# SOME PROBLEMS CONCERNING <br> CLUSTER PROCESSES AND OTHER POINT PROCESSES 

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Unless otherwise stated, the material presented in this thesis is the product of my own original research.

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obtain finer results. We believe that we have already partly countered the first objection, in that generality leads to clarity, which is a necessary first step to building more realistic generalized models, and, in addition, we feel that some idea of robustness of already available models is desirable. We fully concur with the second objection, but frequently the general method obtains optimal results (often more simply). Even if it does not, it is worth knowing what it can achieve, in order to elucidate the advantages of tools which exploit the distinctive features of particular models.

The transformation referred to in the opening paragraph will usually take the form of a mapping between measure spaces. Sometimes (rarely) measurability of this mapping is all the 'smoothness' we will require of it. Other times we will demand that it be more coherent, in the sense of preserving the component structure, and as this is usually not provided, we must impose conditions guaranteeing it. Typically we will find an approximating 'coherent' mapping, and the conditions will arise in demonstrating that the differences are irrelevant.

The question arises as to whether our approach can be generalized still further: if there are two or more component processes, can their independence be removed? The answer to this, in some cases, is in the affirmative. For example, it will be clear from the proof of Theorem 4.4.1 that mixing of cluster processes will follow from joint mixing of the centre and subsidiary processes. However, in those situations in which 'coherency' conditions on our mapping must be imposed, it becomes correspondingly more difficult to interpret their meaning, and probably also to verify. For this reason too (ease of writing down conditions), we have required throughout the thesis that our processes be stationary, although often, but not always, this assumption can, with care, be removed.

Chapter One, as well as providing a brief historical introduction and a definition of point processes, aims at setting down those concepts and
classes of point processes which are frequently employed in the body of the dissertation. This still does not render the thesis self-contained, but minimizes to some extent the amount of external referencing required later.

Chapter Two investigates weak forms of asymptotic independence and (ordinary) limit laws for the number of busy servers process associated with the $G / G / \infty$ queue. In particular we demonstrate that stationarity, mixing, strong laws, the central limit theorem and rates of functional convergence derive from those of the arrival and server processes under reasonable conditions. Ergodicity does not, and rates of ordinary convergence seem to depend on a factor associated with the transformation operation.

Functional limit laws for cluster processes are examined in Chapter Three. A sufficient condition for existence of our generalized cluster process is given, and the measurability of the mapping defining the cluster process is investigated. Then, as in Chapter Two, the functional central limit theorem for the cluster process is shown, under suitable conditions predictably more severe than those of Daley (1972) for ordinary convergence, to flow from those of its components; the same is true for functional laws of the iterated logarithm for processes with right-hand clusters, but an unsolved problem is the extension to processes with left-hand clusters. Related topics (functional strong laws, the law of the iterated logarithm for the $G / G / \infty$ queue, and limit laws for doubly stochastic Poisson processes) are then mentioned.

Chapter Four deals with the preservation of strong forms of asymptotic independence under the clustering operation, initially attempted with a view to weakening the theorems in Chapter Three. This has not eventuated, however, but the theorems are of interest in their own right. Strong, $\phi$ and complete mixing possess an increasing degree of uniformity of their asymptotic independence, and this turns out to be a significant factor in their preservation. In particular, it is indicated that $\phi$-mixing may be maintained only under very stringent conditions (bounded clusters), whereas
strong mixing is maintained under considerably milder conditions. It is also argued that the probability generating functional is not an applicable tool in these circumstances, even if the subsidiary processes are independent.

Chapter Five consists of characterization problems associated with renewal point processes, of which the main contribution is contained in a forthcoming paper (Laslett (1975)), in which we conclude that the output process of a finite capacity $G I / M / 1$ queue (i.e., with renewal input) is never renewal. This is not unexpected, since the property of being a renewal process is very much a local one. This subsection is written in the language of queueing theory, unlike the rest of the thesis. We also look briefly at the problem of $n(\geq 2)$ independent, identically distributed point processes being superposed to produce a renewal process. It is conjectured that all processes must then be Poisson, and proved in the case of the superposed processes being alternating renewal; counter-examples in the non-identical case have come from this area.

Chapter Six provides a list of unsolved and partially solved problems and generalizations. In addition to those associated with the bulk of the thesis, it also exhibits identifiability of the cluster structure of a stationary Poisson cluster process from a complete centre process-cluster process record. The problem is included in this chapter, because it is not solved in the generality required of the rest of this thesis, although it is conjectured that it can be.

For record and referral purposes, we have included a relatively extensive bibliography. It will be apparent that some of these references are only of an incidental nature, whereas others are more extensively applicable, in that they are required in the original sections of the thesis.

Once the Introduction and the definition of a cluster process (Sections 3.2 and 4.4) have been assimilated, the chapters may, apart from the unsolved
problems, be read independently. The section 'Symbols and Abbreviations' should be perused before commencing on the thesis itself.

## SYMBOLS AND ABBREVIATIONS

Some of the details which follow are almost conventional and are given simply to remove any possible ambiguity.

At times we write $a b / c d$ rather than the more cumbersome $(a b) /(c d)$.

| $B(R)$ | the usual Borel $\sigma$-field of the real line $R$ |
| :---: | :---: |
| $B(X)$ | the Borel $\sigma$-field generated by the open sets of a |
|  | specified topology on the space $X$ |
| card\{ $\cdot\}$ | cardinality of the set $\{\cdot\}$ |
| ch.fl. | characteristic functional |
| $\operatorname{cov}\{\cdot, \bullet\}$ | the covariance with respect to the appropriate probability |
|  | measure |
| $E\{\cdot\}$ | expectation with respect to the appropriate probability |
|  | measure |
| i.i.d. | independent and identically distributed |
| $P$ | (the probability measure of) a point process or other |
|  | specified process |
| $P r$ | probability measure on the appropriate measure space |
| p.g.fl. | probability generating functional |
| $\Pi_{i}{ }_{i}$ | product of real or complex numbers $z_{i}$ |
| $R$ | the real line |
| $R_{+}$ | $[0, \infty)$ |
| $R(X)$ | a specified ring of subsets of a set $X$ |
| $\sigma(X)$ | a specified $\sigma$-field of subsets of a set $X$ |
| $\sigma(X)^{T}$ | the restricted product $\sigma$-field on $X^{T^{\prime}}$, where |
|  | $T \subset T^{\prime} \subset R$, i.e., the smallest $\sigma$-field containing sets |
|  | of the form $\left\{x(\cdot) \in X^{T^{\prime}}: x\left(t_{1}\right) \in C_{1}, \ldots, x\left(t_{n}\right) \in C_{n}\right\}$ for |
|  | $n \in Z_{+}, t_{1}, \ldots, t_{n} \in T$ and $C_{i} \in \sigma(X)$. Unless |

otherwise specified $T^{\prime}=T$, and then this corresponds to the usual product $\sigma$-field

| $\sigma_{\text {gen }}$. ${ }^{(C)}$ | the minimal $\sigma$-field generated by sets in the class $C$ |
| :---: | :---: |
| $\sigma\left(x_{1}\right) \times \sigma\left(x_{2}\right)$ | the product $\sigma$-field of $\sigma\left(x_{1}\right)$ and $\sigma\left(X_{2}\right)$, i.e., |
|  | $\sigma_{\text {gen. }}\left\{C_{1} \times C_{2}: C_{1} \in \sigma\left(X_{1}\right), c_{2} \in \sigma\left(X_{2}\right)\right\}$. |
| supp.f | the support of the function $f$ |
| $\operatorname{Var}\{\cdot\}$ | the variance with respect to the appropriate probability |
|  | measure |
| $X^{T}$ | the set of mappings $T \rightarrow X$. We will occasionally abuse |
|  | this notation e.g. if $\operatorname{card}\{T\}=n$, we may write $X^{n}$ and |
|  | the members of $X^{n}$ may be indexed other than by |
|  | 1, ..., $n$. |
| 2 | the set of all integers $\{0, \pm 1, \ldots\}$ |
| $2_{+}$ | the set of non-negative integers $\{0,1, \ldots\}$ |
| $\bar{Z}_{+}$ | $Z_{+} \cup\{\infty\}$ |
| $A+t$ | $\{x+t: x \in A\}$ where $A$ is a subset of $R, t \in R$ |
| $A^{c}$ | complement of the set $A$ |
| $\bar{A}$ | closure of, the set $A$ |
| $\Sigma_{i} A_{i}$ | disjoint union of sets $A_{i}$ |
| $1_{A}(\cdot)$ | indicator function of the set $A$ |
| $\|A\|$ | Lebesgue measure of the set $A$, where $A \in B(R)$ |
| $\|z\|$ | absolute value of the real or complex number $z$ |
| a.b | scalar product of vectors $a$ and $b$ |
| $f \circ g$ | composition of mappings $f$ and $g$ |



## CHAPTER 1

## INTRODUCTION

### 1.1. Historical background

It is by no means an accident that many of the major works on point processes do not contain an historical account of the subject - the task of compiling such a survey would be formidable indeed. Hence we will attempt only a brief outline of the development of the subject; further references specific to the topics discussed in this thesis will be cited in the appropriate places. In this introduction we will not hesitate to draw on D. Vere-Jones' (1973) excellent account of the early history of point processes.

Fluctuations in the counts of objects in various situations seem to have been recognized as a stochastic problem since the latter half of the nineteenth century, although distributions other than the simple Poisson do not occur until the early 1920 's, where they appear in problems on accident proneness (Greenwood and Yule (1920)) and contagion (Polya (1931)) Neyman (1939) introduced the important idea of clustering.

Population processes are intricately part of counting processes, and probably originate from the Bienaymé-Galton-Watson process. Development of the theory (for finite populations) begins with Feller (1950) and Bartlett (1954), and is extended by Moyal (1962) and Harris (1963). The study of particle showers (Bhabha (1950), Ramakrishnan (1950)) should also be mentioned as an early stimulus.

Considerations of distributions of interval lengths can be traced back as far as life-tables themselves, so that, sumprisingly, renewal theory seems to have a longer history than counting problems. Attention seems to have been paid to the problem of relating the two properties, counts and intervals, of renewal processes on the line well before the 1940's (see

Lotka (1939)). The first treatment of processes with correlated intervals is by Wold (1949), to whom the term "point process" is due. He also made the first progress with processes with an infinite number of points.

Queueing theory has provided great impetus to the theory, particularly since Khinchin's monograph (1955). Palm (1943) attempted (without complete success) the first limit theorem for point processes (Poisson processes as the limit of $n$ superposed processes, $n \rightarrow \infty$ ), as well as probing the question of Palm measures and functions. An account of the theory of superpositions may be read in Çinlar (1972); Palm functions and measures have been studied by Khinchin (1955), Slivnyak (1962), Ryll-Nardzewski (1961) and others.

During the $1960^{\prime}$ 's and 1970's the theory has evolved rapidly, with expositions on stationarity and general properties (Matthes (1963a), Beutler and Leneman (1966)), statistical analysis (Cox and Lewis (1966)), spectral theory (Bartlett (1963), Daley (1971), Vere-Jones (1974)), Palm-Khinchin theory (Neveu (1968), Papangelou (1974), Leadbetter (1972a), Jagers (1973)), infinitely divisible point processes (Kerstan and Matthes (1964), Lee (1967), (1968)), cluster models (Neyman and Scott (1958), Lewis (1964a,b), (1969)) and probability generating functionals (Vere-Jones (1968), Westcott (1972)). Generalizations to multivariate (Milne (1971)), multidimensional (Fisher (1972) for a survey) and more abstract point processes (Mecke (1967), Kallenberg (1973), Jagers (1974) and Krickeberg (1974)) have been effected.

Of course in the above list we have omitted many topics and references (which are not intended to indicate priority), and more complete coverage can be obtained from Lewis (1972) and Kerstan, Matthes and Mecke (1974) and references therein. Daley and Milne (1973) have compiled a bibliography.

The rest of this chapter sets out the basic definitions and properties of point processes and weak convergence in a more coherent manner than their piecemeal introduction in the text of the thesis allows. After the formal
definition in Section 1.2, we catalogue in Section 1.3 some known properties for easy reference later. Many of these have become part of the folklore, in which case reference to their source is omitted. The definitions of a few special processes needed in this thesis are given in Section 1.4, while weak convergence is summarised in Section 1.5.

### 1.2. The definition of a point process

Let $N$ denote the set of non-negative integer-valued measures on the real line $R$ which are finite on compact sets, and $\sigma(N)$ the $\sigma$-field generated by $\left\{\{N(A) \leq k\}, k \in Z_{+}\right.$, bounded $\left.A \in B(R)\right\}$. Then we define a point process to be a probability measure on $(N, \sigma(N))$. Throughout this thesis we will reserve $N$ (possibly subscripted) for members of $N ; P$ (possibly subscripted) will mean a point process, or other (specified) random process.

Denote by ( $\Omega, F, P r$ ) an (arbitrary) probability space, and $\eta$ (a random " $N$ ") a measurable mapping from ( $\Omega, F, P r$ ) to ( $N, \sigma(N)$ ) . Then $P=P \eta^{-1}$ specifies a point process, but on some occasions, in particular those in which we wish to emphasize that a fixed underlying probability space is necessary for our arguments to be meaningful, we will refer to $\eta$ as the point process.

The question of existence of measures $P$ arises. As for ordinary stochastic processes, the problem is solved by extending the finite dimensional distributions of the process.

THEOREM 1.2.1 (Moyal (1962), Harris (1963), p. 53; Nawrotzki (1962)). Given a set of functions
$\left\{p\left(A_{1}, \ldots, A_{8} ; k_{1}, \ldots, k_{8}\right): s \geq 1\right.$ and integral, $k_{i} \in Z_{+}$, $A_{i}$ bounded $\left.\in B(R), 1 \leq i \leq s\right\}$,
(i) $p\left(A_{1}, \ldots, A_{8} ; k_{1}, \ldots, k_{8}\right)=p\left(A_{i_{1}}, \ldots, A_{i_{8}} ; k_{i_{1}}, \ldots, k_{i_{s}}\right)$ for every permutation $\left(i_{1}, \ldots, i_{8}\right)$ of $(1, \ldots, s)$;
(ii) $p\left(A_{1}, \ldots, A_{8} ; k_{1}, \ldots, k_{8}\right) \geq 0$ and

$$
\begin{aligned}
\sum_{k_{s}=0}^{\infty} p\left(A_{1}, \ldots, A_{s-1}, A_{s}\right. & \left.; k_{1}, \ldots, k_{s-1}, k_{s}\right) \\
& =p\left(A_{1}, \ldots, A_{s-1} ; k_{1}, \ldots, k_{s-1}\right) ;
\end{aligned}
$$

(iii) $\sum_{k_{1}=0}^{\infty} p\left(A_{1}, k_{1}\right)=1$;
(iv) whenever $A_{1}, \ldots, A_{s}$ are disjoint,

$$
p\left(\bigcup_{i=1}^{s} A_{i} ; k\right)=\sum_{k_{1}+\ldots+k_{s}=k} p\left(A_{1}, \ldots, A_{s} ; k_{1}, \ldots, k_{s}\right)
$$

and

$$
p\left(\bigcup_{i=1}^{s} A_{i}, A_{1}, \ldots, A_{s} ; k, k_{1}, \ldots, k_{s}\right)=0
$$

unless $\sum_{i=1}^{s} k_{i}=k$, when it equals
$p\left(A_{1}, \ldots, A_{s} ; k_{1}, \ldots, k_{s}\right) ;$
(v) $p\left(A_{k} ; 0\right) \uparrow 1$ whenever $A_{k} \downarrow \varnothing$;
then there exists a unique probability measure $P$ on $(N, \sigma(N))$ for which

$$
P\left\{N: N\left(A_{1}\right)=k_{1}, \ldots, N\left(A_{s}\right)=k_{s}\right\}=p\left(A_{1}, \ldots, A_{s} ; k_{1}, \ldots, k_{s}\right) .
$$

It will be required to characterise $P$ more finely (Chapter 6), Since the class of bounded half-open intervals with rational endpoints generates $B(R)$, we may restrict the $A_{i}$ 's to this class. Further, suppose we have functions $p_{0}\left(A_{1}, \ldots, A_{8} ; k_{1}, \ldots, k_{8}\right)$ defined whenever the $A^{\prime}$ 's are disjoint; then the functions $p_{0}$ may be regarded as defining a joint distribution for random variables $\xi\left(A_{1}\right), \ldots, \xi\left(A_{g}\right)$ defined on a space
( $\Omega, F, P_{r}$ ) . Suppose ( $i v$ ) is modified as follows:
(iv)' Let $A_{1}, \ldots, A_{s}$ be any mutually disjoint sets in $B(R)$, and suppose $A_{i}=\sum_{k=1}^{n_{i}} A_{i k}, A_{i k} \in B(R)$, for some $n_{i}=1,2, \ldots$; then the joint distribution of $\xi\left(A_{1}\right), \ldots, \xi\left(A_{8}\right)$ is the same as the joint distribution of $\sum_{k=1}^{n_{1}} \xi\left(A_{\perp k}\right), \ldots, \sum_{k=1}^{n_{s}^{s}} \xi\left(A_{s k}\right)$.

Then
THEOREM 1.2.2 (Harris (1963), p. 54). A set of functions $\left\{p_{0}\left(A_{1}, \ldots, A_{s} ; k_{1}, \ldots, k_{s}\right\}: s \geq 1\right.$ and integral, $k_{i} \in Z_{+}$,

$$
\text { bounded diajoint } \left.A_{i} \in B(R), \perp \leq i \leq s\right\}
$$

satisfying (i) and (ii) of Theorem 1.2.1 whenever the $A_{i}$ 's are disjoint, and (iii), (iv)' and (v) can be uniquely extended to functions $p\left(A_{1}, \ldots, A_{8} ; k_{1}, \ldots, k_{8}\right)$ satisfying (i)-(v) of Theorem 1.2.1, and agreeing with the $p_{0}(\cdot ; \cdot)$ whenever the $A_{i}$ are disjoint.

In particular, if our functions $p_{0}$ are generated by a point process, then the requirements of Theorem 1.2 .2 are met. Also, we only need to know $p_{0}$ for $A_{i}$ 's being bounded half-open intervals with rational endpoints.

Although our processes will only be on the real line, there is no difficulty in extending the preceding concepts to $R^{n}$ (e.g. Fisher (1972)) and Polish spaces (Jagers (1974)).

We may associate with each $N \in N$ the distance $t_{j}(N)$ of the $j$ th point of $N$ from the origin (commonly referred to as the $j$ th epoch) defined as follows:

$$
\begin{align*}
t_{j} \equiv t_{j}(N) & =\inf \{y>0: N(0, y] \geq j\}, j=1,2, \ldots \\
& =\sup \{y \leq 0: N[y, 0]>-j\}, j=0,-1, \ldots \tag{2.1}
\end{align*}
$$

Note that $t_{j}: N \rightarrow R$ is measurable, since

$$
\begin{align*}
& \left\{N: t_{j} \leq x\right\}=\{N: N(0, x] \geq j\} \in \sigma(N), j=1,2, \ldots, \\
& \left\{N: t_{j}<x\right\}=\{N: N[x, 0] \leq-j\} \in \sigma(N), j=0,-1,-2, \ldots \tag{2.2}
\end{align*}
$$

For any $N \in N$, the set $\left\{t_{j}(N)\right\}$ is a finite or countably infinite subset of $R$, multiple points included but with no finite limit points, and satisfying the inequalities

$$
\begin{equation*}
\cdots \leq t_{-1} \leq t_{0} \leq 0<t_{1} \leq t_{2} \cdots \cdots \tag{2.3}
\end{equation*}
$$

(This is not the conventional indexing (see Daley and Vere-Jones (1972), p. 308), but is employed to simplify the notation in Chapters 2 and 3.) If we denote the class of all such subsets of $R$ by $R_{t}^{Z}$, then we can define $n: R_{t}^{Z} \rightarrow N$ by

$$
\begin{equation*}
\eta(A)(\mathrm{t})=\operatorname{card}\left\{j ; t_{j} \in \mathrm{t} \cap A\right\}, A \in B(R), \mathrm{t} \in R_{t}^{Z} ; \tag{2.4}
\end{equation*}
$$

if also we introduce the $\sigma$-algebra $\sigma\left(R_{t}^{2}\right)$ generated by the sets $\{t: \eta(A)(t) \leq k\}, A \in B(R), k \in Z_{+}$, i.e., such that the mappings $\eta$ are measurable, then we may define a point process as a probability measure on $\left(R_{t}^{Z}, \sigma\left(R_{t}^{Z}\right)\right)$. It is intuitively clear that we may specify a point process also via the intervals $\left\{\Pi_{j}\right\} \equiv\left\{t_{j}-t_{j-1}\right\}$ and $t_{0}$ and $t_{1}$, but, throughout this thesis, unless explicitly stated otherwise, we will consider our event epochs as generated in (2.1). We do this for two reasons: not only is $(N, \sigma(N))$ more amenable to generalization, but at times it will be the convenient to introduce the space $N_{\infty}$, the set of all non-negative integer-valued measures on $R$ which may also be infinite on compact sets. Finally we mention that Matthes and others have extended the scheme to marked point processes, where each point $t_{i}$ is associated with a mark $k_{i}$
from a fixed measure space $[K, \sigma(K)]$. As this idea will only be employed circumspectly in this thesis, we will not expand on it. An account may be found in Kerstan, Matthes and Mecke (1974).

### 1.3. Basic properties of point processes

Define the translation operator $T_{y}: N \rightarrow N, y \in R$ by $T_{y} N(\cdot)=N(\cdot+y)$. A point process $P$ is (strictly) stationary if

$$
\begin{equation*}
P\left(T_{y} C\right)=P(C), \text { all } \quad C \in \sigma(N) \tag{3.1}
\end{equation*}
$$

If $P$ is stationary, $P\{N(R)=0$ or $\infty\}=1$, and $P$ has no atoms in the sense that $P\{N(\{x\})>0\}=0$ for all singleton point sets $\{x\}$ (RyllNardzewski (1961)).

A stationary point process $P$ is ergodic if any member $G$ of the invariant $\sigma$-field $T_{1}=\left\{G \in \sigma(N): T_{y}^{-1} G=G\right\}$ satisfies $P(G)=0$ or 1 . Rosenblatt (1962) demonstrates that ergodicity may be charactized by

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} \tau^{-1} \int_{0}^{\tau} P\left(C \cap T_{y} D\right) d y=P(C) P(D) \text {, all } C, D \in \sigma(N) \text {. } \tag{3.2}
\end{equation*}
$$

$P$ is weakly mixing ${ }^{2}$ if

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} \tau^{-1} \int_{0}^{\tau}\left|P\left(C \cap T_{y} D\right)-P(C) P(D)\right| d y=0, \text { all } C, D \in \sigma(N), \tag{3.3}
\end{equation*}
$$

and mixing if

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} P\left(C \cap T_{\tau} D\right)=P(C) P(D) \text {, all } C, D \in \sigma(N) \text {. } \tag{3.4}
\end{equation*}
$$

Note that mixing $\Rightarrow$ weak mixing $\Rightarrow$ ergodicity $\Rightarrow$ stationarity.
We will also use the following stronger forms of asymptotic independence. Let

$$
\begin{equation*}
\sigma(N(B)) \equiv \sigma_{\text {gen. }}\left\{\{N(A)=k\}, A \in B(R) \cap B, k \in Z_{+}\right\}, B \in B(R) . \tag{3.5}
\end{equation*}
$$

A stationary point process $P$ is strong mixing with rate $\alpha(\tau)$ if

[^0]\[

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \alpha(\tau) \tag{3.6}
\end{equation*}
$$

\]

for all $C \in \sigma(N(-\infty, t]), D \in \sigma(N(t+\tau, \infty)), \quad t \in R, \tau \geq 0$. Here $\alpha:[0, \infty) \rightarrow[0,1]$ is a monotone decreasing function satisfying
$\lim \alpha(\tau)=0$. Let $C$ and $D$ be as in (3.6). Then $P$ is $\phi-m i x i n g$ if $\tau \rightarrow \infty$

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \phi(\tau) P(C) \tag{3.7}
\end{equation*}
$$

and completely mixing (with rate $\gamma(\tau)$ ) if

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \gamma(\tau) P(C) P(D) \tag{3.8}
\end{equation*}
$$

where $\phi$ and $\gamma$ have the same properties as $\alpha$ in (3.6). Clearly

$$
\text { complete mixing } \Rightarrow \phi \text {-mixing } \Rightarrow \text { strong mixing } \Rightarrow \text { mixing. }
$$

For a real valued function $g: N \rightarrow R$, we will define the expectation $E$ of $g$ by

$$
\begin{equation*}
E\{g\} \equiv E_{P}\{g\}=\int_{N} g(N) d P(N) \tag{3.9}
\end{equation*}
$$

Hence we can define the first moment measure $M$ of a point process by $M(\cdot)=E\{N(\cdot)\}$, where existence of $M$ is taken to mean $M(A)<\infty$ for all bounded Borel sets $A$. Clearly $M$ is indeed a measure. If $P$ is stationary, then easily $M(A)=m|A|$, where $m=E N(0,1]$ is called the intensity of the point process. Higher moments are defined as

$$
\begin{equation*}
M_{r}\left(A_{1} \times \ldots \times A_{r}\right)=E\left\{N\left(A_{1}\right), \ldots N\left(A_{r}\right)\right\}, A_{i} \in B(R), 1 \leq i \leq r \tag{3.10}
\end{equation*}
$$

and are easily shown to be measures in $R^{r}$. The cumulant measure $C_{2}$
exists if $M_{2}$ does, and is defined by

$$
\begin{equation*}
C_{2}\left(A_{1} \times A_{2}\right)=M_{2}\left(A_{1} \times A_{2}\right)-M\left(A_{1}\right) M\left(A_{2}\right) \equiv \operatorname{cov}\left(N\left(A_{1}\right), N\left(A_{2}\right)\right) \tag{3.11}
\end{equation*}
$$

$C_{2}$ may be a signed measure. If $P$ is stationary, and $M_{r}(\cdot)$ exists, then $M_{r}(\cdot)$ is stationary in the sense that $M_{p}\left\{\left(A_{1}+x\right) \times \ldots \times\left(A_{p}+x\right)\right\}$ is independent of $x \in R$.

