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# AR(1) PROCESSES WITH GIVEN MOMENTS OF MARGINAL DISTRIBUTION 

JIŘí ANDĚL

Let $X_{t}$ be an $\operatorname{AR}(1)$ process given by $X_{t}=b X_{t-1}+e_{t}$ where $b \in(-1,1)$ and $e_{t}$ is a strict white noise. Sometimes $X_{t}$ must satisfy also some additional conditions, e.g. $X_{t} \geqq 0$ or $C \leqq$ $\leqq X_{t} \leqq D$. The problem solved in the paper is how to find a distribution of $e_{t}$ such that the moments $\mathrm{E} X_{t}^{k}(k=1, \ldots, n)$ have given values.

## 1. HISTORY OF THE PROBLEM

Simulation procedures are frequently used for demonstrating theoretical results as well as for modelling real situations. In some cases it is necessary to generate pseudorandom numbers not only with a given marginal distribution but also with a given covariance structure. Chamitov [7] writes that the problem was formulated already in 1963. It is particularly important in simulating dynamical systems and in hydrological applications. The first theoretical results concerning an AR (1) process were published by Gaver and Lewis [8]. A review of methods of this kind can be found in [2]. (See also Section 2 of the present paper.) The disadvantage of this method is that the procedure gives explicit results only in a few special cases. Moreover, it cannot be generalized to the AR models of higher order.
Two direct methods were published in 1983. Chamitov [7] publishes some empirical tables. It is not clear, however, which of the tabulated covariance functions should be chosen when one has a given empirical covariance function calculated from a sample or when a covariance function is given analytically. Further, the class of one-dimensional distributions as well as the class of covariance functions is rather limited. The method is based on the following argumentation. If $\xi_{1}, \ldots, \xi_{k}$ is a sequence of independent identically distributed random variables, then each of them has, of course, the same distribution. But the members of the ordered sample have distributions different from this one. This point was obviously overlooked.

Sondhi [13] starts with a Gaussian white noise $e_{t}$ and with a linear filter $H$. The
output $Z_{t}$ has a normal distribution function $G_{0}$ instead of the distribution function $F$, which is wanted. If we put $X_{t}=F^{-1}\left[G_{0}\left(Z_{t}\right)\right]$ then it is clear that $X_{t}$ has the distribution function $F$. Unfortunately, after the application of the non-linear filter $F^{-1} G_{0}$ the process $X_{t}$ has a covariance function which is different from that belonging to the filter $H$. The procedure suggested by the author is to use Mehler's expansion of the bivariate normal density. It leads to an approximation and the computation is rather complicated. Moreover, there is no possibility how to recognize whether the problem has a solution. Sondhi writes: "Unfortunately, to the best. of our knowledge, no tractable procedure is known to decide whether a covariance function is consistent with a given probability distribution function."

For solving the problem a financial support was offered (see [11], p. 193).
In practical situations, the given one-dimensional stationary distribution is usually derived from some of its moments. Anděl [3] proposed a method how to find a distribution of the white noise such that the resulting linear process $X_{t}$ has given moments $E X_{t}^{k}(k=1, \ldots, n)$. This approach was generalized to some non-linear processes in [4] and applied to real hydrological data in [6].

## 2. INTRODUCTION

Consider an $\operatorname{AR}(1)$ process $X_{t}$ given by

$$
\begin{equation*}
X_{t}=b X_{t-1}+e_{t} \tag{2.1}
\end{equation*}
$$

where $b \in(-1,1)$ and $e_{t}$ are independent random variables with the same distribution function $G$. The problem is to find a $G$ such that all the variables $X_{t}$ have a given stationary distribution with a distribution function $F$. The main tool for solving the problem was introduced in [8]. Let

$$
\omega(u)=\mathrm{E} \exp \left(\mathrm{i} u X_{t}\right) \text { and } \psi(u)=\mathrm{E} \exp \left(\mathrm{i} u e_{t}\right)
$$

be the characteristic function of $X_{t}$ and $e_{t}$, respectively. Since $e_{t}$ is independent of $X_{t-1}$, (2.1) yields

$$
\begin{equation*}
\omega(u)=\omega(b u) \psi(u) \tag{2.2}
\end{equation*}
$$

From here $\psi(u)$ is calculated. This procedure (based on the moment generating function instead of the characteristic function) was applied in [8] for finding $G$ such that $X_{t}$ has an exponential distribution (see Section 7 of our paper).

