# R-POSITIVITY, QUASI-STATIONARY DISTRIBUTIONS AND RATIO LIMIT THEOREMS FOR A CLASS OF PROBABILISTIC AUTOMATA 

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#### Abstract

We prove that certain (discrete time) probabilistic automata which can be absorbed in a "null state" have a normalized quasi-stationary distribution (when restricted to the states other than the null state). We also show that the conditional distribution of these systems, given that they are not absorbed before time $n$, converges to an honest probability distribution; this limit distribution is concentrated on the configurations with only finitely many "active or occupied" sites.

A simple example to which our results apply is the discrete time version of the subcritical contact process on $\mathbb{Z}^{d}$ or oriented percolation on $\mathbb{Z}^{d}$ (for any $d \geq 1$ ) as seen from the "leftmost particle." For this and some related models we prove in addition a central limit theorem for $n^{-1 / 2}$ times the position of the leftmost particle (conditioned on survival until time $n$ ).

The basic tool is to prove that our systems are $R$-positive-recurrent.


1. Introduction and principal results. Let $\left\{X_{n}\right\}_{n \geq 0}$ be a Markov chain on a countable state space $S$, with an absorbing state $s_{0}$. We shall deal exclusively with the discrete time case in this paper, but we believe that all the results and proofs have analogues in the continuous time case. We shall write $S_{0}$ for $S \backslash\left\{s_{0}\right\}$. As usual, we denote the probability measure governing $X$. when conditioned to start at $X_{0}=x$ by $P_{x}$. Then for any probability measure $\nu$ on $S$,

$$
P_{\nu}=\sum_{x \in S} \nu(x) P_{x}=\nu P
$$

is the measure which governs $X$. when the initial distribution is $\nu$. Let $E_{x}$ and $E_{\nu}$ denote expectation with respect to $P_{x}$ and $P_{\nu}$, respectively. The transition probabilities are

$$
\begin{equation*}
P(x, y)=P_{x}\left\{X_{1}=y\right\} \quad \text { and } \quad P^{n}(x, y)=P_{x}\left\{X_{n}=y\right\} . \tag{1.1}
\end{equation*}
$$

It is convenient to introduce the restriction $\widehat{P}$ of $P$ to $S_{0} \times S_{0}$. Because $s_{0}$ is absorbing, we then also have

$$
(\widehat{P})^{n}(x, y)=P_{x}\left\{X_{n}=y\right\}=P^{n}(x, y) \quad \text { for } x, y \in S_{0} .
$$

[^0]Throughout we assume that
$\widehat{P}$ is irreducible; that is, for all $x, y \in S_{0}$ there exists an
$n=n(x, y)$ with $P^{n}(x, y)>0$

The absorption time is

$$
\begin{equation*}
T=T\left(s_{0}\right)=\inf \left\{n \geq 0: X_{n}=s_{0}\right\} \tag{1.3}
\end{equation*}
$$

and we assume that absorption is certain. This means that for some $x \in S_{0}$ (and hence all $x$ )

$$
\begin{equation*}
P_{x}\{T<\infty\}=1 \tag{1.4}
\end{equation*}
$$

The fact that $s_{0}$ is absorbing means of course that

$$
\begin{equation*}
X_{n}=s_{0} \quad \text { for all } n \geq T \quad \text { and } \quad P^{n}\left(s_{0}, x\right)=\delta\left(s_{0}, x\right), \quad n \geq 0 \tag{1.5}
\end{equation*}
$$

A normalized quasi-stationary distribution for $X$ is a probability measure $\nu$ on $S_{0}$ which satisfies the invariance condition

$$
\begin{equation*}
\nu \widehat{P}^{n}=r(n) \nu \tag{1.6}
\end{equation*}
$$

where, necessarily,

$$
\begin{equation*}
r(n)=\sum_{y \in S_{0}} \nu \widehat{P}^{n}(y)=P_{\nu}\{T>n\} \tag{1.7}
\end{equation*}
$$

Note that these distributions are conditionally invariant, in the sense that

$$
P_{\nu}\left\{X_{n}=y \mid T>n\right\}=\nu(y) \quad \text { for all } y \in S_{0}
$$

A normalized quasi-stationary distribution $\nu_{0}$ is called minimal if

$$
\begin{equation*}
E_{\nu_{0}} T=\inf \left\{E_{\nu} T: \nu \text { a normalized quasi-stationary distribution }\right\} \tag{1.8}
\end{equation*}
$$

The interest in normalized quasi-stationary distributions arises from the fact that if for some initial distribution $\nu$ on $S_{0}$,

$$
\begin{equation*}
\frac{P_{\nu}\left\{X_{n}=y\right\}}{P_{\nu}\{T>n\}} \rightarrow \mu(y), \quad y \in S_{0} \tag{1.9}
\end{equation*}
$$

for some probability distribution $\mu$ on $S_{0}$, then $\mu$ is necessarily a normalized quasi-stationary distribution. [See Seneta and Vere-Jones (1966), Theorem 4.1]. If the limit in (1.9) exists, we shall call it a Yaglom limit, because Yaglom (1947) proved the existence of this limit for subcritical branching processes (when $\nu$ is concentrated on one point $x$ ). Such a limit $\mu$ is also called a conditional limit distribution in the literature. There is an extensive literature discussing the existence of normalized quasi-stationary distributions and the Yaglom limit; see Ferrari, Kesten, Martínez and Picco (1995) for some references. An additional recent reference is Roberts and Jacka (1994). Recently, Ferrari, Kesten, Martínez and Picco (1995) gave a necessary and sufficient
condition for the existence of a normalized quasi-stationary distribution for chains $X$. which satisfy, in addition to (1.2) and (1.4), the condition

$$
\begin{equation*}
\lim _{x \rightarrow \infty} P_{x}\{T \leq t\}=0 \text { for all fixed } t<\infty \tag{1.10}
\end{equation*}
$$

This of course means that for all $t<\infty$ and $\varepsilon>0, P_{x}\{T \leq t\} \leq \varepsilon$ for all but finitely many $x$. [Actually Ferrari, Kesten, Martínez and Picco (1995) deal with the continuous time case, but their results carry over to discrete time; see Kesten (1995), Theorem A, for a statement of the discrete time result.] Unfortunately, (1.10) is rarely fulfilled for chains $X$. which describe interacting particles which also have a spatial position or which can be of infinitely many types. The absorbing state is the state in which no particle is present. Typically in such models, the probability of absorption in unit time from any of the infinitely many states in which only one particle is present is bounded away from 0-infinitely many states, because the single particle can have infinitely many positions or types; the above phenomenon occurs if the probability for a particle to die in one time unit is uniformly bounded away from 0 . A simple special example (which was, in fact, the principal motivation for this paper) is the subcritical contact process or oriented percolation on $\mathbb{Z}^{d}$ as seen from the "leftmost" particle (with state space $S$ a certain collection of finite subsets of $\mathbb{Z}^{d}$ containing $\{0\}$, plus the empty set). Our aim here is to prove the existence of a normalized quasi-stationary distribution and a Yaglom limit for a class of interacting particle systems and probabilistic automata which includes oriented percolation on $\mathbb{Z}^{d}, d \geq 1$. This will be done by proving those chains $R$-positive-recurrent. (A more detailed description of some of these examples and the role of the leftmost particle is given before the statement of Theorem 2 below; full details are in Section 4.) Pakes (1995) investigated some other examples in which (1.10) may fail.

We remind the reader of some basic facts which hold solely under the assumption (1.2) [see Vere-Jones (1967)]:

The period $p:=$ g.c.d. $\left\{n: \widehat{P}^{n}(x, x)>0\right\}$ is finite and is the same for all $x \in S_{0}$;
$S_{0}$ can be decomposed into $p$ disjoint subclasses $S_{0,0}, \ldots, S_{0, p-1}$ so that $\widehat{P}(x, y)>0$ only if $x \in S_{0, j}, \quad y \in S_{0, j+1}$ for some $0 \leq$ $j \leq p-1$ (here and in the sequel we take $S_{0, j_{1}}=S_{0, j_{2}}$ when $j_{1} \equiv j_{2}(\bmod p)$; the $S_{0, j}$ are referred to as the cyclically moving subclasses or periodic subclasses);
if $x \in S_{0, i}$ and $y \in S_{0, j}$, then $1 / R:=\lim _{n \rightarrow \infty}\left[\widehat{P}^{n p+j-i}(x, y)\right]^{1 / n p}$ exists and the value of $R$ is independent of $x, y, i$ and $j$;

$$
\begin{equation*}
1 \leq R<\infty . \tag{1.14}
\end{equation*}
$$

The chain $X$. and $\widehat{P}$ are called $R$-recurrent if for some $x \in S_{0}$,

$$
\begin{equation*}
\sum_{n=0}^{\infty} R^{n p} \widehat{P}^{n p}(x, x)=\infty \tag{1.15}
\end{equation*}
$$

and it is called $R$-positive-recurrent if, in addition, for some $x \in S_{0}$,

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} R^{n p} \widehat{P}^{n p}(x, x)>0 \tag{1.16}
\end{equation*}
$$

Again (1.15) and (1.16) either hold for all $x$ in $S_{0}$ simultaneously or for no $x$. If (1.15) and (1.16) hold, then there exist functions $f: S_{0} \rightarrow(0, \infty), \mu: S_{0} \rightarrow$ $(0, \infty)$ which satisfy

$$
\begin{align*}
& \widehat{P}^{n} f(x)=\sum_{y \in S_{0}} \widehat{P}^{n}(x, y) f(y)=R^{-n} f(x), \quad x \in S_{0}, n \geq 0,  \tag{1.17}\\
& \mu \widehat{P}^{n}(y)=\sum_{x \in S_{0}} \mu(x) \widehat{P}^{n}(x, y)=R^{-n} \mu(y), \quad y \in S_{0}, n \geq 0,  \tag{1.18}\\
& \sum_{x \in S_{0, j}} f(x) \mu(x)=1, \quad 0 \leq j \leq p-1, \quad \sum_{x \in S_{0}} f(x) \mu(x)=p \tag{1.19}
\end{align*}
$$

and for $x \in S_{0, i}, y \in S_{0, j}$ for some $0 \leq i, j \leq p-1$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} R^{n p+j-i} \widehat{P}^{n p+j-i}(x, y)=f(x) \mu(y) . \tag{1.20}
\end{equation*}
$$

The functions $f$ and $\mu$ are uniquely determined up to multiplicative constants by (1.17) and (1.18). Equation (1.18) does not say that $\mu$ is a (multiple of a) normalized quasi-stationary distribution, because $\sum_{x \in S_{0}} \mu(x)$ may diverge. In order to obtain a normalized quasi-stationary distribution and for the Yaglom limit relation (1.9) we need

$$
\begin{equation*}
\sum_{x \in S_{0}} \mu(x)<\infty \tag{1.21}
\end{equation*}
$$

[see Seneta and Vere-Jones (1966), Theorem 3.1]. It is an important step in our proof to show that (1.21) indeed holds under the conditions of Theorem 1 below, which is our principal result.

Theorem 1. Assume that the Markov chain $\left\{X_{n}\right\}_{n \geq 0}$ satisfies (1.2) and (1.4) as well as the following conditions:
there exist a nonempty set $\mathscr{U}_{1} \subset S_{0}$, an $\varepsilon_{0}>0$ and a constant
$C_{1}$ such that for all $x \in \mathscr{U}_{1}$ and all $n \geq 0, P_{x}\{T>n$, but $X_{l} \notin \mathscr{U}_{1}$ for all $\left.1 \leq l \leq n\right\} \leq C_{1}\left(R+\varepsilon_{0}\right)^{-n}$,
there exist a state $x_{0} \in \mathscr{U}_{1}$ and a constant $C_{2}$ such that for all $x \in \mathscr{U}_{1}$ and $n \geq 0, P_{x}\{T>n\} \leq C_{2} P_{x_{0}}\{T>n\}$
and
there exist a finite set $\mathscr{U}_{2} \subset S_{0}$ and constants $0 \leq n_{0}<$ $\infty, C_{3}>0$, such that for all $x \in \mathscr{U}_{1}, P_{x}\left\{X_{n} \in \mathscr{U}_{2}\right.$ for some $\left.n \leq n_{0}\right\} \geq C_{3}$.

Then $X$. is $R$-positive-recurrent, so that (1.17)-(1.20) hold for $p=$ period of $\widehat{P}$ and some strictly positive functions $f$ and $\mu$. Moreover, (1.21) is satisfied and

$$
\begin{equation*}
\tilde{\mu}(x):=\frac{\mu(x)}{\sum_{y \in S_{0}} \mu(y)}, \quad x \in S_{0} \tag{1.25}
\end{equation*}
$$

defines a minimal quasi-stationary distribution for $X$. Also, if $S_{0,0}, \ldots, S_{0, p-1}$ are the cyclically moving subclasses as in (1.12), and $x \in S_{0, i}$, then

$$
\begin{equation*}
\lim _{\substack{n \rightarrow \infty \\ n \equiv j(\bmod p)}} R^{n} P_{x}\{T>n\}=f(x) \sum_{y \in S_{0, i+j}} \mu(y) \tag{1.26}
\end{equation*}
$$

and for $x \in S_{0, i}, y \in S_{0, j}$,

$$
\begin{equation*}
\lim _{\substack{n \rightarrow \infty \\ n \equiv j-i(\bmod p)}} P_{x}\left\{X_{n}=y \mid T>n\right\}=\frac{\mu(y)}{\sum_{v \in S_{0, j}} \mu(v)} . \tag{1.27}
\end{equation*}
$$

Finally, for each $x \in S_{0}$ there exist an $\eta(x)>0$ and a $C_{4}(x)<\infty$ such that

$$
\begin{equation*}
P_{x}\left\{X_{n p}=x \text {, but } X_{s} \neq x \text { for } 1 \leq s \leq n p-1\right\} \leq C_{4}(x) \frac{1}{(R+\eta(x))^{n}} . \tag{1.28}
\end{equation*}
$$

Remark 1. Equation (1.27) can be easily generalized to give convergence of finite-dimensional distributions. That is, for $0<s_{1}<\cdots<s_{l}, x \in S_{0, i}$, $y_{1} \in S_{0, j_{1}}, \ldots, y_{l} \in S_{0, j_{l}}$ it holds that

$$
\begin{align*}
& \lim _{n \rightarrow \infty} P_{x}\left\{X_{\left\lfloor s_{r} n\right\rfloor p+j_{r}-i}=y_{r}, 1 \leq r \leq l \mid T>\left\lfloor s_{l} n\right\rfloor p+j_{l}-i\right\} \\
& \quad=\frac{\mu\left(y_{l}\right)}{\sum_{v \in S_{0, j_{l}}} \mu(v)} \prod_{r=1}^{l-1}\left[f\left(y_{r}\right) \mu\left(y_{r}\right)\right] . \tag{1.29}
\end{align*}
$$

Remark 2. A consequence of (1.29) and (1.19) is that for all $\varepsilon>0$ and $0<s \leq 1$, there exists a finite set $S(\varepsilon, s) \subset S$ such that for all $x \in S_{0}$,

$$
\lim _{n \rightarrow \infty} P_{x}\left\{X_{\lfloor s n\rfloor} \in S(\varepsilon, s) \mid T>n\right\} \geq 1-\varepsilon
$$

In other words, when conditioned on no absorption until time $n$, then the chain $X$. tends to spend "most of its time near $s_{0}$." This contrasts with some of the examples in Ferrari, Kesten, Martínez and Picco (1995). For example, let $\widetilde{X}_{n}$ be the absorbing random walk on $\mathbb{Z}_{+}$with 0 as absorbing state and which moves from $x \geq 1$ to $x+1$ or $x-1$ with probability $p$ and $q=1-p$, respectively, with $0<p<\frac{1}{2}$. Then $\left\{n^{-1 / 2} \widetilde{X}_{\lfloor s n\rfloor}\right\}_{0 \leq s \leq 1}$ conditioned on $T>n$
essentially behaves like a positive Brownian excursion (scaled so that it has length 1). In particular, for fixed $0<s<1$,

$$
P_{x}\left\{\left|\tilde{X}_{\lfloor s n\rfloor}<\varepsilon \sqrt{n}\right| T>n\right\}
$$

can be made small uniformly in $n$, by taking $\varepsilon$ small. Thus $\widetilde{X}$ conditioned on $T>n$ spends most of its time far away from the absorbing state 0 .

