# RANDOMNESS AND NON-ERGODIC SYSTEMS 

JOHANNA N.Y. FRANKLIN AND HENRY TOWSNER


#### Abstract

We characterize the points that satisfy Birkhoff's ergodic theorem under certain computability conditions in terms of algorithmic randomness. First, we use the method of cutting and stacking to show that if an element $x$ of the Cantor space is not Martin-Löf random, there is a computable measure-preserving transformation and a computable set that witness that $x$ is not typical with respect to the ergodic theorem, which gives us the converse of a theorem by V'yugin. We further show that if $x$ is weakly 2 -random, then it satisfies the ergodic theorem for all computable measure-preserving transformations and all lower semicomputable functions.


## 1. Introduction

Random points are typical with respect to measure in that they have no measure-theoretically rare properties of a certain kind, while ergodic theorems describe regular measure-theoretic behavior. There has been a great deal of interest in the connection between these two kinds of regularity recently. We begin by defining the basic concepts in each field and then describe the ways in which they are related. Then we present our results on the relationship between algorithmic randomness and the satisfaction of Birkhoff's ergodic theorem for computable measure-preserving transformations with respect to computable (and then lower semi-computable) functions. Those more familiar with ergodic theory than computability theory might find it useful to first read Section 7, a brief discussion of the notion of algorithmic randomness in the context of ergodic theory.
1.1. Algorithmic randomness in computable probability spaces. Computability theorists seek to calibrate the computational strength of subsets of $\omega$, the nonnegative integers. This calibration is accomplished using the notion of a Turing machine, which can be informally viewed as an idealized computer program (for a general introduction to computability theory, see [22, 23, 28]). We present the main concepts we will need here.

A subset $A$ of the natural numbers $\omega$ is computably enumerable, or c.e., if it is the domain of some Turing machine $P$, that is, the set of numbers

[^0]that, when input into $P$, will result in the program halting and returning an output. One can understand the origin of this terminology intuitively: a set is c.e. if we can generate it by running some Turing machine on more and more inputs and enumerating inputs into our set if the machine returns answers for them. Note that this means that we can list the c.e. sets: since we can enumerate the computer programs $\left\langle P_{i}\right\rangle$, we can enumerate their domains. A set $A$ is computable if both it and its complement are c.e., and we say that a function $f: \omega \rightarrow \omega$ is computable if there is a Turing machine $P$ whose domain is $\omega$ such that for all $n, f(n)=P(n)$. These concepts lead us to the last definition we will need: that of an effectively c.e. sequence. A sequence of c.e. sets is said to be effectively c.e. if there is a computable function $f$ such that the $n^{t h}$ set in the sequence is the $f(n)^{t h}$ c.e. set. We note without ceremony that we can consider c.e. and computable sets of objects other than natural numbers. For instance, there is a computable bijection between $\omega$ and the set of finite binary strings $2^{<\omega}$, and we will often speak of an c.e. subset of $2^{<\omega}$.

Subsets of $\omega$ are often identified in a natural way with infinite binary sequences, or reals: a set $A$ corresponds to the infinite binary sequence whose $(n+1)^{s t}$ bit is 1 if and only if $n$ is in $A$. This is the approach that is most often taken in algorithmic randomness, since some definitions of randomness, such as the initial-segment complexity and the betting strategy definitions, are more naturally phrased in terms of infinite binary sequences than subsets of $\omega$.

We recall the standard notations for sequences that will be used in this paper. We write $2^{\omega}$ for the set of infinite binary sequences, that is, the set of functions from $\omega$ to $\{0,1\}$. As mentioned above, we write $2^{<\omega}$ for the set of finite binary sequences, that is, functions from $[0, n)$ to $\{0,1\}$ for some $n$. We sometimes write finite sequences in the form $\sigma=\left\langle s_{0}, \ldots, s_{n-1}\right\rangle$, where $\sigma$ is the sequence with $\sigma(i)=s_{i}$ for all $i<n$.

If $x$ is a finite or infinite sequence then $|x| \in \omega \cup\{\infty\}$ is the length of the sequence, and if $n \leq|x|$ then $x \upharpoonright n$ is the initial segment of $x$ of length $n$ (that is, the restriction of $x$, as a function, to the domain $[0, n)$ ). We write $x \sqsubseteq y$ if $|x| \leq|y|$ and $y \upharpoonright|x|=x$ (that is, if $x$ is an initial segment of $y$ ) and $x \sqsubset y$ if $x \sqsubseteq y$ and $|x|<|y|$ (that is, if $x$ is a proper initial segment of $y$ ). When $\sigma$ is a finite sequence, $\sigma^{\curvearrowright} y$ is the concatenation of $\sigma$ with $y$-that is, $\sigma^{\curvearrowright} y$ is the sequence with $\left(\sigma^{\frown} y\right)(i)=\sigma(i)$ for $i<|\sigma|$ and $\left(\sigma^{\frown} y\right)(i)=y(i-|\sigma|)$ for $i \geq|\sigma|$. If $\sigma \in 2^{<\omega}$ then $[\sigma] \subseteq 2^{\omega}$ is $\left\{x \in 2^{\omega} \mid \sigma \sqsubset x\right\}$, the set of infinite sequences extending $\sigma$, and if $V \subseteq 2^{<\omega}$ then $[V]=\bigcup_{\sigma \in V}[\sigma]$. We call $[\sigma]$ an interval. We say $V$ is prefix-free if whenever $\sigma, \tau \in V, \sigma \sqsubseteq \tau$ implies $\sigma=\tau$.

We will usually use Greek letters such as $\sigma, \tau, v, \rho, \eta, \zeta, \nu, \theta$ for finite sequences and Roman letters such as $x, y$ for infinite sequences.

For a general reference on algorithmic randomness, see [8, 9, 21]. We will confine our attention to the Cantor space $2^{\omega}$ with the Lebesgue measure $\lambda$. In light of Hoyrup and Rojas' theorem that any computable probability
space is isomorphic to the Cantor space in both the computable and measuretheoretic senses [15], there is no loss of generality in restricting to this case.

We can now present Martin-Löf's original definition of randomness [20].
Definition 1.1. An effectively c.e. sequence $\left\langle V_{i}\right\rangle$ of subsets of $2^{<\omega}$ is a Martin-Löf test if $\lambda\left(\left[V_{i}\right]\right) \leq 2^{-i}$ for every $i$. If $x \in 2^{\omega}$, we say that $x$ is Martin-Löf random if for every Martin-Löf test $\left\langle V_{i}\right\rangle, x \notin \cap_{i}\left[V_{i}\right]$.

It is easy to see that $\lambda\left(\cap_{i}\left[V_{i}\right]\right)=0$ for any Martin-Löf test, and since there are only countably many Martin-Löf tests, almost every point is Martin-Löf random.

In Section 6, we will also consider weakly 2-random elements of the Cantor space. Weak 2 -randomness is a strictly stronger notion than Martin-Löf randomness and is part of the hierarchy introduced by Kurtz in [19].
Definition 1.2. An effectively c.e. sequence $\left\langle V_{i}\right\rangle$ of subsets of $2^{<\omega}$ is a generalized Martin-Löf test if $\lim _{n \rightarrow \infty} \lambda\left(\left[V_{i}\right]\right)=0$. If $x \in 2^{\omega}$, we say that $x$ is weakly ${ }^{2}$-random if for every generalized Martin-Löf test $\left\langle V_{i}\right\rangle, x \notin \cap_{i}\left[V_{i}\right]$.
1.2. Ergodic theory. Now we discuss ergodic theory in the general context of an arbitrary probability space before transferring it to the context of a computable probability space. The following definitions can be found in [14.

Definition 1.3. Suppose $(X, \mu)$ is a probability space, and let $T: X \rightarrow X$ be a measurable transformation.

1) $T$ is measure preserving if for all measurable $A \subseteq X, \mu\left(T^{-1}(A)\right)=$ $\mu(A)$.
2) A measurable set $A \subseteq X$ is invariant under $T$ if $T^{-1}(A)=A$ modulo a set of measure 0 .
3) $T$ is ergodic if it is measure preserving and every $T$-invariant measurable subset of $X$ has measure 0 or measure 1 .

One of the most fundamental theorems in ergodic theory is Birkhoff's Ergodic Theorem:

Birkhoff's Ergodic Theorem. [4] Suppose that ( $X, \mu$ ) is a probability space and $T: X \rightarrow X$ is measure preserving. Then for any $f \in L_{1}(X)$ and almost every $x \in X$,

$$
\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i<n} f\left(T^{i}(x)\right)
$$

converges. Furthermore, if $T$ is ergodic then for almost every $x$ this limit is equal to $\int f d \mu$.

If we restrict ourselves to a countable collection of functions, this theorem gives a natural notion of randomness - a point is random if it satisfies the conclusion of the ergodic theorem for all functions in that collection. In a computable measure space, we can take the collection of sets defined by a computability-theoretic property and attempt to classify this notion in
terms of algorithmic randomness. In particular, we are interested in the following property:

Definition 1.4. Let $(X, \mu)$ be a computable probability space, and let $T$ : $X \rightarrow X$ be a measure-preserving transformation. Let $\mathcal{F}$ be a collection of functions in $L_{1}(X)$. A point $x \in X$ is a weak Birkhoff point for $T$ with respect to $\mathcal{F}$ if for every $f \in \mathcal{F}$,

$$
\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i<n} f\left(T^{i}(x)\right)
$$

converges. $x$ is a Birkhoff point for $T$ with respect to $\mathcal{F}$ if additionally

$$
\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i<n} f\left(T^{i}(x)\right)=\int f d \mu
$$

The definition of a Birkhoff point is only appropriate when $T$ is ergodic; when $T$ is nonergodic, the appropriate notion is that of a weak Birkhoff point.

There are two natural dimensions to consider: the ergodic-theoretic behavior of $T$ and the algorithmic complexity of $\mathcal{C}$. The case where $T$ is ergodic has been largely settled.

A point is Martin-Löf random if and only if the point is Birkhoff for all computable ergodic transformations with respect to lower semi-computable functions [2, 11]. The proof goes by way of a second theorem of ergodic theory:

Poincaré Recurrence Theorem ([24], Chapter 26). Suppose that ( $X, \mu$ ) is a probability space and $T: X \rightarrow X$ is measure preserving. Then for all $E \subseteq X$ of positive measure and for almost all $x \in X, T^{n}(x) \in E$ for infinitely many $n$.

In short, the Poincaré Recurrence Theorem says that an ergodic transformation $T$ returns almost every point to every set of positive measure repeatedly, and Birkhoff's Ergodic Theorem says that it will do so with a well-defined frequency in the limit.

A point $x \in X$ is a Poincaré point for $T$ with respect to $\mathcal{C}$ if for every $E \in \mathcal{C}$ with positive measure, $T^{n}(x) \in E$ for infinitely many $n$. In [18], Kučera proved that a point in the Cantor space is Martin-Löf random if and only if it is a Poincaré point for the shift operator with respect to effectively closed sets. Later, Bienvenu, Day, Mezhirov, and Shen generalized this result and showed that in any computable probability space, a point is MartinLöf random if and only if it is a Poincaré point for computable ergodic transformations with respect to effectively closed sets [2]. The proof that Martin-Löf random points are Poincaré proceeds by showing that a point which is Poincaré for any computable ergodic transformation with respect to effectively closed sets must also be a Birkhoff point for computable ergodic transformations with respect to lower semi-computable functions [2, 11].

| Sets | Transformations |  |
| :---: | :---: | :---: |
|  | Ergodic | Nonergodic |
| Computable | Schnorr | Martin-Löf |
|  | $[13]$ | $[30]+$ Theorem 4.4 |
| Lower semi-computable | Martin-Löf | $?$ |
|  | $[2, ~ 11]$ |  |

Table 1. Randomness notions and ergodicity

Similarly, Gács, Hoyrup, and Rojas have shown that if a point fails to be Schnorr random then there is a computable ergodic transformation where the point fails to be Birkhoff for a bounded computable function [13]. (We say that $x$ is Schnorr random if $x \notin \cap_{i}\left[V_{i}\right]$ for all Martin-Löf tests $\left\langle V_{i}\right\rangle$ where $\lambda\left(\left[V_{i}\right]\right)=2^{-i}$ for all $i$; Schnorr randomness is a strictly weaker notion than Martin-Löf randomness [26].) In fact, the transformation they construct has a stronger property-it is weakly mixing-and they show, conversely, that in a computable weakly mixing transformation every Schnorr random point is Birkhoff. However Rojas has pointed out [25] that this latter result can be strengthened: even in a computable ergodic transformation, every Schnorr random point is Birkhoff for every bounded computable function. Since this last result has not appeared in print, we include it in Section 5 for completeness. Combining these results, we see that a point is Schnorr random if and only if the point is Birkhoff for all computable ergodic transformations with respect to computable functions [13].

In this paper, we consider the analogous situations when $T$ is nonergodic. V'yugin [30] has shown that if $x \in 2^{\omega}$ is Martin-Löf random then $x$ is weakly Birkhoff for any (not necessarily ergodic) computable measure-preserving transformation $T$ with respect to computable functions. Our main result is the converse: that if $x$ is not Martin-Löf random then $x$ is not weakly Birkhoff for some particular transformation $T$ with respect to computable functions (in fact, with respect to computable sets).

These results are summarized in Table 1.
This says that a point is weakly Birkhoff for the specified family of computable transformations with respect to the specified collection of functions if and only if it is random in the sense found in the corresponding cell of the table.

We also begin an analysis of the remaining space in the table; we give an analog of V'yugin's result, showing that if $x$ is weakly 2 -random then $x$ is a weak Birkhoff point for all computable measure-preserving transformations with respect to lower semi-computable functions.

The next two sections will be dedicated to a discussion of the techniques we will use in our construction. Section 2 contains a description of the type of partial transformations we will use to construct the transformation $T$ mentioned above, and Section 3 discusses our methods for building new
partial transformations that extend other such transformations. We combine the material from these two sections to prove our main theorem in Section 4 , while Section 6 contains a further extension of our work and some speculative material on a more relaxed form of upcrossings. Section 7 is a general discussion of algorithmic randomness intended for ergodic theorists.

## 2. Definitions and Diagrams

We will build computable transformations $\widehat{T}: 2^{\omega} \rightarrow 2^{\omega}$ using computable functions $T: 2^{<\omega} \rightarrow 2^{<\omega}$ such that (1) $\sigma \sqsubseteq \tau$ implies $T(\sigma) \sqsubseteq T(\tau)$ and (2) $\widehat{T}(x)=\lim _{n \rightarrow \infty} T(x \upharpoonright n)$ is defined and infinite for all $x \in 2^{\omega}$ outside an $F_{\sigma}$ set with measure 0 .

We will approximate such a $\widehat{T}$ by partial transformations:
Definition 2.1. A partial transformation is a total computable function $T: 2^{<\omega} \rightarrow 2^{<\omega}$ such that if $\sigma \sqsubseteq \tau$ then $T(\sigma) \sqsubseteq T(\tau)$. We write $T \sqsubseteq T^{\prime}$ if for all $\sigma, T(\sigma) \sqsubseteq T^{\prime}(\sigma)$. If $T_{0} \sqsubseteq T_{1} \sqsubseteq \cdots \sqsubseteq T_{n} \sqsubseteq \cdots$ is a sequence of partial transformations, there is a natural limit $T: 2^{<\omega} \rightarrow 2^{<\omega}$ given by $T(\sigma)=\lim _{n} T_{n}(\sigma)$.

We say a computable transformation $\widehat{T}: 2^{\omega} \rightarrow 2^{\omega}$ extends $T$ if for every $\sigma, \widehat{T}([\sigma]) \subseteq[T(\sigma)]$.

The "partial" refers to the fact that we may have $\lim _{n \rightarrow \infty} T(x \upharpoonright n)$ be finite for many or all points. If $T_{0} \sqsubseteq \cdots$ is a uniformly computable sequence of partial transformations with $T=\lim _{n} T_{n}$ and for almost every $x$ the limit $\lim _{n}|T(x \upharpoonright n)|=\infty$, then the transformation $\widehat{T}(x)=\lim _{n} T(x \upharpoonright n)$ is a computable transformation.

We will be exclusively interested in partial transformations which are described finitely in a very specific way:

Definition 2.2. A partial transformation $T$ is proper if there are finite sets $T_{-}, T_{+}$such that:

- $T_{-} \cup T_{+}$is prefix-free,
- $\cup_{\sigma \in T_{-} \cup T_{+}}[\sigma]=2^{\omega}$,
- If there is a $\tau \sqsubseteq \sigma$ such that $\tau \in T_{-}$then $T(\sigma)=T(\tau)$,
- If $\sigma=\tau \frown \rho$ with $\tau \in T_{+}$then $T(\sigma)=T(\tau) \frown \rho$,
- If $\sigma \in T_{-}$then $|T(\sigma)|<|\sigma|$,
- If $\sigma \in T_{+}$then $|T(\sigma)|=|\sigma|$,
- If $\sigma \in T_{+}$and $\sigma \neq \tau \in T_{+} \cup T_{-}$then $T(\tau) \nsupseteq T(\sigma)$,
- If $\sigma \in T_{-}$and $\tau \in T_{+} \cup T_{-}$then $T(\tau) \not \neg T(\sigma)$.

We say $\sigma$ is determined in $T$ if there is some $\tau \sqsubseteq \sigma$ with $\tau \in T_{-} \cup T_{+}$.
When $T$ is a proper partial transformation, we write $T_{-}, T_{+}$for some canonically chosen pair of sets witnessing this fact.

