# ERGODIC PROPERTIES OF BROWNIAN MOTION 

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## 0. Introduction

Since Brownian motion is point recurrent in $\boldsymbol{R}^{1}$, recurrent in $\boldsymbol{R}^{2}$ and transient in $\boldsymbol{R}^{n}$, $n \geqq 3$ (see (7)), it follows that the total time spent in a bounded open set in $\boldsymbol{R}^{1}$ or $\boldsymbol{R}^{2}$ is unbounded. With the following ergodic theorems for Brownian motion in $\boldsymbol{R}^{1}$ and $\boldsymbol{R}^{2}$ as motivation, we examine the rate of convergence in these theorems. Note that there is no ergodic property in $\boldsymbol{R}^{n}$ for $n \geqq 3$ since Brownian motion is not dense there.

Theorem 0.1. If $\{X(t), 0 \leqq t<\infty\}$ is a separable Brownian motion process in $\boldsymbol{R}^{1}$ and if $f$ and $g$ are any two Baire functions with finite integrals $\bar{f}$ and $\bar{g} \neq 0$ respectively over $(-\infty$, $\infty$ ), then

$$
\lim _{T \rightarrow \infty} \frac{\int_{0}^{T} f\{X(t)\} d t}{\int_{0}^{T} g\{X(t)\} d t}=\frac{\bar{f}}{\vec{g}}
$$

with probability one. See (3).
Corollary (Ergodic Theorem for Brownian motion in $\boldsymbol{R}^{1}$ ). Let $\{X(t), 0 \leqq t<\infty\}$ be a separable Brownian motion process in $\boldsymbol{R}^{1}$. If $A$ and $B$ are bounded measurable non-empty subsets of $\boldsymbol{R}^{\mathbf{1}}$, then

$$
\lim _{\tau \rightarrow \infty} \frac{\text { total time spent in } A \text { by } X(t) \text { up to time } T}{\text { total time spent in } B \text { by } X(t) \text { up to time } T}=\frac{|A|}{|B|} \text { a.s. }
$$

Theorem 0.2 (Ergodic Theorem for Brownian motion in $\boldsymbol{R}^{2}$ ). Let $X(t)$ be a Brownian motion process in $\boldsymbol{R}^{2}$. Let $D_{1}$ and $D_{2}$ be bounded open sets in the plane such that $D_{2} \neq \phi$. Then

$$
\lim _{\mathrm{T} \rightarrow \infty} \frac{\text { total time spent in } D_{1} \text { by } X(t) \text { up to time } T}{\text { total time spent in } D_{2} \text { by } X(t) \text { up to time } T}=\frac{m\left(D_{1}\right)}{m\left(D_{2}\right)} \text { a.s., }
$$

where $\boldsymbol{m}\left(D_{i}\right)$ is Lebesgue measure of $D_{i}$ in $\boldsymbol{R}^{2}$. See (8).
In $\boldsymbol{R}^{2}$ there are some independence problems for any sets $D_{i}$ which are overcome by considering a stationary Markov chain determined by the process. We are able to prove that, almost surely,

$$
\int_{0}^{T}\left[m\left(D_{2}\right) \chi_{D_{1}}(X(t))-m\left(D_{1}\right) \chi_{D_{2}}(X(t))\right] d t
$$

which measures the difference in time spent in the sets $D_{1}, D_{2}$ is unbounded as $T \rightarrow \infty$ although the ratio of this difference to

$$
\int_{0}^{T} \chi_{D_{i}}(X(t)) d t
$$

tends to zero at a rate given by a suitable law of iterated logarithm.
Throughout this paper we shall assume that we are dealing with a separable version of Brownian motion process denoted by $X(t)=X(t, \omega)=X_{t}$.
$c, C^{\prime}, C_{0}, c_{1}, \ldots$ will denote a finite positive constant whose value is not important and not necessarily the same at different occurrences. Other notations are
$\chi_{\mathrm{A}} \quad$ for indicator function of set $A$,
a.s. for almost surely,
$O(x)$ for "large order" of $x$,
[ $t$ ] for integer part of $t$,
$A^{c}$ for complement of set $A$,
$\bar{A} \quad$ for closure of set $A$,
$d(x, y)$ for distance between $x$ and $y$,
$\partial A$ for boundary of set $A$,
$P_{x}, E_{x}$ for conditional probability and expectation respectively given $X(0)=x$.

