# Multiple Points for a Process in $R^{2}$ with Stable Components 

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## Section 1: Introduction

The object of this paper is to investigate the existence and Hausdorff dimension of the set of multiple points of the sample paths of a type of Markov process in $R^{2}$ with stationary independent increments. To define the process, let $X_{\alpha_{i}, 1}(t) \equiv X_{i}(t)$ denote a stable process of index $\alpha_{i}$ in $R^{1}$ for $i=1,2$; assume the two processes independent and let $X(t) \equiv\left(X_{1}(t), X_{2}(t)\right)$. We call $X(t)$ a process in $R^{2}$ with stable components. If the $\alpha_{i}$ exceed 1 , each component of $X(t)$ is point recurrent but $X(t)$ can be shown (see [6] or [10]) to be transient. If $\alpha_{1}=\alpha_{2}, X(t)$ is stable. Hence we assume that $1<\alpha_{2}<\alpha_{1} \leqq 2$. The situation for $\alpha_{2} \leqq 1$ or for a process with stable components in $R^{3}$ is discussed briefly in Section 6.

Taylor [12] has studied the existence and Hausdorff dimension of multiple points for symmetric stable processes of index $\alpha$ in $R^{n}$. Our basic argument will resemble his, although significant modifications are necessary. After stating the necessary preliminaries in Section 2 we use some potential theory in Section 3 to obtain estimates of various hitting probabilities with respect to $X(t)$ and then proceed to obtain conditions on $k, \alpha_{1}$ and $\alpha_{2}$ for existence (Section 4) of points of multiplicity $k$. In Section 5 we compute the Hausdorff dimension of the multiple points. The problem of multiple points for processes with more than 2 stable components or for ones with a component in a higher dimension remains open, as does the problem for multiple points (or at least double points) for a general process with stationary and independent increments.

## Section 2: Preliminaries

The $n$-dimensional characteristic function of a stable process $X_{\alpha, n}(t)$ of index $\alpha \neq 1$ in $R^{n}$ has the form $\exp [t \psi(y)]$, where

$$
\psi(y)=i(a, y)-\delta|y|^{\alpha} \int_{S_{n}} w_{\alpha}(y, \theta) \mu(d \theta),
$$

with $a \in R^{n}, \delta>0$,

$$
w_{\alpha}(y, \theta)=[1-i \operatorname{sgn}(y, \theta) \tan \pi \alpha / 2]|(y /|y|, \theta)|^{\alpha},
$$

and $\mu$ a probability measure on the surface of the unit sphere $S_{n}$ in $R^{n}$ [8]. We assume $a=0, \delta=1$, and that $\mu$ is not supported by a proper subspace of $R^{n}$. If $\mu$ is uniform, the process is said to be symmetric.

[^0]When speaking of a stable process $X_{\alpha, n}$ we will always write the two subscripts to indicate index $\alpha$ and dimension $n$ except when the process is actually a component of a process with stable components. When the latter occurs, we use a single subscript $i$ to indicate the $i$-th component. The symbol $X(t)$ will be understood to refer to the process $\left(X_{1}, X_{2}\right)$ as defined in Section 1. We will use the notation and methods of [6] to determine various hitting probabilities in Section 3. To do this we use the facts that stable densities $p_{\alpha, n}(t)$ of the type being considered are positive, continuous, bounded in $x$ for each fixed $t$, and satisfy the scaling property [11]:

$$
\begin{equation*}
p_{\alpha, n}(t, x)=r^{n / \alpha} p_{\alpha, n}\left(r t, r^{1 / \alpha} x\right) \quad \text { for } r>0 \tag{2.1}
\end{equation*}
$$

The density, $p(t, x)$, of $X(t)$ will have the form:

$$
\begin{equation*}
p(t, x)=p_{x_{1}, 1}\left(t, x_{1}\right) p_{\alpha_{2}, 1}\left(t, x_{2}\right), \quad \text { where } x=\left(x_{1}, x_{2}\right) . \tag{2.2}
\end{equation*}
$$

It will be assumed that all processes being considered have been defined so as to have sample functions $X(t, \omega)$ which are right continuous and have left limits everywhere [1]. In addition, we assume that the strong Markov property holds and that $X(0)=0$ with probability one.

We define the hitting probability $\Phi(x, E)$ of a planar Borel set $E$ starting from $x \in R^{2}$ by:

$$
\Phi(x, E)=P^{x}[X(t) \in E \text { for some } t>0]
$$

Often in writing $P^{0}$ [ ] we simply use $P\left[\right.$ ]. The set of points in $R^{2}$ that the sample path $X(t, \omega)$ hits in the time interval $[a, b]$ is denoted by $L(a, b ; \omega)$, where $0 \leqq a<b \leqq \infty$. Likewise, for any positive integer $k \geqq 2$, we denote by $L^{k}(a, b ; \omega)$ the set of points hit $k$ times in the time interval $[a, b]$. Sometimes we use $E_{k}(\omega)$ to denote the set $L^{k}(0, \infty ; \omega)$ of $k$-multiple points of the path $X(t, \omega)$. Positive constants whose values remain fixed throughout the discussion will be introduced in order and denoted by $c_{1}, \ldots, c_{35}$. The letter $c$ will be used as a positive constant which can vary in size from statement to statement or line to line.

## Section 3: Potential Theory and Hitting Probabilities

Taylor [13] gives a brief background of the potential theory which we will need. If we let $p_{1}\left(t, x_{1}\right)$ be the density of $X_{1}(t)$ and $p_{2}\left(t, x_{2}\right)$ that of $X_{2}(t)$, the density of $X(t)$ is their product and will be denoted by $p(t, x)$, where $x=\left(x_{1}, x_{2}\right) . U(y)$, the kernel of the process, is given by

$$
\begin{equation*}
U(y)=\int_{0}^{\infty} p(t, y) d t \tag{3.1}
\end{equation*}
$$

and converges for all $y \neq(0,0)$ in $R^{2}$.
If we let $\mu$ be any measure defined on Borel subsets of compact sets $E$ in the plane, the potential at $x$ of the measure $\mu$ on $E$ is $W_{\mu}(x)=\int_{E} U(y-x) \mu(d y)$. The capacity of $E$ is zero iff $W_{\mu}$ is unbounded for every $\mu$ for which $\mu(E)>0$. If there are some $\mu$ such that $W_{\mu}$ is bounded, we define the capacity of $E$ with respect to $X(t)$ by:

$$
C(E)=\sup \left\{\mu(E): W_{\mu}(x) \leqq 1 \text { for all } x\right\}
$$

When $E$ is compact, this supremum is actually attained for a measure $v$, called the capacitory measure on $E$.

Finally, we denote the hitting probability of a compact set $E$ starting from $x$ by:

$$
\Phi(x, E)=P^{x}\{X(t) \in E \text { for some } t>0\}
$$

Hitting probabilities are then given in terms of the kernel and the capacitory measure:

$$
\Phi(x, E)=\int_{E} U(y-x) v(d y) .
$$

Our method will be to obtain bounds on the kernel, the capacity of rectangular regions, and thereby upon hitting probabilities of rectangles. Some of the results are interesting in themselves, but since they will be used to prove later theorems, they will be stated as lemmas.

Lemma 3.1. Let $y=\left(y_{1}, y_{2}\right)$ and $U(y)$ be the kernal as defined by (3.1). Then positive constants $c_{1}, \ldots, c_{4}$, independent of $y$, can be found such that:

$$
\begin{array}{ll}
\frac{c_{1}}{\left|y_{2}\right|^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}} \leqq U(y) \leqq \frac{c_{2}}{\left|y_{2}\right|^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}} & \text { if }\left|y_{1}\right|^{\alpha_{1}} \leqq\left|y_{2}\right|^{\alpha_{2}} \\
\frac{c_{3}}{\left|y_{1}\right|^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}} \leqq U(y) \leqq \frac{c_{4}}{\left|y_{1}\right|^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}} & \text { if }\left|y_{2}\right|^{\alpha_{2}} \leqq\left|y_{1}\right|^{\alpha_{1}} .
\end{array}
$$

Before proceeding to the proof, observe that once the lemma is established, the two inequalities on the right can be made to hold in the entire plane by using $\left|y_{2}\right|^{-1} \leqq\left|y_{1}\right|^{-\alpha_{1} / \alpha_{2}}$ in the first inequality and $\left|y_{1}\right|^{-1} \leqq\left|y_{2}\right|^{-\alpha_{2} / \alpha_{1}}$ in the second.

Proof. The lower bound estimates are derived (p.265) in [6], and the upper bounds are proven (Lemma 3.1) in [7].

