## Annales de l'I. H. P., Section B

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Annales de l'I. H. P., section B, tome 29, n 2 (1993), p. 199-228

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# Non-polar points for reflected Brownian motion 

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Abstract. - For each measurable function $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$, a reflected Brownian motion $X$ in a halfplane with the variable angle of reflection $\theta(x)$ is constructed. A new class of exceptional points on standard Brownian paths is studied. It is related to the problem of whether $\mathbf{X}$ hits a fixed boundary point.

Key words : Brownian motion, reflected Brownian motion, polar sets.
Résumé. - Pour chaque fonction mesurable $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$, on construit un mouvement brownien réfléchi $X$ dans un demi-plan avec un angle variable de réflection $\theta(x)$. On étudie une nouvelle classe de points exceptionnels sur le chemin brownien usuel, reliée au problème de déterminer si X atteint un point fixé sur le bord.

## 1. INTRODUCTION AND MAIN RESULTS

Our main results are
(i) a new construction of reflected Brownian motion X in a half-plane with non-smooth angle of oblique reflection and

Research supported in part by N.S.F. grants D.M.S.-8901255 and D.M.S.-9002852.
(ii) a theorem on existence of some "exceptional" points on the paths of the standard 2-dimensional Brownian motion.

The link between these two seemingly disparate results will be formed by some theorems related to the following question.
(iii) Which points are hit by the reflected Brownian motion X with positive probability?

In order to state our results we introduce some notation. Let $\mathrm{D}_{*}=\{x \in \mathbb{C}: \operatorname{Im} x>0\}$, identify $\mathbb{R}^{2}$ with $\mathbb{C}$ and $\partial \mathrm{D}_{*}$ with $\mathbb{R}$ and suppose that $\theta: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ is $\mathbf{C}^{2}$. Then we may construct a reflected Brownian motion in $D_{*}$ with the oblique angle of reflection $\theta$ in the following way [see, e. g., Rogers (1991)].

Let $\mathrm{Y}=\mathrm{Y}_{1}+i \mathrm{Y}_{2}$ be a standard 2-dimensional Brownian motion, $\mathrm{Y}(0) \in \overline{\mathrm{D}}_{*}$, and let

$$
\begin{gather*}
\mathrm{L}_{t} \stackrel{\mathrm{df}}{=} \max \left(-\inf _{s \leqq t} \mathrm{Y}_{2}(t), 0\right) \\
\mathrm{X}_{2}(t) \stackrel{\mathrm{df}}{=} \mathrm{Y}_{2}(t)+\mathrm{L}_{t} \tag{1.1}
\end{gather*}
$$

Then the equation

$$
\begin{equation*}
\mathrm{X}_{1}(t)=\mathrm{Y}_{1}(t)+\int_{0}^{t} \tan \theta\left(\mathrm{X}_{1}(s)\right) d \mathrm{~L}_{s} \tag{1.2}
\end{equation*}
$$

has a solution, possibly exploding in a finite time R. The process $\mathrm{X}(t)=\mathrm{X}_{1}(t)+i \mathrm{X}_{2}(t), t \in[0, \mathrm{R})$, is the reflected Brownian motion in $\mathrm{D}_{*}$ with the angle of reflection $\theta(x)$ measured in the clockwise direction from the inward pointing normal at $x$.

We will consider the space of RCLL (right continuous with left limits) functions equipped with the $\mathrm{M}_{1}$ topology defined by Skorohod (1960). The original definition deals with the functions defined on [ 0,1 ] and can be extended in an obvious way to an arbitrary non-random interval $[0, t]$. We have to further extend the definition to functions defined on random time intervals. We will say that processes $X_{k}$ converge a.s. to $X$ in $M_{1}$ topology if
(i) X and $\mathrm{X}_{k}, k \geqq 1$, are RCLL;
(ii) $X$ is defined on a random time interval $[0, R)$, each process $X_{k}$ is defined on an interval $\left[0, \mathrm{R}_{k}\right.$ ) and $\mathrm{R}_{k} \rightarrow \mathrm{R}$ a. s.; and
(iii) for every fixed $\varepsilon>0$ and $a<\infty, X_{k}$ converges to $X$ in the $M_{1}-$ topology of Skorohod (1960) on the interval $[0,0 \vee((R-\varepsilon) \wedge a)]$ a.s.

Note that for every Borel measurable function $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ there is a sequence of $C^{2}$ functions $\theta_{k}: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ which converges to $\theta$ almost everywhere.

Theorem 1.1. - Suppose that $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is Borel measurable and the functions $\theta_{k}: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ are $\mathrm{C}^{2}$ and converge almost everywhere to $\theta$ as $k \rightarrow \infty$. Then, on some probability space, there exists a sequence of processes $\left\{\mathrm{X}_{k}\right\}_{k \geqq 1}$ such that
(i) for each $k, \mathrm{X}_{k}$ is a reflected Brownian motion with the angle of reflection $\theta_{k}$ and
(ii) $\mathrm{X}_{k}$ converge a.s. to a process $\mathbf{X}$ in $\mathrm{M}_{1}$ topology.

Note that if X and $\tilde{\mathrm{X}}$ are the limit processes corresponding to two different approximating sequences $\left\{\theta_{k}\right\}$ and $\left\{\tilde{\theta}_{k}\right\}$ then the distributions of X and $\tilde{\mathrm{X}}$ are identical since the theorem may be applied to the sequence $\left\{\theta_{1}, \tilde{\theta}_{1}, \theta_{2}, \tilde{\theta}_{2}, \ldots\right\}$.

We will consider Theorem 1.1 as the definition of the reflected Brownian motion with an arbitrary measurable angle of reflection $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$. We are not going to present in this paper any other characterization of X as a reflected Brownian motion, for example, as a solution to a submartingale problem or as a solution to equations analogous to (1.1)-(1.2). We plan to discuss these questions in a future article. However, we will briefly indicate in Remark 2.2 how it may be proved that the process X is strong Markov.

The process X is continuous if and only if there are no non-degenerate intervals on which we have almost everywhere $\theta=\pi / 2$ or $\theta=-\pi / 2$; this will be easily seen from our proof of Theorem 1.1.

If the limiting process $X$ is continuous then $X_{k}$ converge to $X$ in the topology stronger than $M_{1}$, namely uniformly on $[0,0 \vee((R-\varepsilon) \wedge a)]$. If $\mathbf{X}$ is discontinuous then the convergence does not hold in $\mathbf{J}_{1}$ or $\mathbf{J}_{2}$ topologies of Skorohod (1960) (which are stronger than $\mathbf{M}_{1}$ ) because continuous functions cannot converge to a discontinuous one in either one of these topologies.

Now we turn to exceptional points on Brownian paths. We will need some more definitions and notation. A domain (i.e. an open and connected set) $\mathrm{D} \subset \mathbb{C}$ will be called monotone if $x+i b \in \mathrm{D}$ whenever $x \in \mathrm{D}$ and $b>0$.

For a Greenian domain D , its Green function will be denoted $\mathrm{G}_{\mathrm{D}}(.,$.$) .$ For a monotone domain D , let

$$
\tilde{\mathrm{D}}=\{z \in \mathrm{D}:|\operatorname{Re} z|<1\} \quad \text { and } \quad \mathrm{I}=\{z \in \mathrm{D}: \operatorname{Re} z=0\}
$$

Let $h=h_{\mathrm{D}}$ be defined by

$$
\begin{equation*}
h(x)=|\operatorname{Re} x|+\int_{i y \in 1} \mathrm{G}_{\tilde{\mathrm{D}}}(x, i y) \frac{d y}{\pi} \quad \text { for } \quad x \in \tilde{\mathrm{D}} \tag{1:3}
\end{equation*}
$$

The function $h$ may be characterized as the (unique) bounded harmonic function in D which is equal to $|\operatorname{Re} x|$ on $\partial \tilde{\mathrm{D}}$ (see Lemma 3.1).

For $\varepsilon>0$ let $d_{\varepsilon}: \mathbb{R} \rightarrow \mathbb{R}$ be the smallest non-increasing non-negative function satisfying

$$
\{x \in \partial \mathrm{D}:|x|<\varepsilon\} \subset\left\{x \in \mathbb{C}:|x|<\varepsilon, \operatorname{Re} x \leqq d_{\varepsilon}(\operatorname{Im} x)\right\} .
$$

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Suppose that D is a monotone domain and $x \in \mathrm{D}$. Then $x+a i$ converges as $a \rightarrow \infty$ to a prime end (independent of $x$ ) which we will call $\infty$. Let $\mathrm{I}_{*} \stackrel{\mathrm{df}}{=}\{z \in \mathbb{C}: \operatorname{Re} z=0, \operatorname{Im} z>0\}$. If $0 \in \partial \mathrm{D}$ and $\mathrm{I}_{*} \subset \mathrm{D}$ then ai converges as $a \downarrow 0$ to a prime end which will be denoted 0 [see Pommerenke (1975) for the definition of prime ends].

Suppose Y is a standard 2-dimensional Brownian motion and D is a monotone domain such that $0 \in \partial \mathrm{D}$ and $\mathrm{I}_{*} \subset \mathrm{D}$. We will say that $\mathrm{Y}\left(t_{0}\right)$ is a D-point if $t_{0}>0$ and there exists $t_{1} \in\left(0, t_{0}\right)$ such that $\operatorname{ReY}\left(t_{0}\right)=0$ and $\mathrm{Y}\left(t_{1}, t_{0}\right) \subset \mathrm{D}+\mathrm{Y}\left(t_{0}\right)$. Here $\mathrm{D}+x=\{z \in \mathbb{C}: \exists y \in \mathrm{D}$ such that $y+x=z\}$. We will call $\mathrm{Y}\left(t_{0}\right)$ a right D -point if $\mathrm{Y}\left(t_{0}\right)$ is a D -point and for some $t_{1} \in\left(0, t_{0}\right)$ we have $\operatorname{Re} \mathrm{Y}(s)>0$ for all $s \in\left(t_{1}, t_{0}\right)$.

Theorem 1.2. - Suppose D is a monotone domain, $0 \in \partial \mathrm{D}$ and $\mathrm{I}_{*} \subset \mathrm{D}$.
(i) The standard 2-dimensional Brownian motion has D-points with positive probability (and, therefore, with probability 1) if and only if

$$
\begin{equation*}
\limsup _{\varepsilon \rightarrow 0+} h(\varepsilon i) / \mathrm{G}_{\mathrm{D}}(i, \varepsilon i)<\infty \tag{1.4}
\end{equation*}
$$

(ii) The standard 2-dimensional Brownian motion has right D-points with positive probability (and, therefore, with probability 1) if and only if for some $\varepsilon>0$

$$
\begin{equation*}
\int_{-1}^{0} d_{\varepsilon}(r) r^{-2} d r<\infty \tag{1.5}
\end{equation*}
$$

We will show in Remark 3.2 below that (1.4) fails if and only if

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0+} h(\varepsilon i) / \mathrm{G}_{\mathrm{D}}(i, \varepsilon i)=\infty \tag{1.6}
\end{equation*}
$$

It is not easy to verify (1.4) for an arbitrary domain so let us point out some cases when this can be done. Suppose that $\mathrm{D}=\{z \in \mathbb{C}: \operatorname{Im} z \geqq a \operatorname{Re} z\}$ for some $a \in \mathbb{R}$. Then D-points do not exist by Theorem 2.1 (iii) of Burdzy (1989); one can also relatively easily check that (1.4) is not satisfied. The property of having D-points is monotone in D so we may want to consider domains of the form

$$
\begin{equation*}
\mathrm{D}=\{z \in \mathbb{C}: \operatorname{Im} z \geqq a \operatorname{Re} z-g(\operatorname{Re} z)\} \tag{1.7}
\end{equation*}
$$

where $a \in \mathbb{R}$ and $g: \mathbb{R} \rightarrow \mathbb{R}$ is a non-negative Lipschitz function such that $g(0)=0$. If

$$
\begin{equation*}
\int_{-1}^{1} g(r) r^{-2} d r<\infty \tag{1.8}
\end{equation*}
$$

then $D$ has an angular derivative at 0 [see Burdzy (1987, Definition 9.1)] and we may show as in the proof of Theorem 1.4 (i), using Theorems 8.2 and 9.2 of Burdzy (1987), that (1.4) does not hold; the converse is
not true, i.e., (1.4) need not be satisfied for domains of the form (1.7) even if (1.8) fails. It would take too much space to present an example to this effect. If $g(r)=a_{1}|r|$ for some $a_{1}>0$ then D-points exist by Theorem 2.1 of Burdzy (1989).

Suppose for a moment that $\tilde{\mathrm{D}}=\mathrm{D}$. The ratio $\mathrm{G}_{\mathrm{D}}(\varepsilon i, y) / \mathrm{G}_{\mathrm{D}}(i, \varepsilon i)$ converges to the Martin kernel $\mathrm{K}(0, y)$ in D as $\varepsilon \rightarrow 0$ [see Doob (1984)]. Hence, by the Fatou lemma, (1.4) implies

$$
\int_{\mathrm{I}} \mathrm{~K}(0, y) d y<\infty
$$

It seems that the converse is false. Let $\mathrm{M}=f^{-1}\left(\left\{z \in \mathrm{D}_{*}: \operatorname{Re} z=0\right\}\right)$ where $f$ is a function defined as in Lemma 2.1. Then

$$
\liminf _{\substack{x \rightarrow 0 \\ x \in M}} \frac{|x|}{\mathrm{G}_{\mathrm{D}}(i, x)}+\int_{\mathrm{I}} \mathrm{~K}(0, y) d y<\infty
$$

seems to be equivalent to (1.4). Our proof is long and complicated and therefore we omit it.

Measurable functions $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ correspond to monotone domains $\mathrm{D} \subsetneq \mathbb{C}$. Lemmas 2.1 and 2.2 describe precisely the nature of this correspondence and we briefly discuss it here. For a monotone domain $\mathrm{D} \subsetneq \mathbb{C}$ there exists a univalent function $f$ mapping $\mathrm{D}_{*}$ onto D and such that for almost all $x \in \partial \mathrm{D}_{*}, f^{\prime}(x)$ exists in the sense of the angular limit and belongs to $[-\pi / 2, \pi / 2]$. We let $\theta(x)=\arg f^{\prime}(x)$. Conversely, if $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is given then there exists a univalent function $f$ which maps $\mathrm{D}_{*}$ onto a monotone domain and such that $\arg f^{\prime}(x)=\theta(x)$ for almost all $x$. It should be emphasized that the correspondence between $\theta$ 's and D's is not one-to-one. Our next result exploits this relationship between O's and D's. Let

$$
\mathrm{T}_{\mathrm{A}}=\mathrm{T}(\mathrm{~A})=\inf \{t>0: \mathrm{X}(t) \in \mathrm{A}\}
$$

with the convention inf $\varnothing=\infty$.
Theorem 1.3. - Suppose that
(a) a measurable function $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is given and a monotone domain D and a function $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ are defined as in Lemma 2.2; if it is possible, D and $f$ are chosen so that $0 \in \partial \mathrm{D}, \mathrm{I}_{*} \subset \mathrm{D}$ and $f(0)=0$; or
(b) a monotone domain D is given; $0 \in \partial \mathrm{D}, \mathrm{I}_{*} \subset \mathrm{D}$ and $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ and $\theta(x)=\arg f^{\prime}(x)$ are chosen as in Lemma 2.1 so that $f(0)=0$.

