# Nonuniformly Expanding 1D Maps 

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#### Abstract

This paper attmpts to make accessible a body of ideas surrounding the following result: Typical families of (possibly multimodal) 1-dimensional maps passing through "Misiurewicz points" have invariant densities for positive measure sets of parameters.


The simplest paradigms of chaotic behavior in dynamical systems are found in uniformly expanding and uniformly hyperbolic (or Anosov) maps. Allowing expanding and contracting behaviors to mix leads to a multitude of new possibilities. In spite of much progress, the analysis of most nonuniformly hyperbolic systems has remained hopelessly difficult. Onedimensional maps are an exception. The situation in 1 dimension is made tractable by the fact that the worst enemy of expansion is the critical set, i.e., the set on which $f^{\prime}$ vanishes, and for typical 1D maps, this set is finite. It has been shown that by controlling the orbits starting from this finite set, the dynamics on the rest of the phase space can be tamed.

A nonuniform theory for 1D maps was developed in a series of papers in the late 1970s and 1980s ([J], [M], [CE], [BC1], [R], [BC2], [NS], and others). These ideas are later exploited in the study of attractors with a single direction of instability, beginning with the Hénon maps ([BC2], [BY] etc.) and culminating recently in a general theory of rank-one attractors that can live in phase spaces of arbitrary dimensions [WY2]. In the course of these developments, some of the original 1D arguments have been extended and improved. This paper is written in response to numerous requests from the dynamical systems community to make more accessible a certain body of ideas in 1 dimension, both for its independent interest and as an introduction to the study of rank-one maps in higher dimensions.

The content of this paper can be summarized as follows: Let $I$ be a closed interval or a circle, and let $C^{2}(I, I)$ denote the set of $C^{2}$ maps from $I$ to itself. We seek to identify a reasonably large class of maps $\mathcal{G} \subset C^{2}(I, I)$ with controlled nonuniform expansion, and to give a description of its dynamical properties. This is carried out in the following 3 steps:
(1) First we identify a set $\mathcal{M} \subset C^{2}(I, I)$ defined by strong expanding conditions.
(2) Our class of "good maps" $\mathcal{G} \subset C^{2}(I, I)$ is obtained by relaxing these conditions. Maps in $\mathcal{G}$ are shown to have absolutely continuous invariant measures.
(3) We show that $\mathcal{G}$ is "large" in the sense that for every typical 1-parameter family $\left\{f_{a}\right\}$ passing through $\mathcal{M}$, the set $\left\{a: f_{a} \in \mathcal{G}\right\}$ has positive Lebesgue measure.

We first cite the main references directly related to (1)-(3): The class $\mathcal{M}$ is a slight generalization of the maps studied in $[\mathrm{M}]$. In the special case of the quadratic family

[^0]$f_{a}(x)=1-a x^{2}$, the existence of absolutely continuous invariant measures for a positive measure set of parameters is the well known theorem of Jakobson [J]; for other proofs of Jakobson's theorem, see [BC1], $[\mathrm{BC} 2]$ and $[\mathrm{R}]$. A key idea used in $[\mathrm{BC} 1]$ and $[\mathrm{BC} 2]$, namely the exponential growth of derivatives along critical orbits, is first introduced in [CE]. An analysis along the lines of (1)-(3) above for unimodal maps was carried out in [TTY]. A version of Jakobson's theorem for multimodal maps is given in [T].

While the results of this paper as stated have not been published before, we do not claim that the ideas of the proofs are new. In outline, our proofs follow those in [BC1] and Section 2 of $[\mathrm{BC} 2]$. The generalization from $f_{a}(x)=1-a x^{2}$ to more general maps is along the lines of [TTY]. We have also borrowed heavily from [WY1] and especially [WY2], both in terms of setting and the way in which the arguments are carried out. More detailed references are given at the end of each section.

Organization of paper The class $\mathcal{M}$ in (1) above is discussed in Section 1; the class $\mathcal{G}$ is introduced in Section 2. The result on invariant measures (Theorem 1) is stated and proved in Section 3. The result on positive measure sets of parameters (Theorem 2) is stated in Sect. 4.1 and proved in Sections 4-7.

## Part I DYNAMICAL PROPERTIES

## 1 The Class $\mathcal{M}$

### 1.1 Definition and expanding property

For $f \in C^{2}(I, I)$, let $C=C(f)=\left\{f^{\prime}=0\right\}$ denote the critical set of $f$, and let $C_{\delta}$ denote the $\delta$-neighborhood of $C$ in $I$. For $x \in I$, let $d(x, C):=\min _{\hat{x} \in C}|x-\hat{x}|$.

Definition 1.1 We say $f \in C^{2}(I, I)$ is in the class $\mathcal{M}$ if the following hold for some $\delta_{0}>0$ : (a) Outside of $C_{\delta_{0}}$ : there exist $\lambda_{0}>0, M_{0} \in \mathbb{Z}^{+}$and $0<c_{0} \leq 1$ such that
(i) for all $n \geq M_{0}$, if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta_{0}}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq e^{\lambda_{0} n}$;
(ii) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta_{0}}$ and $f^{n}(x) \in C_{\delta_{0}}$, any $n$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{0} e^{\lambda_{0} n}$.
(b) Inside $C_{\delta_{0}}$ : (i) $f^{\prime \prime}(x) \neq 0$ for all $x \in C_{\delta_{0}}$;
(ii) for all $\hat{x} \in C$ and $n>0, d\left(f^{n}(\hat{x}), C\right) \geq \delta_{0}$;
(iii) for all $x \in C_{\delta_{0}} \backslash C$, there exists $p_{0}(x)>0$ such that $f^{j}(x) \notin C_{\delta_{0}}$ for all $j<p_{0}(x)$ and $\left|\left(f^{p_{0}(x)}\right)^{\prime}(x)\right| \geq c_{0}^{-1} e^{\frac{1}{3} \lambda_{0} p_{0}(x)}$.

Remark 1. The maps in $\mathcal{M}$ are among the simplest with nonuniform expansion: The phase space is divided into two regions, $C_{\delta_{0}}$ and $I \backslash C_{\delta_{0}}$. Condition (a) in Definition 1.1 says that on $I \backslash C_{\delta_{0}}, f$ is essentially uniformly expanding. (b)(iii) says that for $x \in C_{\delta_{0}} \backslash C$, even though $\left|f^{\prime}(x)\right|$ is small, the orbit of $x$ does not return to $C_{\delta_{0}}$ again until its derivative has regained a definite amount of exponential growth; i.e., if $n$ is the first return time of $x \in C_{\delta_{0}}$ to $C_{\delta_{0}}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq e^{\frac{1}{3} \lambda_{0} n}$. (To see this, use (b)(iii) followed by (a)(ii)).

Remark 2. We identify two properties of the critical orbits of $f \in \mathcal{M}$ that will serve as the basis of the generalization in Section 2. Let $\hat{x} \in C$.
(1) $d\left(f^{n}(\hat{x}), C\right) \geq \delta_{0}$ for all $n>0$, i.e., (b)(ii) in Definition 1.1. (This condition is redundant and is included solely for emphasis; it follows from (b)(iii) together with the observation that $p_{0}(x) \rightarrow \infty$ as $d(x, C) \rightarrow 0$.)
(2) $\left|\left(f^{n}\right)^{\prime}(f \hat{x})\right| \geq c_{0}^{\prime} e^{\lambda_{0} n}$ for all $n>0$ where $c_{0}^{\prime}=\left(\max \left|f^{\prime}\right|\right)^{-M_{0}}$. This follows from (b)(ii) and (a)(i).

We record for future use the following important fact about the behavior of $f \in \mathcal{M}$ outside of $C_{\delta}$ for arbitrary $\delta<\delta_{0}$ :

Lemma 1.1 There exists $c_{0}^{\prime \prime}>0$ depending only on $f$ such that for all $\delta<\delta_{0}$ and $n>0$ :
(a) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{0}^{\prime \prime} \delta e^{\frac{1}{3} \lambda_{0} n}$;
(b) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$ and $f^{n}(x) \in C_{\delta_{0}}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{0} e^{\frac{1}{3} \lambda_{0} n}$.

Proof: Let $x$ be such that $f^{i}(x) \notin C_{\delta}$ for $i \in[0, n)$. We divide $[0, n]$ into maximal time intervals $[i, i+k]$ such that $f^{i+j}(x) \notin C_{\delta_{0}}$ for $0<j<k$, and estimate $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|$ as follows:

Case 1. $f^{i}(x), f^{i+k}(x) \in C_{\delta_{0}} .\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right| \geq e^{\frac{1}{3} \lambda_{0} k}$ by Definition 1.1(a)(ii) and (b)(iii). Case 2. $f^{i}(x) \notin C_{\delta_{0}}, f^{i+k}(x) \in C_{\delta_{0}}$. The estimate is given by Definition 1.1(a)(ii).
Case 3. $f^{i}(x), f^{i+k}(x) \notin C_{\delta_{0}}$. If $k \geq M_{0}$, then $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|>e^{\lambda_{0} k}$ by Definition 1.1(a)(i). If $k<M_{0}$, we let $\hat{k}$ be the smallest integer $>k$ such that $f^{i+\hat{k}}(x) \in C_{\delta_{0}}$. Using Definition 1.1(a)(i) for $\hat{k} \geq M_{0}$ and Definition 1.1(a)(ii) for $\hat{k}<M_{0}$, we conclude that $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|>c_{0}\left(\max \left|f^{\prime}(x)\right|\right)^{-M_{0}} e^{\lambda_{0} k}$.

Case 4. $f^{i}(x) \in C_{\delta_{0}}, f^{i+k}(x) \notin C_{\delta_{0}}$. As in Case 3, with extra factor $\left(\min _{y \in C \delta_{0}}\left|f^{\prime \prime}(y)\right|\right) \delta$. Cases 3 and 4 are relevant only for part (a); each appears at most once in the estimate on $\left|\left(f^{n}\right)^{\prime}(x)\right|$.

In the interest of carrying as few constants around as possible, we write $c_{1}=\min \left\{c_{0}, c_{0}^{\prime}, c_{0}^{\prime \prime}\right\}$.

### 1.2 Examples

Example 1. Let $f \in C^{3}(I, I)$ be such that
(i) $S(f)<0$ where $S(f)$ denotes the Schwarzian derivative of $f,{ }^{3}$
(ii) $f^{\prime \prime}(\hat{x}) \neq 0$ for all $\hat{x} \in C$,
(iii) if $f^{n}(x)=x$, then $\left|\left(f^{n}\right)^{\prime}(x)\right|>1$, and
(iv) for all $\hat{x} \in C, \inf _{n>0} d\left(f^{n}(\hat{x}), C\right)>0$.

Then $f \in \mathcal{M}$. For a proof of this fact, see Lemma 2.5 of [WY1].
We note that (i) and (ii) above are satisfied by all members of the quadratic family $f_{a}(x)=1-a x^{2}, a \in(0,2$ ], and (iii) and (iv) are satisfied by an uncountable number of $a$ including $a=2$.

Example 2. Another situation where maps in $\mathcal{M}$ arise naturally is through scaling. The following is a slight generalization of Lemma 5.3 in [WY3] and has the same proof: Let $f_{a}: S^{1} \rightarrow S^{1}$ be given by

$$
f_{a}(\theta)=\theta+a+L \Phi(\theta)
$$

[^1]where $a, L \in \mathbb{R}$ and $\Phi: S^{1} \rightarrow S^{1}$ is an arbitrary function with nondegenerate critical points (and the right side is to be read mod 1). Then there exists $L_{0}>0$ such that for all $L \geq L_{0}$, there exists an $\mathcal{O}\left(\frac{1}{L}\right)$-dense set of $a$ for which $f_{a} \in \mathcal{M}$.

Reference: Maps of the type in Example 1 are introduced and studied in [M]. Maps of the type in Example 2 appear naturally in [WY3] and [WY4].

## 2 The Class $\mathcal{G}: 3$ Basic Properties

Condition (b)(ii) in Definition 1.1 severely limits the scope of $\mathcal{M}$ as a subset of $C^{2}(I, I)$. We now introduce in a neighborhood of each $f_{0} \in \mathcal{M}$ an admissible set of perturbations $\mathcal{G}\left(f_{0}\right)$. Our set of "good maps" $\mathcal{G}$ is then defined to be $\bigcup_{f_{0} \in \mathcal{M}} \mathcal{G}\left(f_{0}\right)$.

Throughout this section, let $f_{0} \in \mathcal{M}$ be fixed, and let $\delta_{0}, \lambda_{0}, c_{1}$ etc. be the constants in Sect. 1.1 associated with $f_{0}$.

### 2.1 Definition of $\mathcal{G}\left(f_{0}\right)$ and basic properties

For $\lambda, \alpha, \varepsilon>0$ and $f \in C^{2}(I, I)$, we say $f \in \mathcal{G}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ if $\left\|f-f_{0}\right\|_{C^{2}}<\varepsilon$ and the following hold for all $\hat{x} \in C=C(f)$ and $n>0$ :
(G1) $d\left(f^{n}(\hat{x}), C\right)>\min \left\{\frac{1}{2} \delta_{0}, e^{-\alpha n}\right\}$;
(G2) $\left|\left(f^{n}\right)^{\prime}(f(\hat{x}))\right| \geq c_{1} e^{\lambda n}$.
Note that with $\lambda<\lambda_{0}$, (G1) and (G2) are relaxations of the conditions on critical orbits for $f_{0}$ (see Remark 2 in Sect. 1.1). The main result of this section is

Proposition 2.1 Given $f_{0} \in \mathcal{M}, \lambda<\frac{1}{4} \lambda_{0}$ and $\alpha<\frac{1}{100} \lambda$, there exists $\delta=\delta\left(f_{0}, \lambda, \alpha\right)$ and $\varepsilon=\varepsilon\left(f_{0}, \lambda, \alpha, \delta\right)>0$ such that (P1)-(P3) below hold for all $f \in \mathcal{G}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$.

Here $\delta<\frac{1}{2} \delta_{0}$ is an auxiliary constant. For simplicity, we assume $\varepsilon$ is small enough that $d\left(f^{n}(\hat{x}), C\right)>\frac{1}{2} \delta_{0}$ for all $\hat{x} \in C$ and $1 \leq n \leq n_{0}$ where $n_{0}$ is a large integer satisfying $e^{-\alpha n_{0}} \ll \delta$. Consequently, (G1) can be violated only when $f^{n}(\hat{x}) \in C_{\delta}$. Precise requirements on $\delta$ and $\varepsilon$ will become clear in the proofs. In general, $\varepsilon \ll \delta \ll 1$. The arguments are perturbative; some of them will require that $\varepsilon$ be taken very close to 0 .

The set $\mathcal{G}\left(f_{0}\right)$ is defined to be the union of $\mathcal{G}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ as ( $\left.\lambda, \alpha, \varepsilon\right)$ ranges over all triples satisfying the conditions in Proposition 2.1.

We now state (P1)-(P3), introducing some useful language along the way.
(P1) Outside of $C_{\delta}$ : (i) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{1} \delta e^{\frac{1}{4} \lambda_{0} n}$;
(ii) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$ and $f^{n}(x) \in C_{\delta_{0}}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{1} e^{\frac{1}{4} \lambda_{0} n}$.

Let $\hat{x} \in C$, and let $C_{\delta}(\hat{x}):=(\hat{x}-\delta, \hat{x}+\delta)$. For $x \in C_{\delta}(\hat{x}) \backslash\{\hat{x}\}$, we define $p(x)$, the bound period of $x$, to be the largest integer such that $\left|f^{i}(x)-f^{i}(\hat{x})\right| \leq e^{-2 \alpha i} \forall i<p(x)$.
(P2) Partial derivative recovery for $x \in C_{\delta} \backslash C$ : For $x \in C_{\delta}(\hat{x}) \backslash\{\hat{x}\}$,
(i) $\frac{1}{3 \ln \left(\max \left|f^{\prime}\right|\right)} \log \frac{1}{|x-\hat{x}|} \leq p(x) \leq \frac{3}{\lambda} \log \frac{1}{|x-\hat{x}|}$;
(ii) $\left|\left(f^{p(x)}\right)^{\prime}(x)\right|>e^{\frac{\lambda}{3} p(x)}$.
(P2) leads to the following general description of orbits:
Decomposition into "bound" and "free" states: For $x \in I$ such that $f^{i}(x) \notin C$ for all $i \geq 0$ (for example, $x=f(\hat{x})$ for $\hat{x} \in C$ ), let

$$
t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq \cdots
$$

be defined as follows: $t_{1}$ is the smallest $j \geq 0$ such that $f^{j}(x) \in C_{\delta}$. For $k \geq 1$, let $p_{k}$ be the bound period of $f^{t_{k}}(x)$, and let $t_{k+1}$ be the smallest $j \geq t_{k}+p_{k}$ such that $f^{j}(x) \in C_{\delta}$. (Note that an orbit may return to $C_{\delta}$ during its bound periods, i.e. $t_{i}$ are not the only return times to $C_{\delta}$.) This decomposes the orbit of $x$ into segments corresponding to time intervals $\left(t_{k}, t_{k}+p_{k}\right)$ and $\left[t_{k}+p_{k}, t_{k+1}\right]$, during which we describe the orbit of $x$ as being in "bound" and "free" states respectively; $t_{k}$ are called times of free returns.
(P3) is about comparisons of derivatives for nearby orbits. To state what it means for two points to be close to each other, we introduce a partition $\mathcal{P}$ on $I$. First let $\mathcal{P}_{0}=\left\{I_{\mu j}\right\}$ be the following partition on $(-\delta, \delta)$ : Assume $\delta=e^{-\mu_{*}}$ for some $\mu_{*} \in \mathbb{Z}^{+}$. For $\mu \geq \mu_{*}$, let $I_{\mu}=\left(e^{-(\mu+1)}, e^{-\mu}\right)$; for $\mu \leq-\mu_{*}$, let $I_{\mu}$ be the reflection of $I_{-\mu}$ about 0 . Each $I_{\mu}$ is further subdivided into $\frac{1}{\mu^{2}}$ subintervals of equal length called $I_{\mu j}$. For $\hat{x} \in C$, let $\mathcal{P}_{0}^{\hat{x}}$ be the partition on $C_{\delta}(\hat{x})$ obtained by shifting the center of $\mathcal{P}_{0}$ from 0 to $\hat{x}$. The partition $\mathcal{P}$ is defined to be $\mathcal{P}_{0}^{\hat{x}}$ on $C_{\delta}(\hat{x})$; on $I \backslash C_{\delta}$, its elements are intervals of length $\approx \delta$.