Theorem 1 has some value only if its hypotheses can be verified without too much trouble for a reasonable class of examples. The hypothesis which looks the most troubling is (1.22). The next theorem shows that this actually follows directly from known results in many instances of so-called probabilistic cellular automata, which are similar to the subcritical contact process. We begin with a Markov chain $\left\{\widetilde{X}_{n}\right\}_{n \geq 0}$ on a countable state space

$$
\begin{align*}
& \widetilde{S}=\widetilde{S}_{d, \kappa} \text { which is a subset of the collection of all functions } \\
& x: \mathbb{Z}^{d} \rightarrow\{0,1, \ldots, \kappa-1\} \quad \text { with } x(z) \neq 0 \text { for only finitely many } z . \tag{1.30}
\end{align*}
$$

Here $\kappa \geq 2$ is some fixed integer. Note that $\mathbb{Z}^{d}$ can be replaced by other lattices and this should indeed be done to treat the usual oriented percolation. For simplicity we restrict ourselves here to the state space (1.30). Let the state $\underline{0}$ with all components equal to 0 be the absorbing state. Thus

$$
\begin{equation*}
P\left\{\tilde{X}_{n}=x \mid \widetilde{X}_{0}=\underline{0}\right\}=\delta(\underline{0}, x), \quad n \geq 0, x \in \widetilde{S} . \tag{1.31}
\end{equation*}
$$

The space $\widetilde{S}$ and the transition probabilities have to be such that

$$
\begin{equation*}
\text { the transition probability matrix is irreducible on } \widetilde{S}_{0}=\widetilde{S} \backslash\{\underline{0}\} . \tag{1.32}
\end{equation*}
$$

(It is for this reason that $\widetilde{S}$ may have to be taken only as a subset of

$$
\left\{x \in\{0,1, \ldots, \kappa-1\}^{\mathbb{Z}^{d}}: x(z) \neq 0 \text { for only finitely many } z\right\}
$$

rather than the full set.) We further impose the following fairly common conditions on the transition probabilities:

Translation invariance: $P\left\{\tilde{X}_{1}=y \oplus u \mid \tilde{X}_{0}=x \oplus u\right\}=$ $P\left\{\widetilde{X}_{1}=y \mid \widetilde{X}_{0}=x\right\}$ for all $u \in \mathbb{Z}^{d}$, where the state $x \oplus u$ is specified by $(x \oplus u)(z)=x(z+u), z \in \mathbb{Z}^{d}$, and similarly for $y \oplus u ;$

Independence of coordinates: Conditionally on $\tilde{X}_{0}, \ldots, \widetilde{X}_{n}$, the coordinates $\left\{\widetilde{X}_{n+1}(z): z \in \mathbb{Z}^{d}\right\}$ of $\widetilde{X}_{n+1}$ are independent;

Finite range: There exists some $\rho<\infty$ such that the conditional distribution of $\widetilde{X}_{n+1}(z)$, given $\widetilde{X}_{n}$, depends only on $\widetilde{X}_{n}\left(z^{\prime}\right)$ for $\left\|z^{\prime}-z\right\|_{2} \leq \rho$.

It is trivial to check that the discrete time version of the contact process and oriented percolation (see Section 4) satisfy (1.31)-(1.35) for $\kappa=2$. In analogy with the terminology for that chain we shall call sites $z$ with $x(z)=0$ $(\neq 0)$ vacant (and occupied, respectively) in state $x$. It is convenient to think of a particle sitting at each occupied site, the particle having one of the types $1, \ldots, \kappa-1$. The chain $\widetilde{X}$. then describes motion, birth, death and change of type of a finite system of particles.

Theorem 1 will not apply to $\widetilde{X}$ itself. As illustrated by oriented percolation, we cannot expect the limit in (1.27) to have a nonzero value because the occupied sites of $\widetilde{X}_{n}$ will wander all over the space $\mathbb{Z}^{d}$ (even under the condition $\{T>n\}$ ). Also condition (1.24) will fail because particles cannot jump in fixed time $n_{0}$ from far away to a site which is occupied in one of the finitely many states in $\mathscr{U}_{2}$. Theorem 1 will only apply to the relative positions of the occupied sites. To make this precise, let $\xi(n)=\left\{\xi_{1}(n), \ldots, \xi_{\nu}(n)\right\}$ be the occupied sites at time $n$, that is the sites in $\mathbb{Z}^{d}$ with $\widetilde{X}_{n}\left(\xi_{i}(n)\right) \neq 0$. Here $\nu=\nu(n)$ is random. If no sites are occupied at time $n$, we take $\nu(n)=0, \xi(n)=\phi$. If $\nu>0$, the order of the $\xi_{i}(n)$ is arbitrary, except that we take $\xi_{1}$ to be the first among $\xi_{1}, \ldots, \xi_{\nu}$ in the lexicographical ordering on $\mathbb{Z}^{d}$. We call $\xi_{1}(n)$ the leftmost occupied position at time $n$. Finally define

$$
X_{n}=\left\{\begin{array}{l}
\underline{0}, \quad \text { if } \nu(n)=0,  \tag{1.36}\\
\widetilde{X}_{n} \oplus\left(-\xi_{1}(n)\right) .
\end{array}\right.
$$

Thus $X_{n}(z)=j$ precisely when $\widetilde{X}_{n}\left(z+\xi_{1}(n)\right)=j$. From the translation invariance (1.33) it is not hard to see that $\left\{X_{n}\right\}_{n \geq 0}$ is again a Markov chain (this is the only place where translation invariance is strictly needed). The state space $S$ of the $X$-chain is of the same form as that of $\widetilde{X}$. Only $X_{n}(0) \neq 0$ whenever $X_{n} \neq \underline{0}$, and $X_{n}(z) \neq 0$ can only occur at a $z$, which comes after 0 in the lexicographical ordering of $\mathbb{Z}^{d}$. The absorbing state for $X$. is again $\underline{0}$.

Here are some final definitions. For $x \in \widetilde{S}$ or $x \in S$,

$$
\begin{align*}
|x| & =\text { number of occupied sites in } x \\
& =\text { number of } z \in \mathbb{Z}^{d} \text { with } x(z) \neq 0 . \tag{1.37}
\end{align*}
$$

The use of the same symbol \| for points in $\widetilde{S}$ and in $S$ will not lead to any serious confusion, since in both cases $|x|$ represents the number of occupied sites. For similar reasons we shall use $T$ for the absorption time of $\widetilde{X}$ and of $X$. The process $\left\{\widetilde{X}_{n}\right\}$ is called strongly subcritical if

$$
\begin{equation*}
n^{d} \sup _{\substack{x \in \widetilde{S} \\ z \in \mathbb{Z}^{d}}} P_{x}\left\{\tilde{X}_{n}(z) \neq 0\right\} \rightarrow 0 \quad \text { as } n \rightarrow \infty . \tag{1.38}
\end{equation*}
$$

Theorem 2. Let $\widetilde{X}$ be a Markov chain on the space $\widetilde{S}$ of (1.30) whose transition probabilities satisfy (1.31)-(1.35). Let $X_{n}$ be defined by (1.36). Assume
that $\left\{\tilde{X}_{n}\right\}$ is strongly subcritical. Then for each $\gamma>0$ there exists a $K<\infty$ and a $C_{5}=C_{5}(\gamma)<\infty$ such that for all $x$ with $|x| \leq K$,

$$
\begin{equation*}
P_{x}\left\{T>n, \text { but }\left|X_{l}\right|>K \text { for all } 1 \leq l \leq n\right\} \leq C_{5} \gamma^{n}, \quad n \geq 0 \tag{1.39}
\end{equation*}
$$

Clearly (1.39) gives us (1.22) with $\mathscr{U}_{1}=\{x:|x| \leq K\}$ by taking $\gamma<R^{-1}$, say $\gamma=(R+1)^{-1}$. Thus, Theorem 2 reduces checking (1.22) to checking that our process is strongly subcritical. For some processes it is known that if they die out w.p. 1 [i.e., satisfy (1.4)], then they are automatically strongly subcritical. In other examples it should be possible to use a Peierls argument or simple comparisons with other chains to show that the chain $X$. is strongly subcritical, at least for certain choices of the transition probabilities. (See also the discussion at the end of Example A in Section 4.) Condition (1.24) is usually innocuous; as illustrated in the proof of Proposition 1, it just depends on probabilities of particles disappearing being bounded away from 0 . Condition (1.23), which is trivial for the contact process or oriented percolation (see Remark 8 in Section 4), turns out to be nontrivial in general. We verify it in Section 4 for some chains of the form described above. For the special case of the discrete time contact process or oriented percolation, all hypotheses are verified in this same section.

Finally, in Section 5 we derive a central limit theorem for the absolute location of the occupied sites in the $\widetilde{X}$-chain. The conclusion of Theorem 1 tells us that $X_{n}$, conditioned on $\{T>n\}$, has a limit distribution. Thus, for $X_{n}$ as in (1.36), this tells us that the number of particles and their distances to the leftmost occupied position $\xi_{1}(n)$ are tight (under the condition $\{T>n\}$ ). However, the absolute position of the occupied set, by which we simply mean $\xi_{1}(n)$ itself, will not be tight without normalization. In fact our last theorem shows that $\xi_{1}(n) / \sqrt{n}$ is usually asymptotically normally distributed and independent of $X_{n}$ (still under the condition $\{T>n\}$ ). This result follows almost immediately from standard central limit theorems for recurrent Markov chains, by means of the strong recurrence result (1.28).

THEOREM 3. Let $\left\{X_{n}\right\}$ be defined by (1.36) for a chain $\tilde{X}$. with rates satisfying (1.31)-(1.35). Assume that (1.22)-(1.24) are satisfied. Then there exist an $M \in \mathbb{R}^{d}$ and a $d \times d$ matrix $\Sigma$ such that

$$
\begin{align*}
& \frac{\xi_{1}(n)-n M}{\sqrt{n}} \text { conditioned on }\{T>n\} \text { converges in }  \tag{1.40}\\
& \text { distribution to } N(\mathbf{0}, \Sigma)
\end{align*}
$$

More generally let $p$ be the period of (the restriction to $S_{0}$ of ) X. and let $S_{0,0}, \ldots, S_{0, p-1}$ be the cyclically moving subclasses in $S_{0}[c f$. (1.12)]. Then, for $0<s_{1}<s_{2}<\cdots<s_{l}, x \in S_{0, i}, y_{1} \in S_{0, j_{1}}, \ldots, y_{l} \in S_{0, j_{l}}$ and $\gamma_{1}, \ldots, \gamma_{l} \in$
$\mathbb{R}, w \in \mathbb{Z}^{d}$, as $n \rightarrow \infty$,

$$
\begin{gather*}
P\left\{\frac{\xi_{1}\left(\left\lfloor s_{r} n\right\rfloor p+j_{r}-i\right)-s_{r} n p M}{\sqrt{n p}} \leq \gamma_{r}, X_{\left\lfloor s_{r} n\right\rfloor p+j_{r}-i}=y_{r}, 1 \leq r \leq l \mid\right. \\
\left.X_{0}=x, \xi_{1}(0)=w, T>\left\lfloor s_{l} n\right\rfloor p+j_{l}-i\right\}  \tag{1.41}\\
\rightarrow P\left\{G_{s_{r}} \leq \gamma_{r}, 1 \leq r \leq l\right\} \frac{\mu\left(y_{l}\right)}{\sum_{v \in S_{0, j_{l}}} \mu(v)} \prod_{r=1}^{l-1}\left[f\left(y_{r}\right) \mu\left(y_{r}\right)\right],
\end{gather*}
$$

where $\left\{G_{t}\right\}_{t \geq 0}$ is a d-dimensional Gaussian process with mean zero, $G_{0}=\mathbf{0}$ and covariance function

$$
E\left\{G_{s} G_{t}^{*}\right\}=(s \wedge t) \Sigma
$$

( $G^{*}$ is the transpose of the column vector $G$ ). In particular this result is valid for the subcritical discrete time contact process (site or bond version) on $\mathbb{Z}^{d}$ with $M=\mathbf{0}, \Sigma=\sigma^{2} \times\left(\right.$ identitymatrixofdimensiond) for some $\sigma^{2}>0$.

If $\Sigma$ is the zero matrix, then (1.41) means that for all $\varepsilon>0$,

$$
\begin{aligned}
& P\left\{\frac{\left\|\xi_{1}\left(\left\lfloor s_{r} n\right\rfloor p+j_{r}-i\right)-s_{r} n p M\right\|}{\sqrt{n}} \leq \varepsilon, X_{\left\lfloor s_{r} n\right\rfloor p+j_{r}-i}=y_{r}, 1 \leq r \leq l \mid\right. \\
& \left.\quad X_{0}=x, \xi_{1}(0)=w, T>\left\lfloor s_{l} n\right\rfloor p+j_{l}-i\right\} \\
& \rightarrow \frac{\mu\left(y_{l}\right)}{\sum_{v \in S_{0, j_{l}}} \mu(v)} \prod_{r=1}^{l-1}\left[f\left(y_{r}\right) \mu\left(y_{r}\right)\right] .
\end{aligned}
$$

Remark 3. Note that the condition $X_{0}=x, \xi_{1}(0)=w$ fixes the full state of $\widetilde{X}_{0}$. If $X_{0}=x$ and $\xi_{1}(0)=w$, then $\widetilde{X}_{0}(z)=x(z-w)$.

Remark 4. Some more information on $M$ and $\Sigma$ is given in (5.12), (5.13) and succeeding lines.

Remark 5. A central limit theorem for oriented percolation (in a Fourier form) was proven before by Nguyen and Yang (1995). They used lace expansions, which only work in high dimension and seem to be restricted to oriented percolation. It would be difficult to extend their method to the more general class of processes described before Theorem 2. On the other hand, Nguyen and Yang also treat the critical oriented percolation in addition to the subcritical one. Our method does not apply at all to the critical case.

For the more difficult case of ordinary, unoriented percolation, a relatedbut slightly different-central limit result was first proved in Campanino, Chayes and Chayes (1991).

Remark 6. It is not difficult to extend the limit theorem (1.41) for the finite-dimensional distributions of $\left\{\xi_{1}(\lfloor s n\rfloor) / \sqrt{n}\right\}_{s \geq 0}$ conditioned on $\{T>n\}$ to a full invariance principle (of the $\xi_{1}$ coordinate only, though).
2. $\boldsymbol{R}$-positive-recurrence. In this section we shall prove Theorem 1 for general absorbing Markov chains. Equations (1.2), (1.4) and (1.22)-(1.24) will be in force throughout this section.

Here is a brief outline of the proof. Fix a state $y_{0}$. Condition (1.22) [in conjunction with (1.13)] says more or less that, conditionally on $T>n$, the chain cannot stay away from $\mathscr{U}_{1}$ too long. This will be used to show that in fact $X$. will visit $\mathscr{U}_{1}$ at a time near $n$, with a probability bounded away from 0 . By (1.24), $X$. will then also visit $\mathscr{U}_{2}$ and even $y_{0}$ at a time near $n$, with a probability bounded away from 0 . We will actually show that even the probability for $X_{k}=y_{0}$ and $X_{n}=y_{0}$ is on the average (over $k \in[0, n]$ ) not too small, still under the condition $\{T>n\}$. A little renewal theory then shows that $X$. must be $R$-positive-recurrent (see Lemma 5 ). As already stated, VereJones' results then show that there exist $f$ and $\mu$ satisfying (1.17)-(1.20). One can therefore introduce the honest transition probability matrix

$$
P^{*}(x, y)=\frac{R P(x, y) f(y)}{f(x)}, \quad x, y \in S_{0} .
$$

This is positive-recurrent. From the fact that $\left\{X_{n}=y_{0}\right\}$ has not too small a probability, we will obtain that $P_{y_{0}}\{T>n\}$ is of the same order as $P^{n}\left(y_{0}, y_{0}\right)$. Using $P^{*}$ and a Markov chain $X^{*}$ with this transition probability quickly leads to (1.21), (1.26) and (1.27).

To prove (1.28) we appeal to the hypothesis (1.23). Specifically, this is used to show that the function $f$ is bounded on $\mathscr{U}_{1}$. This will allow us to reduce (1.28) to fairly standard exponential bounds or large deviation estimates for the chain $X^{*}$ (these no longer involve $R$ ).

The above steps are separated into various lemmas. The $R$-positiverecurrence and (1.17)-(1.20) are proven in Lemma 5; the fact that the $\tilde{\mu}$ of (1.25) is a minimal quasi-stationary distribution is in Lemma 7. Finally, (1.21) and (1.26)-(1.28) are proven in Lemmas 6 and 8 . We carry out the proof only for the aperiodic case $(p=1)$; the extension to general $p$ is routine.

Throughout $D_{i}$ will stand for a strictly positive, finite constant whose value is of no significance to us; the same symbol $D_{i}$ may represent a different constant in different proofs, but we will not vary the value of $D_{i}$ in a single lemma. Also $C_{1}, C_{2}, \ldots$ will be constants in $(0, \infty)$, but they maintain the same value throughout this section.

We begin with a simple reduction, based merely on the irreducibility assumption (1.2), which is not strictly needed, but which simplifies the formulae.

Lemma 1. For any fixed $y_{0} \in S_{0}$ we may assume that $\mathscr{U}_{2}$ consists of $y_{0}$ only and that the $x_{0}$ of (1.23) equals $y_{0}$. (In particular we may assume that $\mathscr{U}_{1}$ contains $y_{0}$. )

Proof. Fix $y_{0} \in S_{0}$. By the irreducibility of $\widehat{P}$ and the finiteness of $\mathscr{U}_{2}$ there exists an $m_{0}<\infty$ and a $D_{1}>0$ such that for all $y \in \mathscr{U}_{2}$,

$$
\begin{equation*}
P^{k}\left(y, y_{0}\right) \geq D_{1} \quad \text { for some } 0 \leq k \leq m_{0} . \tag{2.1}
\end{equation*}
$$

Then, by (1.24), for $x \in \mathscr{U}_{1}$

$$
\begin{aligned}
& P_{x}\left\{X_{n}=y_{0} \text { for some } n \leq n_{0}+m_{0}\right\} \\
& \geq \sum_{j=0}^{n_{0}} \sum_{y \in \mathscr{U}_{2}} P_{x}\left\{X . \text { first visits } \mathscr{U}_{2} \text { at time } j, X_{j}=y\right\} \\
& \quad \times P_{y}\left\{X_{k}=y_{0} \text { for some } k \leq m_{0}\right\} \\
& \geq D_{1} \sum_{j=0}^{n_{0}} P_{x}\left\{X . \text { first visits } \mathscr{U}_{2} \text { at time } j\right\} \geq D_{1} C_{3} .
\end{aligned}
$$

Thus, by replacing $C_{3}$ by $D_{1} C_{3}$ and $n_{0}$ by $n_{0}+m_{0}$, we may assume (1.24) valid with $\mathscr{U}_{2}=\left\{y_{0}\right\}$.