A point process is defined to be weakly stationary if $M_{1}$ and $M_{2}$
exist and are stationary in the above sense (Daley (1971)). Important for us is the fact that the cumulant measure $C_{2}$ of a weakly stationary point
process may be decomposed into Lebesgue measure and the reduced covariance measure $C(\cdot)$ in the following way (using differential notation)

$$
\begin{equation*}
c_{2}((t+d t) \times(t+u+d u))=C(d u) d t, \tag{3.12}
\end{equation*}
$$

and in particular (Daley and Vere-Jones (1972), p. 323)

$$
\begin{equation*}
\operatorname{Var}(N(0, y])=c_{2}((0, y] \times(0, y])=\int_{-y}^{+y}(y-|u|) C(d u) \tag{3.13}
\end{equation*}
$$

The superposition of $n$ independent point processes $P_{1}, \ldots, P_{n}$ is intuitively the overlaying of all the points on one line, but may be specified rigorously via the mapping $\eta: N^{n} \rightarrow N$ defined by

$$
\begin{equation*}
n\left(\left\{N_{i}\right\}\right)=\sum_{i=1}^{n} N_{i} . \tag{3.14}
\end{equation*}
$$

Then $P=\left(\chi_{i=1}^{n} P_{i}\right) n^{-1}$ is the process of superpositions. Note that $\eta$ may be regarded as the sum of $n$ coordinate mappings $\eta_{i}$ each defined on $N^{n}$. It will be convenient on occasions to conform to the usual practice of supposing the summands to be random measures $\eta_{i}:\left(\Omega, F, P_{r}\right) \rightarrow(N, \sigma(N))$ on a fixed space $\Omega$, and defining their superposition $\eta$ by $n(\omega)=\sum_{i=1}^{n} \eta_{i}(\omega)$. This scheme allows for dependent $\eta_{i}$.

We recall that for functions $f: R \rightarrow R$ of compact support,

$$
\begin{equation*}
\int_{-\infty}^{+\infty} f(t) d N(t)=\Sigma_{i} f\left(t_{i}\right) \tag{3.15}
\end{equation*}
$$

is well defined. By Fubini's theorem

$$
\begin{equation*}
E\left\{\int_{-\infty}^{+\infty} f(t) d N(t)\right\}=\int_{-\infty}^{+\infty} f(t) d M(t) \tag{3.16}
\end{equation*}
$$

if $\int_{-\infty}^{+\infty}|f(t)| d M(t)<\infty$.
Supposing $\log \xi(t)$ to have compact support, we may define the probability generating functional (p.g.fl.) of a point process $P$ by

$$
\begin{equation*}
G[\xi]=E\left\{\exp \int_{-\infty}^{+\infty} \log \xi(t) d N(t)\right\} . \tag{3.17}
\end{equation*}
$$

We will denote by $V$ the class of measurable functions $\xi: R \rightarrow[0,1]$ such that $1-\xi$ has compact support. $G$ exists if $\xi$ is in this class, or if $P\{N(R)<\infty\}=1$. For an extensive discussion of the p.g.fl. see Westcott (1972). Note that if $P_{1}, \ldots, P_{n}$ are $n$ independent point processes with p.g.fls. $G_{1}, \ldots, G_{n}$, then the superposition $P$ has p.g.fi. $G[\xi]=\prod_{i=1}^{n} G_{i}[\xi]$.

A point process $P$ is said to be a.s. orderly (or without multiple points) if

$$
\begin{equation*}
P\{N(\{x\})=0 \text { or } 1 \text { (all } x)\}=1, \tag{3.18}
\end{equation*}
$$

and a stationary process $P$ is orderly (or analytically orderly) if

$$
\begin{equation*}
P\{N(0, h] \geq 2\}=o(h) \quad(h+0) . \tag{3.19}
\end{equation*}
$$

If $P$ is stationary with finite intensity $m$, then analytically orderly $\Leftrightarrow$ a.s. orderly ( $\Leftrightarrow$ is Dobrushin's Lemma). We shall also require Khinchin's existence theorem, namely that for a stationary point process $P$ the rate

$$
\begin{equation*}
\lambda=\lim _{h \ngtr 0} P\{N(0, h]>0\} / h \tag{3.20}
\end{equation*}
$$

exists, although it may be infinite. An interesting relation exists between $m$ and $\lambda$ : for a stationary point process, the 'batch-size' distribution for the number of events in a multiple occurrence exists, and has mean $m / \lambda$. Hence $\lambda \leq m$, and a necessary and sufficient condition for $m=\lambda$ is that the process be orderly.

An important recent discovery is the Rényi-Mönch-Kallenberg Theorem:
if $P$ is a.s. orderly, then it is determined by knowing

$$
\begin{equation*}
\phi(I)=P\{N(I)=0\} \tag{3.21}
\end{equation*}
$$

for all $I$ in the semi-ring generated by half-open intervals. Kurtz (1974) has characterized the functions $\phi$ in terms of a property he calls complete monotonicity.

Further properties may be found in the references cited in Section 1.1.

### 1.4. Some examples of point processes

The Poisson process $\eta$ with parameter $\lambda(\cdot)$ (a non-atomic Borel measure) is defined by
(i) For $A \in B(R)$ such that $\lambda(A)<\infty, \eta(A)$ has a Poisson distribution with mean $\lambda(A)$. Otherwise $\eta(A)=\infty$ with probability one.
(ii) $\eta$ is completely random, i.e., $\eta\left(A_{1}\right), \ldots, \eta\left(A_{n}\right)$ are mutually
independent for all finite collections of disjoint sets
$A_{1}, \ldots, A_{n} \in B(R)$.
The Poisson process has been characterized in many ways: Prékopa (1957a,b) showed that a point process is Poisson if and only if it is atomless, completely random and has no multiple points. On the other hand, if
$\operatorname{Pr}\{n(I)=0\}=e^{-\lambda(I)}$ and $\operatorname{Pr}\{n(I) \geq 2\}=o(\lambda(I)), \lambda(I) \downarrow 0$, (4.1) for some non-atomic $\lambda(\cdot)$ and all $I$ consisting of a finite union of intervals, then $\eta$ is Poisson (Rényi (1967)). If $\lambda(A)=\lambda|A|$, some $\lambda<\infty$, we recover the stationary Poisson process.

If the parameter measure of a Poisson process $\eta$ is taken to be the realization of a random measure $\Lambda(\cdot)$, we obtain the doubly stochastic Poisson process. Kingman (1964) characterized doubly stochastic Poisson processes as stationary Poisson processes with unit parameter subjected to a random change of time independent of the original process, an important and too often neglected result.

The best example of a point process which is generated by its interval properties is the renewal process, which starts at time 0 and has i.i.d. inter-event times. The stationary renewal process has the time $t_{1}$ to the
first event delayed according to the distribution

$$
\begin{equation*}
\operatorname{Pr}\left\{t_{1} \leq u\right\}=\lambda \int_{0}^{u}\{1-F(x)\} d x \tag{4.2}
\end{equation*}
$$

where $F$ is the distribution function of inter-event times, and $\lambda^{-1}$ its (finite) mean.

Alternating renewal processes have independent inter-event times, but their distribution switches successively from one lifetime distribution function $F_{1}$ to another $F_{2}$. The concept may be generalised. Processes with consecutive pairs of points equidistant we will call deterministic, and, as a general rule, are very useful for providing counter-examples.

A point process $P$ is infinitely divisible if for every $n=1,2, \ldots$ it may be represented as the superposition of $n$ i.i.d. point processes $P_{n, i} \quad(i=1, \ldots, n)$. Other definitions are possible. Since such processes are the subject of a book by Kerstan, Matthes and Mecke (1974), they will not be discussed in this work, except when they occur incidentally as Poisson cluster processes (see the cited reference).

Point processes may be generated in many other ways: e.g. as the transition times of Markov processes (Rudemo (1973)), or by level crossings (Leadbetter (1972b)). We shall not attempt to expand on these, as this thesis is not concerned with specific point processes, except for illustrative purposes.

### 1.5. A summary of some weak convergence concepts

If probability measures $P, P_{n}$ on a separable metric space $S$ with metric $d$ (referred to as $(S, d)$ ) and Borel sets $B(S)$ satisfy

$$
\begin{equation*}
\int_{S} f d P_{n} \rightarrow \int_{S} f d P \quad(n \rightarrow \infty) \tag{5,1}
\end{equation*}
$$

for every bounded, continuous real function $f$ on $S$, then we say $P_{n}$
converges weakly to $P$ and write $P_{n} \Rightarrow P$. The theory of weak convergence on metric spaces is outlined in Billingsley (1968), and for most theorems we need we will refer to this text. However, we mention the continuous mapping theorem (Theorem 5.1 of Billingsley (1968)), which we often use. Let $h:(S, d) \rightarrow\left(S^{\prime}, d^{\prime}\right)$ be a measurable mapping with a set of discontinuities $D_{h}$.

CONTINUOUS MAPPING THEOREM. If $P_{n} \Rightarrow P$ and $P\left(D_{h}=0\right)$, then $P_{n} h^{-1} \Rightarrow P h^{-1}$.

Various metrics are associated with weak convergence: since we shall use almost all of them, we give a brief summary (adapted from Whitt (1974b)).

Let $\operatorname{RV} \equiv \operatorname{RV}(S, d)$ denote the set of all random variables $(\Omega, F, \operatorname{Pr}) \rightarrow(S, d)$ defined on a fixed space $\Omega$, and $P(S, d)$ the space of all probability measures on (S, B(S)).

Firstly, for $X_{1}, X_{2} \in R V$, define

$$
\begin{equation*}
\alpha\left(X_{1}, X_{2}\right)=\inf \left\{\varepsilon \geq 0: \operatorname{Pr}\left[d\left(X_{1}, X_{2}\right) \geq \varepsilon\right] \leq \varepsilon\right\} \ldots \tag{5.2}
\end{equation*}
$$

$\alpha$ corresponds to convergence in probability.
Now, for any $\varepsilon>0$, define

$$
\begin{equation*}
A^{\varepsilon}=\{y: d(x, y)<\varepsilon \text { for some } x \in A\} . \tag{5.3}
\end{equation*}
$$

The Prohorov metric $\rho$ which induces the topology of weak convergence on $P$ is given by

$$
\begin{align*}
& \rho\left(P_{1}, P_{2}\right)=\max \left\{\gamma\left(P_{1}, P_{2}\right), \gamma\left(P_{2}, P_{1}\right)\right\} \\
& \gamma\left(P_{1}, P_{2}\right)=\inf \left\{\varepsilon \geq 0: P_{1}(F) \leq \varepsilon+P_{2}\left(F^{c}\right), F \text { closed }\right\} . \tag{5.4}
\end{align*}
$$

Note immediately that $\rho$ may be regarded as acting on $R V$, by setting $\rho\left(X_{1}, X_{2}\right)=\rho\left(\operatorname{PrX}_{1}^{-1}, \operatorname{Pr} X_{2}^{-1}\right)$, but is now only a pseudometric. Clearly $\rho\left(X_{1}, X_{2}\right) \leq \alpha\left(X_{1}, X_{2}\right)$. Dudley (1968) proved that $\gamma\left(P_{1}, P_{2}\right)=\gamma\left(P_{2}, P_{1}\right)$, provided $P_{1}(S)=P_{2}(S)$ and there are no restrictions on closed $F^{\prime}$ s,
conditions satisfied here. The Lévy metric $\lambda$ on $P(R)$ is the Prohorov metric restricted to closed sets of the form $(-\infty, x]$. Clearly $\lambda \leq \rho$, but it is well known that $\lambda$ characterises weak convergence (e.g. Loève (1960), p. 215).

The supremum metric $\sigma$ is defined by

$$
\begin{equation*}
\sigma\left(P_{1}, P_{2}\right)=\sup \left\{\left|P_{1}(A)-P_{2}(A)\right|, A \in B(S)\right\} \tag{5.5}
\end{equation*}
$$

$\sigma$ generates a stronger topology on $P$ than $\rho$. We shall also refer to the restriction of $\sigma$ on $P(R)$ to sets of the form $(-\infty, x]$ as the supremum metric $\nu$. Now $\lambda \leq \nu$, but $\nu$ also metrizes weak convergence to those $P$ in $P(R)$ with continuous distribution functions.

Dudley (1966) demonstrates that the dual bounded Lipschitz metric $\beta$ also induces the weak convergence topology. A function $f:(S, d) \rightarrow R$ is Lipschitz if

$$
\begin{equation*}
\|f\|_{1}=\sup _{x \neq y}\{|f(x)-f(y)| / d(x, y)\}<\infty . \tag{5.6}
\end{equation*}
$$

For any such $f$ we may associate the norm

$$
\begin{equation*}
\|f\|=\|f\|_{1}+\|f\|_{\infty} \tag{5.7}
\end{equation*}
$$

where $\|f\|_{\infty}=\sup \{|f(x)|, x \in S\}$. Then $\beta$ is defined by

$$
\begin{equation*}
\beta\left(P_{1}, P_{2}\right)=\sup \left\{\left|\int f d P_{1}-\int f d P_{2}\right|,\|f\| \leq 1\right\} . \tag{5.8}
\end{equation*}
$$

Also, according to Dudley (1968), $\beta \leq 2 \rho$, so that we may summarize the inter-relationship between the metrics (regarded as pseudo-metrics on RV ) by


Closely related to weak convergence is vague convergence, a concept of great importance to point processes: a sequence $\left\{N_{n}\right\}$ of measures $N_{n} \in N$ converges vaguely to $N \in N$ if

$$
\begin{equation*}
\int f d N_{n} \rightarrow \int f d N \quad(n \rightarrow \infty) \tag{5.9}
\end{equation*}
$$

for every $f \in C_{K}$, the set of continuous functions $f: R \rightarrow R$ with compact support. This convergence generates the vague topology on $N$, a basis for which is given by

$$
\begin{equation*}
\left\{N \in N:\left|\int f_{j} d N-\int f_{j} d N_{1}\right|<\varepsilon, 1 \leq j \leq n\right\} \tag{5.10}
\end{equation*}
$$

$n=1,2, \ldots, f_{j} \in C_{K}, N_{1} \in N, \varepsilon>0$. The vague topology is metrizable (by the Prohorov metric) and renders $N$ separable. Weak convergence of probability measures on $(N, B(N))$ may then be studied (Jagers (1974)). The vague topology is a natural topology for $N$, since, if $B(N)$ denotes the Borel $\sigma$-field generated by the open sets of the vague topology, $B(N)=\sigma(N)$ by Proposition 1.1 of Jagers (1974).

Weak convergence may also be studied on the space $C[0,1]$ of continuous functions on $[0,1]$, where $C[0,1]$ has the uniform topology $(U)$ induced by the metric $\rho$ (no confusion with the Prohorov metric $\rho$ should arise)

$$
\begin{equation*}
\rho(x, y)=\sup _{0 \leq t \leq 1}|x(t)-y(t)|, x, y \in C \tag{5.11}
\end{equation*}
$$

More appropriate to the study of point processes is the space $D[0,1]$ of functions on $[0,1]$ which are right-continuous and have left-hand limits. $D$ will be endowed with the $J_{1}$-topology induced by the Skorokhod metric (using $I \equiv I(t)=t$ and $y \circ \lambda(t)=y(\lambda(t))$ )

$$
\begin{equation*}
d(x, y)=\inf _{\lambda \in \Lambda}\{\max (\rho(\lambda, I), \rho(x, y \circ \lambda))\} \tag{5.12}
\end{equation*}
$$

where $\Lambda$ consists of all continuous strictly increasing maps of $[0,1]$ onto itself. We need note only that $d(x, y) \leq \rho(x, y), x, y \in D$, but $d\left(x_{n}, x\right) \rightarrow 0 \Rightarrow \rho\left(x_{n}, x\right) \rightarrow 0 \quad(n \rightarrow \infty)$ if $x$ is continuous.

Weak convergence of probability measures on $C$ and $D$ is characterize by convergence of the finite-dimensional distributions and tightness (Billingsley (1968), Theorems 8.1 and 15.1). Necessary and sufficient conditions for tightness of a sequence $\left\{P_{n}\right\}$ of probability measures on $C$
may be established via the Arzelà-Ascoli theorem in terms of the modulus of continuity of $x \in C$,

$$
\begin{equation*}
w_{x}(\delta)=\sup _{|s-t|<\delta}|x(s)-x(t)| \tag{5.13}
\end{equation*}
$$

The following elementary inequality exists between $w$ and $\rho$ (Billingsley (1968), p. 220),

$$
\begin{equation*}
\left|w_{x}(\delta)-\omega_{y}(\delta)\right| \leq 2 \rho(x, y), x, y \in D[0,1] . \tag{5.14}
\end{equation*}
$$

Weak convergence on $C[0,1]$ or $D[0,1]$ will be referred to as functional convergence. Jagers ((1974), Proposition 3.3) relates functional and weak convergence of a sequence $\left\{P_{n}\right\}$ of point processes converging to a point process $P$. In this case tightness is unnecessary (Straf (1972), Whitt (1975)), provided $D[0,1]$ has the $J_{1}$-topology if $P$ is orderly, or the $M_{1}$-topology (see Skorokhod (1956)) if $P$ is non-orderly. Hence p.g.fls. may be used to characterize weak convergence of point processes (Westcott (1972)).

## CHAPTER 2

## LIMIT LAWS FOR THE $\dot{G} / G / \infty$ QUEUE

### 2.1. Introduction

The $G / G / \infty$ queue may be regarded as a point process (the process of arrival times) subjected to independent random displacements. An important process associated with the $G / G / \infty$ queue is the number of servers $\phi(t)$ busy at any time $t$. Central limit theorems for the accumulated traffic time $\int_{0}^{t} \phi(s) d s$ have been obtained in the case of the $M / G / \infty$ queue with bulk arrivals by Rao (1966), under the assumption of finite third service time moment, and the $G I / G / \infty$ queue with bulk arrivals by Narasimham (1968), for finite second moment. Iglehart and Kennedy (1970) generalised the model to the case in which the service times may be mutually dependent, and, using functional weak convergence techniques, demonstrated a functional central
limit theorem for $\int_{0}^{t} \phi(s) d s$, which, specializing to i.i.d. service times, required finite $(2+\delta)$-th moment for the service time distribution. By confining ourselves to the $G / G / \infty$ queue in which the arrival time point process is stationary, we obtain in Section 2.5 a central limit theorem for $\int_{0}^{t} \phi(s) d s$ when the service times have finite second moment.

Brown and Ross (1969) studied almost sure convergence of the $M / G / \infty$ queue with bulk arrivals, and in particular examined the strong law of large numbers for this case. In Section 2.3 we look at stationarity, engodicity ans mixing of the $\phi(8)$ process, and in Section 2.4 derive conditions under whicl it obeys the strong law of large numbers. We defer investigation of the law of the iterated logarithm to Chapter 3 , where it will be discussed in
conjunction with cluster processes. Finally, in Sections 2.6 and 2.7 , we enquire into rates of convergence (both ordinary and functional) of the accumulated traffic time to normality.

In each of the above limit theorems our assumptions will be too general to allow the use of specialized techniques, predominant in weak and a.s. convergence applications, such as (with particular reference to the $G / G / \infty$ queue) renewal theory, the Skorkhod representation theorem (in Section 2.7), or the probability generating functional (Vere-Jones (1968), Westcott (1970), Section 3.6). Indeed, in accordance with the philosophy outlined in the Preface to this thesis, we will always ask, in a way to be defined more strictly later, that the limit law we desire of the service expended $\int_{0}^{t} \phi(s) d s$ be obeyed by the component processes - for other samples of this attitude, see Iglehart and Kennedy (1970) and Loulou (1973). In particular, we will retain the assumption of dependence between service times. This particular chapter was precipitated by a paper of Kaplan (1974), in which the various limit theorems (except rates of convergence) were obtained for the $G I / G / \infty$ queue by methods strongly depending on the independence assumptions of that system.

Many of the theorems of this chapter have been stated as assuming stationarity, a condition which seems necessary in Sections 2.3 and 2.4, but may (see Problem 6.3.4) be removed in later sections. In some of the theorems of this chapter, two almost identical results exist, and here we have indicated the differences of the second from the first in brackets in the statement of the theorem.

Infinite server queues have been widely applied. The pure birth-anddeath process is an $M / M / \infty$ queue (Feller (1968), p. 460). Amongst early investigations were those by C. Palm, A.K. Erlang and T.C. Fry in telephone trunking problems (see Feller (ibid.) for references). Benes (1957) obtained the distribution of traffic time average in the $M / M / \infty$ case.
T. Lewis (1961) and Nelsen and Williams ((1968), (1970)) consider regular events under arbitrary i.i.d. translations, i.e., the $D / G / \infty$ queue. Smith (1958) has used the generating function of $\phi(s)$ in the $G I / G / \infty$ queue to illustrate infinite products occurring in renewal theory. Thedeen (1969) discusses the queue as a model for traffic on a long road with free overtaking. For applications to textile research, see Rao (1966) and references there. The queue arises in engineering applications as the "jitter" process (Beutler and Leneman (1966)). Other works on the $G / G / \infty$ queue will be referred to in Section 6.2.

### 2.2. Notation and preliminaries

It will be convenient to regard the arrival process as a probability measure $P_{1}$ on $(N, \sigma(N))$, although we will also refer to it as a random non-negative integer-valued measure $\eta:(\Omega, F, P r) \rightarrow(N, \sigma(N))$, so that $P_{1}=P M^{-1}$. As always, we may define measurable mappings $t_{j}: N \rightarrow R$ giving the arrival times $t_{j}(N)$ for any realisation of the arrival process $N$ 。

The process of service times may be viewed as a probability measure $P_{2}$ on $\left(R_{+}^{Z}, B\left(R_{+}\right)^{Z}\right)$, where $Z$ is the set of integers; alternatively, (in Section 2.7 and Chapter 4 ), we will specify it as a random vector $\left\{V_{j}\right\}$ of service times from a fixed space $(\Omega, F, P r)$ to $\left(R_{+}^{Z}, B\left(R_{+}\right)^{Z}\right)$, so that $P_{2}=\operatorname{Pr}\left\{V_{j}\right\}^{-1}$.

Finally, we form the probability triple $\left(N \times R_{+}^{Z}, \sigma(N) \times B\left(R_{+}\right)^{Z}, P_{1} \times P_{2}\right)$; i.e., we regard the arrival process and service process as independent. We will denote vectors $\left\{x_{j}\right\} \in R_{+}^{Z}$ as $x$ if no ambiguity will arise. Let $\bar{Z}_{+}$denote $Z_{+} u\{\infty\}$. Then $\bar{Z}_{+}^{R}$ will mean the set of
non-negative integer or infinite valued functions on $R$. Let $\sigma\left(\bar{Z}_{+}\right)^{R}$ denote the product $\sigma$-field on $\bar{Z}_{+}^{R}$, where $\sigma\left(\bar{Z}_{+}\right)$consists of all subsets of $\bar{z}_{+}$.

We can then define the 'number of busy servers' process via a mapping $\phi=\phi_{1} \circ \phi_{2}$ from $N \times R_{+}^{2} \rightarrow \overline{2}_{+}^{R}$ by
(i) $\phi_{2}: N \times R_{+}^{Z} \rightarrow R^{2} \times R^{Z}$

$$
\phi_{2}(N, x)=(t(N), x)
$$

(ii)

$$
\begin{align*}
\phi_{I}: & R^{Z} \times R^{Z} \rightarrow \mathbb{Z}_{+}^{R} \\
& \left.\phi_{1}\left(\left\{\omega_{j}\right\},\left\{x_{j}\right\}\right)(t)=\sum_{j=-\infty}^{+\infty} 1_{\left[\omega_{j}, \omega_{j}+x_{j}\right.}\right)^{(t)} . \tag{2.1}
\end{align*}
$$

It is clear that for any joint realization ( $N, x$ ) of the arrival and server processes, $\phi$ will record 1 if an arrival occurs $u p$ to and including $t$ and is still being served after $t, 0$ otherwise and then will add these numbers, therefore corresponding to the number of busy servers at $t$, for each $t \in R$. To prove measurability of a mapping to a product space, we need only prove that each cocrdinate mapping is measurable. Hence it is immediately clear that

THEOREM 2.2.1. $\phi$ is measurable with respect to $\sigma(N) \times B\left(R_{+}\right)^{Z}$.
Throughout this chapter we assume that the number of busy servers is finite, i.e., that $\phi \in z_{+}^{R}, P_{1} \times P_{2}$ - a.s. We then define a probability measure $P_{n s}$ (corresponding to the number of busy servers process) on $\left(\bar{z}_{+}^{R}, \sigma\left(\bar{z}_{+}\right)^{R}\right)$ by

$$
P_{n s}=\left(P_{1} \times P_{2}\right) \phi^{-1}
$$

such that $P_{n s}\left(Z_{+}^{R}\right)=1$. From now on we will refer to $P_{n s}$ as if it were acting on $z_{+}^{R}$. We may similarly define a stochastic process
$m:(\Omega, F, P r) \rightarrow\left(\bar{Z}_{+}^{R}, \sigma\left(\bar{Z}_{+}\right)^{R}\right)$ such that $\operatorname{Prm}^{-1}=P_{n s}$. We will refer to $P_{n s}, \phi$, or $m$ as the "number of busy servers" process.

We will frequently employ the following decompositions.

$$
\begin{align*}
\int_{0}^{\tau} \phi(s) d s & =\sum_{i=1}^{N(0, \tau]} \min \left(x_{i}, \tau-t_{i}\right)+\sum_{i=-\infty}^{0}\left[\min \left(x_{i}+t_{i}, \tau\right)\right]^{+} \\
& \equiv \Phi_{1}(\tau)+\Phi_{2}(\tau), \tag{2.2}
\end{align*}
$$

say. Here we have split the integral into two parts, corresponding to those arriving up to time zero and those arriving later. Also, for some sequence $c(\tau) \leq \tau$, define

$$
\begin{align*}
\Phi^{\prime}(\tau) & \equiv\left|\Phi_{1}(\tau)-\sum_{i=1}^{N(0, \tau]} x_{i}\right|  \tag{2.3}\\
& =\sum_{i=1}^{N(0, \tau]}\left(x_{i}+t_{i}-\tau\right)^{+} \\
& \left.=N(0, \tau-c(\tau)] \sum_{i=1}\left(x_{i}+t_{i}-\tau\right)^{+}+\sum_{i=N(0, \tau-c(\tau)]+1}^{N(0, \tau]} x_{i} t_{i}-\tau\right)^{+} \\
& \leq \sum_{i=1}^{N(0, \tau]}\left(x_{i}-c^{\prime}(\tau)\right)^{+}+{ }_{i=N(0, \tau]}^{\left.\sum_{i-c}(\tau)\right]+1} x_{i} \\
& \equiv \Phi_{0}^{\prime}(\tau)+\Phi_{1}^{\prime}(\tau), \tag{2.4}
\end{align*}
$$

say, where we have used a device due to Iglehart and Kennedy (1970). We will decompose $\int_{0}^{\tau} m(s) d s$ similarly as

$$
\begin{gather*}
\int_{0}^{\tau} m(s) d s=M_{1}(\tau)+M_{2}(\tau)  \tag{2.5}\\
M^{\prime}(\tau) \leq M_{0}^{\prime}(\tau)+M_{1}^{\prime}(\tau) \tag{2.6}
\end{gather*}
$$

Note that (2.2) and (2.4) (similarly with (2.5) and (2.6)) give

$$
\begin{equation*}
\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \leq \Phi_{0}^{\prime}(\tau)+\Phi_{1}^{\prime}(\tau)+\Phi_{2}(\tau) \tag{2.7}
\end{equation*}
$$

Thus we have bounded the difference between the quantity we are
interested in, and another quantity in which the structure of the component processes is more evident.

