# The Lifetime of Conditioned Brownian Motion 

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## 1. Introduction

In this paper we show that Brownian $h$-paths, $h>0$ harmonic, have finite lifetimes in plane domains of finite area. The analogous result for bounded domains in higher dimensions is false - we give an example of a bounded domain in $\mathbb{R}^{3}$ and a positive harmonic function for which the $h$-paths have infinite lifetime almost surely. This difference in behavior is related to the scaling property of Brownian motion.

In what follows, $D$ will be a domain in $\mathbb{R}^{n}$ which has a Green function $G$. Denote by $\Delta$ the minimal Martin boundary of $D$ [7]. Let $\Omega$ be the space of all continuous (in Martin topology) functions $\omega:[0, \infty) \rightarrow \mathrm{D} \cup \Delta$ with the property that $\omega(s) \in \Delta$ implies $\omega(t) \in \Delta$ for $t \geqq s$. Let $X(t, \omega)=\omega(t)$ be the coordinate process. Denote by ( $\Omega, \mathscr{F}, \mathscr{F}_{t}, X(t), \theta_{t}, P_{x}$ ) the standard Brownian motion started at $x \in D$. If $p(t, x, y)$ is the transition density of Brownian motion killed on exiting $D$ and $h>0$ is harmonic in $D$, define

$$
p^{h}(t, x, y)=h(x)^{-1} p(t, x, y) h(y)
$$

and let $P_{x}^{h}$ denote the measure on $\Omega$ induced by $\mathrm{p}^{h}$. These are the $h$-paths of Doob [3].

Theorem 1. If $D$ is a domain in $\mathbb{R}^{2}, h>0$ is harmonic in $D$, and $\tau_{D}$ $=\inf \{t>0: X(t) \notin D\}$, then

$$
\begin{equation*}
E_{x}^{h}\left(\tau_{D}\right) \leqq c|D| \tag{1.1}
\end{equation*}
$$

where $c$ is an absolute constant and $|D|$ is the area of $D$.
Statement (1.1) implies that if $h$ is positive and harmonic on $D$, a plane domain of finite area, then

$$
\begin{equation*}
\lim _{t \rightarrow \infty} E_{x}\left[h(X(t)) ; t<\tau_{D}\right]=0 . \tag{1.2}
\end{equation*}
$$

[^0]The difficulty in proving this directly is that the supermartingale $h(X(t)) I\left(t<\tau_{D}\right)$ need not be uniformly integrable. If it were, then the constant times, $t$, could be replaced by any sequence $T_{n}$ of stopping times increasing to infinity. However, consider the following example: Let $D$ be the unit disc in $\mathbb{R}^{2}$ and $h$ the Poisson kernel with pole at ( 1,0 ). Let $D_{n}=\{z \in D: h(z)>n\}$. ( $D_{n}$ is the interior of a circle tangent to $D$ at $(1,0)$.) Let $T_{n}$ be the hitting time of $D_{n}$ by $X(t)$. Then $h\left(X\left(t \wedge T_{n} \wedge \tau_{D}\right)\right)$ is a uniformly integrable martingale, hence

$$
h(0)=E_{0}\left(h\left(X\left(T_{n} \wedge \tau_{D}\right)\right)\right)=E_{0}\left(h\left(X_{T_{n}}\right) ; T_{n}<\tau_{D}\right) .
$$

Thus (1.2) does not hold with $t$ replaced by $T_{n}$.
Another instructive example to consider is Littlewood's crocodile [5, p.268]. The lifetimes of $h$-paths are finite even when conditioned to go to the bad end of the crocodile. Nothing is gained in terms of delaying the process by throwing up obstacles; that is, in $\mathbb{R}^{2}$.

Finally, let us mention some alternative formulations and consequences of our result. The well-known relationship between Green functions and occupation times yields:
Corollary 1. Let $D \subseteq \mathbb{R}^{2}$ be a domain having a Green function $G(x, y)$. Then there is an absolute constant $c$ such that, for any function $h \geqq 0$ harmonic on $D$,

$$
\int_{D} G(x, y) h(y) d y \leqq c|D| h(x) .
$$

In particular, if $h=1$,

$$
\int_{D} G(x, y) d y \leqq c|D|
$$

Doob [3] shows that if Brownian motion is conditioned using the Green function of $D$ with pole inside $D$ then the conditioned paths have finite lifetimes. Combining this with Theorem 1 and the Riesz decomposition theorem produces:

Corollary 2. If $h>0$ is superharmonic on $D$ then $P_{x}^{h}\left(\tau_{D}<\infty\right)=1$ for any $x \in D$, whenever $D$ is a plane domain of finite area.