The following shorthand is used: We refer to $\pi \in \mathcal{P}$ corresponding to (translated) $I_{\mu j}$ intervals in $\mathcal{P}_{0}^{\hat{x}}$ simply as " $I_{\mu j}$ ". For $\pi \in \mathcal{P}, \pi^{+}$denotes the union of $\pi$ and the two elements of $\mathcal{P}$ adjacent to it. For an interval $\gamma \subset I$, we say $\gamma \approx \pi$ if $\pi \subset \gamma \subset \pi^{+}$. For practical purposes, $\pi^{+}$intersecting $\partial C_{\delta}$ can be treated as "inside $C_{\delta}$ " or "outside $C_{\delta}$ ". ${ }^{4}$ For $\gamma \subset I_{\mu j}^{+}$, we define the bound period of $\gamma$ to be $p(\gamma)=\min _{x \in I_{\mu j}^{+}}\{p(x)\}$. For $x, y \in I,[x, y]$ denotes the segment connecting $x$ and $y$. We say $x$ and $y$ in $I$ have the same itinerary (with respect to $\mathcal{P}$ ) through time $n-1$ if there exist $t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq \cdots \leq n$ such that for every $k, f^{t_{k}}[x, y] \subset \pi^{+}$for some $\pi \subset C_{\delta}, p_{k}=p\left(f^{t_{k}}[x, y]\right)$, and for all $i \in[0, n) \backslash \cup_{k}\left[t_{k}, t_{k}+p_{k}\right)$, $f^{i}[x, y] \subset \pi^{+}$for some $\pi \in \mathcal{P}$ with $\pi \cap C_{\delta}=\emptyset$.
(P3) Distortion estimate: There exists $K_{0}>1$ (depending only on $f_{0}$ and on $\lambda$ ) such that if $x$ and $y$ have the same itinerary through time $n-1$, then

$$
\left|\frac{\left(f^{n}\right)^{\prime}(x)}{\left(f^{n}\right)^{\prime}(y)}\right| \leq K_{0} .
$$

(P1)-(P3) are proved in next subsection. We finish by recording the following corollary of Proposition 2.1.

Corollary 2.1 There exists $K_{1}$ (depending only on $f_{0}$ and on $\lambda$ ) such that for all $x \in I$ with $f^{i}(x) \notin C$ for all $0 \leq i<n$,

$$
\left|\left(f^{n}\right)^{\prime}(x)\right|>K_{1}^{-1} d\left(f^{j}(x), C\right) e^{\frac{1}{4} \lambda n}
$$

[^2]where $j$ is the time of the last free return before $n$. The factor $d\left(f^{j}(x), C\right)$ may be replaced by $\delta$ if $f^{n}(x)$ is free.

Proof: Let $0 \leq t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq \cdots$ be as in the paragraph following (P2). The derivatives on time intervals $\left[t_{k}, t_{k}+p_{k}\right.$ ) and $\left[t_{k}+p_{k}, t_{k+1}\right)$ are given by (P2)(ii) and ( P 1 )(ii) respectively, provided these intervals are completed before time $n$. We assume $\delta$ is sufficiently small so that the constant $c_{1}$ in (P1)(ii) is absorbed into the exponential estimate from the proceeding bound period. If $f^{n}(x)$ is in a bound period initiated at time $j$, then $\left|\left(f^{(n-j)}\right)^{\prime}\left(f^{j}(x)\right)\right| \geq\left|f^{\prime}\left(f^{j}(x)\right)\right| K_{0}^{-1} c_{1} e^{\lambda(n-j-1)}$; see (G2) and (P3). If $t_{k}+p_{k} \leq n<t_{k+1}$ for some $k$, then the derivative between time $t_{k}+p_{k}$ and $n$ is given by ( P 1 )(i).

Remarks on the use of constants: In this article, $K, K_{1}, K_{2}, \cdots$, are reserved for use as system constants, which in Part I are constants that are allowed to depend only on (1) $f_{0}$, by which we included also the constants in Sect. 1.1 associated with $f_{0}$, and (2) our choice of $\lambda$. The more important of these constants, such as $K_{0}$ in (P3), carry a subscript; all others are referred to by the generic name $K$. The value of $K$, therefore, may vary from expression to expression.

Notation: Where no ambiguity arises, i.e. when only one map $f$ is involved, we will sometimes write $x_{i}=f^{i}(x)$ for $i=1,2, \cdots$.

### 2.2 Proofs of (P1)-(P3)

Proof of (P1): First we deduce from Lemma 1.1(a) that there exists $N=N(\delta)$ such that for all $y \in I$, if $y, f_{0}(y), \cdots, f_{0}^{N-1}(y) \notin C_{\frac{1}{2} \delta}\left(f_{0}\right)$, then $\left|\left(f_{0}^{N}\right)^{\prime}(y)\right|>e^{\frac{1}{3.5} \lambda_{0} N}$. We then choose $\varepsilon$ small enough that $f$ is sufficiently close to $f_{0}$ for $N$ iterates in the sense below:
(i) if $x$ and $n$ are as in (P1) and $n \leq N$, then $\left|\left(f^{n}\right)^{\prime}(x)-\left(f_{0}^{n}\right)^{\prime}(x)\right|$ is small enough that the conclusions of ( P 1 ) follow from Lemma 1.1;
(ii) if $f^{i}(y) \notin C_{\delta}$ for $0 \leq i<N$, then $\left|\left(f^{N}\right)^{\prime}(y)\right|>e^{\frac{1}{4} \lambda_{0} N}$.

If $n$ in ( P 1 ) is $>N$, we let $k$ be such that $k N \leq n<(k+1) N$, and estimate $\left|\left(f^{n}\right)^{\prime}(x)\right|$ by the chain rule, comparing $\left(f^{N}\right)^{\prime}\left(f^{i N}(x)\right)$ with $\left(f_{0}^{N}\right)^{\prime}\left(f^{i N}(x)\right)$ for $i \leq k$ using (ii) above, and $\left(f^{n-k N}\right)^{\prime}\left(f^{k N}(x)\right)$ with $\left(f_{0}^{n-k N}\right)^{\prime}\left(f^{k N}(x)\right)$ using (i).

Lemma 2.1 The following holds if $\delta$ and $\varepsilon$ are sufficiently small and suitably related: Let $\hat{x} \in C$, and let $x \in C_{\delta}(\hat{x})$. Then for all $y \in[\hat{x}, x]$ and $k<p(x)$,

$$
\frac{1}{2} \leq \frac{\left(f^{k}\right)^{\prime}\left(y_{1}\right)}{\left(f^{k}\right)^{\prime}\left(\hat{x}_{1}\right)} \leq 2
$$

Proof: First,

$$
\log \frac{\left(f^{k}\right)^{\prime}\left(y_{1}\right)}{\left(f^{k}\right)^{\prime}\left(\hat{x}_{1}\right)} \leq \sum_{j=1}^{k} \frac{\left|f^{\prime}\left(y_{j}\right)-f^{\prime}\left(\hat{x}_{j}\right)\right|}{\left|f^{\prime}\left(\hat{x}_{j}\right)\right|} \leq K \sum_{j=1}^{k} \frac{\left|y_{j}-\hat{x}_{j}\right|}{d\left(\hat{x}_{j}, C\right)}
$$

We choose $h_{0}$ large enough that (i) $\sum_{i>h_{0}}^{\infty} e^{-\alpha j} \ll 1$ and (ii) $e^{-\alpha h_{0}}<\frac{1}{2} \delta_{0}$ (so that $d\left(\hat{x}_{j}, C\right)>e^{-\alpha j}$ for $\left.j>h_{0}\right)$. Next we choose $\delta$ small enough that $\delta \sum_{j=1}^{h_{0}} \frac{2}{\delta_{0}}\left(\max \left|f^{\prime}\right|\right)^{j} \ll 1$. Finally, let $\varepsilon$ be small enough that $d\left(\hat{x}_{j}, C\right)>\frac{1}{2} \delta_{0}$ for all $j \leq h_{0}$. Then

$$
\sum_{j=1}^{k} \frac{\left|y_{j}-\hat{x}_{j}\right|}{d\left(\hat{x}_{j}, C\right)}<\sum_{j=1}^{h_{0}} \frac{2}{\delta_{0}}\left(\max \left|f^{\prime}\right|\right)^{j} \delta+\sum_{j=h_{0}+1}^{k} \frac{e^{-2 \alpha j}}{e^{-\alpha j}} \ll 1 .
$$

Proof of (P2): Suppose $|x-\hat{x}|=e^{-h}$. Then (G2) together with Lemma 2.1 implies that

$$
\left|x_{p}-\hat{x}_{p}\right| \geq \frac{1}{2}\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|\left|x_{1}-\hat{x}_{1}\right| \geq K^{-1} e^{\lambda(p-1)}(x-\hat{x})^{2} .
$$

From this we deduce that $p<\frac{3}{\lambda} h$, assuming that $h$ is sufficiently large (or $\delta$ is sufficiently small). The lower bound on $p$ is obtained by comparing the inequalities $\left|x_{p}-\hat{x}_{p}\right|<$ $\left(\max \left|f^{\prime}\right|\right)^{p-1} e^{-2 h}$ and $\left|x_{p}-\hat{x}_{p}\right| \geq e^{-2 \alpha p}$ (definition of bound period).

To prove (P2)(ii), taking square root of $\left|x_{p}-\hat{x}_{p}\right|<K\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|(x-\hat{x})^{2}$, we obtain

$$
\begin{equation*}
K\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|^{\frac{1}{2}}|x-\hat{x}|>e^{-\alpha p} . \tag{1}
\end{equation*}
$$

Then writing $\left|\left(f^{p}\right)^{\prime}(x)\right|$ as

$$
\left|\left(f^{p}\right)^{\prime}(x)\right|=\left|\left(f^{p-1}\right)^{\prime}\left(x_{1}\right)\right|\left|f^{\prime}(x)\right|>\left(K^{-1}\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|^{\frac{1}{2}}|x-\hat{x}|\right) \cdot\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|^{\frac{1}{2}}
$$

and substituting in (1), we see that $\left|\left(f^{p}\right)^{\prime}(x)\right|>K^{-1} c_{1}^{\frac{1}{2}} e^{\frac{1}{2} \lambda(p-1)} e^{-\alpha p}$, which we may assume is $>e^{\frac{1}{3} \lambda p}$ if $p$ is sufficiently large, or equivalently, $\delta$ is sufficiently small.

Proof of (P3): We write $\sigma_{0}=[x, y], \sigma_{k}=f^{t_{k}} \sigma_{0}$, and assume for definiteness that $\sigma_{0} \subset C_{\delta}$ and $t_{q}+p_{q} \leq n$ where $f^{i} \sigma_{0} \cap C_{\delta}=\emptyset$ for all $t_{q}+p_{q} \leq i<n$. (The proof for the case $t_{q}<n<t_{q}+p_{q}$ is contained in that for $n=t_{q}+p_{q}$.) Then

$$
\log \frac{\left(f^{n}\right)^{\prime}(x)}{\left(f^{n}\right)^{\prime}(y)} \leq \sum_{j=0}^{n-1} \frac{\left|f^{\prime}\left(y_{j}\right)-f^{\prime}\left(x_{j}\right)\right|}{\left|f^{\prime}\left(y_{j}\right)\right|} \leq K \sum_{k=1}^{q}\left(S_{k}^{\prime}+S_{k}^{\prime \prime}\right)
$$

where

$$
S_{k}^{\prime}=\sum_{j=t_{k}}^{t_{k}+p_{k}-1} \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)} \quad \text { and } \quad S_{k}^{\prime \prime}=\sum_{t_{k}+p_{k}}^{t_{k+1}-1} \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)}
$$

except for $S_{q}^{\prime \prime}$, which ends at index $n-1$. Observe that for $k<q$, it follows from (P2)(ii) and (P1)(ii) that $\left|\sigma_{k+1}\right| \geq c_{1} e^{\frac{1}{3} \lambda\left(t_{k+1}-t_{k}\right)}\left|\sigma_{k}\right|$, which we may assume is $\geq \tau\left|\sigma_{k}\right|$ for some $\tau>1$ (the factor $c_{1}$ having been absorbed into the exponential assuming $\delta$ is sufficiently small).
I. Bound on $\sum_{k=1}^{q} S_{k}^{\prime}$

First we estimate $S_{k}^{\prime}$. Suppose $y_{t_{k}} \in C_{\delta}(\hat{x})$. For $t_{k}<j<t_{k}+p_{k}$ we write

$$
\frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)}=\frac{\left|y_{j}-x_{j}\right|}{\left|y_{j}-\hat{x}_{j-t_{k}}\right|} \cdot \frac{\left|y_{j}-\hat{x}_{j-t_{k}}\right|}{d\left(y_{j}, C\right)} .
$$

By Lemma 2.1 and the usual estimates near $\hat{x}$, the first factor on the right is

$$
<K \frac{\left|y_{t_{k}+1}-x_{t_{k}+1}\right|}{\left|y_{t_{k}+1}-\hat{x}_{1}\right|}<K \frac{\left|f^{\prime}\left(x_{t_{k}}\right)\right|\left|y_{t_{k}}-x_{t_{k}}\right|}{\left|y_{t_{k}}-\hat{x}\right|^{2}}<K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)} .
$$

Thus

$$
S_{k}^{\prime} \leq K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)}\left(1+\sum_{j=t_{k}+1}^{t_{k}+p_{k}-1} \frac{\left|y_{j}-\hat{x}_{j-t_{k}}\right|}{d\left(\hat{x}_{j-t_{k}}, C\right)}\right) \leq K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)},
$$

treating the sum inside the parentheses as in Lemma 2.1.
Now let $\mathcal{K}_{\mu}=\left\{k \leq q: \sigma_{k} \subset I_{\mu j}\right.$ for some $\left.j\right\}$. We claim that

$$
\sum_{k \in \mathcal{K}_{\mu}} S_{k}^{\prime}<K \sum_{k \in \mathcal{K}_{\mu}} \frac{\left|\sigma_{k}\right|}{e^{-|\mu|}}<K \frac{1}{\mu^{2}}
$$

The first inequality is from the estimate above. The second follows from (i) $\left|\sigma_{k+1}\right| \geq \tau\left|\sigma_{k}\right|$ and (ii) the term with the largest index is bounded above by $\left|I_{\mu j}^{+}\right|$, which is $\leq \frac{3}{\mu^{2}} e^{-|\mu|}$. To finish, we sum over all $\mu,|\mu| \geq \log \frac{1}{\delta}$, to obtain $\sum S_{k}^{\prime}<K$.