Lemma 3.2. Let $R_{a, b}$ be a rectangle centered at $(0,0)$ with sides of length $2 a(2 b)$ parallel to the $y_{1}\left(y_{2}\right)$ axis. For positive $\delta \leqq 1$ denote by $R(\delta)$ the resulting rectangle when $a=\delta^{\alpha_{2}}$ and $b=\delta^{\alpha_{1}}$, and denote by $S(\delta)$ the square which has $a=b=\delta$. Then positive constants $c_{5}$ and $c_{6}$ whose values do not depend upon a or $b$ can be found such that:

$$
c_{5} \max \left\{a^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}, b^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}\right\} \leqq C\left(R_{a, b}\right) \leqq c_{6} \max \left\{a^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}, b^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}\right\}
$$

In particular we have bounds of the form

$$
\delta^{x_{1}+\alpha_{2}-x_{1} \alpha_{2}} \quad \text { and } \quad \delta^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}
$$

for $C(R(\delta))$ and $C(S(\delta))$ respectively.
Proof. As in the proofs of Lemmas 3.1 and 3.3 of [6] we have:

$$
1 \geqq C\left(R_{a, b}\right) \min _{y \in R_{a, b}} U(y),
$$

so that

$$
C\left(R_{a, b}\right) \leqq\left\{\min _{y \in R_{a, b}} U(y)\right\}^{-1}=\max \left\{c_{1}^{-1} b^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}, c_{3}^{-1} a^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}\right\}
$$

upon application of Lemma 3.1. To obtain the lower bound, let $\mu_{L}$ denote Lebesgue measure in $R^{2}$ and observe that:

$$
\begin{aligned}
\int_{R_{a, b}} U(y-x) \mu_{L}(d y) & \leqq c_{2} \int_{-b}^{b} \int_{-a}^{a}\left|y_{2}-x_{2}\right|^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}} d y_{1} d y_{2} \\
& =2 a c_{2} \int_{-b-x_{2}}^{b-x_{2}}|u|^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}} d u=c a b^{\alpha_{2}-\alpha_{2} / \alpha_{1}}
\end{aligned}
$$

where $c$ is a positive constant independent of $a, b$ and $x$. The last integration can be justified by using the $c_{r}$ inequality (p. 155 of [9]) in the cases $0<-b-x_{2}$ and $b-x_{2}<0$, and the fact that $\alpha_{1}>1$ for the case $\left|x_{2}\right| \leqq b$. Now define the measure $\mu^{*}$ on Borel subsets $E$ of $R_{a, b}$ by:

$$
\mu^{*}(E)=c^{-1} a^{-1} b^{-\alpha_{2}+\alpha_{2} / \alpha_{1}} \cdot \mu_{L}(E)
$$

Then $\int_{R_{a, b}} U(y-x) \mu^{*}(d y) \leqq 1$ for all $x$ and:

$$
C\left(R_{a, b}\right) \geqq \mu^{*}\left(R_{a, b}\right)=c^{-1} b^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}} .
$$

In the same manner we can use the fact that $\alpha_{2}>1$ to show that

$$
C\left(R_{a, b}\right) \geqq c^{-1} a^{1-\alpha_{1}+\alpha_{1} / \alpha_{2}}
$$

Our next lemma gives estimates on the delayed hitting probabilities of the rectangles $R_{a, b}$. For rectangles $R_{a, b}$ as defined in Lemma 3.2, and for $T>0$ and all $x \in R^{2}$ let

$$
Q\left(x, R_{a, b}, T\right)=P^{x}\left[X(t) \in R_{a, b} \text { for some } t \geqq T\right] .
$$

We can then proceed as in the proof of Lemma 3.2 of [6] to establish
Lemma 3.3. Let $Q\left(x, R_{a, b}, T\right)$ be as given above. Then positive constants $c_{7}$ and $c_{8}$, independent of $x, a, b$, and $T$, can be found such that:
(i) $Q\left(x, R_{a, b}, T\right) \leqq c_{7} C\left(R_{a, b}\right) T^{1-1 / \alpha_{1}-1 / \alpha_{2}}$,
(ii) $Q\left(x, R_{a, b}, T\right) \geqq c_{8} C\left(R_{a, b}\right) T^{1-1 / \alpha_{1}-1 / \alpha_{2}}$ if $x \in R_{a . b}$ and $T \geqq \max \left\{(2 a)^{\alpha_{1}},(2 b)^{\alpha_{2}}\right\}$.

We can now use both parts of Lemma 3.3 to reason as in the proof of Lemma 4.2 of [6] to obtain

Lemma 3.4. Choose $T_{1} \geqq \max \left\{(2 a)^{\alpha_{1}},(2 b)^{\alpha_{2}}\right\}$ and $x \in R_{a, b}$. Then positive constants $c_{9}>1$ and $c_{10}$, independent of $a, b$ and $x$, can be found such that when $T_{2} \geqq c_{9} T_{1}$ :

$$
P^{x}\left[X(t) \in R_{a, b} \text { for some } t \in\left[T_{1}, T_{2}\right]\right] \geqq c_{10} C\left(R_{a, b}\right) T_{1}^{1-1 / \alpha_{1}-1 / \alpha_{2}}
$$

Various estimates of hitting probabilities of rectangles can be developed depending upon the location of the rectangle and starting point, but we only require

Lemma 3.5. Suppose that $t \notin R_{a, b}$ and that $|x|+d<1$, where $d$ is the length of the diagonal of $R_{a, b}$. Then positive constants $c_{11}, c_{12}$ and $c_{13}$ (independent of $a, b$ and $x$ ) can be found such that:
(i) $P^{x}\left[X(t) \in R_{a, b}\right.$ for some $\left.t>0\right] \geqq c_{11} C\left(R_{a, b}\right) \cdot(|x|+d)^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}}$,
(ii) $P^{x}\left[X(t) \in R_{a, b}\right.$ for some $t$ in $\left.[0, T]\right] \geqq c_{12} C\left(R_{a, b}\right) \cdot(|x|+d)^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}}$ whenever $T \geqq c_{13}(|x|+d)^{\alpha_{2}}$.

Proof. From Lemma 3.1 we obtain the estimate $U(y) \geqq c_{1}|y|^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}}$ if $|y| \leqq 1$. If $v$ is the capacitory measure on $R_{a, b}$, (i) is a consequence of:

$$
\Phi\left(x, R_{a, b}\right)=\int_{R_{a, b}} U(y-x) v(d y) \geqq c_{1} \int_{R_{a, b}}(|x|+d)^{\alpha_{2}-1-\alpha_{2} / \alpha_{1}} v(d y)
$$

To prove (ii) we take the difference of the estimates (i) in Lemmas 3.5 and 3.3 and choose $c_{13}$ appropriately.

## Section 4: Existence Theorems

We now determine conditions for the existence of $k$-multiple points for $X(t)$. For any positive integer $k$ a point $x$ in $R^{2}$ is said to be a $k$-multiple point of the sample path $X(t, \omega)$ if there are $k$ distinct times $0<t_{1}<\cdots<t_{k}$ such that $X\left(t_{i}, \omega\right)=x$ for $i=1,2, \ldots, k$. In [12], Taylor showed that symmetric stable processes $X_{a, n}(t)$ of index $\alpha<n$ have (with probability one) $k$-multiple points if $k(n-\alpha)<n$ and that (with probability one) such points do not exist if $k(n-\alpha) \geqq n$. In Theorems 1 and 2 below we shall see fairly direct proofs for existence and nonexistence, respectively, of $k$-multiple points except for the critical case $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)=\alpha_{1}+\alpha_{2}$; we conclude with a remark as to how the critical case could be handled.

Theorem 1. $X(t)$ has (with probability one) $k$-multiple points if $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)$ $<\alpha_{1}+\alpha_{2}$.

Proof. We could follow Taylor's (Section 4 of [12]) argument to prove existence, but we postpone the elaborate argument required and state that existence of $k$-multiple points is assured under our hypothesis once we show that this same hypothesis implies that the set of $k$-multiple points has positive Hausdorff dimension. This is done in Theorem 4 of the next section.

Theorem 2. With probability one, $k$-multiple points do not exist if $\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)$ $>\alpha_{1}+\alpha_{2}$.

Proof. Cover $R^{2}$ with abutting rectangles $R_{i}, i=1,2,3, \ldots$, congruent to $R(\delta)$ of Lemma 3.2 and whose long sides are parallel to the horizontal axis. We shall estimate how many of these rectangles are hit by time $t=2$, and then show that the probability of $(k-1)$ delayed returns to at least one such rectangle approaches zero as $\delta \rightarrow 0$.

