# EXTREME VALUES FOR BENEDICKS-CARLESON QUADRATIC MAPS 

JORGE MILHAZES FREITAS AND ANA CRISTINA MOREIRA FREITAS


#### Abstract

We consider the quadratic family of maps given by $f_{a}(x)=1-a x^{2}$ with $x \in$ $[-1,1]$, where $a$ is a Benedicks-Carleson parameter. For each of these chaotic dynamical systems we study the extreme value distribution of the stationary stochastic processes $X_{0}, X_{1}, \ldots$, given by $X_{n}=f_{a}^{n}$, for every integer $n \geq 0$, where each random variable $X_{n}$ is distributed according to the unique absolutely continuous, invariant probability of $f_{a}$. Using techniques developed by Benedicks and Carleson, we show that the limiting distribution of $M_{n}=\max \left\{X_{0}, \ldots, X_{n-1}\right\}$ is the same as that which would apply if the sequence $X_{0}, X_{1}, \ldots$ was independent and identically distributed. This result allows us to conclude that the asymptotic distribution of $M_{n}$ is of Type III (Weibull).


## 1. Introduction

The dynamical systems we consider in this work are the quadratic maps given by $f_{a}(x)=$ $1-a x^{2}$ on $I=[-1,1]$, with $a \in \mathcal{B} C$, where $\mathcal{B C}$ is the Benedicks-Carleson parameter set introduced in [BC85]. The set $\mathcal{B C}$ has positive Lebesgue measure and is built in such a way that, for every $a \in \mathcal{B} C$, the Collet-Eckmann condition holds: there is exponential growth of the derivative of $f_{a}$ along the critical orbit, i.e., there is $c>0$ such that

$$
\left|\left(f_{a}^{n}\right)^{\prime}\left(f_{a}(0)\right)\right| \geq \mathrm{e}^{c n}
$$

for all $n \in \mathbb{N}$. This property guarantees not only the non-existence of an attracting periodic orbit but also the existence of an ergodic $f_{a}$-invariant probability measure $\mu_{a}$ that is absolutely continuous with respect to Lebesgue measure on $[-1,1]$. In fact, these Benedicks-Carleson systems are chaotic and highly sensitive on initial conditions. Actually, after some iterates, the behavior of most orbits becomes erratic and distributed on the set $[-1,1]$ according to the invariant measure $\mu_{a}$. Hence, it is meaningful to study the statistical properties of the orbits of these systems and, here, we are particularly concerned with their extreme type behavior.

A natural way to build a stationary stochastic process associated to $f_{a}$, for some $a \in \mathcal{B} C$, is to consider the random variable (r.v.) $X_{0}$ defined on the probability space ( $[-1,1], \mu_{a}$ ), taking values on $[-1,1]$ with distribution function (d.f.) $G_{a}(x)=\mu_{a}\{(-\infty, x] \cap[-1,1]\}$

[^0]and then iterate it by $f_{a}$, i.e., define
\[

$$
\begin{equation*}
X_{n}=f_{a}\left(X_{n-1}\right)=f_{a}^{n}\left(X_{0}\right), \text { for each } n \in \mathbb{N} . \tag{1.1}
\end{equation*}
$$

\]

This way, we obtain a stationary stochastic process $X_{0}, X_{1}, X_{2}, \ldots$ with common marginal d.f. given by $G_{a}(x)=\mu_{a}\left\{X_{0} \leq x\right\}$. The stationarity results from the $f_{a}$-invariance of the probability measure $\mu_{a}$ (see for example [KT66, Section 15.4]). Our goal is to study the asymptotic distribution of the partial maximum

$$
\begin{equation*}
M_{n}=\max \left\{X_{0}, X_{1}, \ldots, X_{n-1}\right\} \tag{1.2}
\end{equation*}
$$

when properly normalized.
The study of distributional properties of the higher order statistics of a sample, like the maximum of the sample is the purpose of Extreme Value Theory. Its major classical theorem asserts that there are only three types of non-degenerate asymptotic distributions for the maximum of an independent and identically distributed (i.i.d.) sample under linear normalization.

The main result of this work states that the limiting law of $M_{n}$ is the same as if $X_{0}, X_{1}, \ldots$ were independent with the same d.f. $G_{a}$. In fact, we verify that under appropriate normalization the asymptotic distribution of $M_{n}$ is of type III (Weibull). As usual, we denote by $G_{a}^{-1}$ the generalized inverse of the d.f. $G_{a}$, which is to say that $G_{a}^{-1}(y):=\inf \left\{x: G_{a}(x) \geq y\right\}$.

Theorem A. For each $a \in \mathcal{B} C$ and every stationary stochastic process $\left(X_{i}\right)_{i \in \mathbb{N}_{0}}$ given by (1.1), we have:

$$
P\left\{a_{n}\left(M_{n}-1\right) \leq x\right\} \rightarrow H(x)= \begin{cases}e^{-(-x)^{1 / 2}} & , x \leq 0 \\ 1 & , x>0\end{cases}
$$

where $a_{n}=\left(1-G_{a}^{-1}\left(1-\frac{1}{n}\right)\right)^{-1}$.
Haiman Hai03] has obtained a similar asymptotic result for the natural stochastic process associated with the tent map. The arguments used here for quadratic maps would allow us to obtain a different proof of Haiman's Theorem.

Also, a study concerning extremes for dynamical systems, essentially focusing the finite sample behavior of maxima, has already been done by Balakrishnan, Nicolis and Nicolis in BNN95.

## 2. Motivation and Strategy

The study of the limit behavior for maxima of a stationary process can be reduced, under adequate conditions on the dependence structure, to the Classical Extreme Value Theory for sequences of i.i.d. random variables. Hence, to the stationary process $X_{0}, X_{1}, \ldots$ we associate an independent sequence of r.v. denoted by $Z_{0}, Z_{1}, \ldots$ with common d.f. given by $G_{a}(x)=P\left\{X_{0} \leq x\right\}=\mu_{a}\left\{X_{0} \leq x\right\}$. We also set for each $n \in \mathbb{N}$

$$
\begin{equation*}
\hat{M}_{n}=\max \left\{Z_{0}, \ldots, Z_{n-1}\right\} \tag{2.1}
\end{equation*}
$$

Let us focus on the conditions that allow us to relate the asymptotic distribution of $M_{n}$ with that of $\hat{M}_{n}$. Following [LR83] we refer to these conditions as $D\left(u_{n}\right)$ and $D^{\prime}\left(u_{n}\right)$, where $u_{n}$ is a suitable sequence of thresholds converging to $\max _{z \in[-1,1]} X_{0}(z)=1$, as $n$ goes to $\infty$, that will be defined below. $D\left(u_{n}\right)$ imposes a certain type of distributional mixing property. Essentially, it says that the dependence between some special type of events fades away as they become more and more apart in the time line. $D^{\prime}\left(u_{n}\right)$ restricts the appearance of clusters, that is, it makes the occurrence of consecutive 'exceedances' of the level $u_{n}$ an unlikely event.

Consider a sequence of stationary r.v. $Y_{1}, Y_{2} \ldots$ with common d.f. $F$. We say that an exceedance of the level $u_{n}$ occurs at time $i$ if $Y_{i}>u_{n}$. The probability of such an exceedance is $1-F\left(u_{n}\right)$ and so the mean value of the number of exceedances occurring up to $n$ is $n\left(1-F\left(u_{n}\right)\right)$. The sequences of levels $u_{n}$ we consider are such that $n\left(1-F\left(u_{n}\right)\right) \rightarrow \tau$ as $n \rightarrow \infty$, for some $\tau \geq 0$, which means that, in a time period of length $n$, the expected number of exceedances is approximately $\tau$.

Observe that, in the case of our stationary process $X_{0}, X_{1}, \ldots$, we have $X_{j}>u_{n}$ if and only if $\left|X_{j-1}\right|<\sqrt{\frac{1-u_{n}}{a}}$. Since $u_{n} \rightarrow 1$ as $n \rightarrow \infty$, then, for large $n$, an exceedance of the level $u_{n}$ is always preceded by a return to a tight vicinity of the critical point. Hence, as one would expect, there is an intimate relation between exceedances and deep returns to the vicinity of the critical point.
$D\left(u_{n}\right)$ is a type of mixing requirement specially adapted to Extreme Value Theory. In this context, the events of interest are those of the form $\left\{X_{i} \leq u\right\}$ and their intersections. Observe that $\left\{M_{n} \leq u\right\}$ is just $\left\{X_{0} \leq u, \ldots, X_{n-1} \leq u\right\}$. A natural mixing condition in this context is the following. Let $G_{i_{1}, \ldots, i_{n}}\left(x_{1}, \ldots, x_{n}\right)$ denote the joint d.f. of $X_{i_{1}}, \ldots, X_{i_{n}}$ and set $G_{i_{1}, \ldots, i_{n}}(u)=G_{i_{1}, \ldots, i_{n}}(u, \ldots, u)$.

