ -st stage, i.e.,

$$\hat{E}_n := \left\{ \min(I), \max(I) : I \in \bigcup_{i=0}^n \mathcal{P}_i \setminus \mathcal{P}_{n+1} \right\}.$$

Example 3. In Figure 3 we depict \mathcal{P}_n and \hat{E}_n (for $n \leq 7$) for the infinite substitution scheme generated by $\mathcal{P} = \{[0, \frac{1}{2}]\} \cup \{[1 - \frac{1}{2} \cdot 3^{-n}, 1 - \frac{1}{6} \cdot 3^{-n}]\}_{n=0}^{\infty}$.

Example 4. A wilder example: let the intervals in \mathcal{P} be the connected components of the complement of the middle third Cantor set.

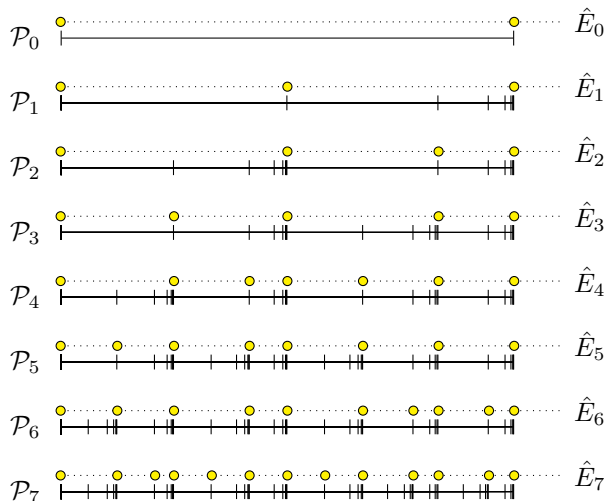


Figure 3: An illustration of $(\mathcal{P}_n)_{n=0}^7$ and $(\hat{E}_n)_{n=0}^7$ for the example generated by the partition $\mathcal{P} = \{[1, \frac{1}{2}]\} \cup \{[1 - \frac{1}{2} \cdot 3^{-n}, 1 - \frac{1}{6} \cdot 3^{-n}]\}_{n=0}^{\infty}$. Here the tick marks (which accumulate on certain points in the interval) denote the elements of E_n and the suspended yellow circles denote the elements of \hat{E}_n .

For simplicity, we restrict our attention to the set of left endpoints, which we shall denote by L_n , although we could equally well have chosen the right endpoints, midpoints, etc.

Our main result is the following generalization of Kakutani's equidistribution theorem. Let $\|I\|$ denote the length of an interval I .

Theorem 1. *Let \mathcal{P} be a countable partition. Then, provided that*

$$-\sum_{I \in \mathcal{P}} \|I\| \log \|I\| < \infty,$$

the set L_n becomes uniformly distributed as $n \rightarrow \infty$.

A brief outline of this note: In section 2, we give a new dynamical viewpoint of the problem and in section 3, we apply renewal theory to prove Theorem 1 in two cases. In section 4 we use a generating function to estimate the discrepancy in the rank one case. In section 5 we use methods of analytic number theory to estimate the discrepancy in the higher rank case, with a generic Diophantine-type assumption.

Our interest in this problem, and the starting point for our analysis, began with the very elegant work of Smilansky [17].

2 Partitions and similarities

Our approach to Theorem 1 is to express the elements of the partition \mathcal{P} in terms of the images of similarities. The refinements into finer partitions, \mathcal{P}_n , can then be expressed in terms of words formed from the index set of \mathcal{P} .

That is, one can write $\mathcal{P} = \{T_i[0, 1]\}_{i \in \mathcal{I}}$, where each $T_i : [0, 1] \rightarrow [0, 1]$ is an orientation preserving similarity with contraction ratio $\alpha_i > 0$. We see that \mathcal{P} being a partition is equivalent to the following:

- $T_i[0, 1] \cap T_j[0, 1] = \emptyset$ for $i \neq j$; and
- $\sum_{i \in \mathcal{I}} \alpha_i = 1$.

Example 5. Given $0 = t_0 < t_1 < \dots$ with $(t_n) \rightarrow 1$, the partition $\{[t_{n-1}, t_n]\}_{n \in \mathbb{N}}$ is equal to $\{T_n[0, 1]\}_{n \in \mathbb{N}}$, where

$$T_n(x) = (t_n - t_{n-1})x + \sum_{k=0}^{n-1} t_k.$$

Example 6. Setting $t_n = 1 - \frac{1}{6} 3^{-n}$ for $n \geq 1$ gives rise to \mathcal{P}_1 in Figure 3.

We now explain how this can be used to give an explicit description of the splitting process.

Definition ((T_i) -refinement). Given a partition $\mathcal{P} = \{S_k[0, 1]\}_k$, where the $\{S_k\}_k$ are orientation preserving similarities, the (T_i) -refinement of \mathcal{P} is the refinement obtained by taking all intervals of maximal length in \mathcal{P} and replacing them by subintervals in the following manner: if $S[0, 1] \in \mathcal{P}$ has maximal length in \mathcal{P} , then it is replaced by the elements of the set

$$\{S \circ T_i[0, 1] \mid i \in \mathcal{I}\}.$$

Definition (Interval substitution scheme). The interval substitution scheme generated by $\{T_i\}_{i \in \mathcal{I}}$ is the sequence of partitions $(\mathcal{P}_n)_{n=0}^\infty$ defined as follows:

- \mathcal{P}_0 is the trivial partition, $\mathcal{P}_0 = \{[0, 1]\}$; and
- \mathcal{P}_{n+1} is the (T_i) -refinement of \mathcal{P}_n .

This gives a convenient presentation of the partitions.

Example 7. The α -Kakutani scheme is the interval substitution scheme generated by the pair $T_1 : x \mapsto \alpha x$ and $T_2 : x \mapsto (1 - \alpha)x + \alpha$.

Example 8. Similarly, the interval substitution scheme generated by the triple $T_1 : x \mapsto x/2$, $T_2 : x \mapsto (x + 3)/6$ and $T_3 : x \mapsto (x + 2)/3$ gives the sequence of partitions depicted in Figure 2.

We now associate to the sequence of partitions $(\mathcal{P}_n)_{n=0}^\infty$ a sequence of families of left endpoints of split intervals, $(L_n)_{n=0}^\infty$.

Definition (L_n). Given an interval substitution scheme $(\mathcal{P}_n)_{n=0}^\infty$ generated by similarities $(T_i)_{i \in \mathcal{I}}$, we define the finite sets L_n ($n \geq 0$) to be

$$L_n = \bigcup_{k=0}^n \bigcup_{I \in \mathcal{P}_k \setminus \mathcal{P}_{k+1}} \min(I).$$

Remark. One can consider a generalisation of the above process by dropping the assumption that the $(T_i)_i$ have to be affine. We will not consider this more general setting, but we note it is easy to give superficial examples where E_n (or L_n) is not uniformly distributed: take, for example, $T_1(x) = \sqrt{x}/2$, $T_2(x) = (x+1)/2$.

Considering the interval substitution scheme generated by $\{T_i\}_{i \in \mathcal{I}}$, it follows inductively that every interval appearing in the process is obtained by applying a sequence of maps from $\{T_i\}_i$ to $[0, 1]$, and so each is naturally described by a finite word in \mathcal{I} . It is convenient to introduce the following notation.

Definition ($W(\mathcal{I})$, $*$, α_v , T_v). Given a countable set \mathcal{I} , the word set $W(\mathcal{I})$ is the semigroup consisting of all words in \mathcal{I} : i.e.

$$W(\mathcal{I}) = \{\emptyset\} \cup \bigcup_{n=1}^{\infty} \mathcal{I}^n,$$

where \emptyset denotes the empty word (unique word of length zero), and the semigroup operation $*$: $W(\mathcal{I}) \times W(\mathcal{I}) \rightarrow W(\mathcal{I})$ denotes concatenation of words, for which \emptyset acts as the identity:

$$(n_1, \dots, n_k) * (m_1, \dots, m_j) = (n_1, \dots, n_k, m_1, \dots, m_j); \quad \mathbf{v} * \emptyset = \emptyset * \mathbf{v} = \mathbf{v}.$$

Furthermore, for ease of notation, we extend the definitions of α_i and T_i to the whole of $W(\mathcal{I})$: For the word $\mathbf{v} = (i_1, \dots, i_k) \in \mathcal{I}^k$, define

$$\alpha_{\mathbf{v}} := \prod_{j=1}^k \alpha_{i_j}, \quad T_{\mathbf{v}} := T_{i_1} \circ \dots \circ T_{i_k},$$

and also define $\alpha_{\emptyset} = 1$ and $T_{\emptyset} = \text{Id}_{[0,1]}$.

That is, given any interval I which appears in the process, $I = T_{\mathbf{v}}[0, 1]$ for some word $\mathbf{v} \in W(\mathcal{I})^*$, and will be split between \mathcal{P}_n and \mathcal{P}_{n+1} (i.e., $I \in \mathcal{P}_n \setminus \mathcal{P}_{n-1}$) precisely when n satisfies

$$\alpha_{\mathbf{v}} = \max_{I \in \mathcal{P}_n} \{\|I\|\},$$

and consequently its left endpoint $T_{\mathbf{v}}(0)$ will appear in L_n , if not already present in L_{n-1} .

Rather than using n in $\{L_n\}_{n \geq 1}$ to parametrise this process, we want to reparameterise this family to reflect the lengths of the maximal intervals, and rewrite it as (X_λ) as follows.