In practice, we will always describe a proper transformation $T$ by describing $T \upharpoonright T_{-} \cup T_{+}$for some particular choice of $T_{-}, T_{+}$, so there is always a canonical choice of $T_{-}$and $T_{+}$. Note that if $\sigma$ is determined then $T(\sigma)$ is


Figure 1. A typical diagram
uniquely defined by the values of $T$ on the finitely many subsequences in $T_{-} \cup T_{+}$. The roles of $T_{-}$and $T_{+}$will be clearer when we introduce a diagrammatic notion for describing transformations. For now, note that once we have $\tau \in T_{+}$, we have entirely determined $\widehat{T} \upharpoonright[\tau]$ for any $\widehat{T}$ extending $T$ : $\widehat{T}\left(\tau^{\frown} x\right)=T(\tau) \subset x$. The requirement that, for such a $\tau,|\tau|=|T(\tau)|$ helps ensure that the resulting transformation is measure preserving.

Throughout this paper, all our partial transformations will be proper.
We will use the method of cutting and stacking, which was introduced by Chacon to produce dynamical systems with specific combinatorial properties [6, 7] Th This method was recently introduced into the study of the computability properties of ergodic theorems by V'yugin [29]. One tries to construct a dynamical system, usually on the real interval $[0,1]$, by specifying the transformation in stages. At a given stage, the interval has been "cut" into a finite number of components, some of which have been "stacked" into "towers" or "ladders." A tower is read upwards, so the interval on the bottom level is mapped by the transformation to the level above, and from that level to the level above that. On the top level of a tower, the transformation is not yet defined. To produce the next stage, the towers are cut into smaller towers and further stacked. By manipulating the order in which the components are stacked, specific properties of the transformation can be enforced. This method has been extensively used in ergodic theory and probability theory to construct examples with specific properties (some overviews of the area are [12, [17, 27]).

A typical diagram is shown in Figure 1. This figure represents that $|\sigma|<$ $\left|\sigma_{0}\right|=\left|\sigma_{1}\right|=\left|\sigma_{2}\right|=\left|\sigma_{3}\right|=\left|\sigma_{4}\right|$ and that for all $v, T\left(\sigma_{i} \smile v\right)=\sigma_{i+1} \smile v$ for $i<4, T\left(\sigma_{4} \smile v\right)=\sigma$, and similarly $T\left(\tau_{i} \smile v\right)=T\left(\tau_{i+1} \frown v\right)$ for $i<2$ while $T\left(\tau_{2}\right)=\langle \rangle$. Although it is not essential to interpret the diagrams, we will

[^1]try to be consistent about the scale of blocks; in Figure 1, the relative width of the blocks suggests that $\left|\tau_{i}\right|=\left|\sigma_{i}\right|+1$-that is, $\lambda\left(\left[\tau_{i}\right]\right)=\lambda\left(\left[\sigma_{i}\right]\right) / 2$; the height of a block does not represent anything, so we draw each block with the same height. The only relevant dimensions are the widths of the blocks and the numbers of blocks in the towers.

In general, a block represents a subset of $2^{\omega}$ of the form $[\sigma]$ for some sequence $\sigma$; by placing the block corresponding to $[\sigma]$ on top of the block corresponding to $[\tau]$, we are indicating that $\tau \in T_{+}$and $T(\tau)=\sigma-$ that is, in the transformation we construct extending $T, T([\tau])=[\sigma]$. (We must, therefore, have $|\sigma|=|\tau|$.) By placing some sequence $\sigma^{\prime}$ with $\left|\sigma^{\prime}\right|<|\sigma|$ on top of the block corresponding to $[\sigma]$, we are indicating that $\sigma \in T_{-}$ and $T(\sigma)=\sigma^{\prime}$ - that is, in the transformation $\widehat{T}$ we construct extending $T$, $\widehat{T}([\sigma]) \sqsubseteq\left[\sigma^{\prime}\right]$.

The roles of $T_{-}$and $T_{+}$in the specification of a proper transformation are now clearer: the elements of $T_{-} \cup T_{+}$are the particular blocks labeled in a given diagram; the elements $\tau \in T_{+}$are those blocks which have another block on top, and therefore we have already defined the value of any extension $\widehat{T}$ of $T$ on every element of $[\tau]$. The elements $\tau \in T_{-}$are topmost blocks of some tower, for which we have (at most) partial information about the ultimate behavior of $\widehat{T}$ on $[\tau]$.

Definition 2.3. We say $\tau$ is blocked if there is any $\sigma$ such that $T(\sigma) \sqsupseteq \tau$. Otherwise we say $\tau$ is unblocked.

An open loop in a partial transformation $T$ is a sequence $\sigma_{0}, \ldots, \sigma_{n}$ such that:

- $\left|\sigma_{0}\right|=\left|\sigma_{1}\right|=\cdots=\left|\sigma_{n}\right|$,
- $T\left(\sigma_{i}\right)=\sigma_{i+1}$ for $i<n$,
- $T\left(\sigma_{n}\right) \sqsubset \sigma_{0}$,
- $\sigma_{0}$ is unblocked.

We refer to $\sigma_{0}$ as the initial element of the open loop $\sigma_{0}, \ldots, \sigma_{n}$ and $\sigma_{n}$ as the final element.

The width of an open loop is the value $2^{-\left|\sigma_{i}\right|}$.
We say $T$ is partitioned into open loops if for every determined $\sigma$ there is an open loop $\sigma_{0}, \ldots, \sigma_{n}$ in $T$ with $\sigma=\sigma_{i}$ for some $i$. (In a proper transformation such an open loop must be unique.) In such a transformation we write $\mathcal{L}_{T}(\sigma)$ for the open loop $\sigma_{0}, \ldots, \sigma_{n}$ such that for some $i, \sigma=\sigma_{i}$. We write $L_{T}(\sigma)$ for $n+1$, the length of the open loop containing $\sigma$.
(We are interested in open loops to preclude the possibility that $T\left(\sigma_{n}\right)=$ $\sigma_{0}$, since we are not interested in-indeed, will not allow the existence of"closed" loops.) Diagrammatically, the requirement that $T$ be partitioned into open loops is represented by requiring that any sequence written above a tower of blocks is a subsequence of the sequence at the bottom of that tower. (For instance, in Figure 1, we require that $\sigma \sqsubset \sigma_{0}$.)

One of the benefits of proper transformations partioned into open loops is that they ensure that our transformation is measure preserving:

Lemma 2.4. Let $T_{0} \sqsubseteq \ldots \sqsubseteq T_{n} \sqsubseteq \ldots$ be a sequence of proper partial transformations partitioned into open loops, let $T=\lim _{n} T_{n}$, and let $\widehat{T}(x)=$ $\lim _{n \rightarrow \infty} T(x \upharpoonright n)$. Suppose that $|\widehat{T}(x)|$ is infinite outside a set of measure 0. Then $\widehat{T}$ is measure preserving.

Proof. It suffices to show that for every $\sigma$ and every $\epsilon>0$,

$$
\left|\lambda\left(\widehat{T}^{-1}([\sigma])\right)-\lambda([\sigma])\right|<\epsilon
$$

Fix $\sigma$ and $\epsilon>0$ and choose $n$ large enough that if $U$ is the set of $\tau \in T_{n,-}$ with $\left|T_{n}(\tau)\right|<|\sigma|, \sum_{\tau \in U} \lambda([\tau])<\epsilon$.

If there is a $\tau \in T_{n,+}$ with $T_{n}(\tau) \sqsubseteq \sigma$ then there is a $\rho$ with $\left|\tau^{\frown} \rho\right|=|\sigma|$ and $T_{n}\left(\tau^{\frown} \rho\right)=\sigma$ and therefore $\widehat{T}([\tau \frown \rho])=[\sigma]$. Since for $m \geq n$ and $v \in T_{m,+} \cup T_{m,-}$ with $v \nsupseteq \tau$ we have $T_{m}(v) \nsupseteq T_{n}(\tau)$, we therefore have $\widehat{T}^{-1}([\sigma])=\left[\tau^{\frown} \rho\right]$ and are done.

Otherwise let $\sigma_{0}, \ldots, \sigma_{k}$ be such that each $\sigma_{i} \in T_{n,+} \cup T_{n,-}$ and $\bigcup_{i \leq k}\left[\sigma_{i}\right]=$ $[\sigma]$. Each $\sigma_{i}$ belongs to an open loop. Let $I$ be the set of $i \leq k$ such that $\sigma_{i}$ is the initial element of $\mathcal{L}_{T_{n}}\left(\sigma_{i}\right)$. If $i \notin I$ then there is a $\tau_{i}$ with $\left|\tau_{i}\right|=\left|\sigma_{i}\right|$ and $T_{n}\left(\tau_{i}\right)=\sigma_{i}$, and therefore $\widehat{T}^{-1}\left(\left[\sigma_{i}\right]\right)=\left[\tau_{i}\right]$.

If $i \in I$, let $v_{i}$ be the final element of $\mathcal{L}_{T_{n}}\left(\sigma_{i}\right)$, so $\left|v_{i}\right|=\left|\sigma_{i}\right|$ but $T_{n}\left(v_{i}\right) \sqsubset \sigma_{i}$. Let $I^{\prime} \subseteq I$ be those $i$ such that $\left|T_{n}\left(v_{i}\right)\right|<|\sigma|$; note that if $i \in I \backslash I^{\prime}$ then $\sigma \sqsubseteq T_{n}\left(v_{i}\right) \sqsubset \sigma_{i}$, so $\left[v_{i}\right] \subseteq \widehat{T}^{-1}([\sigma])$. Also, $\lambda\left(\bigcup_{i \in I^{\prime}}\left[v_{i}\right]\right)<\epsilon$ and

$$
\widehat{T}^{-1}([\sigma]) \supseteq \bigcup_{i \notin I}\left[\tau_{i}\right] \cup \bigcup_{i \in I \backslash I^{\prime}}\left[v_{i}\right]
$$

and so

$$
\begin{aligned}
\lambda\left(\widehat{T}^{-1}([\sigma])\right) & \geq \sum_{i \notin I} \lambda\left(\left[\tau_{i}\right]\right)+\sum_{i \in I \backslash I^{\prime}} \lambda\left(\left[v_{i}\right]\right) \\
& =\sum_{i \notin I} \lambda\left(\left[\tau_{i}\right]\right)+\sum_{i \in I} \lambda\left(\left[v_{i}\right]\right)-\sum_{i \in I^{\prime}} \lambda\left(\left[v_{i}\right]\right) \\
& =\sum_{i \notin I} \lambda\left(\left[\sigma_{i}\right]\right)+\sum_{i \in I} \lambda\left(\left[\sigma_{i}\right]\right)-\sum_{i \in I^{\prime}} \lambda\left(\left[v_{i}\right]\right) \\
& =\sum_{i \leq k} \lambda\left(\left[\sigma_{i}\right]\right)-\sum_{i \in I^{\prime}} \lambda\left(\left[v_{i}\right]\right) \\
& >\lambda([\sigma])-\epsilon .
\end{aligned}
$$

Now consider any $\tau \in T_{n,-}$. If $\left|T_{n}(\tau)\right| \geq|\sigma|$ then we have either $T_{n}(\tau) \sqsupseteq$ $\sigma$, in which case $\tau \in\left\{\tau_{i}\right\} \cup\left\{v_{i}\right\}$, or $T_{n}(\tau) \nexists \sigma$, in which case $T^{-1}([\sigma]) \cap[\tau]=$ $\emptyset$. Therefore we have

$$
\widehat{T}^{-1}([\sigma]) \subseteq \bigcup_{i \notin I}\left[\tau_{i}\right] \cup \bigcup_{i \in I \backslash I^{\prime}}\left[v_{i}\right] \cup \bigcup_{\tau \in U}[\tau]
$$

and so

$$
\begin{aligned}
\lambda\left(\widehat{T}^{-1}([\sigma])\right) & \leq \sum_{i \leq k} \lambda\left(\left[\sigma_{i}\right]\right)+\sum_{\tau \in U} \lambda([\tau]) \\
& <\lambda([\sigma])+\epsilon,
\end{aligned}
$$

completing the proof.
A similar argument shows that also $\lambda(\widehat{T}(A))=\lambda(A)$ for all $A$, but we do not need this fact.

Definition 2.5. When $A \subseteq 2^{\omega}$ and $\sigma \in 2^{<\omega}$, we write $\sigma \in A$ ( $\sigma$ is in $A$ ) if $[\sigma] \subseteq A$. Similarly, we say a sequence $\sigma_{0}, \ldots, \sigma_{k}$ is in $A$ if for each $i \leq k, \sigma_{i}$ is in $A$.

We say $\sigma$ avoids $A$ if $[\sigma] \cap A=\emptyset$. Similarly we say a sequence $\sigma_{0}, \ldots, \sigma_{n}$ avoids $A$ if for each $i \leq k, \sigma_{i}$ avoids $A$.

So $\sigma$ is in $A$ iff $\sigma$ avoids $2^{\omega} \backslash A$.
Definition 2.6. An escape sequence for $\sigma_{0}$ in $T$ is a sequence $\sigma_{1}, \ldots, \sigma_{n}$ such that:

- $\left|\sigma_{1}\right|=\left|\sigma_{2}\right|=\ldots=\left|\sigma_{n}\right|$,
- For all $0 \leq i<n, \sigma_{i+1} \sqsupseteq T\left(\sigma_{i}\right)$,
- If $\sigma_{i+1}$ is blocked then $\sigma_{i+1}=T\left(\sigma_{i}\right)$,
- $T\left(\sigma_{n}\right)=\langle \rangle$,
- All $\sigma_{i}$ are determined.

We say $T$ is escapable if for every determined $\sigma$ with $|T(\sigma)|<|\sigma|$, there is an escape sequence for $\sigma$. If $A, B \subseteq 2^{\omega}$, we say $T$ is $A, B$-escapable if for every determined $\sigma$ in $A$ with $|T(\sigma)|<|\sigma|$, there is an escape sequence for $\sigma$ in $B$.

An escape sequence for $\sigma_{0}$ is reduced if (1) $\sigma_{i} \sqsupseteq T\left(\sigma_{j}\right)$ implies that either $i \leq j+1$ or $\sigma_{i}$ is blocked, and (2) if $i<n$, then $T\left(\sigma_{i}\right) \neq\langle \rangle$.

This is the first of many places where we restrict consideration to determined $\sigma$ with $|T(\sigma)|<|\sigma|$. Note that this is the same as restricting to those $\sigma$ such that there is some $\tau \sqsubseteq \sigma$ with $\tau \in T_{-}$. When $T\left(\sigma_{0}\right)=\langle \rangle$, the empty sequence is a valid escape sequence for $\sigma_{0}$ (and the unique reduced escape sequence).

Escapability preserves the option of extending $T$ in such a way that we can eventually map $\left[\sigma_{0}\right]$ to anything not already in the image of another sequence (although it may require many applications of $T$ ). $A, B$-escapability will be useful at intermediate steps of our construction; typically we want to know that we have $A, B$-escapability so that we can manipulate portions outside of $B$ with interfering with escapability. $\emptyset, \emptyset$-escapable is the same as escapable.

Lemma 2.7. Every escape sequence for $\sigma$ contains a reduced subsequence for $\sigma$.

Proof. We proceed by induction on the length of the sequence. It suffices to show that if $\sigma_{1}, \ldots, \sigma_{n}$ is a nonreduced escape sequence then there is a proper subsequence which is also an escape sequence for $\sigma_{0}$. If for some $i>j+1, \sigma_{i} \sqsupseteq T\left(\sigma_{j}\right)$ with $\sigma_{i}$ unblocked, then $\sigma_{1}, \ldots, \sigma_{j}, \sigma_{i}, \ldots, \sigma_{n}$ is also an escape sequence. If for some $i<n, T\left(\sigma_{i}\right)=\langle \rangle$ then $\sigma_{1}, \ldots, \sigma_{i}$ is also an escape sequence.

Clearly if a sequence is in $B$, any subsequence is as well.
The following lemma is immediate from the definition of an escape sequence.

Lemma 2.8. If $\sigma_{1}, \ldots, \sigma_{n}$ is an escape sequence for $\sigma_{0}$ in $T$ then for every $\rho, \sigma_{1} \frown \rho, \ldots, \sigma_{n} \frown \rho$ is also an escape sequence.

Note that if $\sigma_{1}, \ldots, \sigma_{n}$ is in some set $B$, so is $\sigma_{1} \frown \rho, \ldots, \sigma_{n} \frown \rho$.
If $\tau_{1}, \ldots, \tau_{n}$ is an escape sequence, it is possible to choose an extension $\widehat{T}$ of $T$ and $x \in\left[\tau_{0}\right]$ so that $T^{i}(x) \in\left[\tau_{i}\right]$ for all $i \leq n$. If $\sigma_{0}, \ldots, \sigma_{k}$ is an open loop in $T$ then in every extension $\widehat{T}$ of $T$, whenever $x \in\left[\sigma_{0}\right]$, we have $T^{i}(x) \in\left[\sigma_{i}\right]$ for $i \leq k$. Moreover, because $\widehat{T}$ is measure preserving, if $y \in\left[\sigma_{i+1}\right]$ then $T^{-1}(y) \in\left[\sigma_{i}\right]$. The next lemma shows that these properties interact-an escape sequence can only "enter" an open loop at the beginning, and if this happens, the escape sequence must then traverse the whole open loop in order.
Lemma 2.9. Suppose $\sigma_{0}, \ldots, \sigma_{k}$ is an open loop in $T$ consisting of determined elements, $\tau_{1}, \ldots, \tau_{n}$ is a reduced escape sequence for $\tau_{0}$, and $\left|T\left(\tau_{0}\right)\right|<$ $\left|\tau_{0}\right|$. Then one of the following occurs:

- $\tau_{0} \sqsupseteq \sigma_{k}$ and for $j>0, \tau_{j} \notin \cup_{i \leq k}\left[\sigma_{i}\right]$,
- There is a unique $j>0$ such that for all $i \leq k, \tau_{j+i} \sqsupseteq \sigma_{i}$,
- For all $j, \tau_{j} \notin \cup_{i \leq k}\left[\sigma_{i}\right]$.