## 1. Rate of convergence in $\boldsymbol{R}^{\mathbf{2}}$

Theorem 1.1. For any bounded sets $A$ and $B$ in $\boldsymbol{R}^{2}$,

$$
g(t)=E_{x} \int_{0}^{T} m(B) \chi_{A}(X(t)) d t-E_{x} \int_{0}^{T} m(A) \chi_{B}(X(t)) d t
$$

conver $\simeq 0$ es to a finite limit as $T \rightarrow \infty$; where $m($.$) denotes Lebesgue measure in \boldsymbol{R}^{2}$.
Proof. $g(t)=\int_{0}^{T} \int_{A} \frac{m(B)}{2 \pi t} e^{-|x-y| z / 2 t} d y d t-\int_{0}^{T} \int_{B} \frac{m(A)}{2 \pi t} e^{-|x-z| z / 2 t} d z d t$

$$
=\int_{0}^{T} f_{A}(t)-f_{B}(t) d t \text { say. }
$$

Since $A$ and $B$ are bounded sets, we have

$$
\left|\int_{T_{0}}^{T}\left(f_{A}(t)-f_{B}(t)\right) d t\right|=O\left(\frac{1}{T_{o}}\right)
$$

for sufficiently large $T_{0}$. To complete the proof we apply the Cauchy condition for infinite integrals (see e.g. page 433 of (1)) to the function $f_{A}(t)-f_{B}(t)$.

Theorem 1.2. Let $D_{1}$ and $D_{2}$ be bounded non-empty open sets in $\boldsymbol{R}^{2}$ such that $\bar{D}_{1} \cap \bar{D}_{2}=\varnothing$. Then with probability one,

$$
f(T, \omega)=\int_{0}^{T}\left[m\left(D_{2}\right) \chi_{D_{1}}\left(X_{x}(t, \omega)\right)-m\left(D_{1}\right) \chi_{D_{2}}\left(X_{x}(t, \omega)\right)\right] d t
$$

is unbounded as $T \rightarrow \infty$. However with probability one,

$$
\frac{f(T, \omega)}{\int_{0}^{T} \chi_{D_{i}}\left(X_{x}(t, \omega)\right) d t}=O\left(\left(\frac{\log \log N}{N}\right)^{1 / 2}\right) \rightarrow 0 \quad \text { as } \quad T \rightarrow \infty
$$

where $X_{x}(t, \omega)$ is Brownian motion in $\boldsymbol{R}^{2}$ starting from $x$ and $(N-1)$ is the number of new entries to $D_{1}$ after hitting $D_{2}$, up to time $T$.

Proof. First we obtain uniform upper and lower bounds for $P_{x}$ \{total time spent in $D_{1}$ before hitting $\left.D_{2}>t\right\}$ for all $x \in \bar{D}_{1}$. Let $A, B$ be open sets such that $\bar{A} \subset D_{1}$ and $\bar{D}_{1} \subset B$ such that $D_{2} \subset(\bar{B})^{c}=D$ say. Consider an open subset $A_{1}$ of $A$ such that $d\left(A_{1}, \partial A\right)>0$. Then

$$
P_{x}\left(\sigma_{A}<\sigma_{D}\right) \geqq P_{x}\left(\sigma_{A_{1}}<\sigma_{D}\right) \quad \text { for } \quad x \in \bar{D}_{1} \cap A^{c} ;
$$

where for any Borel set $B$,

$$
\sigma_{B}(\omega)=\sigma_{B}=\left\{\begin{array}{l}
\inf \left(t>0: X_{t} \in B\right) \\
+\infty \text { otherwise }
\end{array}\right.
$$

Since $\sigma_{D}$ is the limit of a monotone increasing sequence of non-negative simple functions $f_{n}$ say, and by Proposition 2.1 of (9)
$P_{x}\left(\sigma_{A_{1}} \leqq f_{n}\right)$ is lower semi-continuous in $x$ for each $n$, it follows that $P_{x}\left(\sigma_{A_{1}}<\sigma_{D}\right)$ is bounded below and assumes its minimum for $x \in \bar{D}_{1} \cap A^{c}$. But $P_{x}\left(\sigma_{A_{1}}<\sigma_{D}\right)>0$, so that $\min _{x} \dot{P_{x}}\left(\sigma_{A_{1}}<\sigma_{D}\right)=C^{\prime}>0$ for $x \in \bar{D}_{1} \cap A^{c}$. Observe that $P_{x}\left(\sigma_{A}<\sigma_{D}\right)=1$ for $x \in A$ and we have proved
there exists a constant $C^{\prime}>0$ such that $P_{x}\left(\sigma_{A}<\sigma_{D}\right)>C^{\prime}$ for all $x \in \bar{D}_{1}$.
Starting from $y \in A$, define

$$
\begin{aligned}
\tau_{\partial D_{1}} & =\text { first passage time out of } D_{1}, \\
\beta & =d\left(y, D_{1}^{c}\right) \\
S(y, \beta) & =\text { circle centre } y \text { and radius } \beta .
\end{aligned}
$$

