# Random Walk in Random Environment: A Counterexample without Potential 

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#### Abstract

We describe a family of random walks in random environment which have exponentially decaying correlations, nearest neighbor transition probabilities which are bounded away from 0 , and are subdiffusive in any dimension $d<\infty$. The random environments have no potential in $d>1$.


KEY WORDS: Random walk; random environment; subdiffusive; exponentially decaying correlations.

## 1. INTRODUCTION

Random walks in random environment have been the subject of considerable attention in recent years. Yet, few rigorous results are known about the behavior in dimensions $d>1$. It has been shown in a momentous forthcoming article ${ }^{(1)}$ that under independent environments and appropriate symmetry conditions, the mean square displacement will be asymptotically linear in time with the scaled distribution approaching that of a normal. It is believed that for models with short-range correlations, the mean square displacement also grows linearly. ${ }^{(26)}$ In ref. 7, a family of models having spatially homogeneous random environments with exponentially decaying correlations and nearest neighbor transition probabilities which are bounded away from 0 was introduced. The random walks on these environments were shown to be subdiffusive in any dimension $d<\infty$. The environments in this family all possess potentials. The models were therefore met with some reservations as valid counterexamples.

The purpose of this article is to construct a family of models with the same features as above, but where the associated random environments do

[^0]not possess potentials (for $d>1$ ). These models are obtained by perturbing the environments in ref. 7 by independent environments so that the random walks retain their subdiffusive behavior. A major part of the construction and proof for these models resembles that in ref. 7; the reader is referred there for additional background.

The models considered in ref. 7 are a special case of random walk on a random hillside. In these systems, one starts with a random function $V: \mathbf{R}^{d} \rightarrow \mathbf{R}$ (the hillside or potential), defines

$$
\begin{equation*}
\alpha(x, y)=\exp [-\beta V((x+y) / 2)] \tag{1}
\end{equation*}
$$

for $x, y \in \mathbf{Z}^{d}$ with $|x-y|=1$, and for convenience sets $\alpha(x, y)=0$ otherwise. The $\alpha(x, y)$ are nonnegative, so if we let

$$
\alpha(x)=\sum_{y} \alpha(x, y)
$$

and

$$
p(x, y)=\alpha(x, y) / \alpha(x),
$$

then

$$
p(x, y) \geqslant 0 \quad \text { and } \quad \sum_{y} p(x, y)=1
$$

i.e., $p$ is a transition probability. From $p$ one constructs a random walk in random environment in the usual way: if $X(n)=x$ (that is, the particle is at $x$ at time $n$ ), then the probability it will jump to $y$ at time $n+1$ is $p(x, y)$ and is independent of what happened before time $n$. The reader should note that the definition of $p$ is unchanged if we replace $\alpha$ by

$$
\begin{equation*}
\bar{\alpha}(x, y)=\exp \{-\beta[V((x+y) / 2)-V(x)]\}, \tag{2}
\end{equation*}
$$

since the extra factor will cancel when one normalizes. The value of $p(x, y)$ therefore depends only on the increments $V((x+\cdot))-V(x)$. Assume that $X(0)=0$.

To construct the potential $V$ used in ref. 7, let $k(z), z \in \mathbb{Z}^{d}$, be independent random variables with

$$
\begin{align*}
& P[k(z)=0]=1-\delta, \\
& P[k(z)=k]=\delta \varepsilon(1-\varepsilon)^{k-1}, \quad k=1,2, \ldots . \tag{3}
\end{align*}
$$

We abbreviate these probabilities by $p_{k}$. Here $0<\delta<1$ and $0<\varepsilon \leqslant 1 / 2$. One may think of $V$ as being the surface of a (random) moon, with $k(z)$
giving the radius of the crater centered at $z$. If one lets $|x|=\left|x_{1}\right|+\cdots+\left|x_{d}\right|$, then the function

$$
\begin{equation*}
\varphi_{k}(x)=\min \{|x|-k, 0\} \tag{4}
\end{equation*}
$$

gives the depth of the (square) crater of radius $k$ centered at 0 . Define the surface of our moon by

$$
\begin{equation*}
V(x)=\min _{z} \varphi_{k(z)}(x-z), \tag{5}
\end{equation*}
$$

where the minimum is taken over $z$ in $\mathbf{Z}^{d}$.
We note that $V$ as defined here has slope $\leqslant 1$, and so, on account of (2),

$$
\begin{equation*}
p(x, y) \geqslant e^{-\beta / 2} / 2 d e^{\beta / 2}=\left(2 d e^{\beta}\right)^{-1} \tag{6}
\end{equation*}
$$

for $|x-y|=1$. From the above definition it is clear that the increments in $V$ have exponentially decaying correlations. The following result from ref. 7 shows that random walks in these environments are subdiffusive.

