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## Products of Random Weights Indexed by Galton-Watson Trees

Publications de l'Institut de recherche mathématiques de Rennes, 1996-1997, fascicule 2 «Fascicule de probabilités », , p. 1-24<br>[http://www.numdam.org/item?id=PSMIR_1996-1997___2_A5_0](http://www.numdam.org/item?id=PSMIR_1996-1997___2_A5_0)

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# Products of random weights indexed by Galton-Watson trees 

June 1997

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Mathematics Subject Classification (1991):
Primary: 60J80; Secondary: 60J15, 60G42, 60G57.
Keywords and phrases:
Galton-Watson trees, self-similar cascades, branching random walks, random measures, martingales, functional equations, moments, tails, continuity.


#### Abstract

We consider the space $\Omega$ of marked Galton-Watson trees. Each $\omega \in \Omega$ corresponds to a tree $T=T(\omega)$. The tree begins with the root $\emptyset \in T(\omega)$ representing the initial ancestor; each node (=individual) $u \in T(\omega)$ gives birth to $N_{u}=N_{u}(\omega) \geq 0$ children, denoted by $u i \in \omega, 1 \leq i \leq N_{u}$, and each child $u i$ is marked with a non-negative number $A_{u i}=A_{u i}(\omega)$. The individus in the same generation behave independently each other, with the same probability law given by the random variable $\left(N, A_{1}, \ldots, A_{N}\right):=$ $\left(N_{\emptyset}, A_{\emptyset_{1}}, \ldots, A_{\emptyset N}\right)$, which is normalized such that $E \sum_{i=1}^{N} A_{i}=1$. Write $X_{u}=A_{u_{1}} A_{u_{1} u_{2}} \cdots A_{u_{1} \ldots u_{n}}$ if $u=u_{1} \ldots u_{n}$, and put $Y_{n}=\sum_{|u|=n} X_{u}$, where the sum is taken over all nodes $u \in T(\omega)$ of length $n$. Then $\left\{Y_{n}: n \geq 1\right\}$ forms a martingale, and converges almost surely to a non-negative random variable, $Z$, as $n \rightarrow \infty$. In the case where the limit is non-degenerate $(P(Z=0)<1)$, we give necessary and sufficient conditions for existence of its moments of given order $p>1$, obtain an equivalence of the tail probabilities $P(Z>x)$ as $x \rightarrow \infty$, and prove that its distribution has a continuous density (with respect to the Lebesgue measure) on ( $0, \infty$ ) under some moment conditions. The results are of applications in the study of: (a) Mandelbrot's self-similar cascades, (b) invariant measures of some infinite particle systems, (c) branching random walks, (d) flows in trees and (e) exact Housdorff measures in random constructions. The proofs make use of the random difference equation


$X=A_{1} X_{1}+B_{1}$, where $X_{1}$ is a random variable independent of $\left(A_{1}, B_{1}\right) \in \mathbb{R}_{+}^{2}$ and has the same law as $X$.

## 1. Introduction and main results

Let $\mathbb{N}^{\star}=\{1,2, \cdots\}$ be the set of positive integers with the discret topology. Put $\mathbb{N}=\{0\} \cup \mathbb{N}^{\star}$ and write

$$
\mathbf{U}=\bigcup_{k=0}^{\infty}\left(\mathbb{N}^{\star}\right)^{k}
$$

for the set of all finite sequences $\mathbf{i}=i_{1} i_{2} \ldots i_{n}\left(i_{k} \in \mathbb{N}^{\star}\right)$, where by convention $\mathbb{N}^{\star 0}=\{\emptyset\}$ contains the null sequence $\emptyset$. Let

$$
\mathrm{I}=\left(\mathbb{N}^{\star}\right)^{\mathbb{N}^{\star}}
$$

be the set of all infinite sequences $\mathbf{i}=i_{1} i_{2} \ldots\left(i_{k} \in \mathbb{N}^{\star}\right)$ with the product topology. If $\mathbf{i}=i_{1} i_{2} . . i_{n}(n \leq \infty)$ is a sequence, we write $|\mathrm{i}|=n$ for its length, and $\mathbf{i} \mid k=i_{1} i_{2} \ldots i_{k}(k \leq n ;)$ for the curtailment of $\mathbf{i}$ after k terms; conventioally, $|\emptyset|=0$ and $\mathbf{i} \mid 0=\emptyset$. If $\mathbf{i} \in \mathbf{U}$ and $\mathbf{j} \in \mathbf{U}$ or $\mathbf{I}$ we write $\mathbf{i j}=(\mathbf{i}, \mathbf{j})$ for the sequence obtained by juxtaposition. In particular $\mathbf{i} \emptyset=\emptyset \mathbf{i}=\mathrm{i}$. We partially order $U$ by writing $i<j$ to mean that for some $i^{\prime} \in U, j=i^{\prime}$, and we use a similar notation if $i \in U$ and $j \in I$. If $i$ and $j$ are two sequences, we write $i \wedge j$ for the common maximal sequence of $i$ and $j$, that is, the maximal sequence $q$ such that $q<i$ and $q<j$.

A tree $T$ is a subset of U satisfying three conditions (cf. Neveu (1986)):
i) $\emptyset \in T$;
ii) if $\mathbf{i j} \in T$, then $\mathbf{i} \in T$;
iii) if $\mathbf{i} \in T$ and $j \in \mathbb{N}^{\star}$, then $\mathbf{i} j \in T$ if and only if $1 \leq j \leq N_{\mathbf{i}}$ for a positive integer $N_{i}$.

We shall write $N$ for $N_{\emptyset}$. The boundary of a tree $T$ is defined as

$$
\partial T=\{\mathbf{i} \in \mathbf{I}: \quad \mathbf{i} \mid n \in T \text { for all } n \in \mathbb{N}\} .
$$

As a subset of $\mathbf{I}$, it is a metrical and compact topological space with

$$
B(\mathbf{i})=\{\mathbf{j} \in \partial \omega: \mathbf{i}<\mathbf{j}\}, \quad \mathbf{i} \in T(\omega)
$$

its topological basis; a possible choice of metric is

$$
d_{c}(\mathbf{i}, \mathbf{j})=c^{-|\mathbf{i} \wedge \mathbf{j}|},
$$

where $c$ is a given number in $(1, \infty)$. The set $B(\mathbf{i})$ is then a ball of radius $c^{-|i|}$.
An element of $T$ is called a node. Each node $u \in T$ is marked with a vector $\eta_{u}=\left(A_{u 1}, A_{u 2}, \ldots\right)$ of $\mathbb{R}_{+}^{\mathbb{N}^{*}}$, where $\mathbb{R}_{+}=[0, \infty)$. If $1 \leq j \leq N_{u}$, we can imagine that the number $A_{u j}$ is associated with the edge ( $u, u j$ ) linking the nodes $u$ and $u j$; the values $A_{u j}$ for $j>N_{u}$ are of no influence for our purpose, and will be taken as 0 for convenience. The marked tree will be denoted by $\left(T,\left(\eta_{u}, u \in T\right)\right)$.

Let $\mathbb{T}$ be the set of all trees, and $\Omega$ be the set of all marked trees $\omega$ (marked as above). An element $\omega$ of $\Omega$ will be written as $\left(T(\omega),\left(\eta_{u}, u \in T(\omega)\right)\right.$ ), where $T(\omega)$ is the underlying tree. We may regard $T$ as the canonical projection from $\Omega$ to $\mathbb{T}$. Thus $T$ may stand for a tree or an operator, according to the context. If $\omega$ is a marked tree and if $\mathbf{i} \in T(\omega)$, we write $T_{\mathbf{i}}(\omega)=\{\mathbf{j} \in \mathbf{U}$ : $\mathrm{ij} \in T(\omega)\}$ for the shifted tree of $T(\omega)$ at i . Note that $T_{\mathrm{i}}(\omega) \in \mathbb{T}$. Denote by

$$
z_{k}(\omega)=\{\mathbf{i} \in T(\omega):|\mathbf{i}|=k\}
$$

the set of nodes of $T(\omega)$ with length $k(k \in \mathbb{N})$, and consider the filtration

$$
\mathcal{F}_{k}=\sigma\left\{\left(N_{\mathbf{i}}, A_{\mathbf{i} 1}, A_{\mathbf{i} 2}, \ldots\right): \quad \mathrm{i} \in z_{k-1}\right\}, \quad k \geq 1
$$

Let $\mathcal{F}:=\sigma\left(\mathcal{F}_{k}, k \geq 1\right)$. For simplicity, we write $\left(N, A_{1}, A_{2}, \ldots\right)$ for $\left(N_{\emptyset}(\omega), A_{\emptyset_{1}}(\omega), A_{\emptyset_{2}}(\omega), \ldots\right)$.