*Similarly, every word $\mathbf{v} \in \mathcal{I}$ gives rise to an interval in some \mathcal{P}_n

Definition. For $\lambda > 1$, let $X_\lambda = \emptyset$, and for $\lambda \in (0, 1]$, let

$$X_\lambda := L_{n(\lambda)}, \quad \text{where } n(\lambda) := \max\{n \geq 0 : \forall I \in \mathcal{P}_n, \|I\| \geq \lambda\}.$$

I.e., given $\lambda \leq 1$, $P_{n(\lambda)+1}$ is the first partition in the process for which all intervals have lengths strictly smaller than λ . From the previous discussion, one obtains a convenient, dynamical formula for X_λ :

$$X_\lambda = \{T_{\mathbf{v}}(0) : \mathbf{v} \in W(\mathcal{I}), \alpha_{\mathbf{v}} \geq \lambda\}.$$

As $\lambda \rightarrow 0^+$, $n(\lambda) \rightarrow \infty$ and so the uniform distribution of L_n as $n \rightarrow \infty$ is equivalent to that of X_λ as $\lambda \rightarrow 0^+$. As an aside, this is equivalent to the convergence of the probability measure μ_λ ,

$$\mu_\lambda = \frac{1}{|X_\lambda|} \sum_{x \in X_\lambda} \delta_x,$$

to the Lebesgue measure (which we will henceforth denote as Leb), as $\lambda \rightarrow 0^+$, where δ_x denotes the Dirac delta measure at x , and $|\cdot|$ denotes cardinality (of a finite set).

3 Proof of Theorem 1

This section is devoted to proving the main result, which we can conveniently rephrase in the following way.

Theorem 1. *Provided that $-\sum_{I \in \mathcal{P}} \|I\| \log \|I\| < \infty$, the measures μ_λ converge to the Lebesgue measure Leb as Borel measures as $\lambda \rightarrow 0^+$, i.e., for any interval $J \subset [0, 1]$, we have $\mu_\lambda(J) \rightarrow \text{Leb}(J) = \|J\|$.*

We first prove the convergence for a given interval of the form $T_{\mathbf{v}}[0, 1]$. For each of these elementary sets, their μ_λ -measure is intimately related to the asymptotics of $|X_\lambda|$ as $\lambda \rightarrow 0^+$. To proceed, we can relate $|X_\lambda|$ to the set of words $|A_\lambda|$, where

$$A_\lambda := \{\mathbf{w} \in W(\mathcal{I}) \mid \alpha_{\mathbf{w}} \geq \lambda\},$$

and exploit a natural renewal equation for the quantity $\lambda|A_\lambda|$.

Lemma 1. *One of the following two cases hold. Either there is a (unique) symbol in $1 \in \mathcal{I}$, say, such that $T_1(0) = 0$; in which case, for all $\lambda > 0$,*

$$|A_\lambda| = \sum_{k=0}^{\infty} |X_{\lambda\alpha_1^{-k}}| \iff |X_\lambda| = |A_\lambda| - |A_{\lambda/\alpha_1}|; \quad (1)$$

or no element of $\{T_{\mathbf{w}}\}_{\mathbf{w} \in W(\mathcal{I})}$ fixes 0, and $|X_\lambda| = |A_\lambda|$ for all $\lambda > 0$.

Proof of Lemma 1. Let $\mathbf{w}, \mathbf{v} \in A_\lambda$ satisfy $T_{\mathbf{w}}(0) = T_{\mathbf{v}}(0)$. If $\mathbf{w} \neq \mathbf{v}$, the disjointness of $T_i[0, 1)$ and the injectivity of the T_i imply inductively that one of these words is obtained from the other by concatenation with a word fixing the identity; i.e., without loss of generality, $\mathbf{w} = \mathbf{v} * \mathbf{j}$, where $\mathbf{j} \in W(\mathcal{I})$ satisfies $T_{\mathbf{j}}(0) = 0$. We now have the two cases:

If there is no symbol $i \in \mathcal{I}$ for which $T_i(0) = 0$, we must have $\mathbf{j} = \emptyset$; which implies the map $\mathbf{v} \mapsto T_{\mathbf{v}}(0)$ is a bijection $A_\lambda \rightarrow X_\lambda$.

For the other case, let $1 \in \mathcal{I}$ be such that $T_1(0) = 0$; it is unique by the disjointness of $T_i[0, 1)$. Then $\mathbf{j} \neq \emptyset$ only if $\mathbf{j} \in \{1\}^k \subset \mathcal{I}^k$ for some $k \in \mathbb{N}$, i.e., \mathbf{j} is a tuple of 1's. It follows that, for each $y \in X_\lambda$, there is a unique word $\mathbf{v}_0(y) \in A_\lambda$ satisfying

- $T_{\mathbf{v}_0(y)}(0) = y$; and
- $T_{\mathbf{v}}(0) = y \implies \mathbf{v} = \mathbf{v}_0(y) * \mathbf{j}$, for some $\mathbf{j} \in \{\emptyset\} \cup \{1\}^k$, $k \in \mathbb{N}$.

In particular, $T_{\mathbf{v}}(0) = y$ implies $\alpha_{\mathbf{v}} = \alpha_{\mathbf{v}_0(y)} \alpha_1^k$, for some $k \in \mathbb{N}_0$. In other words, there is exactly one element of $A_\lambda \setminus A_{\lambda/\alpha_1}$, $\mathbf{v}_0(y)$, which gets mapped onto y . Therefore, $\mathbf{v} \mapsto T_{\mathbf{v}}(0)$ gives a bijection $A_\lambda \setminus A_{\lambda/\alpha_1} \rightarrow X_\lambda$. The right hand side of (1) follows, completing the proof. \square

Now, to continue the proof of the theorem, we combine the previous lemma with the following, which expresses $\mu_\lambda(T_{\mathbf{v}}[0, 1))$ as a ratio involving $|X_\lambda|$.

Lemma 2. *For all $\mathbf{v} \in W(\mathcal{I})$ and $\lambda \in (0, 1]$,*

$$|T_{\mathbf{v}}[0, 1) \cap X_{\lambda\alpha_{\mathbf{v}}}| = |X_\lambda|. \quad (2)$$

In particular, for all $\lambda \in (0, \alpha_{\mathbf{v}}]$,

$$\mu_\lambda(T_{\mathbf{v}}[0, 1)) = \frac{|X_{\lambda/\alpha_{\mathbf{v}}}|}{|X_\lambda|}. \quad (3)$$

*Moreover, if $\mathbf{v} = \emptyset$, or $\mathbf{v} = \mathbf{w} * i$ with $T_i(0) \neq 0$, these hold for all $\lambda > 0$.*

Proof of Lemma 2. Fix $\mathbf{v} \in W(\mathcal{I})$ and consider (2). First let $\lambda \leq 1$. By induction, using the properties of the T_i , if $T_{\mathbf{w}}(0) \in T_{\mathbf{v}}[0, 1)$, then, either $\mathbf{w} = \mathbf{v} * \mathbf{j}$ or $\mathbf{v} = \mathbf{w} * \mathbf{j}$, for some $\mathbf{j} \in W(\mathcal{I})$. Moreover, if $\mathbf{w} \in A_{\lambda\alpha_{\mathbf{v}}}$, the second option gives a contradiction: $\alpha_{\mathbf{w}} > \alpha_{\mathbf{v}} \geq \lambda\alpha_{\mathbf{v}} \geq \alpha_{\mathbf{w}}$. Thus, we have

$$T_{\mathbf{v}}[0, 1) \cap X_{\lambda\alpha_{\mathbf{v}}} = \{T_{\mathbf{v}*\mathbf{j}}(x) : \mathbf{j} \in A_\lambda\}.$$

By injectivity of $T_{\mathbf{v}}$, the right hand side bijectively corresponds to X_λ , giving (2). Applying (2), with $\lambda/\alpha_{\mathbf{v}}$ in place of λ and dividing by $|X_\lambda|$ yields (3).

As for the final remark, now let $\lambda > 1$. This implies the second option holds, $\mathbf{v} = \mathbf{w} * \mathbf{j}$, since the first is satisfied only if $\lambda \leq 1$, by a similar contradiction argument. It follows that $T_{\mathbf{j}}(0) = 0$ and thus, as in the proof of Lemma 1, $\mathbf{j} \in \{1\}^k$ for some $k \in \mathbb{N}$. This contradicts the assumptions of the remark; thus both the left and the right hand side of (2) are empty when $\lambda > 1$. The lemma follows. \square

The significance of relating $|X_\lambda|$ to $|A_\lambda|$ will now become clear from the following renewal equation.

Lemma 3. *The following holds for all $\lambda > 0$.*

$$|A_\lambda| = \sum_{i \in \mathcal{I}} |A_{\lambda/\alpha_i}| + \chi_{\{\lambda \leq 1\}}, \quad (4)$$

where χ is the indicator function. Equivalently, the following renewal equation applies, for $Z(t) := e^{-t}|A_{e^{-t}}|$.

$$Z(t) = \sum_{i \in \mathcal{I}} \alpha_i Z(t - \log(\alpha_i^{-1})) + e^{-t} \chi_{\{t \geq 0\}}, \quad \forall t \in \mathbb{R}.$$

Proof of Lemma 3. The second equation is a restatement of the first, (4). One can obtain (4) by partitioning the non-empty words in A_λ according to their first symbol. That is, one can write the following disjoint union.

$$A_\lambda \setminus \{\emptyset\} = \bigsqcup_{i \in \mathcal{I}} \underbrace{\{i * \mathbf{v} \in W(\mathcal{I}) \mid \alpha_{\mathbf{v}} \leq \lambda/\alpha_i\}}_{\text{in bijection with } A_{\lambda/\alpha_i}}.$$

This gives rise to the sum in (4). Accounting for \emptyset , $\emptyset \in A_\lambda$ if and only if $\lambda \leq 1$, thus providing the indicator term in (4), and completing the proof. \square

To make use of this renewal equation, just as in [2, 8, 17] it is necessary to consider two cases which behave somewhat differently. These cases correspond to, for example, the α -Kakutani schemes for $\alpha = 1/3$ and $\alpha = 1/2$, as described in the introduction.