Then $P_{y}\left\{\tau_{\partial D 1}>t\right\} \geqq P_{y}\left\{\tau_{\partial S(y, \beta)}>t\right\}$. Moreover $P_{y}\left\{\tau_{\partial S(y, \beta)}>t\right\}>C_{2} e^{-c_{1} t}$ by Theorem 2 of (2). Therefore

$$
\begin{equation*}
P_{y}\left\{\text { total time spent in } D_{1} \text { before leaving } B>t\right\} \geqq C_{2} e^{-c_{1} t} ; y \in \bar{A} . \tag{1.2}
\end{equation*}
$$

Next define

$$
Q_{x}=\text { total time spent in } D_{1} \text { before hitting } D_{2} \text {, starting from } x \in \bar{D}_{1}
$$

$R_{x}=\left\{\begin{array}{l}0 \text { if } X_{x}(t) \text { does not hit } A \text { before } \partial B, \\ 1 \text { if } X_{x}(t) \text { hits } A \text { at time } \mu \text { say before hitting } \partial B .\end{array}\right.$
Then $Q_{x} \geqq R_{x} . Q_{X(\mu)}$, where $X(\mu) \in \bar{A}$. By the strong Markov property, (1.1) and (1.2) we obtain the lower bound of

Lemma 1.1. There exist positive constants $c, C^{\prime}, C_{0}, c_{1}$ such that $C_{0} e^{-c_{1} t} \leqq P_{x}\{$ total time spent in $D_{1}$ before hitting $\left.D_{2}>t\right\} \leqq C^{\prime} e^{-c t}$, for all $x \in \bar{D}_{1}$.

For the upper bound let $C_{1}=\inf \left\{P_{x}\left(X(1) \in D_{2}\right): x \in \bar{D}_{1}\right\}$. Then $C_{1}>0$. Define recursively

$$
\tau_{0}=0, \tau_{n}=\text { first hitting time of } \bar{D}_{1} \text { after } \tau_{n-1}+1
$$

Then for any positive integer $n$ and $x \in \bar{D}_{1}, P_{x}$ \{total time spent in $D_{1}$ before hitting $\left.D_{2}>n\right\} \leqq P_{x}\left\{X\left(\tau_{j}+1\right) \notin D_{2}, j=0, \ldots, n-1\right\} \leqq\left(1-C_{1}\right)^{n}<1$.

$$
\begin{aligned}
& \therefore P_{x}\left\{\text { total time spent in } D_{1} \text { before hitting } D_{2}>t\right\} \\
& \leqq P_{x}\left\{\text { total time spent in } D_{1} \text { before hitting } D_{2}>[t]\right\} \\
& \leqq\left(1-C_{1}\right)^{[t]}<1
\end{aligned}
$$

Therefore for suitable constants $C^{\prime}, c$,

$$
P_{x}\left\{\text { total time spent in } D_{1} \text { before hitting } D_{2}>t\right\} \leqq C^{\prime} e^{-c t} \text {, for all } x \in \bar{D}_{1}
$$

This completes the proof of Lemma 1.1.
Let $\alpha_{1}(\omega)=\sigma_{D_{1}}(\omega), \beta_{1}(\omega)=\alpha_{1}(\omega)+\sigma_{D_{2}}\left(\omega_{\alpha_{1}}^{+}\right)$and for $n \geqq 2$,

$$
\alpha_{n}(\omega)=\beta_{n-1}(\omega)+\alpha_{1}\left(\omega_{\beta_{n-1}}^{+}\right), \beta_{n}(\omega)=\alpha_{n}(\omega)+\alpha_{D_{2}}\left(\omega_{\alpha_{n}}^{+}\right)
$$

where $\omega_{t}=X(t, \omega)$ is standard separable Brownian motion, and $\omega_{s}^{+}$is the shifted Brownian motion $\omega_{s}^{+}: t \rightarrow X(t+s, \omega)$. Using our notation we may rewrite, for $x \in D_{1}$,

$$
\int_{0}^{T} \chi_{D 1}\left(X_{x}(t, \omega)\right) d t \text { as } \sum_{i=1}^{n} \int_{\alpha_{i}}^{\beta_{i}} \chi_{D_{1}}\left(X_{x}(t, \omega)\right) d t+R_{D_{1}}^{n}
$$

and

$$
\begin{equation*}
\int_{0}^{T} \chi_{D_{2}}\left(X_{x}(t, \omega)\right) d t \text { as } \sum_{i=1}^{n} \int_{\beta_{i}}^{\alpha_{i+1}} \chi_{D_{2}}\left(X_{x}(t, \omega) d t+R_{D_{2}}^{n}\right. \tag{1.3}
\end{equation*}
$$

