# On the intersections between the trajectories of a normal stationary stochastic process and a high level 

By Harald Cramér

## 1. Introduction

Let $\xi(t)$ be a real-valued, normal and stationary stochastic process with a continuous time parameter $t$, varying over $(-\infty, \infty)$. Suppose that $E \xi(t)=0$ for all $\boldsymbol{t}$, and that the covariance function

$$
r(t-u)=E \xi(t) \xi(u)
$$

satisfies the following two conditions:

$$
\begin{equation*}
r(t)=1-\frac{1}{2!} \lambda_{2} t^{2}+\frac{1}{4!} \lambda_{4} t^{4}+o\left(t^{4}\right) \tag{1}
\end{equation*}
$$

as $t \rightarrow 0$, and

$$
\begin{equation*}
r(t)=O\left(|t|^{-\alpha}\right) \tag{2}
\end{equation*}
$$

for some $\alpha>0$, as $t \rightarrow \pm \infty$. We may, of course, assume $\alpha<1$.
It follows ${ }^{1}$ from (1) that there exists an equivalent version of $\xi(t)$ having, with probability one, a continuous sample function derivative $\xi^{\prime}(t)$, and it will be supposed that $\xi(t)$ has, if required, been replaced by this equivalent version.

Let $u>0$ be given, and consider the intersections between the trajectories $\eta=\xi(t)$ of the $\xi$ process and the horizontal line, or "level", $\eta=u$ during some finite time interval, say $0<t<T$. It follows from a theorem due to Bulinskaja [2] that, with probability one, there are only a finite number of such intersections, and also that, with probability one, there is no point of tangency between the trajectory $\eta=\xi(t)$ and the level $\eta=u$ during the interval $(0, T)$. Thus, with probability one, every point with $\xi(t)=u$ can be classified as an "upcrossing" or a "downcrossing" of the level $u$, according as $\xi^{\prime}(t)$ is positive or negative.

The present paper will be concerned with the upcrossings, and their asymptotic distribution in time, as the level $u$ becomes large. It will be obvious that the case of a large negative $u$, as well as the corresponding problem for the downcrossings, can be treated in the same way.

The upcrossings may be regarded as a stationary stream of random events (cf.,

[^0]
## h. cramér, Stationary stochastic process and high level

e.g. Khintchine's book [5]), and it is well known that the simplest case of such a stream occurs when the successive events form a Poisson process. However, a necessary condition for this case is (Khintchine, l.c., pp. 11-12) that the numbers of events occurring in any two disjoint time intervals should be independent random variables, and it is readily seen that this condition cannot be expected to be satisfied by the stream of upcrossings. On the other hand, it seems intuitively plausible that the independence condition should be at least approximately satisfied when the level $u$ becomes very large, provided that values of $\xi(t)$ lying far apart on the time scale can be supposed to be only weakly dependent.

Accordingly it may be supposed that, subject to appropriate conditions on $\xi(t)$, the stream of upcrossings will tend to form a Poisson process as the level $u$ tends to infinity.

That this is actually so was first proved by Volkonskij and Rozanov in their remarkable joint paper [8]. They assumed, in addition to the condition (1) above, that $\xi(t)$ satisfies the so-called condition of strong mixing. This is a fairly restrictive condition, as can be seen e.g. from the analysis of the strong mixing concept given by Kolmogorov and Rozanov [6]. Moreover, in an actual case it will not always be easy to decide whether the condition is satisfied or not. On the other hand, various interesting properties of the $\xi(t)$ trajectories follow as corollaries from the asymptotic Poisson character of the stream of uperossings (cf. Cramér [3], Cramér and Leadbetter [4]), so that it seems highly desirable to prove the latter property under simpler and less restrictive conditions.

It will be shown in the present paper that it is possible to replace the condition of strong mixing by the condition (2) above. This is considerably more general, and also simpler to deal with in most applications.

I am indebted to Dr. Yu. K. Belajev for stimulating conversations about the problem treated in this paper.

## 2. The main theorem

The number of upcrossings of $\xi(t)$ with the level $u$ during the time interval $(s, t)$ will be denoted by $N(s, t)$. This is a random variable which, by the above remarks, is finite with probability one. When $s=0$, we write simply $N(t)$ instead of $N(0, t)$. From the stationarity and the condition (1) it follows (Bulinskaja [2], Cramér and Leadbetter [4], Ch. 10) that the mean value of $N(s, t)$ is, for $s<t$,

$$
\begin{equation*}
E N(s, t)=E N(t-s)=\mu(t-s), \tag{3}
\end{equation*}
$$

where

$$
\begin{equation*}
\mu=E N(1)=\frac{V \overline{\lambda_{2}}}{2 \pi} e^{-u^{\prime} / 2} \tag{4}
\end{equation*}
$$

The quantity $\mu$ will play an important part in the sequel. We note that $\mu$ is a function of the level $u$, and that $\mu$ tends to zero as $u$ tends to infinity. From (3) we obtain for any $\tau>0$

$$
E N(\tau / \mu)=\tau
$$

For the study of the stream of upcrossings, it will then seem natural to choose $1 / \mu$ as a scaling unit of time, thus replacing $t$ by $\tau / \mu$. We might expect, e.g., that the probability distribution of $N(\tau / \mu)$ will tend to some limiting form when $\mu$ tends to zero while $\tau$ remains fixed. This is in fact the case, as shown by the following theorem, first proved by Volkonskij and Rozanov under more restrictive conditions, as mentioned above.