Theorem A. Suppose that $0<\delta<1,0<\varepsilon \leqslant 1 / 2$, and $N>0$. If $\beta \geqslant 2(N+d+1)$, then

$$
\begin{equation*}
P\left[\max _{1 \leqslant n}|X(j)| \geqslant n^{1 / N}\right] \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty \tag{7}
\end{equation*}
$$

The random walk $X(n)$ has been constructed from the random potential $V(x)$. The presence of this potential can be thought of as placing a longrange constraint on $p(x, y)$. One can therefore consider this random walk as having a "random potential" rather than a "random force." One can, however, modify this example to a random walk $Y(n)$ on $Z^{d}$ with probabilities $p^{\prime}(x, y)$ constructed in terms of $\alpha^{\prime}(x, y)$ in place of $\alpha(x, y)$ as above. Set

$$
\begin{equation*}
\alpha^{\prime}(x, y)=\exp \{-\beta[V((x, y) / 2)-V(x)+W(x, y)]\} \tag{8}
\end{equation*}
$$

for $|x-y|=1$, where $W(x, y)$ are random variables (which are not necessarily independent). Then, as above, set

$$
\begin{equation*}
\alpha^{\prime}(x)=\sum_{y} \alpha^{\prime}(x, y) \tag{9a}
\end{equation*}
$$

and

$$
\begin{equation*}
p^{\prime}(x, y)=\alpha^{\prime}(x, y) / \alpha^{\prime}(x) \tag{9b}
\end{equation*}
$$

Of course,

$$
p^{\prime}(x, y) \geqslant 0 \quad \text { and } \quad \sum_{y} p^{\prime}(x, y)=1
$$

Note that if $0 \leqslant W(x, y) \leqslant M$ for all $x, y$, then

$$
\begin{equation*}
p^{\prime}(x, y) \geqslant e^{\beta M} p(x, y) \geqslant\left(2 d e e^{\beta(M+1)}\right)^{1} . \tag{10}
\end{equation*}
$$

We denote by $Y(n)$ the random walk in random environment corresponding to $p^{\prime}$. Except when specified otherwise, $Y(0)=0$ is assumed.

We prove the following analog of Theorem A .
Theorem 1. Suppose that $0<\delta<1,0<\varepsilon \leqslant 1 / 2$, and $N>0$. If $\beta \geqslant C d(N+1) / \varepsilon$ for appropriate $C$, and $0 \leqslant W(x, y) \leqslant 1 / 4$, then

$$
\begin{equation*}
P\left[\max _{j \leqslant n}|Y(j)| \geqslant n^{1 / N}\right] \rightarrow 0 \quad \text { as } \quad n \rightarrow m \tag{11}
\end{equation*}
$$

The process $Y(j)$ has the properties we desire. As before, ${ }^{\prime} p^{\prime}(x, y)$ : $|x-y|=1\}$ is bounded away from zero. If $W$ is independent of $V$ and has exponentially decaying correlations, so does $V^{\prime}=(V, W)$. Of course it is easy to choose $W$ so that $\alpha^{\prime}$ has no potential if $d>1$ (e.g., $W(x, y)$ i.i.d. for $|x-y|=1$ and $(x, y) \neq\left(x^{\prime}, y^{\prime}\right)$ will suffice $)$.

## 2. DEMONSTRATION OF THEOREM 1

One can prove Theorem 1 by using an argument similar to that for Theorem A. For $X(n)$, the basic plan was motivated by the guess that the largest crater a particle falls into before leaving the ball of radius $r$ is of order $c \log r$, where $c=-2 / \log (1-\varepsilon)$ for $d>1$. $(X(n)$ should visit on the order of $r^{2}$ sites before leaving the ball.) The time it takes to climb out of this crater is of order $e^{\beta c \log r}=r^{\beta c}$. Inverting, one obtains (7), although one actually needs the somewhat stronger assumption $\beta \geqslant 2(N+d+1)$. For $Y(n)$, with $0 \leqslant W(x, y) \leqslant 1 / 4$, the effect of $W$ is compensated by choosing $\beta \geqslant C d(N+1) / \varepsilon$. The particle will tend to fall into the same size craters given by $V$ as before; increasing $\beta$ increases the "pull" of a crater enough to offset $W$.