Next we compare $P_{x_{1}}\{T>n\}$ and $P_{y_{1}}\{T>n\}$ for arbitrary $x_{1}, y_{1} \in S_{0}$. These two probabilities will be of the same order. Indeed, there exists a $k=$ $k\left(x_{1}, y_{1}\right)$ so that

$$
P^{k}\left(y_{1}, x_{1}\right)>0
$$

Since $X_{k}=x_{1} \in S_{0}$ implies $T>k$, we also have

$$
\begin{align*}
P_{y_{1}}\{T>n\} & \geq P^{k}\left(y_{1}, x_{1}\right) P_{x_{1}}\{T>n-k\}  \tag{2.2}\\
& \geq P^{k}\left(y_{1}, x_{1}\right) P_{x_{1}}\{T>n\} .
\end{align*}
$$

In particular, if $y_{0}$ is any fixed state, then the inequality in (1.23) will continue to hold (with a suitable change in $C_{2}$ ) if we replace $x_{0}$ by $y_{0}$. However, (1.23) in its original form required that $x_{0} \in \mathscr{U}_{1}$, and $y_{0}$ may not belong to $\mathscr{U}_{1}$. In the latter case we simply add $y_{0}$ to $\mathscr{U}_{1}$ and verify that (1.22) and (1.24) remain valid after this enlargement of $\mathscr{U}_{1}$. To verify (1.22) we note that the irreducibility of $\widehat{P}$ shows that if $y_{0}$ does not belong to $\mathscr{U}_{1}$, then there exists some $k>0$ and some $y \in \mathscr{U}_{1}$, such that

$$
P_{y}\left\{X_{k}=y_{0}, \quad X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq k\right\}>0 .
$$

Moreover,

$$
\begin{aligned}
& P_{y}\left\{T>n+k, \text { but } X_{l} \notin \mathscr{U}_{1} \text { for all } 1 \leq l \leq n+k\right\} \\
& \geq P_{y}\left\{X_{k}=y_{0}, X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq k\right\} \\
& \quad \times P_{y_{0}}\left\{T>n, \text { but } X_{l} \notin \mathscr{U}_{1} \text { for all } 1 \leq l \leq n\right\} .
\end{aligned}
$$

The first factor in the right-hand side is strictly positive, and the left-hand side is at most $C_{1}\left(R+\varepsilon_{0}\right)^{-n}$. Thus (1.22) also holds for $x=y_{0}$ (after a suitable change of $C_{1}$ ). Now adding $y_{0}$ to $\mathscr{U}_{1}$ can only help for the inequality in (1.22).

Equation (1.24) holds automatically, because if the inequality holds for $x \in$ $\mathscr{U}_{1}$ and for $x=y_{0}$, then it holds for $x \in \mathscr{U}_{1} \cup\left\{y_{0}\right\}$.

From now on we assume that $\mathscr{U}_{2}$ is a singleton, say $\mathscr{U}_{2}=\left\{y_{0}\right\}$, and that this $y_{0}$ belongs to $\mathscr{U}_{1}$ and that $x_{0}$ in (1.23) equals $y_{0}$.

Lemma 2. There exist a constant $C_{6}>0$ and, for all $x \in \mathscr{U}_{1}$, an $n_{1}=n_{1}(x)<$ $\infty$ such that

$$
\begin{equation*}
\sum_{k=0}^{n} P_{x}\left\{X_{k}=y_{0} \mid T>n\right\} \geq C_{6} n, \quad n \geq n_{1} \tag{2.3}
\end{equation*}
$$

Proof. For $n_{0}$ as in (1.24) and $k \leq n-n_{0}$ one has

$$
\begin{aligned}
& \sum_{l=0}^{n_{0}} P_{x}\left\{X_{k+l}=y_{0} \mid T>n\right\} \\
& \quad \geq \frac{1}{P_{x}\{T>n\}} \sum_{u \in \mathscr{U}_{1}} P^{k}(x, u) \sum_{l=0}^{n_{0}} P^{l}\left(u, y_{0}\right) P_{y_{0}}\{T>n-k\} \\
& \quad \geq C_{3} C_{2}^{-1} \sum_{u \in \mathscr{U}_{1}} P^{k}(x, u) \frac{P_{u}\{T>n-k\}}{P_{x}\{T>n\}} \quad[\text { by (1.23) and (1.24)] } \\
& \quad=C_{3} C_{2}^{-1} P_{x}\left\{X_{k} \in \mathscr{U}_{1} \mid T>n\right\} .
\end{aligned}
$$

It therefore suffices to prove that for some $D_{1}>0$,

$$
\begin{equation*}
\sum_{k=0}^{n} P_{x}\left\{X_{k} \in \mathscr{U}_{1} \mid T>n\right\} \geq D_{1} n, \quad x \in \mathscr{U}_{1}, n \geq n_{1} . \tag{2.4}
\end{equation*}
$$

Now, let $0 \leq \sigma_{0}<\sigma_{1}<\cdots$ be the successive times at which $X$. visits $\mathscr{U}_{1}$ and define

$$
\lambda(n)=\max \left\{j: \sigma_{j} \leq n\right\}
$$

We shall take $X_{0}=x \in \mathscr{U}_{1}$ so that $\sigma_{0}=0$ and $\lambda(n)$ will be well defined; $\lambda(n)+1$ will be the number of visits by $X$. to $\mathscr{U}_{1}$ during $[0, n]$. We define

$$
\mathscr{F}_{k}=\sigma \text {-field generated by } X_{0}, \ldots, X_{k}
$$

Then, conditionally on $\mathscr{F}_{\sigma_{j}}, \sigma_{j+1}-\sigma_{j}$ has a defective distribution, which satisfies on the event $\left\{\sigma_{j}<\infty, X_{\sigma_{j}}=y\right\}$ with $y \in \mathscr{U}_{1}$,

$$
\begin{align*}
P_{x}\{r & \left.<\sigma_{j+1}-\sigma_{j}<\infty \mid \mathscr{F}_{\sigma_{j}}\right\} \\
& =P_{y}\left\{r<\sigma_{1}<\infty\right\} \\
& =P_{y}\left\{X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq r, \text { but } X_{k} \text { returns to } \mathscr{U}_{1} \text { for some } k>r\right\}  \tag{2.5}\\
& \leq P_{y}\left\{T>r, \text { but } X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq r\right\} \leq C_{1}\left(R+\varepsilon_{0}\right)^{-r}
\end{align*}
$$

[by (1.22)]. Consequently, for any $\theta \geq 0$ with $e^{\theta}<R+\varepsilon_{0}$, on the event $\left\{\sigma_{j}<\right.$ $\infty\}$,

$$
\begin{aligned}
& \sum_{r=1}^{\infty} e^{\theta r} P_{x}\left\{\sigma_{j+1}-\sigma_{j}=r \mid \mathscr{\mathscr { F }}_{\sigma_{j}}\right\} \\
& \quad \leq e^{\theta}+\left(e^{\theta}-1\right) \sum_{r=1}^{\infty} e^{\theta r} P_{x}\left\{r<\sigma_{j+1}-\sigma_{j}<\infty \mid \mathscr{\mathscr { F }}_{\sigma_{j}}\right\} \\
& \quad \leq D_{2}(\theta)
\end{aligned}
$$

for some $1<D_{2}(\theta)<\infty$. It follows that

$$
\begin{align*}
& P_{x}\left\{j \text { th return to } \mathscr{U}_{1} \text { occurs at time } r\right\} \\
& \quad=P_{x}\left\{\sigma_{j}=r\right\} \\
& \quad \leq e^{-r \theta} E_{x}\left\{\exp \left(\sum_{i=0}^{j-1} \theta\left(\sigma_{i+1}-\sigma_{i}\right)\right) ; \sigma_{j}<\infty\right\}  \tag{2.6}\\
& \quad \leq e^{-r \theta}\left[D_{2}(\theta)\right]^{j} .
\end{align*}
$$

We use this to prove that for some $\alpha>0$ (independent of $x \in \mathscr{U}_{1}$ ) and $n_{1}$,

$$
\begin{equation*}
P_{x}\{\lambda(n) \leq \alpha n \mid T>n\} \leq \frac{1}{2}, \quad x \in \mathscr{U}_{1}, n \geq n_{1}, \tag{2.7}
\end{equation*}
$$

which in turn will imply

$$
\begin{aligned}
\sum_{k=0}^{n} & P_{x}\left\{X_{k} \in \mathscr{U}_{1} \mid T>n\right\} \\
& =E_{x}\left\{\text { number of visits by } X . \text { during }[0, n] \text { to } \mathscr{U}_{1} \mid T>n\right\} \\
& \geq E_{x}\{\lambda(n) \mid T>n\} \geq \alpha n P_{x}\{\lambda(n)>\alpha n \mid T>n\} \geq \frac{1}{2} \alpha n .
\end{aligned}
$$

Thus (2.4), and hence the lemma, will follow from (2.7). However, for $x \in \mathscr{U}_{1}$, by a standard last exit decomposition,

$$
\begin{align*}
& P_{x}\{\lambda(n) \leq \alpha n, T>n\} \\
& =\sum_{0 \leq j \leq \alpha n} \sum_{r=0}^{n} \sum_{y \in \mathscr{U}_{1}} P_{x}\left\{\lambda(n)=j, \sigma_{\lambda(n)}=r, X_{r}=y, T>n\right\} \\
& =\sum_{0 \leq j \leq \alpha n} \sum_{r=0}^{n} \sum_{y \in \mathscr{U}_{1}} P_{x}\left\{\sigma_{j}=r, X_{r}=y\right\} \\
& \quad \times P_{y}\left\{T>n-r, \text { but } X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq n-r\right\}  \tag{2.8}\\
& \leq \sum_{0 \leq j \leq \alpha n} \sum_{r=0}^{n} P_{x}\left\{\sigma_{j}=r\right\} C_{1}\left(R+\varepsilon_{0}\right)^{-n+r \quad[\text { by }(1.22)]} \\
& \leq C_{1} \sum_{0 \leq j \leq \alpha n} \sum_{r=0}^{n} e^{-r \theta}\left[D_{2}(\theta)\right]^{j}\left(R+\varepsilon_{0}\right)^{-n+r} \\
& \leq D_{3}(\theta)\left[D_{2}(\theta)\right]^{\alpha n}\left(R+\varepsilon_{0}\right)^{-n} \sum_{r=0}^{n}\left[\left(R+\varepsilon_{0}\right) e^{-\theta}\right]^{r} .
\end{align*}
$$

Now choose $\theta>0$ such that $e^{\theta}=R+\frac{7}{8} \varepsilon_{0}$, and then $\alpha>0$ so small that

$$
\left[D_{2}(\theta)\right]^{\alpha} e^{-\theta} \leq\left(R+\frac{3}{4} \varepsilon_{0}\right)^{-1}
$$

We then find that the left-hand side of (2.7) is at most

$$
\begin{equation*}
D_{4}(\theta)\left(R+\frac{3}{4} \varepsilon_{0}\right)^{-n}\left[P_{x}\{T>n\}\right]^{-1} \tag{2.9}
\end{equation*}
$$

On the other hand, for $x \in S_{0}$ there exists an $n_{1}<\infty$ such that

$$
\begin{equation*}
P_{x}\{T>n\} \geq \widehat{P}^{n}(x, x) \geq\left(R+\frac{\varepsilon_{0}}{2}\right)^{-n}, \quad n \geq n_{1} \tag{2.10}
\end{equation*}
$$

by virtue of (1.13) (recall that we took $p=1$ ). Equations (2.8)-(2.10) show that if $n_{1}$ is chosen large enough, also

$$
\begin{equation*}
P_{x}\{\lambda(n) \leq \alpha n \mid T>n\} \leq D_{4}(\theta)\left(\frac{R+\varepsilon_{0} / 2}{R+3 \varepsilon_{0} / 4}\right)^{n} \leq \frac{1}{2}, \quad n \geq n_{1} \tag{2.11}
\end{equation*}
$$

as desired.
The preceding lemma showed that there is a reasonable chance for the "average" $k$ that $X_{k}=y_{0}$. We now show that $\left\{X_{k} \in \mathscr{U}_{1}\right.$ for some $k$ "close to $n "\}$ has a high probability (all this conditioned on $\{T>n\}$ ).

Lemma 3. For every $\eta>0$ and $x \in \mathscr{U}_{1}$ there exist $m_{0}=m_{0}(x, \eta)$ and $n_{2}=n_{2}(x, \eta)<\infty$ such that

$$
\begin{equation*}
P_{x}\left\{X_{l} \notin \mathscr{U}_{1} \text { for all } n-m_{0} \leq l \leq n \mid T>n\right\} \leq \eta, \quad n \geq n_{2} \tag{2.12}
\end{equation*}
$$

Proof. Take $n_{1}(x)$ as in (2.10). Then for any $x \in \mathscr{U}_{1}, y \in S_{0}$ and $p \geq n_{1}(x)$

$$
\begin{align*}
P_{x}\{T>n\} & \geq P^{p}(x, x) P_{x}\{T>n-p\} \\
& \geq\left(R+\frac{\varepsilon_{0}}{2}\right)^{-p} P_{x}\{T>n-p\} \tag{2.13}
\end{align*}
$$

and for $m \geq n_{1}(x)$,
$P_{x}\left\{T>n\right.$, last visit to $\mathscr{U}_{1}$ by $X$. occurs before time $\left.n-m\right\}$

$$
\begin{aligned}
& =\sum_{r=0}^{n-m-1} P_{x}\left\{X_{r} \in \mathscr{U}_{1}, X_{l} \notin \mathscr{U}_{1} \text { for } r<l \leq n, T>n\right\} \\
& =\sum_{r=0}^{n-m-1} \sum_{y \in \mathscr{U}_{1}} P^{r}(x, y) P_{y}\left\{X_{l} \notin \mathscr{U}_{1} \text { for } 1 \leq l \leq n-r, T>n-r\right\} \\
& \leq \sum_{p=m+1}^{n} \sum_{y \in \mathscr{U}_{1}} P^{n-p}(x, y) C_{1}\left(R+\varepsilon_{0}\right)^{-p} \quad[\text { by (1.22)] } \\
& \leq C_{1} \sum_{p=m+1}^{n} P_{x}\{T>n-p\}\left(R+\varepsilon_{0}\right)^{-p}
\end{aligned}
$$

$$
\text { (because } X_{n-p}=y \in \mathscr{U}_{1} \text { implies } T>n-p \text { ) }
$$

$$
\leq D_{1}\left(\frac{R+\left(\varepsilon_{0} / 2\right)}{R+\varepsilon_{0}}\right)^{m} P_{x}\{T>n\} \quad[\text { by }(2.13)] .
$$

Thus, for any $\eta>0$, we can fix $m_{0}=m_{0}(x, \eta)$ such that for all large $n$, say $n \geq n_{2}(x, \eta)$, (2.12) holds.