Theorem. Suppose that the normal and stationary process $\xi(t)$ satisfies the conditions (1) and (2). Let $\left(a_{1}, b_{1}\right), \ldots,\left(a_{j}, b_{j}\right)$ be disjoint time intervals depending on $u$ in such a way that, for $i=1, \ldots, j$,

$$
b_{i}-a_{i}=\tau_{i} / \mu,
$$

the integer $j$ and the positive numbers $\tau_{1}, \ldots \tau_{j}$ being independent of $u$. Let $k_{1}, \ldots, k_{j}$ be non--negative integers independent of $u$. Then

$$
\lim _{u \rightarrow \infty} P\left\{N\left(a_{i}, b_{i}\right)=k_{i} \text { for } i=1, \ldots, j\right\}=\prod_{i=1}^{j} \frac{\tau_{i}^{k_{i}}}{k_{i}!} e^{-\tau_{i}}
$$

Thus, when time is measured in units of $1 / \mu$, the stream of upcrossings will asymptotically behave as a Poisson process as the level $u$ tends to infinity.

We shall first give the proof for the case $j=1$, when there is one single interval $(a, b)$ of length $b-a=\tau / \mu$. Owing to the stationarity it is sufficient to consider the interval $(0, \tau / \mu)$. Writing

$$
\begin{equation*}
T=\tau / \mu, \quad P_{k}=P\{N(T)=k\} \tag{5}
\end{equation*}
$$

we then have to prove the relation

$$
\begin{equation*}
\lim _{u \rightarrow \infty} P_{k}=\frac{\tau^{k}}{k!} e^{-\tau} \tag{6}
\end{equation*}
$$

for any given $\tau>0$ and non-negative integer $k$, both independent of $u$. Once this has been achieved, the proof of the general case will follow in a comparatively simple way.

The proof of (6) is rather long, and will be broken up in a series of lemmas. In the following section we shall introduce some notations that will be used in the course of the proof. The lemmas will be given in sections 4 and 5, while section 6 contains the proof of the case $j=1$ of the theorem, and section 7 the proof of the general case.

## 3. Notations

The level $u$ will be regarded as a variable tending to infinity, and we must now introduce various functions of $u$. It will be practical to define them as functions of $\mu$, where $\mu$ is the function of $u$ given by (4). Writing as usual $[x]$ for the greatest integer $\leqslant x$, we define

$$
\left.\begin{array}{rl}
m_{1} & =\left[\mu^{-1}\right], \quad m_{2}  \tag{7}\\
M & =\left[\mu^{\beta-1}\right], \\
M & =\left[T / \mu^{\beta}\right] .
\end{array} \quad n=\left[M /\left(m_{1}+m_{2}\right)\right]+1 .\right\}
$$

## H. Cramér, Stationary stochastic process and high level

Here $\beta$ is any number satisfying the relation

$$
\begin{equation*}
0<(k+4) \beta<\alpha<1, \tag{8}
\end{equation*}
$$

where $\alpha$ is the constant occurring in the condition (2), while $k$ is the integer occurring in (6). We further write

$$
\begin{equation*}
q=T / M, \quad t_{1}=m_{1} q, \quad t_{2}=m_{2} q, \tag{9}
\end{equation*}
$$

and divide the interval $(0, T)$ on the time axis into subintervals, alternatively of length $t_{1}$ and $t_{2}$, starting from the origin. We shall refer to these subintervals as $t_{1^{-}}$and $t_{2}$-intervals respectively, the former being regarded as closed and the latter as open. Each $t_{i}$-interval $(i=1,2)$ consists of $m_{i}$ subintervals of length $q$. The whole interval $(0, T)$, which consists of $M$ intervals of length $q$, is covered by $n$ pairs of $t_{1}$ - and $t_{2}$-intervals, the $n$th pair being possibly incomplete. Any two distinct $t_{1}$-intervals are separated by an interval of length at least equal to $t_{2}$. An important use will be made of this remark in the proof of Lemma 5 below.