Our result also gives a refinement of a result of Lamb [6]. He shows that if $h \geqq 0$ is harmonic on a domain $D$ then $h$ may be decomposed uniquely as $h=h_{1}$ $+h_{2}+h_{3}$ where $h_{1}\left(X\left(t \wedge \tau_{D}\right)\right)$ is a nonnegative uniformly integrable martingale, $h_{2}\left(X\left(t \wedge \tau_{D}\right)\right)$ is a nonnegative martingale with limit 0 , and $h_{3}\left(X\left(t \wedge \tau_{D}\right)\right)$ is a nonnegative supermartingale whose expectation approaches 0 ; moreover $h_{2}$ paths have infinite lifetimes. Thus our result shows that $h_{2}$ vanishes for plane domains of finite area.

In Sect. 2 we give the proof of Theorem 1. In Sect. 3 we construct a bounded region $D$ in $\mathbb{R}^{3}$ and a positive harmonic function $h$ such that $P_{x}^{h}\left(\tau_{D}=\infty\right)=1$ for all $x$ in $D$.

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## 2. Proof of Theorem 1

We begin by establishing two lemmas.
Lemma 2.1. Let $D$ be an open subset of $\mathbb{R}^{2}$. There exists an absolute constant $c$ such that $E_{x}\left(\tau_{D}\right) \leqq c|D|$, for all $x \in D$.

Proof. Fix $x$ in $D$. We may assume $|D|=1$ since if $a>0$ and $a D=\{a y: y \in D\}$ then $|a D|=a^{2} D$ and $E_{a x}\left(\tau_{a D}\right)=a^{2} E_{x}\left(\tau_{D}\right)$. Thus it suffices to show that $E_{x}\left(\tau_{D}\right)$ is less than a constant independent of $D$. Setting $B=B(x, 1)=\{y:|x-y|<1\}$, we have

$$
\begin{aligned}
E_{x}\left(\tau_{D}\right) & =E_{x}\left(\tau_{D} ; \tau_{D}<\tau_{B}\right)+E_{x}\left(\tau_{D} ; \tau_{B} \leqq \tau_{D}\right) \\
& \leqq E_{x} \tau_{B}+E_{x}\left(E_{X\left(\tau_{B}\right)}\left(\tau_{D}\right) ; \tau_{B} \leqq \tau_{D}\right) \\
& \leqq E_{x} \tau_{B}+P_{x}\left(\tau_{B} \leqq \tau_{D}\right) \sup _{y \in D} E_{y}\left(\tau_{D}\right) .
\end{aligned}
$$

Once we have established $P_{x}\left(\tau_{B} \leqq \tau_{T}\right) \leqq p<1$ with $p$ independent of $D$ and $x$, we will have

$$
E_{x}\left(\tau_{D}\right) \leqq \frac{c}{1-p}, \quad \text { with } c=E_{x} \tau_{B}=1 / 2
$$

as desired. Take $r<1$, put $K=D^{c} \cap B(x, r)$ and observe that $|K| \geqq \pi r^{2}-1$. Then if $\lambda$ is the restriction of Lebesgue measure to $K$ and $G$ is the Green function for $B(x, 1)$ set $G^{\lambda}(z)=\int G(z, y) \lambda(d y)$. Then

$$
G^{\lambda}(x)=\int G(x, y) \lambda(d y) \leqq 2 \pi \int_{0}^{1} \log \left(\frac{1}{r}\right) r d r=\frac{\pi}{2}
$$

If $\omega_{z}(K)$ is the harmonic measure of $K$ relative to $B(x, 1)$ then $\omega_{z}(K) \geqq \frac{2}{\pi} G^{\lambda}(z)$. Hence

$$
\omega_{z}(K) \geqq \frac{2}{\pi} G^{\lambda}(x) \geqq \frac{2}{\pi} \log \left(\frac{1}{r}\right)|K| \geqq \frac{2}{\pi} \log \left(\frac{1}{r}\right)\left(\pi r^{2}-1\right) .
$$

Choosing any value of $r<1$ such that the last expression is positive completes the proof. (For a different proof see [2, p. 148].)
Lemma 2.2. Let $h$ be a nonzero minimal harmonic function on $D$. Then for any $x \in D, P_{x}^{h}\left(h\left(X\left(\tau_{D}-\right)=\infty\right)=1\right.$.

