Ulam’s Method for Lasota–Yorke Maps with Holes

Christopher Bose†, Gary Froyland‡, Cecilia González-Tokman‡, and Rua Murray§

Abstract. Ulam’s method is a rigorous numerical scheme for approximating invariant densities of dynamical systems. The phase space is partitioned into a grid of connected sets, and a set-to-set transition matrix is computed from the dynamics; an approximate invariant density is read off as the leading left eigenvector of this matrix. When a hole in phase space is introduced, one instead searches for conditional invariant densities and their associated escape rates. For Lasota–Yorke maps with holes we prove that a simple adaptation of the standard Ulam scheme provides convergent sequences of escape rates (from the leading eigenvalue), conditional invariant densities (from the corresponding left eigenvector), and quasi-conformal measures (from the corresponding right eigenvector). We also immediately obtain a convergent sequence for the invariant measure supported on the survivor set. Our approach allows us to consider relatively large holes. We illustrate the approach with several families of examples, including a class of Lorenz-like maps.

Key words. open dynamical systems, Ulam’s method, Lasota–Yorke maps

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1. Introduction. Dynamical systems $\hat{T}: I \to I$ typically model complicated deterministic processes on a phase space $I$. The map $\hat{T}$ induces a natural action on probability measures $\eta$ on $I$ via $\eta \mapsto \eta \circ \hat{T}^{-1}$. Of particular interest in ergodic theory are those probability measures that are $\hat{T}$-invariant, that is, $\eta$ satisfying $\eta = \eta \circ \hat{T}^{-1}$. If $\eta$ is ergodic, then such $\eta$ describe the time-asymptotic distribution of orbits of $\eta$-almost-all initial points $x \in I$. In this paper, we consider the situation where a “hole” $H_0 \subsetneq I$ is introduced and any orbits of $\hat{T}$ that fall into $H_0$ terminate. The hole induces an open dynamical system $T: X_0 \to I$, where $X_0 = I \setminus H_0$. Because trajectories are being lost to the hole, in many cases, there is no absolutely continuous $T$-invariant probability measure. One can, however, consider conditionally invariant probability measures [28], which satisfy $\eta \circ T^{-1}(I) \cdot \eta = \eta \circ T^{-1}$, where $0 < \eta \circ T^{-1}(I) < 1$ is identified as the escape rate for the open system.
We will study $\tilde{T}$ drawn from the class of Lasota–Yorke maps: piecewise $C^1$ expanding maps of the interval, such that $|D\tilde{T}|^{-1}$ has bounded variation. The hole $H_0$ will be a finite union of intervals. In such a setting, because of the expanding property, one can expect to obtain conditionally invariant probability measures that are absolutely continuous with respect to Lebesgue measure [5, 33, 22]. Such conditionally invariant measures are “natural” when they correspond to the result of repeatedly pushing forward a Lebesgue measure by $\tilde{T}$.

In the next section we will discuss further conditions due to [22] that make this precise: (i) how much of phase space can “escape” into the hole, and (ii) the growth rate of intervals that partially escape relative to the expansion of the map and the rate of escape. These conditions also guarantee the existence of a unique absolutely continuous conditionally invariant probability measure (accim). This accim $\nu$ and its corresponding escape rate $\rho$ are the first two objects that we will rigorously numerically approximate using Ulam’s method. Existence and uniqueness results for subshifts of finite type with Markov holes were previously established by Collet, Martínez, and Schmitt in [8]; see also [6, 7, 17].

One may also consider the set of points $X_\infty \subset I$ that never fall into the hole $H_0$. A probability measure $\lambda$ on $X_\infty$ can be defined as the $n \to \infty$ limit of the accim $\nu$ conditioned on $X_n$. The measure $\lambda$ will turn out to be the unique $\tilde{T}$-invariant measure supported on $X_\infty$, which is absolutely continuous with respect to $\mu$, the quasi-conformal measure for $T$ with escape rate $\rho$. We will also rigorously numerically approximate $\mu$ and thus $\lambda$. Robustness of these objects with respect to Ulam discretizations is essentially due to a quasi–compactness property, and a significant part of the paper is devoted to elaborating on this point.

Our main result, Theorem 3.2, concerns convergence properties of an extension of the well-known construction of Ulam [32], which allows for efficient numerical estimation of invariant densities of closed dynamical systems. The Ulam approach partitions the domain $I$ into a collection of connected sets $\{I_1, \ldots, I_k\}$ and computes single-step transitions between partition sets, producing the matrix

$$\hat{P}_{ij} = \frac{m(I_i \cap \tilde{T}^{-1}I_j)}{m(I_j)}.$$  

Li [21] demonstrated that the invariant density of Lasota–Yorke maps can be $L^1$-approximated by step functions obtained directly from the leading left eigenvector of $\hat{P}$. Since the publication of [21] there have been many extensions of Ulam’s method to more general classes of maps, including expanding maps in higher dimensions [10, 26], uniformly hyperbolic maps [12, 14], nonuniformly expanding interval maps [27, 15], and random maps [13, 18]. Explicit error bounds have also been developed, e.g., in [25, 13, 4].

We will show that in order to handle open systems, the definition of $\hat{P}$ above need only be modified to $P$, having entries

$$P_{ij} = \frac{m(I_i \cap X_0 \cap \tilde{T}^{-1}I_j)}{m(I_j)}.$$  

In analogy with the closed setting, one uses the leading left eigenvector to produce a step function that solves an eigenequation, from which we can easily recover an approximation to

\[ \textup{See Definition 2.3 for the precise meaning, and [22] for a proof of uniqueness.} \]
the accim $\nu$. However, in the open setting, the leading eigenvalue of $P$ also approximates the escape rate $\rho$ of $\nu$, and the right eigenvector approximates the quasi-conformal measure $\mu$. Note that for closed systems, $\rho = 1$ and $\mu = m$.

The literature concerning the analysis of Ulam’s method is now quite large. Early work on Ulam’s method for Axiom A repellers [14] showed convergence of an Ulam-type scheme using Markov partitions for the approximation of pressure and equilibrium states with respect to the potential $-\log |\det D\hat{T}|_{E^u}$. These results apply to the present setting of Lasota–Yorke maps, provided that the hole is Markov and projections are done according to a sequence of Markov partitions. Bahsoun [1] considered non-Markov Lasota–Yorke maps with non-Markov holes and rigorously proved an Ulam-based approximation result for the escape rate. Bahsoun used the perturbative machinery of [20], treating the map $T$ as a small deterministic perturbation of the closed map $\hat{T}$. In contrast, we apply the perturbative arguments of [20] directly to the open map, considering the Ulam discretization as a small perturbation of $T$. The advantage of this approach is that we can obtain approximation results whenever the existence results of [22] apply. The latter make assumptions on the expansivity of $T$ (large enough), the escape rate (slow enough), and the rate of generation of “bad” subintervals (small enough). From these assumptions we construct an improved Lasota–Yorke inequality that allows us to get tight enough constants to make applications plausible. Besides estimating the escape rate, we obtain rigorous $L^1$-approximations of the accim and approximations of the quasi-conformal measure that exploit quasi–compactness and converge weakly to $\mu$. We can treat relatively large holes.

An outline of the paper is as follows. In section 2 we introduce the Perron–Frobenius operator $L$, formally define admissible and Ulam-admissible holes, and develop a strong Lasota–Yorke inequality. Section 3 introduces the new Ulam scheme and states our main Ulam convergence result. Section 4 discusses some specific example maps in detail. Proofs are presented in section 5.

2. Lasota–Yorke maps with holes. The following class of interval maps with holes was studied by Liverani and Maume-Deschamps in [22].

**Definition 2.1.** Let $I = [0, 1]$. We call $T : I \ni a$ Lasota–Yorke map if $T$ is a piecewise $C^1$ map, with finite monotonicity partition $Z$, there exists $\Theta < 1$ such that $\|DT^{-1}\|_{\infty} \leq \Theta$, and $\hat{g} := |DT|^{-1}$ has bounded variation.

The transfer operator for the map $T$ is the bounded linear operator $\hat{L}$, acting on the space $BV$ of functions of bounded variation on $I$, defined by

$$\hat{L}f(x) = \sum_{T(y) = x} f(y)\hat{g}(y).$$

**Definition 2.2.** Let $\hat{T} : I \ni be a Lasota–Yorke map. Let $H_0 \subset I$ be a finite union of closed intervals, and let $X_0 = I \setminus H_0$. Let $T : X_0 \rightarrow I$ be the restriction $T = \hat{T}|_{X_0}$. Both $T$ and the pair $T_0 = (\hat{T}, H_0)$ are referred to as open Lasota–Yorke maps (or, briefly, open systems), and

\[^3\]Throughout this paper, a monotonicity partition $Z$ refers to a partition such that, for every $Z \in Z$, $\hat{T}|_{Z}$ has a $C^1$ extension to $\hat{Z}$.
their associated transfer operator is the bounded linear operator \( \mathcal{L} : BV \cap \) given by

\[
\mathcal{L}(f) = \hat{\mathcal{L}}(1_{X_0} f).
\]

For each \( n \geq 1 \) let \( X_n = \bigcap_{j=0}^{n} T^{-j} X_0 \). Thus, \( X_n \) is the set of points that have not escaped by time \( n \). Also, we denote by \( T^n \) the function \( T^n|_{X_{n-1}} \). One can readily check that

\[
\mathcal{L}^n(f) = \hat{\mathcal{L}}^n(1_{X_{n-1}} f).
\]

**Definition 2.3.** Let \( T \) be an open Lasota–Yorke map. A probability measure \( \nu \) supported on \( X_0 \subset I \) which is absolutely continuous with respect to Lebesgue measure is called an absolutely continuous conditional invariant measure (accim) on \( f \) if there exists a function of bounded variation \( h \) such that \( 1_{X_0} \cdot h = \frac{d
u}{d\nu_0} \) and \( \mathcal{L}h = \rho h \) for some \( 0 < \rho \leq 1 \).