In the general case, it is difficult to determine $G$ from $\psi$. It happens also often that $\psi$ calculated from (2.2) is not a characteristic function, which means that for such an $F$ the problem has no solution. The following example is given in [2]. If we look for a $G$ such that $X_{t}$ have the continuous rectangular distribution on a given interval $[-c, c]$, then using (2.2) we get

$$
\psi(u)=b \sin (c u) / \sin (b c u) .
$$

This is a characteristic function only in the case that $b=1 / n$ for $n= \pm 1, \pm 2, \ldots$.

If $b=1 /(2 n)$ then $\psi$ corresponds to the distribution concentrated at the points $\pm(2 k-1) c /(2 n)$ for $k=1, \ldots, n$ where each point has probability $1 /(2 n)$. If $b=$ $=1 /(2 n+1)$ then $\psi$ is the characteristic function of the distribution concentrated at $\pm 2 k c /(2 n+1), k=0,1, \ldots, n$, where each point has probability $1 /(2 n+1)$. If $b \neq 1 / n$ then $|\psi(u)| \rightarrow \infty$ as $u \rightarrow \pi /(b c)$, and thus $\psi$ cannot be a characteristic function.

In this paper we assume that only some moments of $X_{t}$ are given. The problem if to find a distribution function $G$ (if it exists) such that (2.1) produces a stationary distribution with the given moments. Our procedure can easily be realized on a computer and it shows automatically when no solution exists.

## 3. RELATIONS FOR MOMENTS

Define $m_{k}=\mathrm{E} X_{t}^{k}, s_{k}=\mathrm{E} e_{t}^{k}$ for $k=0,1, \ldots$ It follows from (2.1) that

$$
m_{k}=\sum_{i=0}^{k}\binom{k}{i} b^{i} m_{i} s_{k-i}
$$

and from here we obtain

$$
\begin{equation*}
s_{k}=\left(1-b^{k}\right) m_{k}-\sum_{i=1}^{k-1}\binom{k}{i} b^{i} m_{i} s_{k-i} \tag{3.1}
\end{equation*}
$$

for $k=1,2, \ldots$. If the moments $m_{k}$ are given then (3.1) enables to calculate $s_{k}$ recurrently. For example, the first four moments are

$$
\begin{aligned}
& s_{1}=(1-b) m_{1} \\
& s_{2}=\left(1-b^{2}\right) m_{2}-2 b(1-b) m_{1}^{2} \\
& s_{3}=\left(1-b^{3}\right) m_{3}-3 b(1+2 b)(1-b) m_{1} m_{2}+6 b^{2}(1-b) m_{1}^{3} \\
& s_{4}=\left(1-b^{4}\right) m_{4}-4 b(1-b)\left(1+b+2 b^{2}\right) m_{1} m_{3}+12 b^{2}(1-b)(1+3 b) . \\
& \quad \cdot m_{1}^{2} m_{2}-6 b^{2}\left(1-b^{2}\right) m_{2}^{2}-24 b^{3}(1-b) m_{1}^{4}
\end{aligned}
$$

Sometimes $X_{t}$ must fulfil also other conditions. Three following cases are most important.
(i) $X_{t} \in[C, D]$ for all $t$, where $C<D$ are given numbers.
(ii) $X_{t} \geqq 0$ for all $t$.
(iii) There are no additional restrictions on $X_{t}$.

If the numbers $s_{1}, \ldots, s_{n}$ are calculated, we must decide if there exists a distribution function $G$ such that $s_{k}$ are its moments and that $X_{t}$ fulfil given restrictions.