By a result of Neveu (1986), for each probability distribution $q$ on $\mathbb{N} \times \mathbb{R}_{+}^{\mathbb{N}^{\star}}$, there is a probability law $P=P_{q}$ on $(\Omega, \mathcal{F})$ such that
(i) the distribution of the random variable $\left(N, A_{1}, A_{2}, \ldots\right)$ is $q$;
(ii) given $\mathcal{F}_{k}$, the random variables ( $N_{\mathbf{i}}, A_{\mathbf{i 1}}, A_{\mathbf{i} 2} \ldots$ ), $\mathbf{i} \in z_{k}(\omega)$, are conditionally independent, and their conditional distribution is q .

The property (ii) is referred as the branching property.
Assume that the initial distribution is normalized such that

$$
E\left(\sum_{i=1}^{N} A_{i}\right)=1 .
$$

If $u \notin T(\omega)$, the values $A_{u}(\omega)$ are of no influence for our problem, and may be non-defined; however, for convenience, we set

$$
A_{u}=0 \text { if } u \in \mathbf{U} \backslash T(\omega)
$$

In particular, for all $u \in T(\omega), A_{u i}=0$ if $u>N_{u}$. Put

$$
X_{\emptyset}=1, X_{u}=A_{u_{1}} A_{u_{1} u_{2}} \cdots A_{u_{1} \ldots u_{n}} \text { if } u=u_{1} \ldots u_{n} \in \mathbf{U}
$$

and

$$
Y_{n}=\sum_{|v|=n, v \in T(\omega)} X_{v}, \quad n \geq 1
$$

Then $\left(Y_{n}, \mathcal{F}_{n}\right)(n \geq 1)$ is a (non-negative) martingale, so the limit

$$
Z=\lim _{n \rightarrow} Y_{n}
$$

exists almost surely. Similarly, we put

$$
Z_{u}=\lim _{n \rightarrow \infty} \sum_{v=v_{1} \ldots v_{n} \in T_{u}(\omega)} A_{u v_{1}} A_{u v_{1} v_{2}} \cdots A_{u v_{1} \ldots v_{n}} \text { if } u \in T(\omega)
$$

and $Z_{u}=1$ if $u \in \mathbf{U} \backslash T(\omega)$. Then $Z=Z_{\mathfrak{\theta}}$, and, by the branching property, given $\mathcal{F}_{n}$, the random variables $Z_{u}, u \in z_{n}(\omega)$, are conditionally independent, and their conditional law is the distribution of $Z$. It is easily seen that for almost all $\omega \in \Omega$ and all $u \in T(\omega)$,

$$
X_{u} Z_{u}=\sum_{i=1}^{N_{u}} X_{u i} Z_{u i} .
$$

Therefore for almost all $\omega \in \Omega$, there is a unique Borel measure, $\mu_{\omega}$, defined on $\partial T(\omega)$, such that

$$
\mu_{\omega}(B(u))=X_{u} Z_{u} \text { for all } u \in T(\omega)
$$

We extend this measure as a Borel measure on I by letting $\mu_{\omega}(A)=$ $\mu_{\omega}(A \cap \partial T(\omega))$. Then $\mu_{\omega}$ is a random measure on I with mass $Z(\omega)$.

The preceding identity on $Z_{u}$ shows that the distribution of $Z$ satisfies the equation

$$
\begin{equation*}
Z=\sum_{i=1}^{N} A_{i} Z_{i} \tag{1.1}
\end{equation*}
$$

where the sum is taken to be zero if $N=0$, and, given $\left(N, A_{1}, A_{2}, \ldots\right)$, the random variables $Z_{i}(1 \leq i \leq N)$ are conditionally independent, and their conditional distribution is the law of $Z$. In terms of characteristic functions or Laplace transforms, it reads

$$
\phi(t)=E \prod_{i=1}^{N} \phi\left(A_{i} t\right)
$$

where $\phi(t)=E\left(e^{i t Z}\right)$ or $E\left(e^{-t Z}\right), t \in \mathbb{R}$ or $\mathbb{R}_{+}$, and the empty product is taken to be 1 .

The problem is to study the measure $\mu_{\omega}$ and the distribution of $Z$.
Kahane and Peyrière (1976), and Guivarc'h (1990) studied the problem in the case where $N$ is constant and the $A_{i}(1 \leq i \leq N)$ are i.i.d. Their works were motivated by questions raised by Mandelbrot related to a model of turbulence of Yaglom. Holley and Liggett (1981) studied the same problem in the case where $N$ is constant and the $A_{i}$ are fixed multiples of one random variable, and Durrett and Liggett (1983) considered the more general case where $N$ is constant but the $A_{i}$ have arbitrary joint distribution. Their works were motivated by a number of problems in infinite particle systems. Closely related results are given in Kahane (1987), Ben Nasr (1987), Holley and Waymire (1992), Collet and Koukiou (1992), Chauvin and Rouault (1996), and Liu and Rouault (1996), etc.

If $1<m=E N<\infty$ and $A_{i}=1 / m(1 \leq i \leq N)$, then $Y_{n}$ becomes the well-known martingale $\operatorname{card}\left(z_{n}\right) / m^{n}$ of the Galton-Walton process, where $\operatorname{card}\left(z_{n}\right)$ is the population size at $n-t h$ generation. Similar martingales arise in age-dependent branching processes or branching random walks. Many authors have contributed to the subject, see for example Harris (1948), Kesten and Stigum (1966), Seneta (1968 and 1969), Athreya (1971), Doney (1972 and 1973), Bingham and Doney (1974 and 1975) and Biggins (1977).

The martingale $\left(Y_{n}\right)$ and the equation (1.1), in its various forms, were also used to study some fractal sets or flows in networks, implicitly or directly by Mauldin and Williams (1986), Falconer (1986 and 1987) and Liu (1993 and 1996).

If $E$ is a set or a statement, we write $1_{E}$ or $1\{E\}$ for its indicator function. Let

$$
\tilde{N}:=\sum_{i=1}^{N} 1_{\left\{A_{i}>0\right\}}
$$

be the number of non-zero terms of $A_{i}, 1 \leq i \leq N$. To simplify the discussion, we suppose throughout the paper that

$$
\begin{gather*}
P(\tilde{N}=0 \text { or } 1)<1, \quad P\left(\forall i \in\{1, \ldots, N\}, A_{i}=0 \text { or } 1\right)<1,  \tag{1.2}\\
E \tilde{N}<\infty \text { and } E \sum_{i=1}^{N} A_{i} \log ^{+} A_{i}<\infty, \tag{1.3}
\end{gather*}
$$

where $\log ^{+} x=\max (0, \log x)$. For $x \in[0, \infty)$, write

$$
\begin{equation*}
S(x):=\sum_{i=1}^{N} A_{i}^{x} \text { and } \gamma(x):=E S(x), \tag{1.4}
\end{equation*}
$$

where (and throughout) the sum is taken over all the $i$ 's such that $A_{i}>0$. The function $\gamma$ is defined on $[0, \infty)$ with values in $[0, \infty]$. Then

$$
\gamma(x)<\infty \text { and } \gamma^{\prime}(x)=E \sum_{i=1}^{N} A_{i}^{x} \log A_{i}<\infty
$$

exists for all $x \in(0,1]$ (and $\gamma^{\prime}(1)$ denotes the left derivative); and $\gamma$ is strictly convex on $(0,1)$.

The following results have been known. (We recall that $Y_{1}=\sum_{i=1}^{N} A_{i}$ by our notations.)