Condition $\left(D\left(u_{n}\right)\right)$. We say that $D\left(u_{n}\right)$ holds for the sequence $X_{0}, X_{1}, X_{2}, \ldots$ if for any integers $i_{1}<\ldots<i_{p}$ and $j_{1}<\ldots<j_{k}$ for which $j_{1}-i_{p}>m$, and any large $n \in \mathbb{N}$,

$$
\left|G_{i_{1}, \ldots, i_{p}, j_{1}, \ldots, j_{k}}\left(u_{n}\right)-G_{i_{1}, \ldots, i_{p}}\left(u_{n}\right) G_{j_{1}, \ldots, j_{k}}\left(u_{n}\right)\right| \leq \gamma(m),
$$

where $\gamma(m) \rightarrow 0$ as $m \rightarrow \infty$.
We remark that the actual definition of $D\left(u_{n}\right)$ appearing in [LLR83, Section 3.2] is a weaker requirement but the one considered here suits our purposes.

Condition $\left(D^{\prime}\left(u_{n}\right)\right)$. We say that $D^{\prime}\left(u_{n}\right)$ holds for the sequence $X_{0}, X_{1}, X_{2}, \ldots$ if

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \limsup _{n \rightarrow \infty} n \sum_{j=1}^{[n / k]} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\}=0 \tag{2.2}
\end{equation*}
$$

The sequence $u_{n}$ is such that the average number of exceedances in the time interval $\{0, \ldots,[n / k]\}$ is approximately $\tau / k$, which goes to zero as $k \rightarrow \infty$. However, the exceedances may have a tendency to be concentrated in the time period following the first exceedance at time 0 . Condition 2.2 prevents this from happening, i.e., forbids the concentration of exceedances by bounding the probability of more than one exceedance in the
time interval $\{0, \ldots,[n / k]\}$. This guarantees that the exceedances should appear scattered through the time period $\{0, \ldots, n-1\}$.

The special relevance of both these conditions is connected to Theorem 3.5.2 of [LLR83]: let $a_{n}$ and $b_{n}$ be sequences such that $P\left\{a_{n}\left(\hat{M}_{n}-b_{n}\right) \leq x\right\} \rightarrow H(x)$ for some non-degenerate d.f. $H$; if $D\left(u_{n}\right), D^{\prime}\left(u_{n}\right)$ hold for the stationary sequence $X_{1}, X_{2}, \ldots$, when $u_{n}=x / a_{n}+b_{n}$ for each $x$, then $P\left\{a_{n}\left(M_{n}-b_{n}\right) \leq x\right\} \rightarrow H(x)$. This means that if we are able to show that conditions $D\left(u_{n}\right)$ and $D^{\prime}\left(u_{n}\right)$ are valid for the stationary process $X_{0}, X_{1}, \ldots$, then $M_{n}$ and $\hat{M}_{n}$ share the same asymptotic distribution with the same normalizing sequences. Consequently, our strategy to prove Theorem $A$ is the following:

- Compute the limiting distribution of $\hat{M}_{n}$ and the associated normalizing sequences $a_{n}$ and $b_{n}$.
- Show that conditions $D\left(u_{n}\right)$ and $D^{\prime}\left(u_{n}\right)$ are valid for the stochastic process $X_{0}, X_{1}, X_{2}, \ldots$ defined in (1.1).
The rest of the paper is dedicated to the proof of these assertions and is structured as follows. In Section 3, we describe the properties of the dynamical systems $f_{a}$, with $a \in \mathcal{B} C$, and the Benedicks-Carleson techniques. The validity of condition $D\left(u_{n}\right)$ is a consequence of the very good mixing properties of the systems considered here. Actually, it follows from the fact that these systems possess a weak-Bernoulli generator (see Section 3.8 and Remark (3.1). Then, in Section 4, we study the asymptotic behavior of the maximum in the i.i.d. case, identify the limiting distribution of $\hat{M}_{n}$ and the respective normalizing sequences $a_{n}$ and $b_{n}$. Hence, we are left with the burden of proving $D^{\prime}\left(u_{n}\right)$. In Section 5 , we use the geometric properties of the systems to show Proposition 5.2 that paves the way for the proof of $D^{\prime}\left(u_{n}\right)$ which is finally established in Section 6. In Section 7, we present a small simulation study in order to compare the finite sample behavior of the normalized $M_{n}$ with the asymptotic one.


## 3. Properties of the Benedicks-Carleson parameters

The Benedicks-Carleson Theorem (see [BC85] or Section 2 of [BC91]) states that there exists a positive Lebesgue measure set of parameters, $\mathcal{B C}$, verifying

$$
\begin{align*}
& \text { there is } c>0(c \approx \log 2) \text { such that }\left|D f_{a}^{n}\left(f_{a}(0)\right)\right| \geq \mathrm{e}^{c n} \text { for all } n \geq 0  \tag{EG}\\
& \text { there is a small } \alpha>0 \text { such that }\left|f_{a}^{n}(0)\right| \geq \mathrm{e}^{-\alpha \sqrt{n}} \text { for all } n \geq 1 \tag{BA}
\end{align*}
$$

Before we describe the Benedicks-Carleson strategy let us have an overview of its key ingredients. The critical region is the interval $(-\delta, \delta)$, where $\delta=\mathrm{e}^{-\Delta}>0$ is chosen small but much larger than $2-a$. This region is partitioned into the intervals

$$
(-\delta, \delta)=\bigcup_{m \geq \Delta} I_{m}
$$

where $I_{m}=\left(\mathrm{e}^{-(m+1)}, \mathrm{e}^{-m}\right]$ for $m>0$ and $I_{m}=-I_{-m}$ for $m<0$; then each $I_{m}$ is further subdivided into $m^{2}$ intervals $\left\{I_{m, j}\right\}$ of equal length inducing the partition $\mathcal{P}_{0}$ of $[-1,1]$ into

$$
\begin{equation*}
[-1,-\delta) \cup \bigcup_{m, j} I_{m, j} \cup(\delta, 1] \tag{3.1}
\end{equation*}
$$

Given $J \in \mathcal{P}$, let $n J$ denote the interval $n$ times the length of $J$ centered at $J$ and define $U_{m}:=\left(-\mathrm{e}^{-m}, \mathrm{e}^{-m}\right)$, for every $m \in \mathbb{N}$. In order to study the growth of $D f_{a}^{n}(x)$, for $x \in[-1,1]$ and $a \in \mathcal{B} C$, the orbit is splitted in free periods and bound periods. During the former we are certain that the orbit never visits the critical region. The latter begin when the orbit returns to the critical region and initiates a bound to the critical point, accompanying its early iterates. The behavior of the derivative during these periods is detailed in Subsections 3.1 and 3.2.
3.1. Expansion outside the critical region. There is $c_{0}>0$ and $M_{0} \in \mathbb{N}$ such that
(1) If $x, \ldots, f_{a}^{k-1}(x) \notin(-\delta, \delta)$ and $k \geq M_{0}$, then $\left|D f_{a}^{k}(x)\right| \geq \mathrm{e}^{c_{0} k}$;
(2) If $x, \ldots, f_{a}^{k-1}(x) \notin(-\delta, \delta)$ and $f_{a}^{k}(x) \in(-\delta, \delta)$, then $\left|D f_{a}^{k}(x)\right| \geq \mathrm{e}^{c_{0} k}$;
(3) If $x, \ldots, f_{a}^{k-1}(x) \notin(-\delta, \delta)$, then $\left|D f_{a}^{k}(x)\right| \geq \delta \mathrm{e}^{c_{0} k}$.