## 4. MOMENT PROBLEM ON A FINITE INTERVAL

Theorem 4.1. Let $X_{t}$ be an $\operatorname{AR}(1)$ process defined by (2.1) with $b \in(-1,1)$. Let $C<D$ be given numbers. Define

$$
\begin{aligned}
& H=C(1-b), \quad K=D(1-b) \quad \text { for } \quad b \geqq 0, \\
& H=C-b D, \quad K=D-b C \quad \text { for } \quad b<0 .
\end{aligned}
$$

If $H \leqq e_{t} \leqq K$ for all t , then $C \leqq X_{t} \leqq D$ for all $t$. If $\mathrm{P}\left(e_{t} \notin[H, K]\right)>0$, then $X_{t} \notin[C, D]$ for infinitely many subscripts $t$ with probability 1.
Proof. Let $0 \leqq b<1$. Let $C(1-b) \leqq e_{t} \leqq D(1-b)$ hold for all $t$. Since

$$
X_{t}=e_{t}+b e_{t-1}+b^{2} e_{t-2}+\ldots
$$

we obtain

$$
\begin{aligned}
& X_{t} \leqq D(1-b)\left(1+b+b^{2}+\ldots\right)=D \\
& X_{t} \geqq C(1-b)\left(1+b+b^{2}+\ldots\right)=C
\end{aligned}
$$

Now, consider the case $-1<b<0$. We can write

$$
X_{t}=e_{t}+b^{2} e_{t-2}+b^{4} e_{t-4}+\ldots+b\left(e_{t-1}+b^{2} e_{t-3}+b^{4} e_{t-5}+\ldots\right)
$$

Thus

$$
\begin{aligned}
& X_{t} \leqq(D-b C)\left(1+b^{2}+b^{4}+\ldots\right)+b(C-a D)\left(1+b^{2}+b^{4}+\ldots\right)=D \\
& X_{t} \geqq(C-b D)\left(1+b^{2}+b^{4}+\ldots\right)+b(D-b C)\left(1+b^{2}+b^{4}+\ldots\right)=C
\end{aligned}
$$

A proof of the last assertion of Theorem 4.1 is similar to that of Lemma 10.2.
Theorem 4.2. A sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r}$ is a system of moments on an interval [ $H, K]$ if and only if both the matrices

$$
\mathbf{A}=\left(s_{i+j}\right)_{i, j=0}^{r}, \quad \mathbf{B}=\left((H+K) s_{i+j+1}-H K s_{i+j}-s_{i+j+2}\right)_{i, j=0}^{r-1}
$$

are positive semidefinite.
Proof. See [9], p. 90.
Theorem 4.3. A sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r+1}$ is a system of moments on an interval $[H, K]$ if and only if both the matrices

$$
\mathbf{A}=\left(s_{i+j+1}-H s_{i+j}\right)_{i, j=0}^{r}, \quad \mathbf{B}=\left(K s_{i+j}-s_{i+j+1}\right)_{i, j=0}^{r}
$$

are positive semidefinite.
Proof. See [9], p. 91.
Let us remark that an infinite sequence $\left\{s_{k}\right\}_{0}^{\infty}$ is a system of moments on an interval [ $H, K]$ if and only if the infinite quadratic forms

$$
\sum_{i=0}^{\infty} \sum_{j=0}^{\infty} s_{i+j} x_{i} x_{j}, \quad \sum_{i=0}^{\infty} \sum_{j=0}^{\infty}\left[(H+K) s_{i+j+1}-H K s_{i+j}-s_{i+j+2}\right] x_{i} x_{j}
$$

are positive semidefinite. This is equivalent to the condition that the forms

$$
\sum_{i=0}^{\infty} \sum_{j=0}^{\infty}\left(s_{i+j+1}-H s_{i+j}\right) x_{i} x_{j}, \quad \sum_{i=0}^{\infty} \sum_{j=0}^{\infty}\left(K s_{i+j}-s_{i+j+1}\right) x_{i} x_{j}
$$