where, for

$$
\beta_{n} \leqq T<\alpha_{n+1}, R_{D_{1}}^{n}=0 \text { and } 0 \leqq R_{D_{2}}^{n}<\int_{\beta_{n}}^{\alpha_{n+1}} \chi_{D_{2}}\left(X_{x}(t, \omega)\right) d t
$$

for

$$
\alpha_{n+1} \leqq T<\beta_{n}, R_{D_{2}}^{n}=0 \text { and } 0 \leqq R_{D_{1}}^{n}<\int_{\alpha_{n+1}}^{\beta_{n+1}} \chi_{D_{1}}\left(X_{x}(t, \omega)\right) d t
$$

If $x \notin D_{1}$, there will be an extra initial term in (1.3) which will make no difference to our results. Let

$$
U_{i}=\int_{\alpha_{i}}^{\beta_{i}} X_{D_{1}}\left(X_{x}(t, \omega)\right) d t, \quad V_{i}=\int_{\beta_{i}}^{\alpha_{i+1}} \chi_{D_{2}}\left(X_{x}(t, \omega)\right) d t .
$$

Then $P_{x}\left\{U_{i}>t\right\} \geqq C_{0} e^{-c_{1} t}$ for all $x \in \bar{D}_{1}$ by Lemma 1.1. Therefore for every fixed $\lambda$ no matter how large,

$$
P_{x}\left\{U_{i}<\lambda\right\}<\delta<1 \text { for all } x \in \bar{D}_{1}
$$

By the strong Markov property and the fact that the bounds in Lemma 1.1 are uniform in $x \in \bar{D}_{1}$, we have

$$
P\left\{U_{1}<\lambda, U_{2}<\lambda, \ldots, U_{n}<\lambda\right\}<\delta^{n} \rightarrow 0 \text { as } n \rightarrow \infty .
$$

Hence for any fixed $\lambda$, there is zero probability that all the $U_{i}$ are less than $\lambda$. Since

$$
f(T, \omega)=m\left(D_{2}\right)\left(\sum_{i=1}^{n} U_{i}+R_{D_{1}}^{n}\right)-m\left(D_{1}\right)\left(\sum_{i=1}^{n} V_{i}+R_{D_{2}}^{n}\right)
$$

it is easy to see that, with probability one, $f(T, \omega)>\lambda$ infinitely-often as $n \rightarrow \infty$. Now consider a sequence $\lambda_{n}(=n) \uparrow \infty$ as $n \rightarrow \infty$ and define the event $E_{n}=\left\{\omega: f(T, \omega)>\lambda_{n}(=n)\right.$ infinitely-often $\}$. Then $P\left\{E_{n}\right\}=1$ and since $E_{n}$ is monotone in $\left.n, \underset{n}{ } \underset{n}{\cap} E_{n}\right\}=1$. Therefore for $\omega \in \cap_{n} E_{n}, f(T, \omega)$ is unbounded as $n \rightarrow \infty$. Hence with probability one, $f(T, \omega)$ is unbounded as $T \rightarrow \infty$. For proof of the second part of Theorem 1.2 we require the machinery of stationary Markov chains. First we state a useful result.

Lemma 1.2. Given $X(0)=x$, let $X\left(\alpha_{i}\right)=a_{i}$. Then $\left\{a_{i}\right\}$ is a strictly stationary Markov chain on $\bar{D}_{1}$. Moreover there exists, for $S \subset \bar{D}_{1}$, the unique stationary distribution $p(S)$, for the Markov chain $\left\{a_{i}\right\}$, such that

$$
\left|p^{i}(x, S)-p(S)\right|<c \rho^{i}, x \in \bar{D}_{1}
$$

for constants $c>0,0<\rho<1$; where $p^{i}(x, S)$ is the i-step transition probability for the Markov chain.
For the proof of Lemma 1.2 see Proposition 4.1 in (8) and $\S 5$ in Chapter V of (4).
Precise determination of the stationary distribution for generalised sets $D_{1}, D_{2}$ is difficult. We did obtain $p(S)$ in two special cases (stated below as Remarks 1 and 2) using standard potential theory arguments. Note that similar arguments yield the stationary distribution on any two bounded circles in $\boldsymbol{R}^{2}$.