### 2.3. Stationarity, ergodicity and mixing

To examine stationarity of $P_{n s}$, we must introduce certain translation operators $T_{t}(t \in R), S_{k}(k \in Z)$ and $S_{u}(u \in R)$ as follows.

$$
\begin{align*}
& T_{t}: N \rightarrow N \quad \text { where } \quad T_{t} N(A)=N(A+t), \text { all } A \in B(R),  \tag{3.1}\\
& S_{k}: R^{Z} \rightarrow R^{Z} \text { where } S_{k}\left\{x_{j}\right\}=\left\{x_{k+j}\right\},  \tag{3.2}\\
& S_{u}: Z_{+}^{R} \rightarrow Z_{+}^{R} \text { where } S_{u}\{x(t)\}=\{x(t+u)\} . \tag{3.3}
\end{align*}
$$

The processes $P_{1}, P_{2}$ and $P_{n 8}$ are called stationary if $P_{1}\left(T_{t}^{-1} B\right)=P_{1}(B), P_{2}\left(S_{k}^{-1} B\right)=P_{2}(B), P_{n s}\left(S_{u}^{-1} B\right)=P_{n s}(B)$ for $B$ within the appropriate $\sigma$-algebra, where $T_{t}^{-1} B, S_{k}^{-1} B$ and $S_{u}^{-1} B$ are the corresponding inverse image sets. Mixing and related concepts can be defined in a way analogous to (3.4) of Chapter 1 for any stationary process whose translation operator is indexed by $Z$ or $R$. Before proving our next theorem, we note that, as is well-known from ergodic theory, it is sufficient to prove ergodicity for a class of sets $C$ (i.e. for $C, D \in \mathbb{C}$ ) such that $\left\{\left\{\begin{array}{ll}\bigcup_{i=1}^{n} & C_{i}\end{array}\right\}, C_{i} \in C, n \in Z_{+}\right\}$is an algebra $A$ which generates the $\sigma$-field we require.

THEOREM 2.3.1. (a) If the armival process $P_{1}$ and the service time process $P_{2}$ are stationary, then 80 is the "number of servers" process $P_{n s}$.
(b) If $P_{1}$ is ergodic and $P_{2}$ is mixing, then $P_{n s}$ is ergodic.
(c) If $P_{1}, P_{2}$ are mixing, then so is $P_{n s}$.

Proof. We define $T_{y}^{\prime}: N \times R_{+}^{2} \rightarrow N \times R_{+}^{Z}, y \in R$, by

$$
\begin{array}{rlrl}
T_{y}^{\prime}(N, x) & =\left(T_{y}^{N}, S_{\left.N(0, y]^{x}\right),}\right. & y>0, \\
& =\left(T_{y} N, S_{\left.N[y, 0)^{x}\right),},\right. & y<0, \\
& =(N, x), & & y=0 . \tag{3.4}
\end{array}
$$

For any sets $C, D \in \sigma\left(z_{+}\right)^{R}$, we see that

$$
P_{n s}\left(C \cap S_{u} D\right)=P_{1} \times P_{2}\left\{\left(\phi^{-1} C\right) \cap T_{u}^{r^{-1}}\left(\phi^{-1} D\right)\right\}
$$

where $\phi^{-1} C \equiv\{(N, x): \phi(N, x) \in C\} \in \sigma(N) \times B\left(R_{+}\right)^{Z}$ by Theorem 2.2.1. Hence stationarity, ergodicity and mixing of $P_{n s}$ will follow from stationarity, ergodicity and mixing of $P_{1} \times P_{2}$ with respect to the transformation $T^{\prime}$. We prove only (b): proofs for (a) and (c) are similar.

$$
\text { Let } C^{\prime}=C_{1} \times C_{2}, \quad D^{\prime}=D_{1} \times D_{2} \in \sigma(N) \times B\left(R_{+}\right)^{Z}
$$

Then

$$
\begin{align*}
& \left|\tau^{-1} \int_{0}^{\tau} P_{1} \times P_{2}\left(C^{\prime} \cap T_{u}^{\prime} D^{\prime}\right) d u-P_{1} \times P_{2}\left(C^{\prime}\right) P_{1} \times P_{2}\left(D^{\prime}\right)\right| \\
& =\mid \tau^{-1} \int_{0}^{\tau} \sum_{n=0}^{\infty} P_{1}\left(C_{1} \cap T_{u_{1}} D_{1},\{N(0, u]=n\}\right)\left[P_{2}\left(C_{2} \cap S_{n} D_{2}\right)-P_{2}\left(C_{2}\right) P_{2}\left(D_{2}\right)\right] d u \\
& \\
& \quad+P_{2}\left(C_{2}\right) P_{2}\left(D_{2}\right) \tau^{-1} \int_{0}^{\tau}\left[P_{1}\left(C_{1} \cap T_{u} D_{1}\right)-P_{1}\left(C_{1}\right) P_{1}\left(D_{1}\right)\right] d u \mid \\
& \leq 2 \tau^{-1} \int_{0}^{\tau} P_{1}\{N(0, u] \leq m\} d u+\sup _{n>m}\left|P_{2}\left(C_{2} \cap S_{n} D_{2}\right)-P_{2}\left(C_{2}\right) P_{2}\left(D_{2}\right)\right|  \tag{3.5}\\
& \\
& \quad+\left|\tau^{-1} \int_{0}^{\tau} P_{1}\left(C_{1} \cap T_{u} D_{1}\right) d u-P_{1}\left(C_{1}\right) P_{1}\left(D_{1}\right)\right| .
\end{align*}
$$

But for any $\varepsilon>0$, there exists $m(\varepsilon)$ such that for $n>m(\varepsilon)$,

$$
\sup _{n>m(\varepsilon)}\left|P_{2}\left(C_{2} \cap S_{n} D_{2}\right)-P_{2}\left(C_{2}\right) P_{2}\left(D_{2}\right)\right| \leq \varepsilon,
$$

since $P_{2}$ is mixing. Thus $\underset{\tau \rightarrow \infty}{\lim \sup }(3.5) \leq \varepsilon$, i.e., $\lim _{\tau \rightarrow \infty}(3.5)=0$. By
the remark immediately preceding Theorem 2.3.1, this proves (b) and completes our proof of the theorem.

It is clear from the above proof that if we had defined $T_{y}^{\prime \prime}: N \times R_{+}^{Z} \rightarrow N \times R_{+}^{Z}$ by

$$
\begin{equation*}
T_{y}^{\prime \prime}(N, \mathbf{x})=\left(T_{y} N, S_{[y]} \mathrm{x}\right), \tag{3.6}
\end{equation*}
$$

ergodicity of $P_{1} \times P_{2}$ would have followed from ergodicity of $P_{1}$ and weak mixing (Equation (3.3), Chapter 1) of $P_{2}$. Spectral theory (see e.g. Mackey (1974), p. 213) suggests that this condition is necessary as well as sufficient for $T_{y}^{\prime \prime}$, but we have been unable to weaken Theorem 2.3.1 (b) to include this case for $T_{y}^{\prime}$ (but see Problem 6.3.1). Note that the argument above actually becomes slightly more difficult when modified to the independent service times situation, if it is not realised that then $P_{2}$ is mixing.

### 2.4. The strong law of large numbers

A theorem associated with ergodicity is the strong law of large numbers. Before we can prove it, however, we need some new equipment.

We will suppose as in the previous section that our arrival process $P_{1}$ is stationary, but now with finite intensity, i.e.,

$$
\begin{equation*}
m \equiv E_{1} N(0,1]<\infty ; \tag{4.1}
\end{equation*}
$$

if also $P_{1}$ is orderly (for $P_{1}$ non-orderly, see discussion after Corollary 2.4.2), then it is known (e.g. Slivnyak (1962), (1966)) that uniquely associated with $P_{1}$ is a probability measure $P_{1}^{0}$ (Palm measure) on $(N, \sigma(N))$ (in general, non-stationary) such that the arrival times $t$ satisfy $\ldots<t_{-1}<t_{0}=0<t_{1}<\ldots p_{1}^{0}-$ a.s., and the sequence $\left\{\pi_{j}\right\} \equiv\left\{t_{j+1}-t_{j}\right\}$ of a.s. positive random variables is stationary, i.e., if we
define $\theta_{k}: N \rightarrow N(k \in Z)$ by

$$
\begin{equation*}
\theta_{k} N(\cdot)=N\left(\cdot+t_{k}\right) \tag{4.2}
\end{equation*}
$$

then, if $C \in \sigma_{\text {gen }}\left\{\left\{\pi_{j} \in A\right\}, A \in B(R), j \in Z\right\} \equiv \sigma_{\pi} \subset \sigma(N)$,

$$
\begin{equation*}
P_{1}^{0}\left(\theta_{k}^{-1} C\right)=P_{1}^{0}(C) . \tag{4.3}
\end{equation*}
$$

We will require the expectation

$$
\begin{equation*}
U(x) \equiv E_{1}^{0} N(0, x]<\infty \tag{4.4}
\end{equation*}
$$

for some finite $x$. According to Kaplan (1955), this guarantees finiteness for all finite $x$, and also

$$
\begin{equation*}
\sup _{x>0}\{U(x+u)-U(x)\} \leq 2 U(u)+1, \tag{4.5}
\end{equation*}
$$

whence $U(x)=O(x),(x \rightarrow \infty)$. Under this assumption of orderliness, Daley ((1971), Lemma 9) shows that $U(x) \sim \lambda^{\prime} x$ for some constant $\lambda^{\prime} \geq m$, which is stronger. We will tacitly assume when discussing $P_{1}^{0}$ in the rest of this chapter that the $P_{1}$ to which it corresponds is orderly. We point out, though, that our arguments involving only $P_{1}$ (not $P_{1}^{0}$ or $U(x)$ ) will still go through without this assumption.

Finally we define various invariant $\sigma$-fields: let
$T \equiv\left\{C \in \sigma(N) \times B\left(R_{+}\right)^{Z}: T_{y}{ }^{-1} C=C\right\}, T_{1} \equiv\left\{C \in \sigma(N): T_{y}^{-1} C=C\right\}$,
$T_{2} \equiv\left\{C \in B\left(R_{+}\right)^{2}: S_{k}^{-1} C=C\right\}$ and $T_{1}^{0} \equiv\left\{C \in \sigma_{\pi}: \theta_{k}^{-1} C=C\right\}$, where we have used the definitions at (3.1), (3.2), (3.3) and (4.2).

We may now prove some ergodic theorems for the "busy server" process.
THEOREM 2.4.1. If the arrival process is stationary with $m<\infty$, and the service time process is stationary with
(a) $E_{2}\left\{x_{1}\right\}<\infty$, then

$$
\tau^{-1} \int_{0}^{\tau} \phi(s) d s \rightarrow E_{1} \times E_{2}\{\phi(1) \mid T\}, P_{1} \times P_{2}-a . s
$$

where, also,

$$
E_{1} \times E_{2}\{\phi(1) \mid T\}=E_{1}\left\{N(0,1] \mid T_{1}\right\} E_{2}\left\{x_{1} \mid T_{2}\right\}, P_{1} \times P_{2}-\text { a.s. }
$$

(b) $E_{2}\left\{\left(x_{1} \log x_{1}\right)^{+}\right\}<\infty$, and (4.4) and (4.5) are satisfied, then

$$
\tau^{-1} \int_{0}^{\tau} \phi(s) d s \rightarrow E_{2}\left\{x_{1} \mid T_{2}\right\} / E_{1}^{0}\left\{\Pi_{1} \mid T_{1}^{0}\right\}, \quad P_{1}^{0} \times P_{2}-a . s .,
$$

provided the might-hand side is $P_{1}^{0} \times P_{2}$ - a.s. defined. Proof. (a) Clearly the theorem will be proved if we can show

$$
\begin{equation*}
\tau^{-1}\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \rightarrow 0, P_{1} \times P_{2}-\text { a.s. } \tag{4.6}
\end{equation*}
$$

We use the inequality (2.7). Clearly if $c(\tau)=O(\tau)$, which we assume, $\tau^{-1} \Phi_{1}^{\prime}(\tau) \rightarrow 0$ a.s. Now

$$
\underset{\tau \rightarrow \infty}{\lim \sup } \tau^{-1} \Phi_{0}^{\prime}(\tau) \leq \underset{k \rightarrow \infty}{\lim \sup } \underset{\tau \rightarrow \infty}{\lim \sup } \tau^{-1} \sum_{i=1}^{N(0, \tau]}\left(x_{i}-k\right)^{+}
$$

But $(x-k)^{+}$is a Borel-measurable function of $x$, for any $k \in Z_{+}$, and $E_{2}\left\{\left(x_{1}-k\right)^{+}\right\} \leq E_{2}\left\{x_{1}\right\}<\infty$, and hence by the ergodic theorem

$$
\lim _{\tau \rightarrow \infty} \tau^{-1} \sum_{i=1}^{N(0, \tau]}\left(x_{i}-k\right)^{+}=E_{1}\left\{N(0,1] \mid T_{1}\right\} \cdot E_{2}\left\{\left(x_{1}-k\right)^{+} \mid T_{2}\right\}
$$

But $E_{2}\left\{\left(x_{1}-k\right)^{+} \mid T_{2}\right\} \rightarrow 0$ a.s. $(k+\infty)$, since $\left(x_{1}-k\right)^{+}+0$ a.s. as $k \rightarrow \infty$ (Breiman (1968), Proposition 4.24). To complete the proof of (a), we note only that since $P_{1} \times P_{2}$ is stationary with respect to $T_{y}^{\prime}$, by the ergodic theorem

$$
\tau^{-1} \int_{0}^{\tau} \phi(s) d s \rightarrow E_{1} \times E_{2}\{\phi(1) \mid T\}, P_{1} \times P_{2}-\text { a.s. }
$$

and

$$
\begin{align*}
& \qquad E_{1} \times E_{2}\left\{E_{1} \times E_{2}[\phi(1) \mid T]\right\}=E_{1}\{N(0,1]\} \cdot E_{2}\left\{x_{1}\right\} .  \tag{4.7}\\
& \text { Now, since }\left|\Phi_{1}(\tau)-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \rightarrow 0 \text { a.s., }
\end{align*}
$$

$$
\underset{\tau \rightarrow \infty}{\liminf \tau^{-1}} \int_{0}^{\tau} \phi(s) d s \geq E_{1}\left\{N(0,1] \mid T_{1}\right\} \cdot E_{2}\left\{x_{1} \mid T_{2}\right\}
$$

and, clearly,

$$
\begin{equation*}
E_{1} \times E_{2}\left\{E_{1}\left[N(0,1] \mid T_{1}\right] \cdot E_{2}\left[x_{1} \mid T_{2}\right]\right\}=E_{1}\{N(0,1]\} \cdot E_{2}\left\{x_{1}\right\} . \tag{4.8}
\end{equation*}
$$

(4.7) and (4.8) are consistent only if $\Phi_{2}(\tau) \rightarrow 0$ a.s.
(b) Again we prove $\tau^{-1}\left|\Phi_{1}(\tau)-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \rightarrow 0$ a.s. in exactly the same way. Now we show $\Phi_{2}(\tau) \rightarrow 0$ a.s. Let $F$ denote the distribution function of the service times, i.e., $F(\eta)=P_{2}\left\{x_{i} \leq l\right\}$. Let $r_{k} \equiv\left[c^{k-1}\right]$ for some fixed $a>1$. For a given $\varepsilon>0$, suppose $\Phi_{2}(\tau) \geq \varepsilon \tau$ for some $\tau$. Let $r_{k}$ be the first term of the geometric subsequence larger than $\tau$. Then $\Phi_{2}\left(r_{k}\right) \geq \varepsilon r_{k-1}$. Hence we need only prove $\Phi_{2}(\tau) \rightarrow 0$ a.s. for a geometric subsequence. But

$$
\begin{aligned}
& P_{1}^{0} \times P_{2}\left\{r_{k}^{-1} \Phi_{2}\left(r_{k}\right) \geq \varepsilon\right\} \\
& \leq E_{1}^{0} \times E_{2}\left\{\Phi\left(r_{k}\right)\right\} / \varepsilon r_{k} \\
&=\int_{0}^{\infty} d U(u) \int_{0}^{\infty} \max \left[\min \left(\tau-u, r_{k}\right), 0\right] d F(\eta) / \varepsilon r_{k} \\
&=\left(\varepsilon r_{k}\right)^{-1} \int_{r_{k}}^{\infty} d F(\tau) \int_{\tau-r_{k}}^{\tau}(\tau-u) d U(u)+\varepsilon^{-1} \int_{r_{k}}^{\infty} d F(\tau) \int_{0}^{\tau-r_{k}} d U(u) \\
& \quad+\left(\varepsilon r_{k}\right)^{-1} \int_{0}^{r_{k}} d F(\tau) \int_{0}^{\tau}(z-u) d U(u),
\end{aligned}
$$

so that
$\varepsilon \sum_{k=1}^{\infty} P_{1}^{0} \times P_{2}\left\{r_{k}^{-1} \Phi_{2}\left(r_{k}\right) \geq \varepsilon\right\}$

$$
\begin{aligned}
& \leq 2 B(\log c)^{-1} \int_{1}^{\infty} Z \log Z d F(Z)+(\log c)^{-1} \int_{1}^{\infty} \log \eta d F(Z) \\
&+\left[2 \sum_{k=1}^{\infty} r_{k}^{-1}+B\left(2 c^{2}+2\right)\right] \int_{0}^{\infty} Z d F(Z)+\left(2 c^{2}+1\right)<\infty
\end{aligned}
$$

where we have used $U(Z) \leq B Z+1$ for some constant $B$, bounded the sums by integrals and exchanged the order of integration. Hence $\Phi_{2}(\tau) \rightarrow 0$ a.s. by the Borel-Cantelli lemma. The proof is now completed by observing that

$$
t_{i} / i \xrightarrow{\text { a.s. }} E_{I}^{0}\left\{\Pi_{I} \mid T_{I}^{0}\right\}
$$

and employing $\{N(0, y]<j\}=\left\{t_{j}>y\right\}$ and the monotonicity of $N(0, \tau] . \square$
We would anticipate that Theorem $2.4 .1(b)$ is sub-optimal: $E_{2}\left\{x_{1}\right\}<\infty$ should suffice (see Problem 6.3.2). We deduce immediately from Theorem 2.4.1 the strong law of large numbers.

COROLLARY 2.4.2. If the arrival process is stationary with $m<\infty$ (and (4.4) and (4.5) hold), and satisfies $N(0, n] / n \rightarrow m, P_{1}\left(P_{1}^{0}\right)-a . s .$, and if the service time process is stationary with $\sum_{i=1}^{n} x_{i} / n \rightarrow E_{2}\left\{x_{1}\right\}<\infty$, $P_{2}-a . s .,\left(\right.$ and $\left.E_{2}\left\{\left(x_{1} \log x_{1}\right)^{+}\right\}<\infty\right)$, then

$$
\tau^{-1} \int_{0}^{\tau} \phi(s) d s \rightarrow m E_{2}\left\{x_{1}\right\}, \quad P_{1} \times P_{2}-a . s . \quad\left(P_{1}^{0} \times P_{2}-a . s .\right)
$$

We should compare Corollary 2.4.2 with Theorem 2.3 .1 ( $b$ ) and the comments following Theorem 2.3.1: it would appear (although we have no example) that the strong law of large numbers could hold for $\phi(8)$ without the process being ergodic - simple examples of this behaviour in other contexts are easy to construct (Breiman (1968), pp. 110 and 113, or Hannan (1973), p. 163).

Functional strong laws can be proved for this process in the manner of Iglehart (1971b).

Because queues with bulk arrivals are given much attention in the literature, we wish to make a few points in this context about the Palm measure corresponding to a non-orderly stationary point process. Slivnyak in fact goes further than indicated in the opening paragraphs of this section.

Consider the set $R_{\uparrow}^{2}$ of subsets $t^{\prime}$ of $R$ without finite limit points satisfying

$$
\cdots \leq t_{-1}^{\prime} \leq t_{0}^{\prime}=0 \leq t_{1}^{\prime} \leq \cdots \cdot
$$

Note that if, for a given $t^{\prime}$, there are $k$ points at the origin, then there will be $(k-1)$ other members of $R_{\uparrow}^{Z}$, which will constitute only a relabelling of $t^{\prime}$. Let $\sigma(\uparrow)$ denote the minimal $\sigma$-algebra generated by sets of the form $\left\{t^{\prime}: t_{k}^{\prime} \leq y\right\}, y \in R, k \in Z$.

Slivnyak then proves that uniquely coupled with a stationary point process $P$ of finite intensity $m$ is a probability measure $P^{0}$ on $\left(R_{\uparrow}^{Z}, \sigma(\uparrow)\right)$ such that $\left\{\Pi_{j}^{\prime}\right\} \equiv\left\{t_{j+1}^{\prime}-t_{j}^{\prime}\right\}$ is stationary, i.e., if we define $S_{k}$ as in (3.2),

$$
\operatorname{Pr}^{0}\left(S_{k}^{-1} D\right)=\operatorname{Pr}^{0}(D), \quad k \in Z
$$

where $D \in \sigma_{\text {gen. }}\left\{\left\{\Pi_{j}^{\prime} \in B\right\}, B \in B(R), j \in Z\right\} \subset \sigma(\uparrow)$. Note that a mapping $\eta: R_{\uparrow}^{2} \rightarrow N$ may be defined analogously to (2.4) of Chapter 1 , and will be measurable, e.g. for $u, v>0, k \in Z_{+}$,

$$
\left\{t^{\prime}: \eta(u, u+v] \geq k\right\}=\bigcup_{i=1}^{\infty}\left\{t^{\prime}: t_{i}^{\prime} \leq v, t_{i+1}^{\prime}>u, t_{i+k}^{\prime} \leq u+v\right\}
$$

Therefore $P r^{0}$ induces a probability measure $P^{0}=\operatorname{Pr}^{0} \eta^{-1}$ on $(N, \sigma(N))$ such that if $N_{0}=\{N \in N: N(\{0\})>0\}, P^{0}\left(N_{0}\right)=1$. However, as is readily seen, several t's may map to one $N$, so that an inverse mapping is undefined, i.e., it is no longer meaningful to refer to
$t_{j}(N)=t_{j}^{\prime}\left(t^{\prime}\right)$ in an obvious notation, unless $P$ is a.s. orderly. However, had we adopted Slivnyak's approach, our theorems would hold, with some modifications, for non-orderly processes too.

Let $\operatorname{Pr}_{1}^{0}$ on $R_{\uparrow}^{Z}$ correspond to a given stationary arrival process. Despite the aforementioned difficulties, it is natural to identify the $t^{\prime}$. with arrival epochs. Hence define $\phi^{\prime}: R_{\uparrow}^{Z} \times R_{+}^{Z} \rightarrow \bar{Z}_{+}^{R}$ in a way similar to (2.1). The quantity naturally thrown up from $\int_{0}^{\tau} \phi^{\prime}(s) d s$ corresponding to $\sum_{j=1}^{N(0, \tau]} x_{j}$ is $\sum_{t_{j}^{\prime} \in(0, \tau]} x_{j}$, an expression not in a suitable form for application of random change of time lemmas. However,

$$
\sum_{j=n(\{0\})}^{n(0, \tau]} x_{j} \leq \sum_{t_{j}^{\prime} \in(0, \tau]} x_{j} \leq \sum_{j=1}^{n[0, \tau]} x_{j},
$$

and by showing (use geometric subsequences, (4.9), and Chebyshev's inequality) that for any $\delta>0$,

$$
\tau^{-\delta}\left|\sum_{j=1}^{n[0, \tau]} x_{j}-\sum_{j=\eta(\{0\})}^{\eta(0, \tau]} x_{j}\right| \rightarrow 0, P_{1}^{0} \times P_{2}-a . s .,(\tau \rightarrow \infty),
$$

we could employ the same techniques as for the orderly case to establish the same results. However, we do not consider the notational inconvenience is worth the extra generality.

We would also require generalisations of (4.4) and (4.5). Recall that a finite intensity and stationarity imply that the finite rate $\lambda$ exists (Chapter 1, (3.20)). Again, Kaplan shows that if

$$
U(x) \equiv E^{0} n(0, x] \equiv E_{1}^{0} N(0, x]<\infty
$$

for a finite $x$, then $U(x)$ is finite for all $x$, and

$$
\begin{equation*}
\sup _{x>0}\{U(x+v)-U(x)\} \leq 2 U(v)+m / \lambda, \tag{4.9}
\end{equation*}
$$

whence $U(x)=O(x) \quad(x \rightarrow \infty)$. Our theorems now go through as before.

### 2.5. The central limit theorem

Before we prove the central limit theorem, we need the following simple lemma (which is possibly known, but we could locate no statement or proof). The lemma may also be used to weaken Theorem 3 of Daley (1972).