## II. Bound on $\sum_{k=1}^{q} S_{k}^{\prime \prime}$

For $k<q$ and $t_{k}+p_{k} \leq j \leq t_{k+1}-1$, we have, by (P1)(ii), $\left|\sigma_{k+1}\right| \geq c_{1} e^{\lambda\left(t_{k+1}-j\right)}\left|x_{j}-y_{j}\right|$, so $S_{k}^{\prime \prime} \leq K \frac{\left|\sigma_{k+1}\right|}{\delta}$. This together with $\left|\sigma_{k+1}\right| \geq \tau\left|\sigma_{k}\right|$ gives $\sum_{k=0}^{q-1} S_{k}^{\prime \prime} \leq K \frac{\left|\sigma_{q}\right|}{\delta}<K$. If $n<t_{q+1}, S_{q}^{\prime \prime}$ needs to be estimated separately: Assume the orbit in question visits $C_{\delta_{0}}$ during the time interval $\left[t_{q}+p_{q}, n\right.$ ) (otherwise the estimate is trivial). Let $\hat{n}$ be the time of the last such visit. Then

$$
S_{q}^{\prime \prime}=\sum_{j=t_{q}+p_{q}}^{\hat{n}-1} K \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)}+K \frac{\left|y_{\hat{n}}-x_{\hat{n}}\right|}{d\left(y_{\hat{n}}, C\right)}+\sum_{j=\hat{n}+1}^{n-1} K \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)} .
$$

The first sum is estimated using $\left|y_{\hat{n}-1}-x_{\hat{n}-1}\right|>c_{1} e^{\frac{1}{4} \lambda_{0}(\hat{n}-1-j)}\left|y_{j}-x_{j}\right|$ for $t_{q}+p_{q} \leq j<\hat{n}$ ((P1)(ii)). The last sum is estimated using $\left|y_{n-1}-x_{n-1}\right|>c_{1} \delta_{0} e^{\frac{1}{4} \lambda_{0}(n-j-1)}\left|y_{j}-x_{j}\right|$ for $\hat{n}<j<n-1\left((\mathrm{P} 1)(\mathrm{i})\right.$ with $\delta_{0}$ in the place of $\left.\delta\right)$. It follows that

$$
S_{q}^{\prime \prime} \leq K \frac{\left|y_{\hat{n}-1}-x_{\hat{n}-1}\right|}{\delta}+K \frac{\left|y_{\hat{n}}-x_{\hat{n}}\right|}{\delta}+K \frac{\left|y_{n-1}-x_{n-1}\right|}{\delta_{0}^{2}}<K .
$$

References: Conditions of the type (G1) are first used in [BC1] and [BC2]. (G2) is introduced in [CE]. A version of the material in this section for $f_{0}(x)=1-2 x^{2}$ first appeared in [BC1] and Section 2 of [BC2].

## 3 Absolutely Continuous Invariant Measures

The goal of this section is to prove
Theorem 1 Every $f \in \mathcal{G}$ has an absolutely continuous invariant probability measure.

### 3.1 Growth of segments with bounded distortion

Let $f \in \mathcal{G}$. To prove that $f$ admits an absolutely continuous invariant probability measure (acipm), we need to show that the forward images of Lebesgue measure have certain regularity properties. Expansion is conducive to such regularity, but distortion bounds are also essential. Since (P3) guarantees uniform distortion bounds only on intervals with the same itinerary, we need to show that intervals of this type grow with sufficient regularity. This is the mission of the present subsection.

We begin by introducing what will be referred to as a canonical subdivision by itinerary on an interval $\omega \subset I$. This consists of an increasing sequence of partitions $\mathcal{Q}_{0}<\mathcal{Q}_{1}<\mathcal{Q}_{2}<\cdots$ on $\omega$ defined as follows: Let us say an interval is "short" if it is strictly contained in an element of $\mathcal{P}$. $\mathcal{Q}_{0}$ is defined to be $\left.\mathcal{P}\right|_{\omega}$ except that the end intervals are attached to their neighbors if they are short. We assume inductively that all $\hat{\omega} \in \mathcal{Q}_{i}$ are intervals and all points in $\hat{\omega}$ have the same itinerary through time $i$. To go from $\mathcal{Q}_{i}$ to $\mathcal{Q}_{i+1}$, we consider one $\hat{\omega} \in \mathcal{Q}_{i}$ at a time:

- If $f^{i+1}(\hat{\omega})$ is in a bound period, then it is automatically put into $\mathcal{Q}_{i+1}$. (Observe that if $f^{i+1}(\hat{\omega}) \cap C_{\delta} \neq \emptyset$, then $f^{i+1}(\hat{\omega}) \subset I_{\mu^{\prime} j^{\prime}}^{+}$for some $\mu^{\prime}, j^{\prime}$, i.e. no cutting is needed during bound periods. This is an easy exercise.)
- If $f^{i+1}(\hat{\omega})$ is not in a bound period, but all points in $\hat{\omega}$ have the same itinerary through time $i+1$, we again put $\hat{\omega} \in \mathcal{Q}_{i+1}$.
- If neither of the last two cases holds, then we partition $\hat{\omega}$ into segments $\left\{\hat{\omega}^{\prime}\right\}$ according to their itineraries through time $i+1$, requiring that $f^{i+1}\left(\hat{\omega}^{\prime}\right) \approx \gamma$ for some $\gamma \in \mathcal{P}$ (i.e., no cuts are made that lead to short intervals). The resulting partition is $\left.\mathcal{Q}_{i+1}\right|_{\hat{\omega}}$.

For $x \in \omega$ and $i \geq 0$, let $\mathcal{Q}_{i}(x)$ denote the element of $\mathcal{Q}_{i}$ containing $x$. We introduce the following stopping time on $\omega$ : For $x \in \omega, S(x)$ is the smallest $i>0$ such that $f^{i}\left(\mathcal{Q}_{i-1}(x)\right)$ is not in a bound period and has length $>\delta$. The main result of this subsection is

Proposition 3.1 There exists $K_{2}>0$ such that for any $\omega \subset I$ with $\delta<|\omega|<3 \delta$,

$$
|\{S>n\}|<e^{-K_{2}^{-1} n}|\omega| \quad \text { for } n>K_{2} \log \delta^{-1}
$$

Here $|\cdot|$ denotes the Lebesgue measure of the set. The proof of this proposition follows from a series of lemmas.

Lemma 3.1 Let $\omega \approx I_{\mu j}$. Then $\left|f^{p}(\omega)\right|>e^{-\frac{7 \alpha}{\lambda}|\mu|}$ where $p=p(\omega)$.
Proof: It follows from Lemma 2.1 that for $I_{\mu j} \in \mathcal{P}^{\hat{x}}$,

$$
\begin{aligned}
\left|f^{p}\left(I_{\mu j}\right)\right| & =\frac{\left|f^{p}\left(I_{\mu j}\right)\right|}{\left|f^{p}\left(\left[\hat{x}, \hat{x} \pm e^{-|\mu|}\right]\right)\right|}\left|f^{p}\left(\left[\hat{x}, \hat{x} \pm e^{-|\mu|}\right]\right)\right| \\
& \geq K^{-1} \frac{\left|f\left(I_{\mu j}\right)\right|}{\left|f\left(\left[\hat{x}, \hat{x} \pm e^{-|\mu|}\right]\right)\right|}\left|f^{p}\left(\left[\hat{x}, \hat{x} \pm e^{-|\mu|}\right]\right)\right| \\
& \geq K^{-1} \frac{1}{\mu^{2}} e^{-2 \alpha p} .
\end{aligned}
$$

By (P2)(i), $p<\frac{3}{\lambda}|\mu|$. The lemma then follows assuming $\delta$ is sufficiently small.
For $\omega \approx I_{\mu j}$, we define the extended bound period of $\omega$ to be the largest $n$ such that all points in $\omega$ have the same itinerary (in the sense of (P3)) for $n-1$ iterates. The next lemma follows immediately from Corollary 2.1.

Lemma 3.2 There exists $K$ such that the extended bound period of $I_{\mu j}$ is $<K|\mu|$.
Lemma 3.3 Assume $\delta$ is sufficiently small. Then for $\omega \approx I_{\mu_{0} j_{0}}$,

$$
|\{x \in \omega: S(x)>n\}|<e^{-\frac{1}{2} K^{-1} n}|\omega| \quad \text { for all } n>K \mu_{0}
$$

where $K$ is the constant in Lemma 3.2.
Proof: Let $x \in \omega$ be such that $S(x)>n$. We define the essential return times and addresses of $x$ before $n$ as follows: Let $t$ be the extended bound period of $\omega$. Then either $S=t$, or $f^{t}(\omega) \subset C_{\delta}$ (more precisely $\cup I_{\mu j}^{+}$). In the latter case, we say $t_{1}=t$ is the first essential return time of $x$, and if $f^{t_{1}}\left(\mathcal{Q}_{t_{1}}(x)\right) \approx I_{\mu_{1} j_{1}} \subset C_{\delta}$, then we say $I_{\mu_{1} j_{1}}$ is its first essential return address. If $S$ has not been reached, we continue iterating. Let $t$ be the extended bound period of $f^{t_{1}}\left(\mathcal{Q}_{t_{1}}(x)\right)$. Then either $S(x)=t_{1}+t$, or we define the second essential return time to be $t_{2}=t_{1}+t$ and second return address to be $I_{\mu_{2} j_{2}}$ if $f^{t_{2}}\left(\mathcal{Q}_{t_{2}}(x)\right) \approx I_{\mu_{2} j_{2}} \subset C_{\delta}$, and so on.

Let $A_{q}=\left\{x \in \omega: S(x)>n, f^{i}(x)\right.$ makes $q$ but not $q+1$ essential returns before time $\left.n\right\}$. Then $|\{S>n\}|=\sum_{q}\left|A_{q}\right|$. We write $A_{q}=\cup_{R} A_{q, R}$ where $A_{q, R}=\left\{x \in A_{q}:\right.$ if $\left(\mu_{1}, \cdots, \mu_{q}\right)$ are the $\mu$-coordinates of its first $q$ return addresses, then $\left.\left|\mu_{1}\right|+\left|\mu_{2}\right|+\cdots+\left|\mu_{q}\right|=R\right\}$, and further decompose $A_{q, R}$ into intervals $\sigma$ consisting of points whose first $q$ return addresses are identical. Each such $\sigma$ is equal to $\mathcal{Q}_{t_{q}}(x)$ for some $x$, since the extended bound period of $f^{t_{q}}\left(\mathcal{Q}_{t_{q}}(x)\right)$ is not completed before time $n$ (see above). Writing $Q_{t_{k}}=\mathcal{Q}_{t_{k}}(x)$, we have

$$
\begin{aligned}
|\sigma| & =\frac{\left|Q_{t_{q}}\right|}{\left|Q_{t_{q-1}}\right|} \frac{\left|Q_{t_{q-1}}\right|}{\left|Q_{t_{q-2}}\right|} \cdots \frac{\left|Q_{t_{1}}\right|}{|\omega|}|\omega| \\
& \leq K_{0}^{q} \frac{\left|f^{t_{q-1}+p_{q-1}}\left(Q_{t_{q}}\right)\right|}{\left|f^{t_{q-1}+p_{q-1}}\left(Q_{t_{q-1}}\right)\right|} \frac{\left|f^{t_{q-2}+p_{q-2}}\left(Q_{t_{q-1}}\right)\right|}{\left.\left|f^{t_{q-2}+p_{q-2}\left(Q_{t_{q-2}}\right) \mid} \cdots \frac{\left|f^{p(\omega)}\left(Q_{t_{1}}\right)\right|}{\left|f^{p(\omega)}(\omega)\right|}\right| \omega \right\rvert\,} \\
& \leq c_{1}^{-q} K_{0}^{q} \frac{\left|f^{t_{q}}\left(Q_{t_{q}}\right)\right|}{\left\lvert\, f^{t_{q-1}+p_{q-1}\left(Q_{t_{q-1}}\right) \mid} \frac{\left|f^{t_{q-1}}\left(Q_{t_{q-1}}\right)\right|}{\left.\left|f^{t_{q-2}+p_{q-2}\left(Q_{t_{q-2}}\right) \mid} \cdots \frac{\left|f^{t_{1}}\left(Q_{t_{1}}\right)\right|}{\left|f^{p(\omega)}(\omega)\right|}\right| \omega \right\rvert\,}\right.} \begin{array}{l} 
\\
\end{array} \frac{K^{q} \frac{e^{-\left|\mu_{q}\right|}}{e^{-\frac{7}{\lambda} \alpha\left|\mu_{q-1}\right|}} \frac{e^{-\left|\mu_{q-1}\right|}}{e^{-\frac{7}{\lambda} \alpha\left|\mu_{q-2}\right|}} \cdots \frac{e^{-\left|\mu_{1}\right|}}{e^{-\frac{7}{\lambda} \alpha\left|\mu_{0}\right|}}|\omega| .}{} .
\end{aligned}
$$

Here (P3) is used in the first inequality, (P1)(ii) is used in the second and Lemma 3.1 in the third. With $\alpha<\frac{1}{100} \lambda$, the estimate above gives

$$
|\sigma|<K^{q} e^{-\sum_{k=1}^{q} \frac{9}{10}\left|\mu_{k}\right|+\frac{1}{10}\left|\mu_{0}\right|}|\omega|=K^{q} e^{-\frac{9}{10} R+\frac{1}{10}\left|\mu_{0}\right|}|\omega|:=|\sigma|_{R}
$$

(This estimate is valid only if $f^{i}(\omega)$ has completed its first bound period, which is not a problem since $n>K\left|\mu_{0}\right|$.) We estimate $|\{S>n\}|$ by

$$
|\{S>n\}|=\sum_{q, R}\left|A_{q, R}\right| \leq \sum_{R}\left(\text { number of } \sigma \text { in } \cup_{q} A_{q, R}\right) \cdot|\sigma|_{R}
$$

There are $\binom{R-1}{q}$ ways of decomposing $R$ into a sum of $q+1$ integers. For a fixed $q$-tuple $\left(\mu_{1}, \cdots, \mu_{q}\right), \mu_{i}>0$, we claim there are $\leq 2^{q} \mu_{1}^{2} \mu_{2}^{2} \cdots \mu_{q}^{2}$ possibilities for $\sigma$ with the $\mu$-coordinates of their essential free return addresses being $\left( \pm \mu_{1}, \cdots \pm \mu_{q}\right)$. This is because $f^{t_{k}}\left(\mathcal{Q}_{t_{k}}(\sigma)\right)$ is short enough that it can meet at most one $C_{\delta}(\hat{x})$, which contains $\leq 2 \mu_{k}^{2}$ intervals of the form $I_{ \pm \mu_{k} j}$. Furthermore, for $\left(\mu_{1}, \cdots, \mu_{q}\right)$ with $\left|\mu_{1}\right|+\left|\mu_{2}\right|+\cdots+\left|\mu_{q}\right|=R$, we have $\mu_{1}^{2} \mu_{2}^{2} \cdots \mu_{q}^{2} \leq\left(\frac{R}{q}\right)^{2 q}$.

There is one other piece of information that is crucial to us, namely that all bound periods are $\geq \Delta:=K^{-1} \log \frac{1}{\delta}$. This means that for a given $R$, the only feasible $q$ are $\leq \frac{R}{\Delta}$. For a fixed $R$, then, the number of $\sigma$ in $\cup_{q} A_{q, R}$ is

$$
\leq \sum_{q}\binom{R-1}{q} \cdot 2^{q}\left(\frac{R}{q}\right)^{2 q} \leq \frac{R}{\Delta} \cdot\binom{R}{\frac{R}{\Delta}} \cdot 2^{\frac{R}{\Delta}} \Delta^{2 \frac{R}{\Delta}},
$$

which, by Sterling's formula, is

$$
\sim \frac{R}{\Delta}\left(e^{\epsilon\left(\frac{1}{\Delta}\right)} 2^{\frac{1}{\Delta}} \Delta^{\frac{2}{\Delta}}\right)^{R} \quad \text { where } \quad \epsilon\left(\frac{1}{\Delta}\right) \rightarrow 0 \text { as } \delta \rightarrow 0 .
$$

Calling the expression above $(1+\eta(\delta))^{R}$, we have $\eta(\delta) \rightarrow 0$ as $\delta \rightarrow 0$.
To finish, we note that $n \leq K\left(\left|\mu_{0}\right|+R\right)$ where $K$ is as in Lemma 3.2, since the essential bound period following the $q$ th essential return expires before time $K\left(\left|\mu_{0}\right|+R\right)$. Thus

$$
\begin{aligned}
|\{S>n\}| & <\sum_{R \geq K^{-1} n-\left|\mu_{0}\right|} K^{q}(1+\eta(\delta))^{R} e^{-\frac{9}{10} R+\frac{1}{10}\left|\mu_{0}\right|}|\omega| \\
& <e^{-\frac{4}{5} K^{-1} n+\frac{1}{10}\left|\mu_{0}\right|}|\omega|<e^{-\frac{1}{2} K^{-1} n}|\omega|
\end{aligned}
$$

provided that $n>K\left|\mu_{0}\right|$.
Proof of Proposition 3.1: First there is the trival case where for some $i>0$, all points in $\omega$ have the same itinerary through time $i-1$ and $\left|f^{i}(\omega)\right|>\delta$, so that $\left.S\right|_{\omega}=i$. This case aside, let $t_{0} \geq 0$ be the first time $\mathcal{Q}_{t_{0}}$ contains more than one element. Clearly, $t_{0}<K \log \frac{1}{\delta}$, and $\left|f^{t_{0}}(\omega)\right|>K^{-1} \delta$ by (P1). Let $n>0$ be an arbitrary integer. For each $\omega^{\prime} \in \mathcal{Q}_{t_{0}}$ such that $f^{t_{0}}\left(\omega^{\prime}\right) \approx I_{\mu j},|\mu|<K^{-1} n$, we have $\left|\omega^{\prime} \cap\left\{S>t_{0}+n\right\}\right|<K_{0} e^{-\frac{1}{2} K^{-1} n}\left|\omega^{\prime}\right|$ by Lemma 3.3, $K_{0}$ here being the distortion constant for $f^{t_{0}}$. The measure of the union of $\omega^{\prime} \in \mathcal{Q}_{t_{0}}$ with $|\mu|>K^{-1} n$ is $\leq K e^{-K^{-1} n}$. It follows therefore that

$$
\left|\left\{x \in \omega: S(x)>t_{0}+n\right\}\right|<K e^{-\frac{1}{2} K^{-1} n}|\omega|+K \delta^{-1} e^{-K^{-1} n}|\omega|<e^{-K_{2}^{-1}\left(t_{0}+n\right)}|\omega|
$$

provided $K_{2}$ is sufficiently large and $n+t_{0}>K_{2} \log \delta^{-1}$.