We now combine Lemmas 2 and 3 to get our first real recurrence statement.
Lemma 4. There exist constants $C_{7}>0$ and $n_{3}<\infty$ such that for $n \geq n_{3}$,

$$
\begin{equation*}
\sum_{k=0}^{n} P_{y_{0}}\left\{X_{k}=y_{0} \text { and } X_{n}=y_{0} \mid T>n\right\} \geq C_{7} n . \tag{2.14}
\end{equation*}
$$

Proof. We start from (2.3) with $x=y_{0} \in \mathscr{U}_{1}$. We now take $\eta=C_{6} / 4$ and subtract from (2.3) the following consequence of (2.12):

$$
\begin{aligned}
\sum_{k=0}^{n} P_{y_{0}}\left\{X_{l} \notin \mathscr{U}_{1} \text { for all } n-m_{0}\right. & \leq l \leq n \mid T>n\} \\
\leq \eta(n+1) \leq \frac{1}{2} C_{6} n, \quad n & \geq n_{2}\left(y_{0}, \eta\right)+2
\end{aligned}
$$

$\left[m_{0}=m_{0}\left(y_{0}, \eta\right)\right]$. We obtain that for $n \geq n_{1}\left(y_{0}\right) \vee\left(n_{2}\left(y_{0}, \eta\right)+2\right)$,

$$
\begin{align*}
& \sum_{k=0}^{n} P_{y_{0}}\left\{X_{k}=y_{0} \text { and } X_{l} \in \mathscr{U}_{1} \text { for some } n-m_{0} \leq l \leq n \mid T>n\right\}  \tag{2.15}\\
& \quad \geq \frac{1}{2} C_{6} n .
\end{align*}
$$

The value of $X_{n}$ is still unspecified in (2.15) and we must now make $X_{n}$ itself equal to $y_{0}$. First we note that there exist an $r$ and a $D_{1}>0$ so that

$$
\begin{equation*}
P^{q}\left(y, y_{0}\right) \geq D_{1} \quad \text { for all } y \in \mathscr{U}_{1} \text { and } r \leq q \leq r+m_{0} \tag{2.16}
\end{equation*}
$$

This is so because, by (1.13), there exist an $s$ and a $D_{2}>0$ such that

$$
P^{q}\left(y_{0}, y_{0}\right) \geq D_{2} \quad \text { for all } q \text { with } s \leq q \leq s+n_{0}+m_{0}
$$

Thus for $y \in \mathscr{U}_{1}, s+n_{0} \leq q \leq s+n_{0}+m_{0}$,

$$
\begin{aligned}
P^{q}\left(y, y_{0}\right) & \geq \sum_{l=0}^{n_{0}} P_{y}\left\{\text { first visit to } y_{0} \text { occurs at time } l\right\} P^{q-l}\left(y_{0}, y_{0}\right) \\
& \geq D_{2} \sum_{l=0}^{n_{0}} P_{y}\left\{\text { first visit to } y_{0} \text { occurs at time } l\right\} \geq D_{2} C_{3}
\end{aligned}
$$

This is (2.16) for $r=s+n_{0}$ and $D_{1}=D_{2} C_{3}$. Finally, (2.15) and (2.16) show that

$$
\begin{align*}
& \sum_{k=0}^{n} P_{y_{0}}\left\{X_{k}=y_{0}, X_{n+r}=y_{0}\right\} \\
& \geq \sum_{k=0}^{n} \sum_{l=n-m_{0}}^{n} \sum_{y \in \mathscr{U}_{1}} P_{y_{0}}\left\{X_{k}=y_{0}, \text { first visit to } \mathscr{U}_{1}\right. \\
& \quad  \tag{2.17}\\
& \quad \text { at or after time } n-m_{0} \\
& \left.\quad \text { is at time } l \text { and } X_{l}=y\right\} P^{n+r-l}\left(y, y_{0}\right) \\
& \geq D_{1} \sum_{k=0}^{n} P_{y_{0}}\left\{X_{k}=y_{0}, X_{l} \in \mathscr{U}_{1} \text { for some } n-m_{0} \leq l \leq n\right\} \\
& \geq \frac{1}{2} D_{1} C_{6} n P_{y_{0}}\{T>n\} \geq \frac{1}{2} D_{1} C_{6} n P_{y_{0}}\{T>n+r\}
\end{align*}
$$

Short of replacing $n+r$ by $n$, this is (2.14) with $C_{7}=\frac{1}{4} D_{1} C_{6}$ and $n_{3}=$ $\left(n_{1} \vee n_{2} \vee r\right)+2$.

LEMMA 5. The transition probability matrix $\widehat{P}$ is $R$-positive-recurrent and there exist strictly positive functions $f, \mu: S_{0} \rightarrow(0, \infty)$ which satisfy (1.17)(1.20). Up to multiplicative constants, $f$ and $\mu$ are unique.

Proof. Define

$$
g_{n}=P_{y_{0}}\left\{X_{n}=y_{0}, \quad X_{r} \neq y_{0} \text { for } 1 \leq r \leq n-1\right\}
$$

and

$$
L=\sum_{n=1}^{\infty} g_{n} R^{n}
$$

Then, by Vere-Jones [(1967), Theorem C], $L \leq 1$, and $\widehat{P}$ is not $R$-recurrent is equivalent to

$$
\begin{equation*}
L<1 \tag{2.18}
\end{equation*}
$$

We begin by proving that (2.18) is not consistent with (2.14). To see this, note that (2.14) implies

$$
\begin{align*}
\sum_{k=0}^{n} \widehat{P}^{k}\left(y_{0}, y_{0}\right) \widehat{P}^{n-k}\left(y_{0}, y_{0}\right) & \geq C_{7} n P_{y_{0}}\{T>n\}  \tag{2.19}\\
& \geq C_{7} n \widehat{P}^{n}\left(y_{0}, y_{0}\right), \quad n \geq n_{3}
\end{align*}
$$

(again because $X_{n}=y_{0} \in S_{0}$ implies $T>n$ ). Next introduce i.i.d. random variables $Y, Y_{1}, Y_{2}, \ldots$ with the distribution

$$
\begin{align*}
P\{Y=n\} & =g_{n}^{*}:=g_{n} R^{n}, \quad 1 \leq n<\infty \\
P\{Y=\infty\} & =1-L \tag{2.20}
\end{align*}
$$

[This is really a distribution by the definition of $L$ and (2.18).] Under (2.18) the mass of this distribution on $\infty$ is strictly positive. Now set

$$
T_{l}=\sum_{j=1}^{l} Y_{j}
$$

and observe that

$$
\begin{aligned}
P_{y_{0}}\left\{X_{k}=y_{0} \mid X_{n}=y_{0}\right\} & =\frac{\widehat{P}^{k}\left(y_{0}, y_{0}\right) \widehat{P}^{n-k}\left(y_{0}, y_{0}\right)}{\widehat{P}^{n}\left(y_{0}, y_{0}\right)} \\
& =\frac{R^{k} \widehat{P}^{k}\left(y_{0}, y_{0}\right) R^{n-k} \widehat{P}^{n-k}\left(y_{0}, y_{0}\right)}{R^{n} \widehat{P}^{n}\left(y_{0}, y_{0}\right)} \\
& =\frac{P\left\{k \text { equals some } T_{l}\right\} P\left\{n-k \text { equals some } T_{l}\right\}}{P\left\{n \text { equals some } T_{l}\right\}} \\
& =P\left\{k \text { equals some } T_{l} \mid n \text { equals some } T_{l}\right\} .
\end{aligned}
$$

Therefore, (2.19) shows that for $n \geq n_{3}$,

$$
\begin{aligned}
& C_{7} n \leq \sum_{k=0}^{n} \frac{\widehat{P}^{k}\left(y_{0}, y_{0}\right) \widehat{P}^{n-k}\left(y_{0}, y_{0}\right)}{\widehat{P}^{n}\left(y_{0}, y_{0}\right)} \\
&=E\left\{\text { number of } T_{l} \text { in }[0, n] \mid n \text { equals some } T_{l}\right\} \\
& \leq \frac{1}{2} C_{7} n+(n+1) P\left\{\text { number of } T_{l} \text { in }[0, n] \text { exceeds } \left.\frac{1}{2} C_{7} n \right\rvert\,\right. \\
&\left.n \text { equals some } T_{l}\right\} .
\end{aligned}
$$

However, for any $\eta>0$,

$$
P\left\{n \text { equals some } T_{l}\right\}=R^{n} \widehat{P}^{n}\left(y_{0}, y_{0}\right) \geq(1-\eta)^{n}
$$

for large $n$, by virtue of (1.13). Therefore

$$
\begin{align*}
& P\left\{\text { number of } T_{l} \text { in }[0, n] \text { exceeds } \frac{1}{2} C_{7} n\right\} \\
& \quad \geq \frac{1}{4} C_{7} P\left\{n \text { equals some } T_{l}\right\} \geq \frac{1}{4} C_{7}(1-\eta)^{n} \tag{2.21}
\end{align*}
$$

for all large $n$. This, however, is impossible if

$$
L^{C_{7} / 2}<1-\eta,
$$

because

$$
\begin{aligned}
& P\left\{\text { number of } T_{l} \text { in }[0, n] \text { exceeds } \frac{1}{2} C_{7} n\right\} \\
& \quad \leq P\left\{Y_{j}<\infty \text { for } j \leq \frac{1}{2} C_{7} n+1\right\} \leq L^{C_{7} n / 2} .
\end{aligned}
$$

This is the required contradiction and we conclude that $L=1$ and that $\widehat{P}$ is $R$-recurrent.

Next we show that $\widehat{P}$ is $R$-positive-recurrent. Multiply (2.19) by $s^{n} R^{n}$ and sum over $n$. If we write

$$
U(s)=\sum_{n=0}^{\infty} \widehat{P}^{n}\left(y_{0}, y_{0}\right) s^{n} R^{n},
$$

we obtain

$$
\begin{equation*}
U^{2}(s) \geq C_{7} s U^{\prime}(s)-C_{7}\left(n_{3}+1\right)^{2} R^{n_{3}}, \quad 0 \leq s<1 \tag{2.22}
\end{equation*}
$$

[the $C_{7}\left(n_{3}+1\right)^{2} R^{n_{3}}$ is needed because (2.19) holds only for $n \geq n_{3}$ ]. Now the $R$-recurrence of $\widehat{P}$ simply says that

$$
U(1)=\sum_{0}^{\infty} R^{n} \widehat{P}^{n}\left(y_{0}, y_{0}\right)=\infty .
$$

Therefore, for $\frac{1}{2} \leq t<1$ and some $D_{1}<\infty$,

$$
\frac{1}{U(t)}=\frac{1}{U(t)}-\lim _{s \uparrow 1} \frac{1}{U(s)}=\int_{t}^{1} \frac{U^{\prime}(s)}{U^{2}(s)} d s \leq D_{1}(1-t) .
$$

Consequently,

$$
U(t) \geq \frac{1}{D_{1}(1-t)} \quad \text { as } t \uparrow 1
$$

and necessarily

$$
\limsup _{n \rightarrow \infty} R^{n} \widehat{P}^{n}\left(y_{0}, y_{0}\right)>0
$$

which gives $R$-positive-recurrence.
The existence of $f$ and $\mu$ which satisfy (1.17)-(1.20) and their uniqueness up to multiplicative constants is now guaranteed by Vere-Jones [(1967), Theorem 4.1 and Corollary on page 375].

Lemma 6. The $f$ and $\mu$ of Lemma 5 satisfy (1.21), (1.26) and (1.27).
Proof. Assume that (1.21) fails, so that

$$
\sum_{y \in S_{0}} \mu(y)=\infty .
$$

Then, by (1.20) and Fatou's lemma

$$
\liminf _{n \rightarrow \infty} R^{n} P_{y_{0}}\{T>n\}=\liminf _{n \rightarrow \infty} \sum_{y \in S_{0}} R^{n} P_{y_{0}}\left\{X_{n}=y\right\} \geq f\left(y_{0}\right) \sum_{y \in S_{0}} \mu(y)=\infty .
$$

This, however, is not posible, since (2.14) implies that for $n \geq n_{3}$,

$$
\begin{aligned}
C_{7} \frac{n}{n+1} & \leq \frac{1}{n+1} \sum_{k=0}^{n} P_{y_{0}}\left\{X_{k}=y_{0}, X_{n}=y_{0} \mid T>n\right\} \\
& \leq \frac{P^{n}\left(y_{0}, y_{0}\right)}{P_{y_{0}}(T>n)}=\frac{R^{n} P^{n}\left(y_{0}, y_{0}\right)}{R^{n} P_{y_{0}}\{T>n\}} \\
& \sim \frac{f\left(y_{0}\right) \mu\left(y_{0}\right)}{R^{n} P_{y_{0}}\{T>n\}} \quad[\text { by }(1.20)] .
\end{aligned}
$$

Thus (1.21) must hold.
Equation (1.26) now follows from the dominated convergence theorem. Indeed, by (1.18),

$$
\mu(x) \widehat{P}^{n}(x, y) \leq \mu \widehat{P}^{n}(y)=R^{-n} \mu(y)
$$

so that

$$
R^{n} \widehat{P}^{n}(x, y) \leq \frac{\mu(y)}{\mu(x)}
$$

Thus, by (1.20),

$$
\lim _{n \rightarrow \infty} R^{n} P_{x}\{T>n\}=\lim _{n \rightarrow \infty} \sum_{y \in S_{0}} R^{n} \widehat{P}^{n}(x, y)=f(x) \sum_{y \in S_{0}} \mu(y) .
$$

Finally, (1.27) is immediate from (1.26) and (1.20).
Lemma 7. The measure $\tilde{\mu}$, as defined in (1.25), is the unique minimal normalized quasi-stationary distribution.

Proof. Clearly $\widetilde{\mu}\left(S_{0}\right)=1$. Moreover, by (1.18),

$$
\tilde{\mu} \widehat{P}^{n}=R^{-n} \tilde{\mu},
$$

so that $\mu$ is a normalized quasi-stationary distribution.
To show the minimality of $\widetilde{\mu}$, observe that for any normalized quasistationary distribution $\nu$, (1.6) implies

$$
r(n+m)=r(n) r(m) \quad \text { and } \quad r(n)=[r(1)]^{n}
$$

and hence, by virtue of (1.7),

$$
\begin{aligned}
P_{\nu}\{T>n\} & =\left[P_{\nu}\{T>1\}\right]^{n}, \\
E_{\nu} T & =\sum_{n=0}^{\infty} P_{\nu}\{T>n\}=\left[1-P_{\nu}\{T>1\}\right]^{-1} .
\end{aligned}
$$

Thus, in the set of normalized quasi-stationary distributions the expected absorption times are ordered in the same way as $P_{\nu}\{T>1\}$, and a normalized quasi-stationary distribution $\tilde{\mu}$ is minimal if and only if $P_{\nu}\{T>1\}$ is minimized for $\nu=\widetilde{\mu}$. However, (1.6) implies that $r(1)=\widehat{P}_{\nu}\{T>1\}$ [cf. (1.7)], so that it suffices for the lemma to show that

$$
\begin{equation*}
r(1) \geq R^{-1} \quad \text { for any normalized quasi-stationary distribution } \nu . \tag{2.23}
\end{equation*}
$$

However, (2.23) follows from (1.13) and

$$
\widehat{P}^{n}(x, y) \leq[r(1)]^{n} \frac{\nu(y)}{\nu(x)}
$$

[cf. (1.6) and (1.7)].
Finally, there is only one minimal normalized quasi-stationary distribution because the solution of (1.18) is unique, up to a multiplicative constant [see Vere-Jones (1967), Theorem 4.1].

## LEMMA 8. Inequality (1.28) holds.

Proof. We saw in Lemma 1 that we can take $x_{0}$ equal to any state and that this may be the same state as $y_{0}$. We therefore only have to prove

$$
\begin{equation*}
g_{n}=P_{y_{0}}\left\{X_{n}=y_{0}, \text { but } X_{r} \neq y_{0} \text { for } 1 \leq r \leq n-1\right\} \leq D_{1}\left(R+\eta_{0}\right)^{-n} \tag{2.24}
\end{equation*}
$$

for some $D_{1}<\infty, \eta_{0}>0$, under the assumption that (1.23) holds with $x_{0}$ replaced by $y_{0}$. Now introduce a Markov chain $X_{n}^{*}$ with state space $S_{0}$ and (honest) transition probability matrix

$$
P^{*}(x, y)=\frac{R P(x, y) f(y)}{f(x)}, \quad x, y \in S_{0} .
$$

Analogously to Section 1, write $P_{x}^{*}$ for the distribution of $X^{*}$-paths conditioned on $X_{0}^{*}=x$. It is well known and easy to check that the $R$-positive-recurrence of $X$. implies that $X_{\text {. }}$ is positive-recurrent.

We shall first prove that for some $\alpha>0, \eta_{1}>0$ and $D_{2}<\infty$,

$$
\begin{equation*}
P_{y_{0}}^{*}\left\{X_{r}^{*} \in \mathscr{U}_{1} \text { for fewer than } \alpha n \text { values of } r \text { in }[0, n]\right\} \leq D_{2}\left(1-\eta_{1}\right)^{n} . \tag{2.25}
\end{equation*}
$$

Note the similarity of (2.25) to (2.11). The principal difference is that we now use the measure $P_{y_{0}}^{*}$, instead of the measure $P_{x}$, conditioned on $\{T>n\}$, of (2.11). Even though this will not be needed in the sequel, we point out that
there exists a general relationship between the measures $P_{x}$ conditioned on $\{T>n\}$ and $P_{x}^{*}$. The latter is a limit of the former in the following sense:

$$
\begin{aligned}
\lim _{n \rightarrow \infty} & P_{x}\left\{X_{1}=x_{1}, \ldots, X_{k}=x_{k} \mid T>n\right\} \\
& =\lim _{n \rightarrow \infty} \frac{P_{x_{k}}\{T>n-k\}}{P_{x}\{T>n\}} P_{x}\left\{X_{1}=x_{1}, \ldots, X_{k}=x_{k}\right\} \\
& =R^{k} \frac{f\left(x_{k}\right)}{f(x)} P_{x}\left\{X_{1}=x_{1}, \ldots, X_{k}=x_{k}\right\} \quad[\text { by (1.26)] } \\
& =R^{k} \frac{f\left(x_{k}\right)}{f\left(x_{k-1}\right)} \cdots \frac{f\left(x_{1}\right)}{f(x)} P_{x}\left\{X_{1}=x_{1}, \ldots, X_{k}=x\right\} \\
& =P_{x}^{*}\left\{X_{1}^{*}=x_{1}, \ldots, X_{k}^{*}=x_{k}\right\} .
\end{aligned}
$$

Returning to the proof of (2.25), we have

$$
\begin{align*}
& P_{y_{0}}^{*}\left\{X_{r}^{*} \in \mathscr{U}_{1} \text { for fewer than } \alpha n \text { values of } r \text { in }[0, n]\right\} \\
&  \tag{2.26}\\
& \quad=\sum_{m=n+1}^{\infty} \sum_{y \in \mathscr{U}_{1}} \frac{R^{m} f(y)}{f\left(y_{0}\right)} P_{y_{0}}\left\{\lfloor\alpha n+1\rfloor \text { th visit by } X . \text { to } \mathscr{U}_{1}\right.
\end{align*}
$$

$$
\text { occurs at time } \left.m \text { and } X_{m}=y\right\} .
$$

Moreover, by (1.26) and (1.23), we have for some $D_{3}<\infty$,

$$
\begin{equation*}
\frac{f(y)}{f\left(y_{0}\right)}=\lim _{n \rightarrow \infty} \frac{P_{y}\{T>n\}}{P_{y_{0}}\{T>n\}} \leq D_{3} \quad \text { for all } y \in \mathscr{U}_{1} \tag{2.27}
\end{equation*}
$$

Therefore, the left-hand side of (2.26) is (for large $n$ ) at most

$$
\begin{array}{r}
2 D_{3} \sum_{m=n+1}^{\infty} \sum_{y \in \mathscr{U}_{1}} R^{m} P_{y_{0}}\left\{\lfloor\alpha n+1\rfloor \text {-th visit by } X . \text { to } \mathscr{U}_{1} \text { occurs at time } m,\right. \\
\left.X_{m}=y, T>m\right\}
\end{array}
$$

(again because $X_{m}=y \in \mathscr{U}_{1}$ implies $T>m$ )
$\leq D_{4} \sum_{m=n+1}^{\infty} \sum_{y \in \mathscr{U}_{1}} P_{y_{0}}\left\{\lfloor\alpha n+1\rfloor\right.$ th visit by $X$. to $\mathscr{U}_{1}$ occurs at time $m$,

$$
\left.X_{m}=y \mid T>m\right\} \quad[\text { by (1.26) }]
$$

$\leq D_{4} \sum_{m=n+1}^{\infty} P_{y_{0}}\left\{X_{r} \in \mathscr{U}_{1}\right.$ for fewer than $(\alpha n+2)$ values of $r$ in $[0, m] \mid T>m\}$
$\leq D_{5} \sum_{m=n+1}^{\infty}\left(\frac{R+\varepsilon_{0} / 2}{R+3 \varepsilon_{0} / 4}\right)^{m} \quad[$ by $(2.11)]$.