The proof of Theorem 1 is organized as follows: Lemma 1 and Proposition 1 will give lower bounds on the rate a particle tends to fall into a crater. They correspond to the like-labeled statements in ref. 7. Once it is in a deep crater, we wish for the particle to remain trapped there for a substantial time. Proposition 2 of ref. 7 expresses the time to climb out of a hole (perhaps consisting of many craters) in terms of the equilibrium measure $\alpha(x)$ corresponding to $p(x, y)$; the presence of the potential $V$
allows one to compute $\alpha(x)$. The equilibrium measure for $Y(n)$ is not, however, computable in terms of $V^{\prime}$. (This measure will not in general be "close" to $\alpha(x)$.) So the approach employed in ref. 7 will not work here. We give a different argument in Propositions 2 and 3 . Theorem 1 is then shown using Propositions 1-3.

We continue to use the notation employed in Section 1. We define $\varphi_{k}$, $V^{\prime}, \alpha^{\prime}$, and $p^{\prime}$ as before. Denote by $\mathscr{V}^{\prime}$ the $\sigma$-algebra generated by $V^{\prime}$. As usual, $\Omega$ will denote the probability space and $\omega$ its elements.

Set $B(r)=\left\{x \in \mathbf{Z}^{d}:|x|<r\right\}$. As $r$ increases, $B(r)$ will with high probability contain deeper and more numerous craters. A particle executing the motion $Y(n)$ should on occasion fall into such deep craters. To be more explicit, introduce $a_{i}$ and $b_{i}$ with

$$
\begin{equation*}
a_{i}=[\log i], \quad b_{i}=a_{3}+\cdots+a_{i}, \tag{12}
\end{equation*}
$$

for $i \geqslant 3$, with $[w]$ denoting the integer part of $w \in \mathbf{Z}$. From $b_{i}$, define the sets

$$
\begin{equation*}
B_{i}=\left\{x \in \mathbf{Z}^{d}:|x|<b_{i}\right\} \tag{13}
\end{equation*}
$$

and $A_{i}=B_{1}-B_{i}$. By $\partial B_{i}$, we mean those $x \in \mathbf{Z}^{d}$ with $\operatorname{dist}\left(B_{r}, x\right)=1$. Since we are unable to say much about the motion of $Y(n)$, crude arguments regarding the placement of deep craters are required. In Proposition 1, we give a lower bound on the probability that before leaving $B_{i}, Y(n)$ falls at least to depth $a_{i}$ in a prescribed manner. Although this probability is small, it is not too much smaller than $p_{a}$, and the event will occur with probability close to one for some $B_{i}$ satisfying $B(r / 4) \leqslant B_{i} \leqslant$ $B(r / 2)$, if $r$ is large.

We will find it useful to define

$$
A(x)=\left\{z: \varphi_{k(z)}(x-z)=V(x)\right\}
$$

if $V(x)<0$. We will then say that " $x$ is influenced by $A(x)$." Note that $A(x) \neq \varnothing$, and that for $|y-x|=1, V(y)=V(x)-1$ iff $|y-z|=|x-z|-1$ for some $z \in A(x)$. In this case, $A(y) \subset A(x)$.

Lemma 1. Fix $V, h$, and $x_{0}$, and suppose that $x_{0}$ is influenced by $A$ with $\operatorname{dist}\left(A, x_{0}\right) \geqslant h$. For $\beta \geqslant 4 \log 6 d, 0 \leqslant W(x, y) \leqslant 1 / 4$, and $Y(0)=x_{0}$,

$$
\begin{equation*}
P[V(Y(j))=V(Y(0))-j, j=1, \ldots, h] \geqslant(3 / 4)^{h} . \tag{14}
\end{equation*}
$$

Proof. Let $\mathscr{P}_{m}$ denote the set of paths $\left(x_{0}, \ldots, x_{m}\right)\left(\right.$ i.e., $\left.\left|x_{i}-x_{,},\right|=1\right)$ with $V\left(x_{j}\right)=V\left(x_{0}\right)-j$ for $j=1, \ldots, m$. For given $\left(x_{0}, \ldots, x_{m} \quad 1\right) \in \mathscr{Y}_{m}$, $m \leqslant h$, let

$$
B=\left\{x_{m}:\left(x_{0}, \ldots, x_{m}\right) \in \mathscr{P}_{m}\right\} .
$$

Since $x_{0}$ is influenced by $A$ and $\operatorname{dist}\left(A, x_{0}\right) \geqslant h, B$ is not empty.