The first step is to let $T(\delta, s)$ denote the amount of time the process spends in the rectangle $R(\delta)$ up to time $s$. $T(\delta, s)$ is a random variable whose expectation we estimate by the methods of Pruitt and Taylor [10] (p. 278). Let $s \geqq \delta^{\alpha_{1} \alpha_{2}}$. Then:

$$
\begin{align*}
E[T(\delta, s) & \equiv \int_{0}^{s} P[X(t) \in R(\delta)] d t \\
& =\int_{0}^{s} P\left[\left|X_{1}(1)\right| \leqq \delta^{\alpha_{2}} t^{-1 / \alpha_{1}}\right] P\left[\left|X_{2}(1)\right| \leqq \delta^{\alpha_{1}} t^{-1 / \alpha_{2}}\right] d t  \tag{4.1}\\
& \geqq \int_{0}^{\delta^{\alpha_{1} \alpha_{2}}} P\left[\left|X_{1}(1)\right| \leqq 1\right] P\left[\left|X_{2}(1)\right| \leqq 1\right] d t=c_{14} \delta^{\alpha_{1} \alpha_{2}}
\end{align*}
$$

If we apply the same reasoning of Pruitt and Taylor [10] in their Lemma 6.1 we can estimate $E[M(\delta, s)]$, where $M(\delta, s)$ denotes the number of rectangles from
$\left\{R_{i}\right\}_{i=1}^{\infty}$ that the process enters by time $s$. Their methods, along with (4.1) above, give (for $\delta \leqq 1$ and $s \geqq \delta^{\alpha_{1} \alpha_{2}}$ ):

$$
\begin{equation*}
E[M(\delta, s)] \leqq 2 c_{15} s\{E[T(\delta / 3, s)]\}^{-1} \leqq c_{16} s \delta^{-\alpha_{1} \alpha_{2}} \tag{4.2}
\end{equation*}
$$

Define events $A_{i}$ and $B_{i}$ (independent) by: $A_{i}=\left[X(t)\right.$ enters $R_{i}$ at some time $t \in[1,2]], B_{i}=\left[X(t)\right.$ enters $R_{i}$ at $(k-1)$ times: $\left.t_{1}<\cdots<t_{k-1}\right]$, where $t_{j} \in[2 j+1$, $2 j+2], j=1,2, \ldots, k-1$. Let $I_{i}$ be the indicator function of event $A_{i}$. Then: $P[k-1$ returns (during the specified time intervals) to at least one of the $R_{i}$ ]

$$
\begin{aligned}
& =\sum_{i} P\left[A_{i} \cap B_{i}\right]=\sum_{i} P\left[A_{i}\right] P\left[B_{i}\right] \\
& \leqq c \delta^{\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)(k-1)} \sum_{i} E\left[I_{i}\right] \quad \text { by Lemmas } 3.2 \text { and } 3.3 \\
& \leqq c \delta^{\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)(k-1)} E[M(\delta, 2)] \\
& \leqq c \delta^{\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)(k-1)-\alpha_{1} \alpha_{2}} \quad \text { by }(4.2) .
\end{aligned}
$$

The hypothesis $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)>\alpha_{1}+\alpha_{2}$ guarantees that this final quantity approaches zero as $\delta \rightarrow 0$. Consequently, the event

$$
\left[X\left(t_{j}\right)=X(t) \text { for some } t \in[1,2] \text { and } t_{j} \in[2 j+1,2 j+2], j=1,2, \ldots, k-1\right]
$$

has probability zero. Now apply the same reasoning to any sequence

$$
0<r_{1}<r_{2}<\cdots<r_{2 k} \leqq 1
$$

of $2 k$ distinct rational numbers $r_{j}$ in $[0,1]$, to conclude that for such a sequence with probability one there are not $k$ time instants $t_{1}<\cdots<t_{k}$ for which

$$
t_{j} \in\left[r_{2 j-1}, r_{2 j}\right], \quad 1 \leqq j \leqq k
$$

and:

$$
X\left(t_{1}, \omega\right)=X\left(t_{2}, \omega\right)=\cdots=X\left(t_{k}, \omega\right)
$$

Since the rationals in $[0,1]$ are countable, so is the collection $\left\{K_{\alpha}\right\}$ of all sets of $2 k$ distinct rationals taken from [0, 1]. Thus the $K$ 's can be indexed $\left\{K_{\alpha}\right\}=\left\{K_{i}\right\}_{i=1}^{\infty}$. Define the event $C_{i}, i=1,2,3, \ldots$ by:

$$
C_{i}=\left[X\left(t_{1}\right)=\cdots=X\left(t_{k}\right) ; t_{j} \in\left[r_{2 j-1}, r_{2 j}\right] ; K_{i}=\left\{r_{1}, \ldots, r_{2 k}\right\}\right] .
$$

Then $\bigcup_{i=1}^{\infty} C_{i}=\left[X\left(t_{1}\right)=\cdots=X\left(t_{k}\right)\right.$ for $k$ distinct times in [0,1]] has probability zero since $P\left[C_{i}\right]=0$ for each $i$. Apply the same reasoning to any finite time interval $[0, t]$ to complete the nonexistence proof.

Remark. Our theorems have not taken care of the criterical case

$$
k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)=\alpha_{1}+\alpha_{2} .
$$

It is possible to show non-existence in this case too, but the details are tedious. The methods used by Taylor [12] to settle the critical case for a stable process need to be modified by considering a new type of Hausdorff measure in which only coverings by long thin rectangles of the form used in the proof of Theorem 2 are allowed.

## Section 5: Dimension Theorems

In this section we establish bounds for the Hausdorff dimension of the set of multiple points. We use the terminology $L(a, b ; \omega), L^{k}(a, b ; \omega)$ and $E_{k}(\omega)$ of Section 2 and assume throughout this section that $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)<\alpha_{1}+\alpha_{2}$. In Theorem 3 we use methods similar to those of Theorem 2 to establish upper bounds for $\operatorname{dim} E_{k}$, the Hausdorff dimension of $E_{k}$, and in Theorem 4 we follow Taylor's [12] argument for the lower bound.

Theorem 3. Under the above terminology and hypotheses, with probability one $\operatorname{dim} E_{k}(\omega)$ satisfies:

$$
\operatorname{dim} E_{k}(\omega) \leqq \min \begin{cases}2-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{1} & (=\lambda) \\ 1+\alpha_{1} / \alpha_{2}-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{2} & (=\eta)\end{cases}
$$

Proof. We first show that $\operatorname{dim} E_{k} \leqq \eta$ by showing that $\operatorname{dim} E_{k} \leqq \theta$ whenever $\theta>\eta$. Let such a $\theta$ be chosen and cover $R^{2}$ with abutting rectangles $\left\{R_{i}\right\}_{i=1}^{\infty}$ congruent to $R(\delta)$ as in the proof of Theorem 2. $d_{i, \delta}$, the diagonal of any such rectangle $R_{i}$, is less than $3 \delta^{\alpha_{2}}$ for all $i$. Define events $A_{i}$ and $B_{i}$ as before, so that

$$
Q_{k}(\omega) \equiv \bigcap_{j=1}^{k} L(2 j-1,2 j ; \omega) \subset \bigcup_{i}^{*} R_{i}
$$

where the union is taken over those indices $i$ for which the event $A_{i} \cap B_{i}$ occurs. We now obtain an estimate of $E\left[\sum_{i=1}^{\infty} d_{i, \delta}^{\theta}(\omega)\right]$, the summation extending over the indices in the above union. Let $\mathscr{F}_{2}$ denote the $\sigma$-field of the process up to time $t=2$ and reason exactly as in Theorem 2:

$$
\begin{aligned}
E\left[\sum_{i=1}^{\infty} * d_{i, \delta}^{\theta}(\omega)\right] & =\sum_{i=1}^{\infty} E\left[d_{i, \delta}^{\theta} I_{A_{i}} \cdot I_{B_{i}}\right] \leqq c \delta^{\alpha_{2} \theta} \sum_{i=1}^{\infty} E\left[I_{A_{i}} E\left[I_{B_{i}} \mid \mathscr{F}_{2}\right]\right] \\
& \leqq c \delta^{\alpha_{2} \theta+(k-1)\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)} \cdot E[M(\delta, 2)] \leqq c \delta^{\alpha_{2}(\theta-\eta)}
\end{aligned}
$$

This final quantity approaches zero as $\delta \rightarrow 0$ whenever $\theta>\eta$. Consequently, with probability one there is a sequence of integers $\left\{j_{n}(\omega)\right\}_{n=1}^{\infty}$ such that $\sum_{i=1}^{\infty} d_{i, 1 / j_{n}}^{\hat{\infty}} \rightarrow 0$ as $n \rightarrow \infty$. Hence (with probability one) the $\theta$-Hausdorff measure of $Q_{k}(\omega)$ is zero. Similarly, the set $\bigcap_{j=1}^{k} L\left(r_{2 j-1}, r_{2 j} ; \omega\right)$ has zero $\theta$-Hausdorff measure for any sequence of $2 k$ distinct rational $r_{i}, 0<r_{1}<r_{2}<\cdots<r_{2 k} \leqq 1$. Any point $x$ in $L^{k}(0,1 ; \omega)$ is contained in a set of this form. Since the number of such sets is countable, $L^{k}(0,1 ; \omega)$ has zero $\theta$-Hausdorff measure with probability one. The same proof applies for any finite time interval $[0, t]$. Consequently $\operatorname{dim} E_{k} \leqq \eta$.