LEMMA 2.5.1. Let $\left\{Y_{n}\right\},\left\{Z_{n}\right\}$ and $\left\{N_{n}\right\}$ be sequences of random variables defined on the some fixed probability space such that $Y_{n} \xrightarrow{D} Y$ and $Z_{n} \xrightarrow{D} Z$ for some random variables $Y$ and $Z$ defined on (not necessarily) another probability space. If $\left\{N_{n}\right\}$ consists of non-negative integer-valued random variables such that $N_{n} \xrightarrow{p} \infty$ and $\left\{\left(N_{n}, Z_{n}\right)\right\}$ is independent of $\left\{Y_{n}\right\}$, then

$$
\left(Y_{N_{n}}, Z_{n}\right) \xrightarrow{D}(Y, Z)
$$

Proof. Let ( $y, z$ ) be a continuity point of ( $Y, 2$ ) . Then clearly $y$ is a continuity point of $Y$, and $z$ a continuity point of $Z$. Hence, for any arbitrary $\varepsilon>0$ and for $k \geq k_{0}(y, \varepsilon) \equiv k_{0}$,

$$
\left|\operatorname{Pr}\left\{Y_{k} \leq y\right\}-\operatorname{Pr}\{Y \leq y\}\right| \leq \varepsilon .
$$

(For the purposes of this proof we assume all random variables on a single space ( $\Omega, F, \operatorname{Pr}$ ) ). Thus

$$
\begin{aligned}
& \qquad \begin{aligned}
& \operatorname{Pr}\left\{Y_{N_{n}} \leq y, Z_{n} \leq z\right\} \leq \operatorname{Pr}\left\{N_{n}<k_{0}\right\}+\sum_{k=k_{0}}^{\infty} \operatorname{Pr}\left\{Y_{k} \leq y\right\} \operatorname{Pr}\left\{Z_{n} \leq z, N_{n}=k\right\} \\
& \leq \operatorname{Pr}\left\{N_{n}<k_{0}\right\}+\operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{Z_{n} \leq z\right\}+\varepsilon, \\
& \text { i.e., } \limsup _{n \rightarrow \infty} \operatorname{Pr}\left\{Y_{N_{n}} \leq y, Z_{n} \leq z\right\} \leq \operatorname{Pr}\{Y \leq y, Z \leq z\} .
\end{aligned} \\
& \quad \text { Also, }
\end{aligned}
$$

$$
\begin{aligned}
&{\operatorname{Pr}\left\{Y_{N_{n}} \leq y, z_{n} \leq z\right\}} \\
& \geq \sum_{k=k_{0}}^{\infty} \operatorname{Pr}\left\{Y_{k} \leq y\right\} \operatorname{Pr}\left\{z_{n} \leq z, N_{n}=k\right\} \\
& \geq \operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{z_{n} \leq z\right\}+\operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{N_{n} \geq k_{0}\right\} \\
&-\operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{\left(N_{n} \geq k_{0}\right) \cup\left(z_{n} \leq z\right)\right\}-\varepsilon \\
& \geq \operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{z_{n} \leq z\right\}+\operatorname{Pr}\{Y \leq y\} \operatorname{Pr}\left\{N_{n} \geq k_{0}\right\}-\operatorname{Pr}\{Y \leq y\}-\varepsilon .
\end{aligned}
$$

Hence

$$
\begin{array}{r}
\quad \lim \inf \operatorname{Pr}\left\{Y_{N_{n}} \leq y, Z_{n} \leq z\right\} \geq \operatorname{Pr}\{Y \leq y, z \leq z\}, \\
\text { i.e., } \lim _{n \rightarrow \infty} \operatorname{Pr}\left\{Y_{N_{n}} \leq y, z_{n} \leq z\right\}=\operatorname{Pr}\{Y \leq y, z \leq z\}
\end{array}
$$

We also require the following lemma, which, conforming to a general theme of this chapter (and this thesis), proves that the difference between the quantity we require (the traffic time average) and a quantity which is easy to handle $\left(\sum_{i=1}^{N(0, \tau]} x_{i}\right)$ is small under suitable circumstances.

LEMMA 2.5.2. If the arrival process is stationary with $m<\infty$ (and (4.4) and (4.5) hold), and the sequence of service times is stationary with finite $(2-\delta+\cup)$-th moment, any $u>0,0<\delta \leq 2$, or under any circumstances if $\delta>2$, then

$$
\tau^{-\delta}\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \xrightarrow{P_{1} \times P_{2}\left(P_{1}^{0} \times P_{2}\right)} 0, \quad(\tau \rightarrow \infty) .
$$

Proof. We confine our proof to the $P_{1} \times P_{2}$ case. Arguments for the $P_{1}^{0} \times P_{2}$ case use techniques similar to those in Theorem 2.4.1 (b), and account has to be taken of arrivals at 0 , but no new ideas are involved. Now by (2.2) and (2.3),

$$
\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N(0, \tau]} x_{i}\right|=\Phi_{2}(\tau)+\Phi^{\prime}(\tau) .
$$

For any $\varepsilon>0$,

$$
\begin{align*}
P_{1} \times P_{2}\left\{\tau^{-\delta} \Phi_{2}(\tau) \geq \varepsilon\right\} & \leq E_{1} \times E_{2}\left\{\Phi_{2}(\tau)\right\} / \varepsilon \tau^{\delta} \\
& =\varepsilon^{-1} \tau^{-\delta} \int_{0}^{\infty} d F(\tau) \int_{0}^{\tau} \min (\tau-u, \tau) d u \\
& \leq \varepsilon^{-1} \tau^{-u} \int_{0}^{\tau} \tau^{2-\delta+U} d F(\tau)+\varepsilon^{-1} \int_{\tau}^{\infty} \tau^{2-\delta} d F(\tau) \tag{5.1}
\end{align*}
$$

where $F(Z) \equiv P_{2}\left\{x: x_{i} \leq \mathcal{Z}\right\}$ is the service time distribution (independent of $i$ since $P_{2}$ is stationary). Hence $(5.1) \rightarrow 0,(\tau \rightarrow \infty)$. Similarly, Chebyshev's inequality yields exactly the same upper bound for $P_{1} \times P_{2}\left\{\tau^{-\delta_{\Phi^{\prime}}(\tau)} \geq \varepsilon\right\}$.

This leads to a simple proof of the central limit theorem. We denote convergence in distribution by $D_{1}, D_{1}^{0}, D_{2}, D_{1} \times D_{2}, D_{1}^{0} \times D_{2}$ in an obvious notation.

THEOREM 2.5.3. If the arrival process is stationary with $m<\infty$ (and (4.4) and (4.5) hold), and satisfies

$$
\begin{equation*}
\tau^{-\frac{3}{2}}\{N(0, \tau]-m \tau\} \xrightarrow{D_{1}\left(D_{1}^{0}\right)} X_{1}, \tag{5.2}
\end{equation*}
$$

where $X_{1}$ is some random vamiable, and the service times form a stationary sequence with finite $(3 / 2+\delta)$ th moment, and satisfy

$$
\begin{equation*}
n^{-\frac{z_{2}}{2}}\left\{\sum_{j=1}^{n} x_{j}-n E_{2}\left\{x_{1}\right\}\right\} \xrightarrow{D_{2}} X_{2}, \tag{5.3}
\end{equation*}
$$

where $X_{2}$ is some other random variable defined on the same probability space as $X_{1}$, then

$$
\begin{equation*}
\tau^{-\frac{3}{2}}\left\{\int_{0}^{\tau} \phi(s) d s-m E_{2}\left\{x_{1}\right\} \tau\right\} \xrightarrow{D_{1} \times D_{2}\left(D_{1}^{0} \times D_{2}\right)} E_{2}\left\{x_{1}\right\} X_{1}+m^{\frac{1}{2}} X_{2} \tag{5.4}
\end{equation*}
$$

Proof. We will only prove the $D_{1} \times D_{2}$ convergence: $D_{1}^{0} \times D_{2}$ only requires writing $D_{1}^{0}$ for $D_{1}$ and $P_{1}^{0}$ for $P_{1}$. All of the random variable
in the following indexed by $\tau$ are on the space $N \times R_{+}^{Z}$, so Lemma 2.5.1 may be applied. For this proof, set $\mu=E_{2}\left\{x_{1}\right\}$, and $N_{\tau}=N(0, \tau]$. According to Lemma 2.5.2,
$\tau^{-\frac{\pi}{2}}\left|\left\{\int_{0}^{\tau} \phi(s) d s-m \mu \tau\right\}-\left\{\sum_{i=1}^{N} x_{i}-m \mu \tau\right\}\right|=\tau^{-\frac{\pi}{2}}\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N} x_{i}\right| \xrightarrow{P_{1} \times P_{2}} 0$.
But

$$
\begin{equation*}
\tau^{-\frac{3}{2}}\left\{\sum_{i=1}^{N_{\tau}} x_{i}-m \mu \tau\right\}=\tau^{-\frac{3}{2}} \sum_{i=1}^{N_{\tau}}\left(x_{i}-\mu\right)+\tau^{-\frac{1}{2}} \mu\left\{N_{\tau}-m \tau\right\} . \tag{5.5}
\end{equation*}
$$

By (5.2), $\left(N_{\tau}+1\right) / m \tau \xrightarrow{P_{1}} 1$, and hence by Theorem 4.4 of Billingsley (1968), and the continuous mapping theorem (Section 1.5)

$$
\left(N_{\tau}-m \tau\right) /\left(N_{\tau}+1\right)^{\frac{3}{2}} \xrightarrow{D_{1}} x_{1} / \sqrt{m}
$$

In Lemma 2.5.1 (or a continuous version of it), put

$$
z_{\tau} \equiv \mu\left(N_{\tau}-m \tau\right) /\left(N_{\tau}+1\right)^{\frac{3}{2}}, y_{\tau} \equiv \sum_{i=1}^{[\tau]}\left(x_{i}-\mu\right) /([\tau]+1)^{\frac{3}{2}},
$$

where $[\tau]$ is the integer part of $\tau$, and $N_{\tau} \equiv N(0, \tau]$; then, applying simultaneously the continuous mapping theorem,

$$
\left(N_{\tau}+1\right)^{-\frac{1}{2}}\left\{\sum_{i=1}^{N}\left(x_{i}-\mu\right)+\mu\left(N_{\tau}-m \tau\right)\right\} \xrightarrow{D_{1} \times D_{2}} X_{2}+\mu X_{1} / \sqrt{m}
$$

Again using Theorem 4.4 of Billingsley (1968) and the continuous mapping theorem,

$$
(5.5) \xrightarrow{D_{1} \times D_{2}} \sqrt{m X_{2}}+\mu X_{1} .
$$

In Theorem 2.5.3, we have avoided the question of norming (see Problem 6.3.3)). Hence, writing $\mu=E_{2}\left\{x_{1}\right\}$ again,

LEMMA 2.5.4. If the armival process is stationary, and the process of service times is stationary with $E_{2}\left\{x_{i}^{2}\right\}<\infty$, and

$$
\begin{equation*}
\sum_{i=1}^{\infty}\left|\rho_{i}\right|<\infty \tag{5.6}
\end{equation*}
$$

where

$$
\rho_{i}=E_{2}\left\{\left(x_{0}-\mu\right)\left(x_{i}-\mu\right)\right\}, \text { and }
$$

(a) $\operatorname{Var}_{P_{1}} N(0, u] \sim \lambda_{1} u$, then

$$
\begin{aligned}
& \operatorname{Var}_{P_{1} \times P_{2}}\left\{\sum_{i=1}^{N(0, \tau]} x_{i}\right\} \sim\left(m \operatorname{Var}_{P_{2}}\left(x_{1}\right)+m \rho+\lambda_{1} \mu^{2}\right) \tau, \quad(\tau \rightarrow \infty), \\
& \text { where } \rho \equiv 2 \sum_{i=1}^{\infty} \rho_{i} .
\end{aligned}
$$

(b) $\operatorname{Var}_{P_{1}^{0}} N(0, u] \sim \lambda_{2} u$, and $U(\tau) / \tau \rightarrow \lambda^{\prime}<\infty$, then

$$
\operatorname{Var}_{P_{1}^{0} \times P_{2}}\left\{\sum_{i=1}^{N(0, \tau]} x_{i}\right\} \sim\left(\lambda^{\prime} \operatorname{Var}_{P_{2}}\left(x_{1}\right)+\lambda^{\prime} \rho+\lambda_{2} \mu^{2}\right) \tau, \quad(\tau \rightarrow \infty)
$$

Remark. If the arrival process is orderly, then it is known (Daley (1971)) that $U(\tau) / \tau \rightarrow \lambda^{\prime}<\infty$.

Proof. Again we prove only (a). Write $N_{\tau} \equiv N(0, \tau]$. Easily

$$
\operatorname{Var}\left\{\sum_{i=1}^{N} x_{i}\right\}=m \operatorname{Var}\left\{x_{1}\right\} \tau+\mu^{2} \operatorname{Var} N_{\tau}+2 E_{1}\left\{\sum_{j=1}^{N}\left(N_{\tau}-j\right) \rho_{j}\right\}
$$

Clearly, by (5.6),

$$
\left|E_{1} \sum_{j=1}^{N_{\tau}}\left(N_{\tau}-j\right) \rho_{j}\right| \leq m \sum_{j=1}^{\infty}\left|\rho_{j}\right|<\infty
$$

Now,

$$
\begin{align*}
& \left|E_{1}\left\{\sum_{j=1}^{N} \tau^{-1}\left(N_{\cdot}-j\right) \rho_{j}\right\}-m \sum_{j=1}^{\infty} \rho_{j}\right| \\
& \quad \leq\left|E_{1}\left\{\left(\tau^{-1} N_{\tau}-m\right) \sum_{j=1}^{N_{1}} \rho_{j}\right\}\right|+\tau^{-1} E_{1}\left\{\sum_{j=1}^{N} j\left|\rho_{j}\right|\right\}+m E_{1}\left\{\sum_{j=N_{\tau}+1}^{\infty}\left|\rho_{j}\right|\right\} \tag{5.7}
\end{align*}
$$

By dominated convergence and the fact that $N_{\tau} \rightarrow \infty$ ass. as $\tau \rightarrow \infty$, the final term of $(5.7) \rightarrow 0(\tau \rightarrow \infty)$. The second term has an upper bound of

$$
\begin{align*}
\tau^{-1} E_{1}\left\{\sum_{i=1}^{N}\right. & \left.\sum_{j=i}^{\infty}\left|\rho_{j}\right|\right\} \\
& =\tau^{-1} \int_{\left\{N_{\tau}>k\right\}} \sum_{i=1}^{\tau} \sum_{j=i}^{\infty}\left|\rho_{j}\right| d P_{1}+\tau^{-1} \int_{\left\{N_{\tau} \leq k\right\}} \sum_{i=1} \sum_{j=i}^{N}\left|\rho_{j}\right| d P_{1} \tag{5.8}
\end{align*}
$$

For an arbitrary $\varepsilon>0$, let $k \equiv k(\varepsilon)$ be such that $\sum_{j=k}^{\infty}\left|\rho_{j}\right|<\varepsilon$. Then clearly the second term of $(5.8)$ has an upper bound of

$$
k \sum_{j=1}^{\infty}\left|\rho_{j}\right| P_{1}\left\{N_{\tau} \leq k\right\} / \tau \quad(\rightarrow 0,(\tau \rightarrow \infty)),
$$

and the first term is bounded by

$$
\begin{aligned}
\sum_{i=1}^{k-1} \sum_{j=i}^{\infty}\left|\rho_{j}\right| P_{1}\left\{N_{\tau}>k\right\} / \tau & \\
& +\varepsilon \tau^{-1} \int_{\left\{N_{\tau}>k\right\}}{ }^{N_{\tau} d P_{1}(\leq 0,(\tau \rightarrow \infty))} \text { ( } \quad(m \varepsilon,(\tau \rightarrow \infty))
\end{aligned}
$$

A similar style of argument applied to the first term of (5.8) yields an asymptotic upper bound of $2 m \varepsilon$ for this term: here we need the fact that

$$
\int_{\left\{N_{\tau}>k\right\}}\left(N_{\tau} / \tau-\mu\right) d P_{1}=-\int_{\left\{N_{\tau} \leq k\right\}}\left(N_{\tau} / \tau-\mu\right) d P_{1}
$$

Hence we conclude that $(5.7) \rightarrow 0(\tau \rightarrow \infty)$.
This central limit theorem has been proved here under second moment conditions: however, the functional case (Iglehart and Kennedy (1970)) seems to require that the service times have finite ( $2+\delta$ )-th moment. The generalisation to non-stationary arrival processes seems difficult (see Problem 6.3.4).

### 2.6. Rates of convergence

Let $X$ be a random variable on a space $\left(\Omega_{1}, F_{1}, P r_{1}\right)$, and $Y$ a random variable on another (not necessarily distinct) space $\left(\Omega_{2}, F_{2}, P r_{2}\right)$.

Then we will define the supremum metric for these two random variables by

$$
\nu(X, Y) \equiv \sup _{-\infty<x<\infty}\left|\operatorname{Pr}_{1}\{X \leq x\}-\operatorname{Pr}_{2}\{Y \leq x\}\right|
$$

Then we say that $X_{n} \xrightarrow{D} X$ at rate (at least) $c_{n} \downarrow 0 \quad(n \rightarrow \infty)$ if $v\left(X_{n}, X\right) \leq c_{n}$. Note that $X,\left\{X_{n}\right\}$ do not have to be on the same space, and hence, in searching for rates of convergence results for the traffic time average for the $G / G / \infty$ queue, we can (and will) retain the same notation as in the previous sections. We will require the following result, which is proved as Lemma 1 of Tomko (1972).

LEMMA 2.6.1. Let $X_{t}=Y_{t}+\delta_{t}, t \in R$, be the sum of two random variables, and $N$ be a random variable with the normal distribution. If $\delta_{t} \rightarrow 0$ in probability as $t \rightarrow \infty$ and $\operatorname{Pr}\left\{\left|\delta_{t}\right| \geq \varepsilon_{t}\right\} \leq \beta_{t}$, then

$$
v\left(X_{t}, N\right) \leq \nu\left(Y_{t}, N\right)+\varepsilon_{t}+\beta_{t} .
$$

It is clear that a result similar to Lemma 2.5 .1 may be obtained: however, our metric $v$, unlike some closely related metrics (discussed later) is not continuous under suitably smooth mappings (Whitt (1974b)) and hence we state more specifically:

LEMMA 2.6.2. Let $\left\{y_{\tau}\right\},\left\{z_{\tau}\right\},\left\{N_{\tau}\right\}$ be sequences of random variables on the same space $\left(\Omega, F, P_{r}\right)^{\prime}$ such that $\nu\left(Y_{\tau}, Y\right) \leq \psi(\tau), \nu\left(Z_{\tau}, Z\right) \leq \varphi(\tau)$ Where $Y, Z$ are continuous and are on a common space (not necessarily $\Omega$ ). If $N_{\tau} \xrightarrow{p}+\infty$ and $\left\{\left(N_{\tau}, z_{\tau}\right)\right\}$ is independent of $\left\{Y_{\tau}\right\}$, then

$$
\nu\left(Y_{N_{\tau}}+Z_{\tau}, Y+Z\right) \leq \varphi(\tau)+\int_{0}^{\infty} \psi(x) d P_{r}\left\{N_{\tau} \leq x\right\} .
$$

Proof. Omitted. Similar to the proof of Lemma 2.6.12.
In our applications of this lemma, we will suppose always that $Y$ and $Z$ have a standard normal distribution. As such it can be supplemented by (denoting standard normal random variables by $N$ )

LEMMA 2.6.3. Suppose $\left(N_{\tau}-\kappa \tau\right) / \sigma \sqrt{\tau} \xrightarrow{D} N_{1}$ at rate $\beta(\tau)$, where $N_{\tau}$
is a positive random variable for each $\tau>0$, and $\sigma, \kappa$ are some constants, and, not necessarily independently, $X_{\tau} / \sqrt{\tau} \xrightarrow{D} N_{2}$ at rate $\alpha(\tau)$; then

$$
\operatorname{Pr}\left\{\left|X_{\tau} /\left(K^{-1} N_{\tau}\right)^{\frac{3}{2}}-X_{\tau} / \tau^{\frac{3}{2}}\right| \geq 2 \kappa(\log \tau) / \sigma \sqrt{\tau}\right\} \leq A_{1} \max (\alpha(\tau), \beta(\tau), 1 / \sqrt{\tau})
$$

for some constant $A_{1}$.
Proof. Let $\Phi(x) \equiv \int_{-\infty}^{x} e^{-y^{2} / 2} d y / \sqrt{2 \pi}$; then, for some $y_{\tau}>0$,
$\operatorname{Pr}\left\{\left|X_{\tau} /\left(K^{-1} N_{\tau}\right)^{\frac{3}{2}}-X_{\tau} / \tau^{\frac{3}{2}}\right| \geq \varepsilon_{\tau}\right\}$

$$
\begin{aligned}
& \leq \operatorname{Pr}\left\{\left|X_{\tau} / \tau^{\frac{3}{2}}\right|>y_{\tau}\right\}+\operatorname{Pr}\left\{\left|\left(\tau / \kappa^{-1} N_{\tau}\right)^{\frac{3}{2}}-1\right| \geq \varepsilon_{\tau} / y_{\tau}\right\} \\
& \leq 2 \alpha(\tau)+2\left(1-\Phi\left(y_{\tau}\right)\right)+\operatorname{Pr}\left\{\left|N_{\tau}-\kappa \tau\right| / \sigma \tau \geq \kappa \varepsilon_{\tau} / 2 \sigma y_{\tau}\right\} \\
& \leq 2\{\alpha(\tau)+\beta(\tau)\}+2\left\{1-\Phi\left(y_{\tau}\right)\right\}+2\left\{1-\Phi\left(\kappa \varepsilon_{\tau} \sqrt{\tau} / 2 \sigma y_{\tau}\right)\right\} \cdot
\end{aligned}
$$

To optimise the rate of convergence, we choose $y_{\tau}=\kappa \varepsilon_{\tau} \sqrt{\tau} / 2 \sigma y_{\tau}$.
According to Feller (1968), p. 175, $1-\Phi(x)<e^{-\frac{1}{2} x^{2}}\left[(2 \pi)^{\frac{3}{2}} x\right]^{-1}, x>0$. Hence, since we will require $\varepsilon_{\tau}$ and $1-\Phi\left(y_{\tau}\right)$ to converge to zero at about the same rate, we choose, $y_{\tau}=(\log \tau)^{\frac{7}{2}}$, and obtain

$$
1-\Phi\left(y_{\tau}\right)<(2 \pi \tau)^{-\frac{1}{2}}(\log \tau)^{-1}
$$

and $\varepsilon_{\tau}=2 K(\log \tau) / \sigma \sqrt{\tau}$. Note that we really only require $y_{\tau}$ such that $e^{-\frac{3}{2} y_{\tau}^{2}} / y_{\tau}=1 / \sqrt{\tau}$, and finer solutions are possible, e.g. $y_{\tau}=[\log (\tau / \log \tau)]^{\frac{3}{2}}$.

We will require a statement of the rate at which $\tau^{-\frac{3}{2}} \int_{0}^{\tau} \phi(s) d s$ and $\sum_{i=1}^{N(0, \tau]} x_{i}$ approach each other. Since we will make several such estimates
(Lemmas 2.6.7, 2.6.8 and 2.6.10), we propose
DEFINITION 2.6.4. Suppose

$$
P_{1} \times P_{2}\left\{\tau^{-\frac{1}{2}}\left|\int_{0}^{\tau} \phi(s) d s-\sum_{i=1}^{N(0, \tau]} x_{i}\right| \geq \varepsilon_{\tau}\right\} \leq \beta_{\tau} ;
$$

then we set $X(\tau)=\max \left(\varepsilon_{\tau}, \beta_{\tau}\right)$.
With this definition we have a simple estimate of rate of convergence to normality of the 'number of busy servers' process.

THEOREM 2.6.5. Suppose there are constants $m, \mu, \sigma_{1}, \sigma_{2}$ such that the armival process of the $G / G / \infty$ queue satisfies $(N(0, \tau]-m \tau) / \sigma_{1} \sqrt{\tau} \xrightarrow{D_{1}} N_{1}$ at rate $\varphi(\tau)$, and $\sum_{i=1}^{n}\left(x_{i}-\mu\right) / \sigma_{2} \sqrt{n} \xrightarrow{D_{2}} N_{2}$ at rate $\psi(\tau)$, then

$$
\left(\int_{0}^{\tau} \phi(s) d s-m \mu \tau\right) / \sigma \sqrt{\tau} \xrightarrow{D_{1} \times D_{2}} N
$$

where $N \equiv\left(\mu \sigma_{1} N_{1}+\sqrt{m} \sigma_{2} N_{2}\right) / \sigma$ at rate

$$
\begin{equation*}
A_{2} \max \left[\varphi(\tau), \sum_{k=0}^{\infty} \psi(k) P_{1}\{N(0, \tau]=k\}, \chi(\tau),(\log \tau) / \sqrt{\tau}\right] \tag{6.1}
\end{equation*}
$$

for some constant $A_{2}$, where $\sigma^{2}=\mu^{2} \sigma_{1}^{2}+m \sigma_{2}^{2}$.
Proof. Write $N_{\tau}=N(0, \tau]$. A's will denote constants. We firstly note $P_{1}\left\{\left(N_{\tau}+1\right)^{\frac{7}{2}} / N_{\tau}^{\frac{3}{2}} \geq \tau^{-\frac{1}{2}}\right\} \leq A_{3} \tau^{-\frac{1}{2}}$ for some constant $A_{3}$. By Lemmas 2.6.3 and 2.6.1,

$$
\nu\left(\sqrt{m}\left(N_{\tau}-m \tau\right) / \sigma_{1}\left(N_{\tau}+1\right)^{\frac{3}{2}}, N_{1}\right) \leq A_{4} \max (\varphi(\tau),(\log \tau) / \sqrt{\tau}) ;
$$

hence, by Lemma 2.6.2,

$$
\begin{array}{r}
\nu\left[\mu \sigma_{1} \sigma^{-1} \sqrt{m}\left(N_{\tau}-m \tau\right) / \sigma_{1}\left(N_{\tau}+1\right)^{\frac{3}{2}}+\sqrt{m} \sigma_{2} \sigma^{-1} \sum_{i=1}^{N}\left(x_{i}-\mu\right) / \sigma_{2}\left(N_{\tau}+1\right)^{\frac{3}{2}}, N\right] \\
\leq A_{5} \max \left[\varphi(\tau), \sum_{k=0}^{\infty} \psi(k) P_{1}\{N(0, \tau]=k\},(\log \tau) / \sqrt{\tau}\right] .
\end{array}
$$

where we have defined $N \equiv\left(\mu \sigma_{1} N_{1}+\sqrt{m \sigma_{2}} N_{2}\right) / \sigma$. Thus we obtain by Lemma 2.6.3,
$\nu\left[\left(\sum_{i=1}^{N} x_{i}-m \mu \tau\right] / \sigma \sqrt{\tau}, N\right] \leq A_{6} \max \left[\varphi(\tau), \sum_{k=0}^{\infty} \psi(k) P_{1}\{N(0, \tau]=k\},(\log \tau) / \sqrt{\tau}\right]$.
By Definition 2.6.4 and Lemma 2.6.1 the conclusion follows.
A valid enquiry is into the technical handling of the second term in (6.1). Since many rates of convergence results have a $\psi(k)$ of the form $\psi(k)=\{\log (1+k)\}^{\gamma} /(1+k)^{\delta}, \gamma \geq 0, \delta>0$ we prove

LEMMA 2.6.6. (a) If $\psi(k)=\{\log (1+k)\}^{\gamma} /(1+k)^{\delta}$, some $\gamma \geq 0$, $\delta>0$ and $E_{1}\{N(0, n]\}=O(n)$, then

$$
E_{1}\{\psi(N(0, n])\}=O\left(\max \left[\{\log (1+n)\}^{\gamma} \varphi(n), \psi(n)\right]\right)
$$

(b) If also, $\quad \lim \sup E_{1}|N(0, n]-m m|^{1+\delta} / n<\infty$, then $E_{1}\{\psi(N(0, n])\}=O(\psi(n))$.