This proves (2.25) with

$$
\eta_{1}=\frac{\varepsilon_{0}}{4 R+3 \varepsilon_{0}}
$$

The second ingredient for the proof is the following simple estimate, which holds uniformly for $y \in \mathscr{U}_{1}$ and $k \geq 0$ :

$$
\begin{align*}
P_{y_{0}}^{*} & \left\{X_{r}^{*}=y_{0} \text { for some } k \leq r \leq k+n_{0} \mid X_{k}^{*}=y\right\} \\
& \geq \frac{1}{n_{0}+1} E_{y_{0}}^{*}\left\{\text { number of } r \text { in }\left[k, k+n_{0}\right] \text { with } X_{r}^{*}=y_{0} \mid X_{k}^{*}=y\right\} \\
& =\frac{1}{n_{0}+1} \sum_{j=0}^{n_{0}}\left(P^{*}\right)^{j}\left(y, y_{0}\right) \\
& \geq \frac{f\left(y_{0}\right)}{n_{0}+1}\left[\sup _{y \in \mathscr{U}_{1}} f(y)\right]^{-1} \sum_{j=0}^{n_{0}} P^{j}\left(y, y_{0}\right)  \tag{2.28}\\
& \geq C_{3} \frac{f\left(y_{0}\right)}{n_{0}+1}\left[\sup _{y \in \mathscr{U}_{1}} f(y)\right]^{-1} \quad[\text { by }(1.24)] \\
& \geq D_{6}>0 \quad[\operatorname{by}(2.27)]
\end{align*}
$$

for some constant $D_{6}>0$. If $0 \leq \sigma_{0}^{*}<\sigma_{1}^{*}<\cdots$ are the times of the successive visits by $X^{*}$. to $\mathscr{U}_{1}$ and

$$
\mathscr{F}_{k}^{*}=\sigma \text {-field generated by } X_{0}^{*}, \ldots, X_{k}^{*}
$$

then (2.28) implies that

$$
\begin{aligned}
& P_{y_{0}}^{*}\left\{X_{r}^{*} \text { equals } y_{0} \text { for some } \sigma_{j}^{*} \leq r \leq \sigma_{j+n_{0}}^{*} \mid \mathscr{F}_{\sigma_{j}}^{*}\right\} \\
& \\
& \quad \geq P_{y_{0}}^{*}\left\{X_{r}^{*} \text { equals } y_{0} \text { for some } \sigma_{j}^{*} \leq r \leq \sigma_{j}^{*}+n_{0} \mid \mathscr{F}_{\sigma_{j}}^{*}\right\} \\
& \\
& \quad \geq D_{6}
\end{aligned}
$$

This in turn shows that

$$
\begin{equation*}
P_{y_{0}}^{*}\left\{X_{r}^{*} \neq y_{0} \text { for all } 1 \leq r \leq \sigma_{l}^{*}\right\} \leq\left(1-D_{6}\right)^{(l-1) /\left(n_{0}+1\right)} \tag{2.29}
\end{equation*}
$$

Finally note that

$$
\begin{aligned}
g_{n}^{*} & :=P_{y_{0}}^{*}\left\{X_{n}^{*}=y_{0} \text { but } X_{r}^{*} \neq y_{0} \text { for } 1 \leq r \leq n-1\right\} \\
& =\sum_{x_{1}, \ldots, x_{n-1} \neq y_{0}} P^{*}\left(y_{0}, x_{1}\right) P^{*}\left(x_{1}, x_{2}\right) \cdots P^{*}\left(x_{n-1}, y_{0}\right) \\
& =R^{n} \sum_{x_{1}, \ldots, x_{n-1} \neq y_{0}} P\left(y_{0}, x_{1}\right) \cdots P\left(x_{n-1}, y_{0}\right)=R^{n} g_{n},
\end{aligned}
$$

so that (2.24) is equivalent to

$$
\begin{equation*}
g_{n}^{*} \leq D_{1}\left(\frac{R}{R+\eta_{0}}\right)^{n} \tag{2.30}
\end{equation*}
$$

This follows for suitable $\eta_{0}>0, D_{1}<\infty$ from (2.25) and (2.29) by observing that

$$
\begin{aligned}
g_{n}^{*} \leq & P_{y_{0}}^{*}\left\{X_{r}^{*} \in \mathscr{U}_{1} \text { for fewer than } \alpha n \text { values of } r \text { in }[0, n]\right\} \\
& +P_{y_{0}}^{*}\left\{X_{r}^{*} \neq y_{0} \text { up until the } \alpha n \text {th visit of } X^{*} \text { to } \mathscr{U}_{1}\right\} \\
\leq & D_{2}\left(1-\eta_{1}\right)^{n}+\left(1-D_{6}\right)^{([\alpha n\rfloor-2) /\left(n_{0}+1\right)} .
\end{aligned}
$$

We have now proven all statements in Theorem 1.
3. Cellular automata. In this section we discuss Markov chains of the type discussed before Theorem 2 in the Introduction and prove Theorem 2. As we shall show, the main reason for the validity of Theorem 2 is that the number of occupied sites in its chains does have a very strong downward drift when this number of occupied sites becomes large. This is made precise by the estimate (3.2) below.

PROOF OF THEOREM 2. We shall occasionally write $\widetilde{X}(n, z)$ and $|\widetilde{X}(n)|$ for $\widetilde{X}_{n}(z)$ and $\left|\widetilde{X}_{n}\right|$, respectively, for typographical convenience. We remind the reader that the range $\rho$ is introduced in (1.35). Because $\widetilde{X}$. is strongly subcritical, there exists an $n_{4} \geq 1$ such that for all $x \in \widetilde{S}, z \in \mathbb{Z}^{d}$,

$$
\begin{equation*}
P_{x}\left\{\tilde{X}\left(n_{4}, z\right) \neq 0\right\} \leq \frac{1}{2}(2 \rho+1)^{-d} n_{4}^{-d} . \tag{3.1}
\end{equation*}
$$

We claim that this implies

$$
\begin{equation*}
E_{x}\left\{\left|\widetilde{X}\left(n_{4}\right)\right|\right\} \leq \frac{1}{2}|x| \text { for all } x \in \widetilde{S} . \tag{3.2}
\end{equation*}
$$

To see this, write

$$
\begin{equation*}
\left|\widetilde{X}\left(n_{4}\right)\right|=\sum_{z \in \mathbb{Z}^{d}} I\left[n_{4}, z\right] \tag{3.3}
\end{equation*}
$$

where $I[n, z]$ is the indicator function of the event $\{\tilde{X}(n, z) \neq 0\}$. Now the distribution of $\widetilde{X}(n, z)$ depends on $x$ only through the values of $x(v)$ with $v \in \mathbb{Z}^{d}$ such that $\|v-z\| \leq n \rho$ (where $\|\|$ is short for $\| \|_{2}$ ). For $n=1$ this is just assumption (1.35), and it easily follows for general $n$ by induction on $n$. Consequently, if $x(v)=0$ for all $v$ with $\|v-z\| \leq n \rho$, then

$$
\begin{equation*}
P_{x}\{\tilde{X}(n, z) \neq 0\}=P_{\underline{0}}\{\tilde{X}(n, z) \neq 0\}=0 \tag{3.4}
\end{equation*}
$$

because the state $\underline{0}$, with all components equal to 0 , is absorbing. It follows that if $v_{1}, \ldots, v_{|x|}$ are the sites in $\mathbb{Z}^{d}$ with $x(v) \neq 0$, then a.e. [ $P_{x}$ ],

$$
\begin{equation*}
\left|\tilde{X}\left(n_{4}\right)\right| \leq \sum_{j=1}^{|x|} \sum_{z \in \bar{B}\left(v_{j}, n_{4} \rho\right)} I\left[n_{4}, z\right] \tag{3.5}
\end{equation*}
$$

where

$$
\bar{B}(v, r)=\left\{z \in \mathbb{Z}^{d}:\|z-v\| \leq r\right\} .
$$

Taking expectation with respect to $P_{x}$ in (3.5) gives

$$
\begin{aligned}
E_{x}\left|\tilde{X}\left(n_{4}\right)\right| & \leq \sum_{j=1}^{|x|} \sum_{z \in \bar{B}\left(v_{j}, n_{4} \rho\right)} P_{x}\left\{\tilde{X}\left(n_{4}, z\right) \neq 0\right\} \\
& \leq|x|\left|\bar{B}\left(0, n_{4} \rho\right)\right| \frac{1}{2}(2 \rho+1)^{-d} n_{4}^{-d} \quad[\text { by }(3.1)] \\
& \leq \frac{1}{2}|x|
\end{aligned}
$$

This establishes (3.2).
Now for any integer $m \geq 1$ and $|x| \leq K$,

$$
\begin{align*}
P_{x}\{\mid & \left.\tilde{X}\left(j n_{4}\right) \mid>K \text { for } 1 \leq j \leq m, T>m n_{4}\right\} \\
& \leq \frac{1}{K^{r}} E_{x}\left\{\left|\widetilde{X}\left(m n_{4}\right)\right|^{r} ;\left|\tilde{X}\left(j n_{4}\right)\right|>K, 1 \leq j \leq m\right\} \\
& =\frac{|x|^{r}}{K^{r}} E_{x}\left\{\prod_{j=0}^{m-1} \frac{\left|\tilde{X}\left((j+1) n_{4}\right)\right|^{r}}{\left|\tilde{X}\left(j n_{4}\right)\right|^{r}} ;\left|\tilde{X}\left(j n_{4}\right)\right|>K, 1 \leq j \leq m\right\}  \tag{3.6}\\
& \leq \sup _{|x| \leq K} E_{x}\left(\frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{|x|^{r}}\right)\left[\sup _{|y|>K} E_{y}\left(\frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{|y|^{r}}\right)\right]^{m-1}
\end{align*}
$$

By (3.5),

$$
\left|\widetilde{X}\left(n_{4}\right)\right| \leq|\widetilde{X}(0)|(2 \rho+1)^{d} n_{4}^{d}
$$

so that

$$
\begin{equation*}
\sup _{|x| \leq K} E_{x}\left(\frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{|x|^{r}}\right) \leq(2 \rho+1)^{d r} n_{4}^{d r} \tag{3.7}
\end{equation*}
$$

We next fix $r$ such that

$$
\begin{equation*}
\left(\frac{2}{3}\right)^{r / n_{4}}<\gamma \tag{3.8}
\end{equation*}
$$

[where $\gamma$ is the number appearing in the right-hand side of (1.39)]. We shall show that for all large $K$,

$$
\begin{equation*}
\sup _{|y|>K} E_{y}\left(\frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{|y|^{r}}\right) \leq\left(\frac{2}{3}\right)^{r} \tag{3.9}
\end{equation*}
$$

Now observe that the absorption time $T$ is the same for the two chains $\left\{X_{n}\right\}$ and $\left\{\widetilde{X}_{n}\right\}$. Also $\left|X_{n}\right|=\left|\widetilde{X}_{n}\right|$. Therefore, if (3.9) holds,

$$
\begin{aligned}
& \sup _{x \in S,|x| \leq K} P_{x}\left\{T>n, \text { but }\left|X_{l}\right|>K \text { for all } 1 \leq l \leq n\right\} \\
& \quad=\sup _{x \in \tilde{S},|x| \leq K} P_{x}\left\{T>n, \text { but }\left|\tilde{X}_{l}\right|>K \text { for all } 1 \leq l \leq n\right\} \\
& \quad \leq \sup _{x \in \tilde{S},|x| \leq K} P_{x}\left\{\left|\tilde{X}\left(j n_{4}\right)\right|>K \text { for all } 1 \leq j \leq\left\lfloor n / n_{4}\right\rfloor\right\} \\
& \\
& \quad \leq(2 \rho+1)^{d r} n_{4}^{d r}\left(\frac{2}{3}\right)^{r\left(\left\lfloor n / n_{4}\right\rfloor-1\right)} \quad[\text { by (3.6), (3.7) and (3.9)] } \\
& \\
& \leq C_{5} \gamma^{n} \quad[\text { by }(3.8)]
\end{aligned}
$$

with

$$
C_{5}=\left(\frac{3}{2}\right)^{2 r}(2 \rho+1)^{d r} n_{4}^{d r}
$$

(which is independent of $n, K$ ). Thus (3.9) will imply (1.39) and we now turn to the proof of (3.9).

By (3.3),

$$
\begin{align*}
E_{x}\left|\tilde{X}\left(n_{4}\right)\right|^{r}= & \sum_{z_{1}, \ldots, z_{r}} E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\} \\
= & \sum_{z_{1}, \ldots, z_{r}} \prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]  \tag{3.10}\\
& +\sum_{z_{1}, \ldots, z_{r}}\left\{E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\}-\prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]\right\}
\end{align*}
$$

Also, by (3.2),

$$
\sum_{z_{1}, \ldots, z_{r}} \prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]=\left\{E_{x} \sum_{z} I\left[n_{4}, z\right]\right\}^{r} \leq\left(\frac{1}{2}|x|\right)^{r} .
$$

Moreover,

$$
\left|E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\}-\prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]\right| \leq 1,
$$

so that (3.10) shows

$$
\begin{align*}
E_{x} \frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{|x|^{r}} \leq 2^{-r}+|x|^{-r} & \left(\text { number of } r \text {-tuples } z_{1}, \ldots, z_{r}\right. \text { with }  \tag{3.11}\\
& \left.E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\} \neq \prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]\right) .
\end{align*}
$$

We claim that for any subsets $E_{1}, \ldots, E_{r}$ of $\mathbb{Z}^{d}$, which satisfy

$$
\begin{equation*}
d\left(E_{i}, E_{j}\right)=\inf \left\{\left\|z^{\prime}-z^{\prime \prime}\right\|: z^{\prime} \in E_{i}, z^{\prime \prime} \in E_{j}\right\}>2 n \rho \quad \text { for } \quad i \neq j \tag{3.12}
\end{equation*}
$$

the families $\left\{\tilde{X}(n, z): z \in E_{i}\right\}$ are independent under $P_{x}$. In fact, for $n=1$ all $\widetilde{X}(1, z), \quad z \in \mathbb{Z}^{d}$, are independent under $P_{x}$ by (1.34). For general $n$ our claim follows by induction, by means of (1.34) and (1.35). Indeed if the result is true for $n=1, \ldots, k$ and $E_{1}, \ldots, E_{r}$ satisfy (3.12) with $n=k+1$, then for any bounded functions $f_{i}$, depending on $\left\{\widetilde{X}(k+1, z): z \in E_{i}\right\}, 1 \leq i \leq r$, the Markov property and (1.34) imply that

$$
\begin{equation*}
E_{x}\left\{\prod_{1}^{n} f_{i}\right\}=E_{x}\left\{E_{x}\left\{\prod_{1}^{r} f_{i} \mid \tilde{X}_{k}\right\}\right\}=E_{x}\left\{\prod_{1}^{r} E_{x}\left\{f_{i} \mid \tilde{X}_{k}\right\}\right\} \tag{3.13}
\end{equation*}
$$

(note that the $E_{i}$ are disjoint). Furthermore, by (1.35), $E_{x}\left\{f_{i} \mid \widetilde{X}_{k}\right\}=g_{i}$ for some bounded function $g_{i}$ of the $\widetilde{X}_{k}(v)$ with $v \in E_{i}^{\prime}$, where

$$
E_{i}^{\prime}=\left\{v \in \mathbb{Z}^{d}:\|v-z\| \leq \rho \text { for some } z \in E_{i}\right\}
$$