The quantities defined by (7) and (9) are all functions of $u$. It will be practical to express their order of magnitude for large $u$ in terms of $\mu$. The following relations are easily obtained from (5), (7) and (9):

$$
\left.\begin{array}{ll}
q \sim \mu^{\beta}, & n \sim \tau \mu^{-\beta}  \tag{10}\\
t_{1} \sim \mu^{\beta-1}, & t_{2} \sim \mu^{2 \beta-1}
\end{array}\right\}
$$

We now define a stochastic process $\xi_{q}(t)$ by taking

$$
\xi_{q}(v q)=\xi(v q)
$$

for all integers $\nu$, and determining $\xi_{\alpha}$ by linear interpolation in the interval between two consequitive $\nu q$. To any sample function of the $\xi(t)$ process will then correspond a sample function of $\xi_{q}(t)$, which is graphically represented by the broken line joining the points [ $\nu q, \xi(\nu q)$ ]. For the number of upcrossings of this broken line with the $u$ level we use the notations $N_{q}(s, t)$ and $N_{q}(t)$, corresponding to $N(s, t)$ and $N(t)$. The probability corresponding to $P_{k}$ as defined by (5) is

$$
\begin{equation*}
P_{k}^{(q)}=P\left\{N_{q}(T)=k\right\} . \tag{11}
\end{equation*}
$$

## 4. Lemmas 1-3

Throughout the rest of the paper we assume that the $\xi(t)$ process satisfies the conditions of the above theorem.

Lemma 1. If $T$ and $q$ are given by (5) and (9), $\tau$ and $k$ being fixed as before, we have

$$
\lim _{u \rightarrow \infty}\left(P_{k}-P_{k}^{(\alpha)}\right)=0 .
$$

Evidently $N_{q}(T) \leqslant N(T)$. We shall prove that the non-negative and integervalued random variable $N(T)-N_{q}(T)$ converges in first order mean to zero,
as $u \rightarrow \infty$. It then follows that the probability that $N(T)-N_{q}(T)$ takes any value different from zero will tend to zero, and so the lemma will be proved.

By (3) and (5), the mean value of the number $N(T)$ of upcrossings of $\xi(t)$ in $(0, T)$ is for every $u$

$$
E N(T)=T \mu=\tau
$$

It will now be proved that the mean value $E N_{q}(T)$ tends to the limit $\tau$ as $u \rightarrow \infty$, so that we have

$$
\lim _{u \rightarrow \infty} E\left\{N(T)-N_{q}(T)\right\}=0
$$

Since $N(T)-N_{q}(T)$ is non-negative, this implies convergence in first order mean to zero, so that by the above remark the lemma will be proved.

Consider first the number $N_{q}(q)$ of upcrossings of $\xi_{q}(t)$ in the interval $(0, q)$. This number is one, if $\xi(0)<u<\xi(q)$, and otherwise zero, so that

$$
E N_{q}(q)=P\{\xi(0)<u<\xi(q)\} .
$$

Now $\xi(0)$ and $\xi(q)$ have a joint normal density function, with unit variances and correlation coefficient $r=r(q)$. By (1) and (10) we have

$$
\begin{equation*}
r(q)=1-\frac{1}{2} \lambda_{2} q^{2}+O\left(q^{4}\right) . \tag{12}
\end{equation*}
$$

For the probability that $\xi(0)<u$ and $\xi(q)>u$ we obtain by a standard transformation

$$
E N_{q}(q)=\frac{1}{\sqrt{2 \pi}} \int_{u}^{\infty} e^{-x^{q} / 2} \Phi\left(\frac{u-r x}{\sqrt{1-r^{2}}}\right) d x
$$

where as usual

$$
\Phi(x)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{x} e^{-t^{2} / 2} d t
$$

By some straightforward evaluation of the integral from $u+(1-r)^{\frac{1}{2}}$ to infinity we obtain, using (12), and denoting by $K$ an unspecified positive constant,

$$
E N_{q}(q)=\frac{1}{\sqrt{2 \pi}} \int_{u}^{u+(1-r)^{\frac{1}{2}}} e^{-x^{2} / 2} \Phi\left(\frac{u-r x}{\sqrt{1-r^{2}}}\right) d x+O\left[\exp \left(-e^{K u^{2}}\right)\right]
$$

For the first term in the second member we obtain, using again (12), the expression

$$
\begin{aligned}
\frac{1}{\sqrt{2 \pi}} e^{-u^{2} / 2}(1 & \left.+O\left(q^{\frac{1}{2}} u\right)\right) \int_{u}^{u+(1-r)^{\frac{1}{2}}} \Phi\left(\frac{u-r x}{\sqrt{1-r^{2}}}\right) d x \\
& =\frac{\sqrt{1-r^{2}}}{r \sqrt{2 \pi}} e^{-u^{2} / 2}\left(1+O\left(q^{\frac{1}{2}} u\right)\right)\left(\int_{-\infty}^{0} \Phi(y) d y+O(q u)\right)=\frac{\sqrt{\lambda_{2}}}{2 \pi} q e^{-u^{3} / 2}\left(1+O\left(q^{\frac{1}{2}} u\right)\right),
\end{aligned}
$$

so that

$$
E N_{q}(q)=q \mu+o(q \mu) .
$$

If $v$ denotes an integer, which may tend to infinity with $u$, it then follows from the stationarity that we have

$$
\begin{equation*}
E N_{q}(v q)=\nu q \mu+o(v q \mu) . \tag{13}
\end{equation*}
$$

In particular, taking $v=M$ we obtain from (9) and (5)

$$
E N_{q}(T)=T \mu+o(T \mu) \rightarrow \tau
$$

According to the above remarks, this proves the lemma.
We now consider the number $N_{q}\left(t_{1}\right)$ of $\xi_{q}$ upcrossings in an interval of length $t_{1}=m_{1} q$, observing that by (10) we have $t_{1} \sim \mu^{\beta-1}$.