### 3.2 Proof of Theorem 1

Let $m$ denote Lebesgue measure on $I$, and let $f_{*}^{i}(m)$ be the Borel measure with $f_{*}^{i}(m)(E)=$ $m\left(f^{-i}(E)\right)$. Fix $l:=I_{\mu_{0} j_{0}}$ for some $\mu_{0}, j_{0}$. For $n=1,2, \cdots$, let

$$
\nu_{n}=\frac{1}{n} \sum_{i=0}^{n-1} f_{*}^{i}\left(\left.m\right|_{l}\right)
$$

Clearly, any limit point of $\nu_{n}$ in the weak* topology is $f$-invariant. As we will see, it suffices to show that a positive fraction of these measures is absolutely continuous with respect to $m$ (written " $\ll m$ "). The next lemma helps us "catch" this fraction:

Lemma 3.4 There exist (i) an interval $L \subset I$, (ii) a number $c>0$, (iii) a sequence of integers $n_{1}<n_{2}<\cdots$, and (iv) for each $i=0,1,2, \cdots$, a collection of subsegments $\left\{\omega_{j}^{(i)}\right\}$ of $l$, with the property that the following hold for each $i, j$ :
(a) $f^{i}\left(\omega_{j}^{(i)}\right)=L$;
(b) $\frac{\left(f^{i}\right)^{\prime}(x)}{\left(f^{i}\right)^{\prime}(y)}<K_{0}$ for all $x, y \in \omega_{j}^{(i)}$;
(c) $\frac{1}{n_{k}} \sum_{i=0}^{n_{k}-1} m\left(\cup_{j} \omega_{j}^{(i)}\right) \geq c m(l)$.

We first finish the proof assuming the conclusion of Lemma 3.4: Let

$$
\hat{\nu}_{n_{k}}=\frac{1}{n_{k}} \sum_{i=0}^{n_{k}-1} f_{*}^{i}\left(\left.m\right|_{\left(\cup j_{j} \omega_{j}^{(i)}\right)}\right),
$$

and let $\hat{\nu}$ be a limit point of $\hat{\nu}_{n_{k}}$. It follows from Lemma 3.4(c) that $\hat{\nu}(L)>c m(l)>0$. From Lemma 3.4(a) and (b), we see that if $\rho_{k}$ is the density of $\hat{\nu}_{n_{k}}$, then $\rho_{k}(x) / \rho_{k}(y) \leq K_{0}$ for all $x, y \in L$. These bounds are passed to the limit measure $\hat{\nu}$. In particular, $\hat{\nu} \ll m$.

Now let $\nu$ be a limit point of $\nu_{n_{k}}$. Then $\nu \geq \hat{\nu}$, meaning $\nu-\hat{\nu}$ is a nonnegative measure. We decompose $\nu$ into $\nu_{a c}+\nu_{\perp}$ where $\nu_{a c} \ll m$ and $\nu_{\perp}$ is singular with respect to $m$. Since $f_{*}\left(\nu_{a c}\right) \ll m$ and $f_{*}\left(\nu_{\perp}\right) \perp m$, it follows that both $\nu_{a c}$ and $\nu_{\perp}$ are $f$-invariant. It remains to argue that $\nu_{a c}(I)>0$, which is true since $\nu_{a c} \geq \hat{\nu}$ and $\hat{\nu}(L)>0$.

Proof of Lemma 3.4: We introduce a sequence of stopping times $S_{0}<S_{1}<S_{2}<\cdots$ on $l$ as follows: Let $S_{0}=0$ and $S_{1}=S$ where $S$ is as defined in Sect. 3.1. For $k \geq 1$, let $x \in \omega$ be such that $S_{k}(x)$ has been defined. For definiteness, suppose $x \in \omega \in \mathcal{Q}_{i-1}$ with $\left.S_{k}\right|_{\omega}=i$. We define $\left.S_{k+1}\right|_{\omega}=S_{k}+\left.S\right|_{f^{S_{k}(\omega)}} \circ f^{S_{k}}$, i.e., $S_{k+1}(x)$ is defined to be the smallest $j>i$ such that $f^{j}\left(\mathcal{Q}_{j-1}(x)\right)$ is free and has length $>\delta$. Let $\tilde{\mathcal{Q}}_{i-1}=\left\{\omega \in \mathcal{Q}_{i-1}\right.$ such that $\left.S_{k}\right|_{\omega}=i$ for some $k\}$. It follows from Proposition 3.1 that for all such $\omega$,

$$
\int_{\omega}\left(S_{k+1}-S_{k}\right) d m \leq M|\omega|
$$

for some $M$ possibly dependent on $\delta$ but independent of $\omega, k$ or $i$. Keeping $k$ fixed while summing over $\omega$ and $i$, we obtain $\int_{l}\left(S_{k+1}-S_{k}\right) d m \leq M m(l)$. Summing over $k$ then gives

$$
\int_{l} S_{k} d m \leq M k m(l)
$$

By Chebychev's Inequality,

$$
m\left\{S_{\left[\frac{1}{2 M} N\right]}>N\right\} \leq \frac{1}{2} m(l) .
$$

Hence

$$
\begin{equation*}
\frac{1}{N} \sum_{i \leq N} m\left(\cup\left\{\omega \in \tilde{\mathcal{Q}}_{i-1}\right\}\right) \geq \frac{1}{4 M} m(l) \tag{2}
\end{equation*}
$$

To finish, we partition $I$ into intervals $L_{1}, L_{2}, \cdots, L_{\frac{3}{\delta}}$ of length $\frac{1}{3} \delta$ each. For each $\omega \in \tilde{\mathcal{Q}}_{i-1}$, since $\left|f^{i}(\omega)\right|>\delta$, there exists $n=\psi(\omega)$ such that $f^{i}(\omega) \supset L_{n}$. Let $\hat{\omega}=\omega \cap f^{-i} L_{n}$. By (P3), there exists $K^{\prime}$ such that $m(\hat{\omega})>K^{\prime-1} m(\omega)$. Together with (2), this implies that for each $N$, there exists $n=n(N)$ such that

$$
\frac{1}{N} \sum_{i \leq N} m\left(\cup\left\{\omega \in \tilde{\mathcal{Q}}_{i-1}: \psi(\omega)=n\right\}\right) \geq \frac{\delta}{12 M K^{\prime}} m(l)
$$

Let $n^{*}$ be such that $n^{*}=n(N)$ for infinitely many $N$. We let $L=L_{n^{*}}$, and for each $i$, let $\left\{\omega_{j}^{(i)}\right\}:=\left\{\hat{\omega}: \omega \in \tilde{\mathcal{Q}}_{i-1}, \psi(\omega)=n^{*}\right\}$.

References: A version of the material in Sect. 3.1 is used in [BC2] (in a different context). Sect. 3.2 follows the construction of SRB measures in [BY].

## Part II PARAMETER ISSUES

## 4 Admissible One-parameter Families

### 4.1 Statement of Theorem 2

We say a 1-parameter family of 1D maps $\left\{f_{a} \in C^{2}(I, I), a \in\left(a_{1}, a_{2}\right)\right\}$ is $C^{2}$ if the map $(x, a) \mapsto f_{a}(x)$ is $C^{2}$. If $I$ is an interval, we assume $f_{a}(I)$ is contained in the interior of $I$; small modifications of some statements are needed otherwise. Assuming $a_{1}<0<a_{2}$ and $f_{0} \in \mathcal{M}$, certain orbits of $f_{0}$ have natural continuations to $a$ near 0 . For example:
(i) Continuations $a \mapsto \hat{x}(a)$ of every $\hat{x} \in C\left(f_{0}\right)$ is clearly well defined.
(ii) Let $\Lambda \subset I$ be a closed subset with the property that $f_{0}(\Lambda) \subset \Lambda$ and $\Lambda \cap C\left(f_{0}\right)=\emptyset$. Then $\Lambda$ has a natural continuation $a \mapsto \Lambda(a)$ to a small interval containing 0 . Moreover, for each $x \in \Lambda, a \mapsto x(a)$ is differentiable. (For more detail, see Sect. 4.2).

Definition 4.1 Let $\left\{f_{a}\right\}$ be a 1-parameter family with $f_{0} \in \mathcal{M}$. We say $\left\{f_{a}\right\}$ satisfies the parameter transversality condition (PT) at $f_{0}$ if for every $\hat{x} \in C\left(f_{0}\right)$ and $q=f_{0}(\hat{x})$,

$$
\hat{c}(\hat{x}):=\left.\frac{d}{d a}\left(f_{a}(\hat{x}(a))-q(a)\right)\right|_{a=0} \neq 0
$$

The notation " $q(a)$ " in the displayed formula above is to be interpreted as follows: Since $f_{0} \in \mathcal{M}, q$ is contained in a closed subset $\Lambda$ of the kind in (ii) above. By $q(a)$, we refer to the continuation of $q$ in the sense of (ii).

As $a$ varies, $f_{a}(\hat{x}(a))$ moves with $a$, as does the set $\Lambda(a)$. Roughly speaking, (PT) stipulates that the two move at different speeds. In general, $\Lambda$ moves more slowly, and the trajectory of $a \mapsto f_{a}(\hat{x}(a))$ moves "through" $\Lambda$. This is why we think of (PT) as a transversality condition.

Theorem 2 For every $C^{2}$ family $\left\{f_{a}\right\}$ satisfying (PT) at $f_{0} \in \mathcal{M}$, the set $\left\{a: f_{a} \in \mathcal{G}\right\}$ has positive Lebesgue measure.

The rest of this paper is devoted to the proof of Theorem 2.
To simplify slightly the discussion, we assume from here on that $C\left(f_{a}\right)=C\left(f_{0}\right):=C$ for all $a$. This is easily arranged via $a$-dependent changes of coordinates that do not affect the content of the theorem. As before, critical points will be denoted by "hats" (e.g., $\hat{x} \in C$, while $x$ is an arbitrary point in $I$ ).

## Standing assumptions for Part II:

- $\left(a_{1}, a_{2}\right)$ is an interval with $a_{1}<0<a_{2}$;
- $\left\{f_{a}, a \in\left(a_{1}, a_{2}\right)\right\}$ is a $C^{2}$ family with $C\left(f_{a}\right)=C\left(f_{0}\right)$ for all $a$;
- $f_{0} \in \mathcal{M}$, and (PT) is satisfied at $f_{0}$.

System constants for Part II are allowed to depend only on (i) $f_{0}$ (including the constants in Sect. 1.1 associated with $f_{0}$ ), (ii) the $C^{2}$ norm of the family $\left\{f_{a}\right\}$, (iii) $\hat{c}:=\min _{\hat{x} \in C}|\hat{c}(\hat{x})|$ and (iv) our choice of $\lambda$.

### 4.2 Alternate formulation of (PT)

We begin with some simple facts about the symbolic dynamics of $f=f_{0} \in \mathcal{M}$. Let $\mathcal{J}=\left\{J_{1}, \cdots, J_{q}\right\}$ be the components of $I \backslash C$. For $x \in I$ such that $f^{i} x \notin C$ for all $i \geq 0$, let $\phi(x)=\left(\iota_{i}\right)_{i=0,1, \ldots}$ be given by $\iota_{i}=k$ if $f^{i} x \in J_{k}$.

Lemma 4.1 For $f \in \mathcal{M}$, there exists an increasing sequence of compact sets $\Lambda^{(n)}$ such that (a) $\Lambda^{(n)} \cap C=\emptyset, f\left(\Lambda^{(n)}\right) \subset \Lambda^{(n)}$, and $\left.f\right|_{\Lambda^{(n)}}$ is conjugate to a shift of finite type;
(b) $\cup_{n} \Lambda^{(n)}$ is dense in $I$;
(c) if $\inf _{i \geq 0} d\left(f^{i}(x), C\right)>0$, then $x \in \Lambda^{(n)}$ for some $n$.

Proof: First we argue that $\cup_{i \geq 0} f^{-i} C$ is dense in $I$. If not, there would be an interval $\omega$ with the property that $\phi(x)$ is identical for all $x \in \omega$. Let $\omega$ be a maximal interval of this type. Then either (i) $f^{n+k}(\omega) \subset f^{n}(\omega)$ for some $n, k$, or (ii) $f^{k}(\omega), k=0,1, \cdots$, are pairwise disjoint. Case (i) cannot happen since it implies the presence of a periodic point $x$ with $\left|\left(f^{k}\right)^{\prime} x\right| \leq 1$. Case (ii) is equally absurd, for it implies the existence of $\left\{k_{i}\right\}$ where $f^{k_{i}}(\omega)$ are arbitrarily short. We leave it as an easy exercise to see that this is incompatible with the definition of $\mathcal{M}$.

For definiteness, consider $I=S^{1}$. Let $l_{n}(\hat{x})$ and $r_{n}(\hat{x})$ be the two points in $\cup_{0<i \leq n} f^{-i} C$ closest to $\hat{x} \in C$, and let $\Lambda^{(n)}=\left\{x \in I: f^{i} x \notin \cup_{\hat{x} \in C}\left(l_{n}(\hat{x}), r_{n}(\hat{x})\right) \forall i \geq 0\right\}$. Then for each $n, \Lambda^{(n)}$ is compact and $f\left(\Lambda^{(n)}\right) \subset \Lambda^{(n)}$. That $\cup_{n} \Lambda^{(n)}$ is dense in $I$ follows from (i) $\cup_{0<i \leq n} f^{-i} C \subset \Lambda^{(n)}$ for all $n$ large enough that $\left(l_{n}(\hat{x}), r_{n}(\hat{x})\right) \subset\left(\hat{x}-\delta_{0}, \hat{x}+\delta_{0}\right)$, and (ii) $\cup_{i>0} f^{-i} C$ is dense in $I$. Assertion (c) follows immediately from this construction.

To show that $\left.f\right|_{\Lambda^{(n)}}$ is conjugate to a shift of finite type, let $\mathcal{J}^{(n)}=\left\{J_{i}^{(n)}\right\}$ be the partition of $I$ by $\cup_{0 \leq i \leq n} f^{-i} C$. For $J_{i}^{(n)} \neq\left(l_{n}(\hat{x}), \hat{x}\right)$ or $\left(\hat{x}, r_{n}(\hat{x})\right)$, observe that by construction, $f\left(J_{i}^{(n)}\right)$ is equal to the union of a finite number of elements of $\mathcal{J}^{(n)}$. Let $\Lambda_{i}^{(n)}=\Lambda^{(n)} \cap J_{i}^{(n)}$.

Then the alphabet of the shift in question is $\left\{i: \Lambda_{i}^{(n)} \neq \emptyset\right\}$, and the transition $i \rightarrow j$ is admissible if $f\left(\Lambda_{i}^{(n)}\right) \supset \Lambda_{j}^{(n)}$.

Where $I$ is an interval, $\Lambda^{(n)}$ is as above but restricted to the interval $\left[z_{n}^{1}, z_{n}^{2}\right]$ where $z_{n}^{1}$ and $z_{n}^{2}$ are the two points in $\cup_{0<i \leq n} f^{-i} C$ closest to the endpoints of $I$. The end intervals are excluded because they do not have the required Markov property.

Our next result guarantees that $q(a)$ in Definition 4.1 is well defined.
Corollary 4.1 For $f \in \mathcal{M}$, let $q \in I$ be such that $\delta_{1}:=\inf _{n \geq 0} d\left(f^{n}(q), C\right)>0$. Then for all $g$ with $\|g-f\|_{C^{2}}<\varepsilon$ where $\varepsilon=\varepsilon\left(\delta_{1}\right)$, there is a unique point $q_{g} \in I$ with $\phi_{g}\left(q_{g}\right)=\phi_{f}(q)$.

Proof: Fix $n$ large enough that for all $i \geq 0, f^{i}(q) \notin\left(l_{n}(\hat{x}), r_{n}(\hat{x})\right)$ for all $\hat{x} \in C$, and let $\Lambda=\Lambda^{(n)}$. Let $B=\cup_{i} \partial \bar{\Lambda}_{i}$ where $\bar{\Lambda}_{i}$ is the shortest interval containing $\Lambda_{i}$. Since $B$ is a finite set with $f(B) \subset B$, it consists of preperiodic points. From Lemma 1.1, the periodic points in question are repelling. Thus if $g$ is sufficiently near $f$, there is a unique set $B_{g}$ with $g\left(B_{g}\right) \subset B_{g}$ such that $\left.g\right|_{B_{g}}$ is conjugate to $\left.f\right|_{B}$. Using $B_{g}$, we recover a set $\Lambda_{g}$ on which $g$ is conjugate to $\left.f\right|_{\Lambda}$. The uniqueness of $q_{g}$ follows from the expanding property of $g$ away from $C$ (Lemma 1.1).