Under any of these assumptions we have the following.
(i) The reflected Brownian motion in $\mathrm{D}_{*}$ with the oblique angle of reflection $\theta$ hits 0 with positive probability if and only if $0 \in \partial \mathrm{D}, \mathrm{I}_{*} \subset \mathrm{D}$ and the standard 2-dimensional Brownian motion has D-points a.s.
(ii) The reflected Brownian motion X in $\mathrm{D}_{*}$ with the oblique angle of reflection $\theta$ hits 0 with positive probability and approaches it from the right, i.e.,

$$
\begin{equation*}
\mathrm{P}\left(\mathrm{~T}_{\{0\}}<\infty \text { and } \exists \varepsilon>0: \operatorname{Re} \mathrm{X}(t)>0 \text { for all } t \in\left(\mathrm{~T}_{\{0\}}-\varepsilon, \mathrm{T}_{\{0\}}\right)\right)>0 \tag{1.9}
\end{equation*}
$$

if and only if $0 \in \partial \mathrm{D}, \mathrm{I}_{*} \subset \mathrm{D}$ and the standard 2-dimensional Brownian motion has right D-points a.s.

Smooth 日's correspond to D's with smooth boundaries and vice versa. Hence, it may be interesting to have a look at D's with highly irregular boundaries, for example, at a domain D which lies above a "typical" graph of a 1-dimensional Brownian motion. What can we say about the corresponding $\theta$ and reflected Brownian motion $X$ in this case? Theorem 1.4 (iii) and its proof provide some information about this process.

Suppose that reflected Brownian motion X in $\mathrm{D}_{*}$ with the oblique angle of reflection $\theta$ satisfies $\mathrm{X}(0) \in \mathrm{D}_{*}$ a.s. Let $\mathscr{A}$ be the set of all non-polar points for X , i.e., the set of all points $x \in \partial \mathrm{D}_{*}$ such that

$$
\mathrm{P}\left(\mathrm{~T}_{\{x\}}<\infty\right)>0
$$

How large is the set $\mathscr{A}$ ? The answer depends on $\theta$ and we start our discussion with a few examples.

Example 1.1. - (i) If $\theta(x)=0$ for all $x$ then X has the normal reflection at $\partial \mathrm{D}_{*}$ and it is well known that $\mathscr{A}=\varnothing$.
(ii) Fix some $a>0$ and let $\theta(0)=0, \theta(x)=a$ for $x<0$ and $\theta(x)=-a$ for $x>0$. Then $\mathscr{A}=\{0\}$, by the results of Varadhan and Williams (1985).
(iii) An easy modification of the previous example shows that for some $\theta, \mathscr{A}$ may be equal to the set of all integers.

Can $\mathscr{A}$ be uncountably infinite? Can $\mathscr{A}$ be equal to $\partial \mathrm{D}_{*}$ ?
Let $\lambda$ denote the Lebesgue measure on $\partial \mathrm{D}_{*}$.
Theorem 1.4. - (i) $\lambda(\{x \in \mathscr{A}: \theta(x) \in(-\pi / 2, \pi / 2)\})=0$ for every $\theta$.
(ii) For every $\theta$, the set $\partial \mathrm{D}_{*} \backslash \mathscr{A}$ is dense in $\partial \mathrm{D}_{*}$.
(iii) There exists $\theta$ such that $\lambda\left(\partial \mathrm{D}_{*} \backslash \mathscr{A}\right)=0$. One can choose $\theta$ so that all of the following conditions are satisfied.
(a) For every $a<b$

$$
\begin{gathered}
\lambda(\{x \in(a, b): \theta(x)=-\pi / 2\})>0, \\
\lambda(\{x \in(a, b): \theta(x)=\pi / 2\})>0,
\end{gathered}
$$

and

$$
\lambda(\{x \in(a, b):|\theta(x)| \neq \pi / 2\})=0 .
$$

(b) There is a set $\mathscr{H} \subset \mathscr{A}$ such that $\lambda\left(\partial \mathrm{D}_{*} \backslash \mathscr{H}\right)=0$ and for every $x \in \mathscr{H}$, X can approach $x$ from one side, i. e., for each $x \in \mathscr{H}$ we have either

$$
\mathrm{P}\left(\mathrm{~T}_{\{x\}}<\infty \text { and } \exists \varepsilon>0: \operatorname{Re} \mathrm{X}(t)<x \text { for all } t \in\left(\mathrm{~T}_{\{x\}}-\varepsilon, \mathrm{T}_{\{x\}}\right)\right)>0
$$

or

$$
\mathrm{P}\left(\mathrm{~T}_{\{x\}}<\infty \text { and } \exists \varepsilon>0: \operatorname{Re} \mathrm{X}(t)>x \text { for all } t \in\left(\mathrm{~T}_{\{x\}}-\varepsilon, \mathrm{T}_{\{x\}}\right)\right)>0 .
$$

(c) The process X is recurrent and so
(c1) $\mathrm{P}\left(\mathrm{T}_{\{x\}}<\infty\right)=1$ for every $x \in \mathscr{A}$ and
(c2) $\mathrm{P}\left(\lambda\left(\left\{x \in \partial \mathrm{D}_{*}: \mathrm{T}_{\{x\}}=\infty\right\}\right)=0\right)=1$.
The reflected Brownian motion X discussed in Theorem 1.4 (iii) is continuous despite the fact that $|\theta(x)|=\pi / 2$ a.e.; the reason is that $\theta$ is not constant on any interval.

In view of Theorem 1.4 (iii) ( $c 2$ ) we propose the following
Problem 1.1. - Does there exist $\theta$ such that for some set $\mathscr{B} \subset \partial \mathrm{D}_{*}$ we have $\lambda\left(\partial \mathrm{D}_{*} \backslash \mathscr{B}\right)=0$ and

$$
\mathrm{P}\left(\left\{x \in \mathscr{B}: \mathrm{T}_{\{x\}}=\infty\right\}=\varnothing\right)=1 ?
$$

A construction of reflected Brownian motion with non-smooth angle of reflection was presented by Tsuchiya $(1976,1980)$ under the assumption that $\theta \in(-\pi / 2+\varepsilon, \pi / 2-\varepsilon)$ for some $\varepsilon>0$. In recent articles, Motoo (1989, 1990) presented a more general construction. See also Takanobu and Watanabe (1988).

Theorem 1.2 is a more accurate version of Theorem 2.1 of Burdzy (1989).

The question of whether the reflected Brownian motion in a wedge hits the vertex was treated by Varadhan and Williams (1985) under the assumption of the constant angle of reflection on each side of the wedge. Rogers (1991) attacked the same problem for reflected Brownian motion in a half-plane with a variable angle of reflection of class $\mathrm{C}^{1}$ and obtained some partial results. The complete solution of the last problem was obtained independently by Rogers (1990) and Burdzy and Marshall (1992). A perfectly explicit integral test in terms of $\theta$ determines whether the process hits a fixed boundary point with positive probability. The same test applies in our present context, i.e., when $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is allowed to be any measurable function. We omit the proof as it does not require any new ideas.

Let us briefly discuss the sources of our techniques. We use conformal mappings in the manner of Rogers (1989, 1991). A connection between reflected Brownian motion and standard Brownian motion goes back to Lévy (1948) in the 1 -dimensional case. El Bachir (1983) and Le Gall (1987) adapted the idea to the 2-dimensional processes. The proof of Theorem 1.2 (i) uses an idea of Davis (1983).

We would like to express our gratitude to Chris Bishop, Joe Glover, Ron Pyke, Chris Rogers, Michael Sharpe and Ruth Williams for the most valuable advice.

## 2. A CONSTRUCTION OF REFLECTED BROWNIAN MOTION

The set

$$
\mathscr{S}_{\varphi}=\left\{z \in \mathrm{D}_{*}:|\arg z-\pi / 2|<\varphi\right\}
$$

is called a Stolz angle. Suppose $f$ is a function defined on $\mathrm{D}_{*}$. If the limit

$$
\lim _{z-x \rightarrow 0} f(z)
$$

exists for every $\varphi<\pi / 2$ then it is called the angular (or non-tangential) limit of $f$ at $x$ and denoted $f(x)$. In particular, the angular derivative $f^{\prime}(x)$ is defined as the angular limit of $f^{\prime}$ at $x$, if it exists. If $f: \mathrm{D}_{*} \rightarrow \mathbb{C}$ is analytic and $f(x)$ exists for some $x \in \partial \mathrm{D}_{*}$ then $f^{\prime}(x)$ exists if and only if the angular limit of the function $z \rightarrow \frac{f(z)-f(x)}{z-x}$ exists; if both limits exist then they are identical [Pommerenke (1975, Theorem 10.5)].

Lemma 2.1. - Suppose that D is a monotone domain with non-empty complement. Then there exists a univalent analytic mapping $f$ of $\mathrm{D}_{*}$ onto D
such that $f(x)$ and $f^{\prime}(x)$ exist, $f^{\prime}(x) \neq 0$ and $\theta(x) \stackrel{\text { df }}{=} \arg f^{\prime}(x) \in[-\pi / 2, \pi / 2]$, for almost all $x \in \partial \mathrm{D}_{*}$. If 0 is a prime end of D , we may choose $f$ so that $f(0)=0$ and $\lim _{|z| \rightarrow \infty}|f(z)|=\infty$.

Proof. - Choose regions $\mathrm{D}_{k} \subset \mathrm{D}$ such that $\mathrm{D}_{k} \subset \mathrm{D}_{k+1}, \cup \mathrm{D}_{k}=\mathrm{D}$, and such that $\partial \mathrm{D}_{k}$ is the graph of a smooth function, $k=1,2, \ldots$ Fix $w_{0} \in \mathrm{D}_{1}$ and choose $z_{0} \in \mathrm{D}_{*}$ and conformal maps $f_{k}$ of $\mathrm{D}_{*}$ onto $\mathrm{D}_{k}$ with $f_{k}\left(z_{0}\right)=w_{0}$ and $f_{k}^{\prime}\left(z_{0}\right)>0$. Likewise, let $f$ map $\mathrm{D}_{*}$ conformally onto D with $f\left(z_{0}\right)=w_{0}$ and $f^{\prime}\left(z_{0}\right)>0$. Then, by Theorem IX. 13 of Tsuji (1959), the sequence $\left\{f_{k}\right\}$ converges to $f$ normally on $\mathrm{D}_{*}$ (i.e., uniformly on compact subsets of $\mathrm{D}_{*}$ ). Since $\operatorname{Re} f_{k}^{\prime}(x)>0$ for $x \in \partial \mathrm{D}_{*}$, we conclude $\operatorname{Re} f_{k}^{\prime}(z)>0$ for $z \in \mathrm{D}_{*}$. Taking normal limits, clearly $\operatorname{Re} f^{\prime}(z)>0$ for $z \in \mathrm{D}_{*}$. Thus we can choose a branch of $\arg$ so that $\arg f^{\prime}(z) \in(-\pi / 2, \pi / 2)$ for all $z \in \mathrm{D}_{*}$. By Corollary 2.6, p. 114, of Garnett (1981), $f^{\prime} \in \mathbf{H}^{p}$ for all $p<1$. Thus, by Theorem 2.2, p. 17, of Duren (1970), the angular limits $f^{\prime}(x)$ exist a. e. and by Corollary 4.2, p. 65, of Garnett (1981) are non-zero a.e. By Theorem 5.12, p.88, of Duren (1970), $f \in \mathrm{H}^{q}$ for all $q<\infty$ and thus the angular limits $f(x)$ exist a. e.

As for the conditions $f(0)=0$ and $f(\infty)=\infty$, it is well known that the value of the function $f$ at two boundary points (even three boundary points) may be chosen in an arbitrary way.

In order to avoid trivialities we will assume from now on that $\theta(x)$ is neither identically equal to $-\pi / 2$ nor identically equal to $\pi / 2$.

Lemma 2.2. - Suppose that $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is a Borel measurable function. Then there exists a univalent analytic mapping $f$ of $\mathrm{D}_{*}$ onto a monotone domain D such that $f(x)$ and $f^{\prime}(x)$ exist, $f^{\prime}(x) \neq 0$ and $\arg f^{\prime}(x)=\theta(x)$ for almost all $x \in \mathbb{R}$. We may choose $f$ so that $\lim _{|z| \rightarrow \infty}|f(z)|=\infty$.

Proof. - Let $\theta: \mathrm{D}_{*} \rightarrow \mathbb{R}$ be the bounded continuous harmonic extension of our original function $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ and let $\tilde{\theta}$ be the harmonic conjugate of $\theta$ such that $\tilde{\theta}(i)=0$. Define $f: \mathrm{D}_{*} \rightarrow \mathbb{C}$ by setting $f(i)=i$ and

$$
f^{\prime}(z)=\exp (i(\theta(z)+i \widetilde{\theta}(z)))
$$

Since we assume (2.1), $\theta(z) \in(-\pi / 2, \pi / 2)$ for all $z \in D_{*}$. Hence,

$$
\operatorname{Re} f^{\prime}(z)=e^{-\tilde{\theta}(z)} \cos \theta(z)>0
$$

As in the proof of Lemma 2.1, the angular limits $f(x)$ and $f^{\prime}(x)$ exist and $f^{\prime}(x) \neq 0$ for almost all $x \in \partial \mathrm{D}_{*}$. Let $\gamma(t)=t z+(1-t) w$ where $z, w \in \mathrm{D}_{*}$. Then $\gamma^{\prime}(t)=z-w$ and

$$
\begin{equation*}
f(z)-f(w)=\int_{0}^{1} f^{\prime}(\gamma(t))(z-w) d t=\left[\int_{0}^{1} f^{\prime}(\gamma(t)) d t\right](z-w) \tag{2.2}
\end{equation*}
$$

Since the real part of the last integral is strictly positive, $f(z)=f(w)$ if and only if $z=w$. In other words, the function $f$ is univalent.