Proof. First, suppose some $\tau_{j} \sqsupseteq \sigma_{i}$. Let $j$ be least such that this is the case. If $j=0$ then since $\left|T\left(\tau_{0}\right)\right|<\left|\tau_{0}\right|$, we must have $i=k$.

Suppose $j \neq 0$. If $i \neq 0$ then since $\sigma_{i}$ is blocked, we must have $\tau_{j-1} \sqsubseteq \sigma_{i-1}$, which is impossible by our choice of $j$. So there is a $j>0$ with $\tau_{j} \sqsupseteq \sigma_{0}$. Since each $\sigma_{i}$ satisfies $T\left(\sigma_{i}\right)=\sigma_{i+1}$ with $\left|T\left(\sigma_{i}\right)\right|=\left|\sigma_{i+1}\right|$ and $T$ is proper, for each $i \leq k$, we must have $\tau_{j+i} \sqsupseteq \sigma_{i}$.

So we have shown that if $\tau_{j} \sqsupseteq \sigma_{i}$ for some $i, j$ with $j>0$ then we have a complete copy of the open loop in our escape sequence. We now show that if $j<j^{\prime}$ and $\tau_{j} \sqsupseteq \sigma_{k}$ then we cannot have $\tau_{j^{\prime}} \sqsupseteq \sigma_{i}$; this shows both the second half of the first case and the uniqueness in the second case. For suppose we had $\tau_{j} \sqsupseteq \sigma_{k}, j^{\prime}>j$, and $\tau_{j^{\prime}} \sqsupseteq \sigma_{i}$. By the previous paragraph, we may assume $i=k$. But since $\sigma_{k}$ is determined and $\left|T\left(\sigma_{k}\right)\right|<\sigma_{k}$, we have $T\left(\tau_{j}\right)=T\left(\sigma_{k}\right)=T\left(\tau_{j^{\prime}}\right)$. This means $T\left(\tau_{j}\right) \sqsubseteq \tau_{j^{\prime}+1}$.

Note that $\tau_{j^{\prime}+1}$ cannot be blocked: we have $\left|T\left(\tau_{j^{\prime}}\right)\right|=\left|T\left(\sigma_{k}\right)\right|<\left|\sigma_{k}\right| \leq$ $\left|\tau_{j^{\prime}}\right|=\left|\tau_{j^{\prime}+1}\right|$ since $j^{\prime}>j \geq 0$, so $T\left(\tau_{j^{\prime}}\right) \neq \tau_{j^{\prime}+1}$. So we must have $j^{\prime}+1 \leq$ $j+1$, contradicting the assumption that $j<j^{\prime}$.

## 3. Working with Transformations

We will work exclusively with partial transformations with a certain list of properties which, for purposes of this paper, we call "useful" partial transformations. In this section we describe some basic operations which can be used to manipulate useful transformations. While these operations are ultimately motivated by the construction in the next section, they also provide some intuition for why useful transformations deserve their name.

Definition 3.1. A partial transformation $T$ is useful if:

- $T$ is proper,
- $T$ is partitioned into open loops, and
- $T$ is escapable.

The next lemma illustrates one of the advantages of always working with open loops: we can always modify a useful partial transformation by replacing a open loop with a new open loop of the same total measure but arbitrarily small width.

Lemma 3.2 (Thinning Loops). Let $T$ be a useful partial transformation, let $\sigma_{0}, \ldots, \sigma_{k}$ be an open loop of determined elements and let $\epsilon=2^{-n}$ be smaller than the width of this open loop. Then there is a useful $T^{\prime} \sqsupseteq T$ such that

There is an open loop $\tau_{0}, \ldots, \tau_{k^{\prime}}$ in $T^{\prime}$ of width $\epsilon$ such that $\cup_{j \leq k^{\prime}}\left[\tau_{i}\right]=\cup_{i \leq k}\left[\sigma_{i}\right]$.
Furthermore,

- If $\tau \notin \cup_{i \leq k^{\prime}}\left[\tau_{i}\right]$ is determined in $T^{\prime}$ then $T^{\prime}(\tau)=T(\tau)$ and $L_{T}(\tau)=$ $L_{T^{\prime}}(\tau)$,
- If $T$ is $A, B$-escapable and $\sigma_{0}, \ldots, \sigma_{k}$ is in $B$ then $T^{\prime}$ is $A, B$-escapable as well,
- If $T$ is $A, B$-escapable and $\sigma_{0}, \ldots, \sigma_{k}$ avoids $B$ then $T^{\prime}$ is $A, B$-escapable as well.

Proof. Figure 2 illustrates this lemma. Formally, let the width of $\sigma_{0}, \ldots, \sigma_{k}$ be $2^{-m}$ with $m \leq n$. For any $v$ with $|v|=n-m$, by $v+1$ we mean the result of viewing $v$ as a sequence $\bmod 2$ and adding 1 to it, so $010+1=011$ while $011+1=100$.

Define $T^{\prime} \sqsupseteq T$ by:

- If $\tau=\sigma_{k} \frown \rho$ where $|\rho|=n-m$ and $v$ is not all 1 's then $T^{\prime}(\tau)=$ $\sigma_{0}-(\rho+1)$,
- Otherwise $T^{\prime}(\tau)=T(\tau)$.

Since this is our first construction of this kind, we point out that this is an operation on the description of $T$ as a proper partial transformation: we


Figure 2. Thinning Loops, Lemma 3.2
have $\sigma_{0}, \ldots, \sigma_{k-1} \in T_{+}$but $\sigma_{k} \in T_{-}$. We define

$$
\begin{aligned}
T_{+}^{\prime} & =T_{+} \backslash\left\{\sigma_{i} \mid i<k\right\} \\
& \cup\left\{\sigma_{i} \frown v|i<k,|v|=n-m\}\right. \\
& \cup\left\{\sigma_{k} \frown v| | v \mid=n-m, v \neq\langle 1, \ldots, 1\rangle\right\}
\end{aligned}
$$

and

$$
T_{-}^{\prime}=T_{-} \backslash\left\{\sigma_{k}\right\} \cup\left\{\sigma_{k}\langle 1, \ldots, 1\rangle\right\} .
$$

We have explained how to define $T^{\prime}$ on all elements of $T_{+}^{\prime} \backslash T_{+}$and $T_{-}^{\prime} \backslash T_{-}$, and when $\sigma$ is determined in $T$ but not $T^{\prime}$ we set $T^{\prime}(\sigma)=T(\sigma)$ and so have completely specified the new transformation $T^{\prime}$.

Propriety and the fact that $T^{\prime}$ is partitioned into open loops are trivial.
To see escapability, consider some $v_{0}$ determined in $T^{\prime}$ such that $\left|T^{\prime}\left(v_{0}\right)\right|<$ $\left|v_{0}\right|$ and fix a reduced escape sequence $v_{1}, \ldots, v_{r}$ for $v_{0}$ in $T$. If $v_{0} \in \cup\left[\sigma_{i}\right]$ then by Lemma 2.9 we have $v_{0} \sqsupseteq \sigma_{k}$ and no other element of the escape sequence belongs to $\cup\left[\sigma_{i}\right]$, and therefore $v_{1}, \ldots, v_{r}$ is an escape sequence in $T^{\prime}$ as well.

If $v_{0} \notin \cup\left[\sigma_{i}\right]$ but there is a $j>0$ such that $v_{j+i} \sqsupseteq \sigma_{i}$ for $i \leq k$ then by Lemma 2.8 we may assume that for $j>0,\left|v_{j}\right| \geq n$. Then since both $T\left(v_{j-1}\right) \sqsubseteq v_{j}$ and $\sigma_{0} \sqsubseteq v_{j}, T\left(v_{j-1}\right)$ and $\sigma_{0}$ are comparable. Since $\sigma_{0}$ is unblocked, we must have $T\left(v_{j-1}\right) \sqsubset \sigma_{0}$, and therefore for any $\rho$ of suitable length, the sequence

$$
v_{1}, \ldots, v_{j-1}, \sigma_{0} \frown\langle 0, \ldots, 0\rangle \frown \rho, \ldots, \sigma_{k} \frown\langle 1, \ldots, 1\rangle \frown \rho, v_{j+k+1}, \ldots, v_{n}
$$

is an escape sequence for $v_{0}$ in $T^{\prime}$.
To see that we preserve $A, B$-escapability, if $v_{0}$ is in $A$ then we could have chosen the original escape sequence $v_{1}, \ldots, v_{r}$ in $B$, and therefore (since either $\sigma_{0}, \ldots, \sigma_{k}$ is in $B$ or avoids $B$ ), the same argument shows that in $T^{\prime}$ there is an escape sequence for $v_{0}$ in $B$.
Remark 3.3. In the previous lemma, we actually have slightly more control over escape sequences: for any set $B$ such that $\sigma_{0}, \ldots, \sigma_{k}$ avoids $B$, any
escape sequence in $T$ is an escape sequence in $T^{\prime}$. In particular, in this situation we do not change the lengths of escape sequences.

We also need a modified version of the above lemma where instead of wanting $\cup_{j \leq k^{\prime}}\left[\tau_{i}\right]=\cup_{i \leq k}\left[\sigma_{i}\right]$ we want to have a small amount of the original open loop left alone.
Lemma 3.4. Let $T$ be a useful partial transformation, let $\sigma_{0}, \ldots, \sigma_{k}$ be an open loop of determined elements and let $\epsilon=2^{-n}$ be smaller than the width of this open loop. Then there is a useful $T^{\prime} \sqsupseteq T$ such that

There is an open loop $\tau_{0}, \ldots, \tau_{k^{\prime}}$ in $T^{\prime}$ of width $\epsilon$ such that $\lambda\left(\bigcup_{i \leq k}\left[\sigma_{i}\right] \backslash \bigcup_{j \leq k^{\prime}}\left[\tau_{i}\right]\right)=\epsilon \cdot(k+1)$.

## Furthermore:

- If $\tau \notin \cup_{i \leq k^{\prime}}\left[\tau_{i}\right]$ is determined in $T^{\prime}$ then $T^{\prime}(\tau)=T(\tau)$ and $L_{T}(\tau)=$ $L_{T^{\prime}}(\tau)$,
- If $T$ is $A, B$-escapable and $\sigma_{0}, \ldots, \sigma_{k}$ is in $B$ then $T^{\prime}$ is $A, B \backslash \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]-$ escapable,
- If $T$ is $A, B$-escapable and $\sigma_{0}, \ldots, \sigma_{k}$ avoids $B$ then $T^{\prime}$ is $A, B$-escapable.

Proof. We proceed exactly as above except that we replace the first clause in the definition of $T^{\prime}$ with

If $\tau=\sigma_{k} \frown v$ where $|v|=n-m$ and $v$ is neither all 1's nor all
1 's with a single 0 at the end then $T^{\prime}(\tau)=\sigma_{0} \frown(v+1)$.
Equivalently, we place both $\sigma_{k} \frown\langle 1, \ldots, 1,0\rangle$ and $\sigma_{k} \frown\langle 1, \ldots, 1\rangle$ in $T_{-}^{\prime}$ and all other extensions of $\sigma_{k}$ in $T_{+}^{\prime}$.

We check the stronger escapability condition. If $v_{0}$ is in $A$ and is determined in $T^{\prime}$ with $\left|T^{\prime}\left(v_{0}\right)\right|<\left|v_{0}\right|$, take a reduced escape sequence $v_{1}, \ldots, v_{r}$ for $v_{0}$ in $T$ in $B$. By Lemma 3.2 , we may assume $\left|v_{1}\right| \geq n$. The only non-trivial case is if $v_{0} \notin \bigcup\left[\sigma_{i}\right]$ but there is a $j>0$ such that $v_{j+i} \sqsupseteq \sigma_{i}$ for $i \leq k$. Let $v_{j+i}=\sigma_{i} \subset \rho$ (note that $\rho$ does not depend on $i$ ). For each $i \leq k$, we may replace $v_{j+i}$ with $\sigma_{i} \frown\langle 1, \ldots, 1\rangle \frown \rho^{\prime}$ for any $\rho^{\prime}$ of appropriate length; since $T\left(v_{j-1}\right) \sqsubseteq \sigma_{0}$ (because $\sigma_{0}$ is determined and $T$ is proper), this remains an escape sequence, and the modified escape sequence avoids $B \backslash \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]$.
Remark 3.5. As before, the escape sequences in $T^{\prime}$ promised by the last two conditions in this lemma always have the same length as the ones in $T$.

The next lemma illustrates the use of escape sequences: we take some $\sigma_{0}$ and an escape sequence in $T$ and extend $T$ to a new partial transformation $T^{\prime}$ with the property that $\sigma_{0} \in T_{+}^{\prime}$ and the open loop $\tau_{0}, \ldots, \tau_{k^{\prime}}$ in $T^{\prime}$ which contains $\sigma_{0}$ has the property that $T^{\prime}\left(\tau_{k^{\prime}}\right)=\langle \rangle$. In other words, we can arrange for any $\sigma_{0}$ to belong to a tower which has $\rangle$ on top.
Lemma 3.6 (Escape). Let $T$ be a useful partial transformation, let $\sigma_{0}$ be determined with $\left|T\left(\sigma_{0}\right)\right|<\left|\sigma_{0}\right|$, and let $\sigma_{1}, \ldots, \sigma_{k}$ be a reduced escape sequence for $\sigma_{0}$ such that $\left|\sigma_{0}\right|=\left|\sigma_{1}\right|+1$. Then there is a useful $T^{\prime} \sqsupseteq T$ such that

There is an open loop $\tau_{0}, \ldots, \tau_{k^{\prime}}$ in $T^{\prime}$ with $\left[\sigma_{0}\right] \subseteq \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]$ such that $T^{\prime}\left(\tau_{k^{\prime}}\right)=\langle \rangle$.

## Furthermore,

- If $\tau \notin \cup_{i \leq k^{\prime}}\left[\tau_{i}\right]$ is determined in $T^{\prime}$ then $T^{\prime}(\tau)=T(\tau)$ and $L_{T}(\tau)=$ $L_{T^{\prime}}(\tau)$,
- If $T$ is $A, B$-escapable where $\tau_{0}, \ldots, \tau_{k^{\prime}}$ is in $B$ then $T^{\prime}$ is $A, B$ escapable,
- If $T$ is $A, B$-escapable and $\sigma_{0}$ avoids $B$ then $T^{\prime}$ is $A \backslash \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right], B \backslash$ $\bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]$-escapable.

Proof. Let $T_{-}, T_{+}$witness that $T$ is proper and extend $T$ by defining $T^{\prime}\left(\sigma_{0}\right)=$ $\sigma_{1} \frown\langle 0\rangle$ and for each $i \in(0, k), T^{\prime}\left(\sigma_{i} \frown\langle 0\rangle\right)=\sigma_{i+1} \frown\langle 0\rangle$, and for all $\tau$ which do not extend some $\sigma_{i}$ with $i<k, T^{\prime}(\tau)=T(\tau)$.

To see that $T^{\prime}$ is proper, we need only check that if $\sigma \in T_{+}^{\prime}$ and $\sigma \neq$ $\tau \in T_{+}^{\prime} \cup T_{-}^{\prime}$ then $T^{\prime}(\tau) \nexists T^{\prime}(\sigma)$. Clearly we need only check this for $T^{\prime}(\sigma)=\sigma_{i} \frown\langle 0\rangle$. Since the escape sequence was reduced, we cannot have $\sigma_{i}=\sigma_{j}$ for $i \neq j$, so we can restrict our attention to the $\tau$ such that $T^{\prime}(\tau)=T(\tau)$. If $\sigma_{i}$ was not blocked in $T$ then there is no such $\tau$, and if $\sigma_{i}$ was blocked in $T$ then already $T\left(\sigma_{i-1} \frown\langle 0\rangle\right)=\sigma_{i} \frown\langle 0\rangle$, and the claim follows since $T$ was proper.

It is easy to see that $T^{\prime}$ remains partitioned into open loops.
We check that $T^{\prime}$ is escapable. Let $v_{0}$ be given with $\left|T^{\prime}\left(v_{0}\right)\right|<\left|v_{0}\right|$. Then the same was true in $T$, so $v_{0}$ had an escape sequence $v_{1}, \ldots, v_{r}$ in $T$. We may assume $\left|v_{1}\right| \geq\left|\sigma_{1}\right|+1$. There are a few potential obstacles we need to deal with. First, it could be that for some $i, v_{i} \sqsupseteq \sigma_{0}$. Letting $v_{i}=\sigma_{0} \frown \rho$, we must have that $v_{1}, \ldots, \sigma_{0} \frown \rho, \sigma_{1} \frown\langle 0\rangle \frown \rho, \ldots, \sigma_{k} \frown\langle 0\rangle \frown \rho$ is also an escape sequence for $v_{0}$.