The $E_{i}^{\prime}$ satisfy (3.12) with $n=k$, so that by the induction hypothesis the $g_{i}$ are independent under $P_{x}$. Thus

$$
E_{x}\left\{\prod_{1}^{r} f_{i}\right\}=E_{x}\left\{\prod_{1}^{r} g_{i}\right\}=\prod_{1}^{r} E_{x}\left\{g_{i}\right\}=\prod_{1}^{r} E_{x}\left\{E_{x}\left\{f_{i} \mid \tilde{X}_{k}\right\}\right\}=\prod_{1}^{r} E_{x}\left\{f_{i}\right\}
$$

Thus the families $\left\{\tilde{X}_{k+1}(z): z \in E_{i}\right\}$ are independent, as claimed.
The above independence property shows that

$$
E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\}=\prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]
$$

as soon as $\left\|z_{i}-z_{j}\right\|>2 n_{4} \rho$ for $i \neq j$. Furthermore, by (3.4), $I\left[n_{4}, z_{i}\right]=0$ a.e. [ $P_{x}$ ] unless $z_{i}$ lies within distance $n_{4} \rho$ from the set

$$
\begin{equation*}
\left\{v \in \mathbb{Z}^{d}: x(v) \neq 0\right\} \tag{3.14}
\end{equation*}
$$

Thus

$$
E_{x}\left\{I\left[n_{4}, z_{1}\right] \cdots I\left[n_{4}, z_{r}\right]\right\}=\prod_{i=1}^{r} E_{x} I\left[n_{4}, z_{i}\right]=0
$$

unless each $z_{i}$ lies within distance $n_{4} \rho$ from the set (3.14). Since the set in (3.14) has cardinality $|x|$, the number of $r$-tuples $z_{1}, \ldots, z_{r}$ which lie entirely within distance $n_{4} \rho$ of the set (3.14) and have $\left\|z_{i}-z_{j}\right\| \leq 2 n_{4} \rho$ for some $i \neq j$ is at most

$$
\left[\left(2 n_{4} \rho+1\right)^{d}|x|\right]^{r-1} r^{2}\left(4 n_{4} \rho+1\right)^{d}
$$

This is therefore an upper bound for the number appearing in the right-hand side of (3.11) and, consequently,

$$
E_{x} \frac{\left|\tilde{X}\left(n_{4}\right)\right|^{r}}{\left|x^{r}\right|} \leq 2^{-r}+\frac{1}{|x|} r^{2}\left(4 n_{4} \rho+1\right)^{d r}
$$

Equation (3.9) is now immediate for

$$
|K|>\frac{r^{2}\left(4 n_{4} \rho+1\right)^{d r}}{(2 / 3)^{r}-(1 / 2)^{r}} .
$$

4. Some examples. We shall now use Theorem 2 to give some specific examples to which Theorem 1 applies. These examples therefore have normalized quasi-stationary distributions and the limit relations (1.26) and (1.27) hold for them. All examples are of the form described before Theorem 2, and we allow only the values 0 or 1 for $\widetilde{X}_{n}(z)$ [i.e., $\kappa=2$ in (1.30)]. By the translation invariance (1.33) and the independence property (1.34) the model is completely specified by the probabilities

$$
\begin{equation*}
P\left\{\tilde{X}_{1}(\mathbf{0})=1 \mid \tilde{X}_{0}=x\right\} \tag{4.1}
\end{equation*}
$$

(where $\mathbf{0}$ denotes the origin in $\mathbb{Z}^{d}$ ), and this probability depends only on $x(z)$ for $\|z\| \leq \rho\left[\right.$ by (1.35)]. Note that we have an obvious partial order on $\widetilde{S}_{d, 2}: x^{\prime} \geq$ $x^{\prime \prime}$ if and only if $x^{\prime}(z) \geq x^{\prime \prime}(z)$ for all $z \in \mathbb{Z}^{d}$. We now impose four further conditions on the probabilities in (4.1):
$P\left\{\tilde{X}_{1}(\mathbf{0})=1 \mid \tilde{X}_{0}=x\right\}=0$ if $x(z)=0$ for all $z \in \bar{B}(\mathbf{0}, \rho)$, but
$P\left\{\tilde{X}_{1}\left(z_{0}\right)=1 \mid \widetilde{X}_{0}=x\right\}>0$ for some $z_{0}$ if $x(z)=1$ for some $z \in \bar{B}(\mathbf{0}, \rho) ;$
for any disjoint subsets $A_{1}, \ldots, A_{s}$ of $\bar{B}(\mathbf{0}, \rho)=\{v:\|v\| \leq \rho\}$, it holds that $P\left\{\widetilde{X}_{1}(\mathbf{0})=1 \mid \tilde{X}_{0}(z)=1\right.$ in $\bar{B}(\mathbf{0}, \rho)$ exactly for the $z ' s$ in $\left.\bigcup_{1}^{s} A_{i}\right\} \leq 1-\prod_{i=1}^{s} P\left\{\tilde{X}_{1}(\mathbf{0})=0 \mid \widetilde{X}_{0}(z)=1\right.$ in $\bar{B}(\mathbf{0}, \rho)$ exactly for the $z$ 's in $\left.A_{i}\right\}$.

Condition (4.5) says that the chain starting with only the sites in $\bigcup_{1}^{s} A_{i}$ occupied lies stochastically below the maximum of $s$ independent chains, the $i$ th of which starts with only the sites in $A_{i}$ occupied, $1 \leq i \leq s$. In order to obtain the required irreducibility (1.2) or (1.32), we shall restrict the state space $S$ to

$$
\begin{equation*}
\underline{0} \cup \text { collection of states which can be reached by } X \text {. when } \tag{4.6}
\end{equation*}
$$ starting at $X_{0}=\delta(\mathbf{0})$

$[\delta(\mathbf{0})$ is the state with only the component at the origin equal to 1 and all others equal to 0 ]. We shall prove the following result.

Proposition 1. If (1.31)-(1.35) and (4.2)-(4.5) hold and if $\left\{\tilde{X}_{n}\right\}$ is strongly subcritical, then $\left\{X_{n}\right\}$ [as defined by (1.36) and restricted to the space $S$ of (4.6)] satisfies all hypotheses of Theorem 1. It therefore has a minimal normalized quasi-stationary distribution $\tilde{\mu}$, and the limit relations (1.26) and (1.27) hold as well as the bound (1.28).

Before proving the proposition, we give some very specific examples.
EXAMPLE A (One-dimensional nearest neighbor chains). These chains have $d=1, \rho=1$, and $\widetilde{X}_{n}(z)$ takes on only the values 1 and 0 . Thus, given $\widetilde{X}_{n}$, the distribution of $\widetilde{X}_{n+1}(z)$ depends only on $\widetilde{X}_{n}(z-1), \widetilde{X}_{n}(z), \widetilde{X}_{n}(z+1)$ (remember that $z \in \mathbb{Z}$ now). In fact we shall assume that it depends only on the number of 1's among $\widetilde{X}_{n}(z-1), \widetilde{X}_{n}(z), \widetilde{X}_{n}(z+1)$. The distribution of $\widetilde{X}_{n+1}$ is therefore completely specified by the four probabilities

$$
\pi_{i}:=P\left\{\tilde{X}_{n+1}(z)=1 \mid \tilde{X}_{n}\right\} \text { on the event }
$$

\{exactly $i$ of the values $\tilde{X}_{n}(z-1), \widetilde{X}_{n}(z), \tilde{X}_{n}(z+1)$ equal 1$\}$,

$$
0 \leq i \leq 3 .
$$

In fact, to make $\underline{0}$ absorbing we must take

$$
\begin{equation*}
\pi_{0}=0 \tag{4.7}
\end{equation*}
$$

Equations (4.2)-(4.5) will then hold if we take, in addition to (4.7),

$$
\begin{align*}
0 & <\pi_{1} \leq \pi_{2} \leq \pi_{3}<1  \tag{4.8}\\
\pi_{2} & \leq 1-\left(1-\pi_{1}\right)^{2}=2 \pi_{1}-\pi_{1}^{2} \tag{4.9}
\end{align*}
$$

and

$$
\begin{equation*}
\pi_{3} \leq 1-\left(1-\pi_{1}\right)\left(1-\pi_{2}\right)=\pi_{1}+\pi_{2}-\pi_{1} \pi_{2} \tag{4.10}
\end{equation*}
$$

Indeed (4.2)-(4.4) hold by virtue of (4.8). As for (4.5), $\bar{B}(0,1)=\{-1,0,+1\}$ and the only possible disjoint choices for $A_{1}, \ldots, A_{s}$ are

$$
\begin{equation*}
s=2, \quad A_{1}=\{-1\}, \quad A_{2}=\{0\}, \tag{4.11}
\end{equation*}
$$

or

$$
\begin{equation*}
s=2, \quad A_{1}=\{-1\}, \quad A_{2}=\{0,1\}, \tag{4.12}
\end{equation*}
$$

or

$$
\begin{equation*}
s=3, \quad A_{1}=\{-1\}, \quad A_{2}=\{0\}, \quad A_{3}=\{1\}, \tag{4.13}
\end{equation*}
$$

or equivalent choices obtained from (4.11)-(4.13) by permuting $\{-1,0,1\}$. Condition (4.5) is void for $s=1$. For the situations in (4.11) and (4.12), (4.5) is guaranteed by (4.9) and (4.10), respectively. For the case in (4.13), (4.5) requires

$$
\pi_{3} \leq 1-\left(1-\pi_{1}\right)^{3},
$$

but this is easily seen to be a consequence of (4.9) and (4.10).

We do not have a sharp criterion for this chain to be strongly subcritical. However, this chain lies stochastically below oriented site percolation with $d=1, p=\pi_{3}$, discussed in the next example. Thus if $\pi_{3}$ is sufficiently small, then $\left\{\widetilde{X}_{n}\right\}$ is strongly subcritical.

Example B (The discrete time contact process or oriented percolation). We describe here a discrete time analogue of the contact process. As will be apparent from the description below, this model can also be viewed as percolation on the vertex set $\mathbb{Z}^{d} \times\{1,2, \ldots\}$ with an oriented edge from $y \times\{n\}$ to $z \times\{n+1\}$ if and only if $y$ and $z$ are adjacent on $\mathbb{Z}^{d}$. The usual oriented percolation is essentially the same model but with $\mathbb{Z}^{d}$ replaced by another lattice; we shall not consider that, though. Again we take a nearest neighbor process on $\mathbb{Z}^{d}$, that is, with $\rho=1$. The probability in (4.1) depends only on the number of 1's among the $x(v)$, with $v$ in $\bar{B}(\mathbf{0}, 1) \backslash\{\mathbf{0}\}=\left\{v: v\right.$ adjacent to $\mathbf{0}$ on $\left.\mathbb{Z}^{d}\right\}$. We do allow $d>1$ now, though. There are two versions of this model.
(i) Site version. This version corresponds to the threshold contact process. We take $P\left\{\widetilde{X}_{1}(\mathbf{0})=1 \mid \widetilde{X}_{0}=x\right\}=p$ if at least one $v$ adjacent on $\mathbb{Z}^{d}$ to $\mathbf{0}$ has $x(v)=1$, and $P\left\{\tilde{X}_{1}(\mathbf{0})=1 \mid \tilde{X}_{0}=x\right\}=0$ otherwise. Equation (4.2) is obvious and so are (4.3) and (4.4) if $0<p<1$. We shall not check (4.5) formally, but this is almost obvious from the following percolation construction. Color all sites of $\mathbb{Z}^{d} \times\{1,2, \ldots\}$ white (black) with probability $p$ ( $1-p$, respectively), independently of each other. When $\widetilde{X}_{0}=x$ is given, define $\widetilde{X}_{n}$ for $n \geq 1$ recursively by the following rule. If the site at $z \times\{n\}$ is black, then $\widetilde{X}_{n}(z)=0$ [i.e., $z$ is vacant at time $n$ ]. If $z \times\{n\}$ is colored white and at least one of its neighbors is occupied at time $n-1$ [i.e., $\widetilde{X}_{n-1}(v)=1$ for some neighbor $v$ on $\mathbb{Z}^{d}$ of $z$ ], then $z$ is also occupied at time $n$. If $z$ has no occupied neighbor at time $n-1$, then $z$ is not occupied at time $n$, irrespective of the color of $z \times\{n\}$.

It is known [see for instance Durrett (1988), Section 5a] that this chain has a critical probability $p_{c}=p_{c}(d$, site $) \in(0,1)$ such that for $p \leq p_{c}$ the process becomes extinct, that is, (1.4) holds for any starting configuration with finitely many occupied sites. Moreover, if we write $\delta\left(z_{0}\right)$ for the state $x$ with $x(z)=1$ if and only if $z=z_{0}$, then

$$
\begin{equation*}
P_{\delta\left(z_{0}\right)}\{T>n\} \rightarrow 0 \text { exponentially in } n \tag{4.14}
\end{equation*}
$$

whenever $p<p_{c}$ [the probability in (4.14) is independent of $z_{0}$ ]. This can be seen by redoing Menshikov's proof for the usual unoriented percolation [as given in Grimmett (1989), Section 3.2] or by using Aizenman and Barsky (1987) together with Hammersley's theorem [Theorem 5.1 in Grimmett (1989); the latter's proof in the oriented case is quite easy]. For a proof of (4.14) in a more complicated situation, see also Bezuidenhout and Grimmett (1991).

Equation (4.14) implies that

$$
\begin{equation*}
\text { for } p<p_{c},\left\{\tilde{X}_{n}\right\} \text { is strongly subcritical. } \tag{4.15}
\end{equation*}
$$

To see this note that as in the lines following (3.3), $\widetilde{X}_{n}(\mathbf{0})$ depends only on $x(v)$ with $\|v\| \leq n$. Furthermore, by the above description, $\widetilde{X}_{n}(\mathbf{0})=1$ occurs
if and only if there is a path $v_{0}=\mathbf{0}, v_{1}, v_{2}, \ldots, v_{n}$ on $\mathbb{Z}^{d}$ such that $v_{k} \times\{n-k\}$ is colored white for $0 \leq k \leq n-1$ and $x\left(v_{n}\right)=1$. Therefore,

$$
\begin{align*}
P_{x}\left\{\tilde{X}_{n}(\mathbf{0})=1\right\} & \leq \sum_{\|v\| \leq n} P_{\delta(v)}\left\{\tilde{X}_{n}(\mathbf{0})=1\right\}  \tag{4.16}\\
& \leq(2 n+1)^{d} P_{\delta(\mathbf{0})}\{T>n\}
\end{align*}
$$

and this tends to 0 exponentially fast in $n$. Thus for $p<p_{c}$ all conclusions of Theorem 1 hold for this chain.
(ii) Bond version. We now take

$$
P\left\{\tilde{X}_{1}(z)=1 \mid \tilde{X}_{0}=x\right\}=1-(1-p)^{j}
$$

if $z$ has exactly $j$ occupied neighbors in state $x$. In terms of independent colorings we can describe this model as follows: put an edge or bond between $v \times\{n\} \in \mathbb{Z}^{d} \times\{0,1, \ldots\}$ and $z \times\{n+1\}$ if and only if $v$ and $z$ are adjacent. Color all edges independently white (or black) with probability $p$ (or $1-p$, respectively). Then $\widetilde{X}_{n+1}(z)=1$ if and only if it has a white edge to an occupied neighbor at time $n$, that is, if $\tilde{X}_{n}(v)=1$ for some $v$ adjacent to $z$ for which the edge between $v \times\{n\}$ and $z \times\{n+1\}$ is white. Again this chain has a critical probability $p_{c}=p_{c}(d$, bond $) \in(0,1)$ such that for $p<p_{c}$ all conclusions of Theorem 1 hold. We skip the details.

REMARK 7. A closely related problem to (1.27) for one-dimensional oriented subcritical percolation is the following question: Is there a limit distribution for the configuration of occupied sites at time $n$, as seen from the leftmost occupied site, when in the initial state all sites in $\{0,1,2, \ldots\}$ are occupied? In this case it is not necessary to condition on $\{T>n\}$ because $T=\infty$ almost surely for any initial state with infinitely many occupied sites. This question was investigated by A. Galves, M. Keane and I. Meilijson (private communication). Our results do not apply directly to this problem, but we hope that they will nevertheless be useful.