To see that D is monotone, let $f_{\varepsilon}=f(z+i \varepsilon), \varepsilon>0$. By (2.2),

$$
\operatorname{Re} f(z)-\operatorname{Re} f(w)=\left(\int_{0}^{1} \operatorname{Re} f^{\prime}(\gamma(t)) d t\right)(\operatorname{Re} z-\operatorname{Re} w)
$$

when $\operatorname{Im} z=\operatorname{Im} w$. Thus $\operatorname{Re} f_{\varepsilon}(x)$ is increasing on $\partial \mathrm{D}_{*}$. The curve $\gamma_{\varepsilon} \stackrel{\mathrm{df}}{=}\left\{\operatorname{Re} f_{\varepsilon}(x):-\infty<x<\infty\right\}$ is an analytic Jordan curve. It must divide the plane into two regions, since $f_{\varepsilon}$ extends to be univalent in a neighborhood of $\partial \mathrm{D}_{*}$. The boundary of $f_{\varepsilon}\left(\mathrm{D}_{*}\right)$ must be contained in $\gamma_{\varepsilon} \cup\{\infty\}$, and thus $f_{\varepsilon}\left(\mathrm{D}_{*}\right)$ is monotone. Since $\varepsilon$ is arbitrary, $\mathrm{D}=f\left(\mathrm{D}_{*}\right)$ is also monotone.

Lemma 2.3. - Suppose that $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is measurable and $\theta_{k}: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ are $\mathbf{C}^{2}$-functions which converge almost everywhere to $\theta$ as $k \rightarrow \infty$. Let $\mathrm{D}, f$ correspond to $\theta$ and $\mathrm{D}_{k}, f_{k}$ correspond to $\theta_{k}$ in the same way as in the proof of Lemma 2.2. The monotone domains $\mathrm{D}_{k}$ converge to D in the following sense.
(i) If B is open and such that $\mathrm{B} \cap \partial \mathrm{D} \neq \varnothing$, there is a $k_{0}=k_{0}(\mathrm{~B})$ such that $\mathrm{B} \cap \partial \mathrm{D}_{k} \neq \varnothing$ for all $k \geqq k_{0}$.
(ii) If B is connected and open, with $\mathrm{B} \cap \mathrm{D} \neq \varnothing$ and $\mathrm{B} \subset \mathrm{D}_{k}$ for infinitely many $k$, then $\mathrm{B} \subset \mathrm{D}$.
(iii) If K is compact and $\mathrm{K} \subset \mathrm{D}$ then $\mathrm{K} \subset \mathrm{D}_{k}$ for all $k \geqq k_{0}=k_{0}(\mathrm{~K})$.

Proof. - By the Koebe-1/4 theorem, Theorem 2.3, p.31, of Duren (1983) and Schwarz's lemma, if $g$ is univalent

$$
\operatorname{Im} z\left|g^{\prime}(z)\right| / 2 \leqq \operatorname{dist}\left(g(z), \partial g\left(\mathrm{D}_{*}\right)\right) \leqq 2 \operatorname{Im} z\left|g^{\prime}(z)\right|
$$

for all $z \in \mathrm{D}_{*}$, where dist $(a, \mathrm{~K})$ is the Euclidean distance from $a$ to the set K.

Since $\begin{aligned} & \lim _{k \rightarrow \infty} f_{k}(z)=f(z) \text { and } \lim _{k \rightarrow \infty} f_{k}^{\prime}(z)=f^{\prime}(z), \\ & \operatorname{dist}(f(z), \partial \mathrm{D}) / 4 \leqq \lim _{k \rightarrow \infty} \operatorname{dist}\left(f_{k}(z), \partial \mathrm{D}_{k}\right) \leqq 4 \operatorname{dist}(f(z), \partial \mathrm{D}) .\end{aligned}$
Since $\lim _{k \rightarrow \infty} f_{k}(z)=f(z)$,

$$
\operatorname{dist}(w, \partial \mathrm{D}) / 4 \leqq \lim _{k \rightarrow \infty} \operatorname{dist}\left(w, \partial \mathrm{D}_{k}\right) \leqq 4 \operatorname{dist}(w, \partial \mathrm{D})
$$

for all $w \in \mathrm{D}$. Actually the 4 's can be removed in the above inequalities, for if $\zeta \in \mathrm{D}$, choose $w \in \mathrm{D}$ with $|\zeta-w|=(1-\varepsilon) \operatorname{dist}(\zeta, \partial \mathrm{D})$. By the above inequalities applied to $w$,

$$
(1-3 \varepsilon / 4) \operatorname{dist}(\zeta, \partial \mathrm{D}) \leqq \lim _{k \rightarrow \infty} \operatorname{dist}\left(\zeta, \partial \mathrm{D}_{k}\right) \leqq(1+3 \varepsilon) \operatorname{dist}(\zeta, \partial \mathrm{D})
$$

Thus $\lim \operatorname{dist}\left(\zeta, \partial \mathrm{D}_{k}\right)=\operatorname{dist}(\zeta, \partial \mathrm{D})$ for all $\zeta \in \mathrm{D}$.

$$
k \rightarrow \infty
$$

To prove (i), take $w_{1} \in \mathrm{D} \cap \mathrm{B}$ with $\operatorname{dist}\left(w_{1}, \partial \mathrm{D}\right)<\operatorname{dist}\left(w_{1}, \partial \mathrm{~B}\right) / 2$, and let $f\left(z_{1}\right)=w_{1}$. For $k \geqq k_{0}$, $\operatorname{dist}\left(f_{k}\left(z_{1}\right), \partial \mathrm{D}_{k}\right)<\operatorname{dist}\left(w_{1}, \partial \mathrm{~B}\right) / 2$, and $\left|f_{k}\left(z_{1}\right)-f\left(z_{1}\right)\right|<\operatorname{dist}\left(w_{1}, \partial \mathrm{~B}\right) / 4$, Thus dist $\left(w_{1}, \partial \mathrm{D}_{k}\right)<3 \operatorname{dist}\left(w_{1}, \partial \mathrm{~B}\right) / 4$, proving $\partial \mathrm{D}_{k} \cap \mathrm{~B} \neq \varnothing$ for $k \geqq k_{0}$.

Statement (ii) follows from statement (i) for if $\mathrm{B} \not \ddagger \mathrm{D}$, then $\mathrm{B} \cap \partial \mathrm{D} \neq \varnothing$.
To prove (iii), note that $K_{1} \stackrel{\text { df }}{=} f^{-1}(\mathrm{~K})$ is a compact subset of $\mathrm{D}_{*}$. Choose a large open disc $\Delta$ with $\partial \Delta \subset \mathrm{D}_{*}$ and $\mathrm{K}_{1} \subset \Delta$. The sequence $\left\{f_{k}\right\}$ converges uniformly on $\bar{\Delta}$ to $f$. For some $\varepsilon>0$ and all $k \geqq k_{0}$, $\operatorname{dist}\left(f_{k}(\partial \Delta), \mathrm{K}\right) \geqq \varepsilon$. The winding number $n\left(f_{k}(\partial \Delta), w\right)$ converges to $n(f(\partial \Delta), w)=1$ uniformly for $w \in K$. Since winding numbers are integers, $\mathrm{K} \subset f_{k}\left(\mathrm{D}_{*}\right)$ for all $k \geqq k_{0}$.

Remark 2.1. - (i) It follows from Lemma 2.3 (i) that if $x \in \partial \mathrm{D}$ then there exists a sequence $\left\{x_{k}\right\}$ such that $x_{k} \in \partial \mathrm{D}_{k}$ and $x_{k} \rightarrow x$.
(ii) If $x_{k} \in \partial \mathrm{D}_{k}$ then $x_{k}$ 's cannot converge to any point $x \in \mathrm{D}$, by Lemma 2.3 (iii).
(iii) Lemma 2.3 (ii) is false if the assumption $\mathrm{B} \cap \mathrm{D} \neq \varnothing$ is removed.

Proof of Theorem 1.1.
Step 1. - First we will establish a relationship between reflected Brownian motion in $\mathrm{D}_{*}$, standard Brownian motion and monotone domains. In this step we will assume that $\theta: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ is of class
$\mathrm{C}^{2}$. Let $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ be a mapping corresponding to $\theta$ defined in Lemma 2.2. Since $\theta$ is of class $\mathrm{C}^{2}$, the mapping $f$ is $\mathrm{C}^{2}$ on $\overline{\mathrm{D}}_{*}$. This may be proved as in Step 2 of the proof of Theorem 1.1 of Burdzy and Marshall (1992).

Let $\mathrm{V}(x)=\tan \theta(x)+i$ for $x \in \partial \mathrm{D}_{*}$, that is let $\mathrm{V}(x)$ be the vector of reflection. We may write (1.1)-(1.2) as

$$
\mathrm{X}(t)=\mathrm{Y}(t)+\int_{0}^{t} \mathrm{~V}\left(\mathrm{X}_{s}\right) d \mathrm{~L}_{s}
$$

and the Itô formula yields

$$
\begin{equation*}
f(\mathrm{X}(t))=f(\mathrm{X}(0))+\int_{0}^{t} f^{\prime}(\mathrm{X}(s)) d \mathrm{Y}_{s}+\int_{0}^{t} f^{\prime}(\mathrm{X}(s)) \mathrm{V}(\mathrm{X}(s)) d \mathrm{~L}_{s} \tag{2.3}
\end{equation*}
$$

where $f^{\prime}$ is interpreted as the Jacobian of $f$. The process X spends zero time on $\partial \mathrm{D}_{*}$ and $f$ is analytic outside $\partial \mathrm{D}_{*}$ so $f(\mathrm{X}(0))+\int_{0}^{t} f^{\prime}(\mathrm{X}(s)) d \mathrm{Y}_{s}$ is a time-change of a Brownian motion $Z$, i.e., it is equal to $Z(c(t))$ where

$$
c(t)=\int_{0}^{t}\left(f^{\prime}(\mathrm{X}(s))\right)^{2} d s
$$

Note that $c(0)=0$ and $c$ is continuous and strictly increasing because $f^{\prime}(x) \in \mathbb{C} \backslash\{0\}$ for $x \in \mathrm{D}_{*}$. The local time L does not increase unless X is at the boundary of $\mathrm{D}_{*}$ and $\operatorname{Re}\left(f^{\prime}(x) \mathrm{V}(x)\right)=0$ for $x \in \partial \mathrm{D}_{*}$ so

$$
\int_{0}^{t} f^{\prime}(\mathrm{X}(s)) \mathrm{V}(\mathrm{X}(s)) d \mathrm{~L}
$$

has null real component. Let us time change the process in (2.3); in other words let us define

$$
\begin{gathered}
\mathrm{W}(c(t))=f(\mathrm{X}(t)) \\
\mathrm{A}(c(t))=-i \int_{0}^{t} f^{\prime}(\mathrm{X}(s)) \mathrm{V}(\mathrm{X}(s)) d \mathrm{~L}_{s}
\end{gathered}
$$

Then we have

$$
\mathrm{W}(t)=\mathrm{Z}(t)+i \mathrm{~A}(t)
$$

Hence, W is a reflected Brownian motion in D with the vertical vector of reflection.

We will provide a different representation of $\mathrm{A}(t)$, following Skorohod (1961) [see Lemma 3.6.14 of Karatzas and Shreve (1988)]. Recall that $\mathrm{D}=f\left(\mathrm{D}_{*}\right)$. We will show that

$$
\begin{equation*}
\mathrm{A}(t)=-\sup \{a \leqq 0: \mathbf{Z}[0, t] \subset \mathrm{D}+a i\} \tag{2.4}
\end{equation*}
$$

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First note that

$$
\begin{equation*}
\int_{0}^{\infty} \mathbf{1}_{\mathrm{D}}(\mathrm{~W}(t)) d \mathrm{~A}_{t}=0 \tag{2.5}
\end{equation*}
$$

Let $\tilde{\mathrm{A}}(t)$ be the right hand side of (2.4). It is immediate that both $\mathrm{A}(t)$ and $\tilde{\mathrm{A}}(t)$ are non-decreasing and continuous, they both satisfy (2.5), $\mathrm{A}(0)=\widetilde{\mathrm{A}}(0)$, and $\mathrm{Z}(t)+i \mathrm{~A}(t) \in \overline{\mathrm{D}}$ and $\mathrm{Z}(t)+i \widetilde{\mathrm{~A}}(t) \in \overline{\mathrm{D}}$ for all $t$.

Suppose that $\mathrm{A}\left(t_{1}\right)>\tilde{\mathrm{A}}\left(t_{1}\right)$ for some $t_{1}>0$. Let

$$
t_{0}=\max \left\{t<t_{1}: \mathrm{A}(t)=\widetilde{\mathrm{A}}(t)\right\} .
$$

Then

$$
\operatorname{Im}(\mathrm{Z}(t)+i \mathrm{~A}(t))>\operatorname{Im}(\mathrm{Z}(t)+i \tilde{\mathrm{~A}}(t))
$$

for $t \in\left(t_{0}, t_{1}\right)$ and, since D is monotone, $\mathrm{Z}(t)+i \mathrm{~A}(t) \in \mathrm{D}$ for $t \in\left(t_{0}, t_{1}\right)$. By (2.5), $\mathrm{A}(t)$ does not increase between $t_{0}$ and $t_{1}$, so $\mathrm{A}\left(t_{0}\right)=\mathrm{A}\left(t_{1}\right)$. Then

$$
0<\mathrm{A}\left(t_{1}\right)-\tilde{\mathrm{A}}\left(t_{1}\right)=\mathrm{A}\left(t_{0}\right)-\tilde{\mathrm{A}}\left(t_{1}\right) \leqq \mathrm{A}\left(t_{0}\right)-\tilde{\mathrm{A}}\left(t_{0}\right)=0
$$

This is a contradiction and it shows that $\mathrm{A}(t) \leqq \tilde{\mathrm{A}}(t)$ for all $t$. For similar reasons $\tilde{\mathrm{A}}(t) \leqq \mathrm{A}(t)$ and (2.4) is proved.