Proof. We show first that $h\left(X\left(\tau_{D}-\right)=0, P_{x}-\right.$ a.s. If this were not the case then it is not difficult to show, using the fact that two-dimensional Brownian motion does not hit points, that there exist disjoint Borel subsets $A_{1}$ and $A_{2}$ of the boundary of $D$, such that

$$
P_{x}\left(h\left(X\left(\tau_{D}-\right)\right) I\left(X\left(\tau_{D}\right) \in A_{i}\right)>0\right)>0, \quad i=1,2
$$

Let $g_{i}(x)=E_{x}\left(h\left(X\left(\tau_{D}-\right)\right) ; X\left(\tau_{D}\right) \in A_{i}\right)$. Then $g_{i}(x) \leqq h(x), X \in D$; but since $h$ is minimal, we then must have $g_{i}(x)=c_{i} h(x)$, for constants $c_{i}$, which is clearly impossible.

Next let $D(n)$ be a sequence of open subsets of $D$, whose interiors increase to $D$, and whose closures are compact subsets of $D$. For $M>0$ put $T_{M}$ $=\inf \{t \geqq 0: h(X(t)) \geqq M\}$. Then $h\left(X\left(T_{M} \wedge \tau_{D(n)}\right)\right)$ is a uniformly integrable martingale in $n$. Letting $n$ approach infinity we obtain

$$
h(x)=E_{x} h\left(X\left(T_{M} \wedge \tau_{D}\right)\right)=E_{x}\left(h\left(X\left(T_{M}\right)\right) ; T_{M}<\tau_{D}\right)
$$

This means that $P_{x}^{h}\left(T_{M}<\tau_{D}\right)=1$, and $h\left(X\left(t \wedge \tau_{D}\right)\right)$ is unbounded with $P_{x}^{h}$ - probability 1.

Finally, $h\left(X\left(t \wedge \tau_{D}\right)\right)$ has an infinite limit at $\tau_{D}-$ since $\left.\frac{1}{h(X(t))} I\left(t<\tau_{D}\right)\right)$ is a nonnegative $h$-supermartingale.

Let us now return to the proof of Theorem 1.
We assume without loss of generality that $h$ is minimal. Define

$$
\begin{array}{ll}
D_{n}=\left\{x \in D: 2^{n-1}<h(x)<2^{n+1}\right\}, & n=0, \pm 1, \pm 2, \ldots \\
C_{n}=\left\{x \in D: h(x)=2^{n}\right\}, & n=0, \pm 1, \pm 2, \ldots
\end{array}
$$

so that $D=\bigcup_{n=-\infty}^{\infty} D_{n}$ and $\sum_{n=-\infty}^{\infty}\left|D_{n}\right| \leqq 2|D|$. Consider the stopping times

$$
\begin{aligned}
R(n) & =\tau_{D_{n}}=\inf \left\{t>0: X(t) \notin D_{n}\right\} \\
S(n, 0) & =\inf \left\{t>0: X(t) \in C_{n}\right\} \\
T(n, 1) & =R(n) \circ \theta_{S(n, 0)}+S(n, 0)
\end{aligned}
$$

and if $S(n, 0), \ldots, S(n, k-1), T(n, 1), \ldots, T(n, k-1)$ have been defined, set

$$
\begin{aligned}
& T(n, k)= \begin{cases}R(n) \circ \theta_{S(n, k-1)}+S(n, k-1), & S(n, k-1)<\infty, \\
\infty, & S(n, k-1)=\infty .\end{cases} \\
& S(n, k)= \begin{cases}S(n, 0) \circ \theta_{T(n, k)}+T(n, k), & T(n, k)<\infty \\
\infty, & T(n, k)=\infty .\end{cases}
\end{aligned}
$$

Finally set

$$
N(n)=\inf \{k \geqq 0: S(n, k)=\infty\}=\inf \{k \geqq 0: T(n, k+1)=\infty\},
$$

and

$$
L(n, k)=T(n, k)-S(n, k-1)=R(n) \circ \theta_{S(n, k-1)} .
$$

In order to estimate $E_{x}^{h}(L(n, k-i) ; N(n)=k), i=1,2, \ldots, k-1$, we establish the following observations, each for $n=0,1,2, \ldots$ :

$$
\begin{align*}
& P_{x}^{h}(S(n, 0)<\infty)=1 / 2, x \in C_{n+1},  \tag{2.1}\\
& P_{x}^{h}(S(n, 0)<\infty)=1, x \in C_{n-1},  \tag{2.2}\\
& P_{x}^{h}\left(X(R(n)) \in C_{n-1}\right)=1 / 3, x \in C_{n},  \tag{2.3}\\
& P_{x}^{h}\left(X(R(n)) \in C_{n+1}\right)=2 / 3, x \in C_{n},  \tag{2.4}\\
& E_{x}^{h}\left(R(n) \mid X(R(n)) \in C_{n-1}\right) \leqq 3 E_{x}^{h}(R(n)), \quad \text { for } x \in C_{n}, \tag{2.5}
\end{align*}
$$

and

$$
\begin{equation*}
E_{x}^{h}\left(R(n) \mid X(R(n)) \in C_{n+1}\right) \leqq \frac{3}{2} E_{x}^{h}(R(n)), \quad x \in C_{n} \tag{2.6}
\end{equation*}
$$