Suppose not. There could be some $i$ and $j>0$ such that $v_{i} \sqsupseteq \sigma_{j} \frown\langle 0\rangle$. Let $j$ be least such that this occurs. We cannot have $i=0$ since $\left|T\left(\sigma_{j} \frown\langle 0\rangle \frown \rho\right)\right|=$ $\left|\sigma_{j} \frown\langle 0\rangle \frown \rho\right|$ (unless $j=k$, in which case $T\left(v_{0}\right)=\langle \rangle$ and so the empty sequence is an escape sequence). If $T\left(v_{i-1}\right) \sqsupseteq \sigma_{j} \frown\langle 0\rangle$ then $\sigma_{j}$ was blocked in $T$, so $T\left(\sigma_{j-1}\right)=\sigma_{j}$, and therefore $v_{i-1} \sqsupseteq \sigma_{j-1} \frown\langle 0\rangle$, contradicting the leastness of $j$. In particular, $\sigma_{j}^{\frown}\langle 1\rangle \frown \rho \sqsupseteq T\left(v_{i-1}\right)$ as well. So whenever we have $v_{i}=\sigma_{j} \frown\langle 0\rangle \frown \rho$, we replace it with $v^{\prime}-i=\sigma_{j} \frown\langle 1\rangle \frown \rho$, and the result is still an escape sequence.

Suppose $T$ is $A, B$-escapable, $\tau_{0}, \ldots, \tau_{k^{\prime}}$ is in $B$, and $v_{0}$ is in $A$. Then the argument just given, applied to an escape sequence in $T$ in $B$, gives an escape sequence in $T^{\prime}$ in $B$.

Suppose $T$ was $A, B$-escapable, $\left[\sigma_{0}\right]$ avoids $B$, and $v_{0}$ is in $A$. Then, taking $v_{1}, \ldots, v_{r}$ to be an escape sequence in $T$ in $B$, we cannot have $v_{i} \sqsupseteq \sigma_{0}$, and so we are in the second case of the argument above, which gives an escape sequence in $B \backslash \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]$.


Figure 3.

Remark 3.7. Once again, the last two clauses actually ensure that when we have an escape sequence in $T$ avoiding $B$, we actually have an escape sequence of the same length in $T^{\prime}$ avoiding $B$ or $B \backslash \bigcup_{i \leq k^{\prime}}\left[\tau_{i}\right]$.

The main building block of our construction will combine the steps given by these two lemmas as illustrated in Figure 3 .

In the figure, we have $\sigma \sqsubset \sigma^{\prime}$ and $\tau \sqsubset \tau_{0}$. We begin in a situation where we have an open loop (the one on the right in the "Before" diagram in Figure 3) and we wish to arrange the blocks so that the elements of that open loop belong to an open loop with $\rangle$ on top such as the one on the left (so that we may later place any other open loop on top of it). The sequence $\sigma^{\prime}, \tau_{0}, \tau_{1}, \tau_{2}$ is an escape sequence for $v$ (and the end of an escape sequence for the block below $v$ ). We could simply combine all these open loops - place $\sigma^{\prime}$ on top of $v$ and $\tau_{0}$ on top of $\sigma^{\prime}$-but we don't wish to do so, because we don't want to use up all of the escape sequence; we might be using $\sigma^{\prime}, \tau_{0}, \tau_{1}, \tau_{2}$ as part of an escape sequence for other elements as well, and we need some of it to remain.

So we apply Lemma 3.2 to the open loop on the right, replacing it with a much thinner open loop. Now we can apply Lemma 3.6, which takes subintervals from $\sigma^{\prime}, \tau_{0}, \tau_{1}, \tau_{2}$ and places them above the open loop on the right. Note that, by applying Lemma 3.2 with very small $\epsilon$, we can make the total measure of the shaded portion as small as we like.

## 4. The Main Construction

Our main tool for causing the Birkhoff ergodic theorem to fail at a point is the notion of an upcrossing.

Definition 4.1. Given a measurable, measure-preserving, invertible $\widehat{T}$ : $2^{\omega} \rightarrow 2^{\omega}$, a point $x \in 2^{\omega}$, a measurable $f$, and rationals $\alpha<\beta$, an upcrossing sequence for $\alpha, \beta$ of length $N$ is a sequence

$$
0 \leq u_{1}<v_{1}<u_{2}<v_{2}<\cdots<u_{N}<v_{N}
$$

such that for all $i \leq N$,

$$
\frac{1}{u_{i}+1} \sum_{j=0}^{u_{i}} f\left(\widehat{T}^{j} x\right)<\alpha, \frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}} f\left(\widehat{T}^{j} x\right)>\beta .
$$

$\tau(x, f, \alpha, \beta)$ is the supremum of the lengths of upcrossing sequences for $\alpha, \beta$.

By definition, Birkhoff's ergodic theorem fails at $x$ exactly if $\tau(x, f, \alpha, \beta)=$ $\infty$ for some $\alpha<\beta$. Our plan is to look at an Martin-Löf test $\left\langle V_{j}\right\rangle$ and, as sequences $\sigma$ are enumerated into an appropriate $V_{j}$, ensure that the lower bound on $\tau(x, f, 1 / 2,3 / 4)$ increases for each $x \in[\sigma]$.

While building a transformation as the limit of a sequence of partial transformations, we would like to be able to ensure at some finite stage in the sequence that certain points have many upcrossings. The following notion is the analog of an upcrossing sequence for a partial transformation.

Definition 4.2. Let $T$ be a useful partial transformation, $A \subseteq 2^{\omega}$ and $\tau_{0}, \ldots, \tau_{n}$ an open loop in $T$ such that for each $i$ either $\left[\tau_{i}\right] \subseteq A$ or $\left[\tau_{i}\right] \subseteq$ $2^{\omega} \backslash A$. Let $R=\left\{i \leq n \mid\left[\tau_{i}\right] \subseteq A\right\}$ and $\alpha<\beta$ rationals. A $\tau_{s}$-upcrossing sequence for $\alpha, \beta$ of length $N$ is a sequence

$$
0 \leq u_{1}<v_{1}<u_{2}<v_{2}<\cdots<u_{N}<v_{N} \leq n
$$

such that for all $i \leq N$,

$$
\frac{1}{u_{i}+1} \sum_{j=s}^{u_{i}+s} \chi_{R}(j)<\alpha, \frac{1}{v_{i}+1} \sum_{j=s}^{v_{i}+s} \chi_{R}(j)>\beta .
$$

Note that we shift indices in the sums over by $s$, since we begin with $\tau_{s}$ and count to later elements of the open loop.

Lemma 4.3. Let $T$ be a useful partial transformation, let $A \subseteq 2^{\omega}$, and $\tau_{0}, \ldots, \tau_{n}$ an open loop in $T$ such that for each $i$ either $\left[\tau_{i}\right] \subseteq A$ or $\left[\tau_{i}\right] \subseteq$ $2^{\omega} \backslash A$. Let $0 \leq u_{1}<v_{1}<u_{2}<v_{2}<\cdots<u_{N}<v_{N} \leq n$ be a $\tau_{s}$-upcrossing sequence for $\alpha, \beta$.

Then whenever $\widehat{T}: 2^{\omega} \rightarrow 2^{\omega}$ extends $T$, and $x \in\left[\tau_{s}\right], u_{1}<v_{1}<\cdots<$ $u_{N}<v_{N}$ is an upcrossing sequence for $\alpha, \beta$ in $\widehat{T}$ with the function $\chi_{A}$.

We are finally ready to give our main construction.
Theorem 4.4. Suppose $x \in 2^{\omega}$ is not Martin-Löf random. Then there is a computable set $A$ and a computable transformation $\widehat{T}: 2^{\omega} \rightarrow 2^{\omega}$ such that $x$ is not typical with respect to the ergodic theorem.
Proof. Let $\left\langle V_{j}\right\rangle$ be a Martin-Löf test witnessing that $x$ is not Martin-Löf random, so $x \in \cap_{j}\left[V_{j}\right]$ and $\lambda\left(\left[V_{j}\right]\right) \leq 2^{-j}$. We write $V_{j}=\bigcup_{j, n} V_{j, n}$ where $(j, n) \mapsto V_{j, n}$ is computable. We will assume $V_{j, 0}=\emptyset$ for all $j$ and that $n<m$ implies $V_{j, n} \subseteq V_{j, m}$, and we refer to $V_{j, n+1} \backslash V_{j, n}$ as the portion of $V_{j}$ enumerated at stage $n+1$. It is convenient to assume that there
is at most one element enumerated into any $V_{j}$ at stage $n+1$; that is, $\sum_{j}\left|V_{j, n+1} \backslash V_{j, n}\right| \leq 1$, and so $\sum_{j}\left|V_{j, n}\right|$ is finite for any $n$.

We will construct an increasing sequence of useful partial transformations $T_{0} \sqsubseteq T_{1} \sqsubseteq T_{2} \sqsubseteq \cdots$ so that setting $T=\lim _{n} T_{n}, \widehat{T}(x)=\lim _{n \rightarrow \infty} T(x \upharpoonright n)$ is the desired transformation. We first list the technical requirements on our induction; since they are rather elaborate, we will go through what the intended meanings are before describing the actual construction.

Inductive Specification: We now specify the properties that will be maintained at each stage of our construction. Since they are rather complicated, we recommend that the reader skip them at first; a detailed explanation of their intended meaning is given after.

We will define, as part of our construction of stage 0 , a set $A$, a computable collection of finite sequences which we call componential, and a computable function $d$ defined on componential $\sigma$.

At stage $n$ we will have a partial transformation $T_{n}$, a partition $2^{\omega}=$ $W^{n} \cup \bigcup_{k}\left(A_{k}^{n} \cup B_{k}^{n}\right)$ into components which are finite unions of intervals, a further partition $W^{n}=\bigcup_{k} W_{k}^{n}$ into finite unions of intervals, constants $a_{k}^{n}, b_{k}^{n}$, a function $\rho^{n}$, and for each $i<n / 2$, a finite union of intervals $G_{i} \subseteq 2^{\omega}$, such that the following properties hold:
(1) Structure of components
(1.a) $T_{n}$ is useful,
(1.b) $T_{n} \sqsubseteq T_{n+1}$,
(1.c) $T_{n,+} \subseteq W^{n}$,
(1.d) Each open loop in $T_{n}$ belongs entirely to one component,
(1.e) Each $W_{k}^{n}$ is a union of intervals of the form $[\tau]$ with $\tau$ determined,
(1.f) If $\sigma$ is componential and determined then $[\sigma]$ is contained in a single component in $T_{n}$,
(1.g) $A_{k}^{n} \subseteq A_{k}^{n+1} \cup W_{0}^{n+1}$,
(1.h) $B_{k}^{n} \subseteq B_{k}^{n+1} \cup W_{0}^{n+1}$,
(1.i) $W_{k}^{n} \subseteq W_{k}^{n+1} \cup W_{k+1}^{n+1}$,
(1.j) For each $k, T_{n}$ is $A_{k}^{n}, A_{k}^{n}$-escapable,
(1.k) For each $k, T_{n}$ is $B_{k}^{n}, B_{k}^{n}$-escapable,
(1.1) $T_{n}$ is $W^{n}, W^{n}$-escapable,
(2) Management of upcrossings
(2.a) If $\sigma \in W_{k}^{n}$ is determined then in the open loop in $T_{n}$ containing $\sigma$, there is an upcrossing sequence for $\sigma$ of length $k$ for $2^{\omega} \backslash A$ and $1 / 2,3 / 4$,
(2.b) The domain of $\rho^{n}$ is the elements of $W^{n}$ determined in $T_{n}$,
(2.c) If $\sigma \sqsubseteq \tau$ are both determined and in $W^{n}$ then $\rho^{n}(\sigma)=\rho^{n}(\tau)$,
(2.d) If $\rho^{n}(\sigma)=\rho^{n}(\tau)$ and there is a $k$ with $\sigma, \tau \in W_{k}^{n}$ then $L_{T_{n}}(\sigma)=$ $L_{T_{n}}(\tau)$,
(2.e) If $\sigma \in W_{k}^{n+1} \cap W_{k}^{n}$ is determined in $T_{n+1}$ then
(2.e.i) $\rho^{n+1}(\sigma)=\rho^{n}(\sigma)$, (2.e.ii) $L_{T_{n+1}}(\sigma)=L_{T_{n}}(\sigma)$,
(2.f) For each $k, 0 \leq a_{k}^{n}<\lambda\left(A_{k}^{n}\right)$ and $0 \leq b_{k}^{n}<\lambda\left(B_{k}^{n}\right)$,
(2.g) For each $k$, let $J_{k}^{n}$ be the image of $\rho^{n} \upharpoonright W_{k}^{n}$, and for each $j \in J_{k}^{n}$, let $l_{k, j}$ be the (by (2.d), necessarily unique) value of $L_{T_{n}}(\tau)$ for some $\tau \in W_{k}^{n}$; then
(2.g.i) $\sum_{j \in J_{k}^{n}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j,(n-1) / 2}\right]\right)\right) \leq a_{k}^{n}$,
(2.g.ii) $4 \sum_{j \in J_{k}^{n}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j,(n-1) / 2}\right]\right)\right) \leq b_{k}^{n}$,
(2.h) If $\sigma \in W^{n}$ is determined in $T_{n}$ then $\left[V_{\rho^{n}(\sigma),(n-1) / 2}\right] \cap[\sigma]=\emptyset$,
(3) Almost everywhere defined
(3.a) For each $i<n / 2$ and each $\sigma \in T_{n,+} \cup T_{n,-}, \sigma$ is either in $G_{i}$ or avoids $G_{i}$,
(3.b) For any $i_{0}, \ldots, i_{k-1}, \lambda\left(\bigcap_{j \leq k} G_{i_{j}}\right) \leq 2^{-k}$,
(3.c) If $\sigma \in T_{n,-}$ then $|\sigma| \geq n / 2$,
(3.d) For each componential $\sigma \in T_{n,-},\left|T_{n}(\sigma)\right| \geq \mid\left\{i<n / 2 \mid \sigma\right.$ avoids $\left.G_{i}\right\} \mid-$ $d(\sigma)$.

Explanation of inductive clauses: Initially we will fix a partition into three regions, $W, A, B$. $W$ will be a region known to contain $x$, and $W \cup B$ will be the set demonstrating the failure of the ergodic theorem for $x$. Our strategy will be that when we enumerate some $\tau$ into $V_{j}$ for an appropriate $j$, we will arrange to add an upcrossing by first mapping every element of $[\tau]$ through $A$ for a long time, ensuring that the ergodic average falls below $1 / 2$. We will then have the transformation map those elements through $B$ for a long time to bring the average back up to $3 / 4$. We will do this to each element of $\cap_{j}\left[V_{j}\right]$ infinitely many times, ensuring that elements in this intersection are not typical. ${ }^{2}$

We will have further partitions $A=\bigcup A_{k}$ and $B=\bigcup B_{k} ; A_{k}$ is the portion of $A$ reserved for creating the $k+1$-st upcrossing, and $B_{k}$ the portion of $B$ reserved for the same. We will start with $W_{0}^{0}=W, A_{k}^{0}=A_{k}$, and $B_{k}^{0}=B_{k}$. At later stages, $W_{k}^{n}$ will be the portion of $W^{n}$ known to have at least $k$ upcrossings (2.a). At a given step, we might move intervals from $W_{k}^{n}$ to $W_{k+1}^{n+1}$ (1.i) (because we have created a new upcrossing) and we might move intervals from $A_{k}^{n}$ to $W_{0}^{n+1}$ (because it is part of a newly created upcrossing) and similarly for $B_{k}^{n}(1 . \mathrm{g})(1 . \mathrm{h})$. In any other case, an interval in $W_{k}^{n}$ is in $W_{k}^{n+1}$ and similarly for $A, B$. (It is perhaps slightly confusing that we will use the term "component" to mean either $W^{n}$ or $A_{k}^{n}$ or $B_{k}^{n}$ for some $k$, but that the $W_{k}^{n}$ are not themselves components. However it will become clear that most properties - the behavior of loops and escapability, for instance respect components in the sense in which we are using the term.)

When points are in $A_{k}^{n}$ or $B_{k}^{n}$, they are simply waiting to (possibly) be used in the creation of an upcrossing, so they always belong to open loops

[^2]with length 1 (1.c). All open loops with length longer than 1 wil be entirely within $W^{n}$ (1.d), Similarly, escapability is handled componentwise $(1 . j) \|(1.1)$ Combining these properties, the point is that there is simply no interaction between separate components except when we create upcrossings, and when we create an upcrossing, we put everything in $W^{n}$.

Suppose $\tau$ is enumerated into $V_{j}$, causing us to create a new upcrossing. We now need to look for extensions of $\tau$ to be enumerated into $\bigcup_{j^{\prime}>j} V_{j^{\prime}}$; specifically, we should choose a particular $j^{\prime}>j$ and watch for an extension of $\tau$ to be enumerated into $V_{j^{\prime}}$. We cannot simply choose $j^{\prime}=j+1$, since we may not have enough measure available in $A_{k}^{n}$. Instead we choose a new value for $\rho^{n}(\tau)$, and we will wait for extensions of $\tau$ in $V_{\rho^{n}(\tau)} \cdot(2 . \mathrm{h})$ ensures that if we see some extension of $\tau$ get enumerated into $V_{\rho^{n}(\tau)}$, we will be forced to create a new upcrossing, and conversely (2.e.i) ensures that this is the only time we change $\rho^{n}(\tau)$.