We return now to the $C^{2}$ family $\left\{f_{a}\right\}$ with $f_{0} \in \mathcal{M}$. Let $\hat{x} \in C$ be fixed, and let $q=f_{0}(\hat{x})$. We write $F(x, a):=f_{a}(x)$, reserving $(\cdot)^{\prime}$ for $x$-derivatives, i.e., $\partial_{x} F(x, a)=\left(f_{a}\right)^{\prime}(x)$.

Proposition 4.1 There is an interval $\omega$ in a-space containing 0 in its interior on which $q(a)$ is defined, $a \mapsto q(a)$ is differentiable, and

$$
\begin{equation*}
\frac{d}{d a} q(a)=-\sum_{i=1}^{\infty} \frac{\partial_{a} F\left(f_{a}^{i-1}(q(a)), a\right)}{\left(f_{a}^{i}\right)^{\prime}(q(a))} \tag{3}
\end{equation*}
$$

Proof: Let $q$ be as in the proof of Corollary 4.1. We assume $\omega$ is short enough that for all $a \in \omega, q(a)$ is well defined, $q(a) \cap C_{\frac{1}{2} \delta_{0}}=\emptyset$, and that (P1) holds for $f_{a}$ outside of $C_{\frac{1}{2} \delta_{0}}$ with uniform bounds. In the computations below, we suppress the dependence on $a$, writing $f=f_{a}, q=q(a), \partial_{a} F(\cdot)=\partial_{a} F(\cdot, a)$, and so on.

Continuing to use the notation in Corollary 4.1, we let $\Lambda_{i_{0}, i_{1}, \cdots, i_{m}}=\left\{x \in I: f^{j}(x) \in\right.$ $\left.\Lambda_{i_{j}}, 0 \leq j \leq m\right\}$, and let $\Lambda_{i_{0}, i_{1}, \cdots, i_{m}}(q)$ be the cylinder set containing $q$. For each $m$, choose $q^{(m)} \in \partial \bar{\Lambda}_{i_{0}, i_{1}, \cdots, i_{m}}(q)$. Then $q^{(m)} \rightarrow q$. It suffices to show that as functions of $a, \frac{d}{d a} q^{(m)}$ converges uniformly to the right side of (3). Let $p^{(m)}=f^{m}\left(q^{(m)}\right)$. Differentiating, we obtain

$$
\frac{d}{d a} p^{(m)}=\sum_{i=1}^{m}\left(f^{m-i}\right)^{\prime}\left(f^{i}\left(q^{(m)}\right)\right) \partial_{a} F\left(f^{i-1}\left(q^{(m)}\right)\right)+\left(f^{m}\right)^{\prime}\left(q^{(m)}\right) \frac{d}{d a} q^{(m)}
$$

This gives

$$
\begin{equation*}
\frac{d}{d a} q^{(m)}=\frac{\frac{d}{d a} p^{(m)}}{\left(f^{m}\right)^{\prime}\left(q^{(m)}\right)}-\sum_{i=1}^{m} \frac{\partial_{a} F\left(f^{i-1} q^{(m)}\right)}{\left(f^{i}\right)^{\prime}\left(q^{(m)}\right)} \tag{4}
\end{equation*}
$$

We stress that all the action below takes place outside of $C_{\frac{1}{2} \delta_{0}}$, where $\left|\left(f^{n}\right)^{\prime}\right|$ grows exponentially (with prefactor $\frac{1}{2} \delta_{0} c_{1}$ ).

To estimate (4), observe that since $p^{(m)} \in B$ (see the proof of Corollary 4.1), and $B$ is a finite set, $\frac{d}{d a} p^{(m)}$ is uniformly bounded for all $m$. With $\left|\left(f^{m}\right)^{\prime}\left(q^{(m)}\right)\right|$ growing exponentially, the first term on the right is exponentially small. It remains to check that the second term converges uniformly to the right side of (3). Since the tail of the sum in (3) decreases exponentially (uniformly bounded numerator and exponentially increasing denominator), it suffices to verify that

$$
\begin{aligned}
A & :=\left|\sum_{i=1}^{m} \frac{\partial_{a} F\left(f^{i-1}\left(q^{(m)}\right)\right)}{\left(f^{i}\right)^{\prime}\left(q^{(m)}\right)}-\sum_{i=1}^{m} \frac{\partial_{a} F\left(f^{i-1}(q)\right)}{\left(f^{i}\right)^{\prime}(q)}\right| \\
& \leq \sum_{i=1}^{m} \frac{\left|\partial_{a} F\left(f^{i-1}\left(q^{(m)}\right)\right)-\partial_{a} F\left(f^{i-1}(q)\right)\right|}{\left|\left(f^{i}\right)^{\prime}\left(q^{(m)}\right)\right|}+\sum_{i=1}^{m} \frac{\left|\partial_{a} F\left(f^{i-1}(q)\right)\right|}{\left|\left(f^{i}\right)^{\prime}\left(q^{(m)}\right)\right|} \cdot\left|\frac{\left(f^{i}\right)^{\prime}\left(q^{(m)}\right)}{\left(f^{i}\right)^{\prime}(q)}-1\right|
\end{aligned}
$$

converges to 0 uniformly. Consider the $i$ th term in the first sum. By the expanding property of $f$ outside of $C_{\frac{1}{2} \delta_{0}}$, the numerator is $<$ const $e^{-\frac{1}{4} \lambda_{0}(m-i)}\left|f^{m}\left(q^{(m)}\right)-f^{m}(q)\right|$, while the denominator is $>$ const $e^{\frac{1}{4} \lambda_{0} i}$. The second sum is estimated similarly. Together they give $A<$ const $m e^{-\frac{1}{4} \lambda_{0} m}$.

Let $\hat{x}_{i}(a):=f_{a}^{i}(\hat{x}(a))$, and recall the definition of $\hat{c}(\hat{x})$ in Definition 4.1.

## Corollary 4.2

$$
\hat{c}(\hat{x})=\frac{d \hat{x}_{1}}{d a}(0)+\sum_{i=1}^{\infty} \frac{\partial_{a} F\left(\hat{x}_{i}(0), 0\right)}{\left(f_{0}^{i}\right)^{\prime}\left(\hat{x}_{1}(0)\right)} .
$$

We have thus obtained an equivalent formulation of (PT) that involves only $\partial_{a} F(\cdot, \cdot)$ and properties of $f_{0}$.

### 4.3 Comparability of $x$ - and $a$-derivatives

For $f_{0}, \lambda, \alpha$ and $\varepsilon$ as in Proposition 2.1, we define

$$
\mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right):=\left\{f:\left\|f-f_{0}\right\|_{C^{2}}<\varepsilon \text { and }(G 1),(G 2) \text { hold for all } \hat{x} \in C \text { and } n \leq N\right\} .
$$

Proposition 4.2 Let $\lambda, \alpha$ and $\varepsilon$ be fixed. Then there exist $\hat{\varepsilon}>0$ and $\hat{i} \in \mathbb{Z}^{+}$such that the following holds for all $N \in \mathbb{Z}^{+}$: Let $\Omega_{N} \subset(-\hat{\varepsilon}, \hat{\varepsilon})$ be such that $f_{a} \in \mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ for all $a \in \Omega_{N}$. Then for every $a \in \Omega_{N}$ and $\hat{x} \in C$,

$$
\frac{1}{2}|\hat{c}(\hat{x})|<\frac{\left|\frac{d}{d a} \hat{x}_{i}(a)\right|}{\left|\left(f_{a}^{i-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|}<2|\hat{c}(\hat{x})| \quad \text { for } \hat{i}<i \leq N
$$

Proof: Writing

$$
\frac{d}{d a} \hat{x}_{i}(a)=\left(f_{a}\right)^{\prime}\left(\hat{x}_{i-1}\right) \frac{d}{d a} \hat{x}_{i-1}(a)+\partial_{a} F\left(\hat{x}_{i-1}, a\right),
$$

we obtain inductively

$$
\frac{\frac{d}{d a} \hat{x}_{i}(a)}{\left(f_{a}^{i-1}\right)^{\prime}\left(\hat{x}_{1}\right)}=\frac{d}{d a} \hat{x}_{1}(a)+\sum_{j=1}^{i-1} \frac{\partial_{a} F\left(\hat{x}_{j}, a\right)}{\left(f_{a}^{j}\right)^{\prime}\left(\hat{x}_{1}\right)}
$$

Letting $I(a, i)$ denote the expression on the right side above, we choose $\hat{i}$ large enough that (i) $I(0, \hat{i}) \approx \hat{c}(\hat{x})$ and (ii) for $i>\hat{i}$,

$$
\left|\sum_{j=\hat{i}}^{i-1} \frac{\partial_{a} F\left(\hat{x}_{j}, a\right)}{\left(f_{a}^{j}\right)^{\prime}\left(\hat{x}_{1}\right)}\right| \ll|\hat{c}(\hat{x})| \quad \text { uniformly for all } a \in J
$$

(i) makes sense because $\hat{c}(\hat{x}) \neq 0$ by (PT). (ii) is because $\left|\partial_{a} F\left(\hat{x}_{j}, a\right)\right|<K$ and $\left|\left(f_{a}^{j}\right)^{\prime}\left(\hat{x}_{1}\right)\right|>$ $c_{1} e^{\lambda j}$ from (G2). Since only a finite number of iterates are involved, we may now shrink $\hat{\varepsilon}$ sufficiently so that $|I(a, \hat{i})-I(0, \hat{i})| \ll|\hat{c}(\hat{x})|$ for all $a \in(-\hat{\varepsilon}, \hat{\varepsilon})$.
References: The comparability of $x$ - and $a$-derivatives is first used in [ BC 1$]$; the unimodal case of the material in Sects. 4.2 and 4.3 is in [TTY].

## 5 Evolution of Critical Curves

### 5.1 Logistics

Much of the rest of the analysis revolves around evolutions of the type

$$
a \mapsto \hat{x}_{i}(a), \quad i=0,1,2, \cdots, \quad \text { for } \hat{x} \in C, \quad a \in J,
$$

where $J$ is an interval in $(-\hat{\varepsilon}, \hat{\varepsilon})$. Via the geometry of these curves, we seek to determine what fraction of $J$ corresponds to "bad" parameters, that is to say, what fraction of these curves comes too close to the critical set (therefore violating (G1)), or visits the critical set too often (thereby violating (G2)). Herein lies the dilemma: In order for the curves $a \mapsto \hat{x}_{i}(a)$ to have controlled geometry, the maps $f_{a}$ corresponding to the $a$ 's involved must (individually) be good to start with. On the other hand, by studying curve segments corresponding to good parameters only, how are we to determine what fraction of parameters are bad?

The following discussion motivates our answer to this logistical dilemma.
Properties of $f \in \mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ up to time $\frac{1}{\alpha^{*}} N, \alpha^{*}:=\frac{3}{\lambda} \alpha$
Assume $f \in \mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$, and consider $\hat{x} \in C$.
Suppose for some $n \leq \frac{1}{\alpha^{*}} N, \hat{x}_{n} \in C_{\delta}$ with $d\left(\hat{x}_{n}, C\right)>e^{-\alpha n}$. We claim that (P2) in Sect. 2.1 holds for this return even though $n$ may be greater than $N$. This is because the critical point that will guide $\hat{x}_{n}$ through its partial derivative recovery obeys (G1) and (G2) up to time $N$, and by the proof of (P2), the time it takes to complete this recovery is $<\frac{3}{\lambda} \alpha n \leq N$.

Indeed if we assume (G1) holds for $\hat{x}$ up to time $n, n \leq \frac{1}{\alpha^{*}} N$, then on the time interval $[0, n]$, the orbit of $\hat{x}_{1}$ has the bound/free behavior described in Sect. 2.1. Moreover, by an argument identical to that for Corollary 2.1, we have $\left|\left(f^{j}\right)^{\prime}\left(\hat{x}_{1}\right)\right|>K^{-1} e^{\frac{1}{4} \lambda j}$ for $j \leq n$.

We remark that beyond time $\frac{1}{\alpha^{*}} N$, the dynamical description of $\hat{x}$ in the last paragraph ceases to be valid as soon as a bound period $>N$ is encountered. Conversely, the behaviors of other critical orbits beyond time $N$ do not impact the properties of $\hat{x}$ up to time $\frac{1}{\alpha^{*}} N$.

In view of the discussion above, we modify Proposition 4.2 slightly as follows:
Proposition 4.2' In addition to the hypotheses in Proposition 4.2, we assume that for some $\hat{x} \in C$ and $n \in\left(N, \frac{1}{\alpha^{*}} N\right]$, (G1) holds for $\hat{x}$ up to time $n$. Then the conclusion of Proposition 4.2 holds for this $\hat{x}$ for all $i \leq n$.

### 5.2 Duality between phase-space and parameter-space dynamics

Setting. Let $\lambda<\frac{1}{4} \lambda_{0}$ be as before. To establish the above-mentioned duality, new upper bounds are imposed on $\alpha$ and $\varepsilon$ (or equivalently $\hat{\varepsilon}$ ). Let $\Omega_{N}=\left\{a \in(-\hat{\varepsilon}, \hat{\varepsilon}): f_{a} \in\right.$ $\left.\mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)\right\}$. For the rest of Section 5, we fix $\hat{x} \in C$. All parameters considered are assumed to be in $\Omega_{N}$; all indices considered are assumed to be $\leq \frac{1}{\alpha^{*}} N$, and (G1) is assumed to hold for $\hat{x}$ for all the indices in question. We use the notation $\tau_{i}(a):=\frac{d}{d a} \hat{x}_{i}(a)$.

Our main results are (P1')-(P3'), three properties of $a \mapsto \hat{x}_{i}(a)$ that are the analogs of (P1)-(P3) in Sect. 2.1. We state also two lemmas that lie at the heart of these properties. To avoid disrupting the flow of ideas, proofs are postponed to Sect. 5.3.

Lemma 5.1 Let $n>\hat{i}$ where $\hat{i}$ is as in Proposition 4.2. Then

$$
\left(1-K e^{-\frac{1}{4} \lambda n}\right)\left|\left(f_{a}^{i}\right)^{\prime}\left(\hat{x}_{n}\right)\right| \leq \frac{\left|\tau_{n+i}\right|}{\left|\tau_{n}\right|} \leq\left(1+K e^{-\frac{1}{4} \lambda n}\right)\left|\left(f_{a}^{i}\right)^{\prime}\left(\hat{x}_{n}\right)\right| .
$$

(P1') (Outside of $C_{\delta}$ ): There exists $i_{0} \geq \hat{i}$ such that the following hold for $n \geq i_{0}$ :
(i) If $\hat{x}_{n}$ is free, and $\hat{x}_{n+j} \notin C_{\delta} \forall 0 \leq j<j_{0}$, then $\left|\tau_{n+j}\right|>\frac{1}{2} c_{1} \delta e^{\frac{1}{4} \lambda_{0} j}\left|\tau_{n}\right|$ for $j \leq j_{0}$;
(ii) if in addition $\hat{x}_{n+j_{0}} \in C_{\delta_{0}}$, then $\left|\tau_{n+j_{0}}\right|>\frac{1}{2} c_{1} e^{\frac{1}{4} \lambda_{0} j_{0}}\left|\tau_{n}\right|$.

The reader should think of $i_{0}$ as the time after which $x$ - and $a$-derivatives are sufficiently close in the sense of Lemma 5.1. We assume $\hat{x}_{i} \notin C_{\frac{1}{2} \delta_{0}}$ for all $i<i_{0}$.

Consider next an interval $\omega \subset \Omega_{N}$ with $\hat{x}_{n}(\omega) \subset I_{\mu j}^{+}$. To establish the desired relationship between phase-space and parameter-space dynamics during the bound period, we impose the following additional upper bound on $\alpha$ : Let $L$ be a Lipschitz constant of the $\operatorname{map} G:(x, a) \mapsto\left(f_{a}(x), a\right)$. We assume $\alpha$ is small enough that

$$
\begin{equation*}
L^{\frac{3}{\lambda} \alpha}<e^{\frac{1}{8} \lambda} . \tag{5}
\end{equation*}
$$

For each $a \in \omega$, let $p_{a}$ denote the bound period of $I_{\mu j}$ with respect to $f_{a}$, and let $\operatorname{HD}(\cdot, \cdot)$ denote the Hausdorff distance between two sets.