Theorem 0. (Biggins 1977 and Liu 1997) (a) The following assertions are equivalent: (i) $Z$ is non-degenerate [i.e. $P(Z=0)<1$ ]; (ii) $E\left(Z \mid \mathcal{F}_{n}\right)=Y_{n}$ for all $n \geq 1$; (iii)

$$
\begin{equation*}
E\left[Y_{1} \log ^{+} Y_{1}\right]<\infty \text { and } E \sum_{i=1}^{N} A_{i} \log A_{i}<0 \tag{1.5}
\end{equation*}
$$

(b) Assume only (1.2) and (1.3) [so $\gamma(1)$ is not necessarily equal to 1$]$. Then the equation (1.1) has a non-trivial solution if and only if for some $\alpha \in(0,1]$, $\gamma(\alpha)=1$ and $\gamma^{\prime}(\alpha) \leq 0$; it has a nontrivial solution with finite mean if and only if $\gamma(1)=1$ and (1.5) holds. Moreover, there is at most one solution with mean 1 .

Here, we only consider solutions of (1.1) in the class of probability laws on $[0, \infty)$. We remark that $Z=1$ if and only if $Y_{1}=1$ almost surely. Therefore we suppose throughout that $Y_{1}$ is not a.s.a constant.

The main aim of this paper is to study the moments, the tail probabilities and the absolutely continuity of the distribution of $Z$, in the case where $Z$ is non-degenerate. The following results will be shown.

Theorem 1 (Moments). Assume (1.5). Then For each fixed $p>1$,

$$
E\left(Z^{p}\right)<\infty
$$

if and only if

$$
E\left[Y_{1}^{p}\right]<\infty \text { and } E\left[\sum_{i=1}^{N} A_{i}^{p}\right]<1 .
$$

We remark that if $\gamma(p)<1$ for all $p>1$, then $\left\|\sup _{1 \leq i \leq N} A_{i}\right\|_{p} \leq 1$ for all $p>1$, and so $\left\|\sup _{1 \leq i \leq N} A_{i}\right\|_{\infty} \leq 1$. Therefore, by Theorem $1, Z$ has moments of all orders if and only if

$$
\begin{equation*}
P\left(\forall i \in\{1, \ldots, N\}, A_{i} \leq 1\right)=1 \text { and } E\left[Y_{1}^{p}\right]<\infty \text { for all } p>1 \tag{1.6}
\end{equation*}
$$

The problem will be called lattice if for some $h>0$ and almost all $\omega \in \Omega$, each $\log A_{i}$ is an integer multiple of $h$ whenever $1 \leq i \leq N$ and $A_{i}>0$; the largest such h will be called the span. Otherwise, it is called non-lattice.

## Theorem 2 (Tail probabilities).

(a) (The case where $P\left(\exists i \in\{1, \ldots, N\}, A_{i}>1\right)>0$.) Suppose that for some $\chi>1$,

$$
E\left[\sum_{i=1}^{N} A_{i}^{\chi}\right]=1, E\left[\sum_{i=1}^{N} A_{i}^{\chi} \log ^{+} A_{i}\right]<\infty \text { and } E\left[\left(\sum_{i=1}^{N} A_{i}\right)^{\chi}\right]<\infty .
$$

If the problem is non-lattice, then there is a constant $c \in(0, \infty)$ such that

$$
\lim _{x \rightarrow \infty} x^{x} P(Z>x)=c
$$

If the problem is lattice, then

$$
0<\liminf _{x \rightarrow+\infty} x^{x} P(Z>x) \leq \limsup _{x \rightarrow+\infty} x^{x} P(Z>x)<\infty .
$$

(b)(The case where $P\left(\forall i \in\{1, \ldots, N\}, A_{i} \leq 1\right)=1$.) Suppose that $\|N\|_{\infty}<\infty$ and that for some $x>0,\left\|\sum_{i=1}^{N} A_{i}^{x}\right\|_{\infty} \leq 1$. Let $\rho$ be the least solution in $(1, \infty)$ of the equation

$$
\left\|\sum_{i=1}^{N} A_{i}^{\rho}\right\|_{\infty}=1
$$

(Such a solution certainly exists under the preceding conditions.) Assume also that for some constants $0<\delta<1,0 \leq a<\infty, 0<c<\infty$ and all $0<x<1$ sufficiently small,

$$
P\left(\sum_{i=1}^{N} A_{i}^{p}>1-x \text { and } A_{i} \leq \delta \text { for all } 1 \leq i \leq N\right) \geq c x^{a}
$$

Then for some constants $0<c_{1} \leq c_{2}<\infty$ and all $x>0$ sufficiently large,

$$
\exp \left\{-c_{2} x^{\rho /(\rho-1)}\right\} \leq P(Z \geq x) \leq \exp \left\{-c_{1} x^{\rho /(\rho-1)}\right\}
$$

In the non-lattice case, part (a) of the theorem is due to Guivarc'h (1990) if $N=c \geq 2$ is constant and $A_{i}(1 \leq i \leq c)$ are i.i.d. Our proof develops an idea of Guivarc'h, linking the distributional equation (1.1) with the random difference equation (see sections 2 and 3 ).

Let

$$
\zeta_{1}=\inf \left\{i>0: A_{i}>0\right\}, \text { where } \inf \emptyset=\infty,
$$

be the first (random) index $i$ for which $A_{i}>0$, and put

$$
\zeta_{k+1}=\inf \left\{i>\zeta_{k}: A_{i}>0\right\}, \text { where } \inf \emptyset=\infty, \quad k \geq 1
$$

Define

$$
A_{\zeta_{k}}(\omega)= \begin{cases}A_{i}(\omega) & \text { if } \zeta_{k}(\omega)=i \text { for some } i \in \mathbb{N}^{\star}, \\ 0 & \text { if } \zeta_{k}(\omega)=\infty\end{cases}
$$

Then $A_{\zeta_{k}}(\omega)>0$ if $\tilde{N}(\omega) \geq k \geq 1$. It is easily seen that the functional equation (1.1') is nothing but

$$
\phi(t)=E \prod_{i=1}^{\bar{N}} \phi\left(A_{\zeta_{i}} t\right) .
$$

The advantage of the new equation is that $A_{\zeta_{i}}>0$ for all $1 \leq i \leq \tilde{N}$.

Theorem 3 (Absolute continuity). The distribution of $Z$ is either absolutely continuous or (purely) singularly continuous on $\mathbb{R}_{+}^{\star}=(0, \infty)$. It is absolutely continuous if one of the following conditions holds:
(i) $\tilde{N} \geq 2$ a.s. and, for some $\epsilon>0$,

$$
E\left[A_{\zeta_{1}}^{-\epsilon}+A_{\zeta_{2}}^{-\epsilon}\right]<\infty, E \sum_{i=1}^{N} M_{i}^{-\epsilon}<\infty \text { and } E \sum_{i=1}^{N} A_{i} M_{i}^{-\epsilon}<\infty
$$

where $M_{i}=\max _{j \neq i} A_{j}$.
(ii) $\tilde{N} \geq 1$ a.s. and $A_{\zeta_{1}}$ has an absolutely continuous distribution on $\mathbb{R}_{+}^{\star}$.

Moreover, in case (i), the density function of the distribution of $Z$ is continuous (on $\mathbb{R}_{+}^{*}$ ).

We remark that the condition (i) holds if, for example, $N \geq 2$ almost surely and, given $N$, the random variables $A_{i}, 1 \leq i \leq N$, are conditionally independent, and their conditional law is the law of $A_{1}$, which satisfies the property that $E A_{1}^{-\epsilon}<\infty$ for some $\epsilon>0$.

Informations on the rate of convergence of characteristic functions or Laplace transforms of Z will be given in proofs.

## 2. The random difference equation

In this section, $(\Omega, \mathcal{F}, P)$ denotes an arbitrary probability space, $(A, B)$ and $\left(A_{n}, B_{n}\right)(n \geq 1)$ are i.i.d. random variable defined on $(\Omega, \mathcal{F}, P)$, with values in $\mathbb{R}^{2}$.