When we assign $\rho^{n}(\tau)=j^{\prime}$, we have to consider the worst case scenario: $\lambda\left(\left[V_{j^{\prime}}\right]\right)=2^{j^{\prime}}$ and $\left[V_{j^{\prime}}\right] \subseteq[\tau]$. We therefore need enough measure in $A_{k}^{n}$ and $B_{k}^{n}$ to create the needed upcrossings. The amount of measure we needed is determined not only by the size of [ $V_{j^{\prime}}$ ], but also by the length $L_{T_{n}}(\tau)$ of the loop containing $\tau$. To simplify the calculations, we require that all sequences sharing a value of $\rho^{n}$ have the same length $L_{T_{n}}$ (2.d) and that this value does not change except when we create upcrossings (2.e.ii).

The values $a_{k}^{n}$ and $b_{k}^{n}$ keep track of the assigned portions of $A_{k}^{n}$ and $B_{k}^{n}$ respectively. We should never assign all of $A_{k}^{n}, B_{k}^{n}$ because we need some space for upcrossings (2.f). When we assign $\rho^{n}(\tau)=j^{\prime}$, we will increase $a_{k}^{n}$ and $b_{k}^{n}$ accordingly, and when we create an upcrossing we will use some of $A_{k}^{n}, B_{k}^{n}$, move the corresponding intervals into $W^{n+1}$, and decrease $a_{k}^{n+1}, b_{k}^{n+1}$ accordingly (2.g.i) (2.g.ii).

Finally, we need to make sure that the transformation is defined almost everywhere. (3.d) ensures that if $\sigma$ avoids many of the sets $G_{i}$ then $T_{n}(\sigma)$ is large, and (3.b) ensures that most $\sigma$ avoid many $G_{i}$. A technicality is that we want to extend $T_{n}$ while respecting the fact that escape sequences are supposed to remain within components. We will call $\sigma$ componential if $[\sigma]$ belongs to a single component in $T_{0}$; while there are infinitely many minimal componential elements (each $A_{k}^{0}$, for instance, instance, is an interval $[\tau]$, so $\tau$ is componential but no initial segment is), every $x$ belongs to some $[\tau]$ with $\tau$ componential. $d(\sigma)$ will be the length of the smallest $\tau \sqsubseteq \sigma$ such that $\tau$ is componential. We will first make sure that $T_{n,-}$ is defined on longer and longer sequences (3.c), and then once we reach a componential portion (after at most $d(\sigma)$ steps), we will start making $T_{n}$ longer and longer.

Stage 0: We now give the initial stage of our construction. We assume without loss of generality that $x$ belongs to some interval $W^{0}$ with $\lambda\left(W^{0}\right)<$ 1. For instance, we can suppose we know the first bit of $x$ and let $W^{0}$ equal $[\langle 0\rangle]$ or $[\langle 1\rangle]$ as appropriate. Then, from the remaining measure, we take $A_{k}^{0}$ and $B_{k}^{0}$ to be intervals so that $\lambda\left(B_{k}^{0}\right)=4 \lambda\left(A_{k}^{0}\right)$ for all $k$. We set $W_{0}^{0}=W^{0}$
and take $T_{0}$ to be the trivial transformation (i.e., $T_{0}(\sigma)=\langle \rangle$ for all $\sigma$ ). Choose $j$ large enough that $2^{-j}<\lambda\left(A_{0}^{0}\right)$ and set $\rho^{0}(\sigma)=j$ for all $\sigma$. Set $a_{0}^{0}=2^{-j}, b_{0}^{0}=4 \cdot 2^{-j}$, and for $k>0, a_{k}^{0}=b_{k}^{0}=0$.

We define $\sigma$ to be componential if either $[\sigma] \subseteq W^{0}$ or there is some $k$ such that $[\sigma] \subseteq A_{k}^{0}$ or $[\sigma] \subseteq B_{k}^{0}$. For any componential $\sigma, d(\sigma)$ is the length of the smallest $\tau \sqsubseteq \sigma$ such that $\tau$ is componential.

Odd stages-ensuring $T$ is total: Given $T_{n}$ for an even $n$, we take steps at the odd stage $n+1$ to make sure that $T$ is defined almost everywhere. We will define $T_{n+1,+}=T_{n,+}$ and

$$
T_{n+1,-}=\cup_{\sigma \in T_{n,-}}\left\{\sigma^{\frown}\langle 0\rangle, \sigma^{\frown}\langle 1\rangle\right\} .
$$

Consider some $\sigma \in T_{n,-}$. We set $T_{n+1}\left(\sigma^{\frown}\langle 0\rangle\right)=T_{n}(\sigma)$. If $[\sigma]$ is not componential then $T_{n+1}\left(\sigma^{\frown}\langle 1\rangle\right)=T_{n}(\sigma)$ as well.

If $\sigma$ is componential, let $\tau$ be the initial element of the open loop containing $\sigma$, so $T_{n}(\sigma) \sqsubset \tau$. In particular, there is a $b \in\{0,1\}$ such that $T_{n}(\sigma) \frown\langle b\rangle \sqsubseteq \tau$, and we set $T_{n+1}\left(\sigma^{\frown}\langle 1\rangle\right)=T_{n}(\sigma) \frown\langle b\rangle$. We define $W^{n+1}=$ $W^{n}, A_{k}^{n+1}=A_{k}^{n}, B_{k}^{n+1}=B_{k}^{n}, a_{k}^{n+1}=a_{k}^{n}, b_{k}^{n+1}=b_{k}^{n}$ for all $k$. We define $\rho^{n+1}(\sigma)=\rho^{n}(\sigma)$ for determined $\sigma \in W^{n}$. We define $G_{n / 2}=\bigcup_{\sigma \in T_{n,-}}\left[\sigma^{\frown}\langle 0\rangle\right]$.

Most of the inductive properties follow immediately from the inductive hypothesis since we changed neither the components nor the lengths of any open loops. To check (1.j) $-(1.1)$, it suffices to show that whenever $\sigma \in T_{n,-}$ and $b \in\{0,1\}$, there is an escape sequence for $\sigma_{0}\langle b\rangle$ in $T_{n+1}$ belonging to the same component as $\sigma_{0}$. Given an escape sequence $\sigma_{1}, \ldots, \sigma_{k}$ for $\sigma_{0}$ belonging to the same component, the only potential obstacle is that one or more $\sigma_{i}$ has the form $\sigma_{i}^{-} \frown\langle 1\rangle \frown \rho$ where $\sigma^{-} \in T_{n,-}$. If $\left[\sigma_{i}^{-}\right]$is not componential, there is no obstacle. If $\sigma_{i}^{-}$is componential then everything in $\left[\sigma_{i}^{-}\right]$must belong to the same component; let $\tau_{0}, \ldots, \tau_{r}, \sigma_{i}^{-}$be the open loop containing $\sigma_{i}^{-}$in $T_{n}$, and observe that

$$
\sigma_{1}, \ldots, \sigma_{i}^{-} \frown\langle 1\rangle, \tau_{0} \frown\langle 1\rangle \frown \rho, \ldots, \tau_{r} \frown\langle 0\rangle \frown \rho, \sigma_{i+1}, \ldots, \sigma_{k}
$$

is an escape sequence for $\sigma_{0}$ in $T_{n}$, and since $\left[\sigma_{i}^{-}\right]$belongs to this component and each open loop is contained in one component, this new escape sequence is contained entirely in the component of $\sigma_{0}$. After we have replaced each
 is also an escape sequence in $T_{n+1}$ contained entirely in one component.
(3.a) holds for $i=n / 2$ by construction. To see (3.b) consider the sequence $i_{0}, \ldots, i_{k-1}, n / 2$. By the inductive hypothesis, $\lambda\left(\bigcap_{j \leq k} G_{i_{j}}\right) \leq 2^{-k}$. Each $\sigma \in T_{n,+} \cup T_{n,-}$ is either in $\bigcap_{j \leq k} G_{i_{j}}$ or avoids $\bigcap_{j \leq k} G_{i_{j}}$, so $\bigcap_{j \leq k} G_{i_{j}} \cap G_{n / 2}$ is exactly the union of those $\left[\sigma^{\frown}\langle 0\rangle\right]$ such that $\sigma$ is in $\bigcap_{j \leq k} G_{i_{j}}$. In particular, whenever $\left[\sigma^{\frown}\langle 0\rangle\right] \subseteq \bigcap_{j \leq k} G_{i_{j}} \cap G_{n / 2}$, we must have $[\sigma] \subseteq \bigcap_{j \leq k} G_{i_{j}}$, so $\lambda\left(\bigcap_{j \leq k} G_{i_{j}} \cap G_{n / 2}\right) \leq 2^{-1} \lambda\left(\bigcap_{j \leq k} G_{i_{j}}\right) \leq 2^{-k-1}$.

If $\sigma \frown\langle b\rangle \in T_{n+1,-}$ then $\sigma \in T_{n,-}$ and by the inductive hypothesis, $|\sigma| \geq$ $n / 2$, so $\left|\sigma^{\frown}\langle b\rangle\right|=|\sigma|+1 \geq(n+1) / 2$ since $n$ is even, showing (3.c).

To see (3.d), consider some componential $\sigma^{\complement}\langle b\rangle \in T_{n+1,-}$. If $\sigma$ was componential then we have $\sigma \in T_{n,-}$ and by the inductive hypothesis, $\left|T_{n}(\sigma)\right| \geq \mid\left\{i<n / 2 \mid \sigma\right.$ avoids $\left.G_{i}\right\} \mid-d(\sigma)$. We have

$$
\begin{aligned}
\left|T_{n+1}\left(\sigma^{\frown}\langle 0\rangle\right)\right| & =\left|T_{n}(\sigma)\right| \\
& \geq \mid\left\{i<n / 2 \mid \sigma \text { avoids } G_{i}\right\} \mid-d(\sigma) \\
& =\mid\left\{i<(n+1) / 2 \mid \sigma^{\frown}\langle 0\rangle \text { avoids } G_{i}\right\} \mid-d\left(\sigma^{\frown}\langle 0\rangle\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\left|T_{n+1}\left(\sigma^{\frown}\langle 1\rangle\right)\right| & =\left|T_{n}(\sigma)\right|+1 \\
& \geq \mid\left\{i<n / 2 \mid \sigma \text { avoids } G_{i}\right\} \mid-d(\sigma)+1 \\
& =\mid\left\{i<(n+1) / 2 \mid \sigma^{\frown}\langle 1\rangle \text { avoids } G_{i}\right\} \mid-d\left(\sigma^{\frown}\langle 1\rangle\right) .
\end{aligned}
$$

If $\sigma$ was not componential then $d\left(\sigma^{\circ}\langle b\rangle\right)=\left|\sigma^{\circ}\langle b\rangle\right| \geq(n+1) / 2$, so

$$
\left|T_{n+1}\left(\sigma^{\frown}\langle b\rangle\right)\right| \geq 0 \geq \mid\left\{i<(n+1) / 2 \mid \sigma^{\frown}\langle b\rangle \text { avoids } G_{i}\right\} \mid-d\left(\sigma^{\frown}\langle b\rangle\right) .
$$

Even stages-creating upcrossings; the construction: We now consider the real work. When $n=2 n^{\prime}+1$ is odd, we take steps at the even stage $n+1$ to ensure that the ergodic theorem does not hold for any element of $\cap_{j}\left[V_{j}\right]$ by adding an additional upcrossing to some element enumerated into one of the $V_{j}$ s. As noted above, we assume that there is exactly one $\pi$ in $\bigcup_{j} V_{j, n^{\prime}} \backslash \bigcup_{j} \bigcup_{m<n^{\prime}} V_{j, n^{\prime}}$, so let $\hat{j}$ be such that $\pi \in V_{\hat{j}, n^{\prime}} \backslash \bigcup_{m<n^{\prime}} V_{\hat{j}, m}$. If $\pi$ is not determined, we may partition $[\pi]=\bigcup_{i}\left[\pi_{i}\right]$ where each $\pi_{i}$ is determined, and we will deal with each $\pi_{i}$ separately.

We first consider how to deal with a single element. Let $\tau=\pi$ if $\pi$ is determined or $\tau=\pi_{i}$ for some $i$ is $\pi$ is not determined. We will construct a transformation $T_{n+1}$ and other objects ( $W_{k}^{n+1}$, etc.) satisfying all the conditions above other than (2.h). (We will discuss (2.h) after.) If $\rho^{n}(\tau) \neq \hat{j}$, we will just take $T_{n+1}=T_{n}$ and similarly for all other objects.

So we consider the case where $\rho^{n}(\tau)=\hat{j}$. We have $\tau \in W_{k}^{n}$ for some $k$. We will extend the open loop containing $\tau$ to a longer open loop which will contain a long subsequence coming from $A_{k}^{n}$ and then an even longer subsequence coming from $B_{k}^{n}$, thereby creating a new upcrossing in $T_{n+1}$ for each element of $[\tau]$. We illustrate the intended arrangement in Figure 4. We must ensure that every point in $[\tau]$, a section of fixed total measure, receives a new upcrossing. We must do so while ensuring that the total measure of the portions of $A_{k}^{n}$ used is strictly less than $L_{T_{n}}(\tau) \lambda([\tau])+\left(\lambda\left(A_{k}^{n}\right)-\right.$ $a_{k}^{n}$ ) and the total measure of the portions of $B_{k}^{n}$ used is strictly less than $4 L_{T_{n}}(\tau) \lambda([\tau])+\left(\lambda\left(B_{k}^{n}\right)-b_{k}^{n}\right)$. Finally, the escape sequences will all have a fixed height, which we cannot expect to bound in advance. Our solution will be to thin all the parts other than the escape sequences until the entire tower is so narrow that we can afford the error introduced by the various escape sequences.

We now fix the following:

- Let $\tau_{1}, \ldots, \tau_{t}$ be the open loop containing $\tau$,
- Identify a length $e^{t}$ such that there is a reduced escape sequence for $\tau_{t}$ of length $e^{t}$ contained in $W^{n}$,
- For some finite $V$ and each $i \leq V$, let $\nu_{i} \in B_{k}^{n}$ be such that the $\nu_{i}$ are pairwise distinct, prefix-free, each $\nu_{i}$ is determined, and

$$
\lambda\left(B_{k}^{n}\right)>\sum_{i \leq V} \lambda\left(\left[\nu_{i}\right]\right)>4 \sum_{i \leq t} \lambda\left(\left[\tau_{i}\right]\right)=4 L_{T_{n}}(\tau) \lambda([\tau]) .
$$

Such $\nu_{i}$ must exist because $\lambda\left(B_{k}^{n}\right)>b_{k}^{n} \geq 4 L_{T_{n}}(\tau) \lambda([\tau])$.

- For some finite $U$ and each $i \leq U$, let $v_{i} \in A_{k}^{n}$ be such that the $v_{i}$ are pairwise distinct, prefix-free, each $v_{i}$ is determined, and

$$
\min \left\{\lambda\left(A_{k}^{n}\right), \frac{1}{4} \sum_{i \leq V} \lambda\left(\left[\nu_{i}\right]\right)\right\}>\sum_{i \leq U} \lambda\left(\left[v_{i}\right]\right)>\sum_{i \leq t} \lambda\left(\left[\tau_{i}\right]\right)=L_{T_{n}}(\tau) \lambda([\tau]) .
$$

Such $v_{i}$ must exist because $\lambda\left(A_{k}^{n}\right)>a_{k}^{n} \geq L_{T_{n}}(\tau) \lambda([\tau])$.

- For each $i \leq U$, identify an $e_{i}^{u}$ such that there is a reduced escape sequence for $v_{i}$ of length $e_{i}^{u}$ contained entirely in $A_{k}^{n}$, and set $e_{u}=$ $\sum_{i \leq U} e_{i}^{u}$,
- For each $i \leq V$, identify an $e_{i}^{v}$ such that there is a reduced escape sequence for $\nu_{i}$ of length $e_{i}^{v}$ contained entirely in $B_{k}^{n}$, and set $e_{v}=$ $\sum_{i \leq V} e_{i}^{v}$,
Choose $N, N^{\prime}$ sufficiently large (as determined by the argument to follow). We apply Lemma 3.2 to the open loop $\tau_{1}, \ldots, \tau_{t}$ in $T_{n}$ with $\epsilon=2^{-\left(N+N^{\prime}\right)}$ to obtain a partial transformation $T_{n}^{0,0}$ where the width of the open loop containing $\tau$ is $\epsilon$.

Note that each $v_{i}$ is an open loop in its own right, and similarly for each $\nu_{i}$ (since an open loop containing $v_{i}$ consists entirely of elements in $A_{k}^{n}$, and no element of $A_{k}^{n}$ belongs to $T_{n,+}$ ).

Given $T_{n}^{0, i}$, let $T_{n}^{0, i+1}$ be the result of applying Lemma 3.4 to the open loop $v_{i}$ in $T_{n}^{0, i}$ with $\epsilon=2^{-N}$. In $T_{n}^{0, i+1}$ there is a open loop $v_{i, 0}^{\prime}, \ldots, v_{i, u_{i}^{\prime}}^{\prime}$ with width $2^{-N}, \lambda\left(\left[v_{i}\right] \backslash \bigcup_{j \leq u_{i}^{\prime}}\left[v_{i, j}^{\prime}\right]\right)=2^{-N}$.