Remark 1. Consider two unit discs $A$ and $B$ in $\boldsymbol{R}^{2}$ such that the distance between their centres is $s>2$ units. Then
$p(\psi)=$ stationary probability that Brownian motion enters $B$ through $d \psi \subset \partial B$
and
$p(\phi)=$ stationary $P\{$ Brownian motion enters $A$ through $d \phi \subset \partial A\}$ are given by

$$
p(\psi)=\frac{1-u^{2}}{2 \pi\left(1-2 u \operatorname{Cos} \psi+u^{2}\right)}, p(\phi)=\frac{1-u^{2}}{2 \pi\left(1-2 u \operatorname{Cos} \phi+u^{2}\right)}
$$

where $\psi, \phi$ are the angles between the line joining the centres of the circles and $d \psi, d \phi$ respectively and $u=\frac{1}{2} s-\left(\frac{1}{4} s^{2}-1\right)^{1 / 2}$.

Remark 2. Consider a unit disc $A$ and a circle $B$ with radius $b>1$ unit in $\boldsymbol{R}^{2}$ such that the distance between their centres $s>1+b$. Then
$f(\phi, 1)=$ probability density of stationary distribution on $\partial A$, and
$f(\psi, b)=$ probability density of stationary distribution on $\partial B$
are given by

$$
f(\phi, 1)=\frac{1-U_{A}^{2}}{2 \pi\left(1-2 U_{A} \operatorname{Cos} \phi+U_{A}^{2}\right)}, f(\psi, b)=\frac{b^{2}-U_{B}^{2}}{2 \pi b\left(b^{2}-2 b U_{B} \operatorname{Cos} \psi+U_{B}^{2}\right)}
$$

where

$$
U_{A}=\frac{1+s^{2}-b^{2}}{2 s}-\left(\frac{\left(1+s^{2}-b^{2}\right)^{2}}{4 s^{2}}-1\right)^{1 / 2}, U_{B}=\frac{s^{2}+b^{2}-1}{2 s}-\left(\frac{\left(s^{2}+b^{2}-1\right)^{2}}{4 s^{2}}-b^{2}\right)^{1 / 2}
$$

and $\phi, \psi$ are as in Lemma 1.3.
Next define $\Gamma_{x}(t)=P_{x}\left\{U_{1} \leqq t\right\}, \Gamma_{x}^{i}(t)=P_{x}\left\{U_{i} \leqq t\right\}$. Then

$$
\Gamma_{z}^{i}(t)=P_{z}\left\{U_{i} \leqq t\right\}=\int_{\partial D_{1}} \Gamma_{x}(t) p^{i-1}(z, d x)
$$

and since $\left\{U_{i}\right\}$ inherits stationarity from $\left\{a_{i}\right\}$,
$\Gamma(t)=\int_{\partial D_{1}} \Gamma_{x}(t) p(d x)=\lim _{i \rightarrow \infty} \Gamma_{z}^{i}(t)$ is the asymptotic distribution of $U_{i}$ which does not depend on $i$.

Definition. A strictly stationary sequence $\left\{U_{j}\right\}$ is said to be uniformly mixing if for all $D \in M_{k+n}^{\infty}$,

$$
\left|P\left\{D \mid M_{-\infty}^{k}\right\}-P(D)\right| \leqq \phi(n) \downarrow 0 \text { as } n \rightarrow \infty ;
$$

where the $\sigma$-algebra $M_{k+n}^{\infty}$ describes the future of the sequence $\left\{U_{j}\right\}$ and is generated by $\left\{U_{k+n}, U_{k+n+1}, U_{k+n+2}, \ldots\right\}$ while the $\sigma$-algebra $M_{-\infty}^{k}$ is generated by $\left\{U_{1}, U_{2}, \ldots, U_{k}\right\}$, and $\phi(n)$ is said to be the mixing coefficient.

Lemma 1.3. The sequence $\left\{U_{i}\right\}$ is uniformly mixing.
Proof. Let $X$ be a measure space and let $\mu, v$ be two measures on $X$ such that $|\mu(S)-v(S)|<\varepsilon$ for all $S \subset X$. Let $0 \leqq f(x) \leqq 1$ be a function measurable on $X$. Then $f$ is the limit of a monotone increasing sequence of non-negative simple functions, so that application of Lebesgue's theorem (see page 121 of (6)) to this sequence gives ( $\int f d \mu-$ $\left.\int f d v\right)<\varepsilon$ and $-\varepsilon<\left(\int f d \mu-\int f d v\right)$ separately and hence the result $\left|\int f d \mu(x)-\int f d v(x)\right|<\varepsilon$. Now $\partial D_{1}$ is a measure space on which two measures $p^{i}(x, S), p(S)$ are such that $\left|p^{i}(x, S)-p(S)\right|<c \rho^{i}$ for all $S \subset \partial D_{1}$. Therefore