Consider the random difference equation

$$
\begin{equation*}
X \stackrel{L}{=} A X+B \tag{2.1}
\end{equation*}
$$

where X is a real random variable independent of $(A, B)$, and $\stackrel{L}{\underline{L}}$ denotes equality in law; the law of $X$ is unknown. In terms of characteristic functions, the equation reads

$$
\phi(t)=E\left[e^{i A t} \phi(B t)\right], \quad t \in R .
$$

A probability law, $\mu$, is said to be a solution of (2.1) if there is a random variable $X$ having $\mu$ as its law and satisfying (2.1); when we say that a random variable X is the unique solution of (2.1), we mean that the corresponding law is the unique solution.

Lemma 2.1 [Grintsevichyus (1974), Th. 1 and Prop.1]. If

$$
\begin{equation*}
P(A \neq 0)=1,-\infty<E \log |A|<0 \text { and } E \log ^{+}|B|<\infty \tag{2.2}
\end{equation*}
$$

then (2.1) has a unique solution, and this solution is given by

$$
\begin{equation*}
X=B_{1}+A_{1} B_{2}+A_{1} A_{2} B_{3}+\cdots+A_{1} \ldots A_{n-1} B_{n}+\ldots \tag{2.3}
\end{equation*}
$$

the series being convergent a.s.
Lemma 2.2 [Grintsevichyus (1974), Th.3]. Assume (2.2) and let $\mu$ be the unique solution of (2.1). Then there are only three possible cases: (a) $\mu$ is absolutely continuous; (b) $\mu$ is singularly continuous; (c) $\mu$ is concentrated at some point $c$. The case (c) holds if and only if $P(c=A c+B)=1$.

Lemma 2.3.Assume (2.2) and let $X$ be the unique solution of (2.1). Let $p$ be any fixed number in $(0, \infty)$. If

$$
\begin{equation*}
E\left(|A|^{p}\right)<1 \text { and } E\left(|B|^{p}\right)<\infty \tag{2.4}
\end{equation*}
$$

then $E\left(|X|^{p}\right)<\infty$; the converse holds if additionally $A \geq 0$ and $B \geq 0$ a.s. with $P(B>0)>0$.

Proof. We denote by $\|\cdot\|_{p}$ the norm in $L^{p}, p>0$. By Lemma 2.1, we can suppose that $X$ is given by (2.3). If $\|A\|_{p}<1$ and $\|B\|_{p}<\infty$, then by (2.3) and the triangular inequality in $L^{p}$, we obtain

$$
\|X\|_{p} \leq\|B\|_{p}+\|A\|_{p}\|B\|_{p}+\ldots+\|A\|_{p}^{n-1}\|B\|_{p}+\ldots=\frac{\|B\|_{p}}{1-\|A\|_{p}}<\infty
$$

Conversely if $A \geq 0$ and $B \geq 0$ a.s. and if $X$ is a solution of (2.1) with $E\left(|X|^{p}\right)<\infty$, then $X \geq 0$ by (2.3), and $E\left(B^{p}\right) \leq E\left[(A X+B)^{p}\right]<\infty$ by (2.1); since $P(B>0)>0$, we have also $P(X>0)>0,0<E\left(X^{p}\right)<\infty$ and

$$
E\left[(A X)^{p}\right]<E\left[(A X+B)^{p}\right]=E\left(X^{p}\right)
$$

which implies $E\left[(A)^{p}\right]<1$ by the independence of $A$ and $X$.
Lemma 2.4 [Kesten(1973)- [Grintsevichyus (1975)]. Assume that $P(A>0)=P(B \geq 0)=1, P(B>0)>0$, and that for some $\lambda \in(0, \infty)$,

$$
E\left(A^{\lambda}\right)=1, E\left(A^{\lambda} \log ^{+} A\right)<\infty \text { and } E\left(B^{\lambda}\right)<\infty
$$

Let $X$ be the unique solution of (2.1) and suppose that $X$ is not a.s. a constant. [That is, there does not exist a constant $c$ such that $P(c=A c+$ $B)=1$.] Then
(a) if $\log A$ is not of lattice type, then

$$
\begin{equation*}
\lim _{x \rightarrow \infty} x^{\lambda} P(X>x)=c \tag{2.5}
\end{equation*}
$$

where $c \in(0, \infty)$ is a constant.
(b) if $\log A$ is of lattice type with span $h>0$, then for all real $x$,

$$
\begin{equation*}
\lim _{x \rightarrow \infty} e^{(x+n h) \lambda} P\left(X>e^{(x+n h) \lambda}\right)=c(x) \tag{2.6}
\end{equation*}
$$

where $c(x) \in(0, \infty), x \in R$, is a strictly positive and $h$-periodic function on R. In particular,

$$
\begin{equation*}
0<\liminf _{x \rightarrow+\infty} x^{\lambda} P(X>x) \leq \limsup _{x \rightarrow+\infty} x^{\lambda} P(X>x)<\infty \tag{2.7}
\end{equation*}
$$

## 3. The random difference equation satisfied by $x P_{Z}(d x)$; moments and tails.

For a random variable $X$, we write $P_{X}$ or $L(X)$ for its law.
For $u \in \mathrm{U}$, we use the notations $A_{u}, X_{u}$ and $Z_{u}$ introduced in section 1. For $(\omega, \mathbf{i}) \in \Omega \times \mathbf{I}$, put

$$
\begin{aligned}
\tilde{Z}(\omega, \mathrm{i}) & =Z_{\emptyset}(\omega)=Z(\omega) \\
\tilde{A}_{1}(\omega, \mathbf{i}) & =A_{\mathrm{i} \mid 1}(\omega), \\
\tilde{Z}_{1}(\omega, \mathrm{i}) & =Z_{\mathbf{i} \mid 1}(\omega) .
\end{aligned}
$$

They are measurable functions on $\Omega \times \mathbf{I}$ associated with the product $\sigma$-field of $\mathcal{F}$ and $\mathcal{B}, \mathcal{B}$ being the Borel $\sigma$-field on I. Let $Q$ be the Peyrière's measure on $\Omega \times \mathbf{I}$, defined by

$$
Q(A)=E \int_{\mathbf{I}} 1_{A}(\omega, \mathbf{i}) \mu_{\omega}(d \mathbf{i}), A \in \mathcal{F} \times \mathcal{B}
$$

As $E Z=1, Q$ is a probability measure. We write $E_{Q}[f]$ for the integration of $f$ with respect to $Q$. The following result give the distributions of the
random variables $\tilde{Z},\left(\tilde{A}_{1}, \tilde{B}_{1}\right)$, and shows that the probability law $x P_{Z}(d x)$ satisfies a random difference equation.

Lemma 3.1. For all $(\omega, i) \in \Omega \times I$,

$$
\tilde{Z}(\omega, \mathbf{i})=\tilde{A}_{1}(\omega, \mathbf{i}) \tilde{Z}_{1}(\omega, \mathbf{i})+\tilde{B}_{1}(\omega, \mathbf{i}),
$$

where

$$
\tilde{B}_{1}(\omega, \mathbf{i})=\sum_{i=1}^{N(\omega)} A_{i}(\omega) Z_{i}(\omega) 1\{\mathbf{i} \mid 1 \neq i\}
$$

$\tilde{Z}_{1}$ is independent of $\left(\tilde{A}_{1}, \tilde{B}_{1}\right)$, and has the same distribution as $\tilde{Z}$. Moreover, for all non-negative Borel functions $f, g$ and $h$ (defined on $\mathbb{R}$ or $\mathbb{R}^{2}$ ),

$$
\begin{aligned}
E_{Q}\left[f\left(\tilde{A}_{1}\right)\right] & =E\left[\sum_{i=1}^{N} f\left(A_{i}\right) A_{i}\right], \\
E_{Q}\left[f\left(\tilde{B}_{1}\right)\right] & =E\left[\sum_{1 \leq k \leq N} A_{k} f\left(\sum_{1 \leq i \leq N, i \neq k} A_{i} Z_{i}\right)\right], \\
E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right)\right] & =E\left[\sum_{1 \leq k \leq N} A_{k} h\left(A_{k}, \sum_{1 \leq i \leq N, i \neq k} A_{i} Z_{i}\right)\right], \\
E_{Q}\left[g\left(\tilde{Z}_{1}\right)\right] & =E[g(Z) Z]=E_{Q}[g(\tilde{Z})] .
\end{aligned}
$$