A probability measure \( \mu \) on \( I \) which satisfies \( \mu(\mathcal{L}f) = \rho \mu(f) \) for every function of bounded variation \( f : I \to \mathbb{R} \), with \( \rho \) as above, is called a quasi-conformal measure for \( T \).

**Remark 2.4.** We choose to display \( h \) as opposed to \( 1_{X_0} h \) in the upcoming figures, because our numerical method directly discretizes the eigenfunction \( \mathcal{L}h = \rho h \). Further, the value of \( h \) outside \( X_0 \) illustrates the amount of mass that escapes the open system in one step. For convenience of notation, and despite the fact that the support of \( h \) may intersect \( H_0 \), we will refer to \( h \) as the accim as well.

**Remark 2.5.** It is usual to define \( \nu \) to be an accim if \( \nu(A) = \frac{\nu(T^{-n}A \cap X_n)}{\nu(X_n)} \) for every \( n \geq 0 \) and Borel-measurable set \( A \subset I \). This definition and that of Definition 2.3 are indeed equivalent; see [22, Lemma 1.1] for a proof. The same lemma shows that if \( \mu \) is a quasi-conformal measure for \( T \), then \( \mu \) is necessarily supported on \( X_\infty = \bigcap_{n\geq0} X_n \). It is also usual to require \( \mu \) to satisfy \( \mu(\mathcal{L}f) = \rho \mu(f) \) for continuous functions only. We will see that this makes no difference in our setting, as this weaker requirement implies the stronger one in the previous definition.

### 2.1. Admissible holes and quasi-invariant measures

As in the work of Liverani and Maume-Deschamps [22], we impose some conditions on the open system in order to be able to analyze it. Let us fix some notation.

Let \( (\hat{T}, H_0) \) be an open Lasota–Yorke map, which we also refer to as \( T \). For each \( n \geq 1 \), let \( D_n = \{ x \in I : L^n 1(x) \neq 0 \} \), and let \( D_{\infty} := \bigcap_{n \geq 1} D_n \). In what follows, we assume that \( D_{\infty} \neq \emptyset \).

For each \( \epsilon > 0 \) (not necessarily small), we let \( G_\epsilon = G_\epsilon(T) \) be the collection of finite partitions of \( I \) into intervals such that \( Z_\epsilon \in G_\epsilon(T) \) if (i) the interior of each \( A \in Z_\epsilon \) is either disjoint from or contained in \( X_0 \), and (ii) for each \( A \in Z_\epsilon \), \( \text{var}_A(1_{X_0} | DT^{-1}|) < ||DT^{-1}||_\infty (1 + \epsilon) \). Since \( H_0 \) consists of finitely many intervals, this condition is possible to achieve, as the work of Rychlik [29, Lemma 6] shows. We call \( G_\epsilon \) the collection of \( \epsilon \)-adequate partitions (for \( T \)).

The set of elements of \( Z_\epsilon \), whose interiors are contained in \( X_0 \) is denoted by \( Z_\epsilon^* \). Next, the elements of \( Z_\epsilon^* \) are divided into good and bad. A set \( A \in Z_\epsilon^* \) is good if

\[
\lim_{n \to \infty} \inf_{x \in D_n} \frac{\mathcal{L}^n 1_A(x)}{\mathcal{L}^n 1(x)} > 0.
\]

We point out that it is shown in [22] that the limit above always exists, as the sequence involved is increasing and bounded, and it is clearly nonnegative. The set \( A \) is called bad
when the limit above is 0. We let
\[ Z_{\epsilon,g} = \{ A \in Z_\epsilon^* : A \text{ is good} \}, \]
\[ Z_{\epsilon,b} = \{ A \in Z_\epsilon^* : A \text{ is bad} \}. \]

Finally, two elements of \( Z_\epsilon^* \) are called contiguous if there are no other elements of \( Z_\epsilon^* \) in between them (but there may be elements of \( Z_\epsilon \) that are necessarily contained in \( H_0 \)). We let \( \xi_\epsilon = \xi_\epsilon(T) \) be the infimum over \( \epsilon \)-adequate partitions for \( T \) of the maximum number of contiguous elements in \( Z_{\epsilon,b} \).

In a similar manner, we let \( G_{\epsilon}^{(n)} = G_{\epsilon}^{(n)}(T) \) be the collection of finite partitions of \( I \) into intervals such that \( Z_\epsilon^{(n)} \in G_{\epsilon}^{(n)}(T) \) if (i) the interior of each \( A \in Z_\epsilon^{(n)} \) is either disjoint from or contained in \( X_{n-1} \), and (ii) for each \( A \in Z_\epsilon^{(n)} \), \( \|1_{X_{n-1}}(DT^n)^{-1}\| < \|(DT^n)^{-1}\|_\infty(1 + \epsilon) \).

The partitions \( Z_\epsilon^{(n)}, Z_{\epsilon,g}^{(n)}, Z_{\epsilon,b}^{(n)} \) are defined analogously. We denote by \( \xi_{\epsilon,n} = \xi_{\epsilon,n}(T) \) the infimum over \( \epsilon \)-adequate partitions for \( T^n \) of the maximum number of contiguous elements in \( Z_{\epsilon,b}^{(n)} \); thus \( \xi_\epsilon = \xi_{\epsilon,1} \).

The following quantities are relevant in what follows:
\[
\rho = \rho(T) := \lim_{n \to \infty} \inf_{x \in D_n} \frac{\mathcal{L}^{n+1}(x)}{\mathcal{L}^n(1(x))},
\]
\[
\tilde{\Theta} = \tilde{\Theta}(T) := \exp \left( \lim_{n \to \infty} \frac{1}{n} \log \|DT^n\|^{-1}_\infty \right),
\]
\[
\tilde{\xi}_\epsilon = \tilde{\xi}_\epsilon(T) := \exp \left( \limsup_{n \to \infty} \frac{1}{n} \log(1 + \xi_{\epsilon,n}) \right),
\]
\[
\alpha_\epsilon = \alpha_\epsilon(T) := \|DT^{-1}\|_\infty(2 + \epsilon + \xi_\epsilon).
\]

**Definition 2.6 (admissible holes).** Let \( \tilde{T} : I \to I \) be a Lasota–Yorke map and \( \epsilon > 0 \). We say that \( H_0 \subset I \) is
- an \( \epsilon \)-admissible hole for \( \tilde{T} \) if \( D_\infty \neq \emptyset \) and \( \xi_{\epsilon} \tilde{\Theta} < \rho \),
- an admissible hole for \( \tilde{T} \) if it is \( \epsilon \)-admissible for \( \epsilon = 1 \),
- an \( \epsilon \)-Ulam-admissible hole for \( \tilde{T} \) if \( D_\infty \neq \emptyset \) and \( \alpha_\epsilon < \rho \).

The main result of Liverani and Maume-Deschamps [22] is concerned with the existence of the objects we intend to rigorously numerically approximate. Relevant quasi-compactness properties of \( \mathcal{L} \) are made explicit as follows.

**Theorem 2.7 (see [22, Theorem A and Lemma 3.10]).** Assume that \((\tilde{T}, H_0)\) is an open system with an admissible hole. Then we have the following:
1. There exists a unique accim, \( \nu = 1_{I \setminus X_0} \) for \((\tilde{T}, H_0)\).
2. There exists a unique quasi-conformal measure \( \mu \) for \((\tilde{T}, H_0)\) such that
   \[ \mu(\mathcal{L} f) = \rho \mu(f) \]
   for every \( f \in BV \). Furthermore, this measure is atom-free and satisfies the property that
   \[ \mu(f) = \lim_{n \to \infty} \inf_{x \in D_n} \frac{\mathcal{L}^n f(x)}{\mathcal{L}^n 1(x)} \]

\[ \text{This is the choice made in [22].} \]
for every $f \in BV$, and $\rho = \mu(\mathcal{L}1)$.
3. The measure $\lambda = h\mu$ is, up to scalar multiples, the only $T$-invariant measure supported on $X_\infty$ and absolutely continuous with respect to $\mu$.
4. There exist $\kappa < 1$ and $C > 0$ such that for any function of bounded variation $f$,
   \[
   \left\| \frac{\mathcal{L}^n f}{\rho^n} - h\mu(f) \right\|_\infty \leq C\kappa^n \|f\|_{BV}.
   \]

Remark 2.8. It follows readily from the proof of Theorem 2.7 [22] that the same conclusion can be obtained if the hypothesis of $H_0$ being an admissible hole is replaced by $H_0$ being an $\epsilon$-admissible hole for some $\epsilon > 0$.

To close this section, we present a lemma concerning admissibility of different holes, obtained by enlarging an initial hole $H_0$ to $H_m := I \setminus X_m$. This broadens the applicability of Theorem 3.2 because enlarging the holes may reduce the number of contiguous bad intervals and also reduce the variation remaining on the domain of the open Lasota–Yorke map without decreasing the expansion.

Lemma 2.9 (enlarging holes). Let $T_0 = (\hat{T}, H_0)$ be an open system with an $\epsilon$-admissible hole, and for each $m \geq 0$ let $H_m := I \setminus X_m$. Then, for each $m \geq 0$, $T_m := (\hat{T}, H_m)$ is an open system with an $\epsilon$-admissible hole. Furthermore, let $\rho(T_m)$, $h(T_m)$, and $\mu(T_m)$ be the escape rate, accim, and quasi-conformal measures, respectively, of $T_m$. Then we have the following:
1. $\rho(T_m) = \rho(T_0)$,
2. $\mathcal{L}^m(h(T_m)) = \rho(T_0)^m h(T_0)$, and
3. $\mu(T_m) = \mu(T_0)$.

The proof of Lemma 2.9 is presented in section 5.3.