The proof that $\operatorname{dim} E_{k} \leqq \lambda$ is very similar to the above proof, so we only outline it. This time we cover $R^{2}$ with abutting squares $\left\{R_{i}\right\}_{i=1}^{\infty}$ of side $\delta \leqq 1$. Define $A_{i}$ and $B_{i}$ in analogous fashion. The corollary to Lemma 6.1 of Pruitt and Taylor [10]
tells us that $E[M(\delta, 2)] \leqq c \delta^{-\rho}$, where $\rho=1+\alpha_{2}-\alpha_{2} / \alpha_{1}$. We then obtain (for $\theta>\lambda$ ):

$$
E\left[\sum_{i=1}^{\infty} d_{i, \delta}^{\theta}(\omega)\right] \leqq c \delta^{\theta}\{C(S(\delta))\}^{k-1} E[M(\delta, 2)] \leqq c \delta^{\theta-\lambda}
$$

Reason as before to show that $\operatorname{dim} E_{k} \leqq \lambda$.
We conclude this section with an outline of the modifications of Taylor's [12] argument to establish

Theorem 4. Under the same hypotheses as in Theorem 3, with probability one $\operatorname{dim} E_{k}(\omega)$ satisfies:

$$
\operatorname{dim} E_{k}(\omega) \geqq \min \left\{\begin{array}{l}
2-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{1} \\
1+\alpha_{1} / \alpha_{2}-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{2}
\end{array}\right.
$$

Proof. For ease of reference we formulate and prove a sequence of lemmas and indicate the correspondence of our lemmas with those of [12]. With the exception of the upper bound estimate of Taylor's Lemma 1 we now have (through our Lemmas 3.3-3.5) the counterparts of his Lemmas $1-5$. We now handle this exception with

Lemma 5.1. Let $R_{1}$ and $R_{2}$ be two rectangles (with diagonal $d<\frac{1}{2}$ ) which have the same orientation and shape as $R_{a, b}$, and whose respective centers $e_{i}=\left(c_{i}, d_{i}\right)$, $i=1$ and 2 , are distinct. If $c_{1} \neq c_{2}\left(d_{1} \neq d_{2}\right)$ assume that $\left|c_{1}-c_{2}\right| \geqq 10 a\left(\left|d_{1}-d_{2}\right| \geqq 10 b\right)$. Let $y \in R_{2}$. Then:

$$
P^{y}\left[X(t) \in R_{1} \text { for some } t>0\right] \leqq c_{17} U\left(e_{1}-e_{2}\right) \cdot C\left(R_{a, b}\right)
$$

for some $c_{17}>0$ which is independent of $a, b, e_{1}$ and $e_{2}$.
Proof. Let $v$ be a capacitory measure on $R_{1}$. Then:

$$
\Phi\left(y, R_{1}\right)=\int_{R_{1}} U(x-y) v(d x) \leqq C\left(R_{1}\right) \max _{p_{i} \in R_{i}} U\left(p_{1}-p_{2}\right) .
$$

Now for any $p_{i} \in R_{i}, i=1$ and 2, we can use the bounds for the kernel (Lemma 3.1) to conclude that:

$$
U\left(p_{1}-p_{2}\right) \leqq \min \left\{c_{2} U\left(\left(0, \frac{4}{5}\left|d_{1}-d_{2}\right|\right)\right), c_{4} U\left(\left(\frac{4}{5}\left|c_{1}-c_{2}\right|, 0\right)\right)\right\} \leqq c_{17} U\left(e_{1}-e_{2}\right)
$$

Since the $p_{i}$ are selected arbitrarily in the $R_{i}$, the proof is complete. $\Phi\left(y, R_{2}\right), y \in R_{1}$, can be given the same bound as above once we observe (by Lemma 3.1) that $U\left(e_{2}-e_{1}\right) \leqq c U\left(e_{1}-e_{2}\right)$ for some finite and positive $c$ independent of the $e_{i}$. This symmetrizing property of the kernel will be used as necessary without further comment.

Our next lemma corresponds to Taylor's Lemma 8. Put $M=\max \left\{c_{9}, c_{13}\right\}>1$ (see Lemmas 3.4 and 3.5). This means that if $|x|+d<1$, the estimates of Lemma 3.5 (ii) for entry in $[0, M]$ and Lemma 3.4 for reentry in $[1, M]$ are valid in our next proof.

Lemma 5.2. Let $k$ be a positive integer, and $R_{1}$ and $R_{2}$ a pair of rectangles which satisfy the hypotheses of Lemma 5.1. Also let $\frac{1}{4} \leqq\left|e_{i}\right| \leqq \frac{1}{2}, i=1$ and 2 , and assume that $10 d<\min _{i=1,2}\left|e_{i}\right|$. Denote by $R_{i}^{*}(i=1$ and 2$)$ the event that there are $k$ time instants
$t_{1}, \ldots, t_{k}$ with $0<t_{1} \leqq M ; 1 \leqq t_{j}-t_{j-1} \leqq M, j=2,3, \ldots, k$ such that:

$$
X\left(t_{j}\right) \in R_{i}, \quad j=1,2, \ldots, k
$$

Then positive constants $c_{18}$ and $c_{19}$, independent of the $R_{i}$, exist such that:
(i) $P\left[R_{i}^{*}\right] \geqq c_{18}\left[C\left(R_{i}\right)\right]^{k}, \quad i=1,2$,
(ii) $P\left[R_{1}^{*} \cap R_{2}^{*}\right] \leqq c_{19}\left[C\left(R_{i}\right)\right]^{2 k}\left[U\left(e_{1}-e_{2}\right)\right]^{k}$.

We remark at this stage that our proof of this lemma is close to Taylor's, although we correct a misprint in a combinatorial argument in the proof of (ii). In addition, we shall need to examine carefully the estimate given by $\mu_{2}$ in several places later in our argument, so we provide the details of the proof.

Proof. A lower bound for the probability that $X(t)$ will hit $R_{i}$ in $[0, M]$ is given by (ii) of Lemma 3.5. If the process hits $R_{i}$ in [0,M] let $t_{1}$ be the first entry time; $t_{1}$ is a stopping time. Restart the process at $X\left(t_{1}\right)$ and use the strong Markov property. The conditional probability of reentry into $R_{i}$ in $[1, M]$ can be estimated by Lemma 3.4. Repeat the argument $(k-1)$ times to obtain (i).

To prove (ii) we define some probabilities and give estimates on them:

$$
\begin{aligned}
& \mu_{1}=\Phi\left(0, R_{1} \cup R_{2}\right) \\
& \mu_{2}=\max \left\{\sup _{y \in R_{1}} \Phi\left(y, R_{2}\right) ; \sup _{y \in R_{2}} \Phi\left(y, R_{1}\right)\right\} \\
& \mu_{3}=\sup _{\operatorname{all} y} P^{y}\left[X(t) \in R_{1} \cup R_{2} \text { for some } t \geqq \frac{1}{2}\right] .
\end{aligned}
$$

$\mu_{1}$ is estimated by means of the bounds on $U(y), \mu_{2}$ is estimated by Lemma 5.1, and $\mu_{3}$ by 3.3 (i):

$$
\begin{aligned}
& \mu_{1} \leqq \sum_{i=1}^{2} \int_{R_{i}} U(y-0) v(d y) \leqq \sum_{i=1}^{2} C\left(R_{i}\right) \max _{y \in R_{i}} U(y) \leqq c C\left(R_{i}\right), \quad i=1,2 . \\
& \mu_{2} \leqq c_{17} U\left(e_{1}-e_{2}\right) \cdot C\left(R_{i}\right), \quad i=1,2 . \\
& \mu_{3} \leqq 4 c_{7} C\left(R_{i}\right), \quad i=1,2
\end{aligned}
$$

Now let $\omega \in R_{1}^{*} \cap R_{2}^{*}$, and $N_{k}$ denote the number of ways of selecting $k$ integers out of $2 k$. There must be at least $2 k$ distinct times $t_{1}<\cdots<t_{2 k}$ for which $X\left(t_{j}, \omega\right) \in$ $R_{1} \cup R_{2}$. We may assume that the times corresponding to any two hits of a given rectangle differ by at least one time unit. Thus, for $j=2,3, \ldots, 2 k-1$ we must have $t_{j+1}-t_{j} \geqq \frac{1}{2}$ or $t_{j}-t_{j-1} \geqq \frac{1}{2}$. This implies that there must be at least $(k-1)$ (Taylor writes $k$ here instead of $(k-1)$, though he later corrects himself) integers $j$ for which $t_{j}-t_{j-1} \geqq \frac{1}{2}$.