Definition (Rank). For $n \in \mathbb{N}$ we will say the collection $\{\alpha_i\}_{i \in \mathcal{I}}$ is *rank n* if the smallest additive subgroup of \mathbb{R} containing the set $\{-\log(\alpha_i)\}_{i \in \mathcal{I}}$ is isomorphic to \mathbb{Z}^n . If $\{\alpha_i\}_i$ is not rank n for any $n \in \mathbb{N}$ we will say $\{\alpha_i\}_{i \in \mathcal{I}}$ is *infinite rank*. Also, whenever $\{\alpha_i\}_i$ is not rank one, we say it is *higher rank*.

Example 9. The following examples illustrate the different ranks:

1. $\{1/2^n\}_{n \in \mathbb{N}}$ is rank one;
2. $\{1/2\} \cup \{1/3^n\}_{n \in \mathbb{N}}$ is rank two;
3. $\{1/2\} \cup \{1/3\} \cup \{1/7^n\}_{n \in \mathbb{N}}$ is rank three; and
4. $\{1/n^s\}_{n=2}^\infty$ is infinite rank, where $s \approx 1.728$ satisfies $\zeta(s) = 2$ (here ζ denotes the Riemann zeta function).

3.1 Uniform distribution in the rank one case

In this subsection, we concentrate on the rank one case, also known as the *arithmetic, rationally-related* or *commensurable* case (see [2, 8, 17]). The characteristic feature of this case is that the contraction ratios are all powers of a common number, $\{\alpha_i\}_{i \in \mathcal{I}} \subset \{x^n\}_{n \in \mathbb{N}}$. Thereby, fixing the minimal such $x > 0$, equation (4) defines a discrete renewal equation on the lattice $-\log(x)\mathbb{Z}$, and one can apply the Erdős-Feller-Pollard renewal theorem (see [11]) to obtain the following lemma.

Lemma 4 (Erdős-Feller-Pollard renewal theorem). *Suppose that $\{\alpha_i\}_{i \in \mathcal{I}}$ is rank one, that $x > 0$ is the minimal positive number for which $\{\alpha_i\}_{i \in \mathcal{I}} \subset \{x^n\}_{n \in \mathbb{N}}$, and that $H := -\sum_i \alpha_i \log(\alpha_i) < \infty$, then*

$$Z(-n \log(x)) \rightarrow \frac{1}{H} \quad \text{as } n \rightarrow \infty, \quad n \in \mathbb{N};$$

where $Z(\cdot)$ is defined in Lemma 3.

This lemma yields the following corollary, all but completing the proof of the theorem in this case.

Corollary. *For $\{\alpha_i\}_{i \in \mathcal{I}}$ and H as in the previous lemma, the following hold as $n \rightarrow \infty$, $n \in \mathbb{N}$:*

- $|A_{x^n}| \sim x^{-n}/H$.
- For all $\mathbf{v} \in W(\mathcal{I})$, $\mu_{x^n}(T_{\mathbf{v}}[0, 1]) \rightarrow \alpha_{\mathbf{v}} = \text{Leb}(T_{\mathbf{v}}[0, 1])$.

Proof of Corollary. The first item is simply a restatement of the conclusion of Lemma 4. The second item follows with an application of Lemmas 1 and 2. For $n \in \mathbb{N}$ we have one of the following, according to the two cases of Lemma 1: Either no T_i fixes 0, and

$$\mu_{x^n}(T_{\mathbf{v}}[0, 1]) = \frac{|A_{x^n \alpha_{\mathbf{v}}^{-1}}|}{|A_{x^n}|} \sim \frac{(\alpha_{\mathbf{v}}^{-1} x^n)^{-1}/H}{x^{-n}/H} = \alpha_{\mathbf{v}};$$

or, for $T_1(0) = 0$ and \mathbf{v}_0 the longest subword of \mathbf{v} not ending in a 1,

$$\begin{aligned} \mu_{x^n}(T_{\mathbf{v}}[0, 1]) &= \frac{|A_{x^n \alpha_{\mathbf{v}}^{-1}}| - |A_{x^n (\alpha_{\mathbf{v}} \alpha_1)^{-1}}| + \chi_{\{\alpha_{\mathbf{v}} < x^n \leq \alpha_{\mathbf{v}_0}\}}}{|A_{x^n}| - |A_{x^n \alpha_1^{-1}}|} \\ &\sim \frac{1 - \alpha_1 \alpha_{\mathbf{v}} x^{-n}/H}{1 - \alpha_1 x^{-n}/H} = \alpha_{\mathbf{v}} \end{aligned}$$

as $n \rightarrow \infty$, as required. \square

The last, easy, step in the proof of Theorem 1 concerns packing intervals of the form $T_{\mathbf{v}}[0, 1]$ into a given interval. This will be written as if in the continuous case. For the rank one case, for $\lambda \rightarrow 0^+$, one can read $\lambda = x^n$, $n \rightarrow \infty$; in fact, it is a trivial matter to prove they are equivalent.

Proof of Theorem 1. Let $I \subset [0, 1]$ be an interval, and let $\alpha_{\max} = \max_{i \in \mathcal{I}}(\alpha_i)$. For $n \in \mathbb{N}$, let

$$U_n := \{\mathbf{v} \in \mathcal{I}^n : T_{\mathbf{v}}[0, 1] \subset I\}.$$

We claim the total length of intervals corresponding to U_n approximates the length of I , as follows:

$$\sum_{\mathbf{v} \in U_n} \|T_{\mathbf{v}}[0, 1]\| \geq \|I\| - 2\alpha_{\max}^n. \quad (5)$$

To prove this, take $x \in I \setminus \bigcup_{\mathbf{v} \in U_n} T_{\mathbf{v}}[0, 1]$. Then one of the following hold:

1. $x \in K_n := [0, 1] \setminus \bigcup_{\mathbf{v} \in \mathcal{I}^n} T_{\mathbf{v}}[0, 1]$. An inductive argument shows, for all λ and n , $\mu_\lambda(K_n) = 0 = \text{Leb}(K_n)$.
2. $x \in T_{\mathbf{w}}[0, 1]$ for some $\mathbf{w} \in \mathcal{I}^n \setminus U_n$. Since $T_{\mathbf{w}}[0, 1] \not\subset I$, it is an interval meeting ∂I ; thus there are at most two $\mathbf{w} \in \mathcal{I}^n$ with this property.

In other words, $I \setminus \bigcup_{\mathbf{v} \in U_n} T_{\mathbf{v}}[0, 1] \setminus K_n$ comprises at most two intervals, each with length at most α_{\max}^n . (5) follows.

To apply this, let U be a union of a finitely number of intervals in U_n , and let $\|U\|$ denote its total length (i.e., $\text{Leb}(U)$). Then $\mu_\lambda(I) \geq \mu_\lambda(U)$ and

$$\liminf_{\lambda \rightarrow 0^+} \mu_\lambda(I) \geq \lim_{\lambda \rightarrow 0^+} \mu_\lambda(U) = \|U\|.$$

Taking the supremum over all such finite unions U gives

$$\liminf_{\lambda \rightarrow 0^+} \mu_\lambda(I) \geq \|S\| \geq \|I\| - 2\alpha_{\max}^n.$$

Repeating the above argument for the 0, 1 or 2 intervals comprising $[0, 1] \setminus I$ yields a converse inequality for $\|I\|$:

$$\limsup_{\lambda \rightarrow 0^+} \mu_\lambda(I) \leq \|I\| + 2\alpha_{\max}^n,$$

and the proof is completed by taking $n \rightarrow \infty$. □

3.2 Uniform distribution in the higher rank case

In the remaining, generic case, the proof is very similar to the above. It continues with the following lemma, a convenient application of the Blackwell renewal theorem (see [4]).

Lemma 5 (Blackwell renewal theorem). *Suppose that $\{\alpha_i\}_{i \in \mathcal{I}}$ is not rank one and that $H = -\sum_i \alpha_i \log(\alpha_i) < \infty$. Then one has the continuous limit*

$$Z(t) \rightarrow \frac{1}{H} \quad \text{as } t \rightarrow \infty,$$

where $Z(\cdot)$ is defined in Lemma 3.

This gives the following corollary. The proof is similar to that of the preceding, with the convergence now taking place as $\lambda \rightarrow 0^+$.

Corollary. For $\{\alpha_i\}_{i \in \mathcal{I}}$ as in the previous lemma, the following hold as $\lambda \rightarrow 0^+$.

- $|A_\lambda| \sim 1/\lambda H$,
- For all $\mathbf{v} \in W(\mathcal{I})$, $\mu_\lambda(T_{\mathbf{v}}[0, 1]) \rightarrow \alpha_{\mathbf{v}} = \text{Leb}(T_{\mathbf{v}}[0, 1])$.

The conclusion of Theorem 1 in this case is the same as given in the preceding case, see p.11.

Remark. The same method of proof gives a more general result. Suppose we have a set \mathbb{X} equipped with a probability measure ν , and there is a collection $\{T_i, \alpha_i\}_{i \in \mathcal{I}}$ such that

- a) $T_i : \mathbb{X} \rightarrow \mathbb{X}$ are injective ν -measurable functions;
- b) for all $i, j \in \mathcal{I}$ distinct, $T_i(\mathbb{X}) \cap T_j(\mathbb{X}) = \emptyset$;
- c) for all $\mathbf{v} \in W(\mathcal{I})$, $\nu(T_{\mathbf{v}}(\mathbb{X})) = \alpha_{\mathbf{v}}$;
- d) $\alpha_i > 0$ for all $i \in \mathcal{I}$ and $\sum_{i \in \mathcal{I}} \alpha_i = 1$; and
- e) $-\sum_{i \in \mathcal{I}} \alpha_i \log(\alpha_i) < \infty$.

Then, for any $x \in \mathbb{X} \setminus \bigcup_{i \in \mathcal{I}} T_i(\mathbb{X})$, and for any ν -measurable set S which can be written as

$$S = \bigcup_{\mathbf{v} \in V} T_{\mathbf{v}}(\mathbb{X})$$

(with V is any subset of $w(\mathcal{I})$), we have, as $\lambda \rightarrow 0^+$,

$$\frac{|S \cap X_\lambda(x)|}{|X_\lambda(x)|} \rightarrow \nu(S),$$

where $X_\lambda(x) = \{T_{\mathbf{v}}(x) : \mathbf{v} \in W(\mathcal{I}), \alpha_{\mathbf{v}} \geq \lambda\}$.