The mapping of $\{\mathrm{X}(t), 0 \leqq t<\mathrm{R}\}$ onto $\{\mathrm{Z}(t), 0 \leqq t<\infty\}$ is one-to-one a.s. so the above reasoning may be reversed. Suppose $\theta: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ is $\mathrm{C}^{2}$ and let $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ be the corresponding mapping defined as in Lemma 2.2. Suppose that $\mathbf{Z}(t)$ is a standard 2-dimensional Brownian motion, $\mathrm{Z}(0) \in \overline{\mathrm{D}}$ and $\mathrm{A}(t)$ is defined by (2.4). Let

$$
\begin{gathered}
\mathrm{W}(t)=\mathrm{Z}(t)+i \mathrm{~A}(t) \\
\sigma(t)=\int_{0}^{t}\left(\left(f^{-1}\right)^{\prime}(\mathrm{W}(s))\right)^{2} d s \\
\mathrm{X}(\sigma(t))=f^{-1}(\mathrm{~W}(t))
\end{gathered}
$$

The function $\left(f^{-1}\right)^{\prime}$ is bounded away from 0 and $\infty$ on bounded subsets of D , so $\sigma(t)$ may be infinite only if $\lim _{s \rightarrow t}|\mathrm{~W}(s)|=\infty$. Hence $\mathrm{X}(t)$ is welldefined until possible explosion, i.e., on a random time interval $[0, \mathrm{R})$ such that $\lim _{s \rightarrow \mathbf{R}}|\mathrm{X}(s)|=\infty$. The process X is a reflected Brownian motion in $D_{*}$ with the oblique angle of reflection $\theta$.

Step 2. - Now suppose that $\theta: \mathbb{R} \rightarrow[-\pi / 2, \pi / 2]$ is Borel measurable and that $\theta_{k}: \mathbb{R} \rightarrow(-\pi / 2, \pi / 2)$ are $\mathrm{C}^{2}$ functions which converge almost everywhere to $\theta$ as $k \rightarrow \infty$. One way of constructing such a sequence is as follows. Extend $\theta$ to be bounded, continuous and harmonic way on $\mathrm{D}_{*}$ and let $\theta_{k}(x)=\theta(x+i / k)$ for $x \in \partial \mathrm{D}_{*}$ [recall our assumption (2.1)].

Let $\mathbf{Z}$ be a standard 2-dimensional Brownian motion and let $f, \mathrm{~A}, \mathrm{~W}$, $\sigma$ and $X$ be defined in the same way as in the last part of Step 1. Define $f_{k}, \mathrm{~A}_{k}, \mathrm{~W}_{k}, \sigma_{k}$ and $\mathrm{X}_{k}$ in the analogous way relative to Z and $\theta_{k}$. We do
not know any more whether $\sigma$ takes finite values so at this point $\mathrm{X}(t)$ is well-defined only for some values of $t$.

We will prove in this step that $\mathrm{W}_{k}(t) \rightarrow \mathrm{W}(t)$ a.s. for every $t$. It is obvious that $\operatorname{Re} \mathrm{W}_{k}(t) \rightarrow \operatorname{Re} \mathrm{W}(t)$ as $k \rightarrow \infty$. It will suffice to show that $\mathrm{A}_{k}(t) \rightarrow \mathrm{A}(t)$ a.s. Fix some $t>0$.

Suppose that $\mathrm{A}(t)=\infty$. Let $t_{0}$ be the smallest $s$ such that $\mathrm{A}(s)=\infty$. Note that $t_{0}<t \mathrm{a} . \mathrm{s}$. We must have

$$
\left\{x \in \mathbb{C}: \operatorname{Im} x>a, \operatorname{Re} x \in\left(\operatorname{Re} Z\left(t_{0}\right)-\varepsilon, \operatorname{Re} Z\left(t_{0}\right)+\varepsilon\right)\right\} \cap \partial \mathbf{D} \neq \varnothing
$$

for all $a \in \mathbb{R}, \varepsilon>0$, and, by Lemma 2.3 (i),

$$
\left\{x \in \mathbb{C}: \operatorname{Im} x>a, \operatorname{Re} x \in\left(\operatorname{Re} \mathbf{Z}\left(t_{0}\right)-\varepsilon, \operatorname{Re} \mathbf{Z}\left(t_{0}\right)+\varepsilon\right)\right\} \cap \partial \mathrm{D}_{k} \neq \varnothing
$$

for all $a \in \mathbb{R}, \varepsilon>0$ and $k>k_{0}(a, \varepsilon)$. With probability $1, \mathrm{Z}(0, t)$ contains a closed loop around $\mathrm{Z}\left(t_{0}\right)$. It follows that $\lim \mathrm{A}_{k}\left(t_{0}+\varepsilon\right)=\infty \mathrm{a}$. s. for every $\varepsilon>0$ and since $t_{0}<t$ a. s., we have $\lim _{k \rightarrow \infty} \mathrm{~A}_{k}^{k \rightarrow \infty}(t)=\infty$ a.s.

Now suppose that $\mathrm{A}(t)<\infty$. Then $\mathrm{W}(t) \in \mathrm{D}$ a. s. because $\mathrm{Z}(0, t)$ contains a closed loop around $\mathrm{Z}(t)$. Let

$$
\begin{aligned}
\alpha(t) & =\inf \{s \in[0, t]: \mathrm{W}(s, t] \subset \mathrm{D}\}, \\
\alpha_{k}(t) & =\inf \left\{s \in[0, t]: \mathrm{W}_{k}(s, t] \subset \mathrm{D}\right\} .
\end{aligned}
$$

Note that $\alpha(t)<t$ a.s.
First we will show that $\tilde{\mathrm{A}}(\alpha(t)+) \geqq \mathrm{A}(\alpha(t))$ where $\tilde{\mathrm{A}}(s) \stackrel{\mathrm{df}}{=} \liminf _{k \rightarrow \infty} \mathrm{~A}_{k}(s)$. If $\alpha(t)=0$ then $\widetilde{\mathrm{A}}(\alpha(t)+) \geqq 0=\mathrm{A}(0)=\mathrm{A}(\alpha(t))$. Suppose that $\alpha(t)>0$ and so $x_{1} \stackrel{\text { df }}{=} \mathrm{Z}(\alpha(t))+i \mathrm{~A}(\alpha(t)) \in \partial \mathrm{D}$. Assume that $\tilde{\mathrm{A}}(\alpha(t)+)<\mathrm{A}(\alpha(t))$. Then, for some $\varepsilon_{0}>0$ and all $\varepsilon<\varepsilon_{0}$ and $\delta=\delta(\varepsilon)>0$ the interior of the set

$$
\begin{aligned}
\mathbf{B}(\varepsilon) \stackrel{\text { df }}{=}\{x \in \mathbb{C}: & \operatorname{Re} x= \\
& \operatorname{Re} Z(s), \\
& \operatorname{Im} x \geqq \operatorname{Im} Z(s)+\tilde{\mathrm{A}}(\alpha(t)+)+\varepsilon \text { for some } s \in(\alpha(t), \alpha(t)+\delta)\}
\end{aligned}
$$

is contained in
$\{x \in \mathbb{C}: \operatorname{Re} x=\operatorname{Re} Z(s)$,
$\operatorname{Im} x \geqq \operatorname{Im} Z(s)+\mathrm{A}_{k}(s)$ for some $\left.s \in(\alpha(t), \alpha(t)+\delta)\right\} \subset \mathrm{D}_{k}$ for infinitely many $k$. For each rational $s<t, \mathrm{~W}(s) \in \mathrm{D}$ a. s. because $\mathrm{Z}(0, t)$ contains a loop around $Z(s)$. Hence $B(\varepsilon) \cap D \neq \varnothing$ and, by Lemma 2.3 (ii), $\mathrm{B}(\varepsilon) \subset \mathrm{D}$ and

$$
\mathrm{Z}(\alpha(t)+\delta)+i(\widetilde{\mathrm{~A}}(\alpha(t)+)+\varepsilon) \in \overline{\mathrm{B}}(\varepsilon) \subset \overline{\mathrm{D}} .
$$

It follows that

$$
x_{2} \stackrel{\mathrm{df}}{=} \mathrm{Z}(\alpha(t))+i \tilde{\mathrm{~A}}(\alpha(t)+) \in \overline{\mathrm{D}} .
$$

Therefore the line segment $\mathbf{J}$ joining $x_{1}$ and $x_{2}$ lies in $\partial \mathrm{D}$. By Remark 2.1 (i), there exist vertical line segments $\mathrm{J}_{k} \subset \mathrm{D}_{k}^{c}$ with endpoints converging to $x_{1}$ and $x_{2}$ as $k \rightarrow \infty$.

Now we consider two cases. In the first case the range of the process

$$
\{\operatorname{ReZ}(s), s \in(\alpha(t)-\varepsilon, \alpha(t)+\varepsilon)\}
$$

intersects both intervals

$$
(\operatorname{Re} Z(\alpha(t))-\varepsilon, \operatorname{Re} Z(\alpha(t)))
$$

and

$$
(\operatorname{Re} Z(\alpha(t)), \operatorname{Re} Z(\alpha(t))+\varepsilon)
$$

for every $\varepsilon>0$. Then $\lim \inf \mathrm{A}_{k}(\alpha(t)) \geqq \mathrm{A}(\alpha(t))$ because $\mathrm{W}_{k}(t)$ cannot intersect $\mathrm{J}_{\boldsymbol{k}}$. In the opposite case $\operatorname{Re} Z$ has a local extremum at $\alpha(t)$. This event has probability zero, because $\alpha(t)$ corresponds to a vertical line segment J in $\partial \mathrm{D}$ and there is only a countable number of such line segments. This completes the proof that $\tilde{\mathrm{A}}(\alpha(t)+) \geqq \mathrm{A}(\alpha(t))$. We deduce that

$$
\liminf _{k \rightarrow \infty} \mathrm{~A}_{k}(t) \geqq \tilde{\mathrm{A}}(\alpha(t)+) \geqq \mathrm{A}(\alpha(t))=\mathrm{A}(t) \quad \text { a.s. }
$$

Now suppose that $\lim \sup \mathrm{A}_{k}(t)>\mathrm{A}(t)$. We will prove that this assump-

$$
k \rightarrow \infty
$$

tion leads to a contradiction. By passing to a subsequence if necessary we may assume that $\lim _{k \rightarrow \infty} \alpha_{k}(t) \stackrel{\text { df }}{=} \tilde{\alpha}(t)$ and $\lim _{k \rightarrow \infty} \mathrm{~A}_{k}\left(\alpha_{k}(t)\right) \stackrel{\text { df }}{=} \tilde{\mathrm{A}}(\tilde{\alpha}(t))$ exist. We have

$$
\mathrm{A}(\tilde{\alpha}(t)) \leqq \mathrm{A}(t)<\limsup _{k \rightarrow \infty} \mathrm{~A}_{k}(t)=\underset{k \rightarrow \infty}{\lim \sup } \mathrm{~A}_{k}\left(\alpha_{k}(t)\right)=\tilde{\mathrm{A}}(\tilde{\alpha}(t))
$$

so $\tilde{\mathrm{A}}(\tilde{\alpha}(t))-\mathrm{A}(\tilde{\alpha}(t))>0$. The fact that $\mathrm{W}(\tilde{\alpha}(t)) \in \overline{\mathrm{D}}$ and the definition of A imply that $\mathrm{W}(\tilde{\alpha}(t))+i b \in \mathrm{D}$ for all $b>0$. The sequence $\left\{\mathrm{W}_{k}\left(\alpha_{k}(t)\right)\right\}$ converges to $\mathrm{W}(\tilde{\alpha}(t))+i(\tilde{\mathrm{~A}}(\tilde{\alpha}(t))-\mathrm{A}(\tilde{\alpha}(t)))$ and its elements belong to $\partial \mathrm{D}_{k}$ so its limit belongs to $\mathrm{D}^{c}$, by Remark 2.1 (ii), which is a contradiction. This completes the proof that $\mathrm{W}_{k}(t) \rightarrow \mathrm{W}(t)$ a.s. for every $t$.

Step 3. - In order to simplify the argument, we will assume that $\mathrm{Z}(0) \in \partial \mathrm{D}$. There will be no loss of generality. Fix some $a>0$ and let

$$
\begin{gathered}
\mathrm{T}=\inf \left\{t \geqq 0: \operatorname{Im} f^{-1}(\mathrm{~W}(t)) \geqq a\right\} \\
\mathrm{S}=\sup \left\{s<\mathrm{T}: \operatorname{Im} f^{-1}(\mathrm{~W}(t))=0\right\} .
\end{gathered}
$$

It is possible that $\mathrm{T}=\infty$ with positive probability. First we will assume that $\mathrm{T}<\infty$ a. s . and show that

$$
\begin{equation*}
\underset{k \rightarrow \infty}{\lim \sup } \sup \left\{\operatorname{Im} f^{-1}\left(\mathrm{~W}_{k}(t)\right): 0 \leqq t \leqq \mathrm{~S}\right\}<a \quad \text { a.s. } \tag{2.6}
\end{equation*}
$$

Suppose that (2.6) is not satisfied. By passing to a subsequence if necessary we may assume that there exists a sequence $\left\{t_{k}\right\}$ such that $t_{k} \rightarrow t_{\infty} \in[0, \mathrm{~S}]$ and $\operatorname{Im} f^{-1} \mathrm{~W}_{k}\left(t_{k}\right) \rightarrow a_{1} \geqq a$. By assumption,

$$
\begin{equation*}
\operatorname{Im} f^{-1} \mathrm{~W}\left(t_{\infty}\right)<a \tag{2.7}
\end{equation*}
$$

Recall that the functions $\mathrm{A}_{k}(u)$ are monotone and they converge a.s. to a monotone function $\mathrm{A}(u)$ simultaneously for all rational $u$. It follows that with probability 1 ,

$$
\mathrm{A}\left(t_{\infty}-\right) \leqq \liminf _{k \rightarrow \infty} \mathrm{~A}_{k}\left(t_{k}\right) \leqq \limsup _{k \rightarrow \infty} \mathrm{~A}_{k}\left(t_{k}\right) \leqq \mathrm{A}\left(t_{\infty}+\right)
$$

If A is continuous at $t_{\infty}$ then

$$
\begin{aligned}
\operatorname{Im} f^{-1} \mathrm{~W}\left(t_{\infty}\right) & =\operatorname{Im} f^{-1}\left(\mathrm{Z}\left(t_{\infty}\right)+i \mathrm{~A}\left(t_{\infty}\right)\right) \\
& =\lim _{k \rightarrow \infty} \operatorname{Im} f^{-1}\left(\mathrm{Z}\left(t_{\infty}\right)+i \mathrm{~A}_{k}\left(t_{k}\right)\right) \\
& =\lim _{k \rightarrow \infty} \operatorname{Im} \mathrm{~W}_{k}\left(t_{k}\right)=a_{1}>a
\end{aligned}
$$

and this contradicts (2.7).
If A has a jump at $t_{\infty}$ then $\mathrm{W}\left(t_{\infty}\right) \in \partial \mathrm{D}$ and $\operatorname{Im} f^{-1}\left(\mathrm{~W}\left(t_{\infty}\right)\right)=0$. Since $\lim \sup \mathrm{A}_{k}\left(t_{k}\right) \leqq \mathrm{A}\left(t_{\infty}+\right)$, any cluster point of $\mathrm{W}_{k}\left(t_{k}\right)$ must be contained $k \rightarrow \infty$
in $\mathrm{D}^{c}$. It follows that $\operatorname{Im} f^{-1}\left(\mathrm{~W}_{k}\left(t_{k}\right)\right) \rightarrow 0$ which again is a contradiction. This completes the proof of (2.6).