are positive semidefinite (see [9], p. 92).
If we have numbers $s_{0}=1, s_{1}, \ldots, s_{n}$ then it can be decided using Theorems 4.2 and 4.3 if they are moments or not. If they are, it is important to know at least one distribution having just moments $s_{1}, \ldots, s_{n}$. It is known that there exists a discrete distribution with this property and the points at which the probability is concentrated
are given as the roots of some polynomials. Such points are shortly called "points of concentration". Since this method is very general and not known among the statisticians, we describe it in this paper in detail. Before doing it, let us make the following remark. In the most situations only four moments are investigated. The author believes that then the use of the generalized Tukey's $\lambda$-system (see [12]) can be recommended. The procedure gives smooth densities and the sample can be quickly generated on a computer.

Now, we return to the problem of general $n$.
Theorem 4.4. Let $n=2 r-1$. Let the matrices $\boldsymbol{A}$ and $\boldsymbol{B}$ from Theorem 4.3 for a sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r-1}$ be positive definite. Then the points of concentration of so called lower main representation are the roots of

$$
\left|\begin{array}{lllll}
s_{0} & s_{1} & \ldots & s_{r-1} & 1 \\
s_{1} & s_{2} & \ldots & s_{r} & z \\
\ldots & \cdots & \cdots & \cdots & \cdots
\end{array}\right|=0
$$

The points of concentration of so called upper main representation are the roots of

$$
(K-z)(z-H)\left|\begin{array}{ccccc}
s_{0}^{\prime} & s_{1}^{\prime} & \ldots & s_{r-2}^{\prime} & 1 \\
s_{1}^{\prime} & s_{2}^{\prime} & \ldots & s_{r-1}^{\prime} & z \\
\ldots & \ldots & \ldots & \cdots & \cdots
\end{array}\right| \cdots \cdots \begin{aligned}
& s_{r-1}^{\prime} \\
& s_{r}^{\prime}
\end{aligned} \ldots
$$

where $s_{k}^{\prime}=(H+K) s_{k+1}-H K s_{k}-s_{k+2}$.
Proof. See [9], p. 122.
Theorem 4.5. Let $n=2 r$. Let the matrices $\boldsymbol{A}$ and $\boldsymbol{B}$ from Theorem 4.2 for a sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r}$ be positive definite. Then the points of concentration of the lower main representation are the roots of

$$
(z-H)\left|\begin{array}{lllll}
s_{1}-H s_{0} & s_{2}-H s_{1} & \ldots & s_{r}-H s_{r-1} & 1 \\
s_{2}-H s_{1} & s_{3}-H s_{2} & \ldots & s_{r+1}-H s_{r} & z \\
\ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots c \cdots \cdots \cdots & \\
s_{r+1}-H s_{r} & s_{r+2}-H s_{r+1} & \ldots & s_{2 r}-H s_{2 r-1} & z^{r}
\end{array}\right|=0
$$

The points of concentration of the upper main representation are the roots of

$$
(K-z)\left|\begin{array}{lllll}
K s_{0}-s_{1} & K s_{1}-s_{2} & \ldots K s_{r-1}-s_{r} & 1 \\
K s_{1}-s_{2} & K s_{2}-s_{3} & \ldots K s_{r}-s_{r+1} & z \\
\ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots & \\
\ldots s_{r}-s_{r+1} & K s_{r+1}-s_{r+2} & \ldots K s_{2 r-1}-s_{2 r} & z^{r}
\end{array}\right|=0
$$

Proof. See [9], p. 122.
It is recommended to prefer the lower main representation. One of the reasons for it may be that in the case $n=2 r-1$ it needs a smaller number of the points of concentration.

If $n=2 r$, we find the points $z_{1}, \ldots, z_{r+1}$ of the lower main representation. Then we calculate probabilities $p_{1}, \ldots, p_{r+1}$ from

$$
z_{1}^{k} p_{1}+\ldots+z_{r+1}^{k} p_{r+1}=s_{k} \quad(k=0,1, \ldots, r) .
$$

The solution has the property that $p_{1} \geqq 0, \ldots, p_{r+1} \geqq 0$ and that

$$
z_{1}^{k} p_{1}+\ldots+z_{r+1}^{k} p_{r+1}=s_{k}
$$

holds also for $k=r+1, \ldots, 2 r$. If $n=2 r-1$, the procedure is analogous.