It is also possible to generalise the method of proof to include substitution schemes starting from arbitrary partitions, and one recovers results analogous to those of [2].

These will be discussed in greater depth in the forthcoming doctoral thesis of the second author.

4 Discrepancy estimates

In [18], the author both generalised the method used by Adler and Flatto in [1] to general finite partitions, and posed questions which inspired various other papers. Of particular interest is the behaviour of the *discrepancy*. This corresponds to estimating the speed of convergence (where the I are intervals) of

$$\sup_{I \subset [0,1]} |\mu_\lambda(I) - \|I\|| \rightarrow 0 \text{ as } \lambda \rightarrow 0^+.$$

The Koksma inequality uses such estimates to give convergence rates for integrals of bounded-variation functions. Hence, fast-decaying discrepancies provide potential computationally-efficient numerical integration techniques. The theory of discrepancies (of equidistributing sets and sequences) has naturally received a lot of attention (see, e.g., [9] for an overview). There are also interesting open problems, for example on optimal discrepancy in higher dimensions.

Regarding interval substitution schemes, in the finite-partition case, discrepancy estimates are provided by Drmota–Infusino in [8], extending the application-focused work of Carbone in [5]. We in turn extend this to the context of infinite partitions. The results obtained here are different, depending on whether we are in the rank one case or the higher rank case.

4.1 Discrepancy estimates in the rank one case

In this final section, we extend the analysis of the rank one case to estimate the discrepancy between the measure μ_λ and the Lebesgue measure. More precisely, we have the following result.

Theorem 2. *Suppose that*

1. $\{\alpha_i\}_{i \in \mathcal{I}}$ is rank one;
2. $x > 0$ is the smallest number for which $\{\alpha_i\}_{i \in \mathcal{I}} \subset \{x^n\}_{n \in \mathbb{N}}$; and
3. there is some $\varepsilon > 0$ for which $\sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon} < \infty$.

Then there is an $R^* \in (0, 1)$, made explicit in Lemma 6 below, such that, for all $\rho \in (x/R^*, 1)$, there is a constant $C > 0$ such that, for all $n \in \mathbb{N}$ and all intervals $I \subset [0, 1]$,

$$|\mu_{x^n}(I) - |I|| \leq C\rho^n.$$

The proof of Theorem 2 begins with the following light lemma, defining R^* in terms of a generating function for $|A_{x^n}|$.

Lemma 6. *Given $\{\alpha_i\}_{i \in \mathcal{I}}$ as in Theorem 2, the function formally defined by*

$$g(z) = (z - x) \sum_{n=0}^{\infty} |A_{x^n}| z^n$$

has a holomorphic extension to the open disk of radius R^* about 0, where

$$R^* := \min \left(\{x^{1-\varepsilon}\} \cup \{|z| : z \in \mathbb{C} \setminus \{x\}, \sum_{j \in \mathcal{I}} z^{n_j} = 1\} \right) > x$$

and where $n_i := \log_x(\alpha_i)$. Therefore, denoting by b_n the n th Taylor coefficient of g , given $R < R^*$, there exists $C > 0$ such that, for all $n \in \mathbb{N}$, $b_n \leq CR^{-n}$.

Proof of Lemma 6. From (4), the renewal equation of Lemma 3, one has, for $|z| \leq x^{1-\varepsilon}$ and $z \neq x$,

$$\begin{aligned}
\frac{g(z)}{z-x} &= \sum_{n=0}^{\infty} |A_{x^n}| z^n = \sum_{n=0}^{\infty} \left(\sum_{j \in \mathcal{I}} |A_{x^{n-n_j}}| + 1 \right) z^n \\
&= \sum_{n=0}^{\infty} \sum_{j \in \mathcal{I}} |A_{x^{n-n_j}}| z^n + \frac{1}{1-z} \\
&= \sum_{j \in \mathcal{I}} z^{n_j} \sum_{n=0}^{\infty} |A_{x^{n-n_j}}| z^{n-n_j} + \frac{1}{1-z} \\
&= \sum_{j \in \mathcal{I}} z^{n_j} \frac{g(z)}{z-x} + \frac{1}{1-z},
\end{aligned}$$

which rearranges to

$$g(z) = \frac{z-x}{(z-1)(\sum_{j \in \mathcal{I}} z^{n_j} - 1)}.$$

Therefore, g has a meromorphic expansion on the disk of convergence of

$$z \mapsto \sum_{j \in \mathcal{I}} z^{n_j}; \tag{6}$$

which has radius at least $x^{1-\varepsilon}$, by assumption on the decay of $\{\alpha_i\}$.

An elementary argument (see [11, pp.201–2]) shows that $z = x$ is not only a simple root of (6) with residue $1/H$ (see the statement of Lemmas 4 or 5), but it is also the only root of (6) in the closed disk $\{|z| \leq x\}$. Therefore, g is holomorphic on the open disk of radius $R^* > x$, where R^* is the absolute value of the next smallest root of (6), or equal $x^{1-\varepsilon}$ if no other root exists. \square

Example 10. In certain nice cases, one can say more. For the simplest infinite example, $\alpha_n = 2^{-n}$ ($n \in \mathbb{N}$), $g(z) \equiv 1/2$ is constant.

The final stage of the proofs of Theorems 2 and 3 is similar to that of Theorem 1, but the method of splitting up a general interval needs a little more care. We will only prove the case that none of the T_i fix 0, since the other case is similar but tedious.

Proof of Theorem 2. For simplicity, consider the interval $I = [b, 1)$, for fixed $b \in (0, 1)$, and assume that no $i \in \mathcal{I}$ satisfies $T_i(0) = 0$. For $n \in \mathbb{N}$, let V_n denote the elements of U_n (where U_n is in the proof of Theorem 1) whose interval is not contained in one from $U_{n'}$ for any $n' < n$. More explicitly,

$$V_1 := \{i \in \mathcal{I} : T_i[0, 1) \subset I\} = U_1,$$

and, for $n \geq 2$,

$$\begin{aligned}
V_n &:= \{\mathbf{v} * i \in \mathcal{I}^n : i \in \mathcal{I}, T_{\mathbf{v}*i}[0, 1) \subset I \text{ but } T_{\mathbf{v}}[0, 1) \not\subset I\} \\
&= U_n \setminus U_{n-1}.
\end{aligned}$$

It is simple to show the union over all intervals coming from the V_n ,

$$\bigcup_{n=1}^{\infty} \bigcup_{v \in V_n} T_v[0, 1),$$

is disjoint, contained in I , and differs from I by at most some exceptional set of points—those contained in at most finitely many $T_v[0, 1)$, i.e., a subset of

$$K = [0, 1] \setminus \bigcup_{N \in \mathbb{N}} \bigcap_{n > N} \bigcup_{v \in \mathcal{I}^n} T_v[0, 1).$$

This K , similarly to K_n in the proof of Theorem 1, has both μ_λ and Lebesgue measure zero.

To say more, it is necessary to give a partial description for V_n , involving the itinerary for b .

Definition (Itinerary, $\text{It}(x)$, $\text{It}_n(x)$). This definition has two cases. First suppose that $x \in [0, 1]$ is such that

- x is a left endpoint, $T_v(0)$, for some $v \in W(\mathcal{I})$; or
- x lies in the exceptional set K above.

Then

$$n = \min \left\{ k \in \mathbb{N} : x \notin \bigcup_{v \in \mathcal{I}^k} T_v[0, 1) \right\}$$

exists; and we call $\text{It}(x)$, the *itinerary* of x , the unique word in \mathcal{I}^{n-1} such that

$$x \in T_{\text{It}(x)}[0, 1).$$

and we say x has *finite itinerary*.

Otherwise, we say that x has *infinite itinerary*, and the itinerary $\text{It}(x)$ is the sequence

$$\text{It}(x) := (i_n)_n \in \mathcal{I}^{\mathbb{N}}$$

such that, for each truncation $\text{It}_n(x) := (i_1, \dots, i_n) \in \mathcal{I}^n$, $x \in T_{\text{It}_n(x)}[0, 1)$.

Returning to the proof: From a similar argument to that in Theorem 1, if $\text{It}(b) = (i_1, i_2, \dots, i_n) \in \mathcal{I}^n$ has finite itinerary, then V_k is empty for all $k \geq n+2$, and also for all $k \leq n$, we have

$$\begin{aligned} I_k &= \{(i_1, i_2, \dots, i_{k-1}, i) \in \mathcal{I}^k \mid T_{i_k}(0) < T_i(0)\} \\ &\subset \{\text{It}_k(b)\} \times \mathcal{I} \quad \text{and similarly,} \\ I_{n+1} &\subset \{\text{It}(b)\} \times \mathcal{I}. \end{aligned} \tag{7}$$

Otherwise, if $\text{It}(b) = (i_n)_{n=1}^{\infty}$ is infinite, (7) holds for all $k \in \mathbb{N}$ (it is as if $n = \infty$).