We proceed to prove that $\lim _{k \rightarrow \infty} \sigma_{k}(t)=\sigma(t)$ a.s. for every $t$. Fix some $\varepsilon>0$ and define stopping times

$$
\begin{aligned}
& \mathrm{U}_{1}=0 \\
& \mathrm{~T}_{1}=\inf \left\{t \geqq 0: \operatorname{Im} f^{-1}(\mathrm{~W}(t)) \geqq \varepsilon\right\}, \\
& \mathrm{U}_{k}=\inf \left\{t \geqq \mathrm{~T}_{k-1}: \operatorname{Im} f^{-1}(\mathrm{~W}(t))=0\right\}, \quad k \geqq 2, \\
& \mathrm{~T}_{k}=\inf \left\{t \geqq \mathrm{U}_{k}: \operatorname{Im} f^{-1}(\mathrm{~W}(t)) \geqq \varepsilon\right\}, \quad k \geqq 2 .
\end{aligned}
$$

Let

$$
\sigma(s, t)=\int_{s}^{t}\left(\left(f^{-1}\right)^{\prime}(\mathrm{W}(u))\right)^{2} d u
$$

and define $\Gamma_{k}(s)$ by

$$
\Gamma_{k}\left(\sigma\left(\mathrm{~T}_{k}, t\right)\right)=f^{-1}(\mathrm{~W}(t)), \quad t \in\left(\mathrm{~T}_{k}, \mathrm{U}_{k+1}\right)
$$

Note that $\Gamma_{k}(s)$ is well defined for $s \in\left(0, \sigma\left(\mathrm{~T}_{k}, \mathrm{U}_{k+1}\right)\right)$ since $\sigma\left(\mathrm{T}_{k},.\right)$ is a strictly increasing function on ( $0, \sigma\left(\mathrm{~T}_{k}, \mathrm{U}_{k+1}\right)$ ). By the strong Markov property of W applied at $\mathrm{T}_{\boldsymbol{k}}$ 's and the conformal invariance of Brownian motion, the processes

$$
\left\{\operatorname{Im} \Gamma_{k}(s), 0 \leqq s \leqq \sigma\left(\mathrm{~T}_{k}, \mathrm{U}_{k+1}\right)\right\}, \quad k \leqq 1,
$$

are independent 1 -dimensional Brownian motions starting at $\varepsilon>0$ and killed upon hitting 0 . Let

$$
\tau_{\varepsilon}=\sum_{\substack{k \\ \mathrm{U}_{k+1} \leqq \mathrm{~S}}} \sigma\left(\mathrm{~T}_{k}, \mathrm{U}_{k+1}\right)
$$

It is evident that when $\varepsilon \rightarrow 0$ the distribution of $\tau_{\varepsilon}$ converges to the distribution (say, $\mathscr{L}$ ) of the starting time of the first excursion of height greater or equal to $a$ of reflected 1-dimensional Brownian motion. Moreover, $\tau_{\varepsilon} \rightarrow \sigma(S)$ a.s. as $\varepsilon \rightarrow 0$. Hence $\sigma(S)<\infty$ a.s. and by considering all $a>0$ we see that $\sigma$ is finite for all $t$.

Since $\mathrm{W}_{k}(t) \rightarrow \mathrm{W}(t)$ a.s. for every $t$, the Fatou lemma implies that $\lim \sup \sigma_{k}(t) \geqq \sigma(t)$ a.s. for all $t$. It follows from (2.6) that the starting $k \rightarrow \infty$
points $\mathrm{V}_{k}$ of the first excursion of height greater or equal to $a$ of $\operatorname{Im} \mathrm{X}_{k}$ converge a.s. to $\sigma(\mathrm{S})$ and by the same argument as in the case of $\sigma(\mathrm{S})$, each $\mathrm{V}_{k}$ has the distribution $\mathscr{L}$. This and $\lim \sup \sigma_{k}(\mathbf{S}) \geqq \sigma(\mathbf{S})$ imply that

$$
k \rightarrow \infty
$$

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \sigma_{k}(\mathbf{S})=\sigma(\mathbf{S}) \tag{2.8}
\end{equation*}
$$

If $\lim \sup \sigma_{k}(t) \geqq \sigma(t)+a$ then $\lim \sup \sigma_{k}(s) \geqq \sigma(s)+a$ for all $s \geqq t$. This and (2.8) applied for all $a>0$ show that $\lim _{k \rightarrow \infty} \sigma_{k}(t)=\sigma(t)$ for all $t$ a.s.

The case when $\mathrm{T}=\infty$ with positive probability may be treated in a similar way.
Step 4. - Now we will prove that $\mathrm{X}_{k} \rightarrow \mathrm{X}$ in $\mathrm{M}_{1}$ topology.
Since A is non-decreasing, it has at most a countable number of jumps. We have proved that $\mathrm{W}_{k}(t) \rightarrow \mathrm{W}(t)$ a.s. for every $t$. It follows that $\mathrm{W}_{k}(t) \rightarrow \mathrm{W}(t)$ for almost all continuity points of W , with probability 1 , by the Fubini theorem. Since $\sigma_{k}(t) \rightarrow \sigma(t)$ for all $t$ a.s., we see that $\mathrm{X}_{k}(t) \rightarrow \mathrm{X}(t)$ a.s. for every continuity point of X . This shows that the condition 2.3.3 (a) of Skorohod (1960) is satisfied. In view of his Theorem 2.4.1 all that remains to be shown is condition 2.4.1 (b).

We will assume that the condition fails, i.e., there exist sequences $\left\{t_{k}^{1}\right\}_{k \geqq 1},\left\{t_{k}^{2}\right\}_{k \geqq 1}$ and $\left\{t_{k}^{3}\right\}_{k \geqq 1}$ such that $t_{k}^{1} \leqq t_{k}^{2} \leqq t_{k}^{3},\left|t_{k}^{1}-t_{k}^{3}\right| \rightarrow 0$ and the distance $d_{k}$ from $\mathrm{X}_{k}\left(t_{k}^{2}\right)$ to the line segment $\left[\mathrm{X}_{k}\left(t_{k}^{1}\right), \mathrm{X}_{k}\left(t_{k}^{3}\right)\right]$ does not converge to 0 as $k \rightarrow \infty$. By passing to subsequences, we may assume that the following limits exist:

$$
\begin{gathered}
\lim _{k \rightarrow \infty} t_{k}^{1}=\lim _{k \rightarrow \infty} t_{k}^{2}=\lim _{k \rightarrow \infty} t_{k}^{3} \stackrel{\mathrm{df}}{=} t_{\infty}<\infty \\
\lim _{k \rightarrow \infty} X_{k}\left(t_{k}^{j}\right) \stackrel{\text { df }}{=} x_{j}, \quad j=1,2,3
\end{gathered}
$$

$$
\lim _{k \rightarrow \infty} d_{k} \stackrel{\mathrm{df}}{=} d
$$

If $\sigma^{-1}\left(t_{\infty}\right)$ is a point of continuity of W then it is easy to see that $x_{1}=x_{2}=x_{3}$ and $d=0$.

Suppose that W has a jump at $\sigma^{-1}\left(t_{\infty}\right)$ and let $\mathbf{J}$ be the line segment

$$
\left[\mathrm{Z}\left(\sigma^{-1}\left(t_{\infty}\right)\right)+i \mathrm{~A}\left(\sigma^{-1}\left(t_{\infty}\right)-\right), \mathrm{Z}\left(\sigma^{-1}\left(t_{\infty}\right)\right)+i \mathrm{~A}\left(\sigma^{-1}\left(t_{\infty}\right)\right)\right]
$$

It can be shown using arguments similar to those in the previous steps that any cluster points $z_{j}$ of $\mathrm{W}_{k}\left(\sigma^{-1}\left(t_{k}^{j}\right)\right), j=1,2,3$, can lie only on $\mathrm{J} \subset \partial \mathrm{D}$ and, moreover, by monotonicity of $\mathrm{A}_{k}$ 's, $z_{2}$ lies between $z_{1}$ and $z_{3}$. The mapping $f^{-1}$ takes J into a line segment on $\partial \mathrm{D}_{*}$ and it follows that $d=0$.

Remark 2.2. - Let us indicate how one may prove that X is a strong Markov process.

Let $Z$ and $A$ be defined as in Step 2 of the proof of Theorem 1.1. The Brownian motion Z is strong Markov and so is the vector $(\mathrm{Z}, \mathrm{A})$ since $\mathrm{A}(t)$ is defined in terms of $\mathrm{Z}(s), s \leqq t$. The process W is a function of $(\mathbf{Z}, \mathbf{A})$ and it is clear that the distribution of $\left\{\mathbf{W}(s), s \geqq t_{0}\right\}$ depends only on $\mathrm{W}\left(t_{0}\right)=\mathrm{Z}\left(t_{0}\right)+i \mathrm{~A}\left(t_{0}\right)$ but otherwise it does not depend on the values of $\mathrm{Z}\left(t_{0}\right)$ and $\mathrm{A}\left(t_{0}\right)$. By Theorem 13.5 of Sharpe (1988), W is strong Markov. The process X is obtained from W by first applying a one-toone mapping and then a time-change. According to Theorem 65.9 of Sharpe (1988) the second operation preserves the strong Markov property. In order to make this argument rigorous one would have to verify carefully the assumptions of Sharpe's theorems and analyze the extension of $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ to a univalent mapping of $\overline{\mathrm{D}}_{*}$ onto the Martin compactification of D . We postpone this discussion to a future article.

## 3. EXCEPTIONAL POINTS ON BROWNIAN PATHS

Evidently, the existence of D-points depends only on the local properties of D near the imaginary axis so we will assume without loss of generality that $|\operatorname{Re} z|<1$ for all $z \in \mathrm{D}$, i.e., $\tilde{\mathrm{D}}=\mathrm{D}$.

In this section, $\mathrm{P}^{z}$ and $\mathrm{E}^{z}$ will denote the distribution and expectation corresponding to the standard 2-dimensional Brownian motion starting from $z \in \mathbb{C}$ and Y will denote its trajectories.

Lemma 3.1.- The function $h$ defined in (1.3) satisfies

$$
\begin{equation*}
h(z)=\mathrm{E}^{z}\left(\left|\operatorname{Re} \mathrm{Y}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right)\right|\right) \quad \text { for all } z \in \mathrm{D} \tag{3.1}
\end{equation*}
$$

Proof. - Let $g(z)$ be the function on the right hand side of (3.1). Then $g(z)=h(z)=|\operatorname{Re} z|$ on $\partial \mathrm{D}$. It is easy to see that both $g$ and $h$ are bounded so it will suffice to show that they are harmonic in D. It is a standard fact that $g$ is harmonic in D and $h$ is obviously harmonic in $\mathrm{D} \backslash \mathrm{I}$. All that remains to be shown is that $h$ satisfies the mean-value property at each point $i b_{0} \in \mathrm{I}$. (In this proof $z$ will stand for a complex number and $b_{0}, b, r$ and $\varphi$ for real numbers.)

Fix some $i b_{0} \in \mathrm{I}$ and choose $r>0$ so that $i b_{0}+r e^{i \varphi} \in \mathrm{D}$ for every $\varphi$. Elementary calculations show that

$$
\begin{equation*}
\int_{0}^{2 \pi}\left|\operatorname{Re}\left(i b_{0}+r e^{i \varphi}\right)\right| \frac{d \varphi}{2 \pi}=\frac{2 r}{\pi} \tag{3.2}
\end{equation*}
$$

If $\left|b-b_{0}\right|>r$ then $\mathrm{G}_{\mathrm{D}}(., i b)$ is harmonic in $\left\{z:\left|z-i b_{0}\right| \leqq r\right\}$ and, by the mean value property,

$$
\begin{equation*}
\int_{0}^{2 \pi} \mathrm{G}_{\mathrm{D}}\left(i b_{0}+r e^{i \varphi}, i b\right) \frac{d \varphi}{2 \pi}=\mathrm{G}_{\mathrm{D}}\left(i b_{0}, i b\right) \tag{3.3}
\end{equation*}
$$

If $\left|b-b_{0}\right|<r$, write

$$
\mathrm{G}_{\mathrm{D}}(z, i b)=-\log |z-i b|+f(z)
$$

where $f$ is harmonic in D . Thus for $\left|b-b_{0}\right|<r$

$$
\begin{align*}
& \int_{0}^{2 \pi} \mathrm{G}_{\mathrm{D}}\left(i b_{0}+r e^{i \varphi}, i b\right) \frac{d \varphi}{2 \pi} \\
&= f\left(i b_{0}\right)-\int_{0}^{2 \pi} \log \left|i b_{0}+r e^{i \varphi}-i b\right| \frac{d \varphi}{2 \pi} \\
&=\mathrm{G}_{\mathrm{D}}\left(i b_{0}, i b\right)-\int_{0}^{2 \pi} \log \left|\frac{i b_{0}+r e^{i \varphi}-i b}{i\left(b_{0}-b\right)}\right| \frac{d \varphi}{2 \pi} . \tag{3.4}
\end{align*}
$$

Note that

$$
\log \left|\frac{i b_{0}+r e^{i \varphi}-i b}{i\left(b_{0}-b\right)}\right|=\log \left|-i e^{i \varphi}+\frac{r}{b_{0}-b}\right|
$$

and the function $z \rightarrow \log \left|-i z+\frac{r}{b_{0}-b}\right|$ is harmonic in $\{z:|z|<1\}$ since $r>\left|b_{0}-b\right|$. Hence the last integral in (3.4) is equal to $\log \left|\frac{r}{b_{0}-b}\right|$. We have

$$
\begin{equation*}
\int_{\left\{\left|b-b_{0}\right|<r\right\}} \log \left|\frac{r}{b_{0}-b}\right| \frac{d b}{\pi}=\frac{2 r}{\pi} . \tag{3.5}
\end{equation*}
$$

Now we combine (3.2)-(3.5) to obtain

$$
\begin{aligned}
\int_{0}^{2 \pi} h\left(i b_{0}\right. & \left.+r e^{i \varphi}\right) \frac{d \varphi}{2 \pi} \\
= & \int_{0}^{2 \pi}\left|\operatorname{Re}\left(i b_{0}+r e^{i \varphi}\right)\right| \frac{d \varphi}{2 \pi}+\int_{\mathrm{I}} \int_{0}^{2 \pi} \mathrm{G}_{\mathrm{D}}\left(i b_{0}+r e^{i \varphi}, i b\right) \frac{d \varphi}{2 \pi} \frac{d b}{\pi} \\
= & \frac{2 r}{\pi}+\int_{\mathrm{I}} \mathrm{G}_{\mathrm{D}}\left(i b_{0}, i b\right) \frac{d b}{\pi}-\int_{\left\{\left|b-b_{0}\right|<r\right\}} \log \left|\frac{r}{b_{0}-b}\right| \frac{d b}{\pi} \\
= & \left|\operatorname{Re}\left(i b_{0}\right)\right|+\int_{\mathrm{I}} \mathrm{G}_{\mathrm{D}}\left(i b_{0}, i b\right) \frac{d b}{\pi}=h\left(i b_{0}\right) .
\end{aligned}
$$

Remark 3.1. - The following probabilistic argument may elucidate the relationship between formulae (1.3) and (3.1). We leave it to the reader to supply the rigorous justification of the argument.