Note that

$$
V\left(x_{m}\right) \geqslant V\left(x_{m-1}\right) \quad \text { if } \quad x_{m} \notin B .
$$

So

$$
\alpha^{\prime}\left(x_{m}, x_{m}\right) \leqslant 1 \quad \text { if } \quad x_{m} \notin B
$$

On the other hand,

$$
\alpha^{\prime}\left(x_{m} \quad, \quad, x_{m}\right) \geqslant e^{0 / 4} \quad \text { if } \quad x_{m} \in B
$$

Therefore, if $\beta \geqslant 4 \log 6 d$,

$$
\sum_{m_{m} \in B} p^{\prime}\left(x_{m} \quad, \quad, x_{m}\right) \geqslant|B| e^{\beta / 4} /\left(2 d+|B| e^{\beta / 4}\right) \geqslant 3 / 4 .
$$

Inequality (14) follows by induction.
We will find it convenient to introduce two variants of $V(x)$. Let

$$
\begin{align*}
& V_{i}(x)=\min _{z \in B_{i}} \varphi_{k(z)}(x-z)  \tag{15}\\
& \tilde{V}_{i}(x)=V_{i}(x) \wedge\left(\operatorname{dist}\left(\partial B_{i}, x\right)-a_{i}\right)
\end{align*}
$$

$V_{i}(x)$ measures the potential at $x$ by ignoring the effect of craters outside $B_{i} ; \tilde{V}_{i}(x)$ measures the resulting potential if one in addition includes the effect of a crater of depth $a_{i}$ at a site $z \in \partial B_{i}$ with

$$
\begin{equation*}
|z-x|=\operatorname{dist}\left(\partial B_{i}, x\right) \tag{16}
\end{equation*}
$$

Equations (15) are used in Proposition 1 in the context of $\sigma_{i}$ (defined below). Also, for Proposition 1, let

$$
\begin{equation*}
T_{i}=\min \left\{n:|Y(n)|=b_{i}\right\} \tag{17}
\end{equation*}
$$

and

$$
\begin{align*}
\sigma_{i} & =T_{i} \wedge \min \left\{n:|Y(n)|>b_{i} \quad, \quad V_{i}(Y(n)) \neq \tilde{V}_{i}(Y(n))\right\},  \tag{18}\\
\tau_{i} & =\min \left\{n: V(Y(n)) \leqslant-a_{i}\right\} .
\end{align*}
$$

(If a set is empty, assign the value $\infty$.) The quantity inside min $\{\cdot\}$ in the definition of $\sigma_{i}$ is the first time at which $Y$ visits a site in the annulus $A_{i}$ which would be influenced by a crater of depth $a_{i}$ at a site $z \in \partial B_{i}$ (if it is not already influenced by a yet deeper crater outside $B_{i}$ ). Note that under fixed $V^{\prime}$, these are all stopping times. Lastly, define

$$
\begin{equation*}
G_{i}=\left\{\omega: \tau_{i} \leqslant T_{i}\right\} \tag{19}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathscr{G}_{i}=\sigma\left(G_{1}, \ldots, G_{i}\right), \quad \mathscr{V}_{i}^{\prime}=\sigma\left(V, W,\left\{Y(n): n \leqslant \sigma_{i}\right\}\right) . \tag{20}
\end{equation*}
$$

$G_{i}$ is the event that $Y$ has fallen deeply into a hole before leaving $\bar{B}_{i}=$ $B_{i} \cup \partial B_{i}$. It is easy to check that $\mathscr{G}_{i-1} \subset \mathscr{V}_{i}^{\prime}$.