In particular,

$$
L(\tilde{Z})=x P_{Z}(d x)
$$

Proof. We have

$$
\begin{aligned}
A_{\mathbf{i} \mid 1} Z_{\mathbf{i} \mid 1} & =A_{i} Z_{i} \text { if } \mathbf{i} \mid 1=i \\
& =\sum_{i=1}^{\infty} A_{i} Z_{i} 1\{\mathbf{i} \mid 1=i\} \\
& =\sum_{i=1}^{N} A_{i} Z_{i} 1\{\mathbf{i} \mid 1=i\} \\
Z & =\sum_{i=1}^{N} A_{i} Z_{i} \\
& =\sum_{i=1}^{N} A_{i} Z_{i}[1\{\mathbf{i} \mid 1=i\}+1\{\mathbf{i} \mid 1 \neq i\}]
\end{aligned}
$$

$$
=A_{\mathbf{i} \mid 1} Z_{\mathbf{i} \mid 1}+\sum_{i=1}^{N} A_{i} Z_{i} 1\{\mathbf{i} \mid 1 \neq i\} .
$$

Therefore

$$
\tilde{Z}(\omega, \mathbf{i})=\tilde{A}_{1}(\omega, \mathbf{i}) \tilde{Z}_{1}(\omega, \mathbf{i})+\tilde{B}_{1}(\omega, \mathbf{i}) .
$$

We claim that $\tilde{Z}_{1}$ is independent of $\left(\tilde{A}_{1}, \tilde{B}_{1}\right)$, and has the same distribution as $\tilde{Z}$; by the way, we shall give the distribution of $\left(\tilde{A}_{1}, \tilde{B}_{1}\right)$. In fact, for all non-negative Borel functions $g$ and $h$, we have

$$
\begin{aligned}
E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right) g\left(\tilde{Z}_{1}\right)\right] & =E\left[\sum_{1 \leq k \leq N} h\left(A_{k}, \sum_{1 \leq i \leq N, i \neq k} A_{i} Z_{i}\right) g\left(Z_{k}\right) A_{k} Z_{k}\right] \\
& =E\left[\sum_{1 \leq k \leq N} A_{k} h\left(A_{k}, \sum_{1 \leq i \leq N, i \neq k} A_{i} Z_{i}\right)\right] E[g(Z) Z] .
\end{aligned}
$$

Taking $g=1$ or $h=1$ gives the expressions of $E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right)\right]$ and $E_{Q}\left[g\left(\tilde{Z}_{1}\right)\right]$. Consequently

$$
E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right) g\left(\tilde{Z}_{1}\right)\right]=E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right)\right] E_{Q}\left[f\left(\tilde{Z}_{1}\right)\right] .
$$

This gives the independence of $\left(\tilde{A}_{1}, \tilde{B}_{1}\right)$ and $\tilde{Z}_{1}$. The expressions of $E_{Q}\left[f\left(\tilde{A}_{1}\right)\right]$ and $E_{Q}\left[f\left(\tilde{B}_{1}\right)\right]$ come from $E_{Q}\left[h\left(\tilde{A}_{1}, \tilde{B}_{1}\right)\right]$ by taking $h(x, y)=f(x)$ or $f(y)$. So the proof is finished.

Lemma 3.2. For all $p>1$,

$$
E_{Q}\left[\left(\tilde{B}_{1}\right)^{p-1}\right] \leq E\left[Z^{p-1}\right] E\left[\left(\sum_{k=1}^{N} A_{k}\right)^{p}\right] .
$$

Proof. By lemma 1,

$$
E_{Q}\left[\left(\tilde{B}_{1}\right)^{p-1}\right]=E\left[\sum_{1 \leq k \leq N} A_{k}\left(\sum_{1 \leq i \leq N, i \neq k} A_{i} Z_{i}\right)^{p-1}\right] .
$$

Denote by $I$ the expectation above.

If $p-1 \leq 1$, then

$$
\begin{aligned}
I \leq & E\left\{\sum_{1 \leq k \leq N} A_{k}\left[\sum_{1 \leq i \leq N, i \neq k}\left(A_{i} Z_{i}\right)^{p-1}\right]\right\} \\
& \quad\left(\text { by the inequality }\left(\sum x_{k}\right)^{p-1} \leq \sum x_{k}^{p-1}\right) \\
= & E\left[Z^{p-1}\right] E\left[\sum_{1 \leq k \leq N} A_{k}\left(\sum_{1 \leq i \leq N, i \neq k} A_{i}^{p-1}\right)\right]
\end{aligned}
$$

(by the branching property)

$$
\begin{aligned}
& \leq E\left[Z^{p-1}\right] E\left[\sum_{k=1}^{\infty} A_{k} \sum_{k=1}^{\infty} A_{k}^{p-1}\right] \\
& \leq E\left[Z^{p-1}\right]\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p}\left\|\sum_{k=1}^{\infty} A_{k}^{p-1}\right\|_{p /(p-1)}
\end{aligned}
$$

(by Hölder's inequality)
$\leq E\left[Z^{p-1}\right]\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p} \sum_{k=1}^{\infty}\left\|A_{k}^{p-1}\right\|_{p /(p-1)}$
(by the triangular inequality)
$=E\left[Z^{p-1}\right]\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p}\left(\sum_{k=1}^{\infty} E A_{k}^{p}\right)^{(p-1) / p}$
$=E\left[Z^{p-1}\right]\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p} ;$
Since $\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p} \geq\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{1}=1$, we have $\left\|\sum_{k=1}^{\infty} A_{k}\right\|_{p} \leq E\left[\left(\sum_{k=1}^{\infty} A_{k}\right)^{p}\right]$, so the desired conclusion follows in the case where $p-1 \leq 1$.

If $p-1 \geq 1$, using the inequality

$$
\left(\sum a_{i} z_{i}\right)^{p-1} \leq \sum a_{i} z_{i}^{p-1}, \quad a_{i} \geq 0, \sum a_{i}=1, z_{i} \geq 0
$$

(the convexity of the function $z \mapsto z^{p-1}$ ) for $a_{i}=A_{i} / \sum_{j \neq k} A_{j}$ and $z_{i}=$ $Z_{i}, i \neq k$, we obtain

$$
\left(\sum_{i \neq k} A_{i} Z_{i}\right)^{p-1} \leq\left(\sum_{i \neq k} A_{i}\right)^{p-2}\left(\sum_{i \neq k} A_{i} Z_{i}^{p-1}\right)
$$

(if $\sum_{j \neq k} A_{j}=0$, the inequality is evident). Consequently

$$
I \leq E\left(Z^{p-1}\right) E\left[\sum_{k=1}^{N} A_{k}\left(\sum_{i \neq k} A_{i}\right)^{p-2}\left(\sum_{i \neq k} A_{i}\right)\right]
$$

$$
\leq E\left(Z^{p-1}\right) E\left[\left(\sum_{i=1}^{N} A_{i}\right)^{p}\right]
$$

Proof of Theorem 1. By Lemma 3.1,

$$
E_{Q}\left[\tilde{Z}_{1}^{p-1}\right]=E\left[Z^{p}\right] \text { and } E_{Q}\left[\tilde{A}_{1}^{p-1}\right]=E\left[\sum_{i=1}^{N} A_{i}^{p}\right]
$$

So by lemmas 2.3 and 3.2, we see that for all $p>1$, if

$$
E\left[Z^{p-1}\right]<\infty, E\left[\left(\sum_{k=1}^{N} A_{k}\right)^{p}\right]<\infty \text { and } E\left[\sum_{k=1}^{N} A_{k}^{p}\right]<1
$$

then $E\left[Z^{p}\right]<\infty$. Noting that $E Z<\infty$, an easy induction argument on $n$ shows that, for all $n=2,3, \ldots$, if

$$
p \in(n-1, n], E\left[\left(\sum_{k=1}^{N} A_{k}\right)^{p}\right]<\infty \text { and } E\left[\sum_{k=1}^{N} A_{k}^{p}\right]<1,
$$

then $E\left[Z^{p}\right]<\infty$. This gives the sufficiency of the conditions. The necessity is given in Liu (1997).