REMARK 8. For the discrete time contact process it is almost trivial to verify (1.23) if we take $\mathscr{U}_{1}=\{x:|x| \leq K\}$ for some $K$ and $x_{0}=\delta(\mathbf{0})$. Indeed, for any $x \in \mathscr{U}_{1}$,

$$
\begin{align*}
& P_{x}\left\{X_{n} \text { is not absorbed by time } n\right\} \\
& \quad=P_{x}\{\text { one of the }|x| \text { initial particles survives until time } n\}  \tag{4.17}\\
& \quad \leq|x| P_{\delta(\mathbf{0})}\{T>n\} \leq K P_{\delta(\mathbf{0})}\{T>n\}
\end{align*}
$$

The next lemma shows how to generalize this estimate under (4.2)-(4.5).
LEMMA 9. Let $\left\{\left\{\tilde{X}_{n}^{v}\right\}_{n \geq 0}: v \in \mathbb{Z}^{d}\right\}$ be a family of independent Markov chains, each of which has the same transition probabilities as $\left\{\widetilde{X}_{n}\right\}$, but $\widetilde{X}^{v}$. has initial state $\widetilde{X}_{0}^{v}=\delta(v)$. If (1.31)-(1.35), (4.2) and (4.5) hold, then $\left\{\widetilde{X}_{n}\right\}$,
starting from $\tilde{X}_{0}=x$, is stochastically smaller than the process $\left\{Y_{n}\right\}$ defined by

$$
\begin{equation*}
Y_{n}(z)=\max \left\{\tilde{X}_{n}^{v}(z): x(v)=1\right\} . \tag{4.18}
\end{equation*}
$$

In particular, for any $A \subset \mathbb{Z}^{d}$,

$$
\begin{align*}
P_{x}\left\{\widetilde{X}_{n}(z)=0 \text { for all } z \in A\right\} & \geq P\left\{Y_{n}(z)=0 \text { for all } z \in A\right\} \\
& =\prod_{\substack{v \text { with } \\
x(v)=1}} P\left\{\widetilde{X}_{n}^{v}(z)=0 \text { for all } z \in A\right\} . \tag{4.19}
\end{align*}
$$

Proof. Equation (4.19) is proven by induction on $n$, together with a coupling of $\widetilde{X}$. and all the $\widetilde{X}^{v}$. For $n=1$, (4.19) is immediate from (4.5) with $A_{i}=\left\{v_{i}\right\}$, if $v_{1}, \ldots, v_{s}$ are the sites for which $x(v)=1$. Indeed, by (1.34),

$$
P_{x}\left\{\tilde{X}_{1}(z)=0 \text { for all } z \in A\right\}=\prod_{z \in A} P_{x}\left\{\tilde{X}_{1}(z)=0\right\} .
$$

Moreover, by (1.35), $P_{x}\left\{\tilde{X}_{1}(z)=0\right\}$ depends only on which $v$ 's in $\bar{B}(z, 1)$ have $x(v)=1$ and, by (4.5),

$$
\begin{align*}
P_{x}\left\{\tilde{X}_{1}(z)=0\right\} & \geq \prod_{v_{i} \in \bar{B}(z, 1)} P_{\delta\left(v_{i}\right)}\left\{\tilde{X}_{1}(z)=0\right\} \\
& \geq \prod_{i=1}^{s} P_{\delta\left(v_{i}\right)}\left\{\widetilde{X}_{1}(z)=0\right\}=P\left\{Y_{1}(z)=0\right\} . \tag{4.20}
\end{align*}
$$

Now the $Y_{1}(z), z \in \mathbb{Z}^{d}$, are also independent when $\widetilde{X}_{0}=x$ is given [by virtue of (1.34)]. Thus (4.20) gives (4.19) for $n=1$. It is now easy to couple all the $\widetilde{X}_{1}^{v}, v \in \mathbb{Z}^{d}$, and $\widetilde{X}_{1}$. One merely chooses the $\widetilde{X}_{1}^{v}, v \in \mathbb{Z}^{d}$, independently of each other, with their prescribed distribution. These $\widetilde{X}_{1}^{v}$ determine $Y_{1}$. One now takes $\tilde{X}_{1}(z)=0$ for all $z$ with $Y_{1}(z)=0$ and takes $\tilde{X}_{1}(z)=0$ (respectively, 1) with probability

$$
\frac{P_{x}\left\{\tilde{X}_{1}(z)=0\right\}-P_{x}\left\{Y_{1}(z)=0\right\}}{1-P_{x}\left\{Y_{1}(z)=0\right\}}
$$

(respectively,

$$
\left.\frac{P_{x}\left\{\tilde{X}_{1}(z)=1\right\}}{1-P_{x}\left\{Y_{1}(z)=0\right\}}\right)
$$

when $Y_{1}(z)=1$. These choices of the $\tilde{X}_{1}(z), z \in \mathbb{Z}^{d}$, are conditionally in$\widetilde{\widetilde{X}}^{\text {dependent }}$ (given $x$ and all the $\widetilde{X}_{1}^{v}, v \in \mathbb{Z}^{d}$ ). One easily checks that the $\widetilde{X}_{1}(z), z \in \mathbb{Z}^{d}$, are independent under this construction and, therefore, have the correct distribution. Moreover,

$$
\tilde{X}_{1}(z) \leq Y_{1}(z), \quad z \in \mathbb{Z}^{d} .
$$

Now assume that (4.19) has been proven for $n=1, \ldots, k$ and that we have also constructed joint versions of all $\tilde{X}_{n}, \widetilde{X}_{n}^{v}, v \in \mathbb{Z}^{d}, 0 \leq n \leq k$, such that

$$
\begin{equation*}
\tilde{X}_{n}(z) \leq Y_{n}(z), \quad z \in \mathbb{Z}^{d}, 0 \leq n \leq k \tag{4.21}
\end{equation*}
$$

Let $C_{k}=\left\{w_{1}, \ldots, w_{r}\right\}$ be the collection of sites for which $\tilde{X}_{k}(w)=1$. A fortiori $Y_{k}\left({\underset{\sim}{w}}_{i}\right)=1$ for $1 \leq i \leq r$ and hence we can choose a $v_{i}$ such that $x\left(v_{i}\right)=1, \quad \tilde{X}_{k}^{v_{i}}\left(w_{i}\right)=1$. The $v_{i}, 1 \leq i \leq r$, are not, in general, distinct. Call $w_{i} \sim w_{j}$ if $v_{i}=v_{j}$. Then $C_{k}$ breaks up into disjoint equivalence classes, say $A_{1}, \ldots, A_{s}$. Thus $w_{i}$ and $w_{j}$ belong to the same subclass if and only if $v_{i}=v_{j}$. We now apply (4.5) with these $A_{i}$ or rather with $A_{i} \cap \bar{B}(z, \rho)$ for fixed $z$. Specifically, we first use

$$
P_{x}\left\{\tilde{X}_{k+1}(z)=0 \text { for all } z \text { in } A \mid \tilde{X}_{k}\right\}=\prod_{z \in A} P_{\tilde{X}_{k}}\left\{\tilde{X}_{1}(z)=0\right\}
$$

Then, by (4.5), for fixed $z$,

$$
\begin{array}{r}
P_{\tilde{X}_{k}}\left\{\tilde{X}_{1}(z)=0\right\} \geq \prod_{i=1}^{s} P\left\{\tilde{X}_{1}(z)=0 \mid \tilde{X}_{0}(w)=1 \text { in } \bar{B}(z, \rho)\right. \text { exactly } \\
\text { for the } \left.w \text { 's in } A_{i} \cap \bar{B}(z, \rho)\right\}
\end{array}
$$

Finally, on the event $\left\{\widetilde{X}_{k}^{v_{i}}(w)=1\right.$ for $\left.w \in A_{i}\right\}$,

$$
\begin{aligned}
& P\left\{\tilde{X}_{1}(z)=0 \mid \tilde{X}_{0}(w)=1 \text { in } \bar{B}(z, \rho) \text { exactly for the } w \text { 's in } A_{i} \cap \bar{B}(z, \rho)\right\} \\
& =P\left\{\tilde{X}_{k+1}^{v_{i}}(z)=0 \mid \tilde{X}_{k}^{v_{i}}(w)=1 \text { in } \bar{B}(z, \rho)\right. \text { exactly } \\
& \text { for the } \left.w ' s \text { in } A_{i} \cap \bar{B}(z, \rho)\right\} \\
& \geq P\left\{\widetilde{X}_{k+1}^{v_{i}}(z)=0 \mid \widetilde{X}_{k}^{v_{i}}\right\} \quad[\text { by }(4.2)] .
\end{aligned}
$$

Thus, on the event $\left\{\widetilde{X}_{k}(w)=1\right.$ on $\left.C_{k}=\bigcup_{1}^{s} A_{i}\right\}$, which is contained in

$$
\bigcap_{i=1}^{s}\left\{\widetilde{X}_{k}^{v_{i}}(w)=1 \text { for } w \in A_{i}\right\}
$$

and under (4.21), it holds that

$$
\begin{aligned}
& P_{x}\left\{\tilde{X}_{k+1}(z)=0 \text { for all } z \text { in } A \mid \tilde{X}_{k}, \tilde{X}_{k}^{v}, v \in \mathbb{Z}^{d}\right\} \\
& \quad \geq \prod_{z \in A} \prod_{i=1}^{s} P\left\{\tilde{X}_{k+1}^{v_{i}}(z)=0 \mid \tilde{X}_{k}^{v_{i}}\right\} \\
& \quad=\prod_{i=1}^{s} P\left\{\tilde{X}_{k+1}^{v_{i}}(z)=0 \text { for all } z \in A \mid \widetilde{X}_{k}^{v_{i}}\right\} \\
& \quad \geq P_{x}\left\{Y_{k+1}(z)=0 \text { for all } z \in A \mid \tilde{X}_{k}^{v}, v \in \mathbb{Z}^{d}\right\}
\end{aligned}
$$

As before in the case $n=1$, we can now couple $\widetilde{X}_{k+1}$ and $\widetilde{X}_{k+1}^{v}$ such that (4.21) continues to hold for $n=k+1$. Equation (4.19) for $n=k+1$ is immediate from (4.21) for $n \leq k+1$.

Proof of Proposition 1. We have to check (1.2), (1.4) and (1.22)-(1.24) for the chain $\left\{X_{n}\right\}$ defined in (1.36) [on the space $S$ of (4.6)]. For the set $\mathscr{U}_{2}$ of (1.24) we shall take the singleton $\{\delta(\mathbf{0})\}$, that is, the state which has only one occupied site (which is automatically the leftmost particle and therefore shifted to the origin in $X_{n}$ ). For $\mathscr{U}_{1}$ in (1.23) we shall take the set

$$
\begin{equation*}
\mathscr{U}_{1}(K):=\left\{x \in S_{0}=S \backslash\{\underline{0}\}:|x| \leq K\right\} \tag{4.22}
\end{equation*}
$$

for some large $K$, yet to be determined. Also, $x_{0}$ will be taken equal to $\delta(\mathbf{0})$.
Now first observe that from any state $x \in S$ the transition to $x_{0}$ in one step is possible. This (and several of the succeeding statements) is proven somewhat more simply in terms of the chain $\tilde{X}$. than in terms of $X$. itself. Indeed, if $v_{1}, \ldots, v_{l}$ are the occupied sites of a state $x \in \widetilde{S}$, then with $z_{0}$ as in (4.4),

$$
\begin{aligned}
P_{x}\left\{\left|\tilde{X}_{1}\right|=1\right\} & \geq P_{x}\left\{\tilde{X}_{1}(z)=1 \text { only for } z=v_{1}+z_{0}\right\} \\
& =P_{x}\left\{\tilde{X}_{1}\left(v_{1}+z_{0}\right)=1\right\} \prod_{z \neq v_{1}+z_{0}} P_{x}\left\{\tilde{X}_{1}(z)=0\right\}
\end{aligned}
$$

Now for some constant $D_{1}>0$,

$$
P_{x}\left\{\tilde{X}_{1}\left(v_{1}+z_{0}\right)=1\right\} \geq D_{1}
$$

for all $x$ with $x\left(v_{1}\right)=1$, by (4.2) and (4.4). Also $P_{x}\left\{\tilde{X}_{1}(z)=0\right\}=1$ when $z \notin \bar{B}\left(v_{i}, \rho\right)$ for some $1 \leq i \leq l\left[\right.$ by (1.35) and (4.4)]. For $z \in \bar{B}\left(v_{i}, \rho\right)$,

$$
P_{x}\left\{\tilde{X}_{1}(z)=0\right\}=1-P_{x}\left\{\tilde{X}_{1}(z)=1\right\} \geq D_{2}>0
$$

for some constant $D_{2}$ independent of $x$ and $z[b y$ (1.33), (1.35) and (4.3)]. Thus

$$
\begin{equation*}
P_{x}\left\{\left|\widetilde{X}_{1}\right|=1\right\} \geq D_{1} D_{2}^{l(2 \rho+1)^{d}} \tag{4.23}
\end{equation*}
$$

This proves the irreducibility condition (1.2). It also proves (1.24) (with $n_{0}=1$ ) for our choice of $\mathscr{U}_{1}$ and $\mathscr{U}_{2}$, because the estimate (4.23) depends on $l=|x|$ only; the right-hand side is at least

$$
\begin{equation*}
D_{1} D_{2}^{K(2 \rho+1)^{d}} \tag{4.24}
\end{equation*}
$$

for any $x \in \mathscr{U}_{1}(K)$.
Condition (1.4), the certainty of absorption, follows directly from the assumption that $\widetilde{X}$. is (strongly) subcritical and the fact that $\widetilde{X}_{n}(z)=0$ automatically for all

$$
z \notin \bigcup \bar{B}(v, n \rho)
$$

where the union is over the finitely many $v$ with $\widetilde{X}_{0}(v) \neq 0$ [cf. proof of (3.4)].

Condition (1.23) follows from Lemma 9. Indeed, if $x(v)=1$ exactly when $v \in\left\{v_{1}, \ldots, v_{l}\right\}$, with $l \leq K$, then by (4.19),

$$
\begin{aligned}
P_{x}\{T>n\} & =P_{x}\left\{\tilde{X}_{n}(z)=1 \text { for some } z\right\} \\
& \leq P_{x}\left\{Y_{n}(z)=1 \text { for some } z\right\} \\
& \leq \sum_{i=1}^{l} P_{\delta_{\left(v_{i}\right)}}\left\{\tilde{X}_{n}(z)=1 \text { for some } z\right\} \\
& \leq l P_{\delta(\mathbf{0})}\{T>n\} \leq K P_{\delta(\mathbf{0})}\{T>n\}
\end{aligned}
$$

Finally, condition (1.22) for any sufficiently large choice of $K$ is guaranteed by Theorem 2.
5. A central limit theorem. Finally we prove Theorem 3. We shall prove (1.41) for a single time only; that is, we take, $l=1$ and $s_{l}=1$. The general case follows then easily by induction on $l$, from the Markov property. Also, as before, we restrict ourselves to the aperiodic case (except when we prove $M=0, \Sigma=\sigma^{2} \times$ identity matrix for the discrete time contact process, which has period 2).

For the proof we shall make use of an analogue of the $P^{*}$-measure used in Lemma 8. However, this time we need a measure on the space $\widetilde{S}_{0}:=\widetilde{S} \backslash\{\underline{0}\}$, which is the state space for $\widetilde{X}$. minus its absorbing state. To define this we must first lift the function $f$ of (1.17), which is defined on $S_{0}$, to $\widetilde{S}_{0}$. To this end let $\pi$ be the projection from $\widetilde{S}_{0}$ to $S_{0}$. That is, if $\widetilde{x}(z) \neq 0$ exactly for $z \in\left\{\xi_{1}, \ldots, \xi_{\nu}\right\}$, then we take $\pi(\tilde{x})=\tilde{x} \oplus\left(-\xi_{1}\right)$, which has $\pi(\tilde{x})(z) \neq 0$ if and only if $z \in\left\{\mathbf{0}, \xi_{2}-\xi_{1}, \ldots, \xi_{\nu}-\xi_{1}\right\}$. If $\widetilde{X}_{n}=\tilde{x}$, then $\pi(\tilde{x})$ is the corresponding state of $X_{n}$. We now define

$$
\tilde{f}(\tilde{x})=f(\pi(\tilde{x})), \quad \tilde{x} \in \widetilde{S}_{0}
$$

It is easy to see that for $x=\pi(\tilde{x})$,

$$
\begin{aligned}
\sum_{\tilde{y} \in \tilde{S}} & P\left(\tilde{X}_{n}=\tilde{y} \mid \tilde{X}_{0}=\tilde{x}\right) \tilde{f}(\tilde{y}) \\
& =E\left\{f\left(X_{n}\right) \mid \widetilde{X}_{0}=\tilde{x}\right\}=E\left\{f\left(X_{n}\right) \mid X_{0}=x\right\} \\
& =R^{-n} f(x) \quad[\text { see }(1.17)] \\
& =R^{-n} \tilde{f}(\tilde{x})
\end{aligned}
$$

Thus if we define for $C \in \widetilde{S}_{0}^{n+1}$,

$$
\begin{align*}
& P^{*}\left\{\tilde{X}_{0}, \tilde{X}_{1}, \ldots, \tilde{X}_{n} \in C \mid \tilde{X}_{0}=\tilde{x}\right\} \\
& \quad=\frac{R^{n}}{\tilde{f}(\tilde{x})} P\left\{\tilde{X}_{0}, \ldots, \tilde{X}_{n} \in C \mid \tilde{X}_{0}=\widetilde{x}\right\} \tilde{f}\left(\tilde{X}_{n}\right) \tag{5.1}
\end{align*}
$$

then $P^{*}$ defines an honest distribution on the space of paths $\widetilde{S}_{0}^{\mathbb{Z}_{+}}$. (In order not to overburden notation we shall refrain from attaching asterisks to the $\tilde{X}$
as well, in contrast to what we did in Lemma 8. Note also that the measure induced by $P^{*}$ on the paths of the $X$-process agrees with the $P^{*}$-measure of Lemma 8.)