Remark. If $0<\delta \leq 1$, condition ( $b$ ) holds if
$\lim \sup \operatorname{Var} N(0, n] / n<\infty$, true for many point processes. $n \rightarrow \infty$

Proof. Write $N_{n}=N(0, n], C_{n} \equiv\{N(0, n] \geq n m / 2\}$. Then

$$
\begin{aligned}
\int_{C_{n}} \frac{\left\{\log \left(1+N_{n}\right)\right\}^{\gamma}}{\left(1+N_{n}\right)^{\delta}} d P_{1} & \leq(2 / n m)^{\delta} E_{1}\left\{\log \left(1+N_{n}\right)\right\}^{\gamma} \\
& \leq A_{7}\{\log (1+n)\}^{\gamma} / n^{\delta}
\end{aligned}
$$

by Jensen's inequality, since $\{\log (1+x)\}^{\gamma}$ is a concave function. Also to prove (a),
$\int_{C_{n}^{c}} \frac{\left\{\log \left(1+N_{n}\right)\right\}^{\gamma}}{\left(1+N_{n}\right)^{\delta}} d P_{1} \leq A_{7}^{\prime}\{\log (1+n)\}^{\gamma} P_{1}\left(C_{n}^{c}\right)$

$$
\begin{align*}
& \leq A_{7}^{\prime}\{\log (1+n)\}^{\gamma} P_{[ }\left\{\left|N_{n} / n-m\right|>m / 2\right\}  \tag{6.2}\\
& \leq A_{7}^{\prime}\left(\{\log (1+n)\}^{\gamma} / n^{\delta}\right)\left\{n^{\delta}[1-\Phi(\sqrt{n} / 2 \sigma)]\right\}+\{\log (1+n)\}^{\gamma} \varphi(n)
\end{align*}
$$

as in the proof of Lemma 2.6.2. But

$$
n^{\delta}\{1-\Phi(\sqrt{n} / 2 \sigma)\} \leq(2 \sigma / \sqrt{2 \pi}) n^{\delta-\frac{z_{2}}{2}} e^{-n / 2 \sigma^{2}} \rightarrow 0 \quad(n \rightarrow \infty) .
$$

(b) follows by applying Chebyshev's inequality to (6.2).

We now turn to estimation of $X(\tau)$. Our first result follows easily from Chebyshev's inequality: since all of our results seem to go through most naturally with stationary arrival processes, we make this assumption explicitly from now on.

LEMMA 2.6.7. If the armival process is stationary with $m \equiv E_{1} N(0,1]<\infty$ and the sequence of service times is stationary with finite second moment, then $X(\tau)=O\left(\tau^{-\frac{3}{4}}\right)$.

In our next two estimates we assume the service times to be i.i.d.
LEMMA 2.6.8. If the arrival process is stationary and orderly with $m \equiv E_{1} N(0,1]<\infty$ and (4.4) and (4.5) hold and the service times are i.i.d. with finite third moment, then $\chi(\tau)=O\left(\tau^{-1 / 3}\right)$.

Proof. Once again we use Chebyshev. Following Daley (1972), Theorem 4 , we assert that from the stationarity and orderliness of $P_{1}$,

$$
\begin{equation*}
d u|d v(v)|=P_{1}\{N(u, u+d u]=1, N(u+v, u+v+d v]=1\}+o(d u d v) \tag{6.3}
\end{equation*}
$$

(where $m \equiv 1$ ), so that, using the decompositions (2.2) and (2.3),

$$
\begin{align*}
E_{1} \times E_{2}^{\prime}\left\{\Phi^{\prime}(\tau)^{2}\right\}= & \int_{0}^{\tau} d u \int_{0}^{\infty}\left[(u+\tau-\tau)^{+}\right]^{2} d F(\tau)  \tag{6.4}\\
& +\int_{0}^{\tau} d u \int_{0}^{u} d U(v) \int_{0}^{\infty} d F(\tau) \int_{0}^{\infty}(u+\tau-\tau)^{+}(u-v+m-\tau)^{+} d F(m)  \tag{6.5}\\
& +\int_{0}^{\tau} d u \int_{0}^{\tau-u} d U(v) \int_{0}^{\infty} d F(\imath) \int_{0}^{\infty}(u+\tau-\tau)^{+}(v+m-\tau)^{+} d F(m) \tag{6.6}
\end{align*}
$$

where $F$ is the service time distribution function. Easy manipulation then yields $\int_{0}^{\infty} Z^{3} d F(Z)<\infty$ as an upper bound for (6.4), whereas (6.5) and (6.6) have upper bounds of $A_{8} \int_{0}^{\infty} \tau^{3} d F(\tau) \int_{0}^{\infty} \tau d F(\tau)$. The term $E_{1} \times E_{2}\left\{\Phi_{2}(\tau)^{2}\right\}$ has similar bounds.

Unless the arrival process is renewal, the technique used in Lemma 2.6 .8 cannot readily be extended to higher moments. However, we make

CONJECTURE 2.6.9. If the arrival process is stationary and orderly with $m \equiv E_{1} N(0,1]<\infty$ and (4.4) and (4.5) hold, and the service times are i.i.d. with finite $q$ th moment, then $\chi(\tau)=O\left(\tau^{-\frac{1}{2}(q-1) / q}\right)$.

By making some rough approximations, we can achieve a lower bound for $X(\tau)$. Here, as is often done, we ignore the contributions of arrivals before time zero to the traffic average.

LEMMA 2.6.10. If the arrival process is stationary with finite $p_{1}$ th moment, and the service times are i.i.d. with finite qth moment, then, setting $p_{2} \equiv(\sqrt{1+q}-1), p \equiv \min \left(p_{1}, p_{2}\right), \quad r=\min \left(p_{1}, 2\right)$,

$$
\begin{aligned}
x(\tau) & =O\left(\tau^{\left.-\frac{3}{2}\left(1-p q^{-1}\right) p(p+1)^{-1}\right)}, \quad p>2,\right. \\
& =O\left(\tau^{-\frac{1}{2}\left(1-2 q^{-1}\right) r(r+1)^{-1}}\right), \quad p \leq 2 .
\end{aligned}
$$

Proof. We use the decomposition (2.4) and the inequality of von Bahr and Lisseen (1965) and Dharmadhikari, Fabian and Jogdeo (1968) for
$\left\{X_{n}: n \geq 1\right\}$ a sequence of i.i.d. random variables with $E X_{1}=\nu<\infty$ : if $l \geq 1$, and $S_{k}=\sum_{i=1}^{k} X_{i}$, then

$$
E\left|S_{k}-k \nu\right|^{2} \leq A_{9} k^{\max \left(1, \frac{3}{2} \tau\right)} E\left|X_{1}-\nu\right|^{2},
$$

for some constant $A_{9}$. We also require that $M_{n}(u)^{1 / n} \equiv\left\{E_{1} N(0, u]^{n}\right\}^{1 / n}$ is a sub-additive function of $u$ for stationary arrival processes $P_{1}$ (easily deduced as in Daley (1971)), and hence that for some constant $A_{10}$, $E_{1} N(0, u]^{n} \leq A_{10} u^{n}$.

Thus, by Chebyshev's inequality, with $\mu=E_{2}\left\{x_{1}\right\}$ as usual,
$P_{1} \times P_{2}\left\{\Phi_{1}^{\prime}(\tau) \geq \varepsilon_{\tau}\right\}$

$$
\begin{aligned}
& \leq E_{1} \times E_{2}\left\{\Phi_{1}^{\prime}(\tau)^{\delta}\right\} / \varepsilon_{\tau}^{\delta} \\
& \leq 2^{\delta}\left\{E_{1}\left\{N(\tau-c(\tau), \tau]^{\delta}\right\} \mu^{\delta}+E_{1}\left\{N(\tau-c(\tau), \tau]^{\max \left(1, \frac{7}{2} \delta\right)}\right\} E_{2}\left|x_{1}-\mu\right|^{\delta}\right\} / \varepsilon_{\tau}^{\delta} \\
& \leq A_{11} c(\tau)^{\delta} / \varepsilon_{\tau}^{\delta} .
\end{aligned}
$$

Also,

$$
\left.\begin{array}{l}
P_{1} \times P_{2}\left\{\Phi_{0}(\tau) \geq \varepsilon_{\tau}\right\} \leq A_{12} \tau^{\max \left(1, \frac{3}{2} \delta\right)} E_{2}\left\{\left|\left(x_{1}-c(\tau)\right)^{+}-E_{2}\left(x_{1}-c(\tau)\right)^{+}\right|^{\delta}\right\} / \varepsilon_{\tau}^{\delta} \\
\\
+A_{13} \tau^{\delta}\left[E_{2}\left(x_{1}-c(\tau)\right)^{+}\right]^{\delta /} / \varepsilon_{\tau}^{\delta}
\end{array}\right\} \begin{aligned}
& \leq A_{14}\left[\tau^{\max \left(1, \frac{7}{2} \delta\right)} / \varepsilon_{\tau}^{\delta}\right] E_{2}\left\{\left[\left(x_{1}-c(\tau)\right)^{+}\right]^{\delta}\right\} \\
&  \tag{6.8}\\
& +A_{15}\left(\tau^{\delta} / \varepsilon_{\tau}^{\delta}\right\}\left[E_{2}\left(x_{1}-c(\tau)\right)^{+}\right]^{\delta},
\end{aligned}
$$

In $(6.8)$, if we set $c(\tau)=\tau^{\alpha}, 0<\alpha \leq \frac{1}{2}$, then

$$
\tau c(\tau)^{-1} E_{2}\left(x_{1}-d(\tau)\right)^{+} \leq \tau c(\tau)^{-1} \int_{d(\tau)}^{\infty} \tau d F(\tau) \leq \int_{\tau^{\alpha}}^{\infty} \tau^{1 / \alpha} d F(\tau)
$$

and hence, provided $\int_{0}^{\infty} \tau^{1 / \alpha} d F(\tau)<\infty$,

$$
(6.8) \leq A_{16}\left(\tau^{\alpha} / \varepsilon_{\tau}\right)^{\delta}
$$

Similarly,

$$
(6.7) \leq A_{17}\left(\tau^{\alpha} / \varepsilon_{\tau}\right)^{\delta}
$$

provided $\int_{0}^{\infty} \tau^{(\delta / 2 \alpha)} d F(\tau)<\infty \quad(\delta \geq 2)$, or $\int_{0}^{\infty} \eta^{1 / \alpha} d F(\tau)<\infty \quad(\delta<2)$.
Hence

$$
P_{1} \times P_{2}\left\{\Phi^{\prime}(\tau) / \sqrt{\tau} \geq \varepsilon_{\tau}\right\} \leq A_{18}\left(\tau^{-\left(\frac{1}{2}-\alpha\right)} / \varepsilon_{\tau}\right)^{\delta}
$$

i.e., $X(\tau)=\max \left[\varepsilon_{\tau},\left(\tau^{-\left(\frac{z}{2}-\alpha\right)} / \varepsilon_{\tau}\right)^{\delta}\right]$. Choosing $\varepsilon_{\tau}$ optimally gives $x(\tau)=\tau^{-\left(\frac{1}{2}-\alpha\right)[\delta /(\delta+1)]}$.

Thus the problem becomes to maximize $f(\alpha, \delta)=\left(\frac{1}{2}-\alpha\right)[\delta /(\delta+1)]$ subject to the constraints $0<\delta \leq \min \left(p_{1}, q\right),(1 / \alpha) \leq q$ and $(\delta / 2 \alpha) \leq q$ (for $\delta \geq 2$ ). If $p_{1}<2$, the answer clearly is to choose $\alpha=1 / q$ and $\delta=q_{1}$. If $p_{1} \geq 2$, the situation is not so simple. Drawing an ( $\alpha, \delta$ ) graph including the constraint regions, and observing that $f(\alpha, \delta) \uparrow, \alpha \downarrow$ and $f(\alpha, \delta) \uparrow, \delta \uparrow$, one sees that the optimal $(\alpha, \delta)$ pair lie on the line $\delta=2 \alpha q$, where $2 \leq \delta \leq \min \left(p_{1}, q\right)$. Thus the problem evolves into maximising $g(\delta)=\left(\frac{1}{2}-(\delta / 2 q)\right) \delta(\delta+1)^{-1}$ subject to this constraint. It is easily demonstrated that $g(\delta)$ has two turning points only, at $\delta= \pm \sqrt{q+1}-1$. The one which concerns us, $\sqrt{q+1}-1$, is always a maximum. Thus, writing $p_{2}=\sqrt{q+1}-1, p=\min \left(p_{1}, p_{2}\right)$, we find $\chi(\tau)=\tau^{-g(p)}$
unless $p_{2}<2$, i.e., $q<8$ when it becomes $\chi(\tau)=\tau^{-g(2)}$.
A point to note here is that as $p_{1}, q \rightarrow \infty$ (provided $p_{1} \approx q$ )
$X(\tau) \rightarrow \tau^{-\frac{1}{2}}$, which one would hope for. Indeed, for large $q$, $\chi(\tau) \sim \tau^{-\frac{3}{2}(1-1 / \sqrt{q})}$.

We focus attention now on $\varphi(\tau)$. For point processes specified in terms of counting measures, we will assume $\varphi(\tau)$ can be calculated; however, many point processes on the real line are more easily defined via their inter-arrival times, and for these we provide

LEMMA 2.6.11. Suppose the inter-armival times $\left\{\pi_{i}\right\}$ of a point process satisfy, for some constants $\lambda$ and $\sigma$,

$$
\sup _{y}\left|P\left\{\left(\sum_{i=0}^{n}\left(\pi_{i}-\lambda\right)\right) /(\sigma \sqrt{n}) \leq y\right\}-\Phi(y)\right| \leq \varphi(n) ;
$$

then, for any $\delta<1$,
$\sup \left|P\left\{\left(N(0, \tau]-\lambda^{-1} \tau\right) /\left(\lambda^{-3 / 2} \sigma \sqrt{\tau}\right) \leq y\right\}-\Phi(y)\right|$
$y$

$$
\begin{equation*}
=O\left(\max \left[\varphi\left(\lambda^{-1} \delta \tau\right),(\log \tau) / \sqrt{\tau}\right]\right), \tag{6.9}
\end{equation*}
$$

where $\Phi$ is the cumulative distribution function of the normal distribution.
Remark. For the purposes of this proof, we set $\Pi_{0}=t_{1}$.
Proof. We divide the real line into four segments.
(i) $y>(\log \tau)^{\frac{1}{2}}:$ Let $\gamma_{\tau} \equiv \lambda^{-1} \tau+\lambda^{-3 / 2} \sigma(\tau \log \tau)^{\frac{3}{2}}$, and $[x]$ denote the smallest integer less than $x$; then

$$
\begin{align*}
& \left|P\left\{\left(N(0, \tau]-\lambda^{-1} \tau\right) /\left(\lambda^{-3 / 2} \sigma \sqrt{\tau}\right) \leq y\right\}-\Phi(y)\right|  \tag{6.10}\\
& =\left|P\left\{\left(N(0, \tau]-\lambda^{-1} \tau\right) /\left(\lambda^{-3 / 2} \sigma \sqrt{\tau}\right)>y\right\}-[1-\Phi(y)]\right| \\
& \leq P\left\{N(0, \tau]>\gamma_{\tau}\right\}+e^{-\frac{3}{2} \log \tau} /(2 \pi \log \tau)^{\frac{3}{2}} \\
& =P\left\{\left(\sum_{i=0}^{\left[\gamma_{\tau}\right]}\left(\Pi_{i}-\lambda\right)\right) /\left(\sigma \sqrt{\left[\gamma_{\tau}\right]}\right) \leq\left(\tau-\lambda\left(\left[\gamma_{\tau}\right]+1\right)\right) /\left(\sigma \sqrt{\left[\gamma_{\tau}\right]}\right)\right\}+(2 \pi \tau \log \tau)^{-\frac{3}{2}} \\
& \leq \varphi\left(\left[\gamma_{\tau}\right]\right)+\Phi\left(\left(\tau-\lambda \gamma_{\tau}\right) /\left(\lambda \sqrt{\gamma_{\tau}}\right)\right)+(2 \pi \tau \log \tau)^{-\frac{3}{2}} \\
& \leq \varphi\left(\lambda^{-1} \tau\right)+\Phi\left(-\left(\lambda^{-1} \tau(\log \tau) / \gamma_{\tau}\right)^{\frac{1}{2}}\right)+(2 \pi \tau \log \tau)^{-\frac{3}{2}} \tag{6.11}
\end{align*}
$$

remembering that $\tau-\lambda \gamma_{\tau}$ is negative. In the algebra we have used an inequality on p. 175 of Feller (1968), which also yields an upper bound for the middle term of (6.11) of

$$
\begin{equation*}
\left\{\gamma_{\tau} /\left(\lambda^{-1} \tau \log \tau\right)\right\}^{\frac{1}{2}} \exp \left\{-\frac{1}{2}\left(\lambda^{-1} \tau(\log \tau) / \gamma_{\tau}\right)\right\} . \tag{6.12}
\end{equation*}
$$

Let us call the exponential factor in (6.12) $f(\tau)$. Then

$$
\log (\sqrt{\tau} f(\tau))=\frac{3}{2}(\log \tau)\left\{\lambda^{-\frac{1}{2}} \sigma(\tau \log \tau)^{\frac{3}{2}} / \gamma_{\tau}\right\} \rightarrow 0,(\tau \rightarrow \infty)
$$

so that $f(\tau)=O\left(\tau^{-\frac{3}{2}}\right)$. Since $\gamma_{\tau}=O(\tau),(6.12)$ is $O(\tau \log \tau)^{-\frac{1}{2}}$. Hence (6.10) is

$$
O\left(\max \left[\varphi\left(\lambda^{-1} \tau\right),(\tau \log \tau)^{-\frac{1}{2}}\right]\right)
$$

(ii) $y<-(\log \tau)^{\frac{7}{2}}$ : By similar methods to those of (i)

$$
(6.10)=O\left(\max \left[\varphi\left(\lambda^{-1} \delta \tau\right), \tau^{-\frac{1}{2}}\right]\right)
$$

for any $\delta<1$.
(iii) $0 \leq y \leq(\log \tau)^{\frac{3}{2}}:$ Let $\delta_{\tau} \equiv \delta_{\tau}(y) \equiv \lambda^{-1} \tau+\lambda^{-3 / 2} \sigma \tau^{\frac{3}{2}} y$, and set $\Phi^{*}(y)=1-\Phi(y)$. Then $P\left\{\left(N(0, \tau]-\lambda^{-1} \tau\right) /\left(\lambda^{-3 / 2} \sigma \sqrt{\tau}\right) \leq y\right\}$

$$
\begin{align*}
& =P\left\{\left(\sum_{i=0}^{\left[\delta_{\tau}\right]}\left(\pi_{i}-\lambda\right)\right) /\left(\sigma \sqrt{\left[\delta_{\tau}\right]}\right)>\left(\tau-\lambda\left(\left[\delta_{\tau}\right]+1\right)\right) /\left(\sigma \sqrt{\left[\delta_{\tau}\right]}\right)\right\}  \tag{6.13}\\
& \leq \varphi\left(\left[\delta_{\tau}\right]\right)+\Phi^{*}\left(\left(\tau-\lambda\left(\delta_{\tau}+1\right)\right) /\left(\sigma \lambda^{-\frac{3}{2}} \sqrt{\tau}\right)\right), \tag{6.14}
\end{align*}
$$

remembering that $\tau-\lambda\left(\left[\delta_{\tau}\right]+1\right)$ is negative. Now $|\Phi(x)-\Phi(y)| \leq|x-y| / \sqrt{2 \pi}$, so that by adding and subtracting $\Phi^{*}\left(\left(\tau-\lambda \delta_{\tau}\right) /\left(\sigma \lambda^{-\frac{3}{2}} \sqrt{\tau}\right)\right)$ to the last term in (6.14), we obtain an upper bound of (N.B. $\Phi^{*}(-y)=\Phi(y)$ )

$$
\varphi\left(\lambda^{-1} \tau\right)+\left(2 \pi \lambda \sigma^{2} \tau\right)^{-\frac{3}{2}}+\Phi(y)
$$

In the other direction, $(6.13)$ has a lower bound of
$-\varphi\left(\left[\delta_{\tau}\right]\right)+\Phi *\left(\left(\tau-\lambda \delta_{\tau}\right) / \sigma \sqrt{\delta_{\tau}}\right)$

$$
\begin{aligned}
& \geq-\varphi\left(\lambda^{-1} \tau\right)+\Phi^{*}\left(-y\left(1+\lambda^{-\frac{1}{2}} \sigma y / \sqrt{\tau}\right)^{-\frac{1}{2}}\right) \\
& \geq-\varphi\left(\lambda^{-1} \tau\right)+\Phi *(-y)-(|y| / \sqrt{2 \pi})\left|\left(1+\lambda^{-\frac{1}{2}} \sigma y / \sqrt{\tau}\right)^{-\frac{1}{2}}-1\right| \\
& \geq-\varphi\left(\lambda^{-1} \tau\right)+\Phi(y)-(\log \tau)^{\frac{1}{2}}\left(\lambda^{-\frac{1}{2}} \sigma(\log \tau)^{\frac{7}{2}} / \sqrt{\tau}\right) / \sqrt{2 \pi} .
\end{aligned}
$$

Hence

$$
\begin{aligned}
(6.10) & =O\left(\max \left[\varphi\left(\lambda^{-1} \tau\right),(\log \tau) / \sqrt{\tau}\right]\right) . \\
\text { (iv) }-(\log \tau)^{\frac{3}{2}} \leq y & \leq 0: \text { By methods similar to those of (iii), } \\
(6.10) & =O\left(\max \left[\varphi\left(\lambda^{-1} \delta \tau\right),(\log \tau) / \sqrt{\tau}\right]\right)
\end{aligned}
$$

for any $\delta<1$.
It would interesting to know if the $(\log \tau) / \sqrt{\tau}$ term in (6.9) can be replaced by $1 / \sqrt{\tau}$. Applying our result to the $G I / G / \infty$ queue, we obtain (see Problem 6.3.5)

COROLLARY 2.6.12. If the inter-arrival times and the service times of the $G I / G / \infty$ queue have finite third moments, then

$$
\left.\int_{0}^{\tau} \phi(s) d s-m E_{2}\left\{x_{1}\right\} \tau\right) / \sigma \sqrt{\tau} \xrightarrow{D_{1} \times D_{2}} N
$$

at rate $O\left(\tau^{-1 / 3}\right)$.
As remarked earlier, the metric used up till now is not continuous under Hölder continuous or Lipschitz mappings (Whitt (1974b)). Two suitable metrics are the Prohorov metric $\rho$ and the dual bounded Lipschitz metric $\beta$ (see (5.4) and (5.8) of Chapter 1). Our theorems (except possibly Lemma 2.6.11 and Corollary 2.6.12) hold with these metrics if we prove

LEMMA 2.6.13. Suppose $\beta\left(Y_{\tau}, Y\right) \leq \psi(\tau) \quad\left[\rho\left(Y_{\tau}, Y\right) \leq \psi(\tau)\right]$, and $\beta\left(Z_{\tau}, Z\right) \leq \varphi(\tau) \quad\left[\rho\left(Z_{\tau}, Z\right) \leq \varphi(\tau)\right]$. Then if $N_{\tau}$ is a sequence of positive random variables such that $N_{\tau} \rightarrow+\infty$ in probability, and $\left\{\left(N_{\tau}, z_{\tau}\right)\right\}$ is independent of $\left\{Y_{\tau}\right\}$, then
(a) $\beta\left[\left(Y_{N_{\tau}}, Z_{\tau}\right),(Y, Z)\right] \leq \varphi(\tau)+\int_{0}^{\infty} \psi(x) d P_{r}\left\{N_{\tau} \leq x\right\}$,
and, if in addition,

$$
\begin{aligned}
& \psi(\tau)=(\log \tau)^{\gamma} / \tau \gamma^{\gamma}, \gamma>0, \delta>0 \text { and } \\
& \lim \sup E\left|N_{\tau}-m \tau\right|^{l+\delta} / \tau<\infty \text { for some constant } m \text {, then }
\end{aligned}
$$

(b)

$$
\rho\left[\left(Y_{N_{\tau}}, Z_{\tau}\right),(Y, Z)\right]=O[\max (\varphi(\tau), \psi(\tau))]
$$

Proof. (a) Let $f$ denote an arbitrary bounded real-valued Lipschitz function on $R^{2}$ with norm $\leq 1$; i.e., $\|f\| \equiv \sup _{x \neq y}(|f(x)-f(y)|) /\left(\left[\sum_{i=1}^{2}\left(x_{i}-y_{i}\right)^{2}\right]^{\frac{3}{2}}\right)+\sup \left\{|f(x)|, x \in R^{2}\right\} \leq 1$. $\left(x, y \in R^{2}\right)$

Then the functions $g_{1}(x)=f\left(x, x_{2}\right), x_{2}$ fixed, and $g_{2}(x)=f\left(x_{1}, x\right), x_{1}$ fixed, are clearly bounded real-valued Lipschitz functions on $R^{1}$ with norm $\leq 1$, and hence ( $\int$ will mean $\int_{-\infty}^{+\infty}$ ) $\iint f(y, z) \operatorname{Pr}\left\{\left(Y_{N_{\tau}}, Z_{\tau}\right) \in(d y, d z)\right\}$

$$
\begin{aligned}
& =\int_{0}^{\infty} \iint\left[f(y, z) \operatorname{Pr}\left\{Y_{\imath} \in d y\right\}\right] \operatorname{Pr}\left\{z_{\tau} \in d z, N_{\tau} \in d \tau\right\} \\
& \leq \int_{0}^{\infty} \psi(\tau) \operatorname{Pr}\left\{N_{\tau} \in d \tau\right\}+\int\left[\int f(y, z) \operatorname{Pr}\left\{Z_{\tau} \in d z\right\}\right] \operatorname{Pr}\{Y \in d y\} \\
& \leq E\left\{\psi\left(N_{\tau}\right)\right\}+\varphi(\tau)+\iint f(y, z) \operatorname{Pr}\{(y, z) \in(d y, d z)\} .
\end{aligned}
$$

A similar inequality may be found in the other direction.
(b) Let $F$ be any closed set in $B\left(R^{2}\right)$. Then $F_{x} \equiv\left\{x_{1}:\left(x_{1}, x\right) \in F\right\}$ is a closed set in $B(R)$, and hence

$$
\begin{align*}
& \operatorname{Pr}\left\{\left(Y_{N_{\tau}}, Z_{\tau}\right) \in F\right\} \\
& =\int_{0}^{\infty} \int \operatorname{Pr}\left\{Y_{\tau} \in F_{x}\right\} \operatorname{Pr}\left\{Z_{\tau} \in d x, N_{\tau} \in d \tau\right\} \\
& \leq \int_{0}^{\infty} \psi(l) \operatorname{Pr}\left\{N_{\tau} \in d l\right\}+\int_{0}^{\infty} \int \operatorname{Pr}\left\{y \in F_{x}^{\psi(Z)}\right\} \operatorname{Pr}\left\{Z_{\tau} \in d x, N_{\tau} \in d l\right\} .  \tag{6.15}\\
& \text { Let } \theta<1 \text {; then splitting the } \int_{0}^{\infty} \text { in the second term of (6.15) into } \\
& \int_{(0, m \theta \tau]} \text { and } \int_{(m \theta \tau, \infty)} \text { we obtain an upper bound for it of } \\
& \operatorname{Pr}\left\{N_{\tau} \leq m \theta \tau\right\}+\operatorname{Pr}\left\{\left(Y, Z_{\tau}\right) \in F^{\psi(m \theta \tau)}\right\}  \tag{6.16}\\
& \leq A_{19} / \tau^{\delta}+\operatorname{Pr}\left\{\left(Y, Z_{\tau}\right) \in \overline{F^{\psi(m \theta \tau)}}\right\} \\
& \leq A_{19}(\log \tau)^{\gamma} / \tau^{\delta}+\varphi(\tau)+\operatorname{Pr}\left\{(Y, Z) \in\left(\overline{F^{\psi(m \theta \tau)}}\right)^{\varphi(\tau)}\right\} \\
& \leq A_{19} \psi(\tau)+\varphi(\tau)+\operatorname{Pr}\left\{(Y, Z) \in F^{2 \psi(m \theta \tau)+\varphi(\tau)}\right\}
\end{align*}
$$

using Chebyshev's inequality for the first term in (6.16). For our $\psi(\tau)$, $\psi(m \theta \tau)=O(\psi(\tau))$, and the first term in (6.15) is $\leq A_{20} \psi(\tau)$, by Lemma 2.6.6, so that the conclusion follows.