## 5. SPECIAL DISTRIBUTIONS ON $[-1,1]$

If $X_{t}$ has the continuous rectangular distribution on $[-1,1]$, then its moments are

$$
m_{k}=1 /(k+1) \text { for } k \text { even , } \quad m_{k}=0 \text { for } k \text { odd }
$$

Because $C=-1, D=1$, we get $H=-1+|b|, K=1-|b|$. For simplicity, we consider only the case $b \in(0,1)$. Using the results of Section 3 we get

$$
s_{1}=0, \quad s_{2}=\left(1-b^{2}\right) / 3, \quad s_{3}=0, \quad s_{4}=\left(1-b^{2}\right)\left(3-7 b^{2}\right) / 15
$$

We restrict ourselves to $n=4$. Then

$$
\begin{aligned}
& \boldsymbol{A}=\left(\begin{array}{ccc}
1 & 0 & \left(1-b^{2}\right) / 3 \\
0 & \left(1-b^{2}\right) / 3 & 0 \\
\left(1-b^{2}\right) / 3 & 0 & \left(1-b^{2}\right)\left(3-7 b^{2}\right) / 15
\end{array}\right), \\
& \boldsymbol{B}=(2 / 15)(1-b)(1-2 b)\left(\begin{array}{cc}
5 & 0 \\
0(1+b)(1-3 b)
\end{array}\right)
\end{aligned}
$$

We have $|A|=4\left(1-b^{2}\right)^{2}\left(1-4 b^{2}\right) / 135$. The matrix $A$ is positive definite for $0<b<1 / 2$ and positive semidefinite for $b=\frac{1}{2}$. The matrix B is positive definite for $0<b<\frac{1}{3}$ and positive semidefinite for $b=\frac{1}{2}, b=\frac{1}{3}$. Let $0<b<\frac{1}{3}$. Inserting into the first part of Theorem 4.5 we obtain

$$
\begin{aligned}
& \quad(z-H)\left|\begin{array}{ccc}
-H & s_{2} & 1 \\
s_{2} & -H s_{2} & z \\
-H s_{2} & s_{4} & z^{2}
\end{array}\right|= \\
& =(2 / 45)(1-b)(1-2 b)\left(1-b^{2}\right)(z+1-b) . \\
& \quad \cdot\left[5 z^{2}-2(1+2 b) z-(1-3 b)(1+b)\right] .
\end{aligned}
$$

The roots are

$$
z_{1}=b-1, \quad z_{23}=(1 / 5)\left[1+2 b \mp\left(6-6 b-11 b^{2}\right)^{1 / 2}\right] .
$$

If we insert $b=\frac{1}{2}$ (which is the case not covered by Theorem 4.5), we have $z_{1}=$ $=-0.5, z_{2}=0.3, z_{3}=0.5$. Solving the system of linear equations we get $p_{1}=0.5$, $p_{2}=0, p_{3}=0 \cdot 5$, which is the exact solution (cf. Section 2). Similarly for $b=\frac{1}{3}$ we obtain $z_{1}=-\frac{2}{3}, z_{2}=0, z_{3}=\frac{2}{3}$ and $p_{1}=p_{2}=p_{3}=\frac{1}{3}$ which also agrees with the exact solution.