$$
\left|\int_{\partial D_{1}} \Gamma_{z}(t) p^{i}(x, d z)-\int_{\partial D_{1}} \Gamma_{z}(t) p(d z)\right|<c \rho^{i}
$$

for constants $c>0,0<\rho<1$.
Define $\theta$ as the shift function $\theta\left(U_{1}, U_{2}, \ldots\right)=\left(U_{2}, U_{3}, \ldots\right)$. Then for $D \in M_{k+n}^{\infty}$, $\theta^{-n}(D)$ depends on at most $U_{k}, U_{k+1}, U_{k+2}, \ldots$, and $\theta^{-n-k+1}(D)$ depends on at most $U_{1}$, $U_{2}, \ldots$. By the Markov property of $\left\{a_{i}\right\}$,

$$
\left\{C_{1} \leqq P\left(D \mid a_{k}=z\right) \leqq C_{2}\right\} \quad \Rightarrow \quad\left\{C_{1} \leqq P\left(D \mid M_{-\infty}^{k}\right) \leqq C_{2}\right\} \text { a.s. }
$$

Moreover

$$
P\left(D \mid a_{k}=z\right)=\int_{\partial D_{1}} P\left\{\theta^{-n-k+1}(D) \mid a_{1}=y\right\} p^{n}(z, d y)
$$

by strict stationarity and $P\left\{\theta^{-n-k+1}(D) \mid a_{1}=y\right\}$ is a fixed function of $y$ since it depends neither on $n$ nor $k$. Therefore $P\left\{D \mid a_{k}=z\right\}$ depends on $z$ but not on $k$. If we restart the
sequence at $a_{k+1}=z_{1}$ say, $D$ now depends on a sequence starting from $z_{1}$ and so depends on $U_{n+1}, U_{n+2}, \ldots$.Therefore replacing $k$ by appropriate suffix corresponding to $z_{1}$ gives

$$
\left|P\left\{D \mid M_{-\infty}^{k}\right\}-P(D)\right|=\left|P_{z_{1}}(D)-P(D)\right|=\left|P_{z_{1}}(D)-P_{X\left(\alpha_{1}\right)}(D)\right| .
$$

It is easy to show that $\left|P_{z_{1}}(D)-P_{\mathrm{L}}(D)\right|<c^{\prime} \rho^{n}, 0<\rho<1$, where

$$
P_{L}(D)=\lim _{n \rightarrow \infty} \int_{\partial D_{1}} P_{y_{1}}(D) p^{n}\left(z_{1}, d y_{1}\right)=\int_{\partial D_{1}} P_{y_{1}}(D) p\left(d y_{1}\right) .
$$

Similarly $\left|P_{L}(D)-P_{X\left(\alpha_{1}\right)}(D)\right|<c_{1} \rho^{n+k}, 0<\rho<1$. Therefore

$$
\begin{equation*}
\left|P\left\{D \mid M_{-\infty}^{k}\right\}-P(D)\right|<c \rho^{n}, \text { for constants } c>0,0<\rho<1 ; \tag{1.4}
\end{equation*}
$$

which completes the proof of Lemma 1.3.
Corollary. The sequence $\left\{V_{i}\right\}$ is uniformly mixing.
Now define $Y_{i}=m\left(D_{2}\right) U_{i}-m\left(D_{1}\right) V_{i}$.
It is easy to see, using the ergodic theorem for stationary processes (see e.g. (6)) that $E\left(Y_{i}\right) \rightarrow 0$ as $i \rightarrow \infty$.

Also Lemma 1.1 immediately gives

$$
\begin{equation*}
P_{x}\left\{Y_{i}>t\right\} \leqq c e^{-c_{1} t} \text { for all } x \in \bar{D}_{1} \tag{1.5}
\end{equation*}
$$

Applying the same method as for $\left\{U_{i}\right\}$ we arrive at
Lemma 1.4. The sequence $\left\{Y_{i}-E_{x}\left(Y_{i}\right)\right\}$ is strictly stationary and uniformly mixing, with mixing còefficient given by (1.4).
Next we state two useful results.
Lemma 1.5. A strictly stationary sequence $\left\{U_{i}\right\}$ with $E\left(U_{j}\right)=0$, satisfying the uniform mixing condition, obeys the law of the iterated logarithm if the following conditions are fulfilled:
(i) $E\left|U_{j}\right|^{2+\delta}<\infty, \delta>0$;
(ii) $\sum_{n=1}^{\infty}\{\phi(n)\}^{1 / 2}<\infty$, where $\phi(n)$ is the mixing coefficient;
(iii) $0 \neq \sigma^{2}=E\left[U_{1}^{2}\right]+2 \sum_{j=2}^{\infty} E\left[U_{1} \cap U_{j}\right] . S e e(10)$.