The local time of Y on I [say, L(t)] is also the local time of $|\operatorname{Re} Y|$ at 0 . The process $\mathrm{V}(t) \stackrel{\mathrm{df}}{=}|\operatorname{Re} \mathrm{Y}(t)|-\mathrm{L}(t)$ is a Brownian motion and hence a martingale. By the optional sampling theorem and Theorem V. 1 of Revuz (1970)

$$
\begin{aligned}
\mathrm{E}^{x}\left|\operatorname{Re} \mathrm{Y}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right)\right| & =\mathrm{E}^{x}\left(\left|\operatorname{Re} \mathrm{Y}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right)\right|-\mathrm{L}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right)+\mathrm{L}\left(\mathrm{~T}_{\left.\mathrm{D}^{c}\right)}\right)\right. \\
& =\mathrm{E}^{x} \mathrm{~V}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right)+\mathrm{E}^{x} \mathrm{~L}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right) \\
& =\mathrm{E}^{x} \mathrm{~V}(0)+\mathrm{E}^{x} \mathrm{~L}\left(\mathrm{~T}_{\mathrm{D}^{c}}\right) \\
& =|\operatorname{Re} x|+\int_{\mathrm{I}} \mathrm{G}_{\mathrm{D}}(x, y) \frac{d y}{\pi} .
\end{aligned}
$$

Let $D_{1}=\{z \in \mathbb{C}: \operatorname{Re} z>0\}$. Recall that $D$ stands for a monotone domain and $d_{\varepsilon}$ was defined in Section 1. See Doob (1984) or Burdzy (1987) for the definition of minimal thinness and its relationship to Brownian motion.

Lemma 3.2. - The set $\mathrm{D}^{c} \cap \mathrm{D}_{1}$ is minimal thin in $\mathrm{D}_{1}$ at 0 if and only if there is an $\varepsilon>0$ such that

$$
\begin{equation*}
\int_{-1}^{0} d_{\varepsilon}(r) r^{-2} d r<\infty \tag{3.6}
\end{equation*}
$$

Proof. - (i) Suppose that (3.6) is satisfied for some $\varepsilon>0$ and let

$$
\mathrm{D}_{2}=\left\{x \in \mathbb{C}: \operatorname{Re} x>d_{\varepsilon}(\operatorname{Im} x)\right\} .
$$

By Theorems 9.1 (b) (ii) and 9.2 (i) and (iii) of Burdzy (1987), $D_{1} \backslash D_{2}$ is minimal thin in $D_{i}$ at 0 . Since for some neighborhood $U$ of 0 we have

$$
\mathrm{D}^{c} \cap \mathrm{D}_{1} \cap \mathrm{U} \subset\left(\mathrm{D}_{1} \backslash \mathrm{D}_{2}\right) \cap \mathrm{U}
$$

$D^{c} \cap D_{1}$ is also minimal thin in $D_{1}$ at 0 .
(ii) Now suppose that (3.6) is not satisfied for any $\varepsilon>0$. Let

$$
\mathbf{M}_{k}=\left\{x \in \mathbf{D}^{c} \cap \mathbf{D}_{1}:|x| \leqq 1 / k\right\}
$$

and let $\tilde{d}_{k}: \mathbb{R} \rightarrow \mathbb{R}$ be the smallest non-increasing non-negative function such that

$$
\mathbf{M}_{k} \subset\left\{x \in \mathbb{C}: \operatorname{Re} x \leqq \widetilde{d}_{k}(\operatorname{Im} x)\right\}
$$

Let

$$
\mathrm{D}_{k}=\left\{x \in \mathbb{C}: \operatorname{Re} x>\widetilde{d}_{k}(\operatorname{Im} x)\right\} .
$$

Note that $\tilde{d}_{k}=d_{1 / k}$ so

$$
\int_{-1}^{0} \tilde{d}_{k}(r) r^{-2} d r=\infty
$$

for every $k$. By Theorems 9.1 and 9.2 of Burdzy (1987), for each $k$, the set $D_{1} \backslash D_{k}$ is not minimal thin in $D_{1}$ at 0 .

For $x \in \overline{\mathrm{D}}_{1}, \mathrm{P}_{f}^{x}$ will denote the distribution of a Brownian motion Z in $\mathrm{D}_{1}$ starting from $x$ and conditioned by the harmonic function $f(x) \stackrel{\text { df }}{=} \operatorname{Re} x$. For $t>0$, let

$$
\mathrm{T}(t)=\inf \left\{s>t: \mathrm{Z}(s) \in \mathrm{D}_{k}^{c}\right\} .
$$

Since the set $D_{1} \backslash D_{k}$ is not minimal thin in $D_{1}$ at 0 it follows that $\inf \left\{s>0: \mathrm{Z}(s) \in \mathrm{D}_{k}^{c}\right\}=0, \mathrm{P}_{f}^{0}$-a.s. [Doob (1984, Theorem 3.III.3)] and, therefore, there exists $t=t(k)>0$ such that $\mathrm{P}_{f}^{0}(\mathrm{~T}(t)<\infty)>1 / 2$.

For each $x \in \mathrm{D}_{1} \backslash \mathrm{D}_{k}$, there is an $a \geqq \operatorname{Re} x$ such that

$$
\{z \in \mathbb{C}: \operatorname{Im} z \leqq \operatorname{Im} x, \operatorname{Re} z=a\} \subset \mathbf{D}^{c} \cap \mathbf{M}_{k}
$$

By symmetry, the process Z under $\mathrm{P}_{f}^{x}$, will hit the line $\{z \in \mathbb{C}: \operatorname{Re} z=a\}$ with equal probability below and above $a+i \operatorname{Im} x$. Hence by the strong Markov property applied at $\mathrm{T}(t)$, the $\mathrm{P}_{f}^{0}$-chance of hitting $\mathrm{D}^{c} \cap \mathrm{M}_{k}$ is at least $1 / 4$. Let $k \rightarrow \infty$ to see that

$$
\inf \left\{t>0: \mathrm{Z}(t) \in \mathrm{D}^{c}\right\}=0
$$

with $\mathrm{P}_{f}^{0}$-probability greater or equal to $1 / 4$. By the $0-1$ law, this probability is equal to 1. According to Theorem 3.III. 3 of Doob (1984), the set $D_{1} \backslash D$ is minimal thin in $D_{1}$ at 0 .

Proof of Theorem 1.2. - If D-points exist with positive probability then they exist with probability 1 because Brownian motion returns to the imaginary axis for arbitrarily large times and it is strong Markov. The same remark applies to right D-points.
(i) Step 1. - We start with an estimate of the expectation of a stopping time. Recall that we assume that $|\operatorname{Re} x|<1$ for $x \in \mathrm{D}$. Let $\mathrm{T}_{1}=\min (1, \mathrm{~T}(\mathrm{~J}))$, where J is the imaginary axis. We will show that there
exist constants $0<c_{1} \leqq c_{2}<\infty$ such that for all $x \in \overline{\mathrm{D}}$

$$
\begin{equation*}
c_{1}|\operatorname{Re} x| \leqq \mathrm{E}^{x} \mathrm{~T}_{1} \leqq c_{2}|\operatorname{Re} x| \tag{3.7}
\end{equation*}
$$

The density of $\mathrm{T}(\mathrm{J})$ under $\mathrm{E}^{x}$ is given by [Theorem 7.5.3 of Karlin and Taylor (1975)]

$$
f_{x_{1}}(t)=x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} \exp \left(-x_{1}^{2} / 2 t\right)
$$

where $x_{1}=|\operatorname{Re} x|$. Then

$$
\begin{aligned}
\mathrm{E}^{x} \mathrm{~T}_{1} & \geqq \int_{x_{1} / 2}^{1} x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} \exp \left(-x_{1}^{2} / 2 t\right) t d t \\
& \geqq x_{1} \int_{x_{1}^{2} / 2}^{1}\left(2 \pi t^{3}\right)^{-1 / 2} e^{-1} t d t \\
& \geqq x_{1} \int_{1 / 2}^{1}\left(2 \pi t^{3}\right)^{-1 / 2} e^{-1} t d t=c_{1} x_{1} .
\end{aligned}
$$

On the other hand, with $u=t / x_{1}^{2}$,

$$
\begin{aligned}
\mathrm{E}^{x} \mathrm{~T}_{1} \leqq & \left(\int_{0}^{x_{1}^{2}}+\int_{x_{1}^{2}}^{1}+\int_{1}^{\infty}\right) x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} \exp \left(-x_{1}^{2} / 2 t\right) t d t \\
\leqq & \int_{0}^{x_{1}^{2}} x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} \exp \left(-x_{1}^{2} / 2 t\right) t d t \\
& \quad+\int_{x_{1}^{2}}^{1} x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} t d t+\int_{1}^{\infty} x_{1}\left(2 \pi t^{3}\right)^{-1 / 2} d t \\
= & \int_{0}^{1}(2 \pi)^{-1 / 2} u^{-1 / 2} \exp (-1 / 2 u) x_{1}^{2} d u+\left.x_{1}(2 \pi)^{-1 / 2} 2 t^{1 / 2}\right|_{t=x_{1}^{2}} ^{t=1}+c_{3} x_{1} \\
= & c_{4} x_{1}^{2}+c_{5} x_{1}\left(1-x_{1}\right)+c_{3} x_{1} \leqq c_{2} x_{1} .
\end{aligned}
$$

Step 2. - In this step we will estimate the expectation of another stopping time. Let

$$
\begin{aligned}
& \mathrm{T}_{2}=\inf \left\{t>\mathrm{T}\left(\mathrm{D}^{c}\right): \operatorname{Re} \mathrm{Y}(t)=0\right\}, \\
& \mathrm{T}_{3}=\inf \left(1, \mathrm{~T}_{2}\right), \\
& \mathrm{T}_{4}=\inf \left(\mathrm{T}_{2}, \mathrm{~T}\left(\mathrm{D}^{c}\right)+1\right) .
\end{aligned}
$$

We will prove that for some $0<c_{1}, c_{2}<\infty$ we have

$$
c_{1} h(x) \leqq \mathrm{E}^{x} \mathrm{~T}_{3} \leqq c_{2} h(x)
$$

for $x \in \mathrm{D} \cap \mathrm{J}$. First, observe that $\mathrm{T}_{3} \geqq \mathrm{~T}_{4}-\mathrm{T}\left(\mathrm{D}^{c}\right)$. Use this fact, the strong Markov property applied at $\mathrm{T}\left(\mathrm{D}^{c}\right)$, (3.7) and Lemma 3.1 to see that

$$
\begin{align*}
\mathrm{E}^{x} \mathrm{~T}_{3} & \geqq \mathrm{E}^{x}\left(\mathrm{~T}_{4}-\mathrm{T}\left(\mathrm{D}^{c}\right)\right) \\
& =\mathrm{E}^{x}\left(\mathrm{E}^{\mathrm{Y}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)}(\mathrm{T}(\mathrm{~J}) \wedge 1)\right)  \tag{3.8}\\
& \geqq \mathrm{E}^{x}\left(c_{1}\left|\operatorname{ReY}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)\right|\right)=c_{1} h(x) .
\end{align*}
$$

We also have

$$
\begin{equation*}
\mathrm{T}_{3} \leqq \mathrm{~T}_{4}=\mathrm{T}\left(\mathrm{D}^{c}\right)+\left(\mathrm{T}_{4}-\mathrm{T}\left(\mathrm{D}^{c}\right)\right) \tag{3.9}
\end{equation*}
$$

The same reasoning which leads to (3.8) yields

$$
\begin{align*}
\mathrm{E}^{x}\left(\mathrm{~T}_{4}-\mathrm{T}\left(\mathrm{D}^{c}\right)\right) & =\mathrm{E}^{x}\left(\mathrm{E}^{\mathrm{Y}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)}(\mathrm{T}(\mathrm{~J}) \wedge 1)\right)  \tag{3.10}\\
& \leqq \mathrm{E}^{x}\left(c_{2}\left|\operatorname{Re} \mathrm{Y}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)\right|\right)=c_{2} h(x)
\end{align*}
$$

The process $(\operatorname{Re} \mathrm{Y}(t))^{2}-t$ is a martingale so by the optional sampling theorem we have for $a<\infty$

$$
\mathrm{E}^{x}\left(\operatorname{ReY}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right) \wedge a\right)\right)^{2}=\mathrm{E}^{x}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right) \wedge a\right)
$$

By the dominated convergence theorem and the monotone convergence theorem we can let $a \rightarrow \infty$. Since $|\operatorname{Re} z| \leqq 1$ for $z \in \partial \mathrm{D}$, we have by Lemma 3.1,

$$
\begin{align*}
\mathrm{E}^{x} \mathrm{~T}\left(\mathrm{D}^{c}\right) & =\mathrm{E}^{x}\left(\operatorname{Re} \mathrm{Y}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)\right)^{2}  \tag{3.11}\\
& \leqq \mathrm{E}^{x}\left|\operatorname{ReY}\left(\mathrm{~T}\left(\mathrm{D}^{c}\right)\right)\right|=h(x) .
\end{align*}
$$