If we were capable of keeping the orbit of $x$ away from the critical region then it would be in free period forever and the estimates above would apply. However, it is inevitable that almost every $x \in[-1,1]$ makes a return to the critical region. We say that $n \in \mathbb{N}$ is a return time of the orbit of $x$ if $f_{a}^{n}(x) \in(-\delta, \delta)$. Every free period of $x$ ends with a free return to the critical region. We say that the return had a depth of $m \in \mathbb{N}$ if $f_{a}^{n}(x) \in I_{ \pm m}$, which is equivalent to saying that $m=\left[-\log \left|f_{a}^{n}(x)\right|\right]$. Once in the critical region, the orbit of $x$ initiates a binding with the critical point.
3.2. Bound period definition and properties. Let $\beta=14 \alpha$. For $x \in(-\delta, \delta)$ define $p(x)$ to be the largest integer $p$ such that

$$
\begin{equation*}
\left|f_{a}^{k}(x)-f_{a}^{k}(0)\right|<\mathrm{e}^{-\beta k}, \quad \forall k<p \tag{3.2}
\end{equation*}
$$

Then
(1) $\frac{1}{2}|m| \leq p(x) \leq 3|m|$, for each $x \in I_{m}$;
(2) $\left|D f_{a}^{p}(x)\right| \geq \mathrm{e}^{c^{\prime} p}$, where $c^{\prime}=\frac{1-4 \beta}{3}>0$.

The orbit of $x$ is said to be bound to the critical point during the period $0 \leq k<p$. We may assume that $p$ is constant on each $I_{m, j}$. Note that during the bound period the orbit of $x$ may return to the critical region. These instants are called bound return times.

Roughly speaking, the idea behind the proof of Benedicks-Carleson Theorem is that while the orbit of the critical point is outside the critical region we have expansion (see Subsection (3.1); when it returns we have a serious setback in the expansion but then, by continuity, the orbit repeats its early history regaining expansion on account of (EG). To arrange for the exponential growth of the derivative along the critical orbit (EG) one has to guarantee that the losses at the returns are not too drastic; hence, by parameter elimination, the basic assumption condition ( $(\overline{B A})$ is imposed. The argument is mounted in a very intricate induction scheme that guarantees both the conditions for the parameters
that survive the exclusions. The condition (EG) is usually known as the Collet-Eckmann condition and it was introduced in CE83.
3.3. Bookkeeping, essential and inessential returns. A sequence of partitions $\mathcal{P}_{0} \prec$ $\mathcal{P}_{1} \prec \ldots$ is built so that points in the same element of the partition $\mathcal{P}_{n}$ have the same history up to time $n$. For a detailed description of the construction of this sequence of partitions in the phase space setting we refer to [Fre05, Section 4]. Here, we highlight some of the main aspects of its construction.

For Lebesgue almost every $x \in I,\{x\}=\cap_{n \geq 0} \omega_{n}(x)$, where $\omega_{n}(x)$ is the element of $\mathcal{P}_{n}$ containing $x$. For such $x$ there is a sequence $t_{1}, t_{2}, \ldots$ corresponding to the instants when the orbit of $x$ experiences a free essential return situation, which means $I_{m, k} \subset f_{a}^{t_{i}}\left(\omega_{t_{i}-1}(x)\right)$ for some $|m| \geq \Delta$ and $1 \leq k \leq m^{2}$. We have that $\omega_{n}(x)=\omega_{t_{i-1}}(x)$, for every $t_{i-1} \leq n<t_{i}$ and $f_{a}^{t_{i}}\left(\omega_{t_{i}}(x)\right)=\omega_{0}\left(f^{t_{i}}(x)\right)$, except for the points at the two ends of $f_{a}^{t_{i}}\left(\omega_{t_{i-1}}(x)\right)$ for which it may occur an adjoining to the neighboring interval. If $t_{i}$ is an essential return situation for $x$, then it is either an essential return time for $x$, which means that there exists $m \geq \Delta$ and $1 \leq k \leq m^{2}$ such that $I_{m, k} \subset f_{a}^{t_{i}}\left(\omega_{t_{i}}(x)\right) \subset 3 I_{m, k}$; or an escaping time for $x$, which is to say that $I_{(\Delta-1), 1} \subset f_{a}^{t_{i}}\left(\omega_{t_{i}}(x)\right) \subset(\delta, 1]$ or $I_{-(\Delta-1), 1} \subset f_{a}^{t_{i}}\left(\omega_{t_{i}}(x)\right) \subset[-1,-\delta)$, where $I_{ \pm(\Delta-1), 1}$ is the subinterval of $I_{ \pm(\Delta-1)}$ closest to 0 .

We remark that every point in $\omega \in \mathcal{P}_{n}$ has the same history up to $n$, in the sense that they have the same free periods, return to the critical region simultaneously, with the same depth and their bound periods expire at the same time.

We say that $v$ is a free return time for $x$ of inessential type if $f_{a}^{v}\left(\omega_{v}(x)\right) \subset 3 I_{m, k}$, for some $|m| \geq \Delta$ and $1 \leq k \leq m^{2}$, but $f_{a}^{v}\left(\omega_{v}(x)\right)$ is not large enough to contain an interval $I_{m, k}$ for some $|m| \geq \Delta$ and $1 \leq k \leq m^{2}$.
3.4. Distortion of the derivative. The sequence of partitions described above is designed so that we have bounded distortion in each element of the partition $\mathcal{P}_{n-1}$ up to time $n$. To be more precise, consider $\omega \in \mathcal{P}_{n-1}$. There exists a constant $C$ independent of $\omega, n$ and the parameter $a$ such that for every $x, y \in \omega$,

$$
\begin{equation*}
\frac{\left|D f_{a}^{n}(x)\right|}{\left|D f_{a}^{n}(y)\right|} \leq C . \tag{3.3}
\end{equation*}
$$

See [Fre05, Lemma 4.2] for a proof.
3.5. Growth of returning and escaping components. Let $t$ be an essential return time for $\omega \in \mathcal{P}_{t}$, with $I_{m, k} \subset f_{a}^{t}(\omega) \subset 3 I_{m, k}$ for some $m \geq \Delta$ and $1 \leq k \leq m^{2}$. If $n$ is the next free return situation for $\omega$ (either essential or inessential) then

$$
\begin{equation*}
\left|f_{a}^{n}(\omega)\right| \geq \mathrm{e}^{c o q} \mathrm{e}^{-5 \beta|m|} \tag{3.4}
\end{equation*}
$$

where $q=n-(t+p)$. See Fre05, Lemma 4.1].
Suppose that $\omega \in \mathcal{P}_{t}$ is an escape component. Then, in the next return situation for $\omega$, at time $t_{1}$, we have that

$$
\begin{equation*}
\left|f_{a}^{t_{1}}(\omega)\right| \geq \mathrm{e}^{-\beta \Delta} \tag{3.5}
\end{equation*}
$$

See MS93], Fre06, Lemma 4.2] or [Mor93, Lemma 5.1] for a proof of a similar statement in the space of parameters.
3.6. Existence of absolutely continuous invariant measures. For every $a \in \mathcal{B} C$, the quadratic map $f_{a}$ has an invariant probability measure $\mu_{a}$ that is absolutely continuous with respect to Lebesgue measure on $[-1,1]$. The existence of absolutely continuous invariant measures (a.c.i.m) for a positive Lebesgue measure set of parameters was first proved by Jakobson in Jak81 and others followed. See, for example, CE83, BC85, Now85, Ryc88, BY92] and You92].

The a.c.i.m. $\mu_{a}=\rho_{a} d x$ has the following properties:
(1) $\mu_{a}$ is the only a.c.i.m. of $f_{a}$;
(2) $\left(f_{a}, \mu_{a}\right)$ is exact;
(3) $\rho_{a}=\rho_{1}^{a}+\rho_{2}^{a}$, where $\rho_{1}^{a}$ has bounded variation and $0 \leq \rho_{2}^{a}(x) \leq \operatorname{const} \sum_{j=1}^{\infty} \frac{(1.9)^{-j}}{\sqrt{\left|x-f_{a}^{j}(0)\right|}}$;
(4) The support of $\mu_{a}$ is $\left[f_{a}^{2}(0), f_{a}(0)\right]$ and $\inf _{x \in\left[f_{a}^{2}(0), f_{a}(0)\right]} \rho_{a}(x)>0$.