Lemma 1.6. Suppose the strictly stationary sequence $\left\{U_{j}\right\}$ satisfies the uniform mixing condition. If the random variables $\tau, \eta$ are measurable with respect to $M_{-\infty}^{k}$ and $M_{k+n}^{\infty}$ respectively, and if $E\left(|\tau|^{p}\right)<\infty, E\left(|\eta|^{q}\right)<\infty$ with $p, q>1$ and $1 / p+1 / q=1$, then

$$
|E(\tau \eta)-E(\tau) E(\eta)| \leqq 2\{\phi(n)\}^{1 / p} E^{1 / p}(|\tau| p) E^{1 / q}\left(|\eta|^{q}\right)
$$

where $\phi(n)$ is the mixing coefficient for $\left\{U_{i}\right\}$. See (5).
That conditions (i) and (ii) of Lemma 1.5 are satisfied by $\left\{Y_{i}-E_{x}\left(Y_{i}\right)\right\}$ follows from (1.5) and (1.4) above respectively. Moreover, by the strong Markov property and the uniformity
of the bounds in Lemma 1.1, $P_{x}\left\{Y_{i}>a\right\} \geqq C_{0} e^{-c a}\left(1-C^{\prime} e^{-c_{2} a}\right)$ for $a>0, x \in \bar{D}_{1}$. Therefore for sufficiently large $a, P_{x}\left\{Y_{i}>a\right\} \geqq \varepsilon^{\prime}>0$ for all $i, x \in \bar{D}_{1}$. Similarly, $P_{x}\left\{Y_{i}<-a\right\} \geqq$ $\varepsilon_{1}>0$ for all $i, x \in \bar{D}_{1}$. Therefore there is an $\varepsilon\left(=a^{2}\left(\varepsilon^{\prime}+\varepsilon_{1}\right)\right)$ such that for all integers $i$, the variance $\sigma_{x}^{2}\left(Y_{i}\right)$ of $Y_{i}$ starting at $x$ is at least $\varepsilon$. Hence the limiting variance of $Y_{i}$ starting at $x, \int \sigma_{x}^{2} p(d x)>0$. Conditions (i) and (ii) of Lemma 1.5 are clearly satisfied for the sequence. obtained by requiring that the initial point be random with distribution $p$, from (1.5) and (1.4) above respectively. For condition (iii) we need Lemma 1.6. The conditions of Lemma 1.6 are șatisfied by $Y_{j}-E_{x}\left(Y_{j}\right), j \geqq 2$ and $Y_{1}-E_{x}\left(Y_{1}\right)$ for $p=q=2$. Therefore

$$
\begin{aligned}
2\left|\sum_{j=2}^{\infty} E_{x}\left(\left[Y_{1}-E_{x}\left(Y_{1}\right)\right]\left[Y_{j}-E_{x}\left(Y_{j}\right)\right]\right)\right| \leqq 4 & \sum_{j=2}^{\infty}\{\phi(j-1)\}^{1 / 2} E_{x}\left(\left[Y_{1}-E_{x}\left(Y_{1}\right)\right]^{2}\right) \\
& \leqq \frac{4(c \rho)^{1 / 2}}{1-\rho^{1 / 2}} E_{x}\left(\left[Y_{1}-E_{x}\left(Y_{1}\right)\right]^{2}\right) ; c>0,0<\rho<1
\end{aligned}
$$

Since $\int \sigma_{x}^{2} p(d x)$ is positive, condition (iii) is satisfied if

$$
\begin{equation*}
2\left|\sum_{j=2}^{\infty} E_{x}\left(\left[Y_{1}-E_{x}\left(Y_{1}\right)\right]\left[Y_{j}-E_{x}\left(Y_{j}\right)\right]\right)\right|<E_{x}\left(\left[Y_{1}-E_{x}\left(Y_{1}\right)\right]^{2}\right) \tag{1.6}
\end{equation*}
$$

so that (1.6) holds if

$$
\begin{equation*}
\rho<1 /\left(1+4 c^{1 / 2}\right)^{2} ; c>0,0<\rho<1 \tag{1.7}
\end{equation*}
$$