Lemma 5.2 Let $\omega$ and $\alpha$ be as above. Then the following hold for all $a \in \omega$ :

$$
\operatorname{HD}\left(\hat{x}_{n+j}(\omega), f_{a}^{j}\left(\hat{x}_{n}(\omega)\right)\right) \ll e^{-2 \alpha j} \quad \text { for all } j<p_{a} .
$$

We define the bound period $\hat{p}_{n}(\omega)$ of $\hat{x}_{n}(\omega)$ in parameter-space dynamics to be

$$
\hat{p}_{n}(\omega):=\min \left\{p_{a}: a \in \omega\right\} .
$$

(P2') (Partial derivative recovery): Suppose $\hat{x}_{n}(\omega) \subset I_{\mu j}^{+}$, and let $\hat{p}=\hat{p}_{n}(\omega)$. Then
(a) $\frac{1}{3 \ln \left(\text { max }\left|f^{\prime}\right|\right)}|\mu| \leq \hat{p} \leq \frac{3}{\lambda}|\mu|$;
(b) for $a, a^{\prime} \in \omega$ and $j<\hat{p},\left|\hat{x}_{n+j}(a)-\hat{x}_{n+j}\left(a^{\prime}\right)\right|<2 e^{-2 \alpha j}$;
(c) $\left|\tau_{n+\hat{p}}(a)\right|>e^{\frac{\lambda \hat{p}}{4}}\left|\tau_{n}(a)\right|$ for all $a \in \omega$;
(d) if $\hat{x}_{n}(\omega) \approx I_{\mu j}$, then $\left|\hat{x}_{n+\hat{p}}(\omega)\right| \geq e^{-\frac{8 \alpha}{\lambda}|\mu|}$.

To state (P3'), we divide each orbit in the time interval $\left[i_{0}, n\right]$ into bound and free periods as in Sect. 2.1, and say all $a \in \omega$ have the same itinerary up to time $n$ if (i) their bound and free periods coincide and (ii) whenever $\hat{x}_{i}(\omega)$ is free, it is $\subset \pi^{+}$for some $\pi \in \mathcal{P}$.
(P3') (Global distortion): There exist $i_{1}>i_{0}$ and $K_{3}>1$ such that if $n \geq i_{1}$ and all points in $\omega$ have the same itinerary through step $n-1$, then for all $a, a^{\prime} \in \omega$,

$$
K_{3}^{-1}<\frac{\tau_{n}(a)}{\tau_{n}\left(a^{\prime}\right)}<K_{3} .
$$

Shrinking $\varepsilon$ if necessary, we assume from here on that $\hat{x}_{i} \notin C_{\frac{1}{2} \delta_{0}}$ for all $i<i_{1}$.

### 5.3 Proofs of (P1')-(P3')

Proof of Lemma 5.1: Let $\xi=\hat{x}_{n}$. From

$$
\tau_{n+i}=\left(f_{a}\right)^{\prime}\left(\xi_{i-1}\right) \tau_{n+i-1}+\partial_{a} F\left(\xi_{i-1}, a\right)
$$

we deduce inductively that

$$
\begin{aligned}
\tau_{n+i} & =\left(f_{a}^{i}\right)^{\prime}(\xi) \tau_{n}+\sum_{j=0}^{i-1}\left(f^{i-j-1}\right)^{\prime}\left(\xi_{j+1}\right) \partial_{a} F\left(\xi_{j}, a\right) \\
& =\left(f_{a}^{i}\right)^{\prime}(\xi) \tau_{n}\left(1+\sum_{j=0}^{i-1} \frac{\partial_{a} F\left(\xi_{j}, a\right)}{\left(f_{a}^{j+1}\right)^{\prime}(\xi) \tau_{n}}\right) .
\end{aligned}
$$

Proposition 4.2 ' and Corollary 2.1 then give

$$
\left|\sum_{j=0}^{i-1} \frac{\partial_{a} F\left(\xi_{j}, a\right)}{\left(f_{a}^{j+1}\right)^{\prime}(\xi) \tau_{n}}\right| \leq 2|\hat{c}(\hat{x})| \sum_{j=0}^{i-1}\left|\frac{\partial_{a} F\left(\xi_{j}, a\right)}{\left(f_{a}^{n+j}\right)^{\prime}\left(\hat{x}_{1}(a)\right)}\right| \leq K e^{-\frac{1}{4} \lambda n} .
$$

Proof of (P1'): (P1') follows immediately from (P1) via Lemma 5.1. With $i_{0}$ large enough, the $x$ - and $a$-derivatives are as close as need be.

Proof of Lemma 5.2: It suffices to consider the end points of the segments to be compared. Suppose $\omega=\left[\bar{a}_{1}, \bar{a}_{2}\right]$. Then for $i=1,2$,

$$
\begin{equation*}
\left|\hat{x}_{n+j}\left(\bar{a}_{i}\right)-f_{a}^{j}\left(\hat{x}_{n}\left(\bar{a}_{i}\right)\right)\right|=\left|G^{j}\left(\hat{x}_{n}\left(\bar{a}_{i}\right), \bar{a}_{i}\right)-G^{j}\left(\hat{x}_{n}\left(\bar{a}_{i}\right), a\right)\right| \leq L^{j}\left|\bar{a}_{i}-a\right| \leq L^{j}|\omega| . \tag{6}
\end{equation*}
$$

By Proposition 4.2 and Corollary 2.1, $|\omega| \leq K e^{-\frac{1}{4} \lambda n}$. Also, $j \leq p_{a}$, which, by (P2), is $\leq \frac{3}{\lambda} \alpha n$. Thus $L^{j}|\omega| \ll e^{-2 \alpha j}$ if $\alpha$ satisfies (5).

Proof of (P2'): (a) is true because $\hat{p}=p_{\bar{a}}$ for some $\bar{a} \in \omega$. (b) and (d) are immediate consequences of Lemma 5.2, and (c) follows from Lemma 5.1.

Turning now to the setting of (P3'), we let $a, a^{\prime} \in \omega$, and for some $k$ with $i_{0} \leq k<n$, let $\xi=\hat{x}_{k}(a)$ and $\xi^{\prime}=\hat{x}_{k}\left(a^{\prime}\right)$.

Lemma 5.3 There exists $K>0$ such that for all $i$ such that $k+i \leq n$,

$$
\frac{\left(f_{a}^{i}\right)^{\prime}(\xi)}{\left(f_{a^{\prime}}^{i}\right)^{\prime}\left(\xi^{\prime}\right)} \leq \exp \left\{\sum_{j=0}^{i-1} K \frac{\left|\xi_{j}-\xi_{j}^{\prime}\right|}{d\left(\xi_{j}, C\right)}\right\}
$$

Proof: First we write

$$
\begin{aligned}
\log \frac{\left(f_{a^{\prime}}^{i}\right)^{\prime}\left(\xi^{\prime}\right)}{\left(f_{a}^{i}\right)^{\prime}(\xi)} & \leq \sum_{j=0}^{i-1} \frac{\left|\left(f_{a}\right)^{\prime}\left(\xi_{j}\right)-f_{a^{\prime}}^{\prime}\left(\xi_{j}^{\prime}\right)\right|}{\left|\left(f_{a}\right)^{\prime}\left(\xi_{j}\right)\right|}<\sum_{j=0}^{i-1} K \frac{\left|\xi_{j}-\xi_{j}^{\prime}\right|+\left|a-a^{\prime}\right|}{d\left(\xi_{j}, C\right)} \\
& <K \sum_{j=0}^{i-1}\left(1+\frac{\left|a-a^{\prime}\right|}{\left|\xi_{j}-\xi_{j}^{\prime}\right|}\right) \frac{\left|\xi_{j}-\xi_{j}^{\prime}\right|}{d\left(\xi_{j}, C\right)} .
\end{aligned}
$$

Then we use Proposition $4.2^{\prime}$ and Corollary 2.1 to estimate the quantity in parenthesis: Assuming for definiteness that $a^{\prime}<a$,

$$
\begin{equation*}
\left|\xi_{j}-\xi_{j}^{\prime}\right|=\int_{a^{\prime}}^{a}\left|\tau_{k+j}(s)\right| d s \geq \frac{1}{2}|\hat{c}(\hat{x})| \int_{a^{\prime}}^{a}\left|\left(f_{s}^{k+j}\right)^{\prime}\left(\hat{x}_{1}(s)\right)\right| d s \geq \frac{1}{2}|\hat{c}(\hat{x})| e^{\frac{1}{4} \lambda k}\left|a-a^{\prime}\right| . \tag{7}
\end{equation*}
$$

Proof of (P3'): By Proposition 4.2', it suffices to show

$$
\frac{\left(f_{a^{\prime}}^{n-1}\right)^{\prime}\left(\hat{x}_{1}\left(a^{\prime}\right)\right)}{\left(f_{a}^{n-1}\right)^{\prime}\left(\hat{x}_{1}(a)\right)}<K
$$

the proof of which follows closely that of (P3) in Sect. 2.2. Let $i_{0}<t_{1}<t_{1}+p_{1} \leq t_{2}<$ $t_{2}+p_{2} \leq \cdots$ where $t_{i}$ are free return times and $p_{i}$ are the associated bound periods. For definiteness we assume $t_{q}+p_{q} \leq n \leq t_{q+1}$. Then

$$
\log \frac{\left(f_{a^{\prime}}^{n-1}\right)^{\prime}\left(\hat{x}_{1}\left(a^{\prime}\right)\right)}{\left(f_{a}^{n-1}\right)^{\prime}\left(\hat{x}_{1}(a)\right)} \leq \log \frac{\left(f_{a^{\prime}}^{t_{1}-1}\right)^{\prime}\left(\hat{x}_{1}\left(a^{\prime}\right)\right)}{\left(f_{a}^{t_{1}-1}\right)^{\prime}\left(\hat{x}_{1}(a)\right)}+\sum_{k=1}^{q}\left(S_{k}^{\prime}+S_{k}^{\prime \prime}\right)
$$

where

$$
S_{k}^{\prime}=\log \frac{\left|\left(f_{a^{\prime}}^{p_{k}}\right)^{\prime}\left(\hat{x}_{t_{k}}\left(a^{\prime}\right)\right)\right|}{\left|\left(f_{a}^{p_{k}}\right)^{\prime}\left(\hat{x}_{t_{k}}(a)\right)\right|} \quad \text { and } \quad S_{k}^{\prime \prime}=\log \frac{\left|\left(f_{a^{\prime}}^{t_{k+1}-\left(t_{k}+p_{k}\right)}\right)^{\prime}\left(\hat{x}_{t_{k}+p_{k}}\left(a^{\prime}\right)\right)\right|}{\left|\left(f_{a}^{t_{k+1}-\left(t_{k}+p_{k}\right)}\right)^{\prime}\left(\hat{x}_{t_{k}+p_{k}}(a)\right)\right|},
$$

except for $S_{q}^{\prime \prime}$, which ends at time $n-1$ instead of $t_{q+1}$.
Let $\sigma_{k}=\left[\hat{x}_{t_{k}}(a), \hat{x}_{t_{k}}\left(a^{\prime}\right)\right]$. It follows from (P2')(c) and (P1')(ii) that, for $k<q,\left|\sigma_{k+1}\right| \geq$ $\frac{1}{2} c_{1} e^{\frac{1}{4} \lambda\left(t_{k+1}-t_{k}\right)}\left|\sigma_{k}\right|$, which we may assume is $\geq \hat{\tau}\left|\sigma_{k}\right|$ for some $\hat{\tau}>1$ (the factor $\frac{1}{2} c_{1}$ is again absorbed into the exponential assuming $\delta$ is sufficiently small).

In the proof of (P3), the derivative estimates in $S_{k}^{\prime}$ and $S_{k}^{\prime \prime}$ are converted to distance estimates. This is exactly what we have done in Lemma 5.3. More precisely, if $\xi=\hat{x}_{t_{k}}(a)$ and $\xi^{\prime}=\hat{x}_{t_{k}}\left(a^{\prime}\right)$, then by Lemma 5.3,

$$
S_{k}^{\prime} \leq K \sum_{j=0}^{p_{k}-1} \frac{\left|\xi_{j}-\xi_{j}^{\prime}\right|}{d\left(\xi_{j}, C\right)} \leq K \sum_{j=0}^{p_{k}-1} \frac{\left|f_{a}^{j}(\xi)-f_{a}^{j}\left(\xi^{\prime}\right)\right|}{d\left(\xi_{j}, C\right)}+K \sum_{j=0}^{p_{k}-1} \frac{\left|f_{a}^{j}\left(\xi^{\prime}\right)-f_{a^{\prime}}^{j}\left(\xi^{\prime}\right)\right|}{d\left(\xi_{j}, C\right)} .
$$

The first sum on the right involves only the map $f_{a}$ and is $\leq K \frac{\left|\sigma_{k}\right|}{d(\xi, C)}$ by the corresponding argument in the proof of (P3). For the second sum, we combine the estimate in (6) with the result of (7) to show that it is

$$
\leq K\left(\sum_{j=0}^{p_{k}-1}\left(L e^{\alpha}\right)^{j}\right)\left|a-a^{\prime}\right| \leq K\left(L e^{\alpha}\right)^{\frac{3}{\lambda} \alpha t_{k}} e^{-\frac{1}{4} \lambda t_{k}}\left|\sigma_{k}\right| \ll\left|\sigma_{k}\right| .
$$

After this, the rest of the proof is as before. The sums $S_{k}^{\prime \prime}$ (including $k=q$ ) are estimated similarly, the only difference being that the exponential growth of $\left|\xi_{j}-\xi_{j}^{\prime}\right|$ here is derived from ( $\mathrm{P} 1^{\prime}$ ).

It remains to treat the initial stretch. The estimate from time $i_{0}$ to time $t_{1}$ is identical to that for $S_{k}^{\prime \prime}$, and

$$
\log \frac{\left(f_{a^{\prime}}^{i_{0}-1}\right)^{\prime}\left(x_{1}\left(a^{\prime}\right)\right)}{\left(f_{a}^{i_{0}-1}\right)^{\prime}\left(x_{1}(a)\right)} \leq K^{i_{0}}\left|a-a^{\prime}\right|
$$

where $K$ is related to the $C^{2}$ norm of the 1-parameter family. Since $\left|a-a^{\prime}\right|<K e^{-\frac{1}{4} \lambda i_{1}}$, we know that the contribution of the first $i_{0}$ iterates is $<1$ if $i_{1}$ is sufficiently large.

This completes the proof of ( $\mathrm{P} 3^{\prime}$ ).
Reference: A version of the material in Sects. 5.2 and 5.3 is contained in [BC1].

## 6 Two Parameter Estimates

### 6.1 Processes defined by critical orbits

Let $\lambda, \alpha, \varepsilon$ and $\hat{\varepsilon}$ be as in Sect. 5.2. Associated with each $\hat{x} \in C$ we now introduce the idea of a process $\left\{\gamma_{i}\right\}$ describing the dynamics of $a \mapsto \hat{x}(a)$ combined with a deletion process. The domain of definition of this process is $\Delta_{0}$, a subinterval of $(-\hat{\varepsilon}, \hat{\varepsilon})$. For each $i=0,1,2, \cdots$ and $a \in \Delta_{0}, \gamma_{i}(a)$ is equal to either $\hat{x}_{i}(a)$ or $*$, the meaning of the latter being that the parameter in question will not be considered further, i.e., it is deleted. In particular, if $\gamma_{i}(a)=*$, then $\gamma_{i+k}(a)=*$ for all $k>0$,

Roughly speaking, we seek to identify a decreasing sequence of subsets $\left\{\gamma_{i} \neq *\right\}$ of $\Delta_{0}$ and a sequence of partitions $\mathcal{Q}_{i}$ defined on $\left\{\gamma_{i} \neq *\right\}$ representing the canonical subdivision associated with $a \mapsto \hat{x}(a)$. This subdivision is carried out in a manner analogous to that in Sect. 3.1. As the reader will recall, the construction in Sect. 3.1 relies on the decomposition of orbits into bound and free periods. We have seen in Section 5 that bound/free notions are well defined for $a \mapsto \hat{x}(a)$ under suitable circumstances. The idea is that whenever we are unable to guarantee these circumstances, we will delete the parameter. Later on, we will see that it is useful also to make deletions for other purposes (but will not concern ourselves with that for now). This is the motivation for the definition below.