Let $X_{t}$ have the density

$$
f(x)=2^{-p-q+1} \mathrm{~B}^{-1}(p, q)(x+1)^{p-1}(x-1)^{q-1}, \quad-1<x<1
$$

where $p>0$ and $q>0$ are given parameters and $B$ is the beta function. In the special case when $p=q=1$ we get the rectangular distribution mentioned above. The moments corresponding to $f$ are

$$
m_{k}=(-1)^{k} \sum_{i=0}^{k}\binom{k}{i}(-1)^{i} 2^{i} \frac{p(p+1) \ldots(p+i-1)}{(p+q)(p+q+1) \ldots(p+q+i-1)}
$$

For example, $p=q=0.5$ gives a $U$-distribution of $X_{t}$. Consider its moments $m_{1}, \ldots, m_{6}$. For $b=0.15$ the problem has no solution and for $b=0 \cdot 1$ we obtain the distribution of $e_{t}$ in the form

$$
\begin{array}{lll}
z_{1}=-0.9, & z_{2}=-0.372, & z_{3}=0.333, \\
p_{1}=0.268, & z_{2}=0.895
\end{array}
$$

For $p=0.5$ and $q=3$ using $m_{1}, \ldots, m_{6}$ we come to the conclusion that for $b=0.35$ there is no solution and for $b=0.3$ the distribution of $e_{t}$ is

$$
\begin{array}{llll}
z_{1}=-0.7, & z_{2}=-0.434, & z_{3}=0.060, & z_{4}=0.658 \\
p_{1}=0.599, & p_{2}=0.239, & p_{3}=0.139, & p_{4}=0.023
\end{array}
$$

## 6. MOMENT PROBLEM ON $[0, \infty)$

Theorem 6.1. Let $X_{t}$ be an $\operatorname{AR}(1)$ process defined by (2.1) with $0 \leqq b<1$. If $e_{t} \geqq 0$ for all $t$, then $X_{t} \geqq 0$ for all $t$. If $\mathrm{P}\left(e_{t}<0\right)>0$, then $X_{t}<0$ for infinitely many subscripts $t$ with probability 1.

Proof. The first assertion is clear. If $\mathrm{P}\left(e_{t}<0\right)>0$, then there exist numbers $c>0$ and $q \in(0,1)$ such that $\mathrm{P}\left(e_{t}<-c\right) \geqq q$. Now, we apply Lemma 10.2.

Theorem 6.2. A sequence of numbers $\left\{s_{k}\right\}_{0}^{n}$ is a system of moments on $[0, \infty)$ if and only if the matrices

$$
\boldsymbol{A}=\left(s_{i+j}\right)_{i, j=0}^{[n / 2]}, \quad \mathbf{B}=\left(s_{i+j+1}\right)_{i, j=0}^{[(n-1) / 2]}
$$

are positive semidefinite; [ ] denotes the integer part.
Proof. See [9], p. 237.
An infinite sequence $\left\{s_{k}\right\}_{0}^{\infty}$ is a system of moments on $[0, \infty)$ if and only if the infinite quadratic forms

$$
\sum_{i=0}^{\infty} \sum_{j=0}^{\infty} s_{i+j} x_{i} x_{j}, \quad \sum_{i=0}^{\infty} \sum_{j=0}^{\infty} s_{i+j+1} x_{i} x_{j}
$$

are positive semidefinite (see [1], p. 100).
Theorem 6.3. Let the matrices $\boldsymbol{A}$ and $\boldsymbol{B}$ from Theorem 6.2 be positive definite
for a given sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r}$. Define

$$
\begin{aligned}
& Q_{2 r}(z)=\left|\begin{array}{ccccc}
1 & s_{0} & s_{1} & \ldots & s_{r-1} \\
z & s_{1} & s_{2} & \ldots & s_{r} \\
\ldots & \ldots & \ldots & \ldots & \cdots
\end{array}\right|, \ldots,
\end{aligned}
$$

Then the roots of $Q_{2 r}(z)$ are the points of concentration of the upper main representation and the roots of $Q_{2 r+1}(z)$ are the points of concentration of the lower main representation.

Proof. See [9], p. 260.