As usual, in this theorem it has been tacitly understood that $Y_{\tau}, N_{\tau}$ and $Z_{\tau}$ are all on the same space $\left(\Omega_{1}, F_{1}, P r_{1}\right)$ and $Y$ and $Z$ are on a possibly different space $\left(\Omega_{2}, F_{2}, P_{2}\right)$. In (6.16), we have assumed $Y$ and $Z_{\tau}$ on the same space: to be rigorous, one should work the proof through on the product space with measure $P r_{1} \times P r_{2}$. However, we still need not demand that the limit variables be on the same space as the others in the final analysis.

### 2.7. Rates of functional convergence

Ideas related to speeds of convergence have received a boost since Rosencrantz' illuminating paper in 1968. The results of this section
represent a generalisation of Section 3 of Kennedy (1972b), which is partly based on Rosencrantz' work. According to the usual procedure (and Kennedy's), our random functions will take values in the space $D[0,1]$ of right-continuous functions on [0, 1] with left-hand limits. We will consider rates of convergence in terms of the metric

$$
\begin{equation*}
\alpha(X, Y) \equiv \inf \{\varepsilon \geq 0: \operatorname{Pr}\{\rho(X, Y) \geq \varepsilon\} \leq \varepsilon\}, \tag{7.1}
\end{equation*}
$$

where $\rho(x, y) \equiv \sup _{0 \leq t \leq 1}|x(t)-y(t)|, x, y \in D[0,1]$. Note that this
definition is only meaningful if $X$ and $Y$ are on the same space $(\Omega, F, P r)$, say, a fact which we will assume from now on. $\alpha$ corresponds to convergence in probability (Section 1.5) in the uniform metric on $D[0,1]$ (see Problem 6.3.6). Observe that if $D[0,1]$ is endowed only with the Skorokhod topology (Billingsley (1968), p. 111), it is easy to prove $\{\rho(X, Y) \geq \varepsilon\} \in F$, or, more generally, that $\rho(X, Y)$ is a random variable on ( $\Omega, F$ ) . Hence, the difficulty referred to in Section 2 of Whitt (1974b) of $D$ being non-separable with the uniform topology (Billingsley (1968), p. 150) is avoided.

In this section we will regard the arrival process as a point process $\eta$ and the service times as a sequence $\left\{V_{i}\right\}$ of non-negative random variables on $\left(\Omega, F, P_{r}\right)$, and also suppose that standard Brownian motions $W_{1}$ and $W_{2}$ are defined on $\Omega$, such that $\left(\eta, W_{1}\right)$ and $\left(\left\{V_{i}\right\}, W_{2}\right)$ are independent. We will require the following random functions in $D[0,1]$ :

$$
\begin{align*}
& A_{n}(t) \equiv(n(0, n t]-m n t) /\left(\sigma_{1} \sqrt{n}\right)  \tag{7.2}\\
& S_{n}(t) \equiv\left(\sum_{i=1}^{[n t]}\left(V_{i}-\mu\right)\right) /\left(\sigma_{2} \sqrt{n}\right) \tag{7.3}
\end{align*}
$$

for some appropriate constants $m, \mu, \sigma_{1}$ and $\sigma_{2}$, and

$$
\begin{equation*}
Q_{n}(t) \equiv\left(\sum_{i=1}^{n(0, n t]} V_{i}-m \mu n t\right) /(\sigma \sqrt{n}) \tag{7.4}
\end{equation*}
$$

where $\sigma^{2}=\mu^{2} \sigma_{1}^{2}+m \sigma_{2}^{2}$.

We first require a random change of time lemma.
LEMMA 2.7.1. If for the $G / G / \infty$ queue, $\alpha\left(A_{n}, W_{1}\right) \leq \varphi(n)$ and $\alpha\left(S_{n}, W_{2}\right) \leq \psi(n)$, then

$$
\begin{equation*}
\alpha\left(Q_{n}, W\right) \leq C_{1} \max \left[\varphi(n), \psi(n),(\log n)^{3 / 4} / n^{\frac{1}{4}}\right] \tag{7.5}
\end{equation*}
$$

for some constant $C_{1}$, where $W(\cdot)=\left(\mu \sigma_{1} W_{1}(\cdot)+\sigma_{2} W_{2}(m \cdot)\right) / \sigma$ is also a standard Brownian motion.

Remark. We will choose a scale such that $m<1$. For another remark see Problem 6.3.7.

Proof. Let $\gamma_{n}$ be a sequence $\downarrow 0(n \rightarrow \infty)$, and set $C_{2}=\max \left[2 m \sigma_{2} / \sigma, 16 \sigma_{1} / \sigma\right]$. Then
$\operatorname{Pr}\left\{\rho\left(Q_{n}, W\right) \geq C_{2} \gamma_{n}\right\}$

$$
\begin{align*}
\leq \operatorname{Pr}\left\{\rho\left(A_{n}, W_{1}\right) \geq\right. & \left.\frac{1}{2} C^{2} \sigma\left(\mu \sigma_{2}\right)^{-1} \gamma_{n}\right\}+\operatorname{Pr}\left\{\rho\left(S_{n}(m \cdot), W_{2}(m \cdot)\right) \geq \frac{3}{4} C_{2} \sigma \sigma_{1}^{-1} \gamma_{n}\right\}  \tag{7.6}\\
& +\operatorname{Pr}\left\{\rho\left(\sum_{i=1}^{n(0, n t]}\left(V_{i}-\mu\right) / \sigma_{1} \sqrt{n}, S_{n}(m \cdot)\right) \geq \frac{3}{4} C_{2} \sigma \sigma_{1}^{-1} \Upsilon_{n}\right\} . \tag{7.7}
\end{align*}
$$

Clearly, for $\delta_{n} \downarrow 0(n \rightarrow \infty)$, and provided $\varepsilon_{n} \equiv \delta_{n}+n^{-1} \leq 1-m$ (true for $n \geq$ some $n_{1}$ )

$$
\begin{align*}
\text { (7.7) } \leq & \operatorname{Pr}\left\{\sup _{0 \leq t \leq 1}\left|\frac{n(0, n t]}{n}-m t\right| \geq \delta_{n}\right\} \\
& \left.+\operatorname{Pr}\left\{\sup _{|s-t|<\varepsilon_{n}}| | \sum_{i=[n s]+1}^{[n t]}\left(V_{i}-\mu\right)\right] /\left(\sigma_{1} \sqrt{n}\right) \mid \geq 4 \gamma_{n}\right\} \\
= & \operatorname{Pr}\left\{\rho\left(A_{n}, 0\right) \geq \sqrt{n} \delta_{n} / \sigma_{2}\right\}+\operatorname{Pr}\left\{\omega_{S_{n}}\left(\varepsilon_{n}\right) \geq 4 \gamma_{n}\right\} \\
\leq & \operatorname{Pr}\left\{\rho\left(A_{n}, W_{1}\right) \geq \sqrt{n} \delta_{n} / 2 \sigma_{2}\right\}+\operatorname{Pr}\left\{\rho\left(W_{1}, 0\right) \geq \sqrt{n} \delta_{n} / 2 \sigma_{2}\right\} \\
& +\operatorname{Pr}\left\{\rho\left(S_{n}, W_{2}\right) \geq \gamma_{n}\right\}+\operatorname{Pr}\left\{\omega_{W_{2}}\left(\varepsilon_{n}\right) \geq \gamma_{n}\right\} \quad(7.8) \tag{7.8}
\end{align*}
$$

where $w_{x}(\cdot)$ is the modulus of continuity, and we have used inequality (5.14) of Chapter 1. We choose $\delta_{n}=\left(2 \sigma_{2}^{2} \log n\right)^{\frac{7}{2}} / \sqrt{n}$. Then by equation
(2.4) of Kennedy (1972b), the second term of (7.8) has an upper bound of $4 /\left(\pi n^{2} \log n\right)^{\frac{1}{2}}$. But $n^{\frac{1}{2}} \delta_{n} \rightarrow \infty \quad(n \rightarrow \infty)$, and hence, for some $n_{2}$, the first term of (7.8) has an upper bound of $\varphi(n), n \geq n_{2}$. The final term of (7.8) has, by Lemma 2.4 of Kennedy (1972b), an upper bound of

$$
24(2 \pi)^{-\frac{3}{2}} \gamma_{n}^{-1} \exp \left\{-\gamma_{n}^{2} /\left(18 \varepsilon_{n}\right)\right\} / \varepsilon_{n}^{\frac{3}{2}}
$$

for $n \geq$ some $n_{3}$. Take $\gamma_{n} \geq 6 \delta_{n}^{\frac{3}{2}} \theta$ ( $\geq \sqrt{18} \varepsilon_{n}^{\frac{3}{2}} \theta_{n}$ for $n \geq$ some $\left.n_{4}\right)$ for some $\theta_{n}$. Then the upper bound becomes (go back to (7.8) and reapply Kennedy's Lemma 2.4) $4 \pi^{-\frac{3}{2}} \exp \left\{-\theta_{n}^{2}\right\} / \varepsilon_{n} \theta_{n}$. Hence we take $\theta_{n}=(\log n)^{\frac{3}{2}}$, and obtain $(4 / \sqrt{\pi})(\log n)^{-1} n^{-\frac{3}{2}}$ for this rate, but only if $\gamma_{n} \geq 6 \sqrt{2} \sigma_{2}(\log n)^{3 / 4} / n^{\frac{3}{4}}$. Hence if

$$
\gamma_{n}=\max \left[\varphi(n), \psi(n), 6 \sqrt{2} \sigma_{2}(\log n)^{3 / 4} / n^{\frac{1}{4}}\right],
$$

we obtain an upper bound for (7.6) and (7.7) of

$$
2(\varphi(n)+\psi(n))+C_{3}(\log n)^{3 / 4} / n^{\frac{3}{4}}
$$

for some constant $C_{3}$, provided $n \geq \max \left[n_{1}, n_{2}, n_{3}, n_{4}\right]$. (7.5) then follows, where the constant $C_{1}$ may possibly need to be increased to account for the initial terms.

For our main result we need to define a random function derived from the number of servers process. Let

$$
\begin{equation*}
M_{n}(t) \equiv\left(\int_{0}^{n t} m(s) d s-m \mu n t\right) /(\sigma \sqrt{n}) \tag{7.9}
\end{equation*}
$$

where $m(s)$ is defined in Section 2.2, and $\sigma^{2}=\mu^{2} \sigma_{1}^{2}+m \sigma_{2}^{2}$.
We will also explicitly make the assumption now that the service times are stationary (but see Problem 6.3.4) so that $\mu=E V_{i}$, but we will not assume the arrival process is stationary. Now we can prove

THEOREM 2.7.2. If for the $G / G / \infty$ queue $\alpha\left(A_{n}, W_{1}\right) \leq \varphi(n)$ and $\alpha\left(S_{n}, W_{2}\right) \leq \psi(n)$, and in addition $E \eta(-u, u)=O(u)$ and $E\left(v_{i}^{2 p_{2}}\right)<\infty$ $\left(p_{2}>1\right)$, then

$$
\alpha\left(M_{n}, W\right) \leq C_{4} \max \left[\varphi(n), \psi(n), n^{-\left(p_{2}-1\right) /\left(2 p_{2}+1\right)},(\log n)^{3 / 4} / n^{\frac{3}{4}}\right],
$$

where $W$ is as defined in Lemma 2.7.1.
Proof. Clearly, for $\gamma_{n}+0$, and $C_{5}=\max \left[2 C_{1}, 64,8 / \sigma\right]$, $\operatorname{Pr}\left\{\rho\left(M_{n}, W\right) \geq C_{5} \gamma_{n}\right\} \leq \operatorname{Pr}\left\{\rho\left(M_{n}, Q_{n}\right) \geq C_{5} \gamma_{n} / 2\right\}+\operatorname{Pr}\left\{\rho\left(Q_{n}, W\right) \geq C_{5} \gamma_{n} / 2\right\}$

$$
\begin{aligned}
\leq \operatorname{Pr}\left\{\rho \left(\mid \int_{0}^{n t} m(s) d s-\right.\right. & \left.\left.\sum_{i=1}^{n(0, n t]} V_{i} \mid, 0\right) \geq C_{5} \sigma \gamma_{n} \sqrt{n / 2}\right\} \\
& +C_{1} \max \left[\varphi(n), \psi(n),(\log n)^{3 / 4} / n^{\frac{3}{4}}\right]
\end{aligned}
$$

provided $\quad \gamma_{n} \geq \max \left[\varphi(n), \psi(n),(\log n)^{3 / 4} / n^{\frac{3}{4}}\right]$. Now the first term has an upper bound (for some $c_{n} \uparrow \infty \quad(n \rightarrow \infty)$ ) of

$$
\begin{equation*}
\operatorname{Pr}\left\{M_{2}(n) \geq \frac{3}{4} C_{5} \sigma \gamma_{n} \sqrt{n}\right\}+\operatorname{Pr}\left\{M_{0}^{\prime}(n) \geq C_{5} \sigma \gamma_{n} \sqrt{n} / 8\right\} \tag{7.10}
\end{equation*}
$$

$$
\begin{equation*}
\left.+\operatorname{Pr}\left\{\rho \sum_{i=n\left(0, n t-c_{n}\right]+1}^{n(0, n t]} V_{i} / \sigma \sqrt{n}, 0\right) \geq C_{5} \gamma_{n} / 8\right\} \tag{7.11}
\end{equation*}
$$

in the notation of (2.5) and (2.6). If $E\left(V_{i}^{2}\right)<\infty$, then the first term of (7.10) is $O\left(n^{-\frac{1}{4}}\right)$ if $n^{-\frac{3}{4}}=O\left(\gamma_{n}\right)$ (cf. Lemma 2.6.7), which we require. Also,

$$
\begin{align*}
(7.11) & \leq \operatorname{Pr}\left\{\omega_{Q_{n}}\left(c_{n} / n\right) \geq c_{5} \gamma_{n} / 16\right\}+\operatorname{Pr}\left\{c_{n} / \sqrt{n} \geq C_{5} \gamma_{n} / 16\right\} \\
& \leq \operatorname{Pr}\left\{\rho\left(Q_{n}, W\right) \geq \gamma_{n}\right\}+\operatorname{Pr}\left\{\omega_{W}\left(c_{n} / n\right) \geq \gamma_{n}\right\} \tag{7.12}
\end{align*}
$$

provided $\gamma_{n} \geq c_{n} / \sqrt{n}$. By Chebyshev's inequality, the second term of (7.10) has an upper bound of

$$
\begin{equation*}
C_{6} n^{\frac{3}{2}} \gamma_{n}^{-1} \int_{a_{n}}^{\infty} \tau d F(Z) \tag{7.13}
\end{equation*}
$$

where $F$ is the distribution function of the service times. Hence, for this term to have an upper bound of $\gamma_{n}$, we require
a sufficient condition for which is

$$
\lim _{n \rightarrow \infty} n^{3 / 2} c_{n}^{-2} \int_{c}^{\infty} \eta d F(\tau)<\infty
$$

If $c_{n}=n^{\alpha}$, some $\alpha>0$, this occurs if $E\left(V_{i}^{(3 / 2 \alpha)-1}\right)<\infty$, and hence
$c_{n}=n^{\frac{z_{2}}{2}\left(p_{2}-1\right) /\left(2 p_{2}+1\right)}$, i.e., we require $\gamma_{n} \geq n^{-\left(p_{2}-1\right) /\left(2 p_{2}+1\right)}$. Finally, the second term in (7.12) has an upper bound of $c_{7} n c_{n}^{-3 / 2} \exp \left\{-c_{n} / 18\right\}$, using Lemma 2.4 of Kennedy (1972b), which converges to zero at a rate much faster than $c_{n} / \sqrt{n}$ and hence $\gamma_{n}$.

Applying these result to the $G I / G / \infty$ queue, we obtain
COROLLARY 2.7.3. If the arrival process to the $G / G / \infty$ queue is renewal with $E\left(\Pi_{i}^{2 p_{1}}\right)<\infty$, where $\Pi_{i}$ is an inter-arrival time, and the service times are i.i.d. with $E\left(V^{2 p_{2}}\right)<\infty \quad\left(p_{1}, p_{2}>1\right)$, then if $p=\min \left(p_{1}, p_{2}\right)$,

$$
\alpha\left(M_{n}, W\right) \leq C_{8}\left\{(\log n)^{p} / n^{\min (p-1, p / 2)}\right\}^{(2 p+1)^{-1}} \equiv C_{8} g(n, p)
$$

Proof. Heyde (1969) establishes $\psi(n)=O\left(g\left(n, p_{2}\right)\right)$ and Kennedy (1972b), Lemma 3.4, that $\varphi(n)=O\left(g\left(n, p_{1}\right)\right)$.

Note that by using Lemmas 2.2 and 2.3 of Kennedy (1972b), and employing the techniques of his Lemmas 3.1 and 3.4 and proceeding as in (7.8) here, we can establish $g\left(n, p_{1}\right)$ as the relevant rate of convergence for the arrival process without recourse to Kennedy's Skorokhod representations: we only
require the Skorokhod representation to give us the rate for
$\sum_{i=1}^{[n t]}\left(\Pi_{i}-\lambda\right) / \sigma \sqrt{n}$ (where $\quad \lambda=E\left\{\Pi_{i}\right\}, \quad \sigma^{2}=\operatorname{Var}\left\{\Pi_{i}\right\}$ ) as per Heyde (1969).
This is a conceptually simpler approach.

## CHAPTER 3

## FUNCTIONAL LIMIT LAWS FOR CLUSTER POINT PROCESSES

## 3.1。 Introduction

The cluster process has appeared extensively in the literature, a tribute to its practical interest as well as its theoretical accessibility. It has modelled many processes with some sort of regular "triggering" mechanism: early examples are contagion problems in ecology (Thompson (1955)), and the spatial distribution of galaxies (Neyman and Scott (1952), (1958)), followed by failure patterns in computers (Lewis (1964a), (1964b)), and the occurrence times and energies of earthquakes and aftershocks (VereJones (1970)). It has also been used to investigate "bunching" in traffic flow (Bartlett (1963)). Recently, Hawkes and Oakes (1974) have demonstrated its close alliance to the Hawkes process. Finally, Matthes and others (e.g. Matthes (1963b), Kerstan, Matthes and Mecke (1974)), Goldman (1967) and Lee (1967), (1968) have studied these processes in relation to infinitely divisible point processes, since Poisson cluster processes are equivalent to regular infinitely divisible point processes. Closely related models also exist, such as Neyman and Scott's (1964) branching type model for epidemics. As can be seen, multidimensional cluster processes (e.g. galaxies, epidemics) are clearly important, but in this thesis we confine ourselves to one dimension.

This chapter is concerned with functional limit laws for stationary cluster processes. Weak convergence has come into vogue particularly since the publication of Billingsley's (1968) text on the subject, although the theory in its present form has been available since 1956 (Prohorov (1956), Skorokhod (1956)). Applications of the theory of weak convergence have been too diverse to detail - for a review and references, see Iglehart (1974). We mention only that Jagers (1974) has a theory of weak convergence of
random measures and point processes on Polish spaces. The equivalence of functional central limit theorems for counting processes and the partial sum processes derived from inter-event times has been established in Iglehart and Whitt (1971) and Verwaat (1972). A summary of the theory of superpositions of point processes, including weak convergence problems, is given in Çinlar (1972). Whitt (1973) investigates rates of convergence of superposed processes to a Poisson process. Thinning of point processes is the subject of a paper by Jagers and Lindvall (1974) and also more generally in Kallenberg (1974). In these latter contexts (superpositions, thinning) in which a sequence of point processes converge to a point process, weak convergence and convergence of finite-dimensional distributions coincide, as was first pointed out by Straf (1972) and more generally by Whitt (1975). This equivalence has been generalized by Saunders (1975) to the space $M$ of finite non-negative measures on a complete $\sigma$-compact metric space, where $M$ is endowed with the topology of weak convergence.

The functional version of the law of the iterated logarithm for i.i.d. random variables was presented by Strassen in his famous paper in 1964. It has since been extended in various directions, for example in Heyde and Scott (1973) and Wichura (1973), and references there. Functional strong laws were introduced by Iglehart (197lb).

Section 2 of this chapter defines and examines existence of cluster processes with dependent clusters, as well as establishing notation needed to prove the limit theorems of Sections 3, 4 and 5. As with $G / G / \infty$ queue, we ask for limit theorems for the cluster process in terms of the same properties holding for the components. Also in Section 5, we investigate two processes related to cluster processes, by way of looking at limit laws for doubly stochastic Poisson processes, and the law of the iterated logarithm for the $G / G / \infty$ queue.

### 3.2. Definition and existence of cluster processes

A cluster point process $\eta^{*}$ (compare Daley (1972)) is generated by two independent components, the centre process $\eta$ consisting of points $\left\{t_{j}\right\}(j=0, \pm 1, \ldots)$, say, each of which initiates a subsidiary process $\eta_{j}$ which is a.s. finite and independent of $\eta$. The full process consists of the superposition of the $\eta_{j}^{\prime} s$, i.e.,

$$
\begin{equation*}
\eta^{*}(A)=\sum_{\text {all } j} \eta_{j}\left(A-t_{j}\right), \text { bounded } A \in B(R) \tag{2.1}
\end{equation*}
$$

However, as in the arrival process for the $G / G / \infty$ queue, we mainly refer to the centre process as a probability measure $P_{1}$ on $(N, \sigma(N))$. The subsidiary processes may be considered a probability measure $P_{2}$ on $\left(N^{2}, \sigma(N)^{Z}\right)$, so that again assuming that the centre process and process of subsidiaries are independent we form the probability triple $\left(N \times N^{Z}, \sigma(N) \times \sigma(N)^{Z}, P_{1} \times P_{2}\right)$.

Let $N_{\infty}$ denote the set of all non-negative integer or infinite valued measures on $R$ with $\sigma$-field generated by the sets $\left\{N \in N_{\infty}: N(A) \leq m\right\}$, $A \in B(R), \quad m \in Z_{+} \cup\{\infty\}$. This is the smallest $\sigma$-field such that the mappings $\phi_{A}: N_{\infty} \rightarrow[0, \infty]$ defined by $\phi_{A}(N)=N(A)$ are measurable, all $A \in B(R)$.

We will define our cluster process via the mapping $\eta_{c}=\phi_{3} \circ \phi_{2} \circ \phi_{1}$ from $N \times N^{2} \rightarrow N_{\infty}$ defined by

$$
\begin{equation*}
N \times N^{Z} \xrightarrow{\phi_{1}} R^{Z} \times N^{2} \xrightarrow{\phi_{2}} N^{2} \xrightarrow{\phi_{3}} N_{\infty} \tag{2.2}
\end{equation*}
$$

(1) $(N, N) \stackrel{\phi_{1}}{\longmapsto}\left(\left\{t_{j}(N)\right\}, N\right)$
(2) $\left(\left\{\beta_{j}\right\},\left\{N_{j}\right\}\right) \stackrel{\phi_{2}}{\longmapsto}\left\{N_{j}\left(\cdot-\beta_{j}\right)\right\}$
(3) $\left\{M_{j}\right\} \stackrel{\phi_{3}}{\longmapsto} \sum_{j=-\infty}^{+\infty} M_{j}$,
where the product spaces in (2.2) are equipped with the product $\sigma$-fields, and $\left\{t_{j}\right\}$ are defined as in (2.1) of Chapter l. It should be clear that (2.2) performs the same function as (2.1), namely, sums the contributions of each subsidiary process $N_{j}$ centred at $t_{j}$ to any Borel set $A$ to give the number of cluster points in $A$.

If $\eta_{c}$ is measurable, we can define a probability measure $P_{c}$ on $\left(N_{\infty}, \sigma\left(N_{\infty}\right)\right)$ by

$$
\begin{equation*}
P_{c}=\left(P_{1} \times P_{2}\right) n_{c}^{-1} \tag{2.3}
\end{equation*}
$$

If $P_{c}(N)=1$, then $P_{c}$ is referred to as the cluster point process. Since proving measurability of $\phi_{2}$ seems to be surprisingly difficult, we give a complete proof of

THEOREM 3.2.1. (a) The mapping $\eta_{c}$ defined around (2.2) is measurable.
(b) $\sigma\left(N_{\infty}\right) \cap N=\sigma(N)$.

Proof. (a) Measurability of $\phi_{1}$ follows from that of the sequence of maps $N \mapsto t_{j}(N), j \in Z$, which in turn follows from (2.2) of Chapter 1. Let $\phi_{A}: N_{\infty} \rightarrow[0, \infty]$ be defined by $\phi_{A}(N)=N(A), A \in B(R)$. To prove $\phi_{3}$ is measurable we have to show that $\phi_{A} \circ \phi_{3}$ is measurable for each Borel set. $A$. We also have the factorization $\phi_{A} \circ \phi_{3}=\psi \circ \psi_{A}$ where

$$
\begin{aligned}
& \left\{M_{j}\right\}_{j \in Z} \stackrel{\psi_{A}}{\longmapsto}\left\{M_{j}(A)\right\}_{j \in Z}, \text { i.e., } \psi_{A}: N^{Z} \rightarrow[0, \infty]^{Z}, \\
& \left\{\beta_{j}\right\}_{j \in Z} \stackrel{\psi}{\longmapsto} \sum_{j=-\infty}^{+\infty} \beta_{j} \quad \text { i.e., } \psi:[0, \infty]^{Z} \rightarrow[0, \infty] .
\end{aligned}
$$

We prove that both $\psi$ and $\psi_{A}$ are measurable. $[0, \infty]^{2}$ takes the product $\sigma$-field, so that measurability of $\psi_{A}$ follows immediately from
measurability of the map $M_{j} \mapsto M_{j}(A)$ from $N$ to $[0, \infty]$. Since
$\psi\left(\left\{\beta_{j}\right\}_{j \in Z}\right)=\lim _{n \uparrow \infty} \sum_{|j| \leq n} \beta_{j}$, and each of the maps $\left\{\beta_{j}\right\} \mapsto \sum_{|j| \leq n} \beta_{j}$ is continuous (when $[0, \infty]^{2}$ is given the product topology) and the product $\sigma$-field is precisely the Borel $\sigma$-field induced by the product topology, measurability follows.