For (2.1) we have

$$
P_{x}^{h}(S(n, 0)<\infty)=\frac{2^{n}}{h(x)}=\frac{1}{2} \quad \text { for } x \in C_{n+1} .
$$

Statement (2.2) follows at once from Lemma 2.2. For (2.3) we have

$$
P_{x}^{h}\left(X(R(n)) \in C_{n-1}\right)=\frac{2^{n+1}}{3}\left(\frac{1}{h(x)}-\frac{1}{2^{n+1}}\right)=\frac{1}{3}
$$

for $x \in C_{n}$. Of course (2.4) follows from (2.3).
Turning to (2.5), we have

$$
\begin{align*}
& E_{x}^{h}\left(R(n) \mid X(R(n)) \in C_{n-1}\right) \\
&=\int_{0}^{\infty} P_{x}^{h}\left(R(n)>\lambda \mid X\left(R_{n}\right) \in C_{n-1}\right) d \lambda \\
&=P_{x}^{h}\left(X(R(n)) \in C_{n-1}\right)^{-i} \int_{0}^{\infty} P_{x}^{h}\left(R(n)>\lambda, X(R(n)) \in C_{n-1}\right) d \lambda \\
&\left.\leqq 3 \int_{0}^{\infty} P_{x}^{h}(R(n)>\lambda) d \lambda \quad \quad \text { by }(2.3)\right)  \tag{2.3}\\
&=3 E_{x}^{h}(R(n)) .
\end{align*}
$$

Similarly (2.6) follows from (2.4). Now, for $i=1,2, \ldots, k-1$ we have

$$
\begin{aligned}
& E_{x}^{h}(L(n, k-i) ; N(n)=k) \\
& \quad=E_{x}^{h}\left(E_{x}^{h}(L(n, k-i) ; N(n)=k \mid X(T(n, k-i))) I(T(n, k-i)<\infty)\right) \\
& \quad=E_{x}^{h}\left(E _ { x } ^ { h } \left(L(n, k-i) ; N(n) \geqq k-i \mid X(T(n, k-i)) P_{x}^{h}(N(n)\right.\right. \\
& \quad=k \mid X(T(n, k-i))) I(T(n, k-i)<\infty))
\end{aligned}
$$

by the strong Markov property. Using (2.1)-(2.4)

$$
\begin{aligned}
P_{x}^{h}(N(n) & =k \mid X(T(n, k-i))) I(T(n, k-i)<\infty) \\
& = \begin{cases}\frac{1}{2}\left(\frac{2}{3}\right)^{i}, & X(T(n, k-i)) \in C_{n-1}, \\
\frac{1}{4}\left(\frac{2}{3}\right)^{i}, & X(T(n, k-i))) \in C_{n+1},\end{cases}
\end{aligned}
$$

Also,

$$
\begin{aligned}
& E_{x}^{h}\left(E_{x}^{h}(L(n, k-i), N(n) \geqq k-i \mid X(T(n, k-i)))\right. \\
& \quad \cdot\left.I\left(X(T(n, k-i)) \in C_{n-1}, T(n, k-i)<\infty\right)\right) \\
&= E_{x}^{h}\left(E_{x}^{h}\left(L(n, k-i) ; N(n) \geqq k-i \mid X(T(n, k-i)) \in C_{n-1}, T(n, k-i)<\infty\right)\right. \\
& \quad \cdot I\left(X\left(T(n, k-i) \in C_{n-1}, T(n, k-i)<\infty\right)\right) \\
&= E_{x}^{h}\left(E_{x}^{h}\left(R(n) \circ S(n, k-i-1) \mid X(T(n, k-i)) \in C_{n-1}, T(n, k-i)<\infty\right)\right.
\end{aligned}
$$

$$
\begin{aligned}
& \left.\cdot I\left(X(T(n, k-i)) \in C_{n-1}, N(n) \geqq k-i\right)\right) \\
= & E_{x}^{h}\left(E _ { x } ^ { h } \left(E_{X(S(n, k-i-1))}^{h}(R(n))\left|X\left(R(n) \in C_{n-1}\right)\right|\right.\right. \\
& \left.\cdot X(T(n, k-i)) \in C_{n-1}, T(n, k-i)<\infty\right) \\
& \left.\cdot I\left(X(T(n, k-i)) \in C_{n-1}, N(n) \geqq k-i\right)\right) \\
\leqq & \sup _{x \in C_{n}} E_{x}^{h}\left(R(n) \mid X(R(n)) \in C_{n-1}\right) P_{x}^{h}\left(N(n) \geqq k-i, X\left(T(n, k-i) \in C_{n-1}\right)\right. \\
\leqq & 3\left(\sup _{x \in C_{n}} E_{x}^{h}(R(n))\right)\left(\frac{2}{3}\right)^{k-i-1} \frac{1}{3}, \quad \text { by }(2.5), \\
= & \left(\frac{2}{3}\right)^{k-i-1} \sup _{x \in C_{n}} E_{x}^{h}(R(n)) .
\end{aligned}
$$