The proof of these statements can be found in You92, Theorems 1 and 2].
3.7. Decay of correlations and Central Limit Theorem. The Benedicks-Carleson quadratic maps have good statistical behavior. In fact, L. S. Young proved that these maps have exponential decay of correlations and satisfy the Central Limit Theorem (You92, Theorems 3 and 4]). This was also obtained by Keller and Nowicki in [KN92]. To be more precise, for every $a \in \mathcal{B} C$, there exists $\varsigma \in(0,1)$ such that for all $\varphi, \psi:[-1,1] \rightarrow \mathbb{R}$ with bounded variation, there is $C=C(\varphi, \psi)$ such that

$$
\begin{equation*}
\left|\int \varphi \cdot\left(\psi \circ f_{a}^{n}\right) d \mu_{a}-\int \varphi d \mu_{a} \int \psi d \mu_{a}\right| \leq C \varsigma^{n}, \quad \forall n \geq 0 \tag{3.6}
\end{equation*}
$$

Moreover, if $\int \varphi d \mu_{a}=0$ then for every $x \in \mathbb{R}$ we have

$$
\begin{equation*}
\mu_{a}\left\{\frac{1}{\sqrt{n}} \sum_{i=0}^{n-1} \varphi \circ f_{a}^{i} \leq x\right\} \underset{n \rightarrow \infty}{\longrightarrow} \Phi(x / \sigma) \tag{3.7}
\end{equation*}
$$

where we are assuming that $\sigma:=\lim _{n \rightarrow \infty} \frac{1}{\sqrt{n}}\left[\int\left(\sum_{i=0}^{n-1} \varphi \circ f_{a}^{i}\right)^{2} d \mu_{a}\right]^{1 / 2}>0$ and $\Phi(\cdot)$ denotes the $N(0,1)$ d.f.
3.8. Exponential weak-Bernoulli mixing. Keller Kel94 has obtained a result even sharper than (3.6). Consider the partition of $[-1,1]$ given by $\mathcal{Q}=\{[-1,0),[0,1]\}$. Also, for integers $k<l$, denote by $\mathcal{Q}_{k}^{l}$ the join of partitions $\bigvee_{i=k}^{l} f_{a}^{-i} \mathcal{Q}$ and by $\mathcal{F}_{k}^{l}$ the $\sigma$-algebra generated by $\mathcal{Q}_{k}^{l}$. According to Kel94 the partition $\mathcal{Q}$ is a weak-Bernoulli generator for every $f_{a}$ with $a \in \mathcal{B} C$. This means that the $\sigma$-algebra $\mathcal{F}_{0}^{\infty}$ coincides, up to sets of Lebesgue measure 0 , with the Borel $\sigma$-algebra of sets in $[-1,1]$ and that

$$
\beta_{m}\left(f_{a}, \mathcal{Q}, \mu_{a}\right) \rightarrow 0, \text { as } m \rightarrow \infty,
$$

where

$$
\begin{aligned}
\beta_{m}\left(f_{a}, \mathcal{Q}, \mu_{a}\right): & =2 \sup _{k>0} \int \sup \left\{\left|\mu_{a}\left(A \mid \mathcal{F}_{0}^{k}\right)-\mu_{a}(A)\right|: A \in \mathcal{F}_{k+m}^{\infty}\right\} d \mu_{a} \\
& =\sup _{k \geq 1, L \geq 1} \sum_{A \in \mathcal{Q}_{0}^{k}, B \in \mathcal{Q}_{k+m}^{k+m+L}}\left|\mu_{a}(A \cap B)-\mu_{a}(A) \mu_{a}(B)\right|
\end{aligned}
$$

In fact, Kel94, Theorem 1] states that there are constants $C>0$ and $0<r<1$ such that

$$
\begin{equation*}
\beta_{m}\left(f_{a}, \mathcal{Q}, \mu_{a}\right) \leq C r^{m} \tag{3.8}
\end{equation*}
$$

for all $m \in \mathbb{N}$.
Remark 3.1. We observe that if we refine the partition $\mathcal{Q}$ by adding one point so that $\left\{X_{0}>u\right\} \in \mathcal{F}_{0}$, where $\mathcal{F}_{0}$ is the $\sigma$-algebra generated by $\mathcal{Q}$, then Keller's argument still holds with the same type of estimate as in (3.8). As a consequence, condition $D\left(u_{n}\right)$ is true for every considered sequence $u_{n}$.

## 4. Domain of attraction of the associated i.i.d. Process

We recall that to every stationary stochastic process $X_{0}, X_{1}, X_{2}, \ldots$ defined in (1.1) we associated an i.i.d. sequence of r.v. $Z_{0}, Z_{1}, Z_{2}, \ldots$ with common d.f. given by $G_{a}(x)=$ $P\left\{X_{0} \leq x\right\}=\mu_{a}\{(-\infty, x] \cap[-1,1]\}$ (see Section (2). In this Section we will determine the domain of attraction corresponding to the d.f. $G_{a}$, i.e., we will compute the limiting distribution of $\hat{M}_{n}$, defined in (2.1), when properly normalized. For that purpose one must look at the tail behavior of $1-G_{a}(x)$ as $x$ gets closer to $\sup _{y \in \mathbb{R}}\left\{G_{a}(y)<1\right\}=1$. According to Section 3.6 (3), if $z$ is close to 1 we may write $\mu_{a}\{(z, 1]\}=\mathcal{O}(\sqrt{1-z})$, in the sense that $\frac{\mu_{a}\{(z, 1]\}}{\sqrt{1-z}} \rightarrow c$ for some $c>0$, as $z \rightarrow 1$. Hence, for $s>0$ sufficiently close to 0 we have:

$$
\begin{equation*}
1-G_{a}(1-s)=\mathcal{O}(\sqrt{1-(1-s)})=\mathcal{O}(\sqrt{s}) \tag{4.1}
\end{equation*}
$$

which means that $\lim _{s \rightarrow 0^{+}} \frac{1-G_{a}(1-s)}{\sqrt{s}}>0$. At this point, we apply [LLR83, Theorem 1.6.2] to obtain that $G_{a}$, in this case, belongs to the domain of attraction of type III (Weibull) with parameter $1 / 2$, since for every $x>0$

$$
\lim _{h \rightarrow 0^{+}} \frac{1-G_{a}(1-x h)}{1-G_{a}(1-h)}=\lim _{h \rightarrow 0^{+}} \frac{\sqrt{x h}}{\sqrt{h}}=x^{1 / 2}
$$

Moreover, according to LLR83, Corollary 1.6.3], if we consider the sequences defined for each $n \in \mathbb{N}$ by $b_{n}=1$ and $a_{n}=\left(1-G_{a}^{-1}(1-1 / n)\right)^{-1}$, where $G_{a}^{-1}(y)=\inf \left\{x: G_{a}(x) \geq y\right\}$, then

$$
P\left\{a_{n}\left(\hat{M}_{n}-b_{n}\right) \leq x\right\} \rightarrow H(x)= \begin{cases}\mathrm{e}^{-(-x)^{1 / 2}} & , x \leq 0 \\ 1 & , x>0\end{cases}
$$

as $n \rightarrow \infty$.

## 5. Probability of an essential RETURN REACHING A CERTAIN DEPTH

In the study of extremes, one is mostly interested in the probability of occurrence of exceedances of the level $u_{n}$. As we have already mentioned in Section 2, these events are related with the occurrence of deep returns. Thus, in this section we do some preparatory work by estimating the probability of the returns hitting a given depth.

For each $x \in I$, let $v_{n}(x)$ denote the number of essential return situations of $x$ between 1 and $n, s_{n}(x)$ be the number of those which are actual essential return times and $\mathfrak{S}_{n}$ the number of the latter that correspond to deep essential returns of the orbit of $x$, i.e, with return depths above a threshold $\Theta \geq \Delta$. Observe that $v_{n}(x)-s_{n}(x)$ is the exact number of escaping situations of the orbit of $x$, up to $n$.

Given the integers $0 \leq s \leq \frac{2 n}{\Theta}, s \leq v \leq n$ and the $s$ integers $\gamma_{1}, \ldots, \gamma_{s}$, each greater than or equal to $\Theta$, we define the event:

$$
A_{\gamma_{1}, \ldots, \gamma_{s}}^{v, s}(n)=\left\{x \in I: v_{n}(x)=v, \mathfrak{S}_{n}(x)=s \text { and the depth of the i-th deep essen- } \quad \text { tial return is } \gamma_{i} \forall i \in\{1, \ldots, s\} .\right.
$$

Remark 5.1. Observe that the upper bound $\frac{2 n}{\Theta}$ for the number of deep essential returns up to time $n$ derives from the fact that each deep essential return originates a bound period of length at least $\frac{1}{2} \Theta$ (see Section 3.2). Since during the bound periods there cannot be any essential return, the number of deep essential returns occurring in a period of length $n$ is at most $\frac{n}{\frac{1}{2} \theta}$.
Proposition 5.2. Given the integers $0 \leq s \leq \frac{2 n}{\Theta}$ and $s \leq v \leq n$, consider $s$ integers $\gamma_{1}, \ldots, \gamma_{s}$, each greater than or equal to $\Theta$. If $\Theta$ is large enough, then

$$
\lambda\left(A_{\gamma_{1}, \ldots, \gamma_{s}}^{v, s}(n)\right) \leq\binom{ v}{s} \operatorname{Exp}\left\{-(1-6 \beta) \sum_{i=1}^{s} \gamma_{i}\right\}
$$