Measurability of $\phi_{2}$ will follow from measurability of the coordinate map $\left(\left\{\beta_{j}\right\},\left\{N_{j}\right\}\right) \mapsto N_{j}\left(\cdot-\beta_{j}\right)$, and since this map only depends on $\left(\beta_{j}, N_{j}\right)$, from the measurability of $(\beta, N) \mapsto N(\cdot-\beta)$, or from the (joint) measurability of $(\beta, N) \mapsto N(A-B)$ for any $A \in B(R)$. For this purpose, define

$$
g(\beta, N)=\int h(x-\beta) d N(x),
$$

where $h(x) \in \mathcal{C}_{K}$, the class of continuous functions with compact support. Without loss of generality $N$ may be endowed with the vague topology (see (5.9) and (5.10) of Chapter 1). Now, for any $\varepsilon>0$, there exists a $\delta$ such that $\left|h(x-\beta)-h\left(x-\beta_{0}\right)\right|<\varepsilon$ for $\left|\beta-\beta_{0}\right|<\delta$. Let $K=\left\{x:|x-y| \leq \delta, y \in(\operatorname{supp} h)+\beta_{0}\right\}$. Then clearly $\left.\left|h(x-\beta)-h\left(x-\beta_{0}\right)\right|<\varepsilon\right]_{K}$ for $\left|\beta-\beta_{0}\right|<\delta$. So on the neighbourhood $\left\{N: N(K) \leq N_{0}(K)+1\right\}$ of $N_{0}$,

$$
\left|\int\left[h(x-\beta)-h\left(x-\beta_{0}\right)\right] N(d x)\right| \leq \varepsilon N(K) \leq \varepsilon\left[N_{0}(K)+1\right] .
$$

Also, as $\beta_{0}$ is fixed, $\left\{N:\left|\int h\left(x-\beta_{0}\right) N(d x)-\int h\left(x-\beta_{0}\right) N_{0}(d x)\right|<\varepsilon\right\}$ is a neighbourhood of $\dot{N}_{0}$ in the vague topology on $N$. Thus

$$
\left|g(\beta, N)-g\left(\beta_{0}, N_{0}\right)\right| \leq \varepsilon\left[N_{0}(K)+2\right]
$$

on the intersection of these two neighbourhoods $\times\left\{\beta:\left|\beta-\beta_{0}\right|<\delta\right\}$. Thus $g(\beta, N)$ is a jointly continuous function of $(\beta, N)$. According to Proposition 1.1 of Jagers (1974), $\sigma(N)=B(N)$, the Borel $\sigma$-field generated
by the open sets of the vague topology. Since the product $\sigma$-field on $R \times N$ is the Borel $\sigma$-field generated by the product topology ((usual topology on $R$ ) $\times$ (vague topology on $N)$ ), it follows that $g$ is also (jointly) measurable in this product $\sigma$-field.

Let $G$ be the class of bounded measurable functions $g: R \rightarrow R$ for which $\int f(x-\beta) g(x-\beta) d N(x)$ is a jointly measurable function of $(\beta, N)$ for every $f \in C_{K}$. Then every bounded continuous function is in $G$, since $f(x-\beta) g(x-\beta)$ is then continuous and has support a closed subset of a compact set, i.e., compact support, and $G$ is closed under uniformly bounded pointwise limits. Hence $G$ contains all bounded Borel measurable functions, and, in particular, $1_{A}$ for every Borel set $A$. Thus $\left.\int f(x-\beta)\right]_{A}(x-\beta) d N(x)$ is measurable. Taking $f_{n}(x-\beta) \uparrow 1$, we obtain by the monotone convergence theorem that $\int 1_{A}(x-\beta) N(d x)=N(A-\beta)$ is measurable.
(b) follows from Halmos (1950), Theorem E, p. 25.

The existence of the cluster process (i.e., the condition $P_{c}(N)=1$ ) is a vital question, but a non-stochastic sufficient condition may be given as follows:

THEOREM 3.2.2. If the subsidiary processes of a cluster process are uniformly bounded in the sense $P_{2}\left\{N: N_{i}(I)>0\right\} \leq F(I) \quad(a Z Z \quad i \in Z$, bounded intervals $I$ ) for some set function $F$, then the cluster process exists if

$$
\begin{equation*}
\int F(I-v) d M_{1}(v)<\infty \tag{2.4}
\end{equation*}
$$

for all bounded intervals $I$, provided the centre process has finite first moment measure $M_{1}$.

Proof。

$$
\begin{aligned}
P_{c}\{N & : N(I)<\infty\}=1 \\
& \Leftrightarrow \int_{N} P_{2}\left\{N: \sum_{j=-\infty}^{+\infty} N_{j}\left(I-t_{j}(N)\right)<\infty\right\} P_{l}(d N)=1
\end{aligned}
$$

by（2．3）and Fubini＇s Theorem，
$\Leftrightarrow P_{2}\left\{N: \sum_{j=-\infty}^{+\infty} N_{j}\left(I-t_{j}\right)<\infty\right\}=1, P_{1}$－a．s．
$\Leftrightarrow P_{2}\left\{N: N_{i}\left(I-t_{i}\right)>0\right.$ i．o．$\}=0, P_{1}$－a．s．， since $N_{i}(R)<\infty, P_{2}-$ a．s．，
$\Leftrightarrow \sum_{i=-\infty}^{+\infty} P_{2}\left\{N: N_{i}\left(I-t_{i}\right)>0\right\}<\infty, P_{1}-$ a．s．，
by the Borel－Cantelli Lemma，

$$
\begin{align*}
& \Leftrightarrow \sum_{i=-\infty}^{+\infty} F\left(I-t_{i}\right)<\infty, P_{1}-\text { a.s., }  \tag{2.5}\\
& \Leftrightarrow \int_{-\infty}^{+\infty} F(I-v) d N(v)<\infty, P_{1}-\text { a.s. }, \\
& \Leftrightarrow \int_{-\infty}^{+\infty} F(I-v) d M_{1}(v)<\infty .
\end{align*}
$$

It is clear that the reverse implications at（2．5）and the preceding step are also necessary if the clusters are i。i。d．（see also Westcott（1971）； in fact the condition immediately following（2．5）was already familiar to the German school（see e．g．Kestan，Matthes and Mecke（1974）；Chapters 5 and 6 and references there）：their derivation was via the Borel－Cantelli Lemma）． In fact，（2．4）is known to be necessary for a stationary Poisson cluster process with i．i．d．clusters（Matthes（1963b），Westcott（1971））．Westcott＇s techniques and Mathes＇formulation do not seem to readily extend to non－ independent subsidiaries．

Let us say that the subsidiary processes are stationary if with
$S_{k}: N^{2} \rightarrow N^{2}$ defined by $S_{k}\left\{N_{j}\right\}=\left\{N_{j+k}\right\}$,

## ADDENDUM (at examiner's request)

We will follow Daley (1972) in employing the decomposition of $\eta_{c}$ into a 'coherent' mapping $\eta^{+}$with remainder terms $\eta^{=}$and $\eta^{-}$. Specifically, we write

$$
\begin{align*}
\eta_{c}(0, x] & =\tilde{n}(0, x]+\eta^{=}(0, x] \\
& =\eta^{+}(0, x]+\eta^{=}(0, x]-\eta^{-}(0, x] \tag{2.7}
\end{align*}
$$

where

$$
\begin{align*}
& \eta^{+}(0, x] \equiv \sum_{j=1}^{N(0, x]} N_{j}(R), \\
& \tilde{n}(0, x] \equiv \sum_{j=1}^{N(0, x]} N_{j}\left(-t_{j}, x-t_{j}\right], \\
& \eta^{=}(0, x] \equiv\left[\sum_{j=-\infty}^{0}+\sum_{j=N(0, x]+1}^{\infty}\right) N_{j}\left(-t_{j}, x-t_{j}\right],  \tag{2.8}\\
& n^{-}(0, x] \equiv \sum_{j=1}^{N(0, x]} N_{j}\left(R \backslash\left(-t_{j}, x-t_{j}\right]\right) .
\end{align*}
$$

$$
\begin{equation*}
P_{2}\left(S_{k}^{-1} C\right)=P_{2}(C) \tag{2.6}
\end{equation*}
$$

for any $C \in \sigma(N)^{Z}$, all $k \in Z$. Then just as in Theorem 2.3.1 (a), we can prove

THEOREM 3.2.3. If the centre process and process of subsidiamies of a cluster process are both stationary, then the cluster process itself is stationary.

In the remainder of this chapter, we will require, for notational convenience, that the centre process be stationary with finite intensity $m$, and the process of subsidiaries be stationary with finite first moment, i.e., $\mu \equiv E_{2}\left\{N_{1}(R)\right\}<\infty$. Under these circumstances it is easy to prove, as in Daley (1972), equation (26), that $E_{c} N(0, x]=\mu m x$. We will write $m_{c}=\mu m$. Also it is clear that the cluster process exists in the sense of Theorem 3.2.2.

T We will follow Daley (1972) in employing the decompositions into a 'coherent' mapping $\eta^{+}$:

$$
\begin{equation*}
\eta_{c}(0, x]=\tilde{\eta}(0, x]+\eta^{=}(0, x], \tilde{n}(0, x]=\eta^{+}(0, x]-\eta^{-}(0, x] \tag{2.7}
\end{equation*}
$$

where

$$
\begin{equation*}
\tilde{n}(0, x]=\sum_{j=1}^{N(0, x]} N_{j}\left(-t_{j}, x-t_{j}\right], n^{+}(0, x]=\sum_{j=1}^{N(0, x]} N_{j}(R) . \tag{2.8}
\end{equation*}
$$

Again the device due to Iglehart and Kennedy (1970) will be very useful. Let $c_{n} \uparrow \infty(n \rightarrow \infty)$ : then, as in (2.4) of Chapter 2,
$\sup \eta^{-}(0, n t]$
$0 \leq t \leq s$

$$
\begin{equation*}
\leq \sum_{j=1}^{N(0, n s]} N_{j}\left(c_{n}, \infty\right)+\sum_{j=1}^{N(0, n s]} N_{j}\left(-\infty,-t_{j}\right]+\sup _{0 \leq t \leq s} n^{+}\left(n t-c_{n}, n t\right] . \tag{2.9}
\end{equation*}
$$

Similarly,
$\sup \eta^{=}(0, n t]$
$0 \leq t \leq s$

$$
\begin{equation*}
\leq \sum_{j=1}^{N(0, n s]} N_{j}\left(-\infty,-c_{n}\right]+\eta^{=}(0, n s]+\sup _{0 \leq t \leq s} \eta^{+}\left(n t, n t+c_{n}\right] \tag{2.10}
\end{equation*}
$$

### 3.3. A functional central limit theorem for cluster point processes

In their summary of the theory of point processes, Daley and Vere-Jones (1972) ask for a functional central limit theorem for cluster point processes. We show in Chapter 4 that $\phi$-mixing theorems, and some strong mixing theorems are inapplicable in this context, but in this section we demonstrate that a different proof may be readily assembled from known sources. The proof follows mainly from techniques employed by Iglehart and Kennedy (1970), and by modifying Daley's (1972) proof for one dimension. In the process we remove a second moment condition on the centre process, and find an alternative to a first moment condition on the subsidiaries, but invoke one extra constraint not required in Daley (1972), namely that the weak limit of the normed centre process should have a.s. continuous sample paths (i.e., is in $C\left[0,{ }^{\infty}\right)$; see Whitt (1970)). We shall use $\Rightarrow$ to denote weak convergence in $D[0, \infty)$, the space of functions on $[0, \infty)$ which are right continuous with left-hand limits. The technique is to prove the weak convergence theorem on $D[0, s]$, each $s>0$, and thus on $D[0, \infty)$ (Lindvall (1973), Whitt (1971)).

Before embarking on the main theorem itself, we will prove the following lemma, which is of independent interest. Let

$$
\begin{equation*}
L(t)=P_{2}\left\{N_{1}(-\infty, t]>0\right\}, R(t)=P_{2}\left\{N_{1}(t, \infty)>0\right\} \tag{3.1}
\end{equation*}
$$

LEMMA 3.3.1. If the centre process $P_{1}$ of finite intensity $m$ and process of subsidiaries $P_{2}$ are both stationary,

$$
\int_{0}^{\infty} L(-u) d u<\infty \Rightarrow \sum_{j=1}^{\infty} N_{j}\left(-\infty,-t_{j}\right]<\infty, P_{1} \times P_{2}-a . s .,
$$

$$
\int_{0}^{\infty} R(u) d u<\infty \Rightarrow \sum_{j=-\infty}^{0} N_{j}\left(-t_{j}, \infty\right)<\infty, P_{I} \times P_{2}-a_{0} s .
$$

Proof. The proof seems more natural in the reverse direction. $\Leftrightarrow$ in this proof means 'implies'。)

$$
\begin{aligned}
& \sum_{j=-\infty}^{0} N_{j}\left(-t_{j}, \infty\right)<\infty, P_{1} \times P_{2} \text {-a.s. } \\
& \quad \Leftrightarrow P_{2}\left\{N: N_{j}\left(-t_{j}, \infty\right)>0 \text { i.o. }, j \leq 0\right\}=0, P_{1}-\text { a.s. } \\
& \\
& \Leftrightarrow \sum_{j=-\infty}^{0} P_{2}\left\{N_{j}\left(-t_{j}, \infty\right)>0\right\}<\infty, P_{1}-\text { a.s. }
\end{aligned}
$$

by the Borel-Cantelli Lemma,

$$
\begin{aligned}
& \Leftrightarrow \int_{-\infty}^{0} P_{2}\left\{N_{1}(-v, \infty)>0\right\} d N(v)<\infty, P_{1} \text {-a.s. } \\
& \Leftrightarrow \int_{0}^{\infty} R(v) d v<\infty
\end{aligned}
$$

where we have taken expectations in the last step, and effected the change of variable $v^{\prime}=-v$. The other half of the proof is similar.

We aim for weak convergence to the Wiener process $W(\cdot)$ of the process

$$
\begin{equation*}
Z_{n}(t) \equiv\left(n_{c}(0, n t]-m_{c} n t\right) / \sqrt{n} \tag{3.2}
\end{equation*}
$$

for $t \in[0, \infty)$. We also require the processes $X_{n}: N \rightarrow D[0, \infty)$, $Y_{n}: N^{2} \rightarrow D[0, \infty)$ defined by

$$
\begin{align*}
& X_{n}(t) \equiv(N(0, n t]-m n t) / \sqrt{n},  \tag{3.3}\\
& Y_{n}(t) \equiv \sum_{j=1}^{[n t]}\left(N_{j}(R)-\mu\right) / \sqrt{n} . \tag{3.4}
\end{align*}
$$

THEOREM 3.3.2. If the centre process $P_{1}$ is stationary, and

$$
\begin{equation*}
P_{1} X_{n}^{-1} \Rightarrow \operatorname{Pr}_{r} X^{-1} \tag{3.5}
\end{equation*}
$$

where $\operatorname{Pr}\{X(\cdot) \in C[0, \infty)\}=1$, and the process of subsidiaries $P_{2}$ is stationary with the functions (3.1) satisfying

$$
\begin{equation*}
n^{2}\{L(-n)+R(n)\} \rightarrow 0 \quad(n \rightarrow \infty) \tag{3.6}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{2} Y_{n}^{-1} \Rightarrow P_{r} Y^{-1} \tag{3.7}
\end{equation*}
$$

then, writing $Z(t)=\mu X(t)+Y(m t)$,

$$
\begin{equation*}
\left(P_{1} \times P_{2}\right) Z_{n}^{-1} \Rightarrow P_{r} Z^{-1} \tag{3.8}
\end{equation*}
$$

and $X$ and $Y$ are independent.
Remark. We choose a scale such that $m<1$.
Proof. We use the decompositions (2.8), (2.9) and (2.10). For random functions $\theta_{n}, \theta$, we will abbreviate $\left(P_{1} \times P_{2}\right) \theta_{n}^{-1} \Rightarrow \operatorname{Pr}^{-1}$ as $\theta_{n} \Rightarrow \theta$. Then

$$
\begin{equation*}
\left(n^{+}(0, n t]-m_{c} n t\right) / \sqrt{n} \Rightarrow \mu X(t)+Y(m t) \tag{3.9}
\end{equation*}
$$

in the Skorokhod topology on $D[0, s]$, by Lemma 1 of Iglehart and Kennedy (1970). Let $\rho_{s}$ and $d_{s}$ respectively denote the uniform and Skorokhod metrics for $D[0, s]$ (cf. (5.11) and (5.12) of Chapter 1). To prove that the limit (3.8) occurs, it is sufficient to show that
$\left.\rho_{s}\left({\left(n^{-}\right.}^{(0, n \cdot]+\eta^{=}}(0, n \cdot]\right) / \sqrt{n}, 0\right) \Rightarrow 0$ by Theorem 4.1 of Billingsley (1968), since $\rho_{s} \geq d_{s}$. Provided $c_{n}=o(\sqrt{n}) \quad(n \rightarrow \infty)$, then

$$
\begin{equation*}
\sup _{0 \leq t \leq s} \eta^{+}\left(n t-c_{n}, n t+c_{n}\right] \Rightarrow 0 \tag{3.10}
\end{equation*}
$$

by Lemma 2 of Iglehart and Kennedy (1970). Also, for any $\varepsilon>0$,

$$
\begin{align*}
& P_{1} \times P_{2}\left\{\sum_{j=1}^{N(0, n s]} N_{j}\left(\left(-c_{n}, c_{n}\right]^{c}\right) \geq \varepsilon \sqrt{n}\right\}  \tag{3.11}\\
& \leq \int_{N} N(0, n s] P_{2}\left\{N_{1}\left(\left(-c_{n}, c_{n}\right]^{c}\right)>0\right\} d P_{1}(N) \\
& \leq n s\left\{L\left(-c_{n}\right)+R\left(c_{n}\right)\right\}=s\left\{n c_{n}^{-2}\right) c_{n}^{2}\left\{L\left(-c_{n}\right)+R\left(c_{n}\right)\right\}
\end{align*}
$$

so that if a sequence $d_{n} \uparrow \infty \quad(n \rightarrow \infty)$ exists such that $d_{n} m_{n} \rightarrow 0(n \rightarrow \infty)$, where $m_{n} \equiv n^{2}\{L(-n)+R(n)\}$, then (3.11) $\rightarrow 0(n \rightarrow \infty)$. Clearly $m_{n} \rightarrow 0$ $(n \rightarrow \infty)$ is necessary for this. Taking $d_{n}=m_{n}^{-\frac{1}{2}}$ shows that it is also sufficient.

Consider

$$
\eta^{=}(0, n s] \leq \sum_{j=-\infty}^{0} W_{j}\left(-t_{j}, \infty\right)+\sum_{j=N(0, n s]+1}^{\infty} N_{j}\left(-\infty, n s-t_{j}\right]
$$

But by stationarity of $P_{1} \times P_{2}$ with respect to the transformation $T_{n s}^{\prime}$ defined analogously to (3.4) of Chapter 2,

$$
\begin{equation*}
\sum_{j=N(0, n s]+1}^{\infty} N_{j}\left(-\infty, n s-t_{j}\right] \text { and } \sum_{j=1}^{\infty} N_{j}\left(-\infty,-t_{j}\right] \tag{3.12}
\end{equation*}
$$

have the same distribution. Also, the middle term of (2.9) has an upper bound of $\sum_{j=1}^{\infty} N_{j}\left(-\infty,-t_{j}\right]$, so that, since the conditions of Lemma 3.3.1 are satisfied, dividing by $\sqrt{n}$ means that it converges $P_{1} \times P_{2}$ - a.s. to zero as $n \rightarrow \infty$. Similarly $\eta^{=}(0, n s]$ converges to zero in $P_{1} \times P_{2}$-probability.

Unfortunately the above argument contains one step which depends critically upon stationarity (see Problem 6.4.1). If the clusters are only right-handed, this is unimportant, but a more robust technique is of course to estimate the small probabilities by Chebyshev's inequality. This yields

COROLLARY 3.3.3. Replacing (3.6) in Theorem 3.3.2 by the pair of conditions

$$
\begin{gather*}
n E_{2}\left\{N_{1}\left((-n, n]^{c}\right)\right\} \rightarrow 0 \quad(n \rightarrow \infty) \\
\frac{1}{\sqrt{n}} \int_{0}^{n} E_{2}\left\{N_{1}\left((-u, u]^{c}\right)\right\} d u \rightarrow 0 \quad(n \rightarrow \infty) \tag{3.13}
\end{gather*}
$$

yields the same conclusions under the same conditions.
The conditions (3.13) are predictably more severe than those of Daley (1972) for ordinary convergence, who used this technique. It should be clear though that we can find a weaker version of theorem 3.3.2 for ordinary convergence.

COROLLARY 3.3.4. If the centre process $P_{1}$ is stationary, and

$$
\begin{equation*}
X_{n}(1) \xrightarrow{D_{1}} X(1) \tag{3.14}
\end{equation*}
$$

and the process of subsidiaries $P_{2}$ is stationary with the functions (3.1) satisfying

$$
\begin{equation*}
\int_{0}^{\infty}\{L(-u)+R(u)\} d u<\infty \tag{3.15}
\end{equation*}
$$

and

$$
\begin{equation*}
Y_{n}(1) \xrightarrow{D_{2}} Y(1) \tag{3.16}
\end{equation*}
$$

then $X(1)$ and $Y(1)$ are independent and

$$
Z_{n}(1) \xrightarrow{D_{1} \times D_{2}} Z(1)
$$

In Theorem 3.3.2 and its Corollaries we have avoided the question of norming. Denote the norming constants for $X_{n}(\cdot), Y_{n}(\cdot)$ and $Z_{n}(\cdot)$ by $\sigma_{X}, \sigma_{Y}$ and $\sigma$ respectively; then clearly $\sigma^{2}=\mu^{2} \sigma_{X}^{2}+m \sigma_{Y}^{2}$. If the centre process is, say, $\phi$-mixing with finite first and second moments and totally finite reduced covariance measure $C(\cdot)$ (see Section 1.3), then $\sigma_{X}^{2}=\int_{R} C(d u)(\geq 0$, supposed $>0)$ and $X(\cdot)=\sigma_{X} W_{1}(\cdot)$, where $W_{1}(\cdot)$ is a standard Brownian motion. It is not difficult to prove, via equation Al of Vere-Jones (Appendix to Daley (1971)) that $\lim _{u \rightarrow \infty} \operatorname{Var}(\eta(0, u]) / u=\int_{R} C_{\eta}(d u)$ for any suitaple weakly stationary point process $\eta$. If, also, the clusters are i.i.d. with $\mu_{2}=E_{2}\left\{N_{1}(R)^{2}\right\}<\infty$, then $\sigma_{Y}^{2}=\left(\mu_{2}-\mu^{2}\right)$, and $Y(\cdot)=\sigma_{Y} W_{2}(\cdot)$ for a standard Brownian motion $W_{2}(\cdot)$. Then $Z_{n}(t) \Rightarrow W(t)$ for a standard Brownian motion $W(\cdot)$ by Theorem 3.3.2, and also $\eta_{c}(\cdot)$ has totally finite reduced covariance measure $C_{c}(\cdot)$ satisfying $\int_{R} C_{c}(d u)=\sigma^{2} \quad$ (Theorem 4 of Daley (1972)).
3.4. The law of the iterated logarithm for cluster point processes

### 3.4.1. SOME DEFINITIONS AND LEMMAS

The law of the iterated logarithm (LIL) has been studied extensively for many years. In the context of point processes, the law has been demonstrated for the split times of branching processes by Athreya and Karlin (1967), and in queueing situations by Iglehart (1971a). The equivalence of the LIL for the counting process and the process of interepoch times has been demonstrated by Verwat (1972), from which results for many point processes are immediate.

Firstly we require some lemmas which are mostly known. Denote by $D^{j}[0, \infty) \quad\left(D^{j}[0, s]\right)$ the $j$-fold product space of $D[0, \infty)(D[0, s])$. These spaces will be endowed with the product Skorokhod topology (Whitt (1971)). Generalizing Strassen (1964), we begin with

DEFINITION 3.4.1. Let $K^{j}\left(K_{s}^{j}\right)$ denote the set of absolutely continuous functions $x$ in $D^{j}[0, \infty)\left(D^{j}[0, s]\right)$ such that

$$
\begin{equation*}
x(0)=0 \text { and } \int_{0}^{\infty} \dot{x}(t)^{2} d t \leq 1\left(\int_{0}^{s} \dot{x}(t)^{2} d t \leq 1\right) \tag{4.1.1}
\end{equation*}
$$

where $\dot{x}$ denotes the derivative of $x$ determined almost everywhere with respect to Lepesgue measure, and the square is to be interpreted as inner product.

It follows immediately from the Schwarz inequality that for $x \in k^{j}$ and $0 \leq a \leq b$;

$$
\begin{equation*}
|x(b)-x(a)| \leq(b-a)^{\frac{1}{2}} \tag{4.1.2}
\end{equation*}
$$

In the next definition we state formally the meaning of a functional law of the iterated logarithm (FLIL).

DEFINITION 3.4.2. A process $\left\{X_{n}(t)\right\}_{n \geq 3}$ defined on a probability space $(\Omega, F, P r)$, whose sample paths lie $\operatorname{Pr}-$ a.s。in $D^{j}[0, \infty)$
$\left(D^{j}[0, s]\right)$ satisfies a FLIL if it is a.s. relatively compact in $D^{j}[0, \infty)$ $\left(D^{j}[0, s]\right)$ and has $K^{j}\left(K_{s}^{j}\right)$ as its set of limit points.

Let $\phi(n)=(2 n \log \log n)^{\frac{1}{2}}$. If $\left\{X_{i}\right\}$ is a sequence of zero mean, unit variance i.i.d. random variables, then Strassen's (1964) proof is easily adapted to show that $S_{n}(t)=\sum_{i=1}^{[n t]} X_{i} / \phi(n)$ satisfies a FLIL in $D[0, s]$, any $s>0$. Now we quote the continuous mapping theorem for the FLIL (Strassen (1964), Wichura (1973)) in the form we require.

LEMMA 3.4.3. Let $\left(\Omega, F, P_{r}\right)$ be a probability space and $S, S^{\prime}$ metric spaces with $K$ a compact subset of $S$. Let $f_{n}: \Omega \rightarrow S$ be mappings such that $\left\{f_{n}(\omega)\right\}, \omega \in \Omega$ is a.s. relatively compact in $S$ and has $K$ as its set of limit points. Let $\Phi: S \rightarrow S^{\prime}$ be a continuous mapping; then a.s. the sequence $\Phi\left(f_{n}(\omega)\right)$ is relatively compact in $S^{\prime}$ and its set of limit points is $\Phi(K)$.