Proposition 1. For $G_{i}, \mathscr{G}_{i-1}$ as defined in (19)-(20) and $\beta \geqslant 4 \log 6 d$,

$$
\begin{equation*}
P\left[G_{i} \mid \mathscr{G}_{i-1}\right] \geqslant \delta \varepsilon(3(1-\varepsilon) / 4)^{a_{t}} . \tag{21}
\end{equation*}
$$

The proof of Proposition 1 is identical to that in ref. 7, and is omitted. The main idea is that since $Y$ conditioned on $\mathscr{V}^{\prime}$ is a Markov chain, one can apply the strong Markov property to $Y$ at time $\sigma_{i}$. One can check that either $V\left(Y\left(\sigma_{i}\right)\right) \leqslant-a_{i}$, which implies $\omega \in G_{i}$, or $\sigma_{i} \leqslant T_{i}$. Choose $Z_{i}$ so that $\left|Z_{i}\right|=b_{i}$ and

$$
\begin{equation*}
\left|Z_{i}-Y\left(\sigma_{i}\right)\right|=b_{i}-\left|Y\left(\sigma_{i}\right)\right| \tag{22}
\end{equation*}
$$

On

$$
K_{i}=\left\{\omega: k\left(Z_{i}\right)=a_{i}\right\}
$$

$X\left(\sigma_{i}\right)$ is influenced by $Z_{i}$ or some other point not in $B_{i}$. So one can apply Lemma 1 to show that on $K_{i}$,

$$
P\left[G_{i} \mid \mathscr{V}_{i}^{\prime}\right] \geqslant(3 / 4)^{a_{1}} .
$$

One can check that

$$
P\left[K_{i} \mid \mathscr{G}_{i} 1\right]=\delta \varepsilon(1-\varepsilon)^{a_{i}}
$$

Therefore,

$$
P\left[\begin{array}{ll}
G_{i} \mid \mathscr{G}_{i} & 1
\end{array}\right] \geqslant E\left[1_{K_{i}} P\left[G_{i} \mid \mathscr{V}_{i}^{\prime}\right] \mid \mathscr{S}_{i} \quad 1\right] \geqslant \delta \varepsilon(3(1-\varepsilon) / 4)^{a_{i}} .
$$

In Proposition 1, we gave a lower bound on the probability that $Y(n)$ falls at least to depth $a_{i}$ before leaving $B_{i}$. In (23) of Proposition 2, we give a lower bound on the time required for $Y(n)$ to rise from a given depth under certain regularity assumptions involving the size of nearby craters. (These assumptions ensure that the motion of a particle is locally influenced by only a single crater, which allows a simple computation of the bound.) This provides the upper bound in (24) on how far $Y(n)$ can move by a given time. After the regularity assumptions are examined in Proposition 3, Proposition 1 and (24) will be applied to demonstrate Theorem 1.

Proposition 2. Fix $V$ so that $V(0) \leqslant-f$. Also assume that (i) no craters of depth at least $h$ intersect $\widetilde{B}=\{x:|x|<2 m h\}$ and (ii) at most $m$ craters of depth at least $f / 2$ intersect $\widetilde{B}$, where $f, h, m>0$. If $Y(0)=0$, then

$$
\begin{equation*}
P[V(Y(j)) \geqslant-f / 2 \text { for some } j \leqslant n] \leqslant 4 d n \exp \left\{\frac{\beta}{8}\left(2-\frac{f}{m}\right)\right\} \tag{23}
\end{equation*}
$$

and

$$
\begin{equation*}
P\left[\max _{i \leqslant n}|Y(j)| \geqslant 2 m h\right] \leqslant 4 d n \exp \left\{\frac{\beta}{8}\left(2-\frac{l}{m}\right)\right\} \tag{24}
\end{equation*}
$$

for all $n>0$.
Proof. First note that on account of (i) and (ii), there are no paths connecting 0 with $\widetilde{B}^{c}$ which remain strictly below the level $-f / 2$. For such a path must remain in the above $m$ craters until reaching $\tilde{B}^{\text {}}$, whereas each such crater has diameter at most $2 h$. Consequently, (24) follows from (23).

To demonstrate (23), consider the set $E$ of those depths $g>f / 2$ for which if $x \in \tilde{B}$ with $V(x)=-g$, then there is at least one neighbor $y$ of $x$ with $I(y)=-g-1$. If $g \notin E$, then some $x$ with $V(x)=-g$ is at the center of a crater. So, by (ii),

$$
\left|g \in H^{\prime}: g>f / 2\right| \leqslant m
$$

One can therefore choose an interval $J=\left(g_{0}-f / 2 m, g_{0}\right)$ with $g_{0} \in[f / 2+f / 2 m, f]$ so that $g \in J$ implies that $g \in E$. That is, there is an unbroken sequence of depths at least $f / 2 m$ long so that only perhaps the greatest depth $g_{0} \in E^{c}$.