Now fix $\tilde{x} \in \widetilde{S}_{0}$ and $y \in S_{0}$ for the remainder of the proof. Let $x=\pi(\tilde{x}) \in S_{0}$ and let $B \subset \mathbb{Z}^{d}$. Then

$$
\begin{align*}
& P\left\{\xi_{1}(n) \in B, \quad X_{n}=y \mid \widetilde{X}_{0}=\tilde{x}, T>n\right\} \\
& \quad=\frac{f(x)}{f(y)} \frac{1}{\tilde{f}(\widetilde{x})} P\left\{\xi_{1}(n) \in B, \quad X_{n}=y \mid \widetilde{X}_{0}=\widetilde{x}\right\} f(y)\left[P_{x}\{T>n\}\right]^{-1} \\
& \quad \sim\left[f(y) \sum_{w \in S_{0}} \mu(w)\right]^{-1} \frac{R^{n}}{\tilde{f}(\widetilde{x})} P\left\{\xi_{1}(n) \in B, \quad X_{n}=y \mid \widetilde{X}_{0}=\widetilde{x}\right\} f(y) \tag{5.2}
\end{align*}
$$

[by (1.26)]

$$
=\left[f(y) \sum_{w \in S_{0}} \mu(w)\right]^{-1} P^{*}\left\{\xi_{1}(n) \in B, \pi\left(\tilde{X}_{n}\right)=y \mid \tilde{X}_{0}=\tilde{x}\right\} .
$$

We now define $\tau_{0}=0<\tau_{1}<\cdots$ as the successive times at which $X$. visits $x$ and write

$$
\begin{equation*}
W_{i}=\left(\widetilde{X}_{\tau_{i-1}}, \widetilde{X}_{\tau_{i-1}+1}, \ldots, \widetilde{X}_{\tau_{i}}\right), \quad i \geq 1, \tag{5.3}
\end{equation*}
$$

for the $i$ th excursion between visits to $x$. We already remarked in Lemma 8 , that under $P^{*}, X$. is positive-recurrent, so that a.e. [ $P^{*}$ ] all $\tau_{i}$ are finite and the $W_{i}$ are well defined. Also by the strong Markov property, and the translation invariance property (1.33), the shifted excursions

$$
\begin{align*}
& W_{i} \oplus\left(-\xi_{1}\left(\tau_{i-1}\right)\right) \\
& \qquad:=\left(\widetilde{X}_{\tau_{i-1}} \oplus\left(-\xi_{1}\left(\tau_{i-1}\right)\right), \widetilde{X}_{\tau_{i-1}+1} \oplus\left(-\xi_{1}\left(\tau_{i-1}\right)\right),\right.  \tag{5.4}\\
& \left.\quad \ldots, \widetilde{X}_{\tau_{i}} \oplus\left(-\xi_{1}\left(\tau_{i-1}\right)\right)\right), \quad i \geq 1,
\end{align*}
$$

are i.i.d. Define further
(5.5) $\quad \Lambda_{i}=$ "length" of $W_{i}=\tau_{i}-\tau_{i-1}$,
(5.6) $\quad \Delta_{i}=$ "displacement" of $W_{i}=\xi_{1}\left(\tau_{i}\right)-\xi_{1}\left(\tau_{i-1}\right)$
and

$$
\Gamma_{i}=\text { "diameter" of } W_{i}=\max \left\{\left|\xi_{1}(n)-\xi_{1}\left(\tau_{i-1}\right)\right|: \tau_{i-1} \leq n \leq \tau_{i}\right\} .
$$

These are functions of $W_{i}$ and, hence, $\left(\Lambda_{i}, \Delta_{i}, \Gamma_{i}\right)_{i \geq 1}$ are also i.i.d. Moreover (1.28) shows that the distribution of $\Lambda_{i}$ under $P^{*}$ has an exponentially bounded
tail, as follows:

$$
\begin{aligned}
P^{*}\left\{\Lambda_{i} \geq l \mid \tilde{X}_{0}=\tilde{x}\right\} & =P^{*}\left\{\Lambda_{1} \geq l \mid \tilde{X}_{0}=\tilde{x}\right\} \\
& =\sum_{k=l}^{\infty} P^{*}\left\{\Lambda_{1}=k \mid \widetilde{X}_{0}=\tilde{x}\right\}=\sum_{k=l}^{\infty} P^{*}\left\{\tau_{1}=k \mid \widetilde{X}_{0}=\widetilde{x}\right\} \\
& =\sum_{k=l}^{\infty} R^{k} P\left\{X_{k}=x, X_{r} \neq x, 1 \leq r \leq k-1 \mid X_{0}=x\right\} \\
& \leq \sum_{k=l}^{\infty} C_{4}\left(\frac{R}{R+\eta(x)}\right)^{k} \leq D_{1}(x)\left(\frac{R}{R+\eta(x)}\right)^{l}
\end{aligned}
$$

In addition we claim that for some $D_{2}=D_{2}(x)<\infty$,

$$
\begin{equation*}
\Gamma_{i} \leq 2 \Lambda_{i} \rho+D_{2} \tag{5.8}
\end{equation*}
$$

To see this, note that if $\tilde{X}_{n}$ has occupied sites at $\xi_{1}(n), \ldots, \xi_{\nu(n)}(n)$, then the occupied sites of $\tilde{X}_{n+1}$ must all lie in

$$
\bigcup_{i=1}^{\nu(n)} \bar{B}\left(\xi_{i}(n), \rho\right)
$$

[by virtue of (1.35) and (1.31); compare (3.4)]. Therefore the occupied sites of $\widetilde{X}_{n}$ for any $\tau_{i-1} \leq n \leq \tau_{i}$ have to lie in

$$
\bigcup_{i=1}^{\nu\left(\tau_{i-1}\right)} \bar{B}\left(\xi_{i}\left(\tau_{i-1}\right),\left(\tau_{i}-\tau_{i-1}\right) \rho\right)
$$

and

$$
\begin{aligned}
\Gamma_{i} & \leq 2\left(\tau_{i}-\tau_{i-1}\right) \rho+\max _{1 \leq r \leq \nu}\left\|\xi_{r}\left(\tau_{i-1}\right)-\xi_{1}\left(\tau_{i-1}\right)\right\| \\
& \leq 2 \Lambda_{i} \rho+\max _{r, s}\left\|v_{r}-v_{s}\right\|
\end{aligned}
$$

${\underset{\sim}{w}}_{\text {where }} v_{1}, v_{2}, \ldots$ are the occupied sites of the state $x$ (because $\tilde{X}_{\tau_{i-1}}=$ $\left.\tilde{x}, X_{\tau_{i-1}}=x\right)$ and $\nu=\nu\left(\tau_{i-1}\right)$. Thus (5.8) holds with $D_{2}=\max _{r, s}\left\|v_{r}-v_{s}\right\|$.

Next we define

$$
\theta(n)=\max \left\{i: \tau_{i} \leq n\right\}
$$

Then

$$
\xi_{1}(n)=\xi_{1}(\theta(n))+\xi_{1}(n)-\xi_{1}(\theta(n))=\xi_{1}(0)+\sum_{1}^{\theta(n)} \Delta_{i}+\xi_{1}(n)-\xi_{1}(\theta(n))
$$

and

$$
\left|\xi_{1}(n)-\xi_{1}(\theta(n))\right| \leq \Gamma(\theta(n)+1) \leq 2 \Lambda_{\theta(n)+1} \rho+D_{2}
$$

In addition,

$$
\frac{1}{\sqrt{r}} \Lambda_{r} \rightarrow 0 \quad \text { almost surely }\left[P^{*}\right]
$$

because the $\Lambda_{r}$ are i.i.d. and have all moments under $P^{*}$. Also, by the strong law of large numbers,

$$
\begin{equation*}
\frac{\theta(n)}{n} \rightarrow\left[E^{*}\left\{\Lambda_{1} \mid \widetilde{X}_{0}=\tilde{x}\right\}\right]^{-1} \quad \text { almost surely }\left[P^{*}\right] \tag{5.9}
\end{equation*}
$$

again because the $\Lambda_{r}$ are i.i.d. and $\tau_{i}=\sum_{1}^{i} \Lambda_{r}$. Thus

$$
\begin{equation*}
\frac{\Lambda_{\theta(n)+1}}{\sqrt{n}} \rightarrow 0 \quad \text { almost surely }\left[P^{*}\right] \tag{5.10}
\end{equation*}
$$

The standard proof of the central limit theorem for Markov chains [Chung (1967), Section I.16] now shows that

$$
\begin{equation*}
\frac{\xi_{1}(n)-n M}{\sqrt{n}} \Rightarrow N(\mathbf{0}, \Sigma) \tag{5.11}
\end{equation*}
$$

under $P^{*}$, conditioned on $X_{0}=\tilde{x}$, with

$$
\begin{gather*}
M=\frac{E^{*}\left\{\Delta_{1} \mid \tilde{X}_{0}=\widetilde{x}\right\}}{E^{*}\left\{\Lambda_{1} \mid \widetilde{X}_{0}=\widetilde{x}\right\}}  \tag{5.12}\\
\Sigma(i, j)=\frac{1}{E^{*}\left\{\Lambda_{1} \mid \widetilde{X}_{0}=\tilde{x}\right\}} E^{*}\left\{\left(\Delta_{1, i}-M_{i} \Lambda_{1}\right)\left(\Delta_{1, j}-M_{j} \Lambda_{1}\right) \mid \tilde{X}_{0}=\widetilde{x}\right\} \tag{5.13}
\end{gather*}
$$

where $\Delta_{1, i}$ and $M_{i}$ are the $i$ th component of $\Delta_{1}$ and $M$, respectively. Slightly more explicit expressions for $M$ and $\Sigma$ can be given in the same form as in Theorems I.14.5, I.14.7 and its Corollary in Chung (1967). Note in this connection that

$$
\begin{align*}
& E^{*}\left\{\text { number of indices } 0 \leq n<\tau_{1} \text { with } X_{n}=v \mid \tilde{X}_{0}=\tilde{x}\right\} \\
& \quad=E^{*}\left\{\text { number of indices } 0 \leq n<\tau_{1} \text { with } X_{n}=v \mid X_{0}=x\right\}  \tag{5.14}\\
& \quad=\frac{f(v) \mu(v)}{f(x) \mu(x)}
\end{align*}
$$

because $f(v) \mu(v)$ is the unique invariant probability measure for the $X$. chain under $P^{*}$. That is,

$$
\begin{align*}
& \sum_{x \in S_{0}} f(x) \mu(x) P^{*}\left\{X_{1}=v \mid X_{0}=x\right\} \\
& \quad=R f(v) \sum_{x \in S_{0}} \mu(x) P\left\{X_{1}=v \mid X_{0}=x\right\}=f(v) \mu(v) \tag{5.15}
\end{align*}
$$

[by (1.18)] and (1.19) holds. We shall not pursue more explicit expressions for $M$ and $\Sigma$ here; see, however, the end of this proof.

In view of (5.11) we basically only have to prove the independence of $X_{n}$ and $\left[\xi_{1}(n)-n M\right] / \sqrt{n}$. The arguments for this are at least part of the folklore and we shall therefore be brief. For $\varepsilon>0$ we can choose $n_{5}$ such that for all $n \geq 2 n_{5}$,

$$
\begin{align*}
& \left|P^{*}\left\{X_{n}=y \mid \tilde{X}_{0}=\tilde{x}\right\}-f(y) \mu(y)\right|  \tag{5.16}\\
& \quad=\left|P^{*}\left\{X_{n}=y \mid X_{0}=x\right\}-f(y) \mu(y)\right| \leq \varepsilon
\end{align*}
$$

and

$$
\begin{equation*}
P^{*}\left\{\tau\left(\theta\left(n-2 n_{5}\right)+1\right) \geq n-n_{5} \mid \widetilde{X}_{0}=\widetilde{x}\right\} \leq \varepsilon \tag{5.17}
\end{equation*}
$$

Equation (5.16) is immediate from (1.26) and (1.27) after a translation to the $P^{*}$-measure. In (5.17) we have written $\tau\left(\theta\left(n-2 n_{5}\right)+1\right)$ instead of $\tau_{\theta\left(n-2 n_{5}\right)+1}$ for typographical convenience; the estimate (5.17) itself is straightforward renewal theory, since $\tau\left(\theta\left(n-2 n_{5}\right)+1\right)$ is the smallest $m>n-2 n_{5}$ with $X_{m}=x$ [see Chung (1967), Theorem I.14.2]. It is further easy to see by means of (5.8) and Chung [(1967), Theorem I.14.2] that for fixed $n_{5}$,

$$
\lim _{a \rightarrow \infty} P^{*}\left\{\left|\xi_{1}(n)-\xi_{1}\left(\tau\left(\theta\left(n-2 n_{5}\right)+1\right)\right)\right|>a \mid \tilde{X}_{0}=\tilde{x}\right\}=0
$$

uniformly in $n \geq 2 n_{5}$. Therefore,

$$
\begin{align*}
& \frac{1}{\sqrt{n}}\left|\xi_{1}(n)-\xi_{1}\left(\tau\left(\theta\left(n-2 n_{5}\right)+1\right)\right)\right| \rightarrow 0 \quad \text { in } P^{*} \text {-measure, }  \tag{5.18}\\
& \text { conditioned on } \tilde{X}_{1}(0)=\tilde{x}, \text { as } n \rightarrow \infty
\end{align*}
$$

Finally, a decomposition with respect to the value of $\tau\left(\theta\left(n-2 n_{5}\right)+1\right)$ gives

$$
\begin{aligned}
& \mid P^{*}\left\{\xi_{1}\left(\tau\left(\theta\left(n-2 n_{5}\right)+1\right)\right) \in B, X_{n}=y \mid \tilde{X}_{0}=\tilde{x}\right\} \\
& \quad-P^{*}\left\{\xi_{1}\left(\tau\left(\theta\left(n-2 n_{5}\right)+1\right)\right) \in B \mid \widetilde{X}_{0}=\widetilde{x}\right\} f(y) \mu(y) \mid \\
& \leq P^{*}\left\{\tau\left(\theta\left(n-2 n_{5}\right)+1\right) \geq n-n_{5} \mid \widetilde{X}_{0}=\tilde{x}\right\} \\
& \quad+\sum_{n-2 n_{5}<r<n-n_{5}} P^{*}\left\{\tau\left(\theta\left(n-2 n_{5}\right)+1\right)=r, \xi_{1}(r) \in B \mid \widetilde{X}_{0}=\tilde{x}\right\} \\
& \quad \times\left|P^{*}\left\{X_{n-r}=y \mid X_{0}=x\right\}-f(y) \mu(y)\right|
\end{aligned}
$$

$\leq 2 \varepsilon \quad[b y(5.16)$ and (5.17) $]$.
Since $\varepsilon>0$ is arbitrary, it is now not hard to obtain from (5.11), (5.18) and (5.19) that

$$
\begin{aligned}
& \lim _{n \rightarrow \infty} P^{*}\left\{\frac{\xi_{1}(n)-n M}{\sqrt{n}}<\gamma, X_{n}=y \mid \tilde{X}_{0}=\tilde{x}\right\} \\
& \quad=P\left\{G_{1} \leq \gamma\right\} f(y) \mu(y)
\end{aligned}
$$

Translating this back to the $P$-measure gives (1.41) for $l=1, s_{l}=1$ [by means of (1.26) again].

The last statement of Theorem 3 is about the values of $M$ and $\Sigma$ for the discrete time contact process. Since $M$ and $\Sigma$ are independent of the choice of $\tilde{x}$ [see Chung (1967), Corollary to Theorem I. 15.4 and Corollary 2 to Theorem I.16.1], we can take $\tilde{x}=\delta(\mathbf{0})$. It is then clear from the symmetry of the discrete time contact process that $M=\mathbf{0}$ and $\Sigma_{i j}=0$ for $i \neq j$ [see (5.12) and (5.13)]. Also all $\Sigma_{i, i}, 1 \leq i \leq d$, must have the same value, say $\sigma^{2}$, so the only nontrivial part of our claim is that

$$
\sigma^{2}=\frac{1}{E^{*}\left\{\Lambda_{1} \mid \tilde{X}_{0}=\delta(\mathbf{0})\right\}} E^{*}\left\{\left(\Delta_{1,1}-M_{1} \Lambda_{1}\right)^{2} \mid \tilde{X}_{0}=\delta(\mathbf{0})\right\}>0
$$

or, equivalently, that

$$
P^{*}\left\{\Delta_{1,1}-M_{1} \Lambda_{1} \neq 0 \mid \tilde{X}_{0}=\delta(\mathbf{0})\right\}>0 .
$$

However, $M_{1}=0$ and one easily sees that

$$
\begin{aligned}
& P^{*}\left\{\Delta_{1,1} \neq 0 \mid \widetilde{X}_{0}=\delta(\mathbf{0})\right\} \\
& \quad \geq P^{*}\left\{\tau_{1}=2, \Delta_{1,1}=2 \mid \widetilde{X}_{0}=\delta(\mathbf{0})\right\} \\
& \quad \geq P^{*}\left\{\widetilde{X}_{1}=\delta\left(e_{1}\right)+\delta\left(-e_{1}\right), \widetilde{X}_{2}=\delta\left(2 e_{1}\right) \mid \widetilde{X}_{0}=\delta(\mathbf{0})\right\}>0,
\end{aligned}
$$

where $e_{1}=(1,0, \ldots, 0)$. Indeed the event $\left\{\widetilde{X}(1)=\delta\left(e_{1}\right)+\delta\left(-e_{1}\right)\right\}$ occurs if at time 1 exactly the sites $e_{1}$ and $-e_{1}$ are occupied; similarly $\left\{\widetilde{X}(2)=\delta\left(2 e_{1}\right)\right\}$ is the event that only the site $2 e_{1}$ is occupied at time 2. The description in Section 4 of the discrete time contact process shows that these events have strictly positive probability. Thus $\sigma^{2}>0$ and the proof of Theorem 3 is complete.

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