Lemma 2. We have

$$
\lim _{u \rightarrow \infty} \frac{E\left\{N_{q}\left(t_{1}\right)\left[N_{q}\left(t_{1}\right)-1\right]\right\}}{E N_{q}\left(t_{1}\right)}=0 .
$$

By (3) we have $E N\left(t_{1}\right)=t_{1} \mu$, while (13) gives for $\nu=m_{1}$

$$
E N_{q}\left(t_{1}\right)=t_{1} \mu+o\left(t_{1} \mu\right) \sim E N\left(t_{1}\right) .
$$

Further, since $N_{q}\left(t_{1}\right) \leqslant N\left(t_{1}\right)$,

$$
E\left\{N_{q}(t)\left[N_{q}\left(t_{1}\right)-1\right]\right\} \leqslant E\left\{N\left(t_{1}\right)\left[N\left(t_{1}\right)-1\right\} .\right.
$$

The truth of the lemma will then follow from the corresponding relation with $N_{q}$ replaced by $N$. Now this latter relation is identical with the relation proved by Volkonskij and Rozanov [8] in their Lemma 3.4. It is proved by them without any mixing hypothesis, assuming only that $\xi(t)$ is regular (or purely non-deterministic) and that $r(t)$ has a fourth order derivative at $t=0$. Their proof is valid without any modification whatever, if these conditions are replaced by our conditions (1) and (2). Thus we may refer to their paper for the proof of this lemma. We note that the proof is based on the important work of S. O. Rice [7].

Lemma 3. As $u \rightarrow \infty$, we have

$$
\begin{aligned}
& P\left\{N_{q}\left(t_{1}\right)=0\right\}=1-q+o(q), \\
& P\left\{N_{q}\left(t_{1}\right)=1\right\}=q+o(q), \\
& P\left\{N_{q}\left(t_{1}\right)>1\right\}=o(q) .
\end{aligned}
$$

For any random variable $\nu$ taking only non-negative integral values we have, writing $\pi_{i}=P\{\nu=i\}$ and assuming $E \nu^{2}<\infty$,

$$
\begin{gathered}
E v=\pi_{1}+2 \pi_{2}+3 \pi_{3}+\ldots \\
E v(v-1)=2 \pi_{2}+6 \pi_{3}+\ldots
\end{gathered}
$$

and consequently

$$
\begin{equation*}
E \nu-E v(\nu-1) \leqslant \pi_{1} \leqslant 1-\pi_{0} \leqslant E v \tag{14}
\end{equation*}
$$

Taking $\nu=N_{q}\left(t_{1}\right)$, and observing that by (10) we have $E N_{q}\left(t_{1}\right) \sim t_{1} \mu \sim q$, the truth of the lemma follows directly from Lemma 2.

## 5. Lemmas 4-5

For each $r=1,2, \ldots, n$, we now define the following events, i.e. the sets of all $\xi(t)$ sample functions satisfying the conditions written between the brackets:

$$
\begin{aligned}
c_{r} & =\left\{\text { exactly one } \xi_{q} \text { upcrossing in the } r \text { th } t_{1} \text {-interval }\right\}, \\
d_{r} & =\left\{\text { at least one } \xi_{q} \text { upcrossing in the } r \text { th } t_{1} \text {-interval }\right\} \\
e_{r} & =\left\{\xi(\nu q)>u \text { for at least one } \nu q \text { in the } r \text { th } t_{1} \text {-interval }\right\} .
\end{aligned}
$$

Further, let $C_{k}$ denote the event that $c_{r}$ occurs in exactly $k$ of the $t_{1}$-intervals in ( $0, T$ ), while the complementary event $c_{r}^{*}$ occurs in the $n-k$ others. $D_{k}$ and $E_{k}$ are defined in the corresponding way, using respectively $d_{r}$ and $e_{r}$ instead of $c_{r}$.

Lemma 4. For the probability $P_{k}^{(q)}$ defined by (11) we have

$$
\lim _{u \rightarrow \infty}\left[P_{k}^{(q)}-P\left\{E_{k}\right\}\right]=0
$$

We shall prove that each of the differences $P_{k}^{(q)}-P\left\{C_{k}\right\}, P\left\{C_{k}\right\}-P\left\{D_{k}\right\}$ and $P\left\{D_{k}\right\}-P\left\{E_{k}\right\}$ tends to zero as $u \rightarrow \infty$.