We can now use a standard argument involving martingles. Set

$$
\begin{equation*}
M(j)=\exp \left\{c\left(V(Y(j))+g_{0}\right)\right\} \tag{25}
\end{equation*}
$$

where $c>0$. For $x \in J$, there is at least one neighbor $y$ with $V(y)=$ $V(x)-1$. For such $y, \alpha^{\prime}(x, y) \geqslant e^{\beta / 4}$, whereas for other neighbors, $\alpha^{\prime}(x, y) \leqslant 1$. So for $Y(j) \in \widetilde{B}$ with $V(Y(j)) \in J$,

$$
\frac{E[M(j+1) \mid M(j)]}{M(j)} \leqslant \frac{e^{\beta / 4} e^{c}+2 d e^{2}}{e^{\beta / 4}+2 d}
$$

one can check that for $c=\beta / 4-\log 2 d$, this equals 1 . Set

$$
\tilde{M}(j)=M(j)-e^{\beta / 4} j
$$

The corresponding inequality

$$
\begin{equation*}
E[\tilde{M}(j+1) \mid \tilde{M}(j)] \leqslant \tilde{M}(j) \tag{26}
\end{equation*}
$$

holds for $Y(j) \in \widetilde{B}, V(Y(j)) \in J$. It is easy to check that for this value of $c$, (26) also holds for $V(Y(j)) \leqslant-g_{0}$. So $\tilde{M}(j)$ is a supermartingale for $j \leqslant T$, the first time at which $V(Y(j)) \geqslant-g_{0}+f / 2 m$.

Now, since $V(0) \leqslant-f, \tilde{M}(0) \leqslant 1$. So by the Optional Sampling Theorem,

$$
E[\tilde{M}(n \wedge T)] \leqslant \tilde{M}(0) \leqslant 1
$$

for all $n$. Consequently by Chebychev's inequality,

$$
\begin{aligned}
P[T \leqslant n] & \leqslant e^{-c f / 2 m} E[M(n \wedge T)] \\
& \leqslant\left(e^{\beta / 4} n+1\right) e^{-c f / 2 m} \\
& \leqslant 4 d n \exp \left\{\frac{\beta}{8}\left(2-\frac{f}{m}\right)\right\}
\end{aligned}
$$

In the proof of Theorem 1, we will set

$$
\begin{equation*}
f=\frac{1}{2} \log r, \quad h=\frac{d+1}{\varepsilon} \log r, \quad m=\frac{10 d}{\varepsilon} . \tag{27}
\end{equation*}
$$

We will therefore need to establish upper bounds on the probabilities that conditions (i) and (ii) of Proposition 2 are violated for these values. This is done in Proposition 3. We set $\widetilde{B}(r)=\{x:|x|<2 m h\}$. We denote by $H_{r}$ the set of $V$ for which no crater of depth at least $h$ intersects $B(r)$ and by $F_{r}$ the set of $V$ for which there are at most $m$ craters of depth at least $f / 2$ which intersect $\widetilde{B}(r)+x$ for all $x \in B(r)$. Here, $f, h$, and $m$ are chosen as in (27) and $+x$ denotes translation by $x$.

Proposition 3. (i) $P\left[H_{r}^{c}\right] \leqslant C_{1} / r$ and (ii) $P\left[F_{r}^{*}\right] \leqslant C_{2} / r$ for appropriate $C_{1}$ and $C_{2}$ depending on $\varepsilon$ and $d$.