Let $T_{n}^{1,0}=T_{n}^{0, U+1}$; given $T_{n}^{1, i}$, let $T^{1, i+1}$ be the result of applying Lemma 3.2 to the open loop $v_{i, 0}^{\prime}, \ldots, v_{i, u_{i}^{\prime}}^{\prime}$ in $T_{n}^{1, i}$ with $\epsilon=2^{-\left(N+N^{\prime}\right)}$. In $T_{n}^{1, i+1}$ there is an open loop $v_{i, 0}^{\prime \prime}, \ldots, v_{i, u_{i}^{\prime \prime}}^{\prime \prime}$ with width $2^{-\left(N+N^{\prime}\right)}$ and $\lambda\left(\left[v_{i}\right] \backslash \bigcup_{j \leq u_{i}^{\prime \prime}}\left[v_{i, j}^{\prime \prime}\right]\right)=$ $2^{-N}$.

Let $T_{n}^{2,0}=T_{n}^{1, U+1}$; given $T_{n}^{2, i}$, let $T_{n}^{2, i+1}$ be the result of applying Lemma 3.4 to the open loop $\nu_{i}$ in $T_{n}^{0, i}$ with $\epsilon=2^{-N}$. In $T_{n}^{2, i+1}$ there is an open loop $\nu_{i, 0}^{\prime}, \ldots, \nu_{i, v_{i}^{\prime}}^{\prime}$ with width $2^{-N}$ and $\lambda\left(\left[\nu_{i}\right] \backslash \bigcup_{j \leq v_{i}^{\prime}}\left[\nu_{i, j}^{\prime}\right]\right)=2^{-N}$.

Let $T_{n}^{3,0}=T_{n}^{2, V+1}$; given $T_{n}^{3, i}$, let $T_{n}^{3, i+1}$ be the result of applying Lemma 3.2 to the open loop $\nu_{i, 0}^{\prime}, \ldots, \nu_{i, v_{i}^{\prime}}^{\prime}$ in $T_{n}^{3, i}$ with $\epsilon=2^{-\left(N+N^{\prime}\right)}$. In $T_{n}^{3, i+1}$ there is an open loop $\nu_{i, 0}^{\prime \prime}, \ldots, \nu_{i, v_{i}^{\prime \prime}}^{\prime}$ with width $2^{-\left(N+N^{\prime}\right)}$ and $\lambda\left(\left[\nu_{i}\right] \backslash \bigcup_{j \leq v_{i}^{\prime \prime}}\left[\nu_{i, j}^{\prime \prime}\right]\right)=$ $2^{-N}$.


Figure 4. Construction of $T^{\prime}$
We set $T_{n}^{4,0}=T_{n}^{3, V+1}$. For each $i \leq U$, we choose an escape sequence for $v_{i, u_{i}^{\prime \prime}}^{\prime \prime}, \eta_{1}^{i}, \ldots, \eta_{e_{i}^{u}}^{i}$ contained entirely in $A_{k}^{n} \backslash \bigcup_{i \leq U, j \leq u_{i}^{\prime \prime}}\left[v_{i, j}^{\prime \prime}\right]$. To see that such an escape sequence exists, recall that an escape sequence for $v_{i}$ in $A_{k}^{n}$ of the right length existed $T_{n}$, and since $v_{i, u_{i}^{\prime}}^{\prime \prime} \sqsupseteq v_{i}$, the escape sequence was an escape sequence for $v_{i, u_{i}^{\prime}}^{\prime \prime}$ as well. By Lemma 3.4 and Remark 3.5 the desired escape sequence existed in $T_{n}^{1,0}$, and its existence was preserved by the remaining steps. Further, if $N$ and $N^{\prime}$ were chosen large enough, we may ensure that the collection of escape sequences $\left\{\bigcup_{j \leq e_{i}^{u}}\left[\eta_{j}^{i}\right]\right\}$ is pairwise disjoint.

Given $T_{n}^{4, i}$, let $T_{n}^{4, i+1}$ be the result of applying Lemma 3.6 to the escape sequence $\eta_{1}^{i}, \ldots, \eta_{e_{i}^{u}}^{i}$ in $T_{n}^{4, i}$.

Let $T_{n}^{5,0}=T_{n}^{4, U+1}$. For each $i \leq V$, we choose an escape sequence for $\nu_{i, v_{i}^{\prime \prime}}^{\prime \prime}, \theta_{1}^{i}, \ldots, \theta_{e_{i}^{e}}^{i}$. These escape sequences exist for the same reason as above. Given $T_{n}^{5, i}$, let $T_{n}^{5, i+1}$ be the result of applying Lemma 3.6 to the escape sequence $\nu_{i, v_{i}^{\prime \prime}}^{\prime \prime}, \theta_{1}^{i}, \ldots, \theta_{e_{i}^{v}}^{i}$ in $T_{n}^{5, i}$.

Let $T_{n}^{6}=T_{n}^{5, V+1}$. Choose an escape sequence $\xi_{1}, \ldots, \xi_{e^{t}}$ for $\tau_{t}$. This exists for the same reason as above. Let $T_{n}^{7}$ be the result of applying Lemma 3.6 to the escape sequence $\tau_{t}, \xi_{1}, \ldots, \xi_{e^{t}}$ in $T_{n}^{6}$.
$\bigcup_{i \leq t}\left[\tau_{i}\right]$ is contained in an open loop in $T_{n}^{7}$ whose final element is $\xi_{e^{t}}\ulcorner\langle 0\rangle$. For each $i \leq U,\left[v_{i}\right]$ is almost contained (except for a portion of measure $\leq e_{u} 2^{-N}$ ) in an open loop with initial element $v_{i, 0}^{\prime \prime}$ and final element $\eta_{e_{i}^{u}}^{i} \sim\langle 0\rangle$.

For each $i \leq V,\left[\nu_{i}\right]$ is almost contained (except for a portion of measure $\leq e_{v} 2^{-N}$ ) in an open loop with initial element $\nu_{i, 0}^{\prime \prime}$ and final element $\theta_{e_{i}^{v}}^{i} \frown\langle 0\rangle$. We define $T_{n+1}$ by taking

$$
T_{n+1,+}=T_{n,+}^{7} \cup\left\{\xi_{e^{t}} \frown\langle 0\rangle\right\} \cup\left\{\eta_{e_{i}^{u}}^{i} \frown\langle 0\rangle \mid i \leq U\right\} \cup\left\{\theta_{e_{i}^{v}}^{i} \frown\langle 0\rangle \mid i<V\right\}
$$

We define $T_{n+1}\left(\xi_{e^{t}} \frown\langle 0\rangle\right)=v_{0,0}^{\prime \prime}$; for $i<U, T_{n+1}\left(\eta_{e_{i}^{u}}^{i} \frown\langle 0\rangle\right)=v_{i+1,0}^{\prime \prime} ; T_{n+1}\left(\eta_{e_{U}^{u}}^{U} \frown\langle 0\rangle\right)=$ $\nu_{0,0}^{\prime \prime}$; for $i<V, T_{n+1}\left(\theta_{e_{i}^{v}}^{i} \frown\langle 0\rangle\right)=\nu_{i+1,0}^{\prime \prime}$.

The conclusion of all this is that in $T_{n+1}$, we have a large open loop which consists exactly of

$$
\begin{aligned}
S & =\bigcup_{i \leq t}\left[\tau_{i}\right] \\
& \cup \bigcup_{i \leq e^{t}}\left[\xi_{i} \frown\langle 0\rangle\right] \\
& \cup \bigcup_{i \leq U}\left[\bigcup_{j \leq u_{i}^{\prime \prime}}\left[v_{i, j}^{\prime \prime}\right] \cup \bigcup_{j \leq e_{i}^{u}}\left[\eta_{j}^{i \frown}\langle 0\rangle\right]\right] \\
& \cup \bigcup_{i \leq V}\left[\bigcup_{j \leq v_{i}^{\prime \prime}}\left[\nu_{i, j}^{\prime \prime}\right] \cup \bigcup_{j \leq e_{i}^{v}}\left[\theta_{j}^{i \frown}\langle 0\rangle\right]\right]
\end{aligned}
$$

If $\sigma$ is not in $S$ then $T_{n+1}(\sigma)=T_{n}(\sigma)$.
For $k^{\prime} \neq k$, we set $A_{k^{\prime}}^{n+1}=A_{k^{\prime}}^{n}, B_{k^{\prime}}^{n+1}=B_{k^{\prime}}^{n}$. We set

$$
A_{k}^{n+1}=A_{k}^{n} \backslash S
$$

and

$$
B_{k}^{n+1}=B_{k}^{n} \backslash S
$$

For each $k^{\prime}>0$, we let

$$
W_{k^{\prime}+1}^{n+1}=\left(W_{k^{\prime}+1}^{n} \backslash S\right) \cup\left(W_{k^{\prime}}^{n} \cap S\right)
$$

and

$$
W_{0}^{n+1}=\left(W_{0}^{n} \backslash S\right) \cup\left(A_{k}^{n} \cap S\right) \cup\left(B_{k}^{n} \cap S\right)
$$

We need to define $\rho^{n+1}$. If $\sigma \in W^{n+1} \backslash S$ is determined then set $\rho^{n+1}(\sigma)=$ $\rho^{n}(\sigma)$. For each $k^{\prime}$ such that $W_{k^{\prime}}^{n+1} \nsubseteq W_{k^{\prime}}^{n}$, choose some $j_{k^{\prime}}$ not in the image of $\rho^{n}$ so that $2^{-j_{k^{\prime}}}$ is very small relative to $\lambda\left(A_{k^{\prime}}^{n}\right)-a_{k^{\prime}}^{n}$ and $\lambda\left(B_{k^{\prime}}^{n}\right)-b_{k^{\prime}}^{n}$. For each determined $\sigma \in S \cap W_{k^{\prime}}^{n+1}$, set $\rho^{n+1}(\sigma)=j_{k^{\prime}}$. Let $L$ be the length of $S$. For each $k^{\prime}$, if $W_{k^{\prime}}^{n+1} \nsubseteq W_{k^{\prime}}^{n}$, set $\delta_{k^{\prime}}=L \cdot 2^{-j_{k^{\prime}}}$, and otherwise set $\delta_{k^{\prime}}=0$. For $k \neq k^{\prime}$, define

$$
a_{k^{\prime}}^{n+1}=a_{k^{\prime}}^{n}+\delta_{k^{\prime}}, b_{k^{\prime}}^{n+1}=b_{k^{\prime}}^{n}+4 \delta_{k^{\prime}}
$$

Define

$$
a_{k}^{n+1}=a_{k}^{n}-L_{T_{n}}(\tau) \lambda([\tau])+\delta_{k}
$$

and

$$
b_{k}^{n+1}=b_{k}^{n}-4 L_{T_{n}}(\tau) \lambda([\tau])+4 \delta_{k}
$$

Even stages-creating upcrossings; the verification: We now check that $T_{n+1}$ satisfies the inductive conditions.

For (1.a) $T_{n}^{7}$ is useful since it was produced by a series of applications of Lemma 3.4. Lemma 3.2, and Lemma 3.6. $T_{n+1}$ is proper by construction. It is partitioned into open loops since either $\sigma \notin S$, in which case $\sigma$ belongs to the same open loop it did in $T_{n}^{7}$, or $\sigma \in S$, in which case $S$ is the open loop containing $\sigma$. We will check escapability below, when we check conditions (1.j) $)(1.1)$.
(1.d) is immediate each open loop in $T_{n+1}$ is either an open loop from $T_{n}$ which remains in the same component, or the open loop $S$, which is in $W^{n+1}$. (1.e) is immediate from the inductive hypothesis and the construction. If $\sigma$ is componential and determined in $T_{n+1}$ then either $\sigma$ was in $W^{0}$, so also in $W^{n+1}$, or in some $A_{k}^{0}$ or $B_{k}^{0}$. If $\sigma$ is not in $A_{k}^{n+1}$ or $B_{k}^{n+1}$, respectively, it must be because some portion of $[\sigma]$ is in $S$, and so was moved to $W^{n+1}$. But $S$ is a union of intervals in $T_{n+1,+} \cup T_{n+1,-}$, so if $\sigma$ is determined and $[\sigma] \cap S \neq \emptyset$ then $[\sigma] \subseteq S$, so $\sigma$ is in $W^{n+1}$. In either case, $\sigma$ belongs to a single component, showing (1.f) (1.g) (1.i) are immediate from the definition.

For $k^{\prime} \neq k$, each $T_{n}^{i, j}$ and $T_{n}^{i}$ is $A_{k^{\prime}}^{n}, A_{k^{\prime}}^{n}$ and $B_{k^{\prime}}^{n}, B_{k^{\prime}}^{n}$-escapable, and $T_{n+1}$ is as well, since this property is preserved by each step of the construction. $T_{n}$ was $A_{k}^{n}, A_{k}^{n}$-escapable, and therefore $A_{k}^{n+1}, A_{k}^{n}$-escapable (since $A_{k}^{n+1} \subseteq A_{k}^{n}$ ), and Lemma 3.4 ensures that $T_{n}^{1,0}$ is $A_{k}^{n+1}, A_{k}^{n+1}$-escapable. This is preserved by each remaining step, so $T_{n+1}$ is $A_{k}^{n+1}, A_{k}^{n+1}$-escapable. A similar argument shows that $T_{n}^{3,0}$ is $B_{k}^{n+1}, B_{k}^{n+1}$-escapable, and so $T_{n+1}$ is as well. This shows (1.j) and (1.k).
$T_{n}^{7}$ is $W^{n}, W^{n}$-escapable since $T_{n}$ was and this property is preserved by each step. If $\sigma \in W^{n+1} \backslash W^{n}$ with $\left|T_{n+1}(\sigma)\right|<|\sigma|$ then $\sigma \sqsupseteq \theta_{e_{V}^{v}}^{V} \frown\langle 0\rangle$, and therefore $T_{n+1}(\sigma)=\langle \rangle$. Otherwise $\sigma \in W^{n}$, so let $\sigma_{0}, \ldots, \sigma_{r}$ be an escape sequence in $T_{n}^{7}$ in $W^{n}$ for $\sigma$; we may assume $\left|\sigma_{1}\right| \geq\left|\tau_{0}^{\prime}\right|$. If no $\sigma_{i} \in S$ then this is also an escape sequence in $T_{n+1}$ in $W^{n+1}$. If some $\sigma_{i} \in S$ then by Lemma 2.9 there is an $i$ with $\sigma_{i} \sqsupseteq \xi_{e^{t}}$ and therefore for some $\rho$

$$
\sigma_{0}, \ldots, \sigma_{i}, v_{0,0}^{\prime \prime} \frown \rho, \ldots, \nu_{0,0}^{\prime \prime} \frown \rho, \ldots, \theta_{e_{V}^{v}}^{V} \frown\langle 0\rangle \frown \rho
$$

is an escape sequence in $W^{n+1}$. This shows (1.1).
For (2.a) consider some determined $\sigma \in W_{k^{\prime}+1}^{n+1}$. If $\sigma$ is not in $S$ then $\sigma \in W_{k^{\prime}+1}^{n}$ and the claim follows since it was true in $T_{n}$. The interesting case is when $\sigma$ is in $S$; we first consider some points about the structure of the open loop $S$, which we may write $\zeta_{0}, \ldots, \zeta_{z}$. It is natural to divide $S$ into three pieces, $S \cap W^{n}, S \cap A_{k}^{n}$, and $S \cap B_{k}^{n}$. There are $z_{0}<z_{1}$ so that

$$
\bigcup_{i \leq z_{0}}\left[\zeta_{i}\right]=S \cap W^{n}, \bigcup_{z_{0}<i \leq z_{1}}\left[\zeta_{i}\right]=S \cap A_{k}^{n}, \bigcup_{z_{1}<i \leq z}\left[\zeta_{i}\right]=S \cap B_{k}^{n} .
$$

Furthermore, we have $\lambda\left(S \cap W^{n}\right)=L_{T_{n}}(\tau) \lambda(\tau)+e^{t} 2^{-N+N^{\prime}}$. On the other hand $\lambda\left(S \cap A_{k}^{n}\right)=\left(1-2^{-N}\right) \lambda\left(\bigcup_{i}\left[v_{i}\right]\right)+2^{-N+N^{\prime}} \sum_{i \leq U} e_{u}^{i}$. By choosing $N$ small enough, we ensured that the first term was larger than $\lambda\left(S \cap W^{n}\right)$,
so $\lambda\left(S \cap A_{k}^{n}\right)>\lambda\left(S \cap W^{n}\right)$. In particular, this means $z_{1}>2 z_{0}$. Finally $\lambda\left(S \cap B_{k}^{n}\right)>\left(1-2^{-N}\right) \lambda\left(\bigcup_{i}\left[\nu_{i}\right]\right)>4\left(1-2^{-N}\right) \lambda\left(\bigcup_{i}\left[\nu_{i}\right]\right)$, so again by choosing $N$ and $N^{\prime}$ large enough, we ensured that $\lambda\left(S \cap B_{k}^{n}\right)>4 \lambda\left(S \cap A_{k}^{n}\right)$. In particular, this means $z>4\left(z_{1}-z_{0}\right)>4 z_{0}$. We'll write $\alpha=\left\{i \mid\left[\zeta_{i}\right] \subseteq\right.$ $\left.2^{\omega} \backslash A\right\}$; in particular $\left(z_{1}, z\right] \subseteq \alpha$ and $\left(z_{0}, z_{1}\right] \cap \alpha=\emptyset$.