Proof. Fix $n \in \mathbb{N}$ and take $\omega_{0} \in \mathcal{P}_{0}$. Note that the functions $v_{n}, s_{n}$ and $\mathfrak{S}_{n}$ are constant on each $\omega \in \mathcal{P}_{n}$. Let $\omega \in \omega_{0} \cap \mathcal{P}_{n}$ be such that $v_{n}(\omega)=v$. Then, there is a sequence $1 \leq t_{1} \leq \ldots \leq t_{v} \leq n$ of essential return situations. Let $\omega_{i}$ denote the element of the partition $\mathcal{P}_{t_{i}}$ that contains $\omega$. We have $\omega_{0} \supset \omega_{1} \supset \ldots \supset \omega_{v}=\omega$. For each $j \in\{0, \ldots, v\}$ we define the set:

$$
Q_{j}=\bigcup_{\omega \in \mathcal{P}_{n} \cap \omega_{0}} \omega_{j}
$$

and its partition

$$
\mathcal{Q}_{j}=\left\{\omega_{j}: \omega \in \mathcal{P}_{n} \cap \omega_{0}\right\} .
$$

Let $\omega \in \mathcal{P}_{n}$ be such that $\mathfrak{S}_{n}(\omega)=s$. Then, we may consider $1 \leq r_{1} \leq \ldots \leq r_{s} \leq v$ with $r_{i}$ indicating that the $i$-th deep essential return occurs in the $r_{i}$-th essential return situation. Now, set $V(0)=Q_{0}=\omega_{0}$. Fix $s$ integers $1 \leq r_{1} \leq \ldots \leq r_{s} \leq v$. Next, for each $j \leq v$ we define recursively the sets $V(j)$. Although the set $V(v)$ will depend on the fixed integers
$1 \leq r_{1} \leq \ldots \leq r_{s} \leq v$, we do not indicate this so that the notation is not overloaded. Suppose that $V(j-1)$ is already defined and $r_{i-1}<j<r_{i}$. Then, we set

$$
V(j)=\bigcup_{\omega \in \mathcal{Q}_{j}} \omega \cap f_{a}^{-t_{j}}\left(I-U_{\Theta}\right) \cap V(j-1) .
$$

If $j=r_{i}$ then we define

$$
V(j)=\bigcup_{\omega \in \mathcal{Q}_{j}} \omega \cap f_{a}^{-t_{j}}\left(I_{\gamma_{i}} \cup I_{-\gamma_{i}}\right) \cap V(j-1)
$$

Observe that for every $j \in\{1, \ldots, v\}$ we have $\frac{|V(j)|}{|V(j-1)|} \leq 1$. Therefore, we concentrate in finding a better estimate for $\frac{\left|V\left(r_{i}\right)\right|}{\left|V\left(r_{i}-1\right)\right|}$. Consider that $\omega_{r_{i}} \in \mathcal{Q}_{r_{i}} \cap V\left(r_{i}\right)$ and let $\omega_{r_{i}-1} \in$ $\mathcal{Q}_{r_{i}-1} \cap V\left(r_{i}-1\right)$ contain $\omega_{r_{i}}$. We have to consider two situations depending on whether $t_{r_{i}-1}$ is an escaping situation or an essential return.

Let us suppose first that $t_{r_{i}-1}$ was an essential return with return depth $\eta$. Then,

$$
\begin{aligned}
\frac{\left|\omega_{r_{i}}\right|}{\left|\omega_{r_{i}-1}\right|} & \leq \frac{\left|\omega_{r_{i}}\right|}{\left|\widehat{\omega}_{r_{i}-1}\right|}, \text { where } \widehat{\omega}_{r_{i}-1}=\omega_{r_{i}-1} \cap f_{a}^{-t_{r_{i}}}\left(U_{1}\right) \\
& \leq C \frac{\left|f_{a}^{t_{r_{i}}}\left(\omega_{r_{i}}\right)\right|}{\left|f_{a}^{t_{r_{i}}}\left(\widehat{\omega}_{r_{i}-1}\right)\right|}, \text { by (3.3) } \\
& \leq C \frac{2 \mathrm{e}^{-\gamma_{i}}}{\mathrm{e}^{-5 \beta \eta}}, \text { by (3.4). }
\end{aligned}
$$

Note that when $r_{i-1}=r_{i}-1$ then $\eta=\gamma_{i-1}$. If, on the other hand, $r_{i-1}<r_{i}-1$ then $t_{r_{i}-1}$ is an essential return with depth $\eta<\Theta \leq \gamma_{i-1}$. Thus, in both situations, we have

$$
\frac{\left|\omega_{r_{i}}\right|}{\left|\omega_{r_{i}-1}\right|} \leq 2 C \frac{\mathrm{e}^{-\gamma_{i}}}{\mathrm{e}^{-5 \beta \gamma_{i-1}}} .
$$

When $t_{r_{i}-1}$ is an escape situation, instead of using (3.4), we can use (3.5) and obtain

$$
\frac{\left|\omega_{r_{i}}\right|}{\left|\omega_{r_{i}-1}\right|} \leq 2 C \frac{\mathrm{e}^{-\gamma_{i}}}{\mathrm{e}^{-\beta \Delta}} \leq 2 C \frac{\mathrm{e}^{-\gamma_{i}}}{\mathrm{e}^{-5 \beta \gamma_{i-1}}} .
$$

Observe also that if $\widehat{\omega}_{r_{i}-1} \neq \omega_{r_{i}-1}$ then, because we are assuming that $\omega_{r_{i}} \neq \emptyset$, we have $\left|f_{a}^{t_{r_{i}}}\left(\widehat{\omega}_{r_{i}-1}\right)\right| \geq \mathrm{e}^{-1}-\mathrm{e}^{-\Theta} \geq \mathrm{e}^{-5 \beta \gamma_{i-1}}$, for large $\Theta$.

At this point we may write

$$
\begin{aligned}
\left|V\left(r_{i}\right)\right| & =\sum_{\omega_{r_{i}} \in \mathcal{Q}_{r_{i}} \cap V\left(r_{i}\right)} \frac{\left|\omega_{r_{i}}\right|}{\left|\omega_{r_{i}-1}\right|}\left|\omega_{r_{i}-1}\right| \\
& \leq 2 C \mathrm{e}^{-\gamma_{i}} \mathrm{e}^{5 \beta \gamma_{i-1}} \sum_{\omega_{r_{i}} \in \mathcal{Q}_{r_{i}} \cap V\left(r_{i}\right)}\left|\omega_{r_{i}-1}\right| \\
& \leq 2 C \mathrm{e}^{-\gamma_{i}} \mathrm{e}^{5 \beta \gamma_{i-1}}\left|V\left(r_{i}-1\right)\right| .
\end{aligned}
$$

This yields

$$
|V(v)| \leq(2 C)^{s} \operatorname{Exp}\left\{-(1-5 \beta) \sum_{i=1}^{s} \gamma_{i}\right\} \mathrm{e}^{5 \beta \gamma_{0}}|V(0)|
$$

where $\gamma_{0}$ is given by the interval $\omega_{0} \in \mathcal{P}_{0}$. If $\omega_{0}=I_{\left(\eta_{0}, k_{0}\right)}$ with $\left|\eta_{0}\right| \geq \Delta$ and $1 \leq k_{0} \leq \eta_{0}^{2}$, then $\gamma_{0}=\left|\eta_{0}\right|$. If $\omega_{0}=(\delta, 1]$ or $\omega_{0}=[-1,-\delta)$, then we can take $\gamma_{0}=0$.

Now, we have to take into account the number of possibilities of having the occurrence of the event $V(v)$ implying the occurrence of the event $A_{\gamma_{1}, \ldots, \gamma_{s}}^{v, s}(n)$. The number of possible configurations related with the different values that the integers $r_{1}, \ldots r_{s}$ can take is $\binom{v}{s}$. Hence, it follows that

$$
\begin{aligned}
\lambda\left(A_{\gamma_{1}, \ldots, \gamma_{s}}^{v, s}(n)\right) & \leq(2 C)^{s}\binom{v}{s} \operatorname{Exp}\left\{-(1-5 \beta) \sum_{i=1}^{s} \gamma_{i}\right\} \sum_{\omega_{o} \in \mathcal{P}_{0}} \mathrm{e}^{5 \beta\left|\gamma_{0}\right|}\left|\omega_{0}\right| \\
& \leq(2 C)^{s}\binom{v}{s} \operatorname{Exp}\left\{-(1-5 \beta) \sum_{i=1}^{s} \gamma_{i}\right\}\left(2(1-\delta)+\sum_{\left|\eta_{0}\right| \geq \Delta} \mathrm{e}^{5 \beta \eta_{0}} \mathrm{e}^{-\left|\eta_{0}\right|}\right) \\
& \leq 3(2 C)^{s}\binom{v}{s} \operatorname{Exp}\left\{-(1-5 \beta) \sum_{i=1}^{s} \gamma_{i}\right\}, \quad \text { for } \Delta \text { large enough } \\
& \leq\binom{ v}{s} \operatorname{Exp}\left\{-(1-6 \beta) \sum_{i=1}^{s} \gamma_{i}\right\} .
\end{aligned}
$$

The last inequality results from the fact that $s \Theta \leq \sum_{i=1}^{s} \gamma_{i}$ and the freedom to choose a sufficiently large $\Theta$.