Proof. The left side of (i) is at most

$$
\begin{align*}
& \sum_{i=0}^{\infty}(2(r+[h]+j))^{d} \delta \varepsilon(1-\varepsilon)^{[h]+i} 1 \\
& \quad \leqslant(8 r)^{d} \varepsilon(1-\varepsilon)^{[h]} \quad: \sum_{i=0}^{\infty}(1-\varepsilon)^{j}+8^{d} \varepsilon(1-\varepsilon)^{[h]} \quad \sum_{j=0}^{\infty} j^{d}(1-\varepsilon)^{\prime} . \tag{28}
\end{align*}
$$

The first term on the right side equals

$$
(8 r)^{d}(1-\varepsilon)^{\lceil h\rceil} .
$$

The second term is at most

$$
\begin{gather*}
8^{d} \varepsilon(1-\varepsilon)^{[h]} \quad \sum_{j=1}^{\infty} j(j+1) \cdots(j+d+1)(1-\varepsilon)^{\prime} \\
=d!(8 / \varepsilon)^{d}(1-\varepsilon)^{[h]} \tag{29}
\end{gather*}
$$

Plugging in $h=((d+1) / \varepsilon) \log r$, one can check that their sum is at most $C_{1} / r$ for appropriate $C_{1}$ depending on $\varepsilon$ and $d$.

The left side of (ii) can be evaluated similariy. It is at most

$$
\left.\begin{array}{rl}
(2 r)^{d} & \sum_{n=0}^{\infty} \cdots \sum_{i_{m}=0}^{\infty} \prod_{t=1}^{m}\left(2\left(2 m h+[f / 2]+j_{i}\right)\right)^{d} \delta \varepsilon(1-i)^{|/ / 2|+j} \\
& =(2 r)^{d}\left(\sum_{0}^{\infty}(2(2 m h+[f / 2]+j))^{d} \delta e(1-\varepsilon)^{(/ 2)}+1\right.
\end{array}\right)^{m} .
$$

Reasoning as between (28) and (29) shows that this is

$$
\leqslant(2 r)^{d}\left\{\left((16 m h)^{d}+d!(8 / \varepsilon)^{d}\right)(1-\varepsilon)^{[/ / 2]} \quad\right\}^{m} .
$$

Plugging in $f=\frac{1}{2} \log r$ and $h=((d+1) / \varepsilon) \log r$, one can check that this is

$$
\leqslant C^{\prime}(\log r)^{d m} r^{d \quad u m / 5}
$$

where $C^{\prime}$ depends on $\varepsilon, d$, and $m$. For $m=10 d / \varepsilon$, this is clearly

$$
\leqslant C_{2} / r
$$

for appropriate $C_{2}$ depending on $\varepsilon$ and $d$.
We now prove Theorem 1 by using Propositions $1-3$ as sketched earlier.

Theorem 1. Suppose that $0<\delta<1,0<\varepsilon \leqslant 1 / 2$, and $N>0$. If $\beta \geqslant C d(N+1) / \varepsilon$ for appropriate $C$, and $0 \leqslant W(x, y) \leqslant 1 / 4$, then

$$
\begin{equation*}
P\left[\max _{j \leqslant n}|Y(j)| \geqslant n^{1 / N}\right] \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty \tag{30}
\end{equation*}
$$

Proof. Fix $r$ and set

$$
\begin{aligned}
l & =\min \left\{i: b_{i} \geqslant r / 4\right\}, \\
L & =\max \left\{i: b_{i}<r / 2\right\} .
\end{aligned}
$$

It is easy to check that for $i \leqslant L$,

$$
\begin{equation*}
a_{i} \leqslant \log r \tag{31}
\end{equation*}
$$

and consequently

$$
\begin{equation*}
l \geqslant \frac{r}{4 \log r}, \quad L-l \geqslant \frac{r}{4 \log r}-1 . \tag{32}
\end{equation*}
$$

Now, repeated application of Proposition 1 shows that

$$
\begin{equation*}
P\left[\bigcap_{i=1}^{L} G_{i}^{c}\right] \leqslant \prod_{i=l}^{L}\left[1-\delta \varepsilon(3(1-\varepsilon) / 4)^{a_{i}}\right] \leqslant \exp \left[-\sum_{i=1}^{L} \delta \varepsilon(3(1-\varepsilon) / 4)^{a_{i}}\right] . \tag{33}
\end{equation*}
$$

Plug in (31) and (32), and note that $\log (3(1-\varepsilon) / 4)>-1$ for $\varepsilon \leqslant 1 / 2$. This shows that (33) is at most

$$
\begin{equation*}
\exp \left[-\delta \varepsilon\left(\frac{r}{4 \log r}-1\right) r^{\log [3(1-\varepsilon) / 4]}\right] \leqslant \exp \left(-\delta \varepsilon r^{\eta}\right) \tag{34}
\end{equation*}
$$

for large enough $r$ and appropriate $\eta>0$. Setting $G=\bigcup_{i=1}^{L} G_{i}$, we obtain from (33)-(34) that

$$
\begin{equation*}
P\left[G^{c}\right] \leqslant \exp \left(-\delta c r^{\eta}\right) \tag{35}
\end{equation*}
$$