The $t_{i}$ are stopping times; hence the product of probabilities for any one of the $N_{k}$ sequences of times must contain at least $(k-1)$ factors of $2 \mu_{3}$ (the 2 accounts for the fact that each factor may be inserted in two different ways); the first factor is always $\mu_{1}$; for each of the $k$ remaining $t_{i}-t_{i-1}$ there is a factor of at most $\mu_{2}+\mu_{3}$. Therefore:

$$
\begin{equation*}
P\left[R_{1}^{*} \cap R_{2}^{*}\right] \leqq N_{k} \mu_{1}\left(2 \mu_{3}\right)^{k-1}\left(\mu_{2}+\mu_{3}\right)^{k}, \tag{5.1}
\end{equation*}
$$

from which (ii) follows by use of the estimates upon the $\mu$ 's and the fact that $\mu_{2}+\mu_{3} \leqq c C\left(R_{i}\right) U\left(e_{1}-e_{2}\right)$ for $i=1,2$.

Before stating the next lemma, we introduce some terminology. We will be using rectangular grids of $r_{1} r_{2}$ points formed by the points of intersection of $r_{1}$ equally spaced horizontal lines and $r_{2}$ equally spaced vertical lines. The points $\left\{x_{i}\right\}$ of such a grid will be said to be numbered by columns starting from the lower left if starting with the left-hand column we number consecutively upward from the bottom and proceed column by column to the right. Thus, the lower left grid point is $x_{1}$ and the upper right is $x_{r_{1} r_{2}}$. For convenience of notation in the discussion which follows, we regard, for $\zeta>0$ and integral $r, r^{\zeta}$ as [ $r^{\zeta}$ ], where [] denotes the greatest integer function, whenever such an expression must be integral valued.

Next, let $r$ be a positive integer. Choose a positive number $L<1$ so that a rectangle $R$ with sides $L^{\alpha_{1}}$ parallel to the $y_{2}$ axis and of length $L^{\alpha_{2}}$ parallel to $y_{1}$ can be placed somewhere in the annulus $A=\left\{y: \frac{1}{4} \leqq|y| \leqq \frac{1}{2}\right\}$. Form the grid $G$ of $r^{\alpha_{1}} r^{\alpha_{2}}$ points by locating $r^{\alpha_{j}}$ equally spaced points along the sides of $R$ which are of length $L^{\alpha_{j}}, j=1$ and 2 , and connecting pairs of opposite points by horizontal or vertical lines. Number the points of $G=\left\{x_{v}: 1 \leqq v \leqq r^{\alpha_{1}} r^{\alpha_{2}}\right\}$ by columns starting from the lower left and let rectangles $R_{v}$ congruent to $R(\delta)$ have their centers at $x_{v}$, $1 \leqq v \leqq r^{\alpha_{1}} r^{\alpha_{2}}$. Assume that the rectangles satisfy the spacing requirements:

$$
10 \delta^{\alpha_{i}} \leqq(L / r)^{\alpha_{i}}, \quad i=1,2
$$

Thus, the sides of the rectangles are smaller by a factor of at least 5 than the horizontal and vertical spacing between their centers. Such a grid of rectangles will be said to be spaced according to $\alpha_{1}$ and $\alpha_{2}$. Observe that any pair of rectangles from such a grid satisfy the conditions of Lemma 5.2 for all large $r$.

Our technique for finding a lower bound for $\operatorname{dim} E_{k}$ will involve finding those $\gamma$ for which an independent symmetric stable process $X_{\gamma, 2}(t)$ of index $\gamma$ in $R^{2}$ hits the set $E_{k}$ and then using the following result of Taylor [12]:

Proposition 5.1. Let $A$ be an analytic set in $R^{2}$. Then for any point $x$,

$$
\operatorname{dim} A=2-\inf \left\{\gamma: \Phi_{\gamma, 2}(x, A)>0\right\}
$$

where $\Phi_{\gamma, 2}(x, A)$ denotes the probability that a symmetric stable process of index $\gamma$ in $R^{2}$ starting at $x$ will hit $A$.

For processes of the type given in this proposition we can use the results of [13] (pp. 1235 and 1237) and of our own Lemma 3.2 (when $\alpha_{1}=\alpha_{2}=\gamma>1$ ) to estimate the kernel, $U_{\gamma}()$, and capacity, $C_{\gamma}\left(R_{a, b}\right)$, for various $R_{a, b}$ with respect to $X_{\gamma, 2}$ :

$$
\begin{aligned}
c_{5} \delta^{\alpha_{2}(2-\gamma)} & \leqq C_{\gamma}(R(\delta)) \leqq c_{6} \delta^{\alpha_{2}(2-\gamma)} \quad \text { when } \gamma>1, \\
C_{\gamma}(S(\delta)) & =c_{20} \delta^{2-\gamma} \quad \text { and } \quad U_{\gamma}(x)=c_{21}|x|^{\gamma-2},
\end{aligned}
$$

for positive constants $c_{20}$ and $c_{21}$ independent of $x$ and $\delta$.
Our final two Lemmas, 5.3 and 5.4, are the counterparts of Taylor's Lemma 15 in the respective cases when $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)$ lies in the interval $\left[\alpha_{1}, \alpha_{1}+\alpha_{2}\right)$ and when it is in the interval $\left(0, \alpha_{1}\right)$. Once these lemmas are established the proof of our Theorem 4 can be completed according to the methods of Taylor's paper and we shall know the value of $\operatorname{dim} E_{k}$.

Lemma 5.3. Let $r$ be a large positive integer, assume that $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)$ lies in the interval $\left[\alpha_{1}, \alpha_{1}+\alpha_{2}\right)$ and suppose that for some $\varepsilon>0$ : ( $\varepsilon$ quite small)

$$
2-\gamma=1+\alpha_{1} / \alpha_{2}-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{2}-\varepsilon \quad(<1)
$$

Form a grid of $r^{\alpha_{1}} r^{\alpha_{2}}$ rectangles congruent to $R(\delta)$ spaced according to $\alpha_{1}$ and $\alpha_{2}$. Let $\delta=c_{22} r^{-\mu}$, where

$$
\mu=\left(\alpha_{1}+\alpha_{2}\right) /\left[k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)+\alpha_{2}(2-\gamma)\right] \quad(>1)
$$

and $c_{22}>0$ will be chosen later.
Let $E_{v}, 1 \leqq \nu \leqq r^{\alpha_{1}} r^{\alpha_{2}}$ be the event that there are time instants $0<t_{1} \leqq N$, $1 \leqq t_{j}-t_{j-1} \leqq N(2 \leqq j \leqq k)$, and $0<t_{k+1} \leqq N$ such that:

$$
X_{\gamma, 2}\left(t_{k+1}, \omega\right) \in R_{v} \quad \text { and } \quad X\left(t_{j}, \omega\right) \in R_{v} \quad(1 \leqq j \leqq k)
$$

Then $c_{22}$ and $c_{23}$, positive and independent of $r$, can be chosen such that

$$
P\left[\bigcup_{v=1}^{r^{\alpha_{1}} r^{\alpha_{2}}} E_{v}\right] \geqq c_{23} \quad \text { for all large } r
$$

(In the above, $X_{\gamma, 2}$ and $X$ are assumed independent and defined on the same basic probability space, and $N$ is chosen large enough to ensure that the required estimates are valid for each process.)