For $N \leq \infty$, we say $\left\{\gamma_{i}, i<N\right\}$ is a process associated with $\hat{x}$ if the following hold: Let $\Delta_{0} \subset(-\hat{\varepsilon}, \hat{\varepsilon})$ be an interval containing 0 , and let $i_{1}$ be as in Sect. 5.2.
(1) (a) We assume $\gamma_{i}\left(\Delta_{0}\right) \cap C_{\frac{1}{2} \delta_{0}}=\emptyset$ for all $i \leq i_{1}$. In this time range, $\gamma_{i}$ has no meaning. We set $\gamma_{i}=\hat{x}_{i}$ and $\mathcal{Q}_{i}=\left\{\Delta_{0}\right\}$. No deletions are permitted.
(b) A subdivision must occur before $\gamma_{i}\left(\Delta_{0}\right)$ meets $C_{\delta}$.
(2) At time $i>i_{1}$, we assume that for all $j<i, \gamma_{j}$ are defined, as are $\mathcal{Q}_{j}$, representing a canonical subdivision. We assume also that the notion of bound/free makes sense on each $\omega \in \mathcal{Q}_{j}$. Consider now $\omega \in \mathcal{Q}_{i-1}$ (on which $\gamma_{i-1} \neq *$ ). We first put on it the canonical partition $\mathcal{Q}_{i}$ as defined in Sect. 3.1. On each $\hat{\omega} \in \mathcal{Q}_{i} \mid \omega$, there are 2 options: we either let $\gamma_{i}=\hat{x}_{i}$ on all of $\hat{\omega}$, or we let it $=*$ on all of $\hat{\omega}$. The rules are as follows:
(a) We are free to set $\gamma_{i}=*$ or $\hat{x}_{i}$ on any $\hat{\omega}$ outside of $C_{\delta}$.
(b) For $\hat{\omega} \subset C_{\delta}$, the following conditions must be met if we wish to set $\left.\gamma_{i}\right|_{\hat{\omega}}=\hat{x}_{i}$ :
(i) $\hat{x}_{i}(\hat{\omega}) \cap\left\{d(\cdot, C)<e^{-\alpha i}\right\}=\emptyset$;
(ii) $\hat{\omega} \subset \Omega_{\alpha^{*} i}$, i.e., $f_{a} \in \mathcal{G}_{\alpha^{*} i}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ for all $a \in \hat{\omega}$.

Finally, we set $\gamma_{i}=*$ on $\left\{\gamma_{i-1}=*\right\}$. This completes our definition at the $i$ th step. Paragraph (2) is then repeated with $i+1$ in the place of $i$.

We observe that the process above is well defined. It is clearly well defined initially. When $\gamma_{i}(\hat{\omega}) \subset C_{\delta},(2)(\mathrm{b})$ guarantees that for all $a \in \hat{\omega}$ for which $\gamma_{i}(a) \neq *$, the ensuing bound period is meaningful (see Sect. 5.1). Once this is taken care of, Sect. 5.2 gives the desired resemblance to phase-space dynamics until the next free return.

We remark that even though $\left\{\gamma_{i}\right\}$ is associated with a particular $\hat{x} \in C$, condition (2)(b)(ii) demands good behavior of all critical orbits up to time $\alpha^{*} i$. The fact that this requirement is only up to time $\alpha^{*} i$, which is $\ll i$, is crucial for us. Observe also that on each $\Delta_{0}$, there is a maximal process, referring to the one in which the only deletions are those in (2)(b). All other processes are subordinate to this one, meaning they are defined by rules that demand further deletions.

### 6.2 Deletions due to (G1)

We fix $\hat{x} \in C$ and $\Delta_{0}$, and let $\left\{\gamma_{i}, i<n\right\}$ be a process associated with $\hat{x}$.
Lemma 6.1 There exists $K$ such that for all $\omega \in \mathcal{Q}_{n-1}$, if $\omega_{1}$ is the part of $\omega$ deleted on account of (G1) at step $n$, then

$$
\left|\omega_{1}\right|<K e^{-\frac{1}{2} \alpha n}|\omega| .
$$

Proof: Suppose $\omega_{1} \neq \emptyset$. Let $j_{0}$ be the largest $j<n$ such that (i) $\omega \in \mathcal{Q}_{j}$, (ii) $\gamma_{j}(\omega)$ is free and "long". (Such a $j_{0}$ exists by condition (1)(b) in Sect. 6.1; one may have to go back to the time when $\omega$ is created as a result of a subdivision.) There are two possibilities:

Case 1. $\gamma_{j_{0}}(\omega)$ is outside of $C_{\delta}$, i.e., $\gamma_{j_{0}}(\omega) \approx \pi$ for some $\pi$ with $\pi \cap C_{\delta}=\emptyset$. Then $\left|\gamma_{j_{0}}(\omega)\right| \geq \delta$, and $\left|\gamma_{n}(\omega)\right|>K^{-1} \delta$. Not knowing the location of $\gamma_{n}(\omega)$, we assume the worst-case scenario, namely that $\gamma_{n}(\omega)$ crosses entirely a forbidden region $\left\{d(\cdot, \hat{y})<e^{-\alpha n}\right\}$ for some $\hat{y} \in C$. Thus the fraction of $\omega$ with $d\left(\hat{x}_{n}, C\right)<e^{-\alpha n}$ is $<2 e^{-\alpha n} \cdot K \delta^{-1}$, which we
may assume is $<K e^{-\frac{1}{2} \alpha n}$ (see the paragraph following Proposition 2.2). Here ( $\mathrm{P} 3^{\prime}$ ) is used to transfer the ratio of lengths on $\gamma_{n}(\omega)$ back to $\omega$.

Case 2. $\gamma_{j_{0}}(\omega) \approx I_{\mu j}$. Let $p$ be the bound period initiated at time $j_{0}$. Observe first that $\left|\gamma_{n}(\omega)\right|>K^{-1}\left|\gamma_{j_{0}+p}(\omega)\right|>K^{-1} e^{-\frac{1}{10}|\mu|}$ : Since $\gamma_{n}(\omega)$ is free (otherwise there would be no deletion), $n \geq j_{0}+p$. The first inequality follows from ( $\mathrm{P} 1^{\prime}$ )(ii) combined (possibly) with ( $\mathrm{P} 2^{\prime}$ )(c); the second follows from ( $\mathrm{P} 2^{\prime}$ )(d). Observe also that by design, $|\mu| \leq \alpha j_{0}<\alpha n$, so the fraction of $\omega$ being estimated is again $<K e^{-\alpha n} e^{\frac{1}{10} \alpha n}<K e^{-\frac{1}{2} \alpha n}$.

### 6.3 Deletions on account of (G2)

We begin with an estimate on derivative growth in terms of the time an orbit spends in bound periods initiated at returns to $C_{\hat{\delta}}$ for arbitrary $\hat{\delta}<\delta$. Consider $f \in \mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ and $n \leq \frac{1}{\alpha^{*}} N$. Let $x \in I$ be such that $d\left(x_{i}, C\right) \geq \min \left\{\frac{1}{2} \delta_{0}, e^{-\alpha i}\right\}$ for all $0 \leq i<n$. By the reasoning in Sect. 5.1, the usual bound/free decomposition makes sense for the orbit of $x$ up to time $n$. Let $B(\hat{\delta} ; n)$ denote the total number of $i, 0 \leq i \leq n$, such that $x_{i} \in C_{\hat{\delta}}$ or it is in a bound period initiated from a visit to $C_{\hat{\delta}}$.

Lemma 6.2 Let $f$ and $x$ be as above. Given $\hat{\delta} \leq \delta$ and $\sigma>0$, if $B(\hat{\delta} ; n) \leq \sigma n$, then

$$
\left|\left(f^{n}\right)^{\prime}(x)\right|>K^{-1} \hat{\delta} e^{\left[(1-\sigma) \frac{1}{4} \lambda_{0}-\alpha\right] n}
$$

Proof: Consider first the case where $x_{n}$ is free. Let $\hat{t}_{1}<\hat{t}_{1}+\hat{p}_{1} \leq \hat{t}_{2}<\hat{t}_{2}+\hat{p}_{2} \leq \cdots \leq$ $\hat{t}_{k}+\hat{p}_{k} \leq n$ where $\hat{t}_{1}, \cdots, \hat{t}_{k}$ are the consecutive free return times to $\{d(\cdot, C)<\hat{\delta}\}$. Then

$$
\left(f^{n}\right)^{\prime}(x)=\left(f^{n-\hat{t}_{k}-\hat{p}_{k}}\right)^{\prime}\left(x_{\hat{t}_{k}+\hat{p}_{k}}\right) \cdot\left(f^{\hat{p}_{k}}\right)^{\prime}\left(x_{\hat{t}_{k}}\right) \cdot\left(f^{\hat{t}_{k}-\hat{t}_{k-1}-\hat{p}_{k-1}}\right)^{\prime}\left(x_{\hat{t}_{k-1}+\hat{p}_{k-1}}\right) \cdots\left(f^{\hat{t}_{1}}\right)^{\prime}(x)
$$

We use (P1)(i) for $\left(f^{n-\hat{t}_{k}-\hat{p}_{k}}\right)^{\prime}\left(x_{\hat{t}_{k}+\hat{p}_{k}}\right)$, (P1)(ii) for $\left(f^{t_{i}-\hat{t}_{i-1}-\hat{p}_{i-1}}\right)^{\prime}\left(x_{\hat{t}_{i-1}+\hat{p}_{i-1}}\right)$, and the trivial estimate $\left|\left(f^{\hat{p}_{i}}\right)^{\prime}\left(x_{\hat{t}_{i}}\right)\right|>c_{1}^{-1}$ for growth during bound periods (see (P2)(ii)). This gives $\left|\left(f^{n}\right)^{\prime}(x)\right|>K^{-1} \hat{\delta} e^{\frac{1}{4} \lambda_{0}(1-\sigma) n}$ since $\hat{p}_{1}+\cdots+\hat{p}_{k} \leq \sigma n$ by assumption. The factor $-\alpha n$ is needed if $n$ is not free; see Corollary 2.1.

Corollary 6.1 Let the hypotheses be as in Lemma 6.2, with $x=\hat{x}_{1}$ for some $\hat{x} \in C$. We assume further that $d\left(\hat{x}_{i}, C\right)>\frac{1}{2} \delta_{0}$ for all $i \leq n_{0}$ where $n_{0}$ is sufficiently large depending on $\hat{\delta}$. Then $B(\hat{\delta} ; 0, n)<\sigma n$ implies $\left|\left(f^{n}\right)^{\prime}\left(\hat{x}_{1}\right)\right|>c_{1} e^{\left[(1-\sigma) \frac{1}{4} \lambda_{0}-\alpha\right] n}$.

Proof: The factor $\hat{\delta}$ is absorbed into the initial growth if $n>n_{0}$.
For $f \in \mathcal{G}$, it can be deduced from properties of the invariant measure that $\int \frac{1}{n} B(\hat{\delta} ; n) d \mu$ decreases with $\hat{\delta}$. In light of the duality in Sect. 5.2 , one may expect a similar phenomenon for $a \mapsto \gamma_{i}(a)$. We formulate below a large deviation estimate useful for estimating the measure of parameters deleted on account of (G2).

Let $\left\{\gamma_{i}, i<n\right\}$ be as in Sect. 6.2. For $a$ such that $\gamma_{n}(a) \neq *$, let $B(a, \hat{\delta} ; n)$ be the number $B(\hat{\delta} ; n)$ defined above with $f=f_{a}$ and $x=\hat{x}$.

Proposition 6.1 Given any $\sigma>0$, there exist positive numbers $\hat{\varepsilon}_{1}=\hat{\varepsilon}_{1}(\sigma)$ and $\hat{\delta}=\hat{\delta}(\sigma)$ such that

$$
\mid\left\{a \in \Delta_{0}: \gamma_{n} \neq * \text { and } B(a, \hat{\delta} ; n)>\sigma n\right\}\left|<e^{-\hat{\varepsilon}_{1} n}\right| \Delta_{0} \mid
$$

### 6.4 Large deviation estimate

We first state the analog of Lemma 3.3. Let $\hat{\omega} \in \mathcal{Q}_{j_{0}}$ be such that $\gamma_{j_{0}}(\hat{\omega}) \neq *$ and is free. On $\hat{\omega}$ we define $\hat{S}$, a stopping time starting from $j_{0}$, as follows: We extend the process on $\omega$ beyond time $j_{0}$, and for each $a \in \omega$, let $k=k(a)>j_{0}$ be the first time when $\gamma_{k}\left(\mathcal{Q}_{k-1}(a)\right)$ is not in a bound period and has length $>\delta$. If such a $k$ exists, we set $\hat{S}(a)=k-j_{0}$. If $a$ is deleted before that happens, we set $\hat{S}(a)=0$.

Lemma 6.3 Let $\hat{\omega} \in \mathcal{Q}_{j_{0}}$ be such that $\gamma_{j_{0}}(\hat{\omega})$ is free and $\approx I_{\mu j}$. Then

$$
|\{a \in \hat{\omega}: \hat{S}(a)>m\}|<e^{-\frac{1}{2} K^{-1} m}|\hat{\omega}| \quad \text { for all } \quad m>K \log |\mu| .
$$

The proof is entirely parallel to that of Lemma 3.3 in Sect. 3.1.
Proof of Proposition 6.1: We take a probabilistic viewpoint, with underlying probability space $\left(\Delta_{0}, P\right), P$ being normalized Lebesgue measure on $\Delta_{0}$. Let $\hat{\delta}>0$ be a small number to be determined. Let $n$ be fixed. The idea is to introduce $X_{i}$ dominated by certain exponential random variables such that $B(a):=B(a, \hat{\delta}, n) \leq \sum X_{i}(a)$.

## Step I. Formulation of problem as one involving $\sum X_{i}$

For each $a \in \Delta_{0}$, we define $t_{0}<t_{1}<\cdots$ and $S_{1}, S_{2}, \cdots$ via the following algorithm, with the understanding that the algorithm terminates as soon as $\gamma_{i}(a)=*$ or time $n$ is reached. To get started, let $t_{0}$ be the smallest $j>0$ such that $\gamma_{j}\left(\mathcal{Q}_{j-1}(a)\right) \cap C_{\hat{\delta}} \neq \emptyset$.
(i) After $t_{i}$ is defined, we define $S_{i+1}$ : If $\mathcal{Q}_{t_{i}}(a) \cap C_{\hat{\delta}}=\emptyset$, set $S_{i+1}=0$; if $\mathcal{Q}_{t_{i}}(a) \subset C_{\hat{\delta}}$, let $S_{i+1}=\min \left(\hat{S}, n-t_{i}\right)$ where $\hat{S}$ is the stopping time above starting from $t_{i}$.
(ii) If $S_{i+1}=0$, let $t_{i+1}$ be the smallest $j>t_{i}$ such that $\gamma_{j}\left(\mathcal{Q}_{j-1}(a)\right) \cap C_{\hat{\delta}} \neq \emptyset$; define $t_{i+1}$ the same way if $S_{i+1}>0$ except that $j$ is taken $\geq t_{i}+S_{i+1}$.
Suppose $t_{i}(a)$ is defined. Let $Q=\mathcal{Q}_{t_{i}-1}(a)$. Assuming $\hat{\delta} \ll \delta$, we claim:
(1) $\gamma_{t_{i}}(Q)$ is free;
(2) $\left|\gamma_{t_{i}}(Q)\right|>\hat{\delta}^{\frac{1}{10}}$;
(3) for all $a^{\prime}, a^{\prime \prime} \in Q, \tau_{t_{i}}\left(a^{\prime}\right) / \tau_{t_{i}}\left(a^{\prime \prime}\right)<K$.
(1) is true because trajectories of critical curves in bound periods initiated outside of $C_{\hat{\delta}}$ cannot meet $C_{\hat{\delta}}$. If $S_{i}>0$, it may happen that $t_{i}=t_{i-1}+S_{i}$, in which case $\left|\gamma_{t_{i}}(Q)\right|>\delta$ by definition. Otherwise we back up to time $t$ when $Q$ was first created as an element of some $\mathcal{Q}_{j}$. Then $t_{i-1}+S_{i} \leq t<t_{i}$, and $\gamma_{t}(Q) \cap C_{\hat{\delta}}=\emptyset$. If $\gamma_{t}(Q)$ is outside of $C_{\delta}$, then $\left|\gamma_{t_{i}}(Q)\right|>K^{-1} \delta$ by $\left(\mathrm{P} 1^{\prime}\right)$. If $\gamma_{t}(Q) \approx I_{\mu j}$ for some $I_{\mu j} \subset C_{\delta} \backslash C_{\hat{\delta}}$, then $\left|\gamma_{t_{i}}(Q)\right|>K^{-1} \hat{\delta}^{\frac{8 \alpha}{\lambda}}$ by ( $\mathrm{P} 2^{\prime}$ )(d). In all cases, (2) holds assuming $\hat{\delta} \ll \delta$. (3) follows from ( $\mathrm{P} 3^{\prime}$ ).

Because of (1)-(3), we think of $t_{i}$ as times of dynamical renewal.
In preparation for Step II, we organize some of the information from above as follows. Let $X_{0}=0$. For $i=1,2, \cdots, n$, let $X_{i}: \Delta_{0} \rightarrow \mathbb{Z}^{+}$be such that $X_{i}(a)=S_{i}(a)$ where $S_{i}(a)$ is defined, 0 otherwise. Then $B \leq \sum_{i \leq n} X_{i}$. It suffices to show $P\left\{\sum_{i \leq n} X_{i}>\sigma n\right\}$ decreases exponentially with $n$. We define the following $\sigma$-algebras on $\omega$ : Let $\bar{A}_{i}$ be the set of $a$ for which $t_{i}$ is defined. Then $A_{i} \in \mathcal{F}_{i}$, and for $a \in A_{i}$, the atom of $\mathcal{F}_{i}$ containing $a$ is $\mathcal{Q}_{t_{i}(a)-1}(a)$. For $a \notin A_{i}$, the atom of $\mathcal{F}_{i}$ containing $a$ is $\mathcal{Q}_{k}(a)$ where $k$ is the last step
before the algorithm above is terminated. One verifies that $\mathcal{F}_{i}$ so defined is a $\sigma$-algebra, that $\mathcal{F}_{0}<\mathcal{F}_{1}<\cdots<\mathcal{F}_{n}$, and that $X_{i}$ is measurable with respect to $\mathcal{F}_{i}$.