We restrict our attention to $Y(n)$ on $G$. On this set,

$$
\begin{equation*}
\tau_{i} \leqslant T_{t} \quad \text { for some } \quad l \leqslant i \leqslant L . \tag{36}
\end{equation*}
$$

Denote by $I$ the first such $i$, and set $Y_{I}=Y\left(\tau_{l}\right)$. From the definition of $\tau$ and $a$,

$$
\begin{equation*}
V\left(Y_{i}\right) \leqslant-a_{l} \leqslant-[\log l] . \tag{37}
\end{equation*}
$$

On account of (32), for $r$ not too small, this is

$$
\begin{equation*}
\leqslant-\frac{1}{2} \log r \tag{38}
\end{equation*}
$$

We will denote by $\mu_{V}$, the subprobability measure on $B(r)$ induced by $Y_{1}$ (restricted to $G$ ) for fixed $V^{\prime}$ and by $Z(m)$ a copy of $X(m)$ with initial distribution given by $\mu_{V^{\prime}}$.

Now consider $\omega \in H_{r} \cap F_{r}$. We will plug Proposition 3 into (24) of Proposition 2 with $f=\frac{1}{2} \log r, h=((d+1) / \varepsilon) \log r$, and $m=10 d / \varepsilon$, and with $Y(0)=0$ replaced by $Z(0)=Y_{1}$. It is not hard to verify that the conditions in Proposition 2 are satisfied. On account of (37)-(38), $Z(0) \leqslant-f$. Note that since $Y_{I} \in B(r / 2)$,

$$
\begin{equation*}
\tilde{B}(r)+Y, \subset B(r) \tag{39}
\end{equation*}
$$

for $r$ not too small. Therefore, since $V \in H_{r}$, it follows that (i) of Proposition 2 holds. Also, since $V \in F_{r}$, (ii) of Proposition 2 holds. So, applying (24) and (39), we obtain

$$
\begin{equation*}
P_{\mu}\left[\max _{i \leqslant n}|Z(j)| \geqslant r \mid \psi^{\prime}\right] \leqslant 4 \ln \exp \left\{\frac{\beta}{8}\left(2-\frac{\varepsilon \log r}{20 d}\right)\right\} \tag{40}
\end{equation*}
$$

for $\omega \in H_{r} \cap F_{r}$.
The right side of (40) does not depend on $V^{\prime}$. Setting $n=r^{N}$ with $N$ fixed, one gets

$$
\begin{align*}
& P_{\mu \nu}\left[\max _{1<r^{N}}|Z(j)| \geqslant r \mid \mathscr{V}^{\prime}\right] \\
& \leqslant 4 d e^{\beta / 4} r^{N \quad 1 . / 16(6) d} \leqslant C_{3} / r \tag{41}
\end{align*}
$$

for $\beta \geqslant 160 d(N+1) / \varepsilon$, where $C_{3}$ does not depend on $r$. Conditioned on $\not{ }^{\prime}$, the process $Y$ is strong Markov. The strong Markov property therefore implies that

$$
P\left[\max _{j \leqslant r^{N}}|Y(j)| \geqslant r ; G \mid \mathscr{V}^{\prime}\right] \leqslant C_{3} / r
$$

for $\omega \in H, \cap F_{r}$. Consequently,

$$
P\left[\max _{j \leqslant r^{N}}|Y(j)| \geqslant r ; G \mid \mathscr{V}^{\prime}\right] \leqslant C_{3} / r .
$$

Together with (35) and Proposition 3, this shows that

$$
\begin{align*}
P\left[\max _{j \leqslant r^{N}}|Y(j)| \geqslant r\right] & \leqslant \exp \left(-\delta \& r^{\eta}\right)+\left(C_{1}+C_{2}+C_{3}\right) / r \\
& \rightarrow 0 \quad \text { as } \quad r \rightarrow \infty \tag{42}
\end{align*}
$$

Inverting, one obtains

$$
P\left[\max _{J \leqslant n}|Y(j)| \geqslant n^{1 / N}\right] \rightarrow 0 \quad \text { as } \quad n \rightarrow \infty
$$

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