Proof. We first make four observations:
(i) $\mu>1$, so that the $R(\delta)$ satisfy the spacing requirements:

$$
10 \delta^{\alpha_{i}}=10\left(c_{22} r^{-\mu}\right)^{\alpha_{i}} \leqq(L / r)^{\alpha_{i}}, \quad i=1 \text { and } 2
$$

(ii) The estimates for the capacity $C(R(\delta))$ of the $R_{v}$ for the $\left(X_{1}, X_{2}\right)$ process and $C_{\gamma}(R(\delta))$ for the $X_{\gamma, 2}$ process give:

$$
\{C(R(\delta))\}^{k} C_{\gamma}(R(\delta)) \sim r^{-\left(\alpha_{1}+\alpha_{2}\right)}
$$

(iii) If $\theta<1$ and $r$ is a large positive integer, a positive constant $c_{24}$ (independent of $r$ ) exists such that:

$$
\sum_{j=1}^{r}(1 / j)^{\theta} \leqq c_{24} r^{1-\theta}
$$

(iv) $P\left[\bigcup_{v=1}^{r^{\alpha_{1}} r^{\alpha_{2}}} E_{v}\right] \geqq r^{\alpha_{1}+\alpha_{2}} P\left[E_{1}\right]-\sum^{*} P\left[E_{v} \cap E_{v^{\prime}}\right]-\sum^{* *} P\left[E_{v} \cap E_{v^{\prime}}\right]$, where $\sum^{* *}$ is taken over distinct $v$ and $v^{\prime}$ for which the corresponding grid points $x_{v}$ and $x_{v^{\prime}}$ lie on the same vertical line, and $\Sigma^{*}$ is taken over all remaining pairs of distinct $v$ and $v^{\prime}$. Denote these sums by $S^{* *}$ and $S^{*}$ respectively.

Since the rectangles $R_{v}$ satisfy the conditions of Lemma 5.2 we can use both parts of that lemma and (ii) above to write:

$$
\begin{equation*}
r^{\alpha_{1}+\alpha_{2}} P\left[E_{1}\right]-S^{*} \geqq c_{25}-c_{26} r^{-2\left(\alpha_{1}+\alpha_{2}\right)} \sum^{*}\left\{U\left(x_{v}-x_{v^{\prime}}\right)\right\}^{k} U_{\gamma}\left(x_{v}-x_{v^{\prime}}\right) \tag{5.2}
\end{equation*}
$$

where the summation is over the pairs of distinct $x_{v}$ and $x_{v^{\prime}}$ which are not on the same vertical line. Moreover,

$$
\begin{equation*}
S^{* *} \leqq c_{27}[C(R(\delta))]^{2 k} C_{\gamma}(R(\delta)) \sum^{* *}\left\{U\left(x_{v}-x_{v^{\prime}}\right)\right\}^{k} \tag{5.3}
\end{equation*}
$$

where the summation is over the $x_{v}$ and $x_{v^{\prime}}$ which lie on the same vertical line. Thus, in (5.3) we are bounding $\mu_{2}$ (of Lemma 5.2) above by 1 so that the estimate of (5.1) for $X_{\gamma, 2}$ to hit $R_{1}$ and $R_{2}$ each once is simply $c C_{\gamma}(R(\delta))$. Failure to do this can result in a $\mu_{2}$ estimate (for the $X_{\gamma, 2}$ process) which exceeds one. Our objective now is twofold:
(a) Toshow that the summation (with we henceforth denote by $\sum^{*}$ ) in (5.2) is bounded above by $c_{28} r^{2\left(x_{1}+\alpha_{2}\right)}$ for some $c_{28}>0$ which does not depend upon $r$ and to then select $c_{22}$ (in the formation of $\delta=c_{22} r^{-\mu}$ ) so that the desired constant $c_{23}$ can be shown to exist.
(b) To show that the bound upon $S^{* *}$ goes to zero as $r \rightarrow+\infty$.

Objective (b) is more easily accomplished than (a), so we do (b) first when $k() \equiv k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) \in\left(\alpha_{1}, \alpha_{1}+\alpha_{2}\right)$ and then state how the case $k()=\alpha_{1}$ can be handled. Under the first hypothesis we have:

$$
\begin{aligned}
S^{* *} & \leqq c_{27}[C(R(\delta))]^{2 k} C_{\gamma}(R(\delta)) r^{\alpha_{2}} \sum_{j=1}^{r^{\alpha_{1}}-1} \sum_{i=1}^{r^{\alpha_{1}}-j}\left\{U\left(\left(0, i(L / r)^{\alpha_{1}}\right)\right)\right\} \\
& \leqq c r^{-\left(\alpha_{1}+\alpha_{2}\right)} r^{(1-\mu) k()} r^{\alpha_{2}} \sum_{j=1}^{r_{1}-1} \sum_{i=1}^{r^{\alpha_{1}}-j} i^{-k() / \alpha_{1}} \\
& \leqq c r^{(1-\mu) k( } \rightarrow 0 \quad \text { as } r \rightarrow \infty \text { since } \mu>1 .
\end{aligned}
$$

When $k()=\alpha_{1}$ the summation on $i$ diverges (as $\left.r \rightarrow \infty\right)$ to give a $\ell n\left(r^{\alpha_{1}}-j\right)$ term and Stirling's formula can be used on the resulting $\ln \left(r^{\alpha_{1}}-1\right)$ ! factor to complete the proof of (b).

We now consider the $\sum^{*}$ factor of (5.2):

$$
\sum^{*} \leqq 2 r^{\alpha_{1}+\alpha_{2}} \sum_{v>r^{\alpha_{1}}}\left\{U\left(x_{v}-x_{1}\right)\right\}^{k} U_{\gamma}\left(x_{v}-x_{1}\right)
$$

Hence we can regard the grid $G$ as being repositioned with its sides along the coordinate axes, $x_{1}$ at $(0,0)$ and $x_{r^{\alpha_{1}+\alpha_{2}}}$ at ( $L^{\alpha_{2}}, L^{\alpha_{1}}$ ), and that we are required to sum $\left\{U\left(z_{v}\right)\right\}^{k} U_{\gamma}\left(z_{v}\right)$ over those $v$ such that $Z_{v} \equiv x_{v}-x_{1}$ does not lie on the vertical $\left(y_{2}\right)$ coordinate axis. In the first quadrant of the ( $y_{1}, y_{2}$ ) plane the line $y_{2}=L^{\alpha_{1}-\alpha_{2}} y_{1}$ and the curve $y_{2}^{\alpha_{2}}=y_{1}^{\alpha_{1}}$ intersect at $(0,0)$ and $\left(L^{\alpha_{2}}, L^{\alpha_{1}}\right)$ and divide the rectangle which forms the repositioned grid into three regions which we number (from bottom to top) by I, III and II respectively. Our sums over these regions will be denoted by $S_{\mathrm{I}}, S_{\mathrm{III}}, S_{\mathrm{II}}$ and we will be finished once we show each sum bounded above by $c r^{\alpha_{1}+\alpha_{2}}$. We now do this.

Region I contains the $y_{1}$ coordinate axis, and along the vertical line through $\left(j(L / r)^{\alpha_{2}}, 0\right)$ there are at most $\left(j^{\alpha_{1} / \alpha_{2}}+1\right)$ grid points which lie in $\mathrm{I}, j=1,2, \ldots, r^{\alpha_{2}}$. If we estimate $U_{\gamma}\left(y_{1}, y_{2}\right)$ by $U_{\gamma}\left(y_{1}, 0\right)$ in I we then have (we suppress the argument of $U_{\gamma}$ when it is the same as that of $U$ ):

$$
\begin{aligned}
S_{\mathrm{I}} & \leqq 2 \sum_{j=1}^{r^{\alpha_{2}}} j^{\alpha_{1} / \alpha_{2}}\left\{U\left(j(L / r)^{\alpha_{2}}, 0\right)\right\}^{k} U_{\gamma}(,) \\
& \leqq c r^{k()+\alpha_{2}(2-\gamma)} \sum_{j=1}^{r^{\alpha}}(1 / j)^{k() / \alpha_{2}+(2-\gamma)-\alpha_{1} / \alpha_{2}} \leqq c r^{\alpha_{1}+\alpha_{2}}
\end{aligned}
$$

by using the defining relationship for $(2-\gamma)$ to justify use of observation (iii).