Given the integers $0 \leq s \leq \frac{2 n}{\Theta}, s \leq v \leq n$ and the integers $\gamma_{0}, \gamma_{1}$, both greater than or equal to $\Theta$, we consider the event:

$$
B_{\gamma_{0}, \gamma_{1}}^{v, s}(n)=\left\{x \in I_{\gamma_{0}}: v_{n}(x)=v, \mathfrak{S}_{n}(x)=s \text { and } n \text { is a free deep return } \begin{array}{c}
\text { with depth } \gamma_{1}
\end{array}\right\}
$$

Corollary 5.3. Consider the integers $0 \leq s \leq \frac{2 n}{\Theta}, s \leq v \leq n$ and $\gamma_{0}, \gamma_{1} \geq \Theta$. If $\Theta$ is large enough, then

$$
\lambda\left(B_{\gamma_{0}, \gamma_{1}}^{v, s}(n)\right) \leq\binom{ v}{s} \operatorname{Exp}\left\{-(1-6 \beta)\left(\gamma_{0}+\gamma_{1}\right)\right\}
$$

The proof of this statement follows easily from Proposition 5.2, by observing that, although $n$ may be an inessential deep return time (instead of an essential deep return), the estimates still prevail and for $\Theta>\Delta$ large enough we have $\sum_{\gamma \geq \Theta} \mathrm{e}^{-(1-6 \beta) \gamma} \leq 1$.

## 6. The condition $D^{\prime}\left(u_{n}\right)$

Assume that $X_{0}, X_{1}, \ldots$ is the stationary stochastic process defined in (1.1) with common d.f. $G_{a}$. For $z \in[-1,1]$ the event $\left\{X_{j}(z)>u_{n}\right\}$ corresponds to the set $f_{a}^{-j}\left(\left(u_{n}, 1\right]\right)$. If $n$ is sufficiently large, then $f_{a}^{-1}\left(u_{n}, 1\right]=\left(-\sqrt{\left(1-u_{n}\right) / a}, \sqrt{\left(1-u_{n}\right) / a}\right) \subset U_{\Delta}$. We may define

$$
\begin{equation*}
\Theta=\Theta(n)=\left[-\frac{1}{2} \log \frac{1-u_{n}}{a}\right] . \tag{6.1}
\end{equation*}
$$

This way, if an exceedance occurs at time $j$ then a deep return with depth over the threshold $\Theta$ must have happened at time $j-1$, i.e., if $X_{j}(z)>u_{n}$ then $X_{j-1}(z)=f_{a}^{j-1}(z) \in U_{\Theta}$.

Remember that the sequence $u_{n}$ is such that $n\left(1-G_{a}\left(u_{n}\right)\right) \rightarrow \tau$, as $n \rightarrow \infty$, which we rewrite as $1-G_{a}\left(u_{n}\right)=\mathcal{O}(1 / n)$. Then, by (4.1), we get $u_{n}=1-\mathcal{O}\left(1 / n^{2}\right)$, which according to (6.1) leads to

$$
\begin{equation*}
\Theta=\mathcal{O}(\log n) \tag{6.2}
\end{equation*}
$$

meaning that $\frac{\Theta}{\log n} \rightarrow c$, for some $c>0$, as $n \rightarrow \infty$.
Observe that we are dealing with very small perturbations of $f_{2}$ for which $f_{2}^{j}(0)=-1$ for every $j \geq 2$. Thus, one expects that after a deep return to the critical region (a tight vicinity of 0 ) it should take a considerable amount of time before another deep return should occur. Since exceedances are related with the occurrence of deep returns then one may have a fair amount of belief that condition (2.2) holds for the sequence $X_{0}, X_{1}, \ldots$

Remark 6.1. If the sequence $X_{0}, X_{1}, \ldots$ was independent, then (2.2) would follow easily since

$$
\begin{aligned}
n \sum_{j=1}^{[n / k]} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\} & =n \sum_{j=1}^{[n / k]} P\left\{X_{0}>u_{n}\right\} P\left\{X_{j}>u_{n}\right\}=n \sum_{j=1}^{[n / k]}\left(1-G_{a}\left(u_{n}\right)\right)^{2} \\
& \leq \frac{n^{2}}{k}\left(1-G_{a}\left(u_{n}\right)\right)^{2} \xrightarrow[n \rightarrow \infty]{ } \frac{\tau^{2}}{k} \xrightarrow[k \rightarrow \infty]{ } 0 .
\end{aligned}
$$

Let us give some insight into the argument we use to prove that (2.2) holds for $X_{0}, X_{1}, \ldots$
(1) We use the exponential decay of correlations (see (3.6)) to compute a turning instant $T=T(n)$ such that the dependence between $X_{0}$ and $X_{j}$ with $j>T$ is negligible. This suggests the splitting of the time interval $\{1, \ldots,[n / k]\}$ into $\{1, \ldots, T\}$ and $\{T+1, \ldots,[n / k]\}$, when $n$ is sufficiently large.
(2) During the time interval $\{T+1, \ldots,[n / k]\}$ we use the fact that for $j>T$ the r.v. $X_{j}$ is almost independent of $X_{0}$ and argue like in Remark 6.1.
(3) For $j \in\{1, \ldots, T\}$ we appeal to Corollary 5.3 to bound $P\left\{X_{0}>u_{n}\right.$ and $\left.X_{j}>u_{n}\right\}$ and afterwards we use the fact that, for $n$ large, $T \ll[n / k]$ to finish the proof.

## Step (1)

Taking $\varphi=\psi=\mathbf{1}_{\left(u_{n}, 1\right]}$ in (3.6), we get

$$
\begin{aligned}
\mid \mu_{a}\left\{X_{0}>u_{n}\right. \text { and } & \left.X_{j}>u_{n}\right\}-\left[\mu_{a}\left\{X_{0}>u_{n}\right\}\right]^{2} \mid= \\
& =\left|\int \mathbf{1}_{\left(u_{n}, 1\right]} \cdot \mathbf{1}_{\left(u_{n}, 1\right]} \circ f_{a}^{j} d \mu_{a}-\left(\int \mathbf{1}_{\left(u_{n}, 1\right]} d \mu_{a}\right)^{2}\right| \\
& \leq C \varsigma^{j}
\end{aligned}
$$

where we may assume that $C$ is the same for all $n \in \mathbb{N}$, because $\left\|\mathbf{1}_{\left(u_{n}, 1\right]}\right\|_{\infty}=1$ and the total variation of $\mathbf{1}_{\left(u_{n}, 1\right]}$ is equal to 1 , for every $n \in \mathbb{N}$.

We compute $T=T(n)$ such that for every $j \geq T$ we have

$$
C \varsigma^{j}<\frac{1}{n^{3}} .
$$

Since $C \varsigma^{j}<\frac{1}{n^{3}} \Leftrightarrow j>\frac{1}{\log \varsigma^{-1}}(3 \log n+\log C)$, we simply take, for $n$ sufficiently large,

$$
\begin{equation*}
T=\frac{4}{\log \varsigma^{-1}} \log n \tag{6.3}
\end{equation*}
$$

For fixed $k$ and $n$ sufficiently large, we have that $T<[n / k]$. Hence, we may write

$$
\begin{aligned}
& n \sum_{j=1}^{[n / k]} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\}= \\
& \\
& \qquad=n \sum_{j=1}^{T} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\}+n \sum_{j=T+1}^{[n / k]} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\}
\end{aligned}
$$

In step (2) below we deal with the second term in the sum, leaving the first term for step (3).