In an entirely similar manner we obtain

$$
\begin{aligned}
& E_{x}^{h}\left(E_{x}^{h}(L(n, k-i) ; N(n) \geqq k-i \mid X(T(n, k-i)))\right. \\
& \left.\quad \cdot I\left(X(T(n, k-i)) \in C_{n+1}, T(n, k-i)<\infty\right)\right) \\
& \quad \leqq\left(\frac{2}{3}\right)^{k-i} \sup _{x \in \mathcal{C}_{n}} E_{x}^{h}(R(n)) .
\end{aligned}
$$

Since

$$
\begin{align*}
\tau_{D} & \leqq \sum_{n=-\infty}^{\infty} \sum_{k=1}^{N(n)} L(n, k), \\
E_{x}^{h}\left(\tau_{D}\right) & \leqq \sum_{n=-\infty}^{\infty} E_{x}^{h} \sum_{k=1}^{N(n)} L(n, k)  \tag{2.7}\\
& =\sum_{n=-\infty}^{\infty} \sum_{k=1}^{\infty} E_{x}^{h}\left(\sum_{i=1}^{k} L(n, i) ; N(n)=k\right) \\
& \leqq\left(\sum_{k=1}^{\infty} k\left(\frac{2}{3}\right)^{k-1}\right) \sum_{n=-\infty}^{\infty} \sup _{x \in C_{n}} E_{x}^{h}(R(n)) \\
& =9 \sum_{n=-\infty}^{\infty} \sup _{x \in C_{n}} E_{x}^{h} R(n) .
\end{align*}
$$

Finally, for $x \in C_{n}$

$$
\begin{aligned}
E_{x}^{h}(R(n)) & \left.=\int_{0}^{\infty} P_{x}^{h}(R(n))>\lambda\right) d \lambda \\
& =h(x)^{-1} \int_{0}^{\infty} E_{x}(h(X(\lambda)) ; R(n)>\lambda) d \lambda \\
& \leqq 2 \int_{0}^{\infty} P_{x}(R(n)>\lambda) d \lambda \\
& =2 E_{x} R(n) \leqq 2 c\left|D_{n}\right| \quad \text { by Lemma } 2.1
\end{aligned}
$$

Thus,

$$
E_{x}^{h}\left(\tau_{D}\right) \leqq 18 c \sum_{n=-\infty}^{\infty}\left|D_{n}\right| \leqq 36 c|D|
$$

and the proof is complete.

## 3. An Example

The purpose of this section is to construct an example of a bounded domain $D$ in $\mathbb{R}^{3}$ and a harmonic function $h$ such that $\tau_{D}=\infty$ with $P_{x}^{h}$-probability 1 for any $x$ in $D$. It is necessary to recall some facts concerning the last exit decomposition of Brownian motion (see [4, 8, or 9].) Let $K$ be a Borel subset of $D$ and $\mu$ the last exit time from $K$ of Brownian motion $X(t)$ conditioned with some harmonic function $h$. Then the process $X(\mu+t)$ has continuous paths on $[0, \infty)$ and is strong Markov on $(0, \infty)$ with respect to an appropriate sequence of $\sigma$-fields; moreover, the semigroup on $(0, \infty)$ is that of $X(t)$ conditioned not to hit $K$.

We will also need the following simple result.
Lemma 3.1. Let $B$ denote the ball in $\mathbb{R}^{3}$ centered at 0 with radius 1 , and let $0<\delta<\frac{1}{4}$ be given. There exist absolute constants $c_{1}$ and $c_{2}$ such that

$$
P_{x}^{g}\left(\tau_{B}>c_{1}\right)>c_{2}, \quad|x|<\delta,
$$

for any positive harmonic function $g$.
Proof. By Harnack's inequality we may choose $c_{2}$, so small that $\frac{g(y)}{g(x)} \geqq 2 c_{2}$ for $|y|<2 \delta,|x|<2 \delta$. Define $\tau_{2 \delta}$ by $\tau_{2 \delta}=\inf \{t:|x(t)-x(0)|>2 \delta\}$. We may choose $c_{1}$ so small that $P_{x}\left(\tau_{2 \delta}<c_{1}\right)>\frac{1}{2}$ for $|x|<\delta$. Then for all $x$