By (13) and (14) the probability of at least one $\xi_{q}$ upcrossing in an interval of length $t_{2}$ is at most $E N_{q}\left(t_{2}\right)=t_{2} \mu+o\left(t_{2} \mu\right)$. Thus the probability of at least one $\xi_{q}$ upcrossing in at least one of the $n t_{2}$-intervals in ( $0, T$ ) is by (10)

$$
O\left(n t_{2} \mu\right)=O\left(\mu^{\beta}\right)
$$

and thus tends to zero as $u \rightarrow \infty$. It follows that we have

$$
\begin{equation*}
P_{k}^{(q)}-P\left\{\text { total number of } \xi_{q} \text { upcrossings in all } n t_{1} \text {-intervals }=k\right\} \rightarrow 0 \tag{15}
\end{equation*}
$$

On the other hand, by the stationarity of $\xi(t)$, Lemma 3 remains true if $N_{q}\left(t_{1}\right)$ is replaced by the number of $\xi_{q}$ upcrossings in any particular $t_{1}$-interval. Since the interval $(0, T)$ contains $n$ of these intervals, it follows from (10) that the probability of more than one $\xi_{q}$ upcrossing in at least one of the $t_{1}$-intervals is $o(n q)=o(1)$, and thus tends to zero as $u \rightarrow \infty$.

From (15) and the last remark, it now readily follows that the differences $P_{k}^{(\alpha)}-P\left\{C_{k}\right\}$ and $P\left\{C_{k}\right\}-P\left\{D_{k}\right\}$ both tend to zero as $u \rightarrow \infty$. It thus only remains to show that this is true also for $P\left\{D_{k}\right\}-P\left\{E_{k}\right\}$.

By the definitions of the events $D_{k}$ and $E_{k}$ we have

$$
\left.\begin{array}{l}
P\left\{D_{k}\right\}=\sum P\left\{d_{r_{1}} \ldots d_{r_{k}} d_{s_{1}}^{*} \ldots d_{s_{n-k}}^{*}\right\},  \tag{16}\\
P\left\{E_{k}\right\}=\sum P\left\{e_{r_{1}} \ldots e_{r_{k}} e_{s_{1}}^{*} \ldots e_{s_{n-k}}^{*}\right\},
\end{array}\right\}
$$

the summations being extended over all $\binom{n}{k}$ groups of $k$ different subscripts $r_{1}, \ldots, r_{k}$ selected among the numbers $1, \ldots, n$, while in each case $s_{1}, \ldots, s_{n-k}$ are the remaining $n-k$ subscripts.

Let $v_{r} q$ denote the left endpoint of the $r$ th $t_{1}$-interval, and denote by $g_{r}$ the event

$$
g_{r}=\left\{\xi\left(\nu_{r} q\right) \geqslant u\right\}
$$

of probability

$$
P\left\{g_{r}\right\}=O\left(\frac{1}{u} e^{-u^{2} / 2}\right)
$$

Then for every $r=1, \ldots, n$

$$
d_{\tau} \subset e_{r} \quad \text { and } \quad e_{r}-d_{r} \subset g_{r},
$$

so that

$$
\begin{equation*}
P\left\{e_{r}-d_{r}\right\}=O\left(\frac{1}{u} e^{-u^{2} / 2}\right) \tag{17}
\end{equation*}
$$

Similarly

$$
\begin{aligned}
& e_{r}^{*} \subset d_{r}^{*}, \quad \text { and } \\
& d_{r}^{*}-e_{r}^{*}=e_{r}-d_{r} \subset g_{r}, \\
& P\left\{d_{r}^{*}-e_{r}^{*}\right\}=O\left(\frac{1}{u} e^{-u^{2} / 2}\right) .
\end{aligned}
$$

By a simple recursive argument (16) then yields, using (8) and (10),

$$
P\left\{D_{k}\right\}-P\left\{E_{k}\right\}=O\left(\frac{n^{k+1}}{u} e^{-u^{\star} / 2}\right)=O\left[\exp \left(-\frac{1-(k+1) \beta}{2} u^{2}\right)\right] \rightarrow 0
$$

which proves the lemma.
By definition, $E_{k}$ is composed of certain events $e_{r}$ and $e_{s}^{*}$. Each of these is associated with one particular $t_{1}$-interval, and it has been observed above that any two $t_{1}$-intervals are separated by an interval which is of length $\geqslant t_{2}$, and thus tends to infinity with $u$. By means of the condition (2) it will now be shown that the component events of $E_{k}$ are asymptotically independent, as $u \rightarrow \infty$. Moreover, owing to stationarity, the probability

$$
\begin{equation*}
p=P\left\{e_{r}\right\} \tag{18}
\end{equation*}
$$

is independent of $r$, so that by (16) the asymptotic independence will be expressed by the following lemma.

Lemma 5. The probability $p$ being defined by (18) we have, as $u \rightarrow \infty$,

$$
P\left\{E_{k}\right\}-\binom{n}{k} p^{k}(1-p)^{n-k} \rightarrow 0
$$

In order to prove this lemma, we consider the points $v q$ on the time axis for all integers $v$ such that $v q$ belongs to one of the $t_{1}$-intervals in $(0, T)$. Each $t_{1}$-interval, which we regard as closed, contains $m_{1}+1$ points $v q$, and there are $n-1$ complete and one possibly incomplete such interval in $(0, T)$. If $L$ is the total number of points $v q$ in all $t_{1}$-intervals, we thus have

$$
(n-1)\left(m_{1}+1\right)<L \leqslant n\left(m_{1}+1\right) .
$$

Let $\eta_{1}, \ldots, \eta_{L}$ be the random variables $\xi(\nu \gamma)$ corresponding to all these $L$ points
$\nu q$, ordered according to increasing $\nu$. The $\eta_{i}$ sequence will consist of $n$ groups, each corresponding to one particular $t_{1}$-interval.