## Step (2)

Let us show that $\lim \sup _{n \rightarrow \infty} n \sum_{j=T+1}^{[n / k]} P\left\{X_{0}>u_{n}\right.$ and $\left.X_{j}>u_{n}\right\} \rightarrow 0$, as $k \rightarrow \infty$. By choice of $T$ we have

$$
\begin{aligned}
n \sum_{j=T+1}^{[n / k]} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\} & \leq n\left(\left(1-G_{a}\left(u_{n}\right)\right)^{2}+\frac{1}{n^{3}}\right)[n / k] \\
& \leq \frac{n^{2}}{k}\left(1-G_{a}\left(u_{n}\right)\right)^{2}+\frac{n^{2}}{k n^{3}}
\end{aligned}
$$

Now, $\frac{n^{2}}{k}\left(1-G_{a}\left(u_{n}\right)\right)^{2}+\frac{n^{2}}{k n^{3}} \xrightarrow[n \rightarrow \infty]{ } \frac{\tau^{2}}{k} \underset{k \rightarrow \infty}{\longrightarrow} 0$ and the result follows.

## Step (3)

We are left with the burden of controlling the term $n \sum_{j=1}^{T} P\left\{X_{0}>u_{n}\right.$ and $\left.X_{j}>u_{n}\right\}$. We begin with the following lemma that will enable us to bound the number of exceedances occurring during the time period $\{1, \ldots, T\}$. In what follows, we are always assuming that $n$ is large enough so that $\Theta>\Delta$.

Lemma 6.2. If a deep return occurs at time $t$ (with depth over the threshold $\Theta$ ), then the next deep return can only occur after $t+\Theta / 2$.

Proof. For every $x \in U_{\Theta}$, the bound period associated to $x$ is such that $p(x) \geq \Theta / 2$, by Section 3.2 (11). For all $j \leq[\Theta / 2]$ we have

$$
\begin{aligned}
\left|f_{a}^{j}(x)\right| & \geq\left|f_{a}^{j}(0)\right|-\mathrm{e}^{-\beta j} \stackrel{\mid \overline{B A})}{\geq} \mathrm{e}^{-\alpha \sqrt{j}}-\mathrm{e}^{-\beta j} \geq \mathrm{e}^{-\alpha j}\left(1-\mathrm{e}^{(\alpha-\beta) j}\right) \\
& \geq \mathrm{e}^{-\alpha j}\left(1-\mathrm{e}^{(\alpha-\beta)}\right), \text { since } \alpha-\beta<0 \\
& \geq \mathrm{e}^{-\alpha \Theta / 2}\left(1-\mathrm{e}^{(\alpha-\beta)}\right), \text { since } j \leq \Theta / 2 \\
& \geq \mathrm{e}^{-\alpha \Theta}, \text { if } n \text { is large enough so that } 1-\mathrm{e}^{\alpha-\beta} \geq \mathrm{e}^{-\alpha \Theta / 2} \\
& \geq \mathrm{e}^{-\Theta}, \text { since } \alpha<1 .
\end{aligned}
$$

As a consequence of Lemma 6.2 we have that the maximum number of exceedances up to time $T$ is at most $2 T / \Theta$.

Next lemma shows that in the bound period following a deep return over the level $\Theta$, there cannot occur deep bound returns before time $T$.

Lemma 6.3. For every $a \in \mathcal{B} C$, if $n$ is sufficiently large then for every $x \in U_{\Theta}$ we have that $f_{a}^{j}(x) \notin U_{\Theta}$ for all $j \leq \min \{p(x), T\}$.

Proof. Consider the map $h:[1, \infty) \rightarrow \mathbb{R}$ given by $h(y)=\mathrm{e}^{-\alpha \sqrt{y}}-\mathrm{e}^{-\beta y}$. There is $J=$ $J(\alpha, \beta) \in \mathbb{N}$ such that $h^{\prime}(y)<0$ for every $y \geq J$. Let $D=\mathrm{e}^{-\alpha \sqrt{J}}$.

Having in mind that by (6.2) and (6.3) we have $\Theta=\Theta(n)=\mathcal{O}(\log n)$ and $T=\mathcal{O}(\log n)$, respectively, assume that $n$ is large enough so that:
(1) $\mathrm{e}^{-\Theta}<D / 2$;
(2) by continuity of $f_{a}$, for every $x \in U_{\Theta}$ and all $j \leq J$ we have $\left|f_{a}^{j}(x)-f_{a}^{j}(0)\right|<D / 2$; (3) $\mathrm{e}^{-\alpha \sqrt{T}}-\mathrm{e}^{-\beta T}>\mathrm{e}^{-\Theta}$.

Let $x \in U_{\Theta}$ and consider $j \leq \min \{p(x), J\}$. By assumption (2) we have $\left|f_{a}^{j}(x)-f_{a}^{j}(0)\right|<$ $D / 2$. Using (BA) and hypothesis (1) we get

$$
\left|f_{a}^{j}(x)\right| \geq\left|f_{a}^{j}(0)\right|-\left|f_{a}^{j}(x)-f_{a}^{j}(0)\right| \geq D-D / 2 \geq D / 2>e^{-\Theta}
$$

Now, let $x \in U_{\Theta}$ and consider $J<j \leq \min \{p(x), T\}$. Using the definition of bound period, (BA) and condition (3) we have

$$
\left|f_{a}^{j}(x)\right| \geq\left|f_{a}^{j}(0)\right|-\left|f_{a}^{j}(x)-f_{a}^{j}(0)\right| \geq \mathrm{e}^{-\alpha \sqrt{j}}-\mathrm{e}^{-\beta j} \geq \mathrm{e}^{-\alpha \sqrt{T}}-\mathrm{e}^{-\beta T}>\mathrm{e}^{-\Theta}
$$

where the third inequality derives from the fact that on $[J, \infty)$ the function $h$ is decreasing.

As a consequence, if there is a deep return at time $t$, then we cannot have bound returns during the time period $(t, t+T]$. This way, we may use Corollary 5.3 to estimate $\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right|$, for $j \leq T$.

Observe that for $n$ large enough, we have

$$
\begin{aligned}
P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\} & =\mu_{a}\left\{\left(u_{n}, 1\right] \cap f_{a}^{-j}\left(u_{n}, 1\right]\right\} \\
& =\mu_{a}\left\{f_{a}^{-1}\left(u_{n}, 1\right] \cap f_{a}^{-(j+1)}\left(u_{n}, 1\right]\right\}
\end{aligned}
$$

and $f_{a}^{-1}\left(u_{n}, 1\right] \cap f_{a}^{-j}\left(f_{a}^{-1}\left(u_{n}, 1\right]\right) \subset U_{\Theta} \cap f_{a}^{-j} U_{\Theta}$. In what follows we will use "const" to denote several positive constants independent of $n$.

Note that

$$
U_{\Theta} \cap f_{a}^{-j} U_{\Theta} \subset \bigcup_{s=0}^{2 j / \Theta} \bigcup_{v=s}^{j} \bigcup_{\gamma_{0} \geq \Theta} \bigcup_{\gamma_{1} \geq \Theta} B_{\gamma_{0}, \gamma_{1}}^{v, s}(j) .
$$

Hence, by Corollary 5.3, we have

$$
\begin{aligned}
\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right| & \leq \sum_{s=0}^{2 j / \Theta} \sum_{v=s}^{j} \sum_{\gamma_{0} \geq \Theta} \sum_{\gamma_{1} \geq \Theta}\binom{v}{s} \mathrm{e}^{-(1-6 \beta)\left(\gamma_{0}+\gamma_{1}\right)} \leq \text { const } \sum_{s=0}^{2 T / \Theta} \sum_{v=s}^{T}\binom{v}{s} \mathrm{e}^{-(1-6 \beta) 2 \Theta} \\
& \leq \text { const } \sum_{s=0}^{2 T / \Theta} \sum_{v=s}^{T}\binom{T}{s} \mathrm{e}^{-(1-6 \beta) 2 \Theta} \leq \mathrm{const} \sum_{s=0}^{2 T / \Theta} T\binom{T}{s} \mathrm{e}^{-(1-6 \beta) 2 \Theta}
\end{aligned}
$$

At this point, we estimate $2 T / \Theta$. Recalling (6.2) and (6.3), one easily gets that there exists a constant $C_{1}>0$ such that $2 T / \Theta \leq C_{1}$, for $n$ sufficiently large. So, to proceed with the estimation $\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right|$, assume that $n$ is sufficiently large so that $2 T / \Theta \leq C_{1}$ and $C_{1} \ll T$. Then,

$$
\begin{aligned}
\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right| & \leq \text { const } \sum_{s=0}^{2 T / \Theta} T\binom{T}{s} \mathrm{e}^{-(1-6 \beta) 2 \Theta}
\end{aligned} \leq \text { const } \sum_{s=0}^{2 T / \Theta} T\binom{T}{C_{1}} \mathrm{e}^{-(1-6 \beta) 2 \Theta} .
$$