Now let $V = \bigcup_{k \in \mathbb{N}} V_k$ and $n \in \mathbb{N}$. It follows from the nullity of K that we can write the following, and we divide the sum into two, corresponding to intervals which have or haven't been split at this n th stage:

$$\begin{aligned} \mu_{x^n}(I) - \|I\| &= \sum_{\mathbf{v} \in V} \mu_{x^n}(T_{\mathbf{v}}[0, 1]) - \alpha_{\mathbf{v}} \\ &= \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} \geq x^n}} \frac{|X_{x^n/\alpha_{\mathbf{v}}}| - \alpha_{\mathbf{v}}|X_{x^n}|}{|X_{x^n}|} - \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} < x^n}} \alpha_{\mathbf{v}}. \end{aligned} \quad (8)$$

We estimate the second sum first, corresponding to intervals which have not yet been split up to this value of n . Firstly, we have

$$0 \leq \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} < x^n}} \alpha_{\mathbf{v}} \leq x^{n(1-\varepsilon)} \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} < x^n}} \alpha_{\mathbf{v}}^{1-\varepsilon},$$

and the sum on the right hand side can be bounded uniformly in b , as follows. Recall that, in the infinite itinerary case, the inclusion in equation (7) holds for all $n \in \mathbb{N}$, and we may write

$$\begin{aligned} \sum_{\mathbf{v} \in V} \alpha_{\mathbf{v}}^{1-\varepsilon} &= \sum_{k=1}^{\infty} \sum_{\mathbf{v} \in V_k} \alpha_{\mathbf{v}}^{1-\varepsilon} \leq \sum_{k=1}^{\infty} \alpha_{\text{It}_{n-1}(x)}^{1-\varepsilon} \sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon} \\ &\leq \sum_{n=1}^{\infty} (\alpha_*)^{1-\varepsilon} \sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon} \quad (\alpha_* = \max_{i \in \mathcal{I}} \{\alpha_i\}) \\ &= \frac{(\alpha_*)^{1-\varepsilon}}{1 - (\alpha_*)^{1-\varepsilon}} \sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon} =: C \end{aligned}$$

which is finite by assumption. The finite itinerary case is similar and we obtain the same bound, C : the only difference is that sum in k is finite.

Therefore, the second sum of (8) is bounded above by $Cx^{n(1-\varepsilon)}$, and we may turn our attention to the first.

Consider the term of this first sum corresponding to $\mathbf{v} \in V$. Write $\alpha_{\mathbf{v}} = x^m$ for some $m \in \mathbb{N}$, and consider a generating series for (the numerator of) the corresponding summand, which relates to the g from Lemma 6:

$$\begin{aligned} \sum_{n=1}^{\infty} (|X_{x^{n-m}}| - x^m|X_{x^n}|)z^n &= (z^m - x^m) \left(\frac{g(z)}{z-x} \right) \\ &= (z^{m-1} + xz^{m-2} + \cdots + x^{m-1})g(z). \end{aligned}$$

Recalling b_n as the n th Taylor coefficient of g , equating coefficients on both sides gives, for all $m \geq n$,

$$|X_{x^{n-m}}| - x^m|X_{x^n}| = b_{n-m+1} + xb_{n-m+2} + \cdots + x^{m-1}b_n.$$

Thus, by Lemma 6, given $\rho \in (x/R^*, 1)$, there is some (possibly updated) constant $C > 0$ such that, for all $n \in \mathbb{N}$, $b_n < C(\rho/x)^n$. Applying this to the previous equation gives, for all $n \geq m$,

$$|X_{x^{n-m}}| - x^m |X_{x^n}| \leq C \frac{\rho^{n-m-1} + \rho^{n-m-2} + \dots + \rho^n}{x^{n-m+1}} = C \left(\frac{\rho}{x}\right)^n \frac{(x/\rho)^m - x^m}{x(1-\rho)},$$

and dividing both sides through by $|X_{x^n}| = \mathcal{O}(x^{-n})$ gives

$$\left| \frac{|X_{x^n/\alpha_v}| - \alpha_v |X_{x^n}|}{|X_{x^n}|} \right| \leq \hat{C} \rho^n \left(\left(\frac{x}{\rho}\right)^m - x^m \right) < \hat{C} \rho^n (\alpha_v^{1-\varepsilon} - \alpha_v),$$

where we have used $x/\rho < (R^*)^{-1} \leq x^{1-\varepsilon}$ in the last inequality, for some constant $\hat{C} = \hat{C}(x, \rho) > 0$. Summing over $v \in V$ bounds the first sum of (8):

$$\sum_{\substack{v \in V \\ \alpha_v \geq x^n}} \left| \frac{|X_{x^n/\alpha_v}| - \alpha_v |X_{x^n}|}{|X_{x^n}|} \right| \leq \left(\hat{C} \sum_{\substack{v \in V \\ \alpha_v \geq x^n}} \alpha_v^{1-\varepsilon} - \alpha_v \right) \rho^n \leq \hat{C}(C+1)\rho^n.$$

Since $n \in \mathbb{N}$ was arbitrary, we have the required estimate. \square

4.2 Discrepancy estimates in the higher rank case

When the collection $\{\alpha_j\}_{j \in \mathcal{I}}$ is not rank one, we require not only a strict decay property on the $\{\alpha_j\}$ —as in Theorem 2—, but also a kind of Diophantine condition. To proceed, we need the following definition.

Definition (R -badly approximable). For $R \in [2, \infty)$ we say a number $\gamma \in \mathbb{R}$ is R -badly approximable if there exists a $d > 0$ such that

$$\forall (l, k) \in \mathbb{Z}^2 \text{ s.t. } l \neq 0, \quad \left| \gamma - \frac{k}{l} \right| > \frac{d}{|l|^R}.$$

Remark. Larger values of R correspond to more easily approximable numbers:

- For any $R = 2$, the property is equivalent to γ having bounded continued fraction coefficients. (Such γ comprise a set of measure 0 containing all quadratic algebraic numbers.)
- For $R > 2$, the property holds Lebesgue almost-everywhere: by Jarnik's theorem [12, Thm. 10.3], the Hausdorff dimension of the complementary set is $2/R < 1$.
- Certain transcendental numbers (e.g., *Liouville numbers*) do not satisfy the property for any R whatsoever.

Theorem 3. *Suppose that $\{\alpha_i\}_{i \in \mathcal{I}}$ is not rank one, that there is some $\varepsilon > 0$ such that $\sum_i \alpha_i^{1-\varepsilon} < \infty$, and that there is a pair $\alpha_j, \alpha_k \in \{\alpha_i\}_{i \in \mathcal{I}}$ such that*

$\log(\alpha_j)/\log(\alpha_k)$ is $(2+r)$ -badly approximable, for some $r \in [0, 1/2)$. Then, for all $P \in (0, P^*)$, there exists a constant C such that, for all intervals $I \subset [0, 1]$,

$$|\mu_\lambda(I) - \|I\|| \leq C(-\log(\lambda))^{-P};$$

where

$$P^* = \frac{1-2r}{8(1+r)}.$$

The proof of Theorem 3 requires us to consider the Mellin transform,

$$g(z) = \int_0^\infty t^{-z-1} |A_{1/t}| dt,$$

which has the following explicit form, courtesy of the renewal equation for $|A_\lambda|$. Let \Re, \Im denote the real and imaginary parts of a complex number, respectively.

Lemma 7. For $\Re(z) > 1$, the Mellin transform $g(z)$ takes the form

$$g(z) = \frac{1}{z(\sum_{j \in \mathcal{I}} \alpha_j^z - 1)}.$$

In particular, if $\sum_j \alpha_j^{1-\varepsilon} < \infty$, g has a meromorphic extension to the half-plane $\{\Re(z) > 1 - \varepsilon\}$.

Proof of Lemma 7. The formula follows from standard properties of the Mellin transform, namely $g(\alpha_j z) = \alpha_j^z g(z)$, and the renewal equation of Lemma 3. \square

As is well-known, one can obtain asymptotic information about a function from the distribution of poles of its Mellin transform.

In particular, this depends on the zeros of the almost-periodic function

$$f(z) = \sum_{j \in \mathcal{I}} \alpha_j^z - 1.$$

There is a lot we can say straight away. From our assumptions on $\{\alpha_j\}_j$, we have that $z = 1$ is a simple zero of f , and by the triangle inequality, there are no zeros of f for $\{\Re(z) > 1\}$. Moreover, if there were another zero of f on the line $\{\Re(z) = 1\}$, it would follow that $\{\alpha_j\}_j$ is rank one, a contradiction.

The role of the poles of the Mellin transform is illustrated in the proof of the following result.

Proposition 1. For any given collection of positive numbers $\{\alpha_j\}_{j \in \mathcal{I}}$ which sum to 1 such that $H = -\sum_j \alpha_j \log(\alpha_j) < \infty$, there is no $\varepsilon > 0$ for which

$$|A_\lambda| = \frac{1}{H\lambda} + \mathcal{O}(\lambda^{1-\varepsilon}) \tag{9}$$

as $\lambda \rightarrow 0^+$.

Proof of Proposition 1. Fix a collection $\{\alpha_j\}_{j \in \mathcal{I}}$ as above and assume for contradiction that there is an $\varepsilon > 0$ for which (9) holds. Taking the Mellin transform of this equation yields that

$$\frac{1}{zf(z)} = \frac{1}{H(z-1)} + \int_1^\infty t^{-z-\varepsilon} \mathcal{O}(1) dt.$$

Since this second integral converges absolutely for all z with $\Re(z) > 1 - \varepsilon$, we see that the left hand side has a meromorphic extension to the half plane $\{z \in \mathbb{C} : \Re(z) > 1 - \varepsilon\}$, with only one pole at $z = 1$. To obtain a contradiction, we provide a sequence of zeros $z_n = u_n + iv_n$ with $u_n \rightarrow 1$, $v_n \rightarrow \infty$, using the theory of almost-periodic functions. In particular, we follow the proof of the corollary to [7, Theorem 3.6].

It can quickly be seen that f is almost-periodic, by [7, Cor. to Thm. 3.12], and is bounded on $\{\Re(z) \geq 1 - \varepsilon/2\}$ since, in this half-plane,

$$|f(z)| \leq 1 + \sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon/2}.$$

Also, since $f(1) = 0$, the definition of almost-periodicity provides a sequence of positive numbers $(y_n)_{n=1}^\infty$ for which $f(1 + iy_n) \rightarrow 0$ as $n \rightarrow \infty$. In other words, the holomorphic functions

$$f_n(z) = f(z + iy_n)$$

are bounded on the same half-plane, and $f_n(1) \rightarrow 0$.