Step II. Large deviation estimate for $\sum_{1 \leq i \leq n} X_{i}$
First we compute the conditional distribution of $X_{i+1}$ given $\mathcal{F}_{i}, i \geq 0$. Consider $Q \in$ $\left.\mathcal{F}_{i}\right|_{A_{i}} .\left(\right.$ On $Q \in \mathcal{F}_{i}$ with $Q \cap A_{i}=\emptyset, X_{i+1}=0$.) From (2) and (3) above, we have
(i) $P\left(X_{i+1}=0 \mid Q\right) \geq 1-K \hat{\delta} \frac{9}{10}$.

For $I_{\mu j} \subset C_{\hat{\delta}}$, Lemma 6.3 together with (1) and (3) above give
(ii) $P\left(X_{i+1}>m \mid Q \cap\left\{\gamma_{t_{i}} \in I_{\mu j}\right\}\right)<K e^{-\frac{1}{2} K^{-1} m}$ if $m \geq K|\mu|$; no information otherwise.

Combining the last two estimates, we obtain for all $m \geq 0$,

$$
\begin{equation*}
P\left(X_{i+1}>m \mid Q\right)<K \hat{\delta}^{-\frac{1}{10}} \min \left(\hat{\delta}, e^{-K^{-1} m}\right)+K \hat{\delta}^{\frac{9}{10}} e^{-\frac{1}{2} K^{-1} m} \tag{8}
\end{equation*}
$$

A simple computation then gives $E\left[e^{\rho X_{i+1}} \mid Q\right]<\infty$ if $\rho<\frac{1}{2} K^{-1}$ (where $K$ is as in the exponents above). We note further that by decreasing $\hat{\delta}$ (keeping $\rho$ fixed), $E\left[e^{\rho X_{i+1}} \mid Q\right]$ can be made arbitrarily close to 1 . Let $\eta>0$ be a number to be determined shortly, and choose $\hat{\delta}=\hat{\delta}(\eta)$ sufficiently small that $E\left[e^{\rho X_{i+1}} \mid Q\right]<e^{\eta}$. Observing that the upper bound in (8) and hence that for $E\left[e^{\rho X_{i+1}} \mid Q\right]$ do not depend on $i$ or on $Q$, we conclude that with the choices of $\rho, \eta$ and $\hat{\delta}$ above, $E\left[e^{\rho X_{i+1}} \mid \mathcal{F}_{i}\right]<e^{\eta}$ for every $i \geq 0$.

To finish, we observe that

$$
\begin{aligned}
E\left[e^{\rho \sum_{1 \leq i \leq n} X_{i}}\right] & =E\left[E\left[e^{\rho \sum_{1 \leq i \leq n} X_{i}} \mid \mathcal{F}_{n-1}\right]\right]=E\left[e^{\rho \sum_{1 \leq i \leq n-1} X_{i}} E\left[e^{\rho X_{n}} \mid \mathcal{F}_{n-1}\right]\right] \\
& \leq e^{\eta} E\left[e^{\rho \sum_{1 \leq i \leq n-1} X_{i}}\right]
\end{aligned}
$$

giving inductively $E\left[e^{\rho \sum_{1 \leq i \leq n} X_{i}}\right] \leq e^{n \eta}$. We arrive, therefore, at the estimate

$$
P\{B>\sigma n\}<P\left\{\sum_{1 \leq i \leq n} X_{i}>\sigma n\right\}<e^{\eta n-\rho \sigma n}
$$

This is $<e^{-\frac{1}{2} \rho \sigma n}$ if $\eta$ is chosen $<\frac{1}{4} \rho \sigma$.
References: A version of Sects. 6.2 and 6.3 is used in [BC2]; Sect. 6.4 is taken from [WY2].

## $7 \quad$ Positive measure sets of good parameters

### 7.1 Preliminary definitions and choices

1. We fix $\lambda \leq \frac{1}{5} \lambda_{0}$.
2. Augmented versions of (G1) and (G2). For reasons to become clear, it will be advantageous to put our good maps "deeper inside" $\mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$. We say $\hat{x} \in C$ satisfies (G1)\# and (G2) ${ }^{\#}$ up to time $N$ if for all $1 \leq i \leq N$,

$$
(\mathbf{G} 1)^{\#} d\left(\hat{x}_{i}, C\right)>\min \left(\frac{1}{2} \delta_{0}, 2 e^{-\alpha i}\right)
$$

(G2) ${ }^{\#}\left|\left(f^{i}\right)^{\prime}\left(\hat{x}_{1}\right)\right|>2 c_{1} e^{\lambda_{1} i}$ where $\lambda_{1}=\lambda+\frac{1}{100} \lambda_{0}$,
and say $f \in \mathcal{G}_{N}^{\#}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ if all $\hat{x} \in C$ satisfy (G1) ${ }^{\#}$ and (G2) ${ }^{\#}$ up to time $N$. Clearly, $\mathcal{G}_{N}^{\#}\left(f_{0} ; \lambda, \alpha, \varepsilon\right) \subset \mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$. The proof of the following lemma is straightforward.

Lemma 7.1 There exists $K_{4}>1$ for which the following holds: If $f_{\hat{a}} \in \mathcal{G}_{N}^{\#}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$, then for all $n \leq N, f_{a} \in \mathcal{G}_{n}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ for all $a \in\left[\hat{a}-K_{4}^{-n}, \hat{a}+K_{4}^{-n}\right]$.
3. Choice of $\alpha$. In addition to $\alpha<\frac{1}{100} \lambda$, we impose two upper bounds on $\alpha$ : The first is introduced in (5) in Sect. 5.2; the second is (9) in Sect. 7.2. With $\lambda$ and $\alpha$ fixed, $\mathcal{G}_{N}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ and $\mathcal{G}_{N}^{\#}\left(f_{0} ; \lambda, \alpha, \varepsilon\right)$ will be abbreviated as $\mathcal{G}_{N}$ and $\mathcal{G}_{N}^{\#}$ from here on.
4. Choices of $\sigma$ and $\hat{\delta}$. We need $\sigma$ to be small enough that the exponent in Corollary 6.1, namely $(1-\sigma) \frac{1}{4} \lambda_{0}-\alpha$, is $>\lambda_{1}$. (For example, $\sigma=\frac{1}{100}$ will do.) We then let $\hat{\delta}$ be given by Proposition 6.1 with $\frac{1}{2} \sigma$ in the place of $\sigma$.
5 The start-up interval $\Delta_{0}$. We choose $\Delta_{0} \subset(-\hat{\varepsilon}, \hat{\varepsilon})$ to contain 0 and to be short enough that for some $n_{0}$ sufficiently large, $d\left(\hat{x}_{i}, C\right)>\frac{1}{2} \delta_{0}$ for all $i \leq n_{0}, \hat{x} \in C$ and $a \in \Delta_{0}$. A number of impositions on $n_{0}$ have been made; see, for example, Sects. 2.1 and 5.2, and Corollary 6.1. There will be more in the next two pages.

### 7.2 Inductive construction of $\Delta$

We seek to construct a sequence of sets $\Delta_{0} \supset \Delta_{n_{0}} \supset \Delta_{2 n_{0}} \supset \Delta_{2^{2} n_{0}} \supset \cdots$ in parameter space with the properties that
(i) for each $\ell,\left\{f_{a}, a \in \Delta_{2^{\ell} n_{0}}\right\} \subset \mathcal{G}_{2^{\ell} n_{0}}^{\#}$ and
(ii) $\Delta:=\cap_{\ell \geq 0} \Delta_{2} n_{n}$ has positive Lebesgue measure.

The rules of construction are detailed below; the measure estimate is given in Sect. 7.3.

## Overview of procedure

Let $C:=\left\{\hat{x}^{1}, \hat{x}^{2}, \cdots, \hat{x}^{q}\right\}$. Associated with each $\hat{x}^{k}$, we define a process $\left\{\gamma_{i}^{k}, i<\infty\right\}$ in the sense of Sect. 6.1 with the property that for every $a$ such that $\gamma_{2^{\ell} n_{0}}^{k}(a) \neq *, \hat{x}^{k}(a)$ satisfies (G1) ${ }^{\#}$ and (G2) ${ }^{\#}$ up to time $2^{\ell} n_{0}$. We then let

$$
\Delta_{2^{\ell} n_{0}}^{k}:=\left\{\gamma_{2^{\ell} n_{0}}^{k} \neq *\right\} \quad \text { and } \quad \Delta_{2^{\ell} n_{0}}:=\cap_{1 \leq k \leq q} \Delta_{2^{\ell} n_{0}}^{k} .
$$

It follows that $f_{a} \in \mathcal{G}_{2^{\ell}{ }_{n 0}}^{\#}$ for every $a \in \Delta_{2^{\ell}{ }_{n 0}}$.
The processes $\gamma_{i}^{k}$ are updated in $N$-to- $2 N$ cycles, $N=2^{\ell} n_{0}, \ell=1,2, \cdots$. Within each cycle, we first update each of the $q$ processes individually, i.e., extend $\gamma_{i}^{k}$ from $i=N$ to $i=2 N$. At the end of this updating, we reset some of the values of $\gamma_{2 N}^{k}$ to $*$ to reflect the combined status of all $q$ processes before moving to the next cycle.

Remarks. It is absolutely essential to take inventory of the global picture at regular time intervals (as we do at times $2^{\ell} n_{0}$ ). Other than that, the precise order in which $\gamma_{i}^{k}$ is updated is unimportant. Also, the number " 2 " has little significance: all that is needed is a relation with $\alpha$ that gives (9) below.

Getting started: Let $n_{1}^{k}$ be the smallest $i>0$ such that $\hat{x}_{i}^{k}\left(\Delta_{0}\right) \cap C_{\frac{1}{2} \delta_{0}} \neq \emptyset$. Then $n_{1}^{k}>n_{0}$, and $\left|\hat{x}_{n_{1}^{k}}^{k}\left(\Delta_{0}\right)\right|>\frac{1}{2} \delta_{0}$. Since $\delta \ll \frac{1}{2} \delta_{0}$, the first subdivision occurs at or before time $n_{1}^{k}$. Condition (1)(b) in the definition of a process in Sect. 6.1 is met. Recall that we have assumed $e^{-\alpha n_{0}} \ll \delta$, so that Lemma 6.1 applies to all deletions due to (G1)\#.

Formal procedure from step $N=2^{\ell} n_{0}$ to step $2 N$
At time $N$, we are handed $q$ well defined processes $\left\{\gamma_{i}^{k}, i \leq N\right\}, k=1,2, \cdots, q$, with the property that for each $a \in \Delta_{N}^{k}=\left\{\gamma_{N}^{k} \neq *\right\}$,
(i) $f_{a} \in \mathcal{G}_{2 \alpha^{*} N}$, and
(ii) $\hat{x}^{k}(a)$ satisfies (G1) ${ }^{\#}$ and (G2) ${ }^{\#}$ up to time $N$.

The procedure from $N$ to $2 N$ consists of the following:
Step 1. The following is carried out for each of the processes $\left\{\gamma_{i}^{k}\right\}$.
1a. We extend $\gamma_{i}^{k}$ up to $i=2 N$, deleting all $\omega \in \mathcal{Q}_{i}$ with $\gamma_{i}(\omega) \cap\left\{d(\cdot, C)<2 e^{-\alpha i}\right\} \neq \emptyset$. (Deletions corresponding to (2)(b)(ii) in Sect. 6.1 are not needed because (i) above already puts $\left\{\gamma_{N}^{k} \neq *\right\} \subset \Omega_{2 \alpha^{*} N}$.) We denote the resulting $\gamma_{2 N}^{k}$ by $\gamma_{2 N, \text { tmp } 1}^{k}$ as its values will be reset momentarily.
1b. On each $\omega \in \mathcal{Q}_{2 N}^{k}$, we let $\gamma_{2 N, \text { tmp } 2}^{k}=*$ if $B(a, \hat{\delta} ; 2 N)>\sigma N$ for $a \in \omega$, and let it be $=\gamma_{2 N, \text { tmp1 }}^{k}$ otherwise.
Step 2. On each $\omega \in \mathcal{Q}_{2 N}^{k}$, we let $\gamma_{2 N}^{k}=\gamma_{2 N, \text { tmp2 }}^{k}$ if $\omega \cap \Delta_{N} \neq \emptyset$, otherwise let it $=*$.
It remains to show that these steps lead to (i) and (ii) above with $N$ replaced by $2 N$. For $a \in \Delta_{2 N}^{k}=\left\{\gamma_{2 N}^{k} \neq *\right\}$, Step 1a ensures that $\hat{x}_{i}^{k}(a)$ satisfies (G1) ${ }^{\#}$ up to time $2 N$. By Step 1b, $B(a, \hat{\delta} ; n) \leq \sigma N<\sigma n$. Our choice of $\sigma$ (see Sect. 7.1) together with Corollary 6.1 then gives the lower bound for $\left|\left(f^{n-1}\right)^{\prime}\left(\hat{x}_{1}^{k}\right)\right|$ in (G2) ${ }^{\#}$. It remains to prove $f_{a} \in \mathcal{G}_{4 \alpha^{*} N}$ for $a \in \Delta_{2 N}^{k}$. Let $\omega$ be the element of $\mathcal{Q}_{2 N}^{k}$ containing $a$. Since $\omega$ survived the deletion in Step 2, there must be a point $\hat{a} \in \omega \cap \Delta_{N}$. By Lemma 7.1, it suffices to show that $\omega \subset\left[\hat{a}-K_{4}^{-4 \alpha^{*} N}, \hat{a}+K_{4}^{-4 \alpha^{*} N}\right]$. By Proposition $4.2^{\prime},|\omega|<2\left(\hat{c} c_{1}\right)^{-1} e^{-2 \lambda_{1} N}$. (Observe that the hypotheses of Proposition 4.2' are met: $\omega \subset \Omega_{2 \alpha^{*} N}$, and $\hat{x}_{i}^{k}$ obeys (G1) up to time $2 N$.) We conclude that $\omega \subset \Omega_{4 \alpha^{*} N}$ if

$$
\begin{equation*}
2\left(\hat{c} c_{1}\right)^{-1} e^{-2 \lambda_{1} N}<K_{4}^{-4 \alpha^{*} N} . \tag{9}
\end{equation*}
$$

This is one of the conditions imposed on $\alpha$ in item 3 of Sect. 7.1.

### 7.3 Lower bound on measure of $\Delta$

For each $N$ and $k$, we now estimate the contribution to $\Delta_{N} \backslash \Delta_{2 N}\left(\operatorname{not} \Delta_{N}^{k} \backslash \Delta_{2 N}^{k}!\right.$ ) due to deletions in $\gamma_{i}^{k}$ from $i=N$ to $i=2 N$.

Step 1a: By Lemma 6.1, the total measure deleted is $<\sum_{N<i \leq 2 N} K e^{-\frac{1}{2} \alpha i}\left|\Delta_{0}\right|$.
Step 1b: By Proposition 6.1, the total measure deleted is $<e^{-\hat{\varepsilon}_{1} N}\left|\Delta_{0}\right|$.
Step 2: The sets $\omega$ removed in this step do not meet $\Delta_{N}$, i.e., they have been deleted
earlier. Thus the actions taken in this step do not contribute to $\Delta_{N} \backslash \Delta_{2 N}$. Summing over the $q$ critical points, we obtain

$$
\left|\Delta_{N} \backslash \Delta_{2 N}\right| \leq q\left(K \sum_{N<i \leq 2 N} e^{-\frac{1}{2} \alpha i}+e^{-\hat{\varepsilon}_{1} N}\right)\left|\Delta_{0}\right|
$$

It follows that

$$
|\Delta| \geq\left|\Delta_{0}\right|-\sum\left|\Delta_{2^{\ell} n_{0}} \backslash \Delta_{2^{\ell+1} n_{0}}\right|>\left(1-q \sum_{i=n_{0}}^{\infty}\left(K e^{-\frac{1}{2} \alpha i}+e^{-\hat{\varepsilon} i}\right)\right)\left|\Delta_{0}\right|
$$

which is positive if $n_{0}$ is sufficiently large. This completes the proof of Theorem 2.
Reference: This section follows [WY2], which, together with its precursor [WY1] contain the corresponding material for rank one attractors.

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[^1]:    ${ }^{3}$ We have elected to replace this condition by an explicit description of the dynamics in Definition 1.1 because (1) that is exactly what is used and (2) we have found that maps that arise in applications often do not have negative Schwarzian derivative.

[^2]:    ${ }^{4}$ In particular, if $\pi$ is one of the outermost $I_{\mu j}$ in $C_{\delta}$, then $\pi^{+}$contains an interval of length $\delta$.