By combining (3.9)-(3.11) we obtain $\mathrm{E}^{x} \mathrm{~T}_{3} \leqq c_{3} h(x)$.
Step 3. - Choose $\delta>0$ so small that the set

$$
\mathrm{M} \stackrel{\mathrm{df}}{=}\{x \in \mathrm{D}: \operatorname{Im} x \geqq 1-\delta,|\operatorname{Re} x| \leqq \delta\}
$$

is contained in $\mathbf{D}$. Let $\Theta$ denote the usual shift operator for the Markov process $Y$ and let

$$
\mathrm{A}_{0}=\left\{t \in[0,1]: \operatorname{Re} \mathrm{Y}(t)=0, \mathrm{~T}(\mathrm{M}+\mathrm{Y}(t)){ }^{\circ} \Theta_{t}+1<\mathrm{T}\left(\mathrm{D}^{c}+\mathrm{Y}(t)\right){ }^{\circ} \Theta_{t}\right\}
$$

Fix some $\varepsilon>0$. Let $\mathrm{T}_{0}=0$ and

$$
\begin{gathered}
\mathrm{U}_{k}=\inf \left\{t>\mathrm{T}_{k}: \mathrm{Y}(t) \in \mathrm{D}^{c}+\mathrm{Y}\left(\mathrm{~T}_{k}\right)-\varepsilon i\right\}, \quad k \geqq 0, \\
\mathrm{~T}_{k}=\inf \left\{t>\mathrm{U}_{k-1}: \operatorname{ReY}(t)=\operatorname{ReY}\left(\mathrm{T}_{k-1}\right)\right\} \wedge\left(\mathrm{T}_{k-1}+1\right), \quad k \geqq 1, \\
\mathrm{~B}_{k}=\left\{\mathrm{T}\left(\mathrm{M}+\mathrm{Y}\left(\mathrm{~T}_{k}\right)-\varepsilon i\right) \circ \Theta_{\mathrm{T}_{k}}+1<\mathrm{T}\left(\mathrm{D}^{c}+\mathrm{Y}\left(\mathrm{~T}_{k}\right)-\varepsilon i\right)^{\circ} \Theta_{\mathrm{T}_{k}}\right\}, \\
\mathrm{A}_{\varepsilon}=\bigcup_{k \geqq 0}\left(\left\{\mathrm{~T}_{k} \leqq 1\right\} \cap \mathrm{B}_{k}\right) .
\end{gathered}
$$

A moment's thought confirms that $A_{0} \subset A_{\varepsilon}$ for every $\varepsilon>0$ so $A_{0} \subset \bigcap_{\varepsilon>0} A_{\varepsilon}$. Now suppose that $\bigcap_{\varepsilon>0} A_{\varepsilon}$ holds, in particular, for every positive integer $m$, the event $\mathrm{A}_{1 / m}$ holds. Let $k=k(m)$ be the smallest $k$ such that $\left\{\mathrm{T}_{k} \leqq 1\right\} \cap \mathrm{B}_{k}$ holds with $\varepsilon=1 / m$. By compactness, a subsequence of $\left\{\mathrm{T}_{k(m)}\right\}_{m \geqq 1}$ converges to a point $t_{0} \in[0,1]$. It is easy to verify that $t_{0}$ satisfies the definition of $A_{0}$. We conclude that $\cap A_{\varepsilon} \subset A_{0}$ and, therefore,
$\bigcap_{\varepsilon>0} \mathrm{~A}_{\varepsilon}=\mathrm{A}_{0}$.

Let $\mathrm{B}=\left\{\mathrm{T}(\mathrm{M})+1<\mathrm{T}\left(\mathrm{D}^{c}\right)\right\}$. The functions $x \rightarrow \mathrm{P}^{x}(\mathrm{~B})$ and $x \rightarrow \mathrm{G}_{\mathrm{D}}(i, x)$ are positive and harmonic in $\mathrm{D} \backslash \mathrm{M}$ and vanish on $\partial \mathrm{D}$. Thus, for $x \in \mathrm{D}$ in a neighborhood of 0 ,

$$
c_{1} \mathrm{G}_{\mathrm{D}}(i, x) \leqq \mathrm{P}^{x}(\mathrm{~B}) \leqq c_{2} \mathrm{G}_{\mathrm{D}}(i, x)
$$

To see this, one may map conformally D onto the half-plane $\mathrm{D}_{*}$. The conformal invariance of harmonic functions and the boundary Harnack principle applied in $D_{*}$ [see, e. g., Theorem 4 of Dahlberg (1977) or Theorem 2.1 of Burdzy (1987)] yield the above inequality.

By the strong Markov property applied at $\mathrm{T}_{k}$ and the translation invariance of Brownian motion we have $\mathrm{P}^{0}\left(\mathrm{~B}_{k}\right)=\mathrm{P}^{\varepsilon i}(\mathrm{~B})$ and, consequently,

$$
\begin{equation*}
c_{1} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) \leqq \mathrm{P}^{0}\left(\mathrm{~B}_{k}\right) \leqq c_{2} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) \tag{3.12}
\end{equation*}
$$

Let N be the largest $k$ such that $\mathrm{T}_{k} \leqq 1$.

$$
\begin{aligned}
\mathrm{P}^{0}\left(\mathrm{~A}_{\varepsilon}\right) & =\mathrm{P}^{0}\left(\cup\left(\left\{\mathrm{~T}_{k} \leqq 1\right\} \cap \mathrm{B}_{k}\right)\right) \\
& =\sum_{k \geqq 1} \mathrm{P}^{0}\left(\left\{\mathrm{~T}_{k} \leqq 1\right\} \cap \mathrm{B}_{k}\right) \quad \text { (disjoint events) } \\
& =\sum_{k \geqq 1} \mathrm{P}^{0}\left(\mathrm{~T}_{k} \leqq 1\right) \mathrm{P}^{0}\left(\mathrm{~B}_{k}\right) \quad \text { (independence) } \\
& =\mathrm{E}^{0}(\mathrm{~N}+1) \mathrm{P}^{0}\left(\mathrm{~B}_{k}\right)
\end{aligned}
$$

In view of (3.12),

$$
\begin{equation*}
c_{1} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) \mathrm{E}^{0}(\mathrm{~N}+1) \leqq \mathrm{P}^{0}\left(\mathrm{~A}_{\varepsilon}\right) \leqq c_{2} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) \mathrm{E}^{0}(\mathrm{~N}+1) \tag{3.13}
\end{equation*}
$$

The random variables $T_{k}-T_{k-1}$ are i.i.d., positive and

$$
\mathrm{T}_{\mathrm{N}+1}=\sum_{k=1}^{\mathrm{N}+1}\left(\mathrm{~T}_{k}-\mathrm{T}_{k-1}\right)
$$

By Wald's identity [Feller (1971), XVIII 2]

$$
\mathrm{E}^{0} \mathrm{~T}_{\mathrm{N}+1}=\mathrm{E}^{0}(\mathrm{~N}+1) \mathrm{E}^{0}\left(\mathrm{~T}_{1}-\mathrm{T}_{0}\right)
$$

This and the fact that $1 \leqq \mathrm{~T}_{\mathrm{N}+1} \leqq 2$ a. s . imply that

$$
1 / \mathrm{E}^{0}\left(\mathrm{~T}_{1}-\mathrm{T}_{0}\right) \leqq \mathrm{E}^{0}(\mathrm{~N}+1) \leqq 2 / \mathrm{E}^{0}\left(\mathrm{~T}_{1}-\mathrm{T}_{0}\right)
$$

The $\mathrm{P}^{0}$-distribution of $\mathrm{T}_{1}-\mathrm{T}_{0}$ is the same as the $\mathrm{P}^{\varepsilon i}$-distribution of the stopping time $\mathrm{T}_{3}$ of Step 2 so

$$
c_{3} / h(\varepsilon i) \leqq \mathrm{E}^{0}(\mathrm{~N}+1) \leqq c_{4} / h(\varepsilon i)
$$

This and (3.13) imply that

$$
\begin{equation*}
c_{5} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) / h(\varepsilon i) \leqq \mathrm{P}^{0}\left(\mathrm{~A}_{\varepsilon}\right) \leqq c_{6} \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) / h(\varepsilon i) \tag{3.14}
\end{equation*}
$$

Now we see that (1.4) is equivalent to $\mathrm{P}^{0}\left(\mathrm{~A}_{0}\right)>0$.

Step 4. - It is now easy to finish the proof. First of all, it is evident that $\mathrm{P}^{0}\left(\mathrm{~A}_{0}\right)=0$ if and only if with probability 1 , there are no points $\mathrm{Y}(t)$ such that $\operatorname{Re} \mathrm{Y}(t)=0$ and $\mathrm{Y}\left(t, t_{1}\right) \subset \mathrm{D}+\mathrm{Y}(t)$ for some $t_{1}>t$. Finally, the last condition is equivalent to almost sure non-existence of D-points, by the time reversal.
(ii) We will assume without loss of generality that the Brownian motion Y starts from i. Assume that (1.5) is satisfied. Let

$$
\begin{gathered}
\mathrm{K}=\{z \in \mathbb{C}: \operatorname{Re} z=0, \operatorname{Im} z \leqq-1\}, \\
\mathrm{T}=\sup \left\{t<\mathrm{T}_{\mathrm{K}}: \operatorname{Re} \mathrm{Y}(t)=0\right\} .
\end{gathered}
$$

Note that $\mathrm{T}<\mathrm{T}_{\mathrm{K}}<\infty$ a.s. Hence $\left\{\mathrm{Y}(t), \mathrm{T}<t<\mathrm{T}_{\mathrm{K}}\right\}$ is an excursion of Y from the imaginary axis. Since (1.5) holds, the set $D^{c}+Y\left(T_{K}\right)$ is minimal thin in $\{z \in \mathbb{C}: \operatorname{Re} z>0\}$ at $\mathrm{Y}\left(\mathrm{T}_{\mathrm{K}}\right)$, by Lemma 3.2. By the time reversal, the local path behavior of Brownian excursions is identical at both ends of an excursion, so Lemma 8.1 (ii) of Burdzy (1987) shows that $\left\{\mathrm{Y}(t), \mathrm{T}<t<\mathrm{T}_{\mathrm{K}}\right\}$ does not intersect $\mathrm{D}^{\mathrm{c}}+\mathrm{Y}\left(\mathrm{T}_{\mathrm{K}}\right)$ immediately prior to $\mathrm{T}_{\mathrm{K}}$ a.s. Now standard arguments show that, with positive probability, the whole piece $Y\left[0, T_{K}\right)$ of the Brownian path is disjoint with $D^{c}+Y\left(T_{K}\right)$ and this means that $\mathrm{Y}\left(\mathrm{T}_{\mathrm{K}}\right)$ is a right D -point with positive probability.

Now suppose that (1.5) does not hold. Assume that $\mathrm{Y}\left(t_{0}\right)$ is a right D-point. Then $\mathrm{Y}\left(t_{0}\right)$ is the endpoint of an excursion of Y from the imaginary axis. By Lemma 8.1 (ii) of Burdzy (1987) and Lemma 3.2, every excursion $\left\{\mathrm{Y}(s), s_{1}<s<s_{2}\right\}$ of Y from the imaginary axis must hit $\mathrm{D}^{c}+\mathrm{Y}\left(s_{2}\right)$ immediately prior to $s_{2}$ and this holds for all excursions simultaneously with probability 1 since there is only a countable number of excursions. In particular, Y must hit $\mathrm{D}^{\boldsymbol{c}}+\mathrm{Y}\left(t_{0}\right)$ immediately prior to $t_{0}$ and, therefore, $\mathrm{Y}\left(t_{0}\right)$ is not a right D -point. This proves our assertion.

Remark 3.2. - It is not hard to see that $\mathrm{A}_{\varepsilon_{1}} \subset \mathrm{~A}_{\varepsilon_{2}}$ for $\varepsilon_{1} \leqq \varepsilon_{2}$. Hence, $\mathrm{P}^{0}\left(\mathrm{~A}_{\varepsilon}\right)$ is a monotone function of $\varepsilon$ and (3.14) implies that (1.6) is equivalent to the negation of (1.4).

Proof of Theorem 1.3. - (i) First suppose that the prime end 0 is well defined in D and, therefore, we may find $f$ with $f(0)=0$. Then Theorem 1.3 (i) follows immediately from the construction of reflected Brownian motion given in the proof of Theorem 1.1.

Now suppose that we cannot find D and $f$ as in Lemma 2.2 and such that the prime end 0 is well defined in D . Let $f$ be any function defined as in Lemma 2.2 and let $\gamma$ be the prime end in D which is the image of 0 under $f$. If $\gamma$ corresponded to an $x_{0} \in \partial \mathrm{D}$ such that $x_{0}+a i \in \mathrm{D}$ for all $a>0$ then we could define $f_{1}(y) \stackrel{\mathrm{df}}{=} f\left(y-x_{0}\right)$ and we would have $f_{1}(0)=0$, contrary to our assumption. Hence $\gamma$ must correspond either to a point $x_{0} \in \partial \mathrm{D}$ which is in the middle of a vertical line segment $\mathrm{J} \subset \partial \mathrm{D}$ or $\gamma$ is the limit of $x_{0}+a i$ as $a \rightarrow-\infty$ for some $x_{0} \in \mathrm{D}$. In either case the process

W defined in the proof of Theorem 1.1 does not aproach $\gamma$ a.s. and therefore X does not approach 0 a.s.
(ii) We can prove as in part (i) that the a.s. existence of right D-points is equivalent to

$$
\begin{aligned}
& \mathrm{P}\left(\mathrm{~T}_{\{0\}}<\infty \quad \text { and } \quad \exists \varepsilon>0:\right. \\
& \quad \operatorname{ReX}(t)>0 \text { for all } t \in\left(\mathrm{~T}_{\{0\}}-\varepsilon, \mathrm{T}_{\{0\}}\right) \\
& \quad \text { such that } \operatorname{Im} \mathrm{X}(t)=0)>0 .
\end{aligned}
$$

It has been shown in Step 1 of the proof of Corollary 1.3 of Burdzy and Marshall (1992) that the last condition holds if and only if (1.9) is satisfied.

## 4. SETS OF NON-POLAR POINTS FOR REFLECTED BROWNIAN MOTION

Proof of Theorem 1.4. - (i) Let $f$ be the function defined in Lemma 2.2 and $\mathrm{D}=f\left(\mathrm{D}_{*}\right)$. For $\lambda$-almost all $x \in \partial \mathrm{D}_{*}$ such that $\theta(x) \in(-\pi / 2, \pi / 2)$ both $f(x)$ and $f^{\prime}(x)$ exist and $\arg f^{\prime}(x)=\theta(x)$. Consider a point $x$ with all these properties and assume without loss of generality that $f(x)=0$. In view of Theorem 1.3, in order to show that $x \notin \mathscr{A}$ it will suffice to prove that D points do not exist a.s.