Furthermore, fixing any index $j \in \mathcal{I}$, for $L := -2\pi/\log(\alpha_j)$, we see that $\sup_{|v| \leq L} |f_n(1 + iv)|$ is bounded away from zero uniformly, since every interval of length L contains a v such that

$$\alpha_j^{1+iv} = -\alpha_j \quad \therefore \quad |f(1 + iv)| \geq 1 + \alpha_j - \sum_{k \in \mathcal{I} \setminus \{j\}} \alpha_k = 2\alpha_j.$$

Therefore, on the rectangle $(1-\varepsilon, 1+\varepsilon) + i(-l, l)$, an application of Montel's theorem shows that f_n uniformly converges (passing to a subsequence if necessary) to some analytic function f_∞ . From the last two considerations, $f_\infty(1) = 0$ and f_∞ is not identically zero.

Now, taking a circle about 1 small enough that f_∞ has no zeros on it, by Hurwitz's theorem, for all n sufficiently large, f_n has a zero \hat{z}_n inside this circle, such that $\hat{z}_n \rightarrow 1$ as $n \rightarrow \infty$.

This thus provides us with a sequence of zeros of f accumulating on the line as required, contradicting the statement that f has only one zero in the above-mentioned half-plane. \square

Considering now the proof of Theorem 3, the following lemma uses the $(2+r)$ -badly approximable hypothesis, following [8]. For simplicity in the following two proofs, we write $\alpha = \max(\alpha_j, \alpha_k)$ and $\beta = \min(\alpha_j, \alpha_k)$.

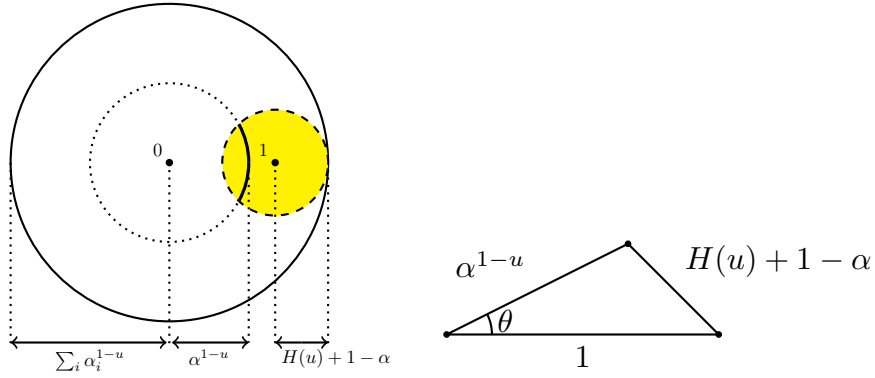


Figure 4: (i) The region in which α^z must lie for $f(z) = 0$ —the bold arc in the shaded circle; (ii) the triangle defining $\theta(u)$, the maximum possible value of $|\eta_\alpha|$.

Lemma 8. *Suppose that $(\alpha_j)_{j \in \mathcal{I}}$ satisfies the conditions of Theorem 3. Then there exists $C > 0$ such that, whenever $z = 1 - u + iv \in \mathbb{C} \setminus \{1\}$ satisfies both $f(z) = 0$ and $u < \varepsilon$, then $u > 0$ and*

$$|v|^{2+2r} \geq \frac{C}{u}.$$

Proof of Lemma 8. The fact that $u > 0$ follows from the discussion preceding the previous proposition. We first show, if $f(z) = 0$, the argument of α^z in $(-\pi, \pi]$ is $\mathcal{O}(\sqrt{u})$ as $u \rightarrow 0^+$. That is, the quantity $\eta_\alpha \in (-\pi, \pi]$, satisfying

$$e^{i\eta_\alpha} = \frac{\alpha^z}{|\alpha^z|} = \frac{\alpha^z}{\alpha^{1-u}},$$

is $\mathcal{O}(\sqrt{u})$. This uses the triangle inequality and a small amount of trigonometry, as we now detail. We have that

$$|\alpha^z - 1| \leq 1 - \alpha + H(u),$$

where

$$H(u) := \sum_{n \in \mathcal{I}} \alpha_n^{1-u} - \alpha^{1-u} - 1 + \alpha.$$

In particular, $H(u) = \mathcal{O}(u)$ as $u \rightarrow 0^+$, by the mean value theorem. Therefore, for u sufficiently small, $H(u) < \alpha$, which gives rise to the picture in Figure 4(i). Consequently, $|\eta_\alpha| < \theta$, where θ is as in 4(ii) and satisfies the following equation.

$$\cos(\theta) = \frac{1 + \alpha^{2-2u} - (1 - \alpha + H(u))^2}{2\alpha^{1-u}} = 1 - \mathcal{O}(u) \quad (u \rightarrow 0),$$

where the constant now depends on the value of α . Using, for example, that

$$\lim_{y \rightarrow 0^+} \frac{\arccos(1-y)}{\sqrt{y}} = \sqrt{2}, \quad \arccos : [-1, 1] \rightarrow [0, \pi],$$

it is clear that $\theta = \mathcal{O}(\sqrt{u})$ as $u \rightarrow 0^+$, hence the same applies to η_α .

We can repeat this argument with β in place of α to bound the analogously defined η_β —i.e., $\eta_\beta = \mathcal{O}(\sqrt{u})$.

Write $v \log(\alpha) = 2\pi k + \eta_\alpha$ and $v \log(\beta) = 2\pi l + \eta_\beta$, supposing $|v| \geq 2\pi / \log(\beta)$ so that k and l are non-zero. Substituting into the definition of $(2+r)$ -badly approximable gives

$$\frac{d}{|l|^{2+r}} \leq \left| \frac{\log(\alpha)}{\log(\beta)} - \frac{k}{l} \right| = \left| \frac{2\pi k + \eta_\alpha}{2\pi l + \eta_\beta} - \frac{k}{l} \right| = \left| \frac{\eta_\alpha}{2\pi l + \eta_\beta} - \eta_\beta \frac{2\pi k + \eta_\alpha}{(2\pi l + \lambda)^2} \right|,$$

where $\lambda \in \mathbb{R}$ is a constant, provided by the mean value theorem, satisfying $0 < |\lambda| < |\eta_\beta| \leq \pi$. Using the triangle inequality on the right hand side, multiplying through by $(2\pi l + \lambda)^2 |l|^r$ and using that $|\eta_\alpha| \leq \pi$, one obtains

$$\begin{aligned} \pi^2 d &\leq 4\pi^2 \left(1 - \frac{\lambda}{l}\right)^2 d \leq 2\pi \left(|\eta_\alpha| \frac{(1 - \lambda/2\pi l)^2}{|1 + \eta_\beta/2\pi l|} + |\eta_\beta| \left| \frac{k}{l} - \frac{\eta_\alpha}{2\pi} \right| \right) |l|^{1+r} \\ &\leq 2\pi \left(9|\eta_\alpha| + |\eta_\beta| \left(\left| \frac{k}{l} \right| + \frac{1}{2} \right) \right) |l|^{1+r} \\ &\leq 2\pi \left(9|\eta_\alpha| + |\eta_\beta| \left(\frac{3}{2} + 2 \frac{\log(\alpha)}{\log(\beta)} \right) \right) |l|^{1+r}, \end{aligned}$$

where the last inequality uses that $|v| > -2\pi / \log(\beta)$. We can divide through by the large bracket on the right hand side and, recalling the \sqrt{u} asymptotic for η_α and η_β , obtain the required inequality for some constant C , for u sufficiently small and $|v|$ sufficiently large.

That the inequality holds in the whole of the specified region (with a possibly different C) follows simply from the fact that zeros of f can only accumulate on the vertical boundary $\{z = 1 - \varepsilon + iv : v \neq 0\}$, and so there is an open neighbourhood of $[1 - \varepsilon, 1]$ containing only one zero of f , at 1. Since there are finitely many zeros of f to cater for (at most), we can adapt C accordingly. \square

The next lemma is a variant on the last and allows us to estimate decay of the Mellin inverse integral inside the zero-free region.

Lemma 9. *Suppose $(\alpha_j)_{j \in \mathcal{I}}$ is a collection of positive numbers as given in Theorem 3. Then there exists $C > 0$ such that, whenever $z = 1 - u + iv \in \mathbb{C}$ with $u \geq 0$ and $\sigma > 0$ sufficiently small,*

$$|f(z)| < \sigma \tag{10}$$

implies one of the following holds: either

$$|v| \leq 2\pi / \log(\beta)$$

or

$$|v|^{2+2r} > \frac{C}{\max(u, \sigma)}.$$

Proof of Lemma 9. The proof is an adaptation of that for Lemma 8. This time, the radius of the circle depicted in Figure 4 is $H(u) + 1 - \alpha + \sigma$ and correspondingly,

$$\cos(\theta) = \frac{1}{2} \frac{1 + \alpha^{2-2u} - (1 - \alpha + H(u) + \sigma)^2}{\alpha^{1-u}} = 1 - \mathcal{O}(\max(u, \sigma))$$

as $\max(u, \sigma) \rightarrow 0$, which gives, for $\max(u, \sigma)$ is sufficiently small),

$$|\theta| \leq \frac{\pi}{2} \sqrt{1 - \cos(\theta)} = \mathcal{O}(\sqrt{\max(u, \sigma)})$$

and the proof continues as before. \square

The next crucial lemma is the analogue to Lemma 6 in the higher rank case.

Lemma 10. *Under the assumptions of Theorem 3,*

$$|A_\lambda| = \frac{1}{H\lambda} + \mathcal{O}(\lambda^{-1}(-\log(\lambda))^{-P}) \quad \lambda \in (0, 1],$$

where $P \in (0, P^*)$ is as given in Theorem 3.