Proof of Theorem 2. By lemmas 3.1 and 3.2, we have

$$
\begin{equation*}
E_{Q}\left[\tilde{A}_{1}^{\chi-1}\right]=E\left[\sum_{i=1}^{N} A_{i}^{\chi}\right]=1 \tag{i}
\end{equation*}
$$

$$
\begin{equation*}
E_{Q}\left[\tilde{A}_{1}^{\chi-1} \log ^{+} \tilde{A}_{1}\right]=E\left[\sum_{i=1}^{N} A_{i}^{\chi} \log ^{+} A_{i}\right]<\infty \tag{ii}
\end{equation*}
$$

and

$$
\begin{equation*}
E_{Q}\left[\tilde{B}_{1}^{\chi-1}\right] \leq E\left[Z^{x-1}\right] E\left[\left(\sum_{i=1}^{N} A_{i}\right)^{\chi}\right] \tag{iii}
\end{equation*}
$$

If the problem is non-lattice, then by Lemma 2.4, the limit

$$
\lim _{t \rightarrow+\infty} t^{-(x-1)} \int_{t}^{\infty} x P_{Z}(d x)=\lim _{t \rightarrow+\infty} t^{-(x-1)} Q(\tilde{Z}>t)
$$

exists and is strictly positive and finite. This implies that the limit

$$
\lim _{t \rightarrow+\infty} t^{-x} P(Z>t)
$$

exists and is strictly positive. In the lattice case, the corresponding conclusions are

$$
0<\liminf _{t \rightarrow+\infty} t^{-(x-1)} Q(\tilde{Z}>t) \leq \limsup _{t \rightarrow+\infty} t^{-(x-1)} Q(\tilde{Z}>t)<\infty
$$

and

$$
0<\liminf _{t \rightarrow+\infty} t^{x} P(Z>t) \leq \limsup _{t \rightarrow+\infty} t^{x} P(Z>t)<\infty
$$

This gives part (a) of the theorem. Part (b) has been proved in Liu (1996).

## 4. The absolute continuity

The discussion will be heavily based on the functional equation (1.1'). We recall that $\tilde{N}=\sum_{i=1}^{N} 1\left\{A_{i}>0\right\}$ is the number of non-zero terms of $\left\{A_{i}: 1 \leq\right.$ $i \leq N\}$. Let

$$
f(x):=\sum_{i=1}^{N} P(\tilde{N}=i) x^{i}, x \in[0, \infty)
$$

be its probability generating function.
Lemma 4.1. Assume $\tilde{N} \geq 1$ a.s. Let $Z$ be any solution of (1.1), and let $\phi(t)=E\left(e^{i t Z}\right)(t \in R)$ be its characteristic function. Then

$$
\limsup _{|t| \rightarrow \infty}|\phi(t)|=0 \text { or } 1
$$

Proof. By (1.1'),

$$
|\phi(t)| \leq E \prod_{i=1}^{N}\left|\phi\left(A_{i} t\right)\right|, \quad t \in R .
$$

Letting $|t| \rightarrow \infty$ and using Fatou's lemma gives

$$
l \leq f(l)
$$

where $l:=\limsup _{|t| \rightarrow \infty} \mid \phi(t)$. Therefore $l=0$ or 1 , noting that $f(x)<x$ if $0<x<1$.

Lemma 4.2. Write

$$
M=\max _{1 \leq i \leq N} A_{i}
$$

and suppose that

$$
E|\log M|<\infty
$$

Let $\phi$ be a solution of (1.1') in the class of characteristic functions, which is not of the form $e^{i t c}$ for some constant $c$. Then

$$
\limsup _{|t| \rightarrow \infty}|\phi(t)|=0
$$

Proof. By Lemma 4.1, it suffices to prove that $\lim \sup _{|t| \rightarrow \infty}|\phi(t)|<1$.
(a) We first prove that for all $t \neq 0,|\phi(t)|<1$. Otherwise, by Lemma 4 of Chap.IV. 1 of Feller, there is some $h>0$ such that $|\phi(h)|=1$ and $|\phi(t)|<1$ if $0<t<h$. By the equation (1.1'),

$$
1=|\phi(h)| \leq E \prod_{i=1}^{N} \mid \phi\left(A_{i} h \mid\right.
$$

Therefore, a.s.

$$
\left|\phi\left(A_{i} h\right)\right|=1 \text { for all } i=1, \ldots, N
$$

Since $P(0<M<1)>0$ (this is necessary for the equation (1.1) to have a non-trivial solution), it follows that for some $0<a<1,|\phi(a h)|=1$, which is a contradiction with the definition of $h$.
(b) We then prove that $\lim \sup _{|t| \rightarrow \infty}|\phi(t)|<1$. By the functional equation, we have

$$
|\phi(t)| \leq E|\phi(M t)| .
$$

We now use ideas of Grintsevichyus (1974) on random walks. Let $M_{i}(i \geq 1)$ be independent copies of $M$, then

$$
|\phi(t)| \leq E\left|\phi\left(M_{1} \ldots M_{n} t\right)\right|, \quad n \geq 1
$$

Fix $t \neq 0$, and write

$$
\phi_{0}=|\phi(t)|, \quad \phi_{n}=\left|\phi\left(M_{1} \ldots M_{n} t\right)\right|, \quad n \geq 1
$$

Then $\left\{\phi_{n}: n \geq 0\right\}$ is a sub-martingale associated with the natural filtration of $\sigma$-fields generated by $\left\{M_{i}: 0 \leq i \leq n\right\}, n \geq 0$, where by convention $M_{0}=1$. Put

$$
S_{0}=1 \text { and } S_{n}=\log M_{1}+\ldots+\log M_{n} \text { for } n \geq 1
$$

If $I$ is an interval, we write

$$
-(I)=\tau_{1}(I)=\inf \left\{n>0: S_{n} \in I\right\}, \text { where } \inf \emptyset=+\infty,
$$

for the first time that the random walk $\left(S_{n}\right)$ hits $I$, and put

$$
\begin{gathered}
\tau_{k+1}(I)=\inf \left\{n>\tau_{k}: S_{n} \in I\right\}, \text { where } \inf \emptyset=+\infty, \quad k \geq 1, \\
U(I)=E\left(\sum_{n=0}^{\infty} 1\left\{S_{n} \in I\right\}\right) .
\end{gathered}
$$

Then

$$
P\left[\tau_{k}(I)<\infty\right]=P\left[\sum_{n=0}^{\infty} 1\left\{S_{n} \in I\right\} \geq k\right], \quad k \geq 1
$$

is the probability that there are at least k hits in $I, \mathrm{U}(\mathrm{I})$ is the expected number of hits in $I$, and

$$
U(I)=\sum_{k=1}^{\infty} P\left[\tau_{k}(I)<\infty\right] .
$$

By the strong Markov's property, it is easily verified that for all $-\infty<a<b<\infty$,

$$
P\left[\tau_{k+1}([a, b])<\infty\right] \leq P[\tau([a, b])<\infty] P\left[\tau_{k}([a-b, b-a])<\infty\right], \quad k \geq 1
$$

Summing for $\mathrm{k}=1,2, \ldots$, we obtain

$$
U([a, b]) \leq P[\tau([a, b])<\infty][1+U([a-b, b-a])] .
$$

Therefore, for all $h>0$,

$$
U([-h, 0]-\log |t|) \leq P[\tau([-h, 0]-\log |t|)<\infty][1+U([-h, h])], \quad k \geq 1
$$