We have the following explicit formula [Doob (1984), 1. VIII (9.3)].

$$
\mathrm{G}_{\mathrm{D}_{*}}(x, y)=\log \frac{|x-y-2 i \operatorname{Im} x|}{|x-y|}=\log \left|1-2 i \frac{\operatorname{Im} x}{x-y}\right|
$$

For every Stolz angle

$$
\mathscr{S}_{\varphi}=\left\{z \in \mathrm{D}_{*}:|z|<1,|\arg z-\pi / 2|<\varphi\right\}
$$

with $\varphi<\pi / 2$ and $x, y \in \mathscr{S}_{\varphi},|y| \geqq 2|x|$, we have

$$
1-2 \frac{|x| / \arccos \varphi}{|y|} \leqq\left|1-2 i \frac{\operatorname{Im} x}{x-y \mid}\right| \leqq 1+2 \frac{|x|}{|y| / 2} .
$$

Hence

$$
c_{1}|x / y| \leqq \mathrm{G}_{\mathrm{D}_{*}}(x, y) \leqq c_{2}|x / y|
$$

for $x, y \in \mathscr{S}_{\varphi},|x / y| \leqq \alpha$ where $\alpha=\min (1 / 8,(\arccos \varphi) / 4)$. We will take $\varphi=\pi / 2-(\pi / 2-|\theta(x)|) / 2$. Then it follows from the definition of the angular derivative and conformal invariance of the Green function that

$$
c_{3}|x / y| \leqq \mathrm{G}_{\mathrm{D}}(x, y) \leqq c_{4}|x / y|
$$

for $x, y \in f\left(\mathscr{S}_{\varphi}\right),|y| \geqq 2|x| / \alpha$, in particular, for $x, y \in \mathrm{I},|y| \geqq 2|x| / \alpha$. Hence, for small $\varepsilon>0$,

$$
h(\varepsilon i)=\int_{\mathrm{I}} \mathrm{G}_{\mathrm{D}}(\varepsilon i, y) d y \geqq \int_{2 \varepsilon / \alpha}^{1} c_{3} \varepsilon / y d y=-c_{3} \varepsilon \log (2 \varepsilon / \alpha)
$$

We see that $h(\varepsilon i) / \mathrm{G}_{\mathrm{D}}(i, \varepsilon i) \geqq-c_{5} \log (2 \varepsilon / \alpha) \rightarrow \infty$ as $\varepsilon \rightarrow 0$ and by Theorem 1.2 D-points do not exist a.s.
(ii) Consider any interval $(a, b), a<b$. If

$$
\lambda(\{x \in(a, b):|\theta(x)|<\pi / 2\})>0
$$

then $(a, b) \backslash \mathscr{A} \neq \varnothing$ by part (i).
If $\theta(x)=-\pi / 2$ for almost all $x \in(a, b)$ then the function $f$ defined in Lemma 2.2 maps ( $a, b$ ) onto a vertical line segment $\mathrm{K} \subset \partial \mathrm{D}$. It follows from the construction of the reflected Brownian motion [see especially (2.4)] that with probability $1, \mathrm{~W}$ never takes values in K , so X does not take values in $(a, b)$ a.s. The same argument takes care of the case when $\theta(x)=\pi / 2$ for $x \in(a, b)$.

Finally, suppose that $|\theta(x)|=\pi / 2$ for almost all $x \in(a, b)$ but $\theta$ is not almost everywhere constant in $(a, b)$. Note that D lies above the graph of a function $g$, possibly taking infinite values. The interval $(a, b)$ is mapped onto a part of $\partial \mathrm{D}$. If $g$ is monotone on any interval between $\operatorname{Re} f(a)$ and $\operatorname{Re} f(b)$ then $\theta$ must take values of one sign on the corresponding interval and the previous argument may be applied. If $g$ has a jump then $\partial \mathrm{D}$ contains a vertical line segment and again we may apply the previous argument. If $g$ is not monotone on a subinterval of $(\operatorname{Re} f(a), \operatorname{Re} f(b))$ and $g$ has no jumps on this interval then $g$ must have a finite local minimum, say at $\operatorname{Re} f\left(x_{0}\right), x_{0} \in(a, b)$. We may suppose without loss of generality that $f\left(x_{0}\right)=0$. One can easily check using Theorem 1.2 that $\mathrm{D}_{*}$-points do not exist; this has been already proved in Burdzy [1989, Theorem 2.1 (iii)]. Since the property of having D-points is monotone in D and $\mathrm{D} \subset \mathrm{D}_{*}$ locally, we conclude that D -points do not exist a.s. Hence, $x_{0} \notin \mathscr{A}$.
(iii) Step 1. - For a domain $\mathrm{D}_{1}, x \in \mathbb{C}$ and $\varepsilon>0$ let $h=h_{x, \varepsilon, \mathrm{D}_{1}, \eta}$ be the smallest Lipschitz function with the Lipschitz constant $\eta>0$ such that

$$
\begin{aligned}
& \left\{y \in \partial \mathrm{D}_{1}:|\operatorname{Re} x-\operatorname{Re} y| \leqq \varepsilon,|\operatorname{Im} x-\operatorname{Im} y| \leqq \varepsilon\right\} \\
& \subset\{y \in \mathbb{C}:|\operatorname{Re} x-\operatorname{Re} y| \leqq \varepsilon,|\operatorname{Im} x-\operatorname{Im} y| \leqq \varepsilon, \operatorname{Re} y \leqq h(\operatorname{Im} y)\}
\end{aligned}
$$

Next, let

$$
\begin{gathered}
\mathrm{K}_{1}\left(x, \varepsilon, \mathrm{D}_{1}, \eta\right)=\{y \in \mathbb{C}:|\operatorname{Re} x-\operatorname{Re} y|<\varepsilon \\
|\operatorname{Im} x-\operatorname{Im} y|<\varepsilon, \operatorname{Re} y>h(\operatorname{Im} y)\} \\
\mathrm{K}\left(x, \varepsilon, \mathrm{D}_{1}, \eta, c\right)=e^{-c i} \mathrm{~K}_{1}\left(e^{c i} x, \varepsilon, e^{c i} \mathrm{D}_{1}, \eta\right)
\end{gathered}
$$

Suppose that D is a monotone domain and let $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ be defined as in Lemma 2.1. Let

$$
\mathrm{M}_{a, b}=\left\{x \in \partial \mathrm{D}_{*}: \arg f^{\prime}(x) \in(a, b)\right\}
$$

and suppose that $\lambda\left(\mathrm{M}_{a, b}\right)>0$ for some $-\pi / 2+\eta_{1}<a<b<\pi / 2-\eta_{1}, \eta_{1}>0$. Then we can argue as in the proof of Theorem 2.6 (i) of Burdzy (1989) to show that for $\eta=\eta_{1} / 2$ there exit $x \in \mathbb{C}, \varepsilon>0$ and $c \in \mathbb{R}$ such that the harmonic measure of

$$
\mathrm{N} \stackrel{\mathrm{df}}{=} f\left(\mathrm{M}_{a, b}\right) \cap \partial \mathrm{K}(x, \varepsilon, \mathrm{D}, \eta, c)
$$

is strictly positive in a connected component of $\mathrm{K}(x, \varepsilon, \mathrm{D}, \eta, c)$. Since $\mathrm{K}(x, \varepsilon, \mathrm{D}, \eta, c)$ is a union of Lipschitz domains, Theorem 1 of Dahlberg (1977) implies that the linear measure of N is positive and, therefore the Lebesgue measure of the projection of N onto $\partial \mathrm{D}_{*}$ is positive, by our choice of $\eta_{1}$ and $\eta$.

Hence, if $\lambda\left(\mathrm{M}_{-\pi / 2, \pi / 2}\right)>0$ then there exists a set $\mathrm{F} \subset \mathbb{R}$ of positive measure such that if $x \in \partial \mathrm{D}$ and $\operatorname{Re} x \in \mathrm{~F}$ then there is $\delta=\delta(x)>0$ such that

$$
\{y \in \mathbb{C}: \arg (y-x) \in(\pi / 2-\delta, \pi / 2+\delta),|y-x|<\delta\} \subset \mathrm{D} .
$$

One can show in a similar way that if for some set $\mathrm{A} \subset \partial \mathrm{D}_{*}$ we have $\lambda(\mathrm{A})>0$ and $\left|\arg f^{\prime}(x)\right|=\pi / 2$ for $x \in \mathrm{~A}$ then the projection of $f(\mathrm{~A})$ on the imaginary axis has a strictly positive measure.

Step 2. - Let $\mathbf{B}^{1}(t)$ and $\mathbf{B}^{2}(t)$ be independent standard 1-dimensional Brownian motions starting from 0 and let D be a (random) monotone domain defined by

$$
\mathrm{D}=\left\{x \in \mathbb{C}: \operatorname{Re} x \geqq 0, \operatorname{Im} x>\mathbf{B}^{1}(\operatorname{Re} x)\right\},
$$

Let $f: \mathrm{D}_{*} \rightarrow \mathrm{D}$ be the function defined in Lemma 2.1 and let X be the reflected Brownian motion in D with the oblique angle of reflection $\theta(x) \stackrel{\mathrm{df}}{=} \arg f^{\prime}(x)$. We will show that with probability 1 (with respect to the distribution of $\mathbf{B}^{1}$ and $\mathbf{B}^{2}$ ) the function $\theta$ satisfies Theorem 1.4 (iii) (a)(b).

For a fixed $t>0, j=1,2$,

$$
\limsup _{\varepsilon \rightarrow 0+} \frac{\mathrm{B}^{j}(t+\varepsilon)-\mathrm{B}^{j}(t)}{\varepsilon}=\limsup _{\varepsilon \rightarrow 0-} \frac{\mathrm{B}^{j}(t+\varepsilon)-\mathrm{B}^{j}(t)}{-\varepsilon}=\infty \quad \text { a.s. }
$$

By the Fubini theorem, the set of all $x \in \partial \mathrm{D}$ such that there is $\delta>0$ with

$$
\{y \in \mathbb{C}: \arg (y-x) \in(\pi / 2-\delta, \pi / 2+\delta),|y-x|<\delta\} \subset \mathrm{D}
$$

has measure 0 , with probability 1. By Step $1,\left|\arg f^{\prime}(x)\right|=\pi / 2$ for almost all $x \in \partial \mathrm{D}_{*}$.

Consider an $x_{0} \in \partial \mathrm{D}_{*}$ such that $\left|\arg f^{\prime}\left(x_{0}\right)\right|=\pi / 2$ and suppose it is mapped by $f$ onto $t_{0}+\mathrm{B}^{1}\left(t_{0}\right) i$ for some $t_{0}>0$ (the argument is analogous for $\left.t_{0}<0\right)$. It follows from the definition of the angular derivative $f^{\prime}\left(x_{0}\right)$ that $\mathrm{B}^{1}(s)<\mathrm{B}^{1}\left(t_{0}\right)$ for all $s \in\left(t_{0}, t_{0}+\varepsilon\right)$ or for all $s \in\left(t_{0}-\varepsilon, t_{0}\right)$ for some $\varepsilon>0$. Hence, $t_{0}$ is the endpoint of an excursion of $\mathrm{B}^{1}$ from $a \stackrel{\text { df }}{=} \mathrm{B}^{1}\left(t_{0}\right)$. For definiteness, suppose that $\mathrm{B}^{1}(s)<\mathrm{B}^{1}\left(t_{0}\right)$ for all $s \in\left(t_{0}, t_{0}+\varepsilon\right)$.

It is well known that excursions of 1 -dimensional Brownian motion have the same local properties as the 3-dimensional Bessel process [see, e. g., Williams (1979)]. By Theorem 3.3 (ii) of Shiga and Watanabe (1973),

$$
\begin{equation*}
\limsup _{\varepsilon \rightarrow 0+} \frac{B^{1}\left(t_{0}+\varepsilon\right)-B^{1}\left(t_{0}\right)}{\varepsilon^{3 / 4}} \leqq-1 \tag{4.1}
\end{equation*}
$$

and since there is only a countable number of excursions of $\mathrm{B}^{1}$ from $a$, the same property holds for all endpoints of all excursion from $a$ a.s. If we move D so that $t_{0}+\mathrm{B}^{1}\left(t_{0}\right) i$ becomes 0 , we see that (4.1) implies that (1.5) (or its mirror image) holds and therefore right (or left) $\mathrm{D}-\left(t_{0}+\mathrm{B}^{1}\left(t_{0}\right) i\right)$-points exist, by Theorem 1.2. Thus the reflected Brownian motion X hits $x_{0}$ with positive probability and approaches it from one side, by Theorem 1.3.

By the Fubini theorem, the set of levels $a$ such that (4.1) is violated for at least one excursion from $a$ has measure 0 . It follows from Step 1 that the set of $x \in \partial \mathrm{D}_{*}$ which cannot be approached from one side with positive probability has measure 0 .

Since $\mathrm{B}^{j}$ is not monotone on any interval with probability 1 , it follows that $\theta(x)=\arg f^{\prime}(x)$ cannot have one sign almost everywhere on any interval $(a, b), a<b$. This completes the proof of $(a)$ and $(b)$.

Suppose that reflected Brownian motions $\mathrm{X}_{1}$ and $\mathrm{X}_{2}$ in $\mathrm{D}_{*}$ have oblique angles of reflection $\theta_{1}$ and $\theta_{2}$ which correspond to two monotone domains $D_{1}$ and $D_{2}$ such that $D_{1} \subset D_{2}$. It follows from our construction of reflected Brownian motion that if $X_{1}$ is recurrent then $X_{2}$ is recurrent as well. If $D_{1}=D_{*}$ then $X_{1}$ is a reflected Brownian motion with the normal reflection on $\partial \mathrm{D}_{*}$ and it is well known that it is recurrent. Now replace in the definition of D Brownian motions $\mathrm{B}^{j}(t), j=1$. 2, with Brownian motions with drift $\widetilde{\mathbf{B}}^{j}(t) \stackrel{\text { df }}{=} \mathrm{B}^{j}(t)-t$. Then $\mathrm{D}_{*}+a i \subset \mathrm{D}$ for some $a>0$ since $\widetilde{\mathbf{B}}^{j}(t) \rightarrow-\infty$ as $t \rightarrow \infty$. The reflected Brownian motion X corresponding to this new domain D is recurrent and it also satisfies $(a)$ and $(b)$ since the constant drift does not change the local properties of the 1-dimensional Brownian motion.

The recurrence of X implies $c 1$ ). Condition ( $c 2$ ) follows from ( $c 1$ ) and the Fubini theorem.

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(Manuscript received July 17, 1991;
revised June 18, 1992.)