Proof of Lemma 10. The proof uses simple complex analysis to estimate the integral

$$F(t) := \frac{1}{2\pi i} \int_{2-i\infty}^{2+i\infty} \frac{t^{z+3}}{z(z+1)(z+2)(z+3)f(z)} dz, \quad \text{for } t > 1. \quad (11)$$

We first relate $F(t)$ to $|A_{1/t}|$. On the line $\{\Re(z) = 2\}$, we see that f is uniformly bounded away from zero,

$$|f(z)| \geq 1 - \sum_{j \in \mathcal{I}} \alpha_j^2 > 0,$$

so $F(t)$ absolutely converges, for all t . Therefore, by the Mellin inversion theorem, the Mellin transform of $t \mapsto F(t)/t^3$ is the denominator of the integrand of F :

$$F^*(z) := \int_0^\infty t^{-z-1} F(t) dt = \frac{1}{z(z+1)(z+2)(z+3)f(z)}.$$

Therefore, using integration by parts, one has $F^{(3)}(t) = |A_{1/t}|$ Lebesgue almost-everywhere.

We now relate the integral in (11) to that over the contour Γ , parametrised by

$$\gamma : \mathbb{R} \rightarrow \mathbb{C}, \quad \gamma(v) = 1 + iv - D \min(1, |v|^{-2-r})$$

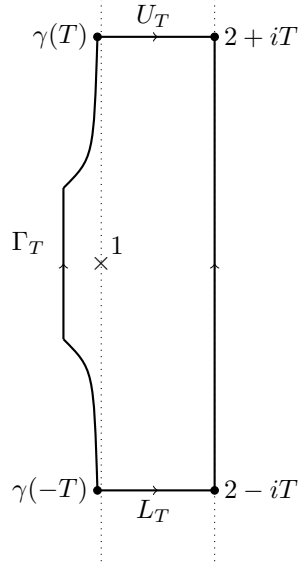


Figure 5: The contour Γ_t used in the proof of Lemma 10.

(see Figure 5), where $D > 0$ is chosen sufficiently small so that the previous two lemmas apply as follows: firstly, the only zero of f which lies on or to the right of Γ is at 1, and secondly, whenever z lies on or to the right of Γ and $|\Im(z)| \geq -2\pi(\log(\beta))^{-1}$, one has

$$|f(z)| \geq D|\Im(z)|^{-2-2r}. \quad (12)$$

Consider, for $T \geq 1$, the contour Γ_T depicted in Figure 5. Cauchy's residue theorem gives

$$\int_{2-iT}^{2+iT} F^*(z) dz = 2\pi i \frac{t^4}{24H} + \int_{\Gamma_T} F^*(z) dz + \int_{U_T} F^*(z) dz + \int_{L_T} F^*(z) dz.$$

Since U_T and L_T have bounded length, a simple application of (12) shows that the corresponding last two integrals are $\mathcal{O}(|T|^{-2+2r})$ as $T \rightarrow \infty$; therefore, taking the limit, one has

$$F(t) = \frac{t^4}{24H} + \frac{1}{2\pi i} \int_{\Gamma} F^*(z) t^{-z-3} dz. \quad (13)$$

From this point, the proof follows along the lines of [10, Theorem 4.6, pp.133–4]. Since $t \mapsto |A_{1/t}|$ is non-decreasing on the real line, so too are the functions F , F' and F'' . Using this, together with repeated applications of the mean value

theorem, gives that, for any $t \in \mathbb{R}$ and $h > 0$, both

$$\begin{aligned} \frac{F(t-3h) - 3F(t-2h) + 3F(t-h) - F(t)}{-h^3} &\leq |A_{1/t}| \quad \text{and} \\ |A_{1/t}| &\leq \frac{F(t+3h) - 3F(t+2h) + 3F(t+h) - F(t)}{h^3} \quad \text{hold.} \end{aligned}$$

Substituting (13) into these expressions yields

$$\begin{aligned} &\frac{F(t \pm 3h) - 3F(t \pm 2h) + 3F(t \pm h) - F(t)}{\pm h^3} = \\ &t \pm \frac{3}{2}h + \frac{1}{2\pi i} \int_{\Gamma} \frac{(t \pm 3h)^{z+3} - 3(t \pm 2h)^{z+3} + 3(t \pm h)^{z+3} - t^{z+3}}{\pm h^3} F^*(z) dz. \end{aligned}$$

Thus, we have the estimate

$$\begin{aligned} |A_{1/t}| &= t + \mathcal{O}(h) + \\ &\mathcal{O} \left(\int_{\Gamma} \left| \frac{(t \pm 3h)^{z+3} - 3(t \pm 2h)^{z+3} + 3(t \pm h)^{z+3} - t^{z+3}}{\pm h^3(z+1)(z+2)(z+3)} \right| \left| \frac{1}{zf(z)} \right| |dz| \right). \end{aligned}$$

From now on, let $h = h(t) \in (0, t)$ be a function of t to be determined later. To begin to estimate the integral, consider $|\Delta_{\pm}(t, h, z)|$ for $t > 1$ and $z \in \Gamma$, where

$$\Delta_{\pm}(t, h, z) := \frac{(t \pm 3h)^z - 3(t \pm 2h)^z + 3(t \pm h)^z - t^z}{\pm h^3(z+1)(z+2)(z+3)}.$$

We estimate $|\Delta_{\pm}|$ in two different ways. For the first, we express $|\Delta_{\pm}|$ as a series of nested integrals:

$$\begin{aligned} |\Delta_{\pm}(t, h, z)| &= \left| \frac{1}{\pm h} \int_t^{t \pm h} \frac{1}{\pm h} \int_{\hat{t}}^{t \pm h} \frac{1}{\pm h} \int_{\hat{\hat{t}}}^{\hat{\hat{t}} \pm h} \hat{\hat{t}}^z \, d\hat{\hat{t}} \, d\hat{t} \, dt \right| \leq (t+3h)^{\Re(z)} \\ &\leq 4t^{\Re(z)}. \end{aligned}$$

For the first inequality, we have taken the absolute value signs inside the integral and applied the mean value theorem three times, noting that $t \mapsto |t^z| = t^{\Re(z)}$ is increasing on $[0, \infty)$, since $\Re(z) > 0$ for $z \in \Gamma$. (The second follows simply from $h < t$ and $\Re(z) < 1$.)

The triangle inequality gives us another estimate: for $t \geq 1$,

$$\begin{aligned} |\Delta_{\pm}(t, h, z)| &\leq \frac{(4t)^{\Re(z)+3} + 3(3t)^{\Re(z)+3} + 3(2t)^{\Re(z)+3} + t^{\Re(z)+3}}{h^3|z+1||z+2||z+3|} \\ &\leq \frac{548t^{\Re(z)+3}h^{-3}}{|z+1||z+2||z+3|}. \end{aligned}$$

From the Bernoulli inequality it follows that, for $z = \gamma(v) \in \Gamma$, the three quantities $|\gamma(v)+1|$, $|\gamma(v)+2|$, $|\gamma(v)+3|$ are all greater than or equal to $(1+|v|)/2$, so altogether we have the following.

$$|\Delta_{\pm}(t, h, \gamma(v))| \leq \min \left(4t^{\Re(\gamma(v))}, \frac{4384 t^{\Re(\gamma(v))+3}}{h^3 (1+|v|)^3} \right).$$

From (12), one can easily deduce that $|\gamma(v)f(\gamma(v))|^{-1} = \mathcal{O}((1+|v|)^{1+2r})$, for all $v \in \mathbb{R}$. Combining the previous three inequalities and using the boundedness of $|\gamma'(v)|$ gives the following.

$$|A_{1/t}| = \frac{t}{H} + \mathcal{O}(h) + \mathcal{O}\left(\int_{-\infty}^{\infty} (1+|v|)^{1+2r} \min\left(t^{\Re(\gamma(v))}, \frac{t^{\Re(\gamma(v))+3}}{(1+|v|)^3}\right) dv\right).$$

Because the integral is symmetric in v , it suffices to estimate the integral from 0 to ∞ , as we now do. Writing $\Re(\gamma(v)) = 1 - \kappa(v)$, where $\kappa(v) = D \min(1, |v|^{-2-2r})$, the previous equation simplifies to the following.

$$\begin{aligned} \frac{|A_{1/t}|}{t} - \frac{1}{H} &= \mathcal{O}\left(\frac{h}{t}\right) + \mathcal{O}\left(\int_0^{\infty} (1+v)^{1+2r} t^{-\kappa(v)} \min\left(1, \frac{(t/h)^3}{(1+v)^3}\right) dv\right) \\ &= \mathcal{O}\left(\frac{h}{t}\right) + \mathcal{O}\left(\int_0^{\infty} (1+v)^{-2+2r} t^{-\kappa(v)} \min\left((1+v)^3, \left(\frac{t}{h}\right)^3\right) dv\right) \\ &= \mathcal{O}\left(\frac{h}{t}\right) + \mathcal{O}\left(\int_1^{\infty} v^{-2+2r} t^{-\kappa(v-1)} \min\left(v^3, \left(\frac{t}{h}\right)^3\right) dv\right). \end{aligned}$$

Now let $\delta \in (0, 1 - 2r)$. For $v, t \geq 1$, we have, since κ is decreasing on $[0, \infty)$,

$$\begin{aligned} v^{-2+2r} t^{-\kappa(v-1)} &\leq v^{-2+2r} t^{-\kappa(v)} \\ &= v^{-2+2r} \exp(-\kappa(v) \log(t)) \\ &= v^{-2+2r+\delta} \exp(-\kappa(v) \log(t) - \delta \log(v)) \\ &= v^{-2+2r+\delta} \exp(-Dv^{-2-2r} \log(t) - \delta \log(v)) \\ &= v^{-2+2r+\delta} \exp(-\xi_{\delta}(t)), \end{aligned}$$

where

$$\begin{aligned} \xi_{\delta}(t) &:= \inf_{v \geq 1} (Dv^{-2-2r} \log(t) + \delta \log(v)) \\ &= \frac{\delta}{2+2r} \left(1 + \log\left(\frac{D(2+2r) \log(t)}{\delta}\right)\right). \end{aligned}$$