3.1. The Ulam scheme. In the case of a closed system $\hat{T}$, the well-known Ulam method introduced in [32] provides a way of approximating the transfer operator with a sequence of finite-rank operators $\hat{L}_k$ defined as in, e.g., [21], each coming from discretizing the interval $I$ into $k$ bins (which may or may not be of equal length). The only requirements are that each bin be a nontrivial interval, and that the maximum diameter of the partition elements, denoted by $\tau_k$, goes to 0 as $k$ goes to infinity. We call such a $k$-bin partition $\mathcal{P}_k$. The operator $\hat{L}_k$ preserves the $k$-dimensional subspace span$\{\chi_j : \chi_j = 1_{I_j}, I_j \in \mathcal{P}_k\}$. The matrix $\hat{P}_k$ defined in the introduction represents the action of $\hat{L}_k$ on this space, with respect to the ordered basis $(\chi_1, \ldots, \chi_k)$ [21].

In the case of an open system $(\hat{T}, H_0)$, one can still follow Ulam’s approach to define a discrete approximation $\hat{L}_k$ to the transfer operator $\mathcal{L}$. For a function $f \in BV$, the operator is defined by $\hat{L}_k(f) = \pi_k(\mathcal{L}f) = \pi_k \hat{L}(1_{X_0} f)$, where $\pi_k$ is given by the formula
   \[
   \pi_k(f) = \sum_{j=1}^k \frac{1}{m(I_j)} \left( \int \chi_j f \, dm \right) \chi_j.
   \]

The entries of the Ulam transition matrix $P_k$ representing $\hat{L}_k$ in the ordered basis $(\chi_1, \ldots, \chi_k)$ are
   \[
   (P_k)_{ij} = \frac{m(I_i \cap X_0 \cap \hat{T}^{-1} I_j)}{m(I_j)}.
   \]
Let \( h \) be densities on \( I \) according to the formula
\[
\chi_k = \sum_{j=1}^{k} [P_k]_{j} \chi_j,
\]
where we adopt the convention that a vector \( \chi \) can be written in component form as \( \chi = ([\chi_1, \ldots, [\chi_k]) \). Nonnegative right eigenvectors \( \psi_k \) of \( P_k \) induce measures \( \mu_k \) on \( I \) according to the formula
\[
\mu_k(E) = \sum_{j=1}^{k} \psi_k j \ m(I_j \cap E).
\]

We conclude the section with the following.

**Lemma 3.1.** Let \( P_k \) be the matrix representation of \( \mathcal{L}_k = \pi_k \circ \mathcal{L} \) with respect to the basis \( \{ \chi_j \} \). If \( P_k \psi_k = \rho_k \psi_k \), then the measure \( \mu_k \) corresponding to \( \psi_k \) satisfies \( \mu_k(\mathcal{L}_k \pi_k \varphi) = \rho_k \mu_k(\varphi) \) for every \( \varphi \in L^1(m) \).

**Proof.** Let \( \varphi \in L^1(m) \), and set \( \varphi_k = \pi_k \varphi \). Then,
\[
\mu_k(\varphi) = \int \varphi \, dm = \sum_{j=1}^{k} \int_{I_j} \varphi \, dm \psi_k j = \sum_{j=1}^{k} \int_{I_j} \pi_k \varphi \, dm \psi_k j \]
\[
= \sum_{j=1}^{k} \int_{I_j} \varphi_k \, dm \, (P_k)_{jj} \psi_k j (\rho_k)^{-1} = \sum_{j=1}^{k} \int_{I_j} \mathcal{L}_k \varphi_k \, dm \psi_k j (\rho_k)^{-1} \]
\[
= (\rho_k)^{-1} \int \mathcal{L}_k \varphi_k \, dm = \rho_k^{-1} \mu_k(\mathcal{L}_k \varphi_k),
\]
where the last equality in the second line follows from the fact that \( P_k \) is the matrix representing \( \mathcal{L}_k \) in the basis \( \{ \chi_j \} \) and acts on densities by right multiplication (i.e., if \( p \) is the vector representing the function \( \varphi_k \), then \( p \mathcal{L}_k \) is the vector representing \( \mathcal{L}_k \varphi_k \)).

### 3.2. Statement of the main result.

The main result of this paper is the following. Its proof is presented in section 5.2.

**Theorem 3.2.** Let \( \tilde{T} : I \supset \circ \) be a Lasota–Yorke map with an \( \epsilon \)-Ulam-admissible hole \( H_0 \). Let \( h \in BV \) be the unique accin for the open system \((\tilde{T}, H_0)\), and \( \mu \) the unique quasi-conformal measure for the open system supported on \( X_\infty \), as guaranteed by Theorem 2.7. Let \( \rho \) be the associated escape rate. For each \( k \in \mathbb{N} \), let \( \rho_k \) be the leading eigenvalue of the Ulam matrix \( P_k \). Let \( h_k \) be densities induced from nonnegative right eigenvectors of \( P_k \) corresponding to \( \rho_k \). Let \( \mu_k \) be measures induced from nonnegative right eigenvectors of \( P_k \) corresponding to \( \rho_k \).

Then, the following hold:

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Note: The above text is a snippet from a scientific article, focusing on the derivation of certain properties and results related to Lasota–Yorke maps and Ulam matrices.
Thus, the elements of \( \{ \epsilon (1 + \alpha) \} \) and there exists \( \eta \in (0, 1) \) such that \( |\rho_k - \rho| \leq O(\tau_k^\eta) \), where \( \tau_k \) is the maximum diameter of the elements of \( P_k \).

(2) \( \lim_{k \to \infty} h_k = h \) in \( L^1(m) \).

(3) \( \lim_{k \to \infty} h_k = \mu \) in the weak-* topology of measures. Furthermore, for every sufficiently large \( k \), \( \supp(\mu) \subseteq \supp(\mu_k) \).

We will also establish a relation between admissibility and Ulam-admissibility of holes.

**Lemma 3.3 (admissibility and Ulam-admissibility).** If \( H_0 \) is an \( \epsilon \)-admissible hole for \( \hat{T} \), there is some \( n \in \mathbb{N} \) such that \( H_{n-1} := I \setminus X_{n-1} \) is \( \epsilon \)-Ulam-admissible for \( \hat{T}^n \).

The proof of this lemma is presented in section 5.4. This result, together with Lemma 2.9, broadens the scope of applicability of Theorem 3.2 by allowing us to (i) replace the map by an iterate (Lemma 3.3), or (ii) enlarge the hole in a dynamically consistent way (Lemma 2.9). It also ensures that several examples in the literature can be treated with our method; in particular, all the examples presented in [22].

**Remark 3.4.** In the case of full-branched maps (see section 4.1 for a precise definition), the value of \( \alpha \) in (2.1) can be replaced by \( \|DT^{-1}\|_\infty (1 + \epsilon + \zeta_\lambda) \), and we can still obtain the conclusions of Theorem 3.2. Essentially this is because, taking the usual approach of considering \( T \) bivalued at the endpoints of the monotonicity partition, in the full-branched case one can regard \( g \) as continuous on monotonicity intervals and hence find a finite partition \( \mathcal{Z}_g \) such that for every \( A \in \mathcal{Z}_g \), \( \text{var}_A(g) \leq \epsilon \|DT^{-1}\|_\infty \). This is instead of the bound \( \text{var}_A(g) \leq (1 + \epsilon)\|DT^{-1}\|_\infty \), which is used in Lemma 5.1 just before (5.3).

4. **Examples.** To illustrate the efficacy of Ulam’s method beyond the small-hole setting, we present some examples of Ulam-admissible open Lasota–Yorke systems. We start with the case of full-branched maps in section 4.1, and treat some more general examples, including \( \beta \)-shifts, in section 4.2. We then analyze Lorenz-like maps, which provide transparent evidence of the scope of the results for open systems, as well as closed systems with repellers. They also illustrate how the admissibility hypothesis may be checked in applications.

4.1. **Full-branched maps.** In the examples ahead, we will use the following notation. Given a Lasota–Yorke map with holes, \( (\hat{T}, H_0) \) with monotonicity partition \( \mathcal{Z} \), we let \( \mathcal{Z}_H = \{ Z \in \mathcal{Z} : Z \subseteq H_0 \} \), \( \mathcal{Z}_f = \{ Z \in \mathcal{Z} : Z \cap H_0 = \emptyset, T(Z) = I \} \), and \( \mathcal{Z}_a = \{ Z \in \mathcal{Z} : Z \notin \mathcal{Z}_H \cup \mathcal{Z}_f \} \). Thus, the elements of \( \mathcal{Z}_f \) are precisely those contained in \( X_0 \) that are full branches for \( T \), and those of \( \mathcal{Z}_a \) are the remaining ones.

**Definition 4.1.** A full-branched map with holes, \( (\hat{T}, H_0) \), is a Lasota–Yorke map with holes such that \( \mathcal{Z}_a = \emptyset \).

For piecewise linear maps, the situation is rather simple.

**Lemma 4.2.** Let \( T_0 = (\hat{T}, H_0) \) be a piecewise linear full-branched map with holes. Then,

\[ \lim_{k \to \infty} \rho_k = \rho, \]

where \( \rho \) is a simple eigenvalue for \( P_k \). Furthermore,

\[ \lim_{k \to \infty} h_k = h \text{ in } L^1(m). \]

\[ \lim_{k \to \infty} h_k = \mu \text{ in the weak-* topology of measures. Furthermore, for every sufficiently large } k, \supp(\mu) \subseteq \supp(\mu_k). \]

\[ \text{We will also establish a relation between admissibility and Ulam-admissibility of holes.} \]

**Lemma 3.3 (admissibility and Ulam-admissibility).** If \( H_0 \) is an \( \epsilon \)-admissible hole for \( \hat{T} \), there is some \( n \in \mathbb{N} \) such that \( H_{n-1} := I \setminus X_{n-1} \) is \( \epsilon \)-Ulam-admissible for \( \hat{T}^n \).