Now consider some $\sigma$ in $S \cap W_{k^{\prime}+1}^{n+1}$. Then $\sigma \in W_{k^{\prime}}^{n}$, so $\sigma \in\left[\zeta_{s}\right]$ for some $s \leq z_{0}$. So it suffices to show that we add an upcrossing to $\zeta_{s}$. In $T_{n}$, there was an upcrossing sequence for $\zeta_{s}$ of length $k^{\prime}-$ say, $0 \leq u_{1}<v_{1}<\cdots<$ $u_{k^{\prime}}<v_{k^{\prime}}$. Then $0 \leq u_{1}<v_{1}<\cdots<u_{k^{\prime}}<v_{k^{\prime}}$ is an upcrossing sequence for $\zeta_{s}$ in $T_{n+1}$. We claim that $0 \leq u_{1}<v_{1}<\cdots<u_{k^{\prime}}<v_{k^{\prime}}<z_{1}-s<z-s$ is an upcrossing sequence in $T_{n+1}$ as well. This is because

$$
\frac{1}{z_{1}-s} \sum_{j=s}^{z_{1}} \chi_{\alpha}(j) \leq \frac{z_{0}-s}{z_{1}-s} \leq \frac{z_{0}}{z_{1}}<1 / 2
$$

while

$$
\frac{1}{z-s} \sum_{j=s}^{z} \chi_{\alpha}(j) \geq \frac{z-z_{1}}{z-s} \geq \frac{z-z_{0}}{z}>3 / 4
$$

(2.b) (2.e.i) are immediate from the definition. For any $\eta \notin S$ determined in $T_{n}^{7}, L_{T_{n}^{7}}(\eta)=L_{T_{n}}(\eta)$ since this is preserved by each step in the construction of $T_{n}^{7}$. The passage from $T_{n}^{7}$ to $T_{n+1}$ only affects the loop $S$, so $L_{T_{n+1}}(\eta)=L_{T_{n}^{7}}(\eta)=L_{T_{n}}(\eta)$, giving (2.e.ii).

For any $k^{\prime}$ we have $a_{k^{\prime}}^{n}<\lambda\left(A_{k^{\prime}}^{n}\right)$ and $A_{k^{\prime}}^{n+1}=A_{k^{\prime}}^{n} \backslash\left(S \cap A_{k^{\prime}}^{n}\right)$ (where $S \cap A_{k^{\prime}}^{n}=\emptyset$ unless $\left.k=k^{\prime}\right)$, so also $a_{k^{\prime}}^{n}-\lambda\left(A_{k^{\prime}}^{n}\right)<\lambda\left(A_{k^{\prime}}^{n+1}\right)$. The same holds for $b_{k^{\prime}}^{n}$. Since we could choose the values $j_{k^{\prime}}$ arbitrarily small, we can make them small enough that (2.f) holds.

We turn to (2.g.i) For $k^{\prime} \neq k$, if there is a $\eta \in W_{k^{\prime}}^{n+1} \backslash S$ then there was an $\eta^{\prime} \in W_{k^{\prime}}^{n}$ with $\eta^{\prime} \sqsubseteq \eta, \rho^{n+1}(\eta)=\rho^{n}(\eta)$, and $L_{T_{n+1}}(\eta)=L_{T_{n}}(\eta)$. (Recall that $L$ is the length of $S$.) Then

$$
\begin{aligned}
\sum_{j \in J_{k}^{n+1}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}}\right]\right)=\right. & \sum_{j \in J_{k}^{n+1} \backslash\left\{j_{k^{\prime}}\right\}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}}\right]\right)\right. \\
& +l_{k^{\prime}, j_{k^{\prime}}}\left(2^{-j_{k^{\prime}}}-\lambda\left(\left[V_{j_{k^{\prime}}, n^{\prime}}\right]\right)\right. \\
\leq & \sum_{j \in J_{k}^{n}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}-1}\right]\right)\right. \\
& +l_{k^{\prime}, j_{k^{\prime}}}\left(2^{-j_{k^{\prime}}}-\lambda\left(\left[V_{j_{k^{\prime}}, n^{\prime}}\right]\right)\right. \\
\leq & a_{k^{\prime}}^{n}+L 2^{-j_{k^{\prime}}} \\
= & a_{k^{\prime}}^{n+1}
\end{aligned}
$$

For $k$ we have to take into account that we have reduced $a_{k}^{n}$ by $L_{T_{n}}(\tau) \lambda([\tau])$, and that this is compensated for by the fact that $\tau \in V_{j, n^{\prime}}$ :

$$
\begin{aligned}
\sum_{j \in J_{k}^{n+1}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}}\right]\right) \leq\right. & \sum_{j \in J_{k}^{n+1} \backslash\left\{j_{k}\right\}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}-1}\right]\right)\right. \\
& -l_{k, j^{\prime}} \lambda\left(\left[V_{\hat{j}, n^{\prime}}\right]\right) \\
& +l_{k, j_{k}}\left(2^{-j_{k}}-\lambda\left(\left[V_{j_{k}, n^{\prime}}\right]\right)\right. \\
\leq & \sum_{j \in J_{k}^{n}} l_{k, j}\left(2^{-j}-\lambda\left(\left[V_{j, n^{\prime}-1}\right]\right)\right. \\
& -L_{T_{n}}(\tau) \lambda([\tau]) \\
& +l_{k, j_{k}}\left(2^{-j_{k}}-\lambda\left(\left[V_{\left.j_{k}, n^{\prime}\right]}\right]\right)\right. \\
\leq & a_{k}^{n}-L_{T_{n}}(\tau) \lambda([\tau])+L 2^{-j_{k}} \\
= & a_{k}^{n+1} .
\end{aligned}
$$

The argument for (2.g.ii) is identical.
Since $n$ is odd, an integer $i<n / 2$ iff $i<(n+1) / 2$, so (3.a) (3.d) follow immediately from the inductive hypothesis.

Even stages-creating upcrossings; undetermined elements: Recall that we enumerated $\pi$ into $V_{j}$ and consider the decomposition of $\pi$ into determined elements $\pi_{i}$. Our last difficulty is ensuring (2.h)

We first observe that if $\pi$ is determined, so $\tau=\pi$ in the construction above, then we have satisfied (2.h) First, in the $\rho^{n}(\pi) \neq \hat{j}$ case, no changes were made and (2.h) followed immediately from the inductive hypothesis. When $\rho^{n}(\pi)=\hat{j}$, if $\sigma$ is not in $S$ then $\rho^{n+1}(\sigma)=\rho^{n}(\sigma)$ and $\sigma$ avoids $\pi$, so

$$
\left[V_{\rho^{n+1}(\sigma), n / 2}\right] \cap[\sigma]=\left[V_{\rho^{n}(\sigma), n / 2}\right] \cap[\sigma]=\left[V_{\rho^{n}(\sigma),(n-1) / 2}\right] \cap[\sigma]=\emptyset
$$

since $\pi$ was the only element in $\bigcup_{j} V_{j,(n-1) / 2} \backslash \bigcup_{j} V_{j,(n-1) / 2-1}$. For any element in $S, \rho^{n+1}(\sigma) \neq \rho^{n}(\sigma)$, and we chose a fresh value for $\rho^{n+1}(\sigma)$, so we may assume we chose $V_{\rho^{n+1}(\sigma), n / 2}$ to be empty.

We now consider the general case where $\pi$ is not determined. We have a partition $[\pi]=\bigcup_{i \leq r}\left[\pi_{i}\right]$ where each $\pi_{i}$ is determined. We let $U_{0}=\{i \mid$ $\left.\rho^{n}\left(\pi_{i}\right)=\hat{j}\right\}$. If $U_{0}=\emptyset$ then we may simply take $T_{n+1}=T_{n}$, so assume $U_{0}$ is non-empty. Pick some $i_{0} \in U_{0}$ such that the final element of the open loop containing $\pi_{i_{0}}$ has an escape sequence disjoint from $\bigcup_{i \in U_{0}}\left[\pi_{i}\right]$. To find such an escape sequence, take any $\pi_{i}$ with $i \in U_{0}$ and take an escape sequence $v_{0}, \ldots, v_{k}$ for the final element of the open loop containing $\pi_{i}$; if this is not already such an escape sequence, let $d \leq k$ be greatest such that $v_{d} \sqsupseteq \tau_{j}$ for some $j \in U$, and let $d^{\prime} \geq d$ be such that $v_{d^{\prime}}$ is contained in the final element of the open loop containing $\pi_{j}$. Then $v_{d^{\prime}+1}, \ldots, v_{k}$ is an escape sequence for the final element of the open loop containing $\pi_{j}$ which, since $d$ was chosen greatest, is disjoint from $\bigcup_{i \in U_{0}}\left[\pi_{i}\right]$.

We may apply the construction above with $\tau=\pi_{i_{0}}$, giving a transformation $T_{n}^{8,0}$ (the transformation that we referred to as $T_{n+1}$ above). For each $i \in U_{0}$, one of two things happens: either $\pi_{i}$ is in the open loop containing $\pi_{i_{0}}$ in $T_{n}$, in which case $\pi_{i}$ may not be determined in $T_{n}^{8,0}$, but any determined $\sigma$ in $\left[\pi_{i}\right]$ satisfies $\rho_{n}^{8,0}(\sigma) \neq \hat{j}$; or $\pi_{i}$ is not in the open loop containing $\pi_{i_{0}}$, in which case $\left[\pi_{i}\right]$ is disjoint from the escape sequence used in the construction of $T_{n}^{8,0}$, and so also disjoint from the open loop $S$, so $\pi_{i}$ is still determined in $T_{n}^{8,0}$. Let $U_{1} \subseteq U_{0}$ be the set of $i \in U_{0}$ such that $\pi_{i}$ is not in the open loop containing $\pi_{i_{0}}$ in $T_{n}$. Clearly $i_{0} \in U_{0} \backslash U_{1}$.

If $U_{1} \neq \emptyset$, we repeat this argument, choosing an $i_{1} \in U_{1}$ such that the final element of the open loop containing $\pi_{i_{1}}$ has an escape sequence disjoint from $\bigcup_{i \in U_{1}}\left[\pi_{i}\right]$, and we apply the construction above to give $T_{n}^{8,1}$. We repeat this argument $r^{\prime} \leq r$ times, giving $T_{n}^{8, r^{\prime}-1}$ so that $U_{r^{\prime}}=\emptyset$, and therefore every determined $\sigma$ in $[\tau]$ satisfies $\rho_{n}^{8, r^{\prime}-1}(\sigma) \neq \hat{j}$. Then we may set $T_{n+1}=T_{n}^{8, r^{\prime}-1}$. Each $T_{n}^{8, i}$ satisfies all inductive clauses except for (2.h) and further satisfies (2.h) for all $\sigma$ avoiding $\pi$ and also for all $\pi_{j}$ with $j \notin U_{i}$. Since $U_{r^{\prime}-1}=\emptyset$, $\overline{T_{n+1}}$ at last satisfies (2.h).
Checking the construction: We have completed the inductive construction. We now show that $\hat{T}$ is the desired transformation. Suppose $x \in \bigcap_{j}\left[V_{j}\right]$; we must show that $\hat{T}$ has infinitely many upcrossings on $x$. It suffices to show that $\hat{T}$ has at least $k$ upcrossings for every $k$. Suppose not. Since $x \in W^{0}$ and each element of $W_{k}^{n}$ has at least $k$ upcrossings, let $k$ be largest such that $x \in W_{k}^{n}$ for some $k$, and pick some large enough $n$ so that $x \in W_{k}^{n}$. Let $\sigma \in T_{n,+} \cup T_{n,-}$ be such that $x \in[\sigma]$. Let $j=\rho^{n}(\sigma)$. Since $x \in \bigcap_{j}\left[V_{j}\right]$ but $[\sigma] \cap\left[V_{j,(n-1) / 2}\right]=\emptyset$, there is some $m>n$ such that $x \in\left[V_{j,(m-1) / 2}\right]$. Let $\tau \in T_{m,+} \cup T_{m,-}$ be such that $x \in[\tau]$. Since $x \in W_{k}^{m}, \rho^{m}(\tau)=\rho^{n}(\sigma)$ by (2.e.i), but this contradicts (2.h)

To see that $|\hat{T}(x)|$ is infinite except on an effective $F_{\sigma}$ set of measure 0 , we claim that if $|\hat{T}(x)|$ is finite then is an $i_{0}$ such that for every $i>i_{0}, x \in G_{i}$. Since $\lambda\left(\bigcap_{i>i_{0}} G_{i}\right)=0$, we can also see that $\lambda\left(\bigcup_{i_{0}} \bigcap_{i>i_{0}} G_{i}\right)=0$, so this shows that the set of such $x$ has measure 0 . Let $|\hat{T}(x)|=k$ be finite. Choose some $n$ and some componential $\sigma \in T_{n,-}$ with $x \in[\sigma]$ and $\left|T_{n}(\sigma)\right|=k$. For each $m>n$, let $\sigma_{m} \in T_{m,-}$ be such that $x \in\left[\sigma_{m}\right]$. Since $\left|T_{m}\left(\sigma_{m}\right)\right|=k$ for all $m>n$, in particular at each odd $m$ we must have $\left[\sigma_{m}\right] \subseteq G_{(m-1) / 2}$. Therefore for all $i>(n-1) / 2$, we have $x \in G_{i}$, as desired.

## 5. The Ergodic Case

As promised above, we now present the strengthening of Gács, Hoyrup and Rojas' result to show that every Schnorr random point is Birkhoff for computable, bounded functions with computable ergodic transformations.

Theorem 5.1. Let $f$ be a bounded computable function and suppose $T$ is an ergodic, computable, measure-preserving transformation. Then every Schnorr random point is Birkhoff for $f$.

Proof. It is convenient to work with a strong Borel-Cantelli test: an effectively c.e. sequence $\left\langle V_{i}\right\rangle$ is a strong Borel-Cantelli test if $\sum_{i} \lambda\left(\left[V_{i}\right]\right)<\infty$ is a computable real number. It is known [13] that $x$ is Schnorr random iff for every strong Borel-Cantelli test $\left\langle V_{i}\right\rangle, x$ is in only finitely many $V_{i}$.

We will write $A_{n}(x)$ for the function $\frac{1}{n} \sum_{i<n} f\left(T^{i} x\right)$. Our main tool is Theorem 5.3 of [1]:

Let $T$ and $f$ be computable, let $f^{*}(x)=\lim _{n \rightarrow \infty} A_{n}(x)$, and suppose that $\left\|f^{*}\right\|_{L^{2}}$ is computable. Then for each $\epsilon>0$, there is an $N(\epsilon)$, computable from $T, f,\left\|f^{*}\right\|_{L^{2}}, \epsilon$ such that

$$
\mu\left(\left\{x\left|\max _{N(\epsilon) \leq m}\right| A_{m}(x)-A_{n}(x) \mid>\epsilon\right\}\right) \leq \epsilon .
$$

Since $T$ is ergodic, $f^{*}$ is the function constantly equal, almost everywhere, to $c=\int f d \lambda$, so $\left\|f^{*}\right\|_{L^{2}}$ is just $c$, which, since $f$ is bounded and computable, is itself a computable real number. Define a sequence $n_{i}=N\left(2^{-i-2}\right)$. Then for each $i$ we take

$$
V_{i}=\left\{x\left|\max _{n_{i} \leq m \leq n_{i+1}}\right| c-A_{m}(x) \mid>2^{-i}\right\} .
$$

If $x \in\left[V_{i}\right]$ then either $x$ is not a Birkhoff point-a set measure 0 by the Birkhoff ergodic theorem-or there is an $m^{\prime} \geq n_{i}=N\left(2^{-i}\right)$ with $\left|c-A_{m^{\prime}}(x)\right|<2^{-i-1}$. In the latter case we have $\left|A_{m}(x)-A_{m^{\prime}}(x)\right| \geq 2^{-i-1}$ and therefore either $\left|A_{n}(x)-A_{m}(x)\right| \geq 2^{-i-2}$ or $\left|A_{n}(x)-A_{m^{\prime}}(x)\right| \geq 2^{-i-2}$, and so $\max _{n_{i} \leq m}\left|A_{m}(x)-A_{n}(x)\right|>2^{-i-2}$. Therefore $\lambda\left(\left[V_{i}\right]\right) \leq 2^{-i-2}$.

Since the $V_{i}$ are computable sets and $\lim _{j \rightarrow \infty} \sum_{i>j} \lambda\left(\left[V_{i}\right]\right)=0$ with a computable rate of convergence, $\sum_{i} \lambda\left(\left[V_{i}\right]\right)$ is a computable real number. Since $\sum_{i} \lambda\left(\left[V_{i}\right]\right) \leq \sum_{i} 2^{-i-2}=1 / 2,\left\langle V_{i}\right\rangle$ is a strong Borel-Cantelli test.

If $x$ is not a Birkhoff point then for some $\delta>0$ there are infinitely many $m$ with $\left|c-A_{m}(x)\right|>\epsilon$. Fix $i$ with $2^{-i}<\delta$, and there are infinitely many $j>i$ with $x \in V_{j}$, so $x$ fails to be Schnorr random.