On the other hand, by the renewal theorem (cf. Feller Chap. 11.9), there is $h>0$ such that

$$
\lim _{|t| \rightarrow \infty} U([-h, 0)-\log |t|)=\frac{h}{-E \log M}>0
$$

(We remark that $E[\log M]<0$ because $E[\log M]=\alpha^{-1} E\left[\log M^{\alpha}\right]<\alpha^{-1}[\log \gamma(\alpha)]=0$, where $\alpha$ is defined in Theorem 0.) So there are numbers $\delta>0$ and $t_{0}>0$ such that

$$
U([-h, 0)-\log |t|) \geq \delta \quad \text { for all }|t|>t_{0}
$$

It follows that for all $|t|>t_{0}$,

$$
P\left[\tau \left([-h, 0]-\log |t| \geq \frac{U([-h, 0)-\log |t|}{1+U([-h, h])} \geq \frac{\delta}{1+U([-h, h])}=\delta_{1}>0\right.\right.
$$

Therefore, writing $\tau=\tau([-h, 0]-\log |t|)$ and

$$
\gamma=\max _{|x| \in\left[e^{-h}, 1\right]}|\phi(x)|
$$

( $\gamma<1$ by part (a) above), and using the stopped martingale theorem, we obtain that, for all $|t|>t_{0}$,

$$
\begin{aligned}
|\phi(t)| & \leq \limsup _{n \rightarrow \infty} E\left|\phi\left(e^{S_{n \wedge \tau} t}\right)\right| \quad(n \wedge \tau=\min (n, \tau)) \\
& \leq E\left[\left|\phi\left(e^{S_{\tau} t}\right)\right| 1\{\tau<\infty\}\right]+P(\tau=\infty) \\
& \leq \gamma P(\tau<\infty)+(1-P(\tau<\infty)) \\
& \leq 1-(1-\gamma) \delta_{1} .
\end{aligned}
$$

So $\lim \sup _{|t| \rightarrow \infty}|\phi(t)|<1$, as desired.
Lemma 4.3. If $\tilde{N} \geq 2$ almost surely and

$$
E\left[A_{\zeta_{1}}^{-a}+A_{\zeta_{2}}^{-a}\right]<\infty
$$

for some $a>0$, then

$$
\phi(t)=O\left(|t|^{-a}\right), \quad|t| \rightarrow \infty
$$

where $\phi(t)=E\left[e^{i t Z}\right], t \in \mathbb{R}$.
Proof. Write

$$
\psi(t)=\sup _{|x| \geq t}|\phi(t)|, \quad t>0
$$

Then $\psi(t)$ is non-increasing, $\psi(0)=1$, and $\lim _{t \rightarrow \infty} \psi(t)=0$ by Lemma 4.2. (We remark that the conditions of lemma imply that $E\left[M^{-a}\right] \leq E\left[A_{\zeta_{1}}^{-a}\right]<\infty$, so that $E[|\log M|]<\infty$.) By the functional equation (1.1'), it is easily seen that

$$
\begin{equation*}
\psi(t) \leq E \prod_{i=1}^{N} \psi\left(A_{i} t\right) \leq E\left[(\psi(A t))^{2}\right] \tag{*}
\end{equation*}
$$

where $A=\min \left(A_{\zeta_{1}}, A_{\zeta_{2}}\right)$, and the last inequality holds since $\tilde{N} \geq 2$ almost surely. As $P(A \leq x) \leq E\left[A^{-a}\right] x^{a} \leq E\left[A_{\zeta_{1}}^{-a}+A_{\zeta_{2}}^{-a}\right] x^{a}$ by Markov's inequality, it follows that for $K=E\left[A_{\zeta_{1}}^{-a}+A_{\zeta_{2}}^{-a}\right]$ and all $0 \leq u \leq 1$ and $t>0$,

$$
\begin{aligned}
\psi(t) & \leq P[A \leq u]+(\psi(u t))^{2} \\
& \leq K u^{a}+(\psi(u t))^{2} \\
& \leq\left[(K u)^{a / 2}+\psi(u t)\right]^{2}
\end{aligned}
$$

Therefore, by Lemma 1.4.1.a of Barral (1997), for all $b<a / 2$,

$$
\psi(t)=O\left(t^{-b}\right), \quad t \rightarrow \infty
$$

Fix $b<a / 2$. So for some constant $C>0$ and all $t>0, \psi(t) \leq C t^{-a / 2}$. So $\psi(A t) \leq C(A t)^{-b}$ and, by the inequality $\left(^{*}\right)$, for all $t>0$,

$$
\psi(t) \leq C^{2} E A^{-2 b}\left(t^{-2 b}\right)
$$

The same argument applies for $b_{1}=2 b$ and $C_{1}=C^{2} E A^{-2 b}$, yielding that for all $t>0$,

$$
\psi(t) \leq C_{1}^{2} E A^{-2 b_{1}}\left(t^{-2 b_{1}}\right)
$$

Taking $b=a / 4$ gives the result desired.
Lemma 4.4. Assume $\tilde{N} \geq 2$ almost surely and put $M_{i}=\max _{j \neq i} A_{j}$. Let $\phi$ be a non-trivial solution of (1.1'). If for some $\epsilon>0, \phi(t)=O\left(|t|^{-\epsilon}\right),|t| \geq 1$,

$$
E \sum_{i=1}^{N} M_{i}^{-\epsilon}<\infty \text { and } E \sum_{i=1}^{N} A_{i} M_{i}^{-\epsilon}<\infty
$$

then for some $\delta>0$,

$$
\phi^{\prime}(t)=O\left(|t|^{-(1+\delta)}\right), \quad|t| \rightarrow \infty .
$$

Proof. From the functional equation (1.1'), we obtain

$$
\phi^{\prime}(t)=E\left[\sum_{i=1}^{N} A_{i} \phi^{\prime}\left(A_{i} t\right) \prod_{j \neq i} \phi\left(A_{j} t\right)\right]
$$

Since $|\phi(t)| \leq 1$, it follows that $\left|\prod_{j \neq i} \phi\left(A_{j} t\right)\right| \leq \phi\left(M_{i} t\right)$ and

$$
\begin{equation*}
\left|\phi^{\prime}(t)\right| \leq E\left[\sum_{i=1}^{N} A_{i}\left|\phi^{\prime}\left(A_{i} t\right)\right|\left|\phi\left(M_{i} t\right)\right|\right] . \tag{**}
\end{equation*}
$$

Because $\left|\phi^{\prime}(t)\right| \leq 1$ and

$$
|\phi(t)| \leq K|t|^{-\epsilon}
$$

for some constant $K>0$ and all $t \in \mathbb{R}^{*}$, the preceding inequality implies that

$$
\left|\phi^{\prime}(t)\right| \leq E\left[\sum_{i=1}^{N} A_{i} K\left|M_{i} t\right|^{-\epsilon}\right] \leq K_{1}|t|^{-\epsilon},
$$

where $K_{1}=K E\left[\sum_{i=1}^{N} A_{i} M_{i}^{-\epsilon}\right]<\infty$.
If $\epsilon \leq 1$, using the inequalities $\left|\phi^{\prime}\left(A_{i} t\right)\right| \leq K_{1}\left|A_{i} t\right|^{-\epsilon}, \phi\left(M_{i} t\right) \leq K\left|M_{i} t\right|^{-\epsilon}$ and (**), we see that

$$
\left|\phi^{\prime}(t)\right| \leq E\left[\sum_{i=1}^{N} A_{i} K_{1}\left|A_{i} t\right|^{-\epsilon} K\left|M_{i} t\right|^{-\epsilon}\right]=K_{2}|t|^{-2 \epsilon}, \quad t \in \mathbb{R}^{\star}
$$

where $K_{2}=K K_{1} E\left[\sum_{i=1}^{N} A_{i}^{1-\epsilon} M_{i}^{-\epsilon}\right]<\infty$. (We remark that the moment conditions in the lemma implies that $E\left[\sum_{i=1}^{N} A_{i}^{x} M_{i}^{-\epsilon}\right]<\infty$ for all $x \in[0,1]$.)

If $2 \epsilon \leq 1$, we continue the procedure, and so on. If $n=\max \{k \geq 0: k \epsilon \leq$ $1\}$, then $n \epsilon \leq 1$ but $(n+1) \epsilon>1$, and the argument as above shows that for some constant $0<K_{n+1}<\infty$ and all $t \in \mathbb{R}^{\star}$,

$$
\left|\phi^{\prime}(t)\right| \leq K_{n+1}|t|^{-(n+1) \epsilon} .
$$

Proof of Theorem 3. By Lemma 2.2, $L(\tilde{Z})$ is either absolutely continuous or singularly continuous on $(0, \infty)$. Since $L(\tilde{Z})=x P_{Z}(d x)$, the same is true for the distribution of $Z$.