The proof of this lemma is presented in section 5.4. This result, together with Lemma 2.9, broadens the scope of applicability of Theorem 3.2 by allowing us to (i) replace the map by an iterate (Lemma 3.3), or (ii) enlarge the hole in a dynamically consistent way (Lemma 2.9). It also ensures that several examples in the literature can be treated with our method; in particular, all the examples presented in [22].

**Remark 3.4.** In the case of full-branched maps (see section 4.1 for a precise definition), the value of \( \alpha \) in (2.1) can be replaced by \( \|DT^{-1}\|_\infty (1 + \epsilon + \zeta_\lambda) \), and we can still obtain the conclusions of Theorem 3.2. Essentially this is because, taking the usual approach of considering \( T \) bivalued at the endpoints of the monotonicity partition, in the full-branched case one can regard \( g \) as continuous on monotonicity intervals and hence find a finite partition \( \mathcal{Z}_g \) such that for every \( A \in \mathcal{Z}_g \), \( \text{var}_A(g) \leq \epsilon \|DT^{-1}\|_\infty \). This is instead of the bound \( \text{var}_A(g) \leq (1 + \epsilon)\|DT^{-1}\|_\infty \), which is used in Lemma 5.1 just before (5.3).

4. **Examples.** To illustrate the efficacy of Ulam’s method beyond the small-hole setting, we present some examples of Ulam-admissible open Lasota–Yorke systems. We start with the case of full-branched maps in section 4.1, and treat some more general examples, including \( \beta \)-shifts, in section 4.2. We then analyze Lorenz-like maps, which provide transparent evidence of the scope of the results for open systems, as well as closed systems with repellers. They also illustrate how the admissibility hypothesis may be checked in applications.

4.1. **Full-branched maps.** In the examples ahead, we will use the following notation. Given a Lasota–Yorke map with holes, \( (\hat{T}, H_0) \) with monotonicity partition \( \mathcal{Z} \), we let \( \mathcal{Z}_H = \{ Z \in \mathcal{Z} : Z \subseteq H_0 \} \), \( \mathcal{Z}_f = \{ Z \in \mathcal{Z} : Z \cap H_0 = \emptyset, T(Z) = I \} \), and \( \mathcal{Z}_a = \{ Z \in \mathcal{Z} : Z \notin \mathcal{Z}_H \cup \mathcal{Z}_f \} \). Thus, the elements of \( \mathcal{Z}_f \) are precisely those contained in \( X_0 \) that are full branches for \( T \), and those of \( \mathcal{Z}_a \) are the remaining ones.

**Definition 4.1.** A full-branched map with holes, \( (\hat{T}, H_0) \), is a Lasota–Yorke map with holes such that \( \mathcal{Z}_a = \emptyset \).

For piecewise linear maps, the situation is rather simple.

**Lemma 4.2.** Let \( T_0 = (\hat{T}, H_0) \) be a piecewise linear full-branched map with holes. Then,
for every \( \epsilon > 0 \) the following holds: \( \xi_{\epsilon}(T_0) = 0 \),

\[
\rho(T_0) = 1 - \text{Leb}(H_0), \quad \text{and} \quad \alpha_\epsilon(T_0) = \max_{Z \in Z_f} \text{Leb}(Z)(1 + \epsilon).
\]

**Proof.** If \( T_0 \) is a piecewise linear full-branched map, then each interval \( Z \in Z_f \) is good. Therefore \( \xi_\epsilon(T_0) = 0 \). Also,

\[
\mathcal{L}_0(1)(x) = \sum_{y \in Z \in Z_f, T_0(y) = x} \frac{1}{|DT_0(y)|} = \sum_{Z \in Z_f} \text{Leb}(Z) = 1 - \text{Leb}(H_0),
\]

which yields the first claim. The second statement follows from Remark 3.4 and the fact that 
\[
\sup_{x \in Z \in Z_f} \frac{1}{|DT_0(x)|} = \max_{Z \in Z_f} \text{Leb}(Z). \]

In fact, in the piecewise linear, full-branched setting, a direct calculation shows that the
Lebesgue measure is an accim for the open system. For perturbations of these systems, explicit
estimates of \( \rho \) and \( \alpha_\epsilon \) are not generally available. However, we have the following bounds.

**Lemma 4.3.** Let \( T_0 = (T, H_0) \) be a full-branched map with holes. Then, for every \( \epsilon > 0 \)
there exists some computable \( m \in \mathbb{N} \) such that \( \xi_{\epsilon}(T_m) = 0 \), where \( T_m := (\hat{T}, H_m) \) is obtained
from \( T_0 \) by enlarging the hole, as in Lemma 2.9. Furthermore,

\[
\rho(T_m) = \rho(T_0) \geq \inf_{x \in I} \sum_{y \in Z \in Z_f, T_0(y) = x} \frac{1}{|DT_0(y)|} =: \rho_0 \quad \text{and} \quad \alpha_{\epsilon}(T_m) \leq \sup_{x \in Z \in Z_f} \frac{1}{|DT_0(x)|}(1 + \epsilon) =: \alpha_{\epsilon,0}.
\]

An immediate consequence is the following.

**Corollary 4.4.** In the setting of Lemma 4.3, if \( \rho_0 > \alpha_{\epsilon,0} \), then \( H_m \) is \( \epsilon \)-Ulam admissible
for \( \hat{T} \). In this case, Lemma 2.9 allows one to approximate the escape rate, accim, and quasiconformal measure
for \( T_0 \) via Theorem 3.2 applied to \( T_m \).

**Proof of Lemma 4.3.** First, let us note that for any map with \( Z_f \neq \emptyset \) we have that \( D_\infty \neq \emptyset \),
as the map has at least one fixed point outside the hole. If \( m \) is sufficiently large, each interval
\( Z \in Z^{(m)} \) is either (i) contained in \( H_{m^{-1}} \), and thus not in \( Z^{*(m)} \), or (ii) \( T_m(Z) = I \) and
\( \text{var}_Z(g_{1X_m}) < \|g_{1X_m}\|_\infty(1 + \epsilon) \). In the latter case, \( Z \) is a good interval for \( T_0 \), because
\( \mu_0(Z) = \rho_0^{-m} \mu_0(\mathcal{L}_0^{m}(1)) > \rho_0^{-m} \|DT_0^m\|_\infty \mu_0(I) > 0 \). Since good intervals for \( T_0 \) and for \( T_m \)
coincide (see the beginning of the proof of Lemma 2.9), we get that \( \xi_{\epsilon}(T_m) = 0 \).

Furthermore,

\[
\rho(T_0) = \rho(T_0)\mu_0(1) = \mu_0(\mathcal{L}_0(1)) \geq \inf_{x \in I} \mathcal{L}_0(1)(x) = \inf_{x \in I} \sum_{y \in Z \in Z_f, T_0(y) = x} \frac{1}{|DT_0(y)|}.
\]

The bound on \( \alpha_{\epsilon}(T_m) \) follows directly from Remark 3.4.

The following is an interesting consequence of Lemmas 4.2 and 4.3.

**Corollary 4.5.** Let \( (T, H_0) \) be a piecewise linear full-branched map with holes and at least
two full branches. Thus, \( \text{Leb}(H_0) < 1 - \max_{Z \in Z_f} \text{Leb}(Z) \). Then, if \( \epsilon > 0 \) is sufficiently small,
\( H_0 \) is \( \epsilon \)-Ulam-admissible for any full-branched map \((\hat{S}, H_0)\) that is a sufficiently small \( C^{1+\text{Lip}} \) perturbation of \((T, H_0)\) (where the \( C^{1+\text{Lip}} \) topology is defined, for example, by the norm given by the maximum of the \( C^{1+\text{Lip}} \) norms of each branch). In particular, Theorem 3.2 applies.

**Proof.** The statement for \((T, H_0)\) follows from Lemma 4.2. For perturbations, the statement follows from Lemma 4.3, by observing that the quantities \( \rho_0 \) and \( \alpha_{\epsilon,0} \), as well as the variation of \( 1/|D\hat{T}| \) on each interval, depend continuously on \( \hat{T} \), with respect to the \( C^{1+\text{Lip}} \) topology.

Corollary 4.5 can apply to maps with arbitrarily large holes, as the next example shows.

**Example 4.6 (arbitrarily large holes).** Let \( \delta > 0, H_0 = [\delta, 1 - \delta] \), and

\[
T_\delta(x) = \begin{cases} 
\delta^{-1}x & \text{if } x < \delta, \\
\delta^{-1}(1 - x) & \text{if } 1 - \delta \leq x \leq 1.
\end{cases}
\]

Then, \( \text{Leb}(H_0) = 1 - 2\delta < 1 - \delta = 1 - \max_{Z \in \mathcal{Z}} \text{Leb}(Z) \), and the hypotheses of Corollary 4.5 are satisfied. Thus, Ulam’s method converges for sufficiently small \( C^{1+\text{Lip}} \) perturbations of \( T_\delta \) that are full-branched.

**Remark 4.7.**

(I) It is worth noting that if in Example 4.6 the hole is enlarged to \([\delta, 1]\), neither the hypotheses of Corollary 4.5 nor the results of [22] apply. This corresponds to a degenerate setup where the survivor set consists of a single point. In this case, the Ulam method could still be implemented. The leading left eigenvectors would successfully approximate an accim, which is uniform with escape rate \( \delta \). However, the corresponding measures induced from the right eigenvectors would converge in the weak-* topology of \( C(I) \) to an (invariant) atomic measure at 0, instead of to a quasi-conformal measure, as the partition is refined. This simple example illustrates that there are obstacles to applying Theorem 3.2 if the hypotheses are weakened.