Further, let $f_{1}\left(y_{1}, \ldots, y_{L}\right)$ be the $L$-dimensional normal probability density of $\eta_{1}, \ldots, \eta_{L}$, and let $\Lambda_{1}$ be the corresponding covariance matrix. (Our reasons for using the subscript I here and in the sequel will presently appear.) From (16) we obtain

$$
\begin{equation*}
P\left\{E_{k}\right\}=\int_{E_{k}} f_{1} d y=\sum \int_{e_{r_{1}} \ldots e_{s_{1}}^{*} \ldots} f_{1} d y \tag{19}
\end{equation*}
$$

where the abbreviated notation should be easily understood, the summation being extended as explained after (16).

Let us now consider one particular term of the sum in the last member of (19), say the term where the group of subscripts $r_{1}, \ldots, r_{k}$ coincides with the integers $1, \ldots, k$. It will be readily seen that any other term can be treated in the same way as we propose to do with this one. This term is

$$
F(1)=\int_{G} f_{1} d y
$$

where $G$ denotes the set

$$
G=e_{1} \ldots e_{k} e_{k+1}^{*} \ldots e_{n}^{*}
$$

$F(1)$ may be regarded as a function of the covariances which are elements of the matrix $\Lambda_{1}$. Let us consider in particular the dependence of $F(1)$ on those covariances $\varrho_{i j}=E \eta_{i} \eta_{j}$ which correspond to variables $\eta_{i}$ and $\eta_{j}$ belonging to different $t_{1}$-intervals. If all covariances $\varrho_{i j}$ having this character are replaced by $\lambda_{i j}=h \varrho_{i j}$, with $0 \leqslant h \leqslant 1$, while all other elements of $\Lambda_{1}$ remain unchanged, the resulting matrix will be

$$
\begin{equation*}
\Lambda_{h}=h \Lambda_{1}+(\mathbf{l}-h) \Lambda_{0} \tag{20}
\end{equation*}
$$

while the density function $f_{1}$ will be replaced by a certain function $f_{h}$. Evidently $f_{0}$, corresponding to the covariance matrix $\Lambda_{0}$, will be the normal density function that would apply if the groups of variables $\eta_{i}$ belonging to different $t_{1^{-}}$ intervals were all mutually independent, while the joint distribution within each group were the same as before.

Thus $\Lambda_{1}$ and $\Lambda_{0}$ are both positive definite, and it then follows from (20) that the same is true for $\Lambda_{h}$, so that $f_{h}$ is always a normal probability density. Writing

$$
F(h)=\int_{G} f_{h} d y
$$

it follows from the remarks just made that we have

$$
F(0)=\int_{e_{1}} f_{0} d y \ldots \int_{e_{k}} t_{0} d y \int_{e_{k+1}^{*}} f_{0} d y \ldots \int_{e_{n}^{*}} f_{0} d y=P\left\{e_{1}\right\} \ldots P\left\{e_{k}\right\} P\left\{e_{k+1}^{*}\right\} \ldots P\left\{e_{n}^{*}\right\}
$$

By stationarity this reduces to

$$
\begin{equation*}
F(0)=p^{k}(1-p)^{n-k} \tag{21}
\end{equation*}
$$

where $p$ is given by (18).

## H. Cramer, Stationary stochastic process and high level

We shall now evaluate the difference $F(1)-F(0)$ by a development of a method used by S. M. Berman [1]. We note that for any normal density function $f\left(x_{1}, \ldots, x_{n}\right)$ with zero means and covariances $r_{i j}$ we have

$$
\frac{\partial f}{\partial r_{i j}}=\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}
$$

In our case, $f_{h}\left(y_{1}, \ldots, y_{L}\right)$ is a normal density, depending on $h$ through the covariances $\lambda_{i j}=h \varrho_{i j}$. Hence

$$
\begin{equation*}
F^{\prime}(h)=\int_{G} \frac{d f_{h}}{d h} d y=\sum \varrho_{i j} \int_{G} \frac{\partial f_{h}}{\partial \lambda_{i j}} d y=\sum \varrho_{i j} \int_{G} \frac{\partial^{2} f_{h}}{\partial y_{i} \partial y_{j}} d y \tag{22}
\end{equation*}
$$

the summation being extended over all $i, j$ such that $\eta_{i}$ and $\eta_{j}$ belong to different $t_{1}$-intervals. With respect to the integral over the set $G=e_{1} \ldots e_{k} e_{k+1}^{*} \ldots e_{n}^{*}$ occurring in the last sum in (22), we have to distinguish three different cases.