$$
\begin{aligned}
P_{x}^{g}\left(\tau_{B}>c_{1}\right) & =g(x)^{-1} E_{x}\left(g\left(X\left(c_{1}\right)\right) ; \tau_{B}>c_{1}\right) \\
& \geqq g(x)^{-1} E_{x}\left(g\left(X\left(c_{1}\right)\right) ; \tau_{2 \delta}>c_{1}\right) \\
& \geqq 2 c_{2} P_{x}\left(\tau_{2 \delta}>c_{1}\right), \quad>c_{2},
\end{aligned}
$$

and the proof is complete.
Note that the result holds for $B$ a ball of radius $\eta$ if $\delta$ is replaced by $\delta \eta$ and $c_{1}$ by $\eta^{2} c_{1}$.

Consider for $D$ the following region in $\mathbb{R}^{3}$ : For $n=2,3, \ldots$ place $n^{6}$ nonintersecting balls of radius $\frac{1}{2} n^{-3}$ in a square array with centers on $\left\{\left(x, y, n^{-2}\right)\right.$; $0 \leqq x \leqq 3,0<y \leqq 3\}$. Thus we have infinitely many levels with balls on level $n$ arranged in $n^{3}$ rows (of $n^{3}$ balls each) parallel to the $y$ axis. Order the balls on the $n$-th level lexicographically according to the $(x, y)$ coordinates of their centers and denote them by $B_{n}^{k}, k=1,2, \ldots, n^{6}$. Now on each level connect each ball to the next in order by a thin cylindrical tube. On all levels but the first connect similarly the first ball to the last ball on the level directly above and on all levels connect the last ball to the first ball on the level directly below. Use smooth tubes of finite total volume. For definiteness let us make the following additional assumptions about the connecting tubes: each tube intersects its balls in circular caps subtending solid angle $\varepsilon<\frac{\pi}{2}$, and the two caps in each ball (except the first) are antipodal. That completes the description of $D$.

Some subsequence of the sequence of centers of the balls comprising $D$ is fundamental in the sense of Martin, and hence converges to a certain Martin boundary point $\xi$. The corresponding harmonic function $K(\cdot, \xi)$ may or may
not be minimal, but in any case there is some minimal function $h$ corresponding to a point in the support of the canonical representing measure of $K(\cdot, \xi)$. The function $h$ vanishes identically on the Euclidean boundary of $D$. Therefore, by Lemma 2.2, the $h$-paths must thread their way through the entire maze of $D$. Fix $x_{0}$ in $B_{2}^{1}$. Since $h$ is minimal $P_{x_{0}}^{h}\left(\tau_{D}=\infty\right)$ must be either 0 or 1 . We show it is 1 .

Define $T(n, k)$ to be the time between the last exit, $\mu(n, k)$, from the tube joining $B_{n}^{k-1}$ to $B_{n}^{k}$ and the first exit from $B_{n}^{k}$. Let $L(n)$ be the total time spent in the layer determined by $z=n^{-2}$. We need only show that there are constants $\lambda_{1}$ and $\lambda_{2}$ such that

$$
\begin{equation*}
P_{x_{0}}^{h}\left(L(n)>\lambda_{1}\right)>\lambda_{2} . \tag{3.1}
\end{equation*}
$$

Indeed it then follows that

$$
P_{x_{0}}^{h}\left(L(n)>\lambda_{1} \text { infinitely often }\right)>0
$$

and hence $P_{x_{0}}^{h}\left(\tau_{D}=\infty\right)>0$, which implies the desired conclusion.
Fix $n$ and let $V$ denote the random vector of last exit positions $X(\mu(n, k)), k$ $=1,2, \ldots, n^{6}$. Choose $0<\delta<\frac{1}{4}$ and let $\beta_{n}^{k}$ denote the ball of radius $\frac{1}{2} n^{-3} \delta$ concentric with $B_{n}^{k}$. We will show that there is an absolute constant $c_{3}$ such that

$$
\begin{equation*}
P_{x_{0}}^{h}\left(X(\mu(n, k)+t \wedge T(n, k)) \text { hits } \beta_{n}^{k} \mid \mathrm{V}\right) \geqq \mathrm{c}_{3} \text {, a.s. } \tag{3.2}
\end{equation*}
$$