(II) Example 4.6 displays the potential misalignment between statistical and topological features of open dynamical systems: as \( \delta \) is varied, the maps \( T_\delta \) are all topologically conjugate to one another, yet each \( \delta \) has a unique natural escape rate. (As \( \delta \to 1/2 \) these rates approach 0.) Nonetheless, each map also supports an uncountable number of accims for each \( \rho \in (0, 1) \) [9, section 3], but the densities of these measures do not have bounded variation and are therefore undetectable by our methods.

Other examples of this type may be found in [1] and [2]. Bahsoun [1] established rigorous computable bounds for the errors in the Ulam method, which allowed him to find rigorous bounds on the escape rate for open Lasota–Yorke maps. Bahsoun and Bose [2] related the escape rate to the Lebesgue measure of the hole. Both results rely on the existence of Lasota–Yorke-type inequalities, relating \( BV \) and \( L^1(m) \) norms. Such inequalities may be obtained by exploiting the full-branched structure of the map.

**Example 4.8 (Bahsoun [1]).** Let

\[
\hat{T}(x) = \begin{cases} 
2.08x & \text{if } x < \frac{1}{2}, \\
2 - 2x & \text{if } x \geq \frac{1}{2}.
\end{cases}
\]

In this case, Corollary 4.5 applies. In fact, \( \text{Leb}(H_0) = \frac{0.8}{1.6} \) and \( 1 - \max_{Z \in \mathcal{Z}} \text{Leb}(Z) = \frac{1}{2} \). We
note that $\rho$ controls the rate of mass loss, which is slower than 4.08/4.16, while $\alpha_c$ is related to the relaxation rate on the survivor set.

4.2. Nearly piecewise linear maps with enough full branches. When nonfull branches are present, the dynamics is typically non-Markovian. Thus, even in the piecewise linear setting there may not be direct ways to find the various objects of interest (escape rates, accims and conformal measures) exactly. We show that Ulam’s method provides rigorous approximations in specific systems. The following example is closely related to [22, Definition 6.2 and Lemma 6.3].

Lemma 4.9. Let $T = (\hat{T}, H_0)$ be a piecewise linear Lasota–Yorke map with holes, and assume $Z_f \neq \emptyset$. Let $c_u$ be the maximum number of contiguous elements in $Z_u$. If $\|DT^{-1}\|_\infty (3 + c_u) < \rho$, then $H_0$ is $(1 + \epsilon)$-Ulam-admissible for $\hat{T}$ for every $\epsilon > 0$ sufficiently small. Thus, the hypotheses of Theorem 3.2 are satisfied.

Proof. For any map with $Z_f \neq \emptyset$ we have that $D_\infty \neq \emptyset$, as the map has at least one fixed point outside the hole. Furthermore, for each $Z \in Z$ one has that $\var Z(g) \leq 2\|g\|_\infty$, so $Z$ is a $(1 + \epsilon)$-adequate partition for $T$. Also, it follows from the definition of $Z_0$ that $Z_f \subseteq Z_g$. Thus, $Z_h \subseteq Z_u$ and $\epsilon_{1+\epsilon} \leq c_u$. Therefore, $\alpha_c \leq \|DT^{-1}\|_\infty (3 + \epsilon + c_u) < \rho$, provided that $\epsilon > 0$ is sufficiently small. \hfill \blacksquare

A concrete example to which the previous lemma applies is that of $\beta$-shifts.

Example 4.10. Let $\beta > 1$, and let $\hat{T}_\beta$ be the $\beta$-shift, $\hat{T}_\beta(x) = \beta x \pmod{1}$. Let $H_0 \subset I$ be a finite union of closed intervals, and let $f$ be the number of full branches of $\hat{T}_\beta$ outside $H_0$. Then, for the open system $(\hat{T}_\beta, H_0)$, we have that $\rho \geq \frac{f}{\beta}$. Then, the hypotheses of Lemma 4.9 are satisfied, provided $f > 3 + c_u$. This happens, for example, when $\beta \geq 5$ and $H_0$ is a single interval of the form $[\frac{[\beta]}{\beta}, y]$ or $[y, 1]$, with $\frac{[\beta]}{\beta} < y < 1$. It also happens when $\beta \geq 6$ and $H_0$ is a single interval contained in $[\frac{[\beta]}{\beta}, 1]$, or when $\beta \geq 7$ and $H_0$ is any interval leaving at least seven full branches in $X_0$. (Recall from subsection 2.1 that two bad elements of $Z_h$ are contiguous if there are no good elements of $Z_f \cup Z_u$ between them, but there may be elements of $Z_h$ in between.)

We include Figures 1–3, obtained from numerical experiments for $\beta = 5.9$ and two different choices of holes. They include approximations to the densities of accims and cumulative distribution functions of the quasi-conformal measures for systems with a hole, as well as the accim and conformal measure for the closed system.

Remark 4.11. Using lower bounds on $\rho$ such as those of Lemma 4.3, one can extend the conclusion of Lemma 4.9 as in Corollary 4.5, to cover small $C^{1+\text{Lip}}$ perturbations of piecewise linear maps that respect the partition $Z_h \cup Z_f \cup Z_u$.

4.3. Lorenz-like maps. Let us consider the following two-parameter family of maps of $I = [-1, 1]$: 

\begin{equation}
T_{c,\alpha}(x) = \begin{cases} 
\alpha x^\alpha - 1 & \text{if } x > 0, \\
1 - c|x|^\alpha & \text{if } x < 0,
\end{cases}
\end{equation}

where $c > 0$, $\alpha \in (0, 1)$. When $c > 2$, the system is open, and the hole is implicitly defined as $H_{0, c, \alpha} = T_{c, \alpha}^{-1}(\mathbb{R} \setminus [-1, 1])$. 


This family of maps has been studied in connection with the famous Lorenz equations,

\begin{align}
\dot{x} &= \sigma(y - x), \\
\dot{y} &= rx - y - xz, \\
\dot{z} &= -bz + xy.
\end{align}

(4.2)

We take a relatively standard point of view [19, 30, 16], regarding \(\sigma = 10\) and \(b = 8/3\) as fixed and \(r\) as a parameter. The chaotic attractor discovered by Lorenz [23] at \(r = 28\) has since been proved to exist by Tucker [31] (via computer-assisted methods). Its formation is now well understood: A homoclinic explosion occurs at \(r_{\text{hom}} \approx 13.9265\), giving rise to
a chaotic saddle. As \( r \) increases through \( r_{het} \approx 24.0579 \), heteroclinic connections between \((0,0,0)\) and a symmetric pair of periodic orbits \( \Gamma^\pm \) appear, and the chaotic saddle becomes an attractor \( \Omega \). The orbits \( \Gamma^\pm \) disappear in subcritical Hopf bifurcations at \( r_{Hopf} \approx 24.7368 \) (parameter values from [11]). For \( r < r_{het} \) almost all orbits are asymptotic to one of two fixed points; for \( r_{het} < r < r_{Hopf} \) orbits may approach one of these fixed points or the attractor \( \Omega \); for \( r > r_{Hopf} \) almost all orbits are attracted to \( \Omega \).

Maps like (4.1) model this situation via the following reductions. First, solutions to the ODEs (4.2) induce a flow on \( \mathbb{R}^3 \); from this, a return map to the section \( \Sigma = \{(x,y,z) : z = r - 1\} \) may be constructed. This two-dimensional map is an open dynamical system, since not all orbits return to \( \Sigma \).\(^7\) For “preturbulent” \( r \in (r_{hom},r_{het}) \), the chaotic saddle admits a strong stable foliation; the return map to \( \Sigma \) may be further reduced by identifying points on the same stable leaf, resulting in one-dimensional models. We illustrate our results with the much-studied family (4.1) (see [16]). The discontinuity at \( x = 0 \) corresponds to the intersection of the stable manifold of \((0,0,0)\) with \( \Sigma \); the exponent \( 0 < \alpha < 1 \) is derived from the eigenvalues of the linearization of the system at the origin, \( \alpha = |\lambda_s|/\lambda_u \). The parameter \( c \) controls how “open” the map is: when \( c \leq 2 \), the system is closed, and when \( c > 2 \), the one-step survivor set \( X_0 \) has the form \( X_0 = [-x_{c,\alpha},x_{c,\alpha}] \), where \( x_{c,\alpha} = (2/c)^{1/\alpha} \); this is illustrated in red in Figure 4.

The escape rates of the system \( T_{c,\alpha} \) for parameters \( 0 < \alpha < 1 \), \( 2 < c < 3 \) are illustrated in Figure 5. Figure 6(left) illustrates the cumulative distribution functions of the quasi-conformal measures, \( \mu_{c,\alpha} \), for \( c = 2.01 \) and various values of \( \alpha \). The densities of the accims with respect to Lebesgue are illustrated in Figure 6 for several \( \alpha \) values. For \( \alpha < 0.5 \), the densities become concentrated near the endpoints, as the \( \alpha = 0.45 \) plot in Figure 6(right) illustrates.

The escape rate results for these one-dimensional maps can be interpreted coherently with respect to the behavior of the Lorenz system (4.2) (although the scenarios differ according to

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\(^7\)For example, the stable manifold to the fixed point \((0,0,0)\) intersects \( \Sigma \), and some orbits of the flow travel directly to \((0,0,0)\) after leaving \( \Sigma \).
Regarding $T_{c,\alpha}$ as a map on $\mathbb{R}$, for each value of $\alpha \in (0, 1)$ and $c > 2$ there are two pairs of fixed points: repellers at $\pm y_{c,\alpha} \in (-1, 1)$ (illustrated in green in Figure 4) and an attracting outer pair $\pm z_{c,\alpha}$ with $|z_{c,\alpha}| > 1$ (beyond the domain of Figure 4). The inner points $\pm y_{c,\alpha}$ correspond to the periodic orbits $\Gamma^\pm$ from the Lorenz flow, and the outer pair correspond to the attracting fixed points of the flow.