## 7. EXPONENTIAL DISTRIBUTION

Let $b \in(0,1)$. If $X_{t}$ has the exponential distribution $E x(\lambda)$ with the density

$$
f(x)=\lambda^{-1} \exp (-x / \lambda), \quad x>0
$$

then in the model (2.1) the $e_{t}$ must be a variable which equals to zero with probability $b$ and which has $E x(\lambda)$ with probability $1-b$ (see [8]). The moments of $e_{t}$ are $s_{k}=k!(1-b) \lambda^{k}$. Let us consider the case of four moments. Both $\boldsymbol{A}$ and $\boldsymbol{B}$ are positive definite for $b \in(0,1)$, the roots of $Q_{5}(z)$ are

$$
z_{1}=0, \quad z_{2}=\left(3-3^{1 / 2}\right) \lambda, \quad z_{3}=\left(3+3^{1 / 2}\right) \lambda
$$

and the corresponding probabilities are

$$
p_{1}=(1+2 b) / 3, \quad p_{2}=(1-b)\left(2+3^{1 / 2}\right) / 6, \quad p_{3}=(1-b)\left(2-3^{1 / 2}\right) / 6
$$

Numerically, if $\lambda=1$ and $b=0 \cdot 5$, we have

$$
\begin{array}{lll}
z_{1}=0, & z_{2}=1.268, & z_{3}=4.732 \\
p_{1}=0.667, & p_{2}=0.311, & p_{3}=0.022
\end{array}
$$

## 8. MOMENT PROBLEM ON $(-\infty, \infty)$

This case is described in [3] and so we introduce briefly only main results.
Theorem 8.1. A sequence of numbers $\left\{s_{k}\right\}_{0}^{2 r}$ is a system of moments on $(-\infty, \infty)$ if and only if the matrix

$$
\boldsymbol{A}=\left(s_{i+j}\right)_{i, j=0}^{r}
$$

is positive semidefinite.
Proof. See [9], p. 246.

A sequence $\left\{s_{k}\right\}_{0}^{\infty}$ is a system of moments on $(-\infty, \infty)$ if and only if the infinite quadratic form

$$
\sum_{i=0}^{\infty} \sum_{j=0}^{\infty} s_{i+j} x_{i} x_{j}
$$

is positive semidefinite (see [9], p. 246).
Theorem 8.2. If $s_{1}, \ldots, s_{2 r}$ are given numbers such that the matrix $\boldsymbol{A}$ given in Theorem 8.1 is positive definite, then the points of concentration of so called canonical representation which includes a given point $w$ are the roots of

$$
\left.Q(z)=\left\lvert\, \begin{array}{cccccc}
s_{0} & s_{1} & \ldots & s_{r-1} & 1 & 1 \\
s_{1} & s_{2} & \ldots & s_{r} & w & z \\
\ldots & \cdots & \cdots & \cdots & \cdots & \cdots
\end{array}\right.\right] \cdots \cdots,
$$

Proof. See [9], p. 247.

## 9. A SPECIAL CASE ON $(-\infty, \infty)$

We proved in Section 5 that for $b \in\left(\frac{1}{3}, \frac{1}{2}\right)$ there exists no distribution of $e_{t}$ such that $X_{t}$ belongs to $[-1,1]$ and has the moments $m_{1}=0, m_{2}=\frac{1}{3}, m_{3}=0, m_{4}=\frac{1}{5}$. However, if we do not insist on the condition $X_{t} \in[-1,1]$, then $\boldsymbol{A}$ being positive definite for $b \in\left(0, \frac{1}{2}\right)$ allows to find a solution on the real line. The polynomial $Q(z)$ with $w=0$ has the roots

$$
z_{13}=\mp\left[\left(3-7 b^{2}\right) / 5\right]^{1 / 2}, \quad z_{2}=0
$$

and the corresponding probabilities are

$$
p_{1}=p_{3}=\frac{5}{6}\left(1-b^{2}\right) /\left(3-7 b^{2}\right), \quad p_{2}=\frac{4}{3}\left(1-4 b^{2}\right) /\left(3-7 b^{2}\right)
$$

For example, if $b=0.4$ then

$$
\begin{array}{lll}
z_{1}=-0.613, & z_{2}=0, & z_{3}=0.613 \\
p_{1}=0.372, & p_{2}=0.256, & p_{3}=0.372
\end{array}
$$

In this special case we obtain $-1.532 \leqq X_{t} \leqq 1.532$.