This last equality holds for all t sufficiently large, by elementary calculus. Hence

$$\frac{|A_{1/t}|}{t} - \frac{1}{H} = \mathcal{O}\left(\frac{h}{t}\right) + \mathcal{O}\left(e^{-\xi_{\delta}(t)} \int_1^{\infty} v^{-2+2r+\delta} \min\left(v^3, \left(\frac{t}{h}\right)^3\right) dv\right).$$

Now, writing $\omega = vh/t$ and substituting, the integral becomes

$$\left(\frac{t}{h}\right)^{2+2r+\delta} \int_{h/t}^{\infty} \omega^{-2+2r+\delta} \min(\omega^3, 1) d\omega,$$

and this can be split into two parts,

$$\int_{h/t}^1 \omega^{1+2r+\delta} d\omega + \int_1^{\infty} \omega^{-2+2r+\delta} d\omega,$$

both of which are finite, since $2r + \delta \in (0, 1)$. Therefore, for all $t \geq 1$,

$$\frac{|A_{1/t}|}{t} - \frac{1}{H} = \mathcal{O}\left(\frac{h}{t}\right) + \mathcal{O}\left(e^{-\xi_\delta(t)} \left(\frac{t}{h}\right)^{2+2r+\delta}\right).$$

Finally, choosing $h(t) = t \exp(\frac{-\xi_\delta(t)}{3+2r+\delta})$ ensures both terms have the same order of magnitude, and the previous equation simplifies to the required expression:

$$|A_{1/t}| = \frac{t}{H} + \mathcal{O}(t \log(t)^{-P}),$$

where

$$P = P(r, \delta) = \frac{\delta}{(2+2r)(3+2r+\delta)} \in \left(0, \frac{1-2r}{8(1+r)}\right). \quad \square$$

The concluding stages of the proof of Theorem 3 are similar to those of Theorem 2; but with some notable differences.

Proof of Theorem 3. For simplicity of writing, we again consider only the case that none of the T_i fix 0. Let $\lambda \in (0, 1)$ and recall $I = [b, 1)$, V_k and V from the proof of Theorem 2 (p.14 onwards).

Similarly to that proof, we may write the discrepancy in terms of three sums,

$$\begin{aligned} \mu_\lambda(I) - \|I\| &= \sum_{\mathbf{v} \in V} \mu_\lambda(T_{\mathbf{v}}[0, 1)) - \alpha_{\mathbf{v}} \\ &= \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} \geq \lambda/\alpha_{\max}}} \frac{|X_{\lambda/\alpha_{\mathbf{v}}}|}{|X_\lambda|} - \alpha_{\mathbf{v}} \\ &\quad + \sum_{\substack{\mathbf{v} \in V \\ \lambda \leq \alpha_{\mathbf{v}} < \lambda/\alpha_{\max}}} \frac{|X_{\lambda/\alpha_{\mathbf{v}}}|}{|X_\lambda|} - \sum_{\substack{\mathbf{v} \in V \\ \alpha_{\mathbf{v}} < \lambda/\alpha_{\max}}} \alpha_{\mathbf{v}}, \end{aligned} \quad (14)$$

where $\alpha_{\max} = \max_{i \in \mathcal{I}}(\alpha_i)$.

The second and third sum of (14) decay much faster than the first, which can be deduced from the summability of $\alpha_{\mathbf{v}}^{1-\varepsilon}$ alone. For the third sum, the argument from Theorem 2 (p.16) applies to bound it by a multiple of $\lambda^{1-\varepsilon}$. The second sum of (14) corresponds to newly split intervals: if $\alpha_{\mathbf{v}} < \lambda/\alpha_{\max}$, then $\alpha_{\mathbf{v}}\alpha_i < \lambda$ for any $i \in \mathcal{I}$, so no proper subinterval of $T_{\mathbf{v}}[0, 1)$ has been split for

this value of λ , and $|X_{\lambda/\alpha_{\mathbf{v}}}| = 1$. Therefore, assuming b has an infinite itinerary,

$$\begin{aligned}
\sum_{\substack{v \in V \\ \lambda \leq \alpha_v < \lambda/\alpha_{\max}}} |X_{\lambda/\alpha_v}| &= |\{\mathbf{v} \in V : \lambda \leq \alpha_{\mathbf{v}} < \lambda/\alpha_{\max}\}| \\
&= \sum_{n \in \mathbb{N}} |\{\mathbf{v} \in V_n : \lambda \leq \alpha_{\mathbf{v}} < \lambda/\alpha_{\max}\}| \\
&\leq \sum_{n \in \mathbb{N}} |\{i \in \mathcal{I} : \lambda \leq \alpha_{\text{It}_{n-1}(b)} \alpha_i < \lambda/\alpha_{\max}\}| \\
&= \sum_{n \in \mathbb{N}} |\{i \in \mathcal{I} : \lambda/\alpha_{\text{It}_{n-1}(b)} \leq \alpha_i < \lambda/(\alpha_{\max} \alpha_{\text{It}_{n-1}(b)})\}| \\
&\leq \sum_{n \in \mathbb{N}} |\{i \in \mathcal{I} : \lambda/\alpha_{\text{It}_{n-1}(b)} \leq \alpha_i < \lambda/\alpha_{\text{It}_n(b)}\}| \\
&= |\{i \in \mathcal{I} : \lambda \leq \alpha_i\}| \\
&\leq \sum_{i \in \mathcal{I}} \alpha_i^{1-\varepsilon} \lambda^{\varepsilon-1}.
\end{aligned}$$

Thus, the second sum of (14) is bounded by a multiple of λ^ε . The finite itinerary case is similar, involving a finite sum in n .

It remains to bound the first sum of (14), using the asymptotics for $|A_{\lambda/\alpha_{\mathbf{v}}}|$ provided by the previous lemma. One finds that there exists C, C' such that, for all $\mathbf{v} \in V$ with $\alpha_{\mathbf{v}} \geq \lambda/\alpha_{\max}$,

$$\left| \frac{|X_{\lambda/\alpha_{\mathbf{v}}}|}{|X_\lambda|} - \alpha_{\mathbf{v}} \right| \leq C \alpha_{\mathbf{v}} ((-\log(\lambda/\alpha_{\mathbf{v}}))^{-P}) \leq C' \alpha_{\mathbf{v}} ((-\log(\lambda))^{-P}),$$

the last inequality following from the fact that $x \mapsto \log(\lambda)/\log(\lambda/x)$ is uniformly bounded on $[\alpha_{\max}, \infty)$. Summing over \mathbf{v} gives the required estimate. \square

References

- [1] Adler, R. L., Flatto, L. (1977). Uniform distribution of Kakutani's interval splitting procedure. *Z. Wahrscheinlichkeitstheorie und Verw. Gebiete* **38**, no. 4, 253–259.
- [2] Aistleitner, C., Hofer, M. (2013). Uniform distribution of generalized Kakutani's sequences of partitions. *Ann. Mat. Pura Appl. (4)* **192**, no. 4, 529–538.
- [3] Aistleitner, C., Hofer, M., Ziegler, V. (2014). On the uniform distribution modulo 1 of multidimensional LS-sequences. *Ann. Mat. Pura Appl. (4)* **193**, no. 5, 1329–1344.
- [4] Blackwell, D. (1948). A renewal theorem. *Duke Math. J.* **15**, 145–150.

- [5] Carbone, I. (2012). Discrepancy of LS-sequences of partitions and points. *Ann. Mat. Pura Appl. (4)* **191**, no. 4, 819–844.
- [6] Carbone, I., Iacò, M. R., Volčič, A. (2014). A dynamical system approach to the Kakutani-Fibonacci sequence. *Ergodic Theory Dynam. Systems* **34**, no. 6, 1794–1806.
- [7] Corduneanu, C. (1961). *Almost Periodic Functions*. Interscience Tracts in Pure and App. Math. **22**. Wiley Interscience: London.
- [8] Drmota, M., Infusino, M. (2012). On the discrepancy of some generalized Kakutani’s sequences of partitions. *Unif. Distrib. Theory* **7**, no. 1, 75–104.
- [9] Drmota, M., Tichy, R. (1997). *Sequences, Discrepancies and Applications*. Lecture notes in math. 1651. Springer: Berlin.
- [10] Ellison, F., Ellison, W. (1985). *Prime Numbers*. Wiley Interscience: London.
- [11] Erdős, P., Feller, W., Pollard, H. (1949). A property of power series with positive coefficients. *Bull. Amer. Math. Soc.* **55**, 201–204.
- [12] Falconer, K. (2014). *Fractal geometry: Mathematical foundations and applications*. Third edition. John Wiley & Sons, Ltd.: Chichester.
- [13] Hardy, G. H., Littlewood, J. E. (1914) Some problems of diophantine approximation II. *Acta Math.* **37**, no. 1, 193–239.
- [14] Hofer, M. Iacò, M. R., Tichy, R. (2015). Ergodic properties of β -adic Halton sequences. *Ergodic Theory Dynam. Systems* **35**, no. 3, 895–909.
- [15] Iacò, M. R., Ziegler, V. (2017). Discrepancy of generalized LS-sequences. *Unif. Distrib. Theory* **12**, no. 2, 37–63.
- [16] Kakutani, S. (1976). A problem of equidistribution on the unit interval $[0, 1]$. *Measure theory (Proc. Conf., Oberwolfach, 1975)*. Lecture Notes in Math. **541**. Springer: Berlin. 369–375.
- [17] Smilansky, Y. (2020). *Uniform Distribution of Kakutani Partitions Generated By Substitution Schemes*, Israel J. Math., **240**, 667–710
- [18] Volčič, A. (2011) A generalization of Kakutani’s splitting procedure. *Ann. Mat. Pura Appl. (4)* **190**, no. 1, 45–54.
- [19] Weiss, C. (2018). On the classification of LS-sequences. *Unif. Distrib. Theory* **13**, no. 2, 83–92.
- [20] Weyl, H. (1916). Über die Gleichverteilung von Zahlen mod. Eins. *Math. Ann.* **77**, no. 3, 313–352.