Since the exact values of $c$ and $\rho$ are unknown it is not possible to determine whether or not (1.7) is satisfied. A way out of this difficulty is to consider the subsequences $\left\{Y_{k i+j}\right\}$ $\left.E_{x}\left(Y_{k i+j}\right)\right\}, i=0,1,2, \ldots$; for a FIXED integer $k$ whose value will be determined later. There are $k$ such subsequences of $\left\{Y_{i}-E_{x}\left(Y_{i}\right)\right\}$ namely

$$
\left\{Y_{j}-E_{x}\left(Y_{j}\right)\right\},\left\{Y_{k+j}-E_{x}\left(Y_{k+j}\right)\right\},\left\{Y_{2 k+j}-E_{x}\left(Y_{2 k+j}\right)\right\}, \ldots \text { for } j=1,2, \ldots k ;
$$

such that

$$
\sum_{j=1}^{k} \sum_{i=0}^{n-1}\left(Y_{k i+j}-E_{x}\left(Y_{k i+j}\right)\right)=\sum_{i=1}^{n k}\left(Y_{i}-E_{x}\left(Y_{i}\right)\right) .
$$

For a typical subsequence, condition (iii) of Lemma 1.5 is satisfied if

$$
\begin{equation*}
\rho^{k}<1 /\left(1+4 c^{1 / 2}\right)^{2} \tag{1.8}
\end{equation*}
$$

Because $c, \rho$ are fixed for $\left\{Y_{i}-E_{x}\left(Y_{i}\right)\right\}$ we can choose $k$ such that $k$ is the smallest integer for which (1.8) holds. For this value of $k$, it is clear that all three conditions of Lemma 1.5 hold. Therefore the sequence $\left\{Y_{k i+j}-E_{x}\left(Y_{k i+j}\right)\right\}$ for $k$ given as above, obeys the law of the iterated logarithm. The tail of the distribution of $Y_{i}$ starting at $x$ has a negative exponential upper bound (application of Lemma 1.1). Moreover $E_{x}\left\{Y_{i}\right\} \rightarrow 0$ as $i \rightarrow \infty$. Theorem 1.1 therefore implies

$$
\begin{equation*}
\left|\sum_{i=1}^{N} Y_{i}\right|=\left|\sum_{i=1}^{n k} Y_{i}\right| \leqq \sum_{j=1}^{k}\left|\sum_{i=0}^{n-1} Y_{k i+j}\right|<k\left(c N \sigma^{2} \log \log N\right)^{1 / 2} \tag{1.9}
\end{equation*}
$$

for sufficiently large N , where $\sigma^{2}=\max \left\{\sigma_{j}^{2} ; j=1,2, \ldots, k\right\}$ where $\sigma_{j}^{2}$ is equivalent to $\sigma^{2}$ in Lemma 1.5 for each subsequence $\left\{Y_{k i+j}-E_{x}\left(Y_{k i+j}\right)\right\}$. If
(a) $\beta_{N} \leqq T<\alpha_{N+1}, f(T, \omega)=\sum_{i=1}^{N} Y_{i}-m\left(D_{1}\right) R_{D_{2}}^{N}, 0 \leqq R_{D_{2}}^{N}<V_{N}$,
(b) $\alpha_{N+1} \leqq T<\beta_{N+1}, f(T, \omega)=\sum_{i=1}^{N} Y_{i}+m\left(D_{2}\right) R_{D_{1}}^{N}, 0 \leqq R_{D_{1}}^{N}<U_{N+1}$,
where ( $N-1$ ) is the number of new entries to $D_{1}$ from $D_{2}$ up to time $T$. Since the tail of the distribution of $m\left(D_{2}\right) U_{N+1}$ has a negative exponential upper bound, there exists an $N_{0}$ such that $m\left(D_{2}\right) U_{N+1} \leqq N^{1 / 2}$ for all $N \geqq N_{0}$. Then in both cases (a) and (b) above we have, from (1.9), that

$$
\begin{equation*}
|f(T, \omega)|=O\left((N \operatorname{loglog} N)^{1 / 2}\right) \text { a.s. as } N \rightarrow \infty . \tag{1.10}
\end{equation*}
$$

But

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \int_{0}^{T} \chi_{D_{i}}\left(X_{x}(t, \omega)\right) d t=c^{\prime} \text { a.s. }
$$

(see e.g. (8)). Therefore with probability one,

$$
\frac{f(T, \omega)}{\int_{0}^{T} \chi_{D_{i}}\left(X_{x}(t, \omega)\right) d t}=O\left(\left(\frac{\log \log N}{N}\right)^{1 / 2}\right) \rightarrow 0 \text { as } N \rightarrow \infty
$$

where ( $N-1$ ) is the number of new entries to $D_{1}$ after hitting $D_{2}$. The fact that $N \rightarrow \infty$ as $T \rightarrow \infty$ completes the proof of Theorem 1.2.

Remark 3. The same result in $\boldsymbol{R}^{1}$ is substantially easier because the hitting point of the interval is unique which implies that the sequence $\left\{Y_{i}\right\}$ of random variables are independent and identically distributed. This allows an application of the standard law of iterated logarithm.

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