- At some $c = c_*(\alpha) \leq 2$ the inner and outer pairs coalesce in a saddle-node bifurcation,
and for \( c < c_* \) the only attractor is a chaotic absolutely continuous invariant measure supported on \([-1, 1]\).

\((\alpha > 1/2)\) Each \( T_{c,\alpha} \) is uniformly expanding on \( X_0 \) for \( c > 2 \). For \( c > 2 \) there is a chaotic repeller in \( X_0 \) and a fully supported acim on \([-1,1]\). Lebesgue almost every orbit escapes and is asymptotic to one of the “outer fixed points.” At \( c = 2 \) the points \( \pm x_{c,\alpha} = \pm 1 = \pm y_{c,\alpha} \) become fixed points, with \( T'(\pm 1) = 2\alpha > 1 \). The open system thus “closes up” as \( c \) decreases to 2; this corresponds to the bifurcation point \( T_{\text{bif}} \) in the Lorenz flow (where the origin connects to \( \Gamma^\pm \)). For values of \( c < 2 \), \( T_{c,\alpha} \) admits an acim (which can be accessed numerically by Ulam’s method), and the quasi-conformal measure is simply the Lebesgue measure. The approach of \( \rho_k \) to 1 as \( c \to 2 \) can be seen in Figure 5, and the close agreement of \( \mu_{2,0.01,\alpha} \) with the Lebesgue measure can be seen in Figure 6(left) for \( \alpha = 0.95 \).

\((\alpha < 1/2)\) For \( c > 2 \), \( T_{c,\alpha} \) is open on \([-1,1]\), but the uniform expansion property fails for \( c \) sufficiently close to 2. Indeed, when \( c = 2 \) the fixed points at \( \pm 1 \) are the \emph{outer pair} \( \pm z_{c,\alpha} \) and \( T'(\pm 1) < 1 \). For \( c \in (c_*,2) \), these attractors \( \pm z_{c,\alpha} \in [-1,1] \) and \emph{coexist with a chaotic repeller in} \([-y_{c,\alpha},y_{c,\alpha}]\). Fortunately, for \( c > c_* \) the open system \( T_{c,\alpha} \) with hole \( I \setminus [-y_{c,\alpha},y_{c,\alpha}] \) is a Lasota–Yorke map with holes, because it is piecewise expanding. Corollary 4.4 shows that \( \varepsilon \)-Ulam admissibility of the open system is implied if \( \varepsilon \) is sufficiently small and

\[
|T'_{c,\alpha}(y_{c,\alpha})|^{-1} = \sup_{x \in [-y_{c,\alpha},y_{c,\alpha}]} |T'_{c,\alpha}(x)|^{-1} < \inf L_{c,\alpha} \mathcal{L}(x).
\]

This condition can be verified directly via elementary calculus. Thus, for \( c > 2 \), our main theorem holds for the application of Ulam’s method to \( T_{c,\alpha} \) on \([-y_{c,\alpha},y_{c,\alpha}]\). However, it is simple to extend this result to \([-1,1]\): all points in the intervals \( \pm (y_{c,\alpha},1) \) escape in finitely many iterations, and corresponding cells of the partitions used in Ulam’s method are “transient.” The leading eigenvalue from Ulam’s method and approximate quasi-conformal measure on \([-y_{c,\alpha},y_{c,\alpha}]\) agree with those computed on
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The approximate accims agree (modulo scaling) between ±yc,α; the only difference is that the different X0’s lead to different concentrations of mass on preimages of the hole. The approximated escape rates are displayed in Figure 5, and concentration of accim on the hole (neighborhoods of ±1) is evident in Figure 6(right). Note also that Figure 6(left) depicts some approximate quasi-conformal measures for c = 2.01 and α < 0.45.

Remark 4.12. Recent work on Lorenz-like systems [24] has focused on Lorenz maps with less regularity, such as piecewise C1+ϵ. We expect that our approach could be extended to this setting, although some technical modifications would be necessary.

5. Proofs.

5.1. Auxiliary lemmas. Under the assumptions of Theorem 2.7, the quasi-conformal measure μ of (ˆT,H0) satisfies some further properties that will be exploited in our approach. The measure μ can be used to define a useful cone of functions in BV. For each a > 0 let

\[ C_a = \{0 \leq f \in BV : \text{var}(f) \leq a\mu(f)\}. \]

Combining the result of Lemmas 4.2 and 4.3 from [22] with the argument in the proof of Lemma 3.7 (therein), the conditions on T imply the existence of a constant a1 > 0 such that for any a > a1 there are an \( \epsilon_a > 0 \) and \( N \in \mathbb{N} \) such that

\[ (5.1) \quad \mathcal{L}^N C_a \subseteq C_a - \epsilon_a. \]

The values of N, a1, and \( \epsilon_a \) are all computable in terms of the constants associated with T. We present a modified version of these arguments, based on the classical work of Rychlik [29], that specialize to the case \( N = 1 \) and allow us to improve some of the constants involved in the estimates of [22]. Most notably, the value of \( \alpha_\epsilon \) below is smaller than that in [22], a fact which will allow us to treat a larger class of open systems.

Lemma 5.1. Let \( (\hat{T},H_0) \) be a Lasota–Yorke map with an \( \epsilon \)-Ulam-admissible hole. Then there exists \( K_\epsilon > 0 \) such that for every \( f \in BV \),

\[ \text{var}(\hat{L}f) \leq \alpha_\epsilon \text{var}(f) + K_\epsilon \mu(|f|). \]

Furthermore, there is a constant a1 > 0 such that for any a > a1 there is an \( \epsilon_a > 0 \) such that

\[ (5.2) \quad \mathcal{L} C_a \subseteq C_a - \epsilon_a. \]

Proof. We address the general case first; the particular full-branched case will be addressed at the end of the proof. In the general case, we recall that \( \alpha_\epsilon = \|DT^{-1}\|_\infty(2 + \epsilon + \xi_\epsilon) \).

Let \( \mathcal{Z} \) be the monotonicity partition for \( \hat{T} \). Define \( \hat{g} : I \to \mathbb{R} \) by \( \hat{g}(x) = |D\hat{T}(x)|^{-1} \) for every \( x \in (I \setminus \bigcup_{Z \in \mathcal{Z}} \partial Z) \cup \{0,1\} \), and \( \hat{g}(x) = 0 \) otherwise. We obtain the following Lasota–Yorke inequality by adapting the approach of Rychlik [29, Lemmas 4-6]. Let \( \mathcal{Z}_\epsilon \in \mathcal{G}_\epsilon \). Then,

\[ \text{var}(\hat{L}f) \leq \text{var}(f\hat{g}) \leq (2 + \epsilon)\|D\hat{T}^{-1}\|_\infty \text{var}(f) + \|D\hat{T}^{-1}\|_\infty(1 + \epsilon) \sum_{A \in \mathcal{Z}_\epsilon} \inf_A |f|. \]
We slightly modify \( \hat{g} \) to account for the jumps at the hole \( H_0 \), and define \( g : I \to \mathbb{R} \) by \( g = 1_{X_0} \hat{g} \). Now, only elements of \( \mathcal{Z}_e^* \) contribute to the variation of \( \mathcal{L}f \), and we get

\[
\text{var}(\mathcal{L}f) = \text{var}(\mathcal{L}(1_{X_0}f)) \leq \text{var}(f(1_{X_0} \hat{g})) = \sum_{A \in \mathcal{Z}_e^*} \text{var}(f(1_{X_0} \hat{g})) \\
\leq \sum_{A \in \mathcal{Z}_e^*} \text{var}(f) \|1_{X_0} \hat{g}\|_\infty + \|1_Af\|_\infty \text{var}(1_{X_0} \hat{g}) \\
\leq \sum_{A \in \mathcal{Z}_e^*} \text{var}(f) \|DT^{-1}\|_\infty + \left(\inf_A |f| + \text{var}(f)\right) \text{var}(g).
\]

Thus, since for every \( A \in \mathcal{Z}_e^* \), \( \text{var}_A(g) \leq \|DT^{-1}\|_\infty (1 + \epsilon) \), one has that

\[
(5.3) \quad \text{var}(\mathcal{L}f) \leq (2 + \epsilon)\|DT^{-1}\|_\infty \text{var}(f) + \sum_{A \in \mathcal{Z}_e^*} \|DT^{-1}\|_\infty (1 + \epsilon) \inf_A |f|.
\]

Now we proceed as in the proof of [22, Lemma 2.5] and observe that there exists \( \delta > 0 \) such that if \( A \in \mathcal{Z}_{\epsilon,g} \), then

\[
(5.4) \quad \inf_A |f| \leq \delta^{-1} \mu(1_A|f|),
\]

whereas if \( A \in \mathcal{Z}_{\epsilon,b} \), we let \( A' \in \mathcal{Z}_{\epsilon,g} \) be the nearest good partition element,\(^8\) and get

\[
\inf_A |f| \leq \inf_{A'} |f| + \text{var}_{I(A,A')} (f),
\]

where \( I(A,A') \) is an interval that contains \( A \) and has as an endpoint \( x_{A'} \in A' \), fixed in advance, such that, after possibly redefining \( f \) at the discontinuity points of \( f \), \( |f(x_{A'})| = \inf_{A'} |f| \). Notice that either \( I(A,A') \subseteq I_-(A') \) or \( I(A,A') \subseteq I_+(A') \), where \( I_+(A') \) is the union of \( A'_+ := A' \cap \{x : x \geq x_{A'}\} \) with the contiguous elements of \( \mathcal{Z}_{\epsilon,b} \) on the right of \( A' \), and \( I_-(A') \) is defined in a similar manner. Thus,

\[
(5.5) \quad \sum_{A \in \mathcal{Z}_{\epsilon,b}} \inf_A |f| \leq \xi_{\epsilon} \text{var}(f) + 2\xi_{\epsilon} \sum_{A' \in \mathcal{Z}_{\epsilon,g}} \inf_{A'} |f|,
\]

where the factor 2 appears due to the fact that a single good interval could have at most \( \xi_{\epsilon} \) bad intervals on the left and \( \xi_{\epsilon} \) bad intervals on the right. Combining (5.4) and (5.5), we get

\[
\sum_{A \in \mathcal{Z}_{\epsilon}^*} \inf_A |f| \leq \xi_{\epsilon} \text{var}(f) + \delta^{-1}(1 + 2\xi_{\epsilon}) \sum_{A' \in \mathcal{Z}_{\epsilon,g}} \mu(1_A|f|).
\]