## 10. APPENDIX

Lemma 10.1. Let $A_{1}, A_{2}, \ldots$ and $B_{1}, B_{2}, \ldots$ be two sequences of events satisfying the following conditions.
(i) The events $A_{1}, A_{2}, \ldots$ are independent.
(ii) For each $n$, the events $A_{n}$ and $B_{n}$ are independent.
(iii) $\sum \mathrm{P}\left(A_{n}\right)=\infty$.
(iv) $\mathrm{P}\left(B_{n}\right) \rightarrow 1$.

Then infinitely many events $A_{n} \cap B_{n}$ occur with probability one.
Proof. See [5].
Lemma 10.2. Let $X_{t}$ be an $\operatorname{AR}(1)$ process defined by (2.1) with $0 \leqq b<1$. If there exist numbers $c>0$ and $q \in(0,1)$ such that $\mathrm{P}\left(e_{t}<-c\right) \geqq q$, then $X_{t}<0$ for infinitely many subscripts with probability one.

Proof. Let $j_{k}$ be the smallest integer such that

$$
j_{k} q^{k} \geqq 1 \text { for } k=1,2, \ldots
$$

Arrange the positive integers into the subsets $S_{1}, S_{2}, \ldots$ in the following way. Let $S_{1}=\left\{1, \ldots, j_{1}\right\}$. Let $S_{2}$ contain the elements of $j_{2}$ couples

$$
\left(j_{1}+1, j_{1}+2\right),\left(j_{1}+3, j_{1}+4\right), \ldots,\left(j_{1}+2 j_{2}-1, j_{1}+2 j_{2}\right)
$$

let $S_{3}$ contain elements of $j_{3}$ triples starting with $\left(j_{1}+2 j_{2}+1, j_{1}+2 j_{2}+2\right.$, $j_{1}+2 j_{2}+3$ ) and so on. The numbers $1, \ldots, j_{1}$, the last numbers of couples, the last numbers of triples etc. denote $n_{1}, n_{2}, \ldots$. If $n_{i} \in S_{k}$ then we use the decomposition

$$
X_{n_{i}}=U_{n_{i}}+Z_{n_{i}}
$$

where

$$
\begin{gathered}
U_{n_{i}}=e_{n_{i}}+b e_{n_{i}-1}+\ldots+b^{k-1} e_{n_{i}-k+1} \\
Z_{n_{i}}=b^{k} e_{n_{i}-k}+b^{k+1} e_{n_{i}-k-1}+\ldots
\end{gathered}
$$

Define $\mu=\mathrm{E} e_{t}, \sigma^{2}=\operatorname{var} e_{t}$. Then

$$
E Z_{n_{i}}=\mu b^{k} /(1-b), \quad \operatorname{var} Z_{n_{i}}=\sigma^{2} b^{2 k} /\left(1-b^{2}\right)
$$

and thus

$$
\mathrm{E} Z_{n_{i}}^{2}=\operatorname{var} Z_{n_{i}}+\left(\mathrm{E} Z_{n_{i}}\right)^{2} \rightarrow 0 \quad \text { as } \quad i \rightarrow \infty
$$

Introduce events

$$
A_{i}=\left\{U_{n_{i}}<-c\right\}, \quad B_{i}=\left\{Z_{n_{i}}<c\right\} .
$$

Since

$$
\mathrm{P}\left(Z_{n_{i}} \geqq c\right) \leqq c^{-2} \mathrm{E} Z_{n_{i}}^{2} \rightarrow 0
$$

we have $\mathrm{P}\left(Z_{n_{i}}<c\right) \rightarrow 1$. It is clear that

$$
\sum_{i=1}^{\infty} \mathrm{P}\left(A_{i}\right)=\underbrace{q+\ldots+q}_{j_{1}}+\underbrace{q^{2}+\ldots+q^{2}}_{j_{2}}+\ldots=\infty
$$

The events $A_{1}, A_{2}, \ldots$ are independent and for each $i$ the events $A_{i}$ and $B_{i}$ are also independent. Thus the assertion follows from Lemma 10.1.
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