## 6. Upcrossings

Throughout this section, we will take $T$ to be a computable, measurepreserving transformation.

Recall the following theorem of Bishop [5]:

## Theorem 6.1.

$$
\int \tau(x, f, \alpha, \beta) d x \leq \frac{1}{\beta-\alpha} \int(f-\alpha)^{+} d x .
$$

This is easily used to derive the following special case of a theorem of V'yugin:

Theorem 6.2 (30). If $x$ is Martin-Löf random and $f$ is computable then $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ converges.

Proof. Suppose $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ does not converge. Then there exist $\alpha<\beta$ such that $\frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ is infinitely often less than $\alpha$ and also infinitely often greater than $\beta$. Equivalently, $\tau(x, f, \alpha, \beta)$ is infinite. But observe that when $f$ is computable, $\tau(x, f, \alpha, \beta)$ is lower semi-computable, so in particular,

$$
V_{n}=\{x \mid \tau(x, f, \alpha, \beta) \geq n\}
$$

is computably enumerable and $\mu\left(V_{n}\right) \leq \frac{1}{n(\beta-\alpha)} \int(f-\alpha)^{+} d x$. Therefore an appropriate subsequence of $\left\langle V_{n}\right\rangle$ provides a Martin-Löf test, and $x \in \cap_{n} V_{n}$, so $x$ is not Martin-Löf random.

We now consider the case where $f$ is lower semi-computable. We will have a sequence of uniformly computable increasing approximations $f_{i} \rightarrow f$, and we wish to bound the number of upcrossings in $f$. The difficulty is that $\tau\left(x, f_{i}, \alpha, \beta\right)$ is not monotonic in $i$ : it might be that an upcrossing sequence for $f_{i}$ ceases to be an upcrossing sequence for $f_{i+1}$.

In order to control this change, we need a suitable generalization of upcrossings, where we consider not only the upcrossings for $f$, but for all functions between $f$ and $f+h$ where $h$ is assumed to be small.

Definition 6.3. A loose upcrossing sequence at $x$ for $\alpha, \beta, f, h$ is a sequence

$$
0 \leq u_{1}<v_{1}<u_{2}<v_{2}<\cdots<u_{N}<v_{N}
$$

such that for all $i \leq N$,

$$
\frac{1}{u_{i}+1} \sum_{j=0}^{u_{i}} f\left(T^{j} x\right)<\alpha, \frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}}(f+h)\left(T^{j} x\right)>\beta .
$$

$v(x, f, h, \alpha, \beta)$ is the supremum of the lengths of loose upcrossing sequences for $\alpha, \beta, f, h$.

Loose upcrossings are much more general than we really need, and so the analog of Bishop's theorem is correspondingly weak. For instance, consider the case where $T$ is the identity transformation, $f=\chi_{A}$, and $h=\chi_{B}$ with $A$ and $B$ disjoint (so $f+h=\chi_{A \cup B}$ ). Then $v(x, f, h, \alpha, \beta)=\infty$ for every $x \in B$ whenever $0<\alpha<\beta<1$. Nonetheless, we are able to show the following:
Theorem 6.4. Suppose $h \geq 0, \int h d x<\epsilon$ and $\beta-\alpha>\delta$. There is a set $A$ with $\mu(A)<4 \epsilon / \delta$ such that

$$
\int_{X \backslash A} v(x, f, h, \alpha, \beta) d x
$$

is finite.
Proof. By the usual pointwise ergodic theorem, there is an $n$ and a set $A^{\prime}$ with $\mu\left(A^{\prime}\right)<2 \epsilon / \delta$ such that if $x \notin A^{\prime}$ then for all $n^{\prime}, n^{\prime \prime} \geq n$,

$$
\left|\frac{1}{n^{\prime}+1} \sum_{j=0}^{n^{\prime}} h\left(T^{j} x\right)-\frac{1}{n^{\prime \prime}+1} \sum_{j=0}^{n^{\prime \prime}} h\left(T^{j} x\right)\right|<\delta / 2
$$

Consider those $x \notin A^{\prime}$ such that, for some $n^{\prime} \geq n$,

$$
\frac{1}{n^{\prime}+1} \sum_{j=0}^{n^{\prime}} h\left(T^{j} x\right) \geq \delta
$$

We call this set $A^{\prime \prime}$. Then for all $n^{\prime} \geq n$, such an $x$ satisfies

$$
\frac{1}{n^{\prime}+1} \sum_{j=0}^{n^{\prime}} h\left(T^{j} x\right) \geq \delta / 2
$$

and in particular,

$$
\int_{A^{\prime \prime}} h d x \geq \delta \mu\left(A^{\prime \prime}\right) / 2
$$

Therefore $\mu\left(A^{\prime \prime}\right) \leq 2 \epsilon / \delta$. If we set $A=A^{\prime} \cup A^{\prime \prime}$, we have $\mu(A)<4 \epsilon / \delta$.
Now suppose $x \notin A$. We claim that any loose upcrossing sequence for $\alpha, \beta, f, h$ with $n \leq u_{1}$ is already an upcrossing sequence for $\alpha, \beta-\delta$. If $n \leq$ $u_{1}<v_{1}<\cdots<u_{N}<v_{N}$ is a loose upcrossing sequence, we automatically satisfy the condition on the $u_{i}$. For any $v_{i}$, we have

$$
\beta<\frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}}(f+h)\left(T^{j} x\right)=\frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}} f\left(T^{j} x\right)+\frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}} h\left(T^{j} x\right)
$$

Since $\frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}} h\left(T^{j} x\right) \leq \delta$, it follows that $\frac{1}{v_{i}+1} \sum_{j=0}^{v_{i}} f\left(T^{j} x\right)>\beta-\delta$ as desired. Therefore

$$
\int_{X \backslash A} v(x, f, h, \alpha, \beta) d x \leq \mu(X \backslash A) \int_{X \backslash A} n+\tau(x, f, \alpha, \beta-\delta) d x
$$

is bounded.
Theorem 6.5. If $x$ is weakly 2 -random and $f$ is lower semi-computable then $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ converges.

Proof. Suppose $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ does not converge. Then there exist $\alpha<\beta$ such that $\frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ is infinitely often less than $\alpha$ and also infinitely often greater than $\beta$. Equivalently, $\tau(x, f, \alpha, \beta)$ is infinite. Let $f_{n} \rightarrow f$ be the sequence of computable functions approximating $f$ from below.

For each $n$, we set

$$
V_{n}=\left\{x \mid \exists m \geq n v\left(x, f_{n}, f_{m}-f_{n}, \alpha, \beta\right) \geq n\right\}
$$

By construction, $x \in \cap_{n} V_{n}$. To see that $V_{n+1} \subseteq V_{n}$, observe that if

$$
v\left(x, f_{n+1}, f_{m}-f_{n+1}, \alpha, \beta\right) \geq n+1
$$

then there is a loose upcrossing sequence witnessing this, and it is easy to check (since the $f_{n}$ are increasing) that this is also a loose upcrossing sequence witnessing

$$
v\left(x, f_{n}, f_{m}-f_{n}, \alpha, \beta\right) \geq n+1>n
$$

We must show that $\mu\left(V_{n}\right) \rightarrow 0$. Fix $\delta<\beta-\alpha$ and let $\epsilon>0$ be given. Choose $n$ to be sufficiently large that $\left\|f_{n}-f\right\|<\delta \epsilon / 4$. Then, since the $f_{n}$ approximate $f$ from below, clearly $v\left(x, f_{m}, f_{m^{\prime}}-f_{m}, \alpha, \beta\right) \leq v\left(x, f_{m}, f-\right.$ $\left.f_{m}, \alpha, \beta\right)$ for any $m^{\prime} \geq m$. By the previous theorem, there is a set $A$ with $\mu(A)<\epsilon / 2$ such that $\int_{X \backslash A} v\left(x, f_{m}, f-f_{m}, \alpha, \beta\right) d x$ is bounded. We may choose $n^{\prime} \geq n$ sufficiently large that

$$
B=\mu\left(\left\{x \notin A \mid v\left(x, f_{m}, f-f_{m}, \alpha, \beta\right) \geq n^{\prime}\right\}\right)<\epsilon / 2
$$

Then $V_{n^{\prime}} \subseteq A \cup B$, so $\mu\left(V_{n^{\prime}}\right) \leq \epsilon$.
6.1. Room for Improvement. It is tempting to try to improve Theorem 6.4. The premises of that theorem are too general and the proof is oddly "half-constructive"-we mix the constructive and nonconstructive pointwise ergodic theorems. One would think that by tightening the assumptions and using Bishop's upcrossing version of the ergodic theorem in both places, we could prove something stronger.

In the next theorem, we describe an improved upcrossing property which, if provable, would lead to a substantial improvement to Theorem 6.5: balanced randomness would guarantee the existence of this limit. (Recall that a real is balanced random if it passes every balanced test, or sequence $\left\langle V_{i}\right\rangle$ of r.e. sets such that for every $i, \mu\left(\left[V_{i}\right]\right) \leq 2^{-i}$ and $V_{i}=W_{f(i)}$ for some function $f$ with a recursive approximation that has at most $2^{n}$ mind-changes for the value of $f(n)[10]$.) The property hypothesized seems implausibly strong, but we do not see an obvious route to ruling it out.

Theorem 6.6. Suppose the following holds:
Let $f$ and $\epsilon>0$ be given, and let $0 \leq h_{0} \leq h_{1} \leq \cdots \leq h_{n}$ be given with $\left\|h_{n}\right\|_{L^{\infty}}<\epsilon$. Then

$$
\int_{X} \sup _{n} \tau\left(x, f+h_{n}, \alpha, \beta\right) d x<c\left(\|f\|_{L^{\infty}, \epsilon}\right)
$$

where $c\left(\|f\|_{\left.L^{\infty}, \epsilon\right)}\right.$ is a computable bound depending only on $\|f\|_{L^{\infty}}$ and $\epsilon$.
Then whenever $x$ is balanced random and $f$ is lower semi-computable then $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ converges.
Proof. We assume $\|f\|_{L^{2}} \leq 1$ (if not, we obtain this by scaling). Suppose $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ does not converge. Then there exist $\alpha<\beta$ such that $\frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ is infinitely often less than $\alpha$ and also infinitely often greater than $\beta$. Equivalently, $\tau(x, f, \alpha, \beta)$ is infinite. Let $f_{n} \rightarrow f$ be the sequence of computable functions approximating $f$ from below.

We define the set

$$
V_{(n, k)}=\left\{x \mid \exists m \geq n \tau\left(x, f_{m}, \alpha, \beta\right) \geq k\right\}
$$

We then define the function $g\left(n, n^{\prime}\right)$ to be least such that $\forall m \in\left[n, n^{\prime}\right] \| f_{n^{\prime}}-$ $f_{m} \|<2^{-n}$ and $g(n)=\lim _{n^{\prime}} g\left(n, n^{\prime}\right)$. Since the sequence $f_{m}$ converges to $f$ from below, $g(n)$ is defined everywhere, and $|\{s \mid g(n, s+1) \neq g(n, s)\}|<2^{n}$
for all $n$. Indeed, $g(n)$ is the least number such that $\forall m \geq g(n)\left\|f-f_{m}\right\| \leq$ $2^{-n}$.

Observe that $\mu\left(V_{(n, k)}\right)<\frac{c\left(\|f\|_{L} \infty, 2^{-n}\right)}{k}$. Choose $h(n)$ to be a computable function growing quickly enough that $\frac{c\left(\|f\|_{L^{\prime}, 2^{-n}}\right)}{h(n)} \leq 2^{-n}$ for all $n$. If $x \in$ $V_{(g(n+1), h(n+1))}$ then there is some $m \geq g(n+1)$ so that $\tau\left(x, f_{m}, \alpha, \beta\right) \geq$ $h(n+1)$. Since $g(n+1) \geq g(n)$ and $h(n+1) \geq h(n)$, we also have that $x \in V_{(g(n), h(n))}$. Therefore $\left\langle V_{(g(n), h(n))}\right\rangle$ is a balanced test.

But since $\tau(x, f, \alpha, \beta)$ is infinite, we must have $x \in \cap V_{(g(n), h(n))}$. This contradicts the assumption that $x$ is balanced random, so $\lim _{n \rightarrow \infty} \frac{1}{n+1} \sum_{j=0}^{n} f\left(T^{j} x\right)$ converges.

In fact, the test $\left\langle V_{(g(n), h(n))}\right\rangle$ has an additional property: if $s_{0}<s_{1}<s_{2}$ with $g\left(n+1, s_{0}\right) \neq g\left(n+1, s_{1}\right) \neq g\left(n+1, s_{2}\right)$ then $g\left(n, s_{0}\right) \neq g\left(n, s_{2}\right)$. This means that $\left\langle V_{(g(n), h(n))}\right\rangle$ is actually an Oberwolfach test [3], and so we can weaken the assumption to $x$ being Oberwolfach random.

## 7. Discussion for Ergodic Theorists

In the context of analytic questions like the ergodic theorem, matters of computability are mostly questions of continuity and uniformity: the computability of a given property usually turns on whether it depends in an appropriately uniform way on the inputs. Algorithmic randomness gives a precise way of characterizing how sensitive the ergodic theorem is to small changes in the underlying function.

The paradigm is to distinguish different sets of measure 0 , viewed as an intersection $A=\cap_{i} A_{i}$, by characterizing how the sets $A_{i}$ depend on the given data (in the case of the ergodic theorem, the function $f$ ). The two main types of algorithmic randomness that have been studied in this context thus far are Martin-Löf randomness and Schnorr randomness. In both cases, we ask that the sets $A_{i}$ be unions $A_{i}=\cup_{j} A_{i, j}$ of sets where $A_{i, j}$ is determined based on a finite amount of information about the orbit of $f$ (in particular, the dependence of $A_{i, j}$ on $f$ and $T$ should be continuous). (To put it another way, we ask that the set of exceptional points which violate the conclusion of the ergodic theorem be contained in a $G_{\delta}$-set which depends on $f$ in a uniform way.) The distinction between the two notions is that in Schnorr randomness, $\mu\left(A_{i}\right)=2^{-i}$, while in Martin-Löf randomness, we only know $\mu\left(A_{i}\right) \leq 2^{-i}$. This means that in the Schnorr random case, a finite amount of information about the orbit of $f$ suffices to limit the density of $A_{i}$ outside of a small set (take $J$ large enough that $\mu\left(\cup_{j \leq J} A_{i, j}\right)$ is within $\epsilon$ of $2^{-i}$; then no set disjoint from $\cup_{i \leq J} A_{i, j}$ contains more than $\epsilon$ of $A_{i}$ ). In the Martin-Löf random case, this is not possible: if $\mu\left(A_{i}\right) \leq 2^{-i}-\epsilon$, no finite amount of information about the orbit of $f$ can rule out the possibility that some $A_{i, j}$ with very large $j$ will add a set of new points of measure $\epsilon$. In particular, while we can identify sets which do belong to $A_{i}$, finite information about the orbit of $f$ does not tell us much about which points are not in $A_{i}$.

The two classes of functions discussed in this paper are the computable and the lower semi-computable ones; these are closely analogous to the continuous and lower semi-continuous functions. Unsurprisingly, both the passage from computable to lower semi-computable functions and the passage from ergodic to nonergodic transformations make it harder to finitely characterize points violating the conclusion of the ergodic theorem. Perhaps more surprising, both changes generate precisely the same result: if a point violates the conclusion of the ergodic theorem for a computable function with a nonergodic transformation, we can construct a lower semi-computable function with an ergodic transformation for which the point violates the conclusion of the ergodic theorem, and vice versa.

The remaining question is this: What happens when we make both changes? What characterizes the points which violate the conclusion of the ergodic theorem for lower-semi computable functions with nonergodic transformations? The answer is likely to turn on purely ergodic theoretic questions about the sensitivity of upcrossings, such as the hypothesis we use above.

Question 7.1. Let $(X, \mu)$ be a metric space and let $T: X \rightarrow X$ be measure preserving. Let $\epsilon>0$ be given. Is there a bound $K$ (depending on $T$ and on $\epsilon)$ such that for any $f$ with $\|f\|_{L^{\infty}} \leq 1$ and any sequence $0 \leq h_{0} \leq h_{1} \leq$ $\cdots \leq h_{n}$ with $\left\|h_{n}\right\|_{L^{\infty}}<\epsilon$,

$$
\int \sup _{n} \tau\left(x, f+h_{n}, \alpha, \beta\right) d x<K
$$

where $\tau(x, g, \alpha, \beta)$ is the number of upcrossings from below $\alpha$ to above $\beta$ starting with the point $x$ ?

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Department of Mathematics, University of Connecticut U-3009, 196 Auditorium Road, Storrs, CT 06269-3009, USA

E-mail address: johanna.franklin@uconn.edu
URL: www.math.uconn.edu/~franklin
Department of Mathematics, University of Pennsylvania, 209 South 33rd Street, Philadelphia, PA 19104-6395, USA

E-mail address: htowsner@math.upenn.edu
URL: http://www.math.upenn.edu/~htowsner


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[^1]:    ${ }^{1}$ Actually, according to [12], the method was first used several decades earlier by von Neumann and Kakutani, but not published until later [16.

[^2]:    ${ }^{2}$ It is not possible to ensure that every element of the set $V_{j}$ receives $j$ upcrossings, since this would imply that the theorem holds for every $x$ which failed to be even Demuth random, which would contradict V'yugin's theorem.