Plugging this back into (5.3), we get

\[
\text{var}(\mathcal{L}f) \leq \|DT^{-1}\|_\infty (2 + \epsilon + \xi_{\epsilon}) \text{var}(f) + \|DT^{-1}\|_\infty (1 + \epsilon) \delta^{-1}(1 + 2\xi_{\epsilon}) \mu(|f|).
\]

\(^8\)It is shown in [22, Lemma 2.4] that whenever \( (T,H_0) \) is an open system with an admissible hole, then \( \mathcal{Z}_{\epsilon,g} \neq \emptyset \).
We get the first part of the lemma by choosing $K_\epsilon = \| DT^{-1} \|_\infty (1 + \epsilon) \delta^{-1} (1 + 2 \xi_\epsilon)$. For the second part, we recall that $\mu(\mathcal{L} f) = \rho \mu(f)$, so for every $f \in \mathcal{C}_a$ we have that
\[
\frac{\text{var}(\mathcal{L} f)}{\mu(\mathcal{L} f)} \leq \frac{\alpha_\epsilon}{\rho} a + \frac{K_\epsilon}{\rho}.
\]
Thus, $\mathcal{L} f \in \mathcal{C}_a$, provided that $a > \frac{K_\epsilon}{\rho} \alpha_\epsilon =: a_1$.

Moving toward a $BV, L^1(\text{Leb})$ Lasota–Yorke inequality, we have the following.

**Lemma 5.2.** Let $\zeta > 0$ be given. Then there is a constant $B_\zeta < \infty$ such that
\[
\mu(f) \leq B_\zeta |f|_1 + \zeta \text{var}(f)
\]
for $0 \leq f \in BV(I)$.

**Proof.** Let $Z^{(n)}$ be the $n$-fold monotonicity partition for $T_0$, where $n$ is such that $\mu(Z) < \frac{\zeta}{2}$ for all $Z \in Z^{(n)}$. This choice is possible in view of [22, Lemma 3.10]. Choose $k$ such that every subinterval of size $\frac{1}{k}$ intersects at most two such $Z$. Then, if $Y$ is any subinterval of length $1/k$, there are elements $Z_1, Z_2 \in Z^{(n)}$ such that $Y \subset Z_1 \cup Z_2$; hence $\mu(Y) < \zeta = \zeta k m(Y)$. Now let $\xi$ be a partition of $I$ into subintervals of length $1/k$, and set
\[
F = \sum_{Y \in \xi} \text{ess sup}_Y f 1_Y.
\]
Then, $f \leq F$ and $F - f \leq \sum_{Y \in \xi} V_Y(f) 1_Y$, where $V_Y(f)$ denotes the variation of $f$ inside the interval $Y$. Thus,
\[
|F - f|_1 \leq \sum_{Y \in \xi} V_Y(f) m(Y) \leq V_1(f)/k.
\]
We now estimate
\[
\int f \, d\mu \leq \int F \, d\mu = \sum_{Y \in \xi} \text{ess sup}_Y f \mu(Y) \leq \sum_{Y \in \xi} \text{ess sup}_Y f \zeta = \zeta k |F|_1 = \zeta k |f|_1 + \zeta k |F - f|_1 \leq \zeta k |f|_1 + \zeta V_1(f).
\]
Setting $B_\zeta = \zeta k$ completes the proof.

A direct consequence of Lemmas 5.1 and 5.2 is the following.

**Corollary 5.3.** Let $\alpha_\epsilon < \alpha < \rho$, where $\alpha_\epsilon$ is defined in (2.1). Then, there exists $K > 0$ such that
\[
\text{var}(\mathcal{L} f) \leq \alpha \text{var}(f) + K |f|_1. \tag{5.6}
\]

\[\text{For convenience, we have dropped the dependence of } \epsilon \text{ on } \alpha \text{ and } K. \text{ This should cause no confusion in what follows, as } \epsilon \text{ is fixed throughout the section.}\]
Proof. Let \( \zeta' = \frac{\alpha - \omega}{K' \alpha} \), where \( K' \) comes from Lemma 5.1. Let \( B_{\zeta'} \) be given by Lemma 5.2. Then, Lemma 5.1 ensures

\[
\var(\mathcal{L} f) \leq \alpha \var(f) + K'(B_{\zeta'} |f|_1 + \zeta' \var(f)) = \alpha \var(f) + K |f|_1.
\]

Another useful result regarding the relation between the Ulam approximations and the accm and quasi-conformal measure is the following.

Lemma 5.4. There exists \( n > 0 \) such that \( (P^n)_{ij} > 0 \) for all \( i, j \) satisfying \( \mu(I_i) > 0 \) and \( \int_I h \ dm > 0 \).

Proof. Fix \( i, j \) satisfying the hypotheses. By Theorem 2.7,

\[
\lim_{n \to \infty} \| (\mathcal{L}^n I_i) / \rho^n - \mu(I_i) h \|_\infty = 0.
\]

Choose \( n_{ij} \) large enough so that \( \int_I \mathcal{L}^N I_i \ dm > 0 \) for all \( N \geq n_{ij} \). Because there are a finite number of \( I_i \) and \( I_j \), we can set \( n = \max_{i,j} n_{ij} \) and obtain \( \int_I \mathcal{L}^n I_i \ dm > 0 \) for all \( i, j \) satisfying the hypotheses. Note that this implies \( \int_I (\pi_k \mathcal{L})^n I_i \ dm > 0 \) because the support of the integrand is possibly enlarged by taking Ulam projections. This now implies \( (P^n)_{ij} > 0 \). \( \square \)

5.2. Proof of the main result. The lemmas presented in section 5.1 allow us to derive parts (1) and (2) of Theorem 3.2 via the perturbative approach from [20]. Indeed, Theorem 2.7 shows that \( \rho > \alpha \) is the leading eigenvalue of \( \mathcal{L} \), and that it is simple. Furthermore, \( \mathcal{L}_k \) is a small perturbation of \( \mathcal{L} \) for large \( k \), in the sense that \( \sup_{\|f\|_{BV}=1} |(\mathcal{L}_k - \mathcal{L}) f|_1 \to 0 \) as \( k \to \infty \).

Indeed,

\[
\sup_{\|f\|_{BV}=1} |(\mathcal{L}_k - \mathcal{L}) f|_1 = \sup_{\|f\|_{BV}=1} |(\pi_k - I_d) \mathcal{L} f|_1 \leq \sup_{\|f\|_{BV}=\|\mathcal{L}\|_{BV}} |(\pi_k - I_d) f|_1 \\
\leq \|\mathcal{L}\|_{BV} \max_{I_j, \mathcal{L}k} m(I_j),
\]

and the latter is proportional to \( \tau_k \), the diameter of the partition, which tends to 0 as \( k \to \infty \).

Since \( \pi_k \) decreases variation [21], Corollary 5.3 implies the uniform inequality

\[
\text{var}(\mathcal{L}_k f) \leq \alpha \var(f) + K |f|_1 \quad \forall k \in \mathbb{N},
\]

which is the last hypothesis to check to be in the position to apply the perturbative machinery of [20]. In particular, this implies quasi-compactness and hence a spectral decomposition of \( \mathcal{L}_k \) acting on \( BV \). This result ensures that, for sufficiently large \( k \), \( \mathcal{L}_k \) has a simple eigenvalue \( \rho_k \) near \( \rho \), and its corresponding eigenvector \( h_k \in BV \) converges to \( h \) in \( L^1(\text{Leb}) \), giving the convergence statements in (1) and (2).

In order to show (3), we consider the operator \( \mathcal{L}_k := \mathcal{L}_k \circ \pi_k \). In view of Lemma 3.1, \( \mathcal{L}_k \mu_k = \rho_k \mu_k \) and \( \mathcal{L}_k h_k = \rho_k h_k \). As in the previous paragraph, one can check that \( \mathcal{L}_k \) is a small perturbation of \( \mathcal{L} \). In fact,

\[
\sup_{\|f\|_{BV}=1} |(\mathcal{L}_k - \mathcal{L}) f|_1 \leq 2 \max_{I_j, \mathcal{L}k} m(I_j) = 2 \tau_k.
\]
Also, the Lasota–Yorke inequality (5.6) holds with \(\mathcal{L}\) replaced by \(\mathcal{L}_k\). Thus, [20, Corollary 1] (see (iii) below) shows that, for large \(k\), \(\rho_k\) is the leading eigenvalue of \(\mathcal{L}_k\).

Let \(\Pi_k\) be the spectral projectors defined by

\[
\Pi_k := \frac{1}{2\pi i} \oint_{\partial B_\delta(\rho)} (z - \mathcal{L}_k)^{-1} \, dz,
\]

where \(\delta\) is small enough to exclude all spectra of \(\mathcal{L}\) apart from the peripheral eigenvalue \(\rho\). Also let

\[
\Pi_0 := \frac{1}{2\pi i} \oint_{\partial B_\delta(\rho)} (z - \mathcal{L})^{-1} \, dz.
\]

Then, [20, Corollary 1] provides \(K_1, K_2 > 0\), and \(\eta \in (0, 1)\) for which

(i) \(\| (\Pi_k - \Pi_0) f \|_1 \leq K_1 \tau_k^\eta \| f \|_{BV}\),

(ii) \(\| \Pi_k f \|_{BV} \leq K_2 \| \Pi_0 f \|_1\),

(iii) for large enough \(k\), \(\text{rank}(\Pi_k) = \text{rank}(\Pi_0)\).