Region II contains the $y_{2}$ axis, and along the horizontal line through $\left(0, j(L / r)^{\alpha_{1}}\right)$ there are at most $j r^{\alpha_{2}-\alpha_{1}}$ grid points which lie in II, $r^{\alpha_{1}-\alpha_{2}} \leqq j \leqq r^{\alpha_{1}}$. To see this, determine the points of intersection of the grid lines with the line $y_{2}=L^{\alpha_{1}-\alpha_{2}} y_{1}$. If we estimate $U_{\gamma}\left(y_{1}, y_{2}\right)$ by $U_{\gamma}\left(0, y_{2}\right)$ in II we then have:

$$
\begin{aligned}
& S_{\mathrm{II}} \leqq \sum_{j=r^{\alpha_{1}-\alpha_{2}}}^{r^{\alpha_{1}}} j r^{\alpha_{2}-\alpha_{1}}\left\{U\left(0, j(L / r)^{\alpha_{1}}\right)\right\}^{k} U_{\gamma}(,) \\
& \quad \leqq c r^{\alpha_{2}-\alpha_{1}+k()+\alpha_{1}(2-\gamma) \sum_{j=1}^{r \alpha_{1}}(1 / j)^{k()) / \alpha_{1}+(2-\gamma)-1} \leqq c r^{\alpha_{1}+\alpha_{2}}}
\end{aligned}
$$

by reasoning as before and using the fact that $k() \geqq \alpha_{1}$.
Region III is slightly more complicated. Along horizontal lines in III we estimate $U\left(y_{1}, y_{2}\right)$ by $U\left(0, y_{2}\right)$ and $U_{y}\left(y_{1}, y_{2}\right)$ by $U_{\gamma}\left(y_{1}, 0\right)$. Hence

$$
\begin{aligned}
& S_{\text {III }} \leqq \sum_{j=1}^{r^{\alpha_{1}}} \sum_{i=1}^{j^{\alpha_{2} / \alpha_{1}}}\left\{\underset { r ^ { \alpha _ { 1 } } } { } \left\{\left(i(L / r)^{\alpha_{2}}, \underset{j^{\alpha} 2 / \alpha_{1}}{\left.j(L / r)^{\alpha_{1}}\right)}\right\}^{k} U_{\gamma}(,)\right.\right. \\
& \leqq r^{\substack{k()+\alpha_{2}(2-\gamma)}} \sum_{j=1}^{r \alpha_{1}}(1 / j)^{k() / \alpha_{1}} \sum_{i=1}^{j^{\alpha_{2} / \alpha_{1}}}(1 / i)^{2-\gamma} \leqq c r^{\alpha_{1}+\alpha_{2}}
\end{aligned}
$$

by noting that $(2-\gamma)<1$ and applying observation (iii) to both summations.
Finally, note that a positive power of the constant $c_{22}$ used to define $\delta$ enters as a factor of the $c_{25}$ in (5.2) and as a squared factor in $c_{26}$. Hence the derived $c_{23}$ can be found and Lemma 5.3 is proven. The conclusion of Theorem 4 when $k() \in$ [ $\alpha_{1}, \alpha_{1}+\alpha_{2}$ ) can now be established by using Taylor's argument (Section 7 of [12]) and noting the relationship between the two functions $f_{1}(x)=2-x / \alpha_{1}$ and $f_{2}(x)=1+\alpha_{1} / \alpha_{2}-x / \alpha_{2} 0<x<\alpha_{1}+\alpha_{2}$ on the interval $\left[\alpha_{1}, \alpha_{1}+\alpha_{2}\right.$ ). We now complete the proof of Theorem 4 by handling the case $k() \in\left(0, \alpha_{1}\right)$ in our final lemma.

Lemma 5.4. Let $r$ be a large positive integer, assume that $k$ () lies in the interval $\left(0, \alpha_{1}\right)$ and suppose that for some $\varepsilon>0$ : ( $\varepsilon$ quite small)

$$
2-\gamma=2-k() / \alpha_{1}-\varepsilon \quad(>1)
$$

Form a grid of $r^{2}$ squares of side $\delta$ centered at the points of intersection of $r$ equally spaced vertical lines along the $L^{\alpha_{2}}$ side and $r$ equally spaced horizontal lines along the $L^{\alpha_{1}}$ side of a rectangle of the type used in Lemma 5.3. Number the centers of the squares by columns starting from the lower left and let $\delta=c_{29} r^{-\mu}$, where

$$
\mu=2 \alpha_{1} /\left[k()+\alpha_{1}(2-\gamma)\right] \quad(>1)
$$

and $c_{29}>0$ will be chosen later.
Let $E_{v}, 1 \leqq v \leqq r^{2}$ be the event that there are time instants $0<t_{1} \leqq N$,

$$
1 \leqq t_{j}-t_{j-1} \leqq N \quad(2 \leqq j \leqq k)
$$

and $0<t_{k+1} \leqq N$ such that

$$
X_{\gamma}\left(t_{k+1}, \omega\right) \in R_{v} \quad \text { and } \quad X\left(t_{j}, \omega\right) \in R_{v} \quad(1 \leqq j \leqq k)
$$

Then $c_{29}$ and $c_{30}$, positive and independent of $r$, can be chosen such that

$$
P\left[\bigcup_{v=1}^{r^{2}} E_{v}\right] \geqq c_{30}
$$

for all large $r$.

Proof. We again make some observations:
( $i^{\circ}$ ) $\mu>1$, so that the squares satisfy the requirements

$$
10 \delta=10\left(c_{29} r^{-\mu}\right) \leqq L^{\alpha_{i} / r}, \quad i=1 \text { and } 2
$$

(iio) $[C(S(\delta))]^{k} C_{y}(S(\delta)) \sim \delta^{k() / \alpha_{1}+(2-\gamma)}$.
(iii ${ }^{\circ}$ ) If $\theta>1$ and $r$ is a large positive integer, a positive constant $c_{31}$ (independent of $r$ ) exists such that:

$$
\sum_{j=r}^{\infty}(1 / j)^{\theta}<c_{31} r^{1-\theta}
$$

(iv $\left.{ }^{\mathrm{o}}\right) P\left[\bigcup_{v=1}^{r^{2}} E_{v}\right] \geqq r^{2} P\left[E_{1}\right]-\sum^{\mathrm{o}} P\left[E_{v} \cap E_{v^{\prime}}\right]-\sum^{\mathrm{oo}} P\left[E_{v} \cap E_{v^{\prime}}\right]$, where $\sum^{\mathrm{oo}}$ is taken over distinct $v$ and $v^{\prime}$ for which the corresponding grid points $x_{v}$ and $x_{v^{\prime}}$ lie on the same horizontal line, and $\sum^{\circ}$ is taken over all remaining pairs of distinct $v$ and $v^{\prime}$. Denote these sums by $S^{o 0}$ and $S^{0}$ respectively.

Now argue as in (5.2) to establish

$$
r^{2} P\left[E_{1}\right] \geqq c_{32}-c_{33} r^{-2} \sum^{0}\left\{U\left(x_{v}-x_{v^{\prime}}\right)\right\}^{k} U_{\gamma}\left(x_{v}-x_{v^{\prime}}\right)
$$

where the summation is over distinct $x_{v}$ and $x_{v^{\prime}}$ which are not on the same horizontal line. Moreover,

$$
S^{00} \leqq c_{34}[C(S(\delta))]^{k}\left[C_{\gamma}(S(\delta))\right]^{2} \sum^{\circ o} U_{\gamma}\left(x_{v}-x_{v^{\prime}}\right)
$$

where the summation is over distinct $x_{v}$ and $x_{v^{\prime}}$ on the same horizontal line. This time we are bounding $\mu_{2}$ above by 1 so that the estimate of $(5.1)$ for $X(t)$ to hit $R_{1}$ and $R_{2} k$ times apiece is simply $c[C(S(\delta))]^{k}$. Again we have a twofold objective:
$\left(\mathrm{a}^{0}\right)$ Show the summation (denoted by $\left.\sum^{\circ}\right)$ in $\left(5.2^{\circ}\right)$ is bounded above by $c_{35} r^{2}$ and select $c_{29}$ appropriately.
( $\mathrm{b}^{\circ}$ ) Show that the bound upon $S^{00}$ goes to zero as $r \rightarrow+\infty$.
First consider ( $\mathrm{b}^{\circ}$ ). By (iii ${ }^{\circ}$ ) above and $\left(5.3^{\circ}\right)$ we have:

$$
\begin{aligned}
S^{\mathrm{oo}} & \leqq c(1 / r)^{2+\mu(2-\gamma)} r \sum_{j=1}^{r-1} \sum_{i=1}^{r-j} U_{\gamma}\left(\left(i L^{\alpha_{2}} / r, 0\right)\right) \\
& \leqq c r^{(2-\gamma)(1-\mu)} \rightarrow 0 \quad \text { as } r \rightarrow+\infty(\text { since } \mu>1 \text { and } 2-\gamma>1) .
\end{aligned}
$$

We now consider the $\sum^{0}$ factor of $\left(5.2^{\circ}\right)$ :

$$
\sum^{0} \leqq 2 r^{2} \sum\left\{U\left(x_{v}-x_{1}\right)\right\}^{k} U_{\gamma}\left(x_{v}-x_{1}\right) \quad(v \neq j r+1,0 \leqq j \leqq r-1) .
$$