Case $A$. When $\eta_{i}$ and $\eta_{j}$ both belong to $t_{1}$-intervals of subscripts $>k$, say to the $t_{1}$-intervals of subscripts $k+1$ and $k+2$ respectively, integration with respect to $y_{i}$ and $y_{j}$ has to be performed over $e_{k+1}^{*}$ and $e_{k+2}^{*}$ respectively. By definition of the sets $e_{r}$, both $y_{i}$ and $y_{j}$ thus have to be integrated over $(-\infty, 0)$, and so we obtain by direct integration with respect to $y_{i}$ and $y_{j}$

$$
\begin{equation*}
\int_{G} \frac{\partial^{2} f_{h}}{\partial y_{i} \partial y_{j}} d y=\int_{G} f_{h}\left(y_{i}=y_{j}=u\right) d y^{\prime} \tag{23}
\end{equation*}
$$

The notation used in the last integral is to be understood so that we have to take $y_{i}=y_{j}=u$ in $f_{h}$, and then integrate with respect to all $y$ 's expect $y_{i}$ and $y_{j}$. As $t_{h}>0$ always, we have

$$
0<\int_{G} \frac{\partial^{2} f_{n}}{\partial y_{i} \partial y_{j}} d y<\int_{-\infty}^{\infty} \ldots \int_{-\infty}^{\infty} f_{h}\left(y_{i}=y_{j}=u\right) d y^{\prime} .
$$

The last integral, where all the $y$ 's except $y_{i}$ and $y_{j}$ are integrated out, yields the joint density function of the random variables corresponding to $\eta_{i}$ and $\eta_{i}$ in the normal distribution with covariance matrix $\Lambda_{h}$, for the values $y_{i}=y_{j}=u$, so that

$$
\begin{equation*}
0<\int_{G} \frac{\partial^{2} f_{h}}{\partial y_{i} \partial y_{j}} d y<\frac{1}{2 \pi\left(1-h^{2} \varrho_{i j}^{2}\right)^{\frac{1}{2}}} \exp \left[-u^{2} /\left(1+h\left|\varrho_{i j}\right|\right)\right] . \tag{24}
\end{equation*}
$$

Case B. Let now $\eta_{i}$ and $\eta_{j}$ both belong to $t_{1}$-intervals of subscripts $\leqslant k$, say to those of subscripts 1 and 2 respectively. Then integration with respect to each of the groups of variables to which $y_{i}$ and $y_{j}$ belong has to be performed over $e_{1}$ and $e_{2}$ respectively. By definition, $e_{1}$ is the set of all points in the $y$ space such that at least one of the $y$ 's associated with the first $t_{1}$-interval exceeds $u$, and correspondingly for $e_{2}$. Some reflection will then show that the integration indicated in the first member of (23) can still be carried out directly, and yields the same result, with the only difference that in the second member
of (23) the integration has to be performed over a set $G^{\prime}$, obtained from $G$ by replacing $e_{1}$ and $e_{2}$ by $e_{1}^{*}$ and $e_{2}^{*}$ respectively. It follows that the inequality (24) still holds.

Case C. Finally we have the case when $\eta_{i}$ and $\eta_{j}$ belong to $t_{1}$-intervals of different kinds, say to the first and the $(k+1)$ st respectively. As before the integration in the first member of (23) can be carried out directly. In this case, however, we obtain the relation (23) with a changed sign of the second member, and $e_{1}$ replaced by $e_{1}^{*}$ in the expression of the domain of integration. In this case we thus obtain the inequality (24) with changed inequality signs.

Thus in all three cases we have the inequality

$$
\begin{equation*}
\left|\int_{G} \frac{\partial^{2} f_{h}}{\partial y_{i} \partial y_{j}} d y\right|<\frac{1}{2 \pi\left(1-h^{2} \varrho_{i j}^{2}\right)^{\frac{1}{2}}} \exp \left[-u^{2} /\left(1+h\left|\varrho_{i j}\right|\right)\right] . \tag{25}
\end{equation*}
$$

Now $\varrho_{i j}$ is the covariance between the variables $\eta_{i}=\xi\left(v_{i} q\right)$ and $\eta_{j}=\xi\left(\nu_{j} q\right)$, where the points $v_{i} q$ and $v_{j} q$ belong to different $t_{1}$-intervals, and are thus separated by an interval of length at least equal to $t_{2}$. By the condition (2) we then have

$$
\left|\varrho_{i j}\right|=\left|r\left(v_{i} q-v_{j} q\right)\right|<K t_{2}^{-\alpha},
$$

where as usual $K$ denotes an unspecified positive constant. Further, there are less than $L^{2} \leqslant n^{2}\left(m_{1}+1\right)^{2}$ covariances $\varrho_{i j}$. Owing to stationarity some of the $\varrho_{i j}$ are equal, but it is easily seen that this does not affect our argument. It then follows from (22) and (25), using (7) and (10), that we have

$$
\begin{aligned}
& \left|F^{\prime}(h)\right|<K n^{2} m_{1}^{2} t_{2}^{-\alpha} e^{-u^{2}}<K \mu^{\alpha-4 \beta}, \\
& |F(1)-F(0)|=\left|\int_{0}^{1} F^{\prime}(h) d h\right|<K \mu^{\alpha-4 \beta}
\end{aligned}
$$

This holds for any of the $\binom{n}{k}$ terms in the last member of (19), and $F(0)$ will in all cases be given by (21), so that we finally obtain

$$
\left|P\left\{E_{k}\right\}-\binom{n}{k} p^{k}(1-p)^{n-k}\right|<K\binom{n}{k} \mu^{\alpha-4 \beta}<K \mu^{\alpha-(k+4) \beta}
$$

By (8) we have $(k+4) \beta<\alpha$, so that the last member tends to zero as $u \rightarrow \infty$, and the lemma is proved.