Since \(\rho\) is simple and isolated, this setup implies that for large enough \(k\), each \(\Pi_k\) is a bounded, rank-1 operator on \(BV\),

\[
\Pi_k = \mu_k(\cdot) h_k,
\]

where each \(h_k \in BV\), \(\mathcal{L}_k h_k = \rho_k h_k\), and \(\rho_k \in B_\delta(\rho)\). Since \(h_k = \Pi_k h_k\), we can choose \(|h_k|_1 = 1\) so that \(\| h_k \|_{BV} \in [1, K_2]\). Now, let \(g \in BV\). Then, by the above,

\[
|\mu_k(g) - \mu(g)| = |(\mu_k(g) - \mu(g)) h_k|_1 \leq |\mu_k(g) h_k - \mu(g) h_k|_1 + |\mu(g) h_k - h_k|_1
\]

\[
= |\Pi_k(g) - \Pi_0(g)|_1 + |\mu(g)| |h_k - h_k|_1 \to 0 \quad \text{as } k \to \infty.
\]

Since \(\mu\) and \(\mu_k\) are in fact measures, the above is enough to show that \(\mu_k \to \mu\) in the weak-* topology.

In particular, there is a \(k_0\) such that \(\mu_k(h) > 0\) for all \(k \geq k_0\). To show the last claim of (3), we will show that if \(\mu_k(h) > 0\), then \(\text{supp}(\mu) \subseteq \text{supp}(\mu_k)\). Let \(\psi_k\) be a leading right eigenvector of \(\mathcal{L}_k\) such that \(P_k \psi_k = \rho_k \psi_k\) and \(|\psi_k|_1 \geq \frac{\mu(I_i)}{m(I_i)} (l = 1, \ldots, k)\). Choose \(i\) such that \(\mu(I_i) > 0, j\) such that \(|\psi_k|_j = \int_{I_j} h \, dm = \int_{I_j} h \, d\mu_k > 0\), and \(n \geq n_{ij}\) as in Lemma 5.4. Then,

\[
|\psi_k|_l = \rho^{-(n)} |P^n_k \psi_k|_l \geq \rho^{-(n)} |P^n_k|_{ij} |\psi_k|_j > 0.
\]

This establishes that \(\mu_k(I_i) > 0\) and hence that \(\text{supp}(\mu) \subseteq \cup \{I_i : \mu(I_i) > 0\} \subset \text{supp}(\mu_k)\), as claimed.

For the quantitative statement of (1), note that for every \(f \in BV\), \(0 = (\mathcal{L} - \rho I) h = (\mathcal{L} - \rho I) \Pi_0 f\), so that

\[
(\rho_k - \rho) h_k = (\mathcal{L}_k - \mathcal{L}) h_k + (\mathcal{L} - \rho)(\Pi_k - \Pi_0) h_k.
\]

Hence,

\[
|\rho_k - \rho| |h_k|_1 \leq 2 \tau_k \| h_k \|_{BV} + (|\mathcal{L}|_1 + |\rho|) K_1 \tau_k^\eta \| h_k \|_{BV}
\]

\[
\leq 2 \tau_k (1 + |\rho|) K_1 \tau_k^\eta K_2 |h_k|_1,
\]

where \(K_1, K_2, \text{and } \eta\) are as above. This gives the error bound \(|\rho_k - \rho| \leq O(\tau_k^\eta)\). 

\[\blacksquare\]
5.3. Proof of Lemma 2.9. Let $\mathcal{L}_m$ be the transfer operator associated with $T_m$. That is, $\mathcal{L}_m(f) = \hat{\mathcal{L}}(1_{X_m} f)$. Then, $\mathcal{L}_m(f) = \hat{\mathcal{L}}^n(1_{X_{m+n-1}} f)$, and therefore,

$$\hat{\mathcal{L}}^m \circ \mathcal{L}_m^n = \mathcal{L}_m^{m+n}.$$  

(5.8)

Hence, an interval is good for $T_0$ if and only if it is good for $T_m$ for every $m$. In the rest of this proof we will say that an interval is good if it is good for either (and therefore all) $T_m$.

Let $Z_0 = Z \vee H_0$, where $H_0$ is the partition of $H_0$ into intervals, and we recall that $Z$ is the monotonicity partition of $T$. Let $G_\epsilon$ be an $\epsilon$-adequate partition for $T_0$. Then, a partition $G_{\epsilon,m}$ may be constructed by cutting each element of $G_\epsilon \vee Z_0^{(m)}$ into at most $K$ pieces, where $K$ is independent of $m$, in such a way that the variation requirement $\max_{Z \in G_{\epsilon,m}} \var_Z(\hat{g}1_{X_m}) \leq \|DT_m^{-1}\|_\infty (1 + \epsilon)$ is satisfied, and thus $G_{\epsilon,m}$ is an $\epsilon$-adequate partition for $T_m$. Indeed, $K = 2 + \|\hat{g}\|_\infty / \essinf(\hat{g})$ is a possible choice; the term 2 allows one to account for possible jumps at the boundary points of $H_m$, as there are at most two of them in each $Z \in G_\epsilon \vee Z_0^{(m)}$. The term $M = \|\hat{g}\|_\infty / \essinf(\hat{g})$ allows one to split each interval $Z \in G_\epsilon \vee Z_0^{(m)}$ into at most $M$ subintervals $Z_1, \ldots, Z_M$, in such a way that, for every $1 \leq j \leq M$, $\var_{\text{int}(Z_j)}(\hat{g}1_{X_m}) \leq (1 + \epsilon)\|\hat{g}1_{X_m}\|_\infty$. The chosen value of $M$ is necessary to account for the possible discrepancy between $\|\hat{g}1_{X_0}\|_\infty$ and $\|\hat{g}1_{X_m}\|_\infty$. (Recall also that $\hat{g}$ is continuous on each $\text{int}(Z_j)$.)

Now, let $b = \#Z_0$. Then, each bad interval of $G_\epsilon$ gives rise to at most $Kb^m$ (necessarily bad) intervals in $G_{\epsilon,m}$. When a good interval of $G_\epsilon$ is split, it also gives rise to at most $Kb^m$ intervals in $G_{\epsilon,m}$. In this case some of the intervals may be bad, but it is guaranteed that at least one of them remains good, as being good is equivalent to having nonzero $\mu$ measure. Thus, the number of contiguous bad intervals in $G_{\epsilon,m}$ is at most $Kb^m(B + 2)$, where $B$ is the number of contiguous bad intervals in $G_\epsilon$. Therefore, $\hat{\xi}_e(T_m) = \exp \left( \limsup_{n \to \infty} \frac{1}{n} \log (1 + \xi_{e,n}(T_m)) \right) \leq \hat{\xi}_e(T_0)$.

Clearly, $\hat{\Theta}(T_m) \leq \hat{\Theta}(T_0)$. Finally, we will show that $\rho(T_0) \leq \rho(T_m)$. Recall that $\rho_j$ is the leading eigenvalue of $\mathcal{L}_j$. Let $f \in BV$ be nonzero and such that $\mathcal{L}_0 f = \rho_0 f$. We claim that $\mathcal{L}_m(1_{X_{m-1}} f) = \rho_0 1_{X_{m-1}} f$, which yields the inequality, because necessarily $1_{X_{m-1}} f \neq 0$. Indeed,

$$\rho_0 1_{X_{m-1}} f = 1_{X_{m-1}} \mathcal{L}_0 f = 1_{X_{m-1}} \mathcal{L}_m f = \mathcal{L}_m(1_{X_{m-1}} f),$$

where the second equality follows from the fact that $\mathcal{L}_0(1_{H_m})$ is supported on $T(H_m) = H_{m-1}$; the third one, from the fact that $\mathcal{L}_m f$ is supported on $T(X_m) \subseteq X_{m-1}$; and the last one, because $\mathcal{L}_m(1_{H_{m-1}} f) = 0$.

The first statement of the lemma follows. The relations between escape rates, accims, and quasi-conformal measures follow from comparing via (5.8) the statements of part (4) of Theorem 2.7 applied to $T_0$ and $T_m$. \hfill\blacksquare

5.4. Proof of Lemma 3.3. Assume that $H_0$ is an $\epsilon$-admissible hole for $\hat{T}$. Then, $T^n := (\hat{T}^n, H_{n-1})$ is an open Lasota–Yorke map. Fix $\Theta < \eta < \rho$ so that for all $n$ sufficiently large,

$$\exp \left( \frac{1}{n} \log \|(DT^n)^{-1}\|_\infty \right) \exp \left( \frac{1}{n} \log (1 + \xi_{e,n}) \right) < \eta.$$
Then, \( \| (D \tau)^{n-1} \|_\infty \xi_{x,n} < \eta^n \). By possibly making \( n \) larger, we can assume that \( (2 + \epsilon) \| (D \tau)^{n-1} \|_\infty < \eta^n \), and that \( 2\eta^n < \rho^n \). Then, \( \| (D \tau)^{n-1} \|_\infty (2 + \epsilon + \xi_{x,n}) < \rho^n (n) \).

We remark that \( \xi(x) = \xi_{x,n}(\tau) \). Thus \( \alpha(x) = \| (D \tau)^{n-1} \|_\infty (2 + \epsilon + \xi_{x,n}) \). Furthermore, in view of Theorem 2.7, \( \rho(x) = \lim_{n \to \infty} \inf_{x \in D_{\text{ann}}} \frac{L^{n+1}(x)}{L^{2n+1}(x)} = \mu(L^{n+1}) = \rho^n \). \( \blacksquare \)

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