However, let us first show how (3.1) follows. Conditioned on $V$ the $T(n, k)$ are independent random variables. (This follows from the conditional independence of past and future at a last exit time. See [4 or 9].) Let $c_{1}$ be as in Lemma 3.1 and let $A_{k}$ be the event in (3.2). Then

$$
\begin{aligned}
P_{x_{0}}^{h}\left(T(n, k)>n^{-6} c_{1} \mid V\right) & \geqq P_{x_{0}}^{h}\left(T(n, k)>n^{-6} c_{1}, A_{k} \mid V\right) \\
& \geqq E_{x_{0}}^{h}\left(P_{X(\sigma)}^{h}\left(\tau_{B_{n}^{k}}>n^{-6} c_{1} \mid V\right) ; A_{k} \mid V\right),
\end{aligned}
$$

where $X(\sigma)$ denotes the hitting position on $\beta_{n}^{k}$. The conditional probability inside may be collapsed to a single conditioning with a new harmonic function on $B_{n}^{k}$. More specifically, note that for $B \subseteq \operatorname{supp} P_{x}^{h} X(\mu(n, k+1))^{-1}$

$$
k(B, x)=P_{x}^{h}(X(\mu(n, k+1) \in B)
$$

is $h$-harmonic in $B_{n}^{k}$ (use Dynkin's formula, for example) and $k(\cdot, x) \ll k\left(\cdot, y_{0}\right)$ where $y_{0}$ is any fixed point in $B_{n}^{k}$ (use maximum principle on the harmonic function $k(B, \cdot) h(\cdot))$. Then $k_{z}(x)=\lim _{B \rightarrow\{z\}} \frac{k(B, x)}{k\left(B, y_{0}\right)}, B$ shrinking nicely to $z$, defines an $h$-harmonic function on $B_{n}^{k}$ as the convergence is uniform on compact subsets of $B_{n}^{k}$. Using now the above stated conditional independence

$$
\begin{aligned}
E_{x_{0}}^{h}\left(P _ { X ( \sigma ) } ^ { h } \left(\tau_{B_{n}^{k}}\right.\right. & \left.\left.>n^{-6} c_{1} \mid V\right) ; A_{k} \mid V\right)(\omega) \\
& =E_{x_{0}}^{h}\left(P_{X(\sigma)}^{h}\left(\tau_{B_{n}^{k}}>n^{-6} c_{1} \mid X(\mu(n, k+1))=z\right) ; A_{k} \mid V\right)(\omega)
\end{aligned}
$$

where $z=z(\omega)=X(\mu(n, k+1)(\omega),(\omega))$. Furthermore,

$$
\begin{gathered}
E_{x_{0}}^{h}\left(P_{X(\sigma)}^{h}\left(\tau_{B_{n}^{k}}>n^{-6} c_{1}\right) k_{z}\left(X\left(\tau_{B_{n}^{k}}\right)\right) k_{z}(X(\sigma))^{-1} ; A_{k} \mid V\right)(\omega) \\
=E_{x_{0}}^{h}\left(P_{X(\sigma)}^{h k_{z}}\left(\tau_{B_{n}^{k}}>n^{-6} c_{1}\right) ; A_{k} \mid V\right)(\omega) .
\end{gathered}
$$

It then follows from Lemma 3.1 and (3.2) that

$$
P_{x_{0}}^{h}\left(T(n, k)>n^{-6} c_{1} \mid V\right) \geqq c_{2} c_{3}, \quad \text { a.s. }
$$

Let $N$ denote the number of the $T(n, k)$ in the triangular arrays of random variables, independent in each row, which are greater than $n^{-6} c_{1}$. It follows from the weak law of large numbers (see e.g. [1]), that there is a constant $c_{4}$ such that $P\left(\left.N>\frac{n^{6}}{2} c_{2} c_{3} \right\rvert\, V\right)>c_{4}$. Thus

$$
P\left(\left.\sum_{k=1}^{n^{6}} T(n, k)>\frac{1}{2} c_{1} c_{2} c_{3} \right\rvert\, V\right)>c_{4} .
$$

Since $L(n) \geqq \sum_{k=1}^{n^{6}} T(n, k)$, (3.1) follows with $\lambda_{1}=\frac{1}{2} c_{1} c_{2} c_{3}$ and $\lambda_{2}=c_{4}$.
There remains only to show (3.2). Let $C$ denote the intersection with $\partial B_{n}^{k}$ of the outgoing tube. Let $Y(t)$ denote the process $X(\mu(n, k)+(t \wedge T(n, k)))$. Then there is a nonnegative function $g$, harmonic in $B_{n}^{k}$ with boundary values supported in $C$, such that the semigroup of $Y(t)$ on $(0, \infty)$ is the same as that of $X\left(t \wedge \tau_{B \hbar}\right)$ conditioned by $g$. With these remarks in mind it is easy to see that (3.2) follows from