## 6. Proof of the case $j=1$ of the theorem

By (18), $p=P\left\{e_{r}\right\}$ is defined as the probability that at least one of the random variables $\xi(0), \xi(q), \xi(2 q), \ldots, \xi\left(m_{1} q\right)$ takes a value exceeding $u$. According to (17), this differs from the probability $P\left\{d_{r}\right\}$ of at least one $\xi_{q}$ upcrossing in the first $t_{1}$-interval by a quantity of the order

$$
O\left(\frac{1}{u} e^{-u^{v / 2}}\right)
$$

## h. Cramér, Stationary stochastic process and high level

By Lemma 3, the latter probability is

$$
P\left\{d_{r}\right\}=q+o(q) .
$$

Thus we obtain from Lemma 5 , observing that by (10) we have $1 / u e^{-u^{2 / 2}}=o(q)$,

$$
P\left\{E_{k}\right\}-\binom{n}{k}[q+o(q)]^{k}[1-q+o(q)]^{n-k} \rightarrow C_{1}
$$

By (10) we have $n q \rightarrow \tau$, and thus

$$
\lim _{u \rightarrow \infty} P\left\{E_{k}\right\}=\frac{\tau^{k}}{k!} e^{-\tau} .
$$

Lemmas 1 and 4 then finally give the relation (6) that was to be proved:

$$
\lim _{u \rightarrow \infty} P_{k}=\frac{\tau^{k}}{k!} e^{-\tau} .
$$

Thus we have proved the simplest case of the theorem, when $j=1$, so that there is only one interval.

## 7. Proof of the general case

The generalization to the case of an arbitrary number $j>1$ of intervals is now simple.

For any $\varepsilon>0$, it follows from the result just proved that, for every $i=1, \ldots, j$, the random variable

$$
N\left(a_{i}+\varepsilon / \mu, b_{i}-\varepsilon / \mu\right),
$$

where $b_{i}-a_{i}=\boldsymbol{\tau}_{i} / \mu$, will be asymptotically Poisson distributed with parameter $\tau_{i}-2 \varepsilon$. In the same way as in the proof of Lemma 5 it is shown that these $j$ variables are asymptotically independent, so that we have

$$
\begin{equation*}
P\left\{N\left(a_{i}+\varepsilon / \mu, b_{i}-\varepsilon / \mu\right)=k_{i} \text { for } i=1, \ldots, j\right\} \rightarrow \prod_{i=1}^{j} \frac{\left(\tau_{i}-2 \varepsilon\right)^{k_{i}}}{k_{i}!} e^{-\left(\tau_{i}-2 \varepsilon\right)} . \tag{26}
\end{equation*}
$$

From the asymptotic Poisson distributions of the variables

$$
N\left(a_{i}, a_{i}+\varepsilon / \mu\right) \quad \text { and } \quad N\left(b_{i}-\varepsilon / \mu, b_{i}\right)
$$

it further follows that, with a probability exceeding $1-2 j \varepsilon$, these variables will ultimately be zero for all $i=1, \ldots, j$. Since $j$ is fixed, and $\varepsilon>0$ is arbitrarily small, the truth of the theorem then follows from (26).

## REFERENCES

1. Berman, S. M., Limit theorems for the maximum term in stationary sequences. Ann. Math. Stat. 35, 502 (1964).
2. Bulinskaja, E. V., On the mean number of crossings of a level by a stationary Gaussian process. Teor. Verojatnost. i Primenen. 6, 474 (1961).
3. Cramér, H., A limit theorem for the maximum values of certain stochastic processes. Teor. Verojatnost. i Primenen. 10, 137 (1965).
4. Craver, H., and Leadbetter, M. R., Stationary and Related Stochastic Processes. To be published by Wiley and Sons, New York.
5. Khintchine, A. Y., Mathematical Methods in the Theory of Queueing. Griffin and Co., London 1960.
6. Kolmogorov, A. N., and Rozanov, Yu. A., On strong mixing conditions for stationary Gaussian processes. Teor. Verojatnost. i Primenen. 5, 222 (1960).
7. Rice, S. O., Distribution of the duration of fades in radio transmission. Bell Syst. Techn. Journ. 37, 581 (1958).
8. Volkonskiv, V. A., and Rozanov, Yu. A., Some limit theorems for random functions. Teor. Verojatnost. i Primenen. 4, 186 (1959) and 6, 202 (1961).

[^0]:    ${ }^{1}$ With respect to the general theory of the normal stationary process and its sample functions we refer to the forthcoming book [4] by Cramér and Leadbetter.