Let $\phi$ be the characteristic function of $Z$.
If (i) holds, then by lemmas 4.3 and 4.4 ,

$$
\phi^{\prime}(t)=O\left(|t|^{-(1+\delta)}\right), \quad|t| \rightarrow \infty
$$

for some $\delta>0$. Since $-i \phi^{\prime}(t)$ is the characteristic function of the probability measure $x P_{Z}(d x)$, by the inverse formula of Fourier transform, the measure $x P_{Z}(d x)$ has a continuous density on $\mathbb{R}$, so $P_{Z}(d x)$ has a continuous density on $(0, \infty)$.

Assume that (ii) holds. Without loss of generality, we can suppose that $N \geq 1$ and $A_{1}>0$ if $1 \leq i \leq N$ a.s. (Otherwise, we consider ( $\tilde{N}, A_{\zeta_{1}}, \ldots, A_{\zeta_{\tilde{N}}}$ ) instead of ( $N, A_{1}, \ldots, A_{N}$ ); see the discussion before the statement of Theorem
3.) Therefore $A_{1}=A_{\varsigma_{1}}$ is absolutely continuous on $\mathbb{R}$ and $P\left(Z_{1}=0\right)=0$. So by the lemma at page 166 of Grintsevichyus (1974), $L(Z)=L\left(A_{1} Z_{1}+B\right)$ is absolutely continuous on $\mathbb{R}$, where $B=\sum_{2 \leq i \leq N} A_{i}$. (The sum is taken to be zero if $N<2$.)

Acknowledgment. The author is very grateful to Yves Guivarc'h for many valuable discussions.

## References

Athreya, K.B. (1971) : A note on a functional equation arising in GaltonWatson branching processes. J. Appl. Prob., 8, 589-598.

Athreya, K.B. and Ney, P.E. (1972) : Branching processes. Springer, Berlin.
Barral J. (1997): Continuté, Moments d'ordres négatifs et Analyse Multifractale de cascades multiplicatives de Mandelbrot. Thèse, Univ. Paris-Sud, Orsay.

Ben Nasr, F. (1987) : Mesures aléatoires de Mandelbrot associées à des substitutions. CRAS, Sér. I, 304, 255-258.

Biggins, J.D. (1977) : Martingale convergence in the branching random walk. J. Appl. Prob. 14, 25-37.

Biggins, J.D. and Bingham, N.H. (1993) : Large deviations in the supercritical branching process. Adv. Appl. Prob. 25, 757-772.

Bingham, N.H. and Doney,R.A. (1974): Asymptotic properties of supercritical branching processes I: The Galton-Watson process. Adv. Appl. Prob., 6, 711-731.

Bingham, N.H. and Doney,R.A. (1975): Asymptotic properties of supercritical branching processes II: Crump-Mode and Jirina processes. Adv. Appl. Prob., 7, 66-82.

Bingham, N.H., Goldie,C.M. and Teugels,J.L. (1987) : Regular variation. Cambridge: Cambridge University Press.

Chauvin, B. and Rouault, A. (1996) Boltzmann-Gibbs weights in the branching random walk. IMA Congress on Branching Processes $n^{\circ} 84$. Lecture Notes in Maths.

Collet, P. and Koukiou, F. (1992) : Large deviations for multiplicative chaos. Comm. Math. Phys. 147, 329-342.

Crump, K. and Mode, C.J. (1968-1969) : A general age-dependent branching process J. Math.Anal. Appl., 24, 497-508 and 25, 8-17.

Doney, R.A. (1972) : A limit theorem for a class of supercritical branching processes. J. Appl. Prob., 9, 707-724 .

Doney, R.A. (1973) : On a functional equation for general branching processes. J. Appl. Prob., 10, 497-508 .

Durrett, R. and Liggett, T. (1983) : Fixed points of the smoothing transformation. Z. Wahrsch. verw. Gebeite, 64, 275-301.

Falconer, K.J. (1986) : Random fractals. Math.Proc.Camb.Phil.Soc., 100, 559-582.

Falconer, K.J. (1987) : Cut set sums and tree processes. Proc. Amer. Math. Soc. (2), 101, 337-346.

Feller, W. (1970): An introduction to probability theory and its applications, vol.II, 2nd. ed. John Wiley \& Sons, New York.

Graf, S. Mauldin, R.D. and Williams, S.C. (1988) : The exact Hausdorff dimension in random recursive constructions. Mem. Amer. Math. Soc. 71, 381.

Grintsevichyus,A.K.(1974): On the continuity of a sum of dependent variables connected with independent walks on lines. Theory Prob. Appl. 19, 163-168.

Grintsevichyus (1975): One limit distribution for a random walk on the line. Lithunian Math. Trans. 15, 580-589.

Guivarc'h, Y. (1990) : Sur une extension de la notion de loi semi-stable. Ann. IHP 26, 261-285.

Harris, T.E. (1948) : Branching processes, Ann. Math. Stat. 19, 474-494.
Holley, R. and Liggett, T. (1981) : Generalized potlach and smoothing processes. Z. Wahrsch. verw. Gebeite, 55, 165-195.

Holley, R. and Waymire, E.C. (1992) : Multifractal dimensions and scaling exponents for strongly bounded cascades. Ann. Appl. Prob. 2, 819-845.

Kahane, J.P. (1987) : Multiplications aléatoires et dimension de Hausdorff. Ann. IHP, Sup. au no. 2, vol. 23, 289-296 .

Kahane, J.P. and Peyrière (1976) : Sur certaines martingales de Benoit Mandelbrot. Adv. Math., 22, 131-145.

Kesten,H. (1973): Random difference equations and renewal theory for products of random matrices. Acta Math.131, 207-248.

Kesten,H and Stigum,B.P.(1966): A limit theorem for multidimensional GaltonWatson processes. Ann.Math.Statist.37, 1211-1223.

Liu, Q. (1993) : Sur quelques problèmes à propos des processus de branchement, des flots dans les réseaux et des mesures de Hausdorff associées. Thèse, Université Paris 6.

Liu, Q. (1996): The exact Hausdorff dimension of a branching set. Prob. Th. Rel. Fields, 104; 1996, 515-538.

Liu, Q. (1996): The growth of an entire characteristic function and the tail probabilities of the limit of a tree martingale. In Trees. Progress in Probability, vol.40, pp 51-80, 1996, Birkhöuser: Verlag Basel. Eds.: B.Chauvin, S.Cohen, A.Rouault.

Liu, Q. (1997a) : Sur une équation fonctionnelle et ses applications : une extension du théorème de Kesten-Stigum concernant des processus de branchement. To appear in Adv. Appl. Prob. 1997.

Liu, Q. (1997b) : Fixed points of a generalized smoothing transformation and applications to branching processes. To appear in Adv. Appl. Prob. 1997.

Liu, Q and Rouault, A.(1996): On two measures defined on the boundary of a branching tree. In "Classical and modern branching processes", ed. K.B. Athreya and P. Jagers. IMA Volumes in Mathematics and its applications, vol. 84, 1996, pp.187-202. Springer-Verlag.

Mandelbrot, B. (1974) : Multiplications aléatoires et distributions invariantes par moyenne pondérée aléatoire. CRAS Paris vol. 278, 325-346 et 355-358.

Mauldin, R.A. and Williams, S.C. (1986) : Random constructions, asymptotic geometric and topological properties. Trans. Amer. Math. Soc. 295, 325-346.

Seneta, E. (1968) : On recent theorems concerning the supercritical GaltonWatson process. Ann. Math. Stat. 39, 2098-2102.

Seneta, E. (1969) : Functional equations and the Galton-Watson process. Adv. Appl. Prob. 1, 1-42.

Waymire, E.C. and Williams, S.C. (1995) : Multiplicative cascades : dimension spectra and dependence. J. of Fourier Analysis and Appl., Kahane Special Issue.