Reposition the grid as before, form the three regions I, III and II, and consider the sum of $\left\{U\left(z_{v}\right)\right\}^{k} U_{\gamma}\left(z_{v}\right)$ for $z_{v}$ lying in these regions (denoted by $S_{1}^{\mathrm{o}}, S_{\text {III }}^{\mathrm{o}}$ and $S_{\text {II }}^{\mathrm{o}}$ respectively). Once we show that each of these sums is bounded above by $c r^{2}$, the lemma (and hence the theorem) will be established. In regions I and II we use observation (ii) of our previous lemma, while III requires (iio) of the present lemma. To use (ii) in I we also require the condition $k()<\alpha_{1}$. The desired bounds follow from:

$$
\begin{aligned}
& S_{\mathrm{I}}^{\mathrm{o}} \leqq c r^{1-\alpha_{1} / \alpha_{2}} \sum_{j=1}^{r} j^{\alpha_{1} / \alpha_{2}}\left\{U \left(\left(j L^{\left.\left.\left.\alpha_{2} / r, 0\right)\right)\right\}^{k} U_{\gamma}()}\right.\right.\right. \\
& S_{\mathrm{II}}^{\mathrm{o}} \leqq c \sum_{j=1}^{\mathrm{r}} j\left\{U\left(\left(0, j L^{\alpha_{1}} / r\right)\right)\right\}^{k} U_{\gamma}()
\end{aligned}
$$

and

$$
\begin{aligned}
S_{\mathrm{III}}^{\mathrm{o}} & \leqq \sum_{j=1}^{r} \sum_{i=j}^{r}\left\{U\left(\left(i L^{\alpha_{2}} / r, j L^{\left.\alpha_{1} / r\right)}\right)\right\}^{k} U_{\gamma}(,)\right. \\
& \leqq r^{k() / \alpha_{1}+(2-\gamma)} \sum_{j=1}^{r}(1 / j)^{k(1) / \alpha_{1}} \sum_{i=j}^{r}(1 / i)^{2-\gamma} \\
& \leqq r^{k() / \alpha_{1}+(2-\gamma) \sum_{j=1}^{r}(1 / j)^{k() / \alpha_{1}+(2-\gamma)-1}}
\end{aligned}
$$

## Section 6: Remarks

(i) In [4] we have recently shown that for fixed Borel sets $E \subset[0,1]$ and for $X(t)$ of the type being considered here we have (with probability one):

$$
\operatorname{dim} X(E)=\left\{\begin{array}{l}
\alpha_{1} \operatorname{dim} E \quad \text { for } 0 \leqq \operatorname{dim} E \leqq 1 / \alpha_{1} \\
1+\alpha_{2} \operatorname{dim} E-\alpha_{2} / \alpha_{1} \quad \text { for } 1 / \alpha_{1} \leqq \operatorname{dim} E \leqq 1
\end{array}\right.
$$

where $X(E)$ denotes the set of points in $R^{2}$ hit at some time $t \in E$. Intuitively, we can think of using up the $X_{1}$ component on sets $E$ of small dimension $\left(\operatorname{dim} X_{1}(E)=\right.$ $\alpha_{1} \operatorname{dim} E$ when $\operatorname{dim} E \leqq 1 / \alpha_{1}$ according to the results of Blumenthal and Getoor [2] in the stable case) and then using up the $X_{2}$ component.

Theorems 3 and 4 show us hat:

$$
\operatorname{dim} E_{k}= \begin{cases}2-k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right) / \alpha_{1} & \text { when } 0<k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)<\alpha_{1} \\ 1+\alpha_{1} / \alpha_{2}-k() / \alpha_{2} & \text { when } \alpha_{1} \leqq k()<\alpha_{1}+\alpha_{2} .\end{cases}
$$

For stable processes $X_{\alpha, 2}$ of index $\alpha>1$ in the plane Taylor's result is: $\operatorname{dim} E_{k}=$ $2-k(2-\alpha)$. Thus, we are again led to think (intuitively) of successively using up the two components. In fact, this is what led to the discovery of our $\operatorname{dim} E_{k}$ results.
(ii) The proofs of our dimension theorems required considerable effort to uncover, and it did not seem possible to avoid using the two kinds of grids and rectangles. This is no doubt a reflection of some of the basic properties of the process.
(iii) If we consider the quantity $k\left(\alpha_{1}+\alpha_{2}-\alpha_{1} \alpha_{2}\right)$ we can determine a maximum multiplicity possible. For the processes being considered double points must exist, but it is possible (by choosing $\alpha_{1}$ and $\alpha_{2}$ close to 1) that triple points do not exist. Moreover, given an integer $k>2$, it is possible (by choosing $\alpha_{1}$ and $\alpha_{2}$ near 2 for large $k$ ) to have points of multiplicity $k$ but that $(k+1)$-multiple points will fail to exist.
(iv) The problem of multiple points in $R^{3}$ when one of the components is planar or if 3 independent one-dimensional components are present remains unsolved. Earlier work on these processes ( $[4,6,10]$ ) suggests that the present case is the most interesting one, but Fristedt's [3] extension of Taylor's [12] result suggests that the problem may involve some detailed calculations. We conjecture that there are basically two cases to consider, depending upon whether $\alpha_{1}>d_{1}\left(\right.$ dimension of $X_{1}$ ) or $\alpha_{1} \leqq d_{1}$. Taylor's work in the stable case ([12] and [13]) requires each of the stable components to be type $A$, and to have stable index greater than $\frac{1}{2}$ when the component is linear. When $\alpha_{1}>d_{1}$ we suspect the results
to be identical to those obtained in the present paper. When $\alpha_{1} \leqq d_{1}$ we predict that the behaviour of the process is similar to that of $X_{1}$. Analogous results in the cases $\alpha_{1}>d_{1}$ and $\alpha_{1} \leqq d_{1}$ were obtained in the previous studies.
(v) The above conjecture has been found to be true for $\left(X_{1}, X_{2}\right)$ in $R^{2}$ when $\frac{1}{2}<\alpha_{2} \leqq 1<\alpha_{1}$ and $X_{2}$ is type A (and in addition symmetric if $\alpha_{2}=1$ ). One must only check that the estimates of Lemma 3.1 for the kernel and of Lemma 3.2 for the capacity of the rectangles actually used in the argument hold. In Lemma 3.1 we can estimate the kernel in the same way as before. The lower bound poses no difficulty, while the upper bound as done in Lemma 3.1 of [7] does seem to require $\alpha_{2} \geqq \frac{1}{2}$ (which we assume anyway). In Lemma 3.2 we thus have the same upper bounds upon $C\left(R_{a, b}\right)$ as before. The lower bounds for $C(R(\delta))$ and $C(S(\delta))$ which are obtained by only assuming $\alpha_{1}>1$, namely $c^{-1} b^{1-\alpha_{2}+\alpha_{2} / \alpha_{1}}$, are the same as the upper bounds for the capacity of these two types of rectangles. Hence it suffices in all of our theorems to allow $\frac{1}{2}<\alpha_{2} \leqq 1<\alpha_{1}$ when the process is planar.

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## References

1. Blumenthal, R.M., Getoor, R.K.: Markov Processes and Potential Theory. New York: Academic Press 1968
2. Blumenthal, R.M., Getoor, R.K.: The dimension of the set of zeros and the graph of a symmetric stable process. Illinois J. Math. 6, 308-316 (1962)
3. Fristedt, B.: An extension of a theorem of S. J. Taylor concerning the multiple points of the symmetric stable process. Z. Wahrscheinlichkeitstheorie verw. Gebiete 9, 62-64 (1967)
4. Hendricks, W.J.: A dimension theorem for processes with stable components. Ann. Probab. (To appear)
5. Hendricks, W.J.: Hausdorff dimension in a process with stable components-an interesting counterexample. Ann. Math. Statis. 43, 690-694 (1972)
6. Hendricks, W.J.: Lower envelopes near zero and infinity for processes with stable components. Z. Wahrscheinlichkeitstheorie verw. Gebiete 16, 261-278 (1970)
7. Jain, N., Pruitt, W.E.: Collisions of stable processes. Illinois J. Math. 13, 241-248 (1969)
8. Lévy, P.: Theorie de l'addition des variables aleatoires. Paris 1937
9. Loève, M.: Probability Theory. Princeton: Van Nostrand 1963
10. Pruitt, W.E., Taylor, S.J.: Sample path properties of processes with stable components. Z. Wahrscheinlichkeitstheorie verw. Gebiete 12, 267-289 (1969)
11. Skorohod, A.V.: Asymptotic formulas for stable distribution laws. Selected Translations in: Math. Statist. Probab. 1, 157-162 (1961)
12. Taylor, S. J.: Multiple points for the sample paths of the symmetric stable processes. Z. Wahrscheinlichkeitstheorie verw. Gebiete 5, 247-264 (1966)
13. Taylor, S.J.: Sample path properties of a transient stable process. J. Math. Mech. 16, 1229-1246 (1967)

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