Lemma 3.2. Let $0<\delta<\frac{1}{4}$ and $B$ and $B_{\delta}$ be the balls of radius 1 and $\delta$ respectively, centered at the origin. For $0<\varepsilon<\frac{1}{2}$ define $I_{\varepsilon}$ as the intersection with $\partial B$ of the ball of radius $\varepsilon$ centered at $(0,1,0)$, and $J_{\varepsilon}$ as the intersection with $B$ of a similar ball centered at $(0,-1,0)$. Let $g$ be any nonnegative function harmonic on $B$ with boundary values supported in $I_{\varepsilon}$. Finally, let $Y(t)$ be a process defined on a probability space $(\Omega, \mathscr{F}, P)$ and having the following properties:
(3.3) $Y(t)$ has continuous paths on $[0, \infty)$.

$$
\begin{equation*}
P\left(Y(0) \in J_{\varepsilon / 2}\right)=1 \tag{3.4}
\end{equation*}
$$

(3.5) $Y(t)$ is strong Markov on $(0, \infty)$.
(3.6) The semigroup of $Y(t)$ on $(0, \infty)$ is the same as that of Brownian $g$-paths.

Then there is a constant $c_{4}$, depending only on $\varepsilon$ and $\delta$, such that

$$
P\left(Y(t) \text { hits } B_{\delta}\right) \geqq c_{4} .
$$

Proof. Let $P(z, w)$ denote the Poisson kernel for the annular region $B \backslash B_{\delta}$ with pole at $w$ and let $Q(z, w)$ denote the Poisson kernel for $B$ with pole at $w$. Then there is an absolute constant $c_{5}$ such that

$$
\begin{equation*}
\frac{P(z, w)}{Q(z, w)} \leqq c_{5}<1 \quad \text { for } z \in J_{\varepsilon}, w \in I_{\varepsilon} \tag{3.7}
\end{equation*}
$$

To see this, first note that it is enough to show that

$$
\lim _{x \uparrow 1} \frac{P((0,-x, 0), w)}{Q((0,-x, 0), w)}<1
$$

for each $w$ in $I_{2 \varepsilon}$. Fix such a $w$ and let $q(x)$ and $p(x)$ stand for the denominator and numerator above and $\omega_{x}$ for harmonic measure relative to $B \backslash B_{\delta}$ at the point $(0,-x, 0)$. Then

$$
p(x)=q(x)-\int_{2 B_{\delta}} Q(z, w) \omega_{x}(d z)
$$

Since $Q(z, w)$ is bounded below by some positive absolute constant on $B_{\delta}$ it $\begin{aligned} & \text { suffices to show that } \lim _{x \uparrow 1} \frac{\omega_{x}\left(B_{\delta}\right)}{q(x)}>0 \text {. This follows at once from the explicit ex- } \\ & \text { pressions }\end{aligned}$

$$
\omega_{x}\left(B_{\delta}\right)=\frac{\delta(1-x)}{x(1-\delta)}
$$

and

$$
q(x)=c\left(1-x^{2}\right)|w-(0,-x, 0)|^{-3} .
$$

Now it follows from (3.3) and (3.4) that there is a strictly positive stopping time $\tau$ such that $Y(\tau)$ belongs to $J_{\varepsilon}$ with probability 1. Fix $z$ in $J_{\varepsilon}$. By (3.6) we have

$$
P_{z}\left(Y(t) \text { hits } B_{\delta}\right)=\frac{u(z)}{g(z)}
$$

where $u$ is a function harmonic in $B \backslash B_{\delta}$ with boundary values $g$ on $B_{\delta}$ and 0 on $\partial B$. Now

$$
g(z)=\int_{I_{c}} g(w) Q(z, w) d w
$$

and

$$
u(z)=g(z)-\int_{I_{\varepsilon}} g(w) P(z, w) d w .
$$

But

$$
\int_{I_{\varepsilon}} g(\mathrm{w}) \mathrm{P}(\mathrm{z}, \mathrm{w}) \mathrm{dw} \leqq \mathrm{c}_{5} \int_{I_{\varepsilon}} \mathrm{g}(\mathrm{w}) \mathrm{Q}(\mathrm{z}, \mathrm{w}) \mathrm{dw}
$$

by (3.7). Thus

$$
\frac{u(z)}{g(z)} \geqq 1-c_{5}>0 .
$$

Finally, by (3.5) we have

$$
P\left(Y(t) \text { hits } B_{\delta}\right)=E P_{Y(\tau)}\left(Y(t) \text { hits } B_{\delta}\right) \geqq 1-c_{5}
$$

and the proof is complete.

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