So far we managed to obtain an estimate for $\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right|$. Let us see how too derive one for $P\left\{X_{0}>u_{n}\right.$ and $\left.X_{j}>u_{n}\right\} \leq \mu_{a}\left\{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right\}$. By assertion (3) in Section 3.6 we
have

$$
\begin{aligned}
\mu_{a}\left\{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right\}= & \int_{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}} \rho_{1}^{a}(x) d x+\int_{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}} \rho_{2}^{a}(x) d x \\
\leq & \text { const. }\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right|+\sum_{j=1}^{\Theta^{2}} \int_{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}} \frac{(1.9)^{-j}}{\sqrt{\left|x-f_{a}^{j}(0)\right|}} d x \\
& +\sum_{j=\Theta^{2}}^{\infty} \int_{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}} \frac{(1.9)^{-j}}{\sqrt{\left|x-f_{a}^{j}(0)\right|}} d x .
\end{aligned}
$$

Using (BA) we get $\left|f_{a}^{j}(0)\right|>\mathrm{e}^{-\alpha \Theta}$, for all $j \leq \Theta^{2}$. Thus, for all $x \in U_{\Theta}$ and $\Theta$ large enough, we have $\sqrt{\left|x-f_{a}^{j}(0)\right|}>\sqrt{\mathrm{e}^{-\alpha \Theta}-\mathrm{e}^{-\Theta}}>1 / 2 . \mathrm{e}^{-\alpha \Theta / 2}$, which implies that

$$
\begin{aligned}
\frac{1}{\sqrt{\left|x-f_{a}^{j}(0)\right|}} \leq 2 \mathrm{e}^{\alpha \Theta / 2} & \left.\leq 2 \mathrm{e}^{\beta \Theta} \text { (recall that } \beta>\alpha\right) . \text { Consequently, for sufficiently large } \Theta \\
\mu_{a}\left\{U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right\} & \leq \text { const. }\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right|+\text { const. }\left|U_{\Theta} \cap f_{a}^{-j} U_{\Theta}\right| . \mathrm{e}^{\beta \Theta}+\text { const. } \mathrm{e}^{\Theta^{2} \log 1.9} \\
& \leq \text { const. } T^{C_{1}+1} \mathrm{e}^{-(1-7 \beta) 2 \Theta}
\end{aligned}
$$

Now, since $\mathrm{e}^{-\Theta} \leq \operatorname{const}\left(1-G_{a}\left(u_{n}\right)\right)$, we finally conclude:

$$
\begin{aligned}
n \sum_{j=1}^{T} P\left\{X_{0}>u_{n} \text { and } X_{j}>u_{n}\right\} & \leq \text { const } \cdot n \sum_{j=1}^{T} T^{C_{1}+1}\left(1-G_{a}\left(u_{n}\right)\right)^{2(1-7 \beta)} \\
& \leq \text { const } \cdot n T^{C_{1}+2}\left(1-G_{a}\left(u_{n}\right)\right)^{2(1-7 \beta)} \\
& \leq \text { const } \cdot n(\log n)^{C_{1}+2}\left(1-G_{a}\left(u_{n}\right)\right)^{2-14 \beta} \\
& \leq \text { const } \cdot\left[n\left(1-G_{a}\left(u_{n}\right)\right)\right]^{3 / 2}\left(1-G_{a}\left(u_{n}\right)\right)^{1 / 2-14 \beta}
\end{aligned}
$$

for sufficiently large $n$. The result follows because

$$
\lim _{n \rightarrow \infty}\left[n\left(1-G_{a}\left(u_{n}\right)\right)\right]^{3 / 2}\left(1-G_{a}\left(u_{n}\right)\right)^{1 / 2-14 \beta}=\tau^{3 / 2} \cdot 0=0
$$

## 7. Simulation Study

In this section we present a small simulation study illustrating the finite sample behavior of the normalized $M_{n}$, defined in (1.2), for the Benedicks-Carleson quadratic map $f_{2}(x)=$ $1-2 x^{2}$. We have that $P\left\{X_{0} \leq x\right\}=G_{2}(x)=1 / 2+\arcsin (x) / \pi$. According to Theorem, the normalizing sequence is $a_{n}=(1-\cos (\pi / n))^{-1}$, for each $n \in \mathbb{N}$, and the theoretical limiting distribution for $a_{n}\left(M_{n}-1\right)$ is

$$
H(x)= \begin{cases}\mathrm{e}^{-(-x)^{1 / 2}} & , x \leq 0 \\ 1 & , x>0\end{cases}
$$

We performed the following experiment. We picked at random (according to the d.f. $G_{2}$ ) a point $z$ in the interval $[-1,1]$, computed its orbit up to time $n$ and calculated
$M_{n}(z)=\max \left\{z, f_{2}(z), \ldots, f_{2}^{n-1}(z)\right\}$. We repeated the process $m$ times to obtain a sample $\left\{M_{n}\left(z_{1}\right), \ldots, M_{n}\left(z_{m}\right)\right\}$ and approximated, for certain values of $x, P\left\{a_{n}\left(M_{n}-b_{n}\right) \leq x\right\}$ by

$$
\begin{equation*}
\frac{1}{m} \sum_{i=1}^{m} \mathbf{1}_{\left\{a_{n}\left(M_{n}\left(z_{i}\right)-b_{n}\right) \leq x\right\}}, \tag{7.1}
\end{equation*}
$$

where $1_{\left\{a_{n}\left(M_{n}\left(z_{i}\right)-b_{n}\right) \leq x\right\}}=\left\{\begin{array}{ll}1 & \text { if } a_{n}\left(M_{n}\left(z_{i}\right)-b_{n}\right) \leq x \\ 0 & \text { if } a_{n}\left(M_{n}\left(z_{i}\right)-b_{n}\right)>x\end{array}\right.$, for each $1 \leq i \leq m$.
In Table 1 we present the results obtained by realizing the above experiment, considering different values of $x, n$ and $m$ and compare them with the theoretical limiting ones given by $H(x)$.

|  |  | $n=1000$ |  | $n=10000$ |  | $n=20000$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $x$ | $H(x)$ | $m=1000$ | $m=10000$ | $m=10000$ | $m=20000$ | $m=20000$ |
| -0.001 | 0.9689 | 0.976 | 0.9671 | 0.9719 | 0.9677 | 0.9708 |
| -0.01 | 0.9048 | 0.894 | 0.9079 | 0.9056 | 0.9076 | 0.9052 |
| -0.1 | 0.7289 | 0.724 | 0.7263 | 0.7323 | 0.7303 | 0.7335 |
| -0.3 | 0.5783 | 0.569 | 0.5773 | 0.5823 | 0.5840 | 0.5813 |
| -0.5 | 0.4931 | 0.491 | 0.4906 | 0.5012 | 0.4984 | 0.4941 |
| -0.7 | 0.4332 | 0.407 | 0.4272 | 0.4403 | 0.4374 | 0.4338 |
| -1 | 0.3679 | 0.388 | 0.3631 | 0.3663 | 0.3729 | 0.3678 |
| -3 | 0.1769 | 0.164 | 0.1731 | 0.1748 | 0.1833 | 0.1729 |
| -5 | 0.1069 | 0.124 | 0.1024 | 0.1092 | 0.1108 | 0.1056 |
| -8 | 0.0591 | 0.059 | 0.0510 | 0.0557 | 0.0617 | 0.0580 |
| -10 | 0.0423 | 0.049 | 0.0350 | 0.0435 | 0.0438 | 0.0414 |
| -30 | 0.0042 | 0.002 | 0.0031 | 0.0033 | 0.0048 | 0.0041 |
| -50 | 0.00085 | 0.001 | 0.0007 | 0.0009 | 0.0007 | 0.0009 |

TABLE 1. Simulation results

As one may verify the results of the experiment are quite close to the asymptotic theoretical ones and there is a general tendency of getting better as $n$ increases which is precisely the behavior we were expecting. It is also noticeable that there is an improvement when $m$ increases, which is also predictable since our approximation (7.1) gets to be more accurate.

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Jorge Milhazes Freitas, Centro de Matemática da Universidade do Porto, Rua do Campo Alegre 687, 4169-007 Porto, Portugal

E-mail address: jmfreita@fc.up.pt
URL: http://www.fc.up.pt/pessoas/jmfreita
Ana Cristina Moreira Freitas, Centro de Matemática \& Faculdade de Economia da Universidade do Porto, Rua Dr. Roberto Frias, 4200-464 Porto, Portugal

E-mail address: amoreira@fep.up.pt


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