We wish to relate the FLIL on $D^{j}[0, \infty)$ and $D^{j}[0, s], s>0$, to each other. This can be done in the same way as Whitt (1971) for weak convergence. Let the continuous mapping $r_{s}: D^{j}[0, \infty) \rightarrow D^{j}[0, s]$, $s>0$, be defined by $r_{s}(x)(t)=x(t), 0 \leq t \leq s$, for arbitrary $x \in D^{j}[0, \infty)$. Then we obtain

LEMMA 3.4.4. A random vector $\left\{X_{n}(t)\right\}$ in $D^{j}[0, \infty)$ satisfies a FLIL if and only if the random sequences $\left\{r_{s}\left(x_{n}\right)(t)\right\}$ in $D^{j}[0, s]$ satisfy a FLIL for each $s>0$.

Proof. The forward direction follows immediately from Lemma 3.4.3 and observing that $r_{s}\left(K^{\dot{j}}\right)=K_{s}^{j}$. Conversely, the relative compactness of $\left\{X_{n}(t)\right\}$ follows from that of $\left\{r_{s}\left(X_{n}\right)(t)\right\}$, all $s>0$, by an easy
extension of Theorem 2.3 of Whitt (1971) to $D^{j}[0, \infty)$. That the derived set is $K^{j}$ ensues easily from the form of $K_{s}^{j}$.

We require the following simple result in Section 3.4.2. Let $0<\lambda \leq 1, \alpha \geq 0, \beta \geq 0, \lambda \alpha^{2}+\beta^{2}=1$ and define $g: D^{2}[0, s] \rightarrow D[0, s]$ by

$$
\begin{equation*}
g(y, z)=\alpha y(\lambda \cdot)+\beta z(\cdot) . \tag{4.1.3}
\end{equation*}
$$

LEMMA 3.4.5. $g\left(K_{s}^{2}\right)=K_{s}^{1}$.
Proof. Let $x \in K_{s}^{l}$, and define $z(t) \in D[0, s]$ by $z(t)=\beta x(t)$. Also define $y(t) \in D[0, s]$ by

$$
\begin{aligned}
y(t) & =\alpha \lambda_{x}(t / \lambda), \quad 0 \leq t \leq \lambda s \\
& =\alpha \lambda x(s) \quad, \quad \lambda s \leq t \leq s .
\end{aligned}
$$

Then clearly $(y(0), z(0))=(0,0)$, and $(y, z)$ are absolutely continuous with respect to Lebesgue measure. Also

$$
\begin{aligned}
\int_{0}^{s}\left[\dot{y}(t)^{2}+\dot{z}(t)^{2}\right] d t & =\int_{0}^{\lambda s} \alpha^{2} \lambda^{2}\left[\left.\frac{d x(u)}{d u}\right|_{u=t / \lambda}\right]^{2} d t+\beta^{2} \int_{0}^{s} \dot{x}(t)^{2} d t \\
& =\left(\lambda \alpha^{2}+\beta^{2}\right) \int_{0}^{s} \dot{x}(t)^{2} d t \leq 1
\end{aligned}
$$

where we have made the change of variable $t^{\prime}=t / \lambda$ in the final step. Since $x(t)=\alpha y(\lambda t)+\beta z(t),(y, z) \in K_{s}^{2}, K_{s}^{1} \subset g\left(K_{s}^{2}\right)$.

Now to prove $g\left(K_{s}^{2}\right) \subset K_{s}^{1}$. Let $(y, z) \in K_{s}^{2}$. Clearly $g(y, z)(0)=0$,
and $g(y, z)$ is absolutely continuous with respect to Lebesgue measure.

$$
\begin{aligned}
\int_{0}^{s}\{\alpha & \left.\frac{d y(\lambda t)}{d t}+\beta \dot{z}(t)\right\}^{2} d t \\
& \left.=\int_{0}^{s}\left(\alpha^{2}+\beta^{2} \lambda^{-1}\right)\left[\frac{d y(\lambda t)}{d t}\right]^{2}+\left(\lambda \alpha^{2}+\beta^{2}\right) \dot{z}(t)^{2}-\left\{\alpha \sqrt{\lambda z}(t)-\beta \lambda^{-\frac{1}{2}} \frac{d y(\lambda t)}{d t}\right\}^{2} d t \right\rvert\, \\
& \leq\left(\alpha^{2}+\beta^{2} \lambda^{-1}\right) \lambda \int_{0}^{\lambda s} \dot{y}(t)^{2} d t+\int_{0}^{s} \dot{z}(t)^{2} d t \\
& \leq \int_{0}^{s}\left(\dot{y}(t)^{2}+\dot{z}(t)^{2}\right) d t \leq 1 .
\end{aligned}
$$

For $x, y \in D[0, s], s>0$, define the supremum metric by

$$
\begin{equation*}
\rho_{s}(x, y) \equiv \sup _{0 \leq t \leq s}|x(t)-y(t)| \tag{4.1.4}
\end{equation*}
$$

and the modulus of continuity of $x$ by

$$
\begin{equation*}
w_{s}(x, \delta) \equiv \sup _{\substack{0 \leq t, u \leq s \\|t-u| \leq \delta}}|x(t)-x(u)| \tag{4.1.5}
\end{equation*}
$$

In Lemma 3.4.7 and Theorem 3.4.8 we will need
LEMMA 3.4.6. If a process $\left\{X_{n}(t)\right\}$ on $D[0, \infty)$ satisfies a FLIL, then for any $s>0$,
(a) $\rho_{s}\left(X_{n}, K_{s}^{1}\right) \rightarrow 0$ a.s. $(n \rightarrow \infty)$,
(b) $\lim _{\delta \downarrow 0} \lim _{n \rightarrow \infty} \sup _{s}\left(X_{n}, \delta\right)=0$ a.s.

Proof. (a) Since $K_{s}^{1}$ is the derived set in the Skorokhod topology (metric $d_{s}$ ) on $D[0, s], d_{s}\left(X_{n}, K_{s}^{l}\right) \rightarrow 0$ as. Hence for each $\omega \in \Omega$, $\varepsilon>0$, there exists an $n_{0}(\omega, \varepsilon)$ such that for each $n \geq n_{0}(\omega, \varepsilon)$ there is $z_{n}(\omega, \varepsilon) \in K_{s}^{1}$ satisfying

$$
\begin{equation*}
d_{s}\left(X_{n}, z_{n}\right)<\varepsilon . \tag{4.1.6}
\end{equation*}
$$

Now if $\Lambda_{s}$ is the set of increasing invertible maps $[0, s] \rightarrow[0, s]$

$$
\left|X_{n}(t)-z_{n}(t)\right| \leq\left|X_{n}(t)-z_{n}\left(\lambda_{n}(t)\right)\right|+\left|z_{n}\left(\lambda_{n}(t)\right)-z_{n}(t)\right|
$$

for any $\lambda_{n} \in \Lambda$. Since $z_{n} \in K_{s}^{1},\left|z_{n}(a)-z_{n}(b)\right| \leq(a-b)^{\frac{1}{2}}$, and by
(4.1.6), $\sup _{t}\left|\lambda_{n}(t)-t\right| \leq \varepsilon$ for some $\lambda_{n}$. Hence $\rho_{s}\left(X_{n}, K_{s}^{1}\right) \leq \varepsilon+\varepsilon^{\frac{1}{2}}$, i.e., (a) holds.
(b) can then be proved as in Iglehart (1971a), Theorem 3.3 (see Problem 6.4.2). For completeness, we give it here. Clearly for $z \in K_{s}^{\perp}$,
$\omega_{s}(z, \delta) \leq \delta^{\frac{3}{2}}$. But for $\varepsilon>0, \omega \in \Omega, n \geq n_{0}(\omega, \varepsilon)$, there exists (as in (a)) a $z_{n} \in K_{s}^{I}$ such that

$$
w_{s}\left(X_{n}, \delta\right) \leq w_{s}\left(z_{n}, \delta\right)+2 \rho_{s}\left(X_{n}, z_{n}\right) \leq \delta^{\frac{3}{2}}+2 \varepsilon
$$

using (5.13) of Chapter 1 , and the conclusion follows.
This proof depends strongly on the properties of $K_{S}$. It is not simply a result of the relative compactness of processes satisfying the FLIL (i.e., of the Arzelà-Ascoli Theorem (Billingsley (1968), p. 211)), although for processes 'sufficiently like those in $C[0, s]$ ' it is (since the Skorokhod topology relativized to $C$ coincides with the uniform topology there). For instance if $S_{n}(t)=\sum_{j=1}^{[n t]} X_{j} / \sigma \phi(n)$ satisfies a FLIL, where $\left\{X_{j}\right\}$ is stationary, $E X_{1}=0$ and $\sigma$ is some positive constant, and if we associate with $S_{n}(\cdot)$ its linear approximation

$$
S_{n}^{\prime}(t)=\left[\sum_{j=1}^{[n t]} X_{j}+(n t-k) X_{k+1}\right] / \sigma \phi(n), \quad k s \leq n t \leq k s+1, \quad 0 \leq k \leq n s-1,
$$

then

$$
\begin{equation*}
\rho_{s}\left(S_{n}, S_{n}^{\prime}\right) \leq \max _{1 \leq k \leq n s}\left|X_{k}\right| / \sigma \phi(n) \tag{4.1.7}
\end{equation*}
$$

Clearly the RHS of $(4.1 .7) \rightarrow 0$ a.s. if $\left|X_{n}\right| / \phi(n) \rightarrow 0$ a.s. which follows from the Borel-Cantelli Lemma if $E X_{1}^{2}<\infty$ (see Problem 6.4.3). This type of argument may also be applied to $\eta_{n}(t)=(n(0, n t]-m n t) / \sigma \phi(n)$ for a suitable point process $\eta$. The point is, though, that if the derived set is not $K^{l}$ (e.g., Theorem 4.2 of Iglehart (197la)) for some process $X_{n}(t)$ which is relatively compact in $D$, and $X_{n}$ is not 'sufficiently like
a process in $C[0, s]^{\prime}$, Lemma 3.4.6 will have to be reproved before it can be used.

Finally we require a 'random change of time' lemma: this result has been observed before by Freedman (1967) and Iglehart (1971b), although not quite in this setting.

LEMMA 3.4.7. Let $\left\{X_{n}(t)\right\}_{n \geq 3}$ and $\left\{N_{t}\right\}(t \geq 0)$ be processes defined on a space $(\Omega, F, P r)$ such that $N_{t} \geq 0, N_{t} \uparrow \infty(t \rightarrow \infty)$ and $N_{t} / t \rightarrow m, 0<m<1$ a.s. $(t \rightarrow \infty)$; if $\left\{X_{n}(\cdot)\right\}$ satisfies a FLIL, then upon setting

$$
\begin{aligned}
\Phi_{n}(t) & \equiv N_{n t} / n \\
\rho_{s}\left(X_{n}\left(\Phi_{n}(\cdot)\right), X_{n}(m \cdot)\right) & \rightarrow 0 \quad \text { a.s., each } s>0
\end{aligned}
$$

Remark. This of course means that $X_{n}\left(\Phi_{n}(\cdot)\right)$ and $X_{n}\left(m^{\bullet}\right)$ obey the same FLIL on $D[0, s]$, any $s>0$, and hence on $D[0, \infty$ ) by Lemma 3.4.4.

Proof. The process $N_{t}$ obeys a functional strong law of large numbers (as per Iglehart (1971b), Theorem 3.1); hence, a.s. for each $\omega \in \Omega$, $\delta>0$, there exists an $n_{0} \equiv n_{0}(\omega, \delta)$ such that for $n \geq n_{0}$,

$$
\rho_{s}\left(\Phi_{n}(\cdot), m \cdot\right) \leq \delta
$$

Hence, choosing $\delta \leq s(1-m)$, for $n \geq n_{0}$,

$$
\rho_{s}\left(X_{n}\left(\Phi_{n}(\cdot)\right), X_{n}(m \cdot)\right) \leq w_{s}\left(X_{n}, \delta\right)
$$

and the conclusion follows from Lemma 3.4.6 (b).

### 3.4.2. THE LIL FOR PROCESSES WITH RIGHT HAND CLUSTERS

We are now in a position to prove our theorem for cluster processes wit right-hand clusters. We must suppose the centre process and process of subsidiaries to be defined on the same space $N \times N^{Z}$.

Once again we suppose our processes stationary. Let

$$
\begin{equation*}
X_{n}(t) \equiv(N(0, n t]-m n t) / \sigma_{1} \phi(n), \tag{4.2.1}
\end{equation*}
$$

$$
\begin{equation*}
Y_{n}(t) \equiv \sum_{j=1}^{[n t]}\left(N_{j}(R)-\mu\right) / \sigma_{2} \phi(n) \tag{4.2.2}
\end{equation*}
$$

for some appropriate constants $\sigma_{1}, \sigma_{2}$, where $\phi(n)=(2 n \log \log n)^{\frac{7}{2}}$.
Then setting $\sigma^{2}=\mu^{2} \sigma_{1}^{2}+m \sigma_{2}^{2}$, we define

$$
\begin{equation*}
Z_{n}(t) \equiv\left(n_{c}(0, n t]-m_{c} n t\right) / \sigma \phi(n) \tag{4.2.3}
\end{equation*}
$$

THEOREM 3.4.8. If $\left(X_{n}, Y_{n}\right)$ jointly satisfy a FLIL, i.e., if $\left(X_{n}, Y_{n}\right)$ is relatively compact in $D^{2}[0, \infty) P_{1} \times P_{2}-$ a.s. and has $K^{2}$ as its set of limit points, and if either

$$
\begin{equation*}
\int_{e}^{\infty} E_{2}\left\{N_{1}(u, \infty)\right\} /(\log \log u)^{\frac{3}{2}} d u<\infty \tag{4.2.4}
\end{equation*}
$$

or, if, in the notation of (3.1),

$$
\begin{equation*}
\int_{0}^{\infty} u R(u) d u<\infty \tag{4.2.5}
\end{equation*}
$$

then $P_{1} \times P_{2}-$ a.s. $Z_{n}$ is relatively compact in $D[0, \infty)$ with limit set $K^{1}$.

Remarks. The condition on $\left(X_{n}, Y_{n}\right)$ and (4.2.4) and (4.2.5) will be discussed in Section 3.4.3. We will take $m<1$, a condition easily achieved by scaling (see Problem 6.4.4).

Proof. Clearly (recalling that $m_{c}=\mu m$ ),
$Z_{n}(t)=\sigma^{-1} \mu \sigma_{1} X_{n}(t)+\sigma^{-1} \sigma_{2} Y_{n}(N(0, n t] / n)$

$$
+\left\{\sum_{j=1}^{N(0, n t]} N_{j}\left(n t-t_{j}, \infty\right) / \sigma \phi(n)\right\}+\left\{\sum_{j=-\infty}^{0} N_{j}\left(-t_{j}, n t-t_{n}\right] / \sigma \phi(n)\right\} .
$$

Define a continuous mapping $g_{s}: D^{2}[0, s] \rightarrow D[0, s]$ by

$$
\begin{equation*}
g_{s}(x, y)=\left(\sigma_{1} \mu x(\cdot)+\sigma_{2} y(m \cdot)\right) / \sigma \tag{4.2.7}
\end{equation*}
$$

Expressing the first two terms of $(4.2 .6)$ in terms of $g_{s}$, we find
$\rho_{s}\left(g_{s}\left[X_{n}(\cdot), Y_{n}(N(0, n \cdot] / m n)\right], g_{s}\left(X_{n}, Y_{n}\right)\right)$

$$
=\sigma^{-1} \sigma_{2} \rho_{s}\left(Y_{n}(N(0, n \cdot] / n), Y_{n}(m \cdot)\right) \rightarrow 0, P_{1} \times P_{2}-a_{0} s
$$

by Lemma 3.4.7. Hence the FLIL for the first two terms of (4.2.6) is the same as for $g_{s}\left(X_{n}, Y_{n}\right)$, which, by Lemma 3.4 .3 , is relatively compact in $D[0, s]$ and has $g_{s}\left(K_{s}^{2}\right)=K_{s}^{1}$ (Lemma 3.4.5) as its set of limit points。 Then, defining $g: D^{2}[0, \infty) \rightarrow D[0, \infty)$ as in (4.2.7), $g\left(X_{n}, Y_{n}\right)$ is almost surely relatively compact in $D\left[0, \infty\right.$ ) with limit set $K^{1}$ (Lemma 3.4.4).

We now prove that the condition (4.2.4) guarantees that the remainder terms in (4.2.6) converge almost surely (i.e., $\rho_{s}$ (remainder terms, 0 ) $\rightarrow 0$, $P_{1} \times P_{2}$ - a.s.) to zero.

For any arbitrary fixed $s>0$, define random functions in $D[0, s]$ by

$$
\begin{aligned}
& \gamma_{n}(t) \equiv\left(n^{+}(0, n t]-m_{c} n t\right) / \sigma \phi(n) \\
& \theta_{n}(t) \equiv \sum_{j=1}^{N(0, n t]} N_{j}\left(c_{m n s}, \infty\right) / \phi(n),
\end{aligned}
$$

where $c_{n} \uparrow \infty(n \rightarrow \infty)$. Note that $\gamma_{n}(\cdot)$ represents the first two terms of (4.2.6). Now
$\rho_{s}\left(\sum_{j=1}^{N(0, n t]} N_{j}\left(n t-t_{j}, \infty\right) / \phi(n), 0\right)$

$$
\begin{align*}
& \leq \theta_{n}(s)+\rho_{s}\left(\eta^{+}\left(n t-c_{m n s}, n t\right], 0\right) \\
& \leq \theta_{n}(s)+\sigma \rho_{s}\left(\gamma_{n}(t), \gamma_{n}\left(t-c_{m n s} / n\right)\right)+\mu c_{m n s} / \phi(n) . \tag{4.2.8}
\end{align*}
$$

We take $c_{n}=o(\phi(n))$, so that the final term of (4.2.8) converges to zero. The second term of (4.2.8) has an upper bound of $\sigma w_{s}\left(\gamma_{n}, c_{m n s} / n\right)$, which, by Lemma 3.4.6, converges to zero almost surely since we have proved $\gamma_{n}(t)$ satisfies a FLIL. We now tackle the first term. This term can be handled using a geometric subsequence argument, but it seems worthwhile to
record the following alternative argument, which yields exactly the same condition. Let

$$
\theta_{n}^{\prime}(t) \equiv \sum_{j=1}^{[m n t]} N_{j}\left(c_{m n s}, \infty\right) / \phi(n)
$$

Since $X_{n}$ satisfies a FLIL, there exists, for any $\delta>0$ and for each $N \in N\left(P_{1}\right.$-a.s. $)$, a $n_{0}(N, \delta)$ such that for $n \geq n_{0}(N, \delta)$,

$$
[n(m t-\delta)] \leq N(0, n t] \leq[n(m t+\delta)], 0 \leq t \leq s
$$

Hence, for $n \geq n_{0}(N, \delta)$,

$$
\rho_{s}\left(\theta_{n}, \theta_{n}^{\prime}\right) \leq w_{s}\left(\theta_{n}^{\prime}, \delta\right) \leq 2 \rho_{s}\left(\theta_{n}^{\prime}, 0\right)=2 \theta_{n}^{\prime}(s)
$$

and $\theta_{n}(s) \rightarrow 0$ a.s. if $\theta_{n}^{\prime}(s) \rightarrow 0$ a.s. But

$$
\begin{equation*}
\theta_{n}^{\prime}(s) \leq\left\{\sum_{j=1}^{[m n s]} N_{j}\left(c_{j}, \infty\right) / \phi([m n s])\right\} \cdot\{\phi([m n s]) / \phi(n)\} \tag{4.2.9}
\end{equation*}
$$

and the RHS of (4.2.9) converges to zero if the first factor converges to zero, i.e., if $\sum_{j=1}^{n} N_{j}\left(c_{j}, \infty\right) / \phi(j) \rightarrow$ a.s. by Kronecker's Lemma. But by monotonicity, we only require a subsequence to converge, i.e.,

$$
\begin{equation*}
\sum_{j=1}^{N(0, n]} N_{j}\left(c_{j}, \infty\right) / \phi(j) \rightarrow P_{1} \times P_{2}-\text { a.s. } \tag{4.2.10}
\end{equation*}
$$

We require $c_{j}=o(\phi(j))$, so we take $c_{j}=m^{-\frac{1}{2}}(j \log \log \log j)^{\frac{7}{2}}$; sincethe LHS of (4.2.10) will converge if it has an asymptotic upper bound wh converges, and since $t_{j} / j \rightarrow m^{-1}$ a.s. (cf. Proof of Theorem 2.4.1 (b)), we require only the convergence of

$$
\begin{equation*}
\sum_{j=1}^{N(0, n]} N_{j}\left(\sqrt{t_{j}}, \infty\right) / \phi\left(t_{j}\right) \tag{4.2.11}
\end{equation*}
$$

By the sub-martingale convergence theorem (Breiman (1968), p. 89), a sufficient condition is

$$
\underset{n \rightarrow \infty}{\lim \sup } \int_{e}^{n} E_{2}\left\{N_{1}(\sqrt{u}, \infty)\right\} /(u \log \log u)^{\frac{1}{2}} d u<\infty
$$

Changing variables $u^{\prime}=\sqrt{u}$ gives (4.2.4). Note that a finer condition can be obtained by refining the $\sqrt{t}$ factor in (4.2.11), but since the resultant expression is untidy, and there is a grosser approximation in the argument (see Problem 6.4.5), we have not given it.

The final term of (4.2.6) may be dealt with by a standard geometric sybsequence argument, which we will not elaborate on. It yields the condition (use Chebyshev's inequality, the Borel-Cantelli Lemma, bound sums by integrals, exchange integrals of. proof of Theorem 2.4.1 (b))

$$
\int_{0}^{\infty} E_{2}\left\{N_{1}(u, \infty)\right\} / \sqrt{u} d u<\infty
$$

which is weaker than (4.2.4).
To obtain the condition (4.2.5), we now employ a geometric subsequence argument on $\theta_{n}(s)$. Let $\alpha>1$, and consider $r_{k}=\left[\alpha^{2 k}\right], k \in Z_{+}$. For an arbitrary fixed $\varepsilon>0$, writing $d(n)=c_{m n s}$, define
$B_{k}=\left\{\sum_{j=1}^{N\left(0, r_{k} t\right]} N_{j}\left(d\left(r_{k-1}\right), \infty\right)>\varepsilon \phi\left(r_{k-1}\right)\right\}$. Then clearly
$P_{1} \times P_{2}\left\{\theta_{n}(s)>\varepsilon \phi(n)\right.$ i.o. $\} \leq P_{1} \times P_{2}\left\{B_{k}\right.$ i.o. $\}$, since, if for some $Z$, $\theta_{n}(Z)>\varepsilon \phi(Z)$, then, setting $r_{k}(Z)$ to be the next term in the geometric subsequence greater than $Z, B_{k}(Z)$ holds (in an obvious notation). We require $P_{1} \times P_{2}\left\{B_{k}\right.$ i.o. $\}$ to be zero, which occurs if $\sum P_{1} \times P_{2}\left\{B_{k}\right\}<\infty$ by the Borel-Cantelli Lemma. Arguing as in (3.11), and choosing $d(n)=\sqrt{n}$ (we require $d(n)=o(\phi(n)$ ), and choosing finer $d(n)$ 's leads to intractab expressions), this occurs if

$$
\begin{aligned}
\sum_{k=3}^{\infty} r_{k} R\left(r_{k-1}\right) & \leq \int_{3}^{\infty} \alpha^{2 \theta} R\left(\alpha^{\theta-3}\right) d \theta \\
& \leq\left(\alpha^{6} / \log \alpha\right) \int_{1}^{\infty} u R(u) d u<\infty
\end{aligned}
$$

where we change variables $u=\alpha^{\theta-3}$ in the last step.

The final term of (4.2.6) converges to zero a.s. under the condition (4.2.5) by Lemma 3.3.1.

COROLLARY 3.4.9 (Ordinary Law of the Iterated Logarithm). Under the conditions of Theorem 3.4.8,

$$
\lim _{n \rightarrow \infty} \sup _{n} Z_{n}(1)=+1, \quad \lim _{n \rightarrow \infty} \inf _{n} Z_{n}(1)=-1
$$

Proof. As indicated by Strassen (1964).

### 3.4.3. COMMENTS

(a) The LIL for processes with double-sided clusters

If the clusters can occur on the left-hand side of their initiating points, we obtain more remainder terms which we must prove converge to zero almost surely. Most of these can be handled as before, but

$$
\begin{equation*}
\sum_{j=N(0, n]+1}^{\infty} N_{j}\left(-t_{j}, n-t_{j}\right] / \phi(n) \tag{4.3.1}
\end{equation*}
$$

proves difficult. If the left-hand clusters have bound $d$, then $\infty$ in (4.3.1) may be replaced by $N(0, n+d]$, and this converges to zero a.s. as does $\gamma_{n}(t)$ in (4.2.8). In the general case, if we replace $N(0, n]+1$ in (4.3.1) by $N\left(0, \alpha^{2} n\right]+1$ for some $\alpha>1$, then we can prove convergence using the standard geometric subsequence argument and Chebyshev's inequality under the condition

$$
\int_{1}^{\infty} E_{2}\left\{N_{1}(-\infty,-u]\right\} / \sqrt{u} d u<\infty
$$

Conditions ensuring that the remainder (the sum from $N(0, n]+1$ to $N\left(0, \alpha^{2} n\right]$ ) converge to zero a.s. seem elusive. The reasoning around (3.12) has also defied extension to the a.s. convergence case.
(b) The conditions in Theorem 3.4.8.

It is clear that if the clusters have only one member each, then (4.2.4) is much weaker than (4.2.5). On the other hand, it is easy to construct examples in which (4.2.5) is satisfied but (4.2.4) not, but mostly
these involve $N_{1}(R)$ not having a first moment. The following example indicates that under fairly reasonable circumstances (4.2.4) may be much weaker. A general comparison does not seem feasible.

The cluster members of a Bartlett-Lewis process form a finite renewal process of length $S$, where the inter-point times are i.i.d. with common distribution $F$ and independent of $S$. Let
$R_{\boldsymbol{s}}(r)=\sum_{j=r+1}^{\infty} \operatorname{Pr}\{S=j\}$. Then, using Lawrance (1972), Equation (4.3.7),
(4.2.4) becomes

$$
\begin{equation*}
\int_{e}^{\infty} \sum_{i=0}^{\infty} R_{s}(i)\left(1-F^{(i+1)^{*}}(u)\right) /(\log \log u)^{\frac{1}{2}} d u<\infty \tag{4.3.2}
\end{equation*}
$$

Suppose, for example, that $F$ has regularly varying tails with exponent $\alpha>0$ (see Feller (1966), p. 268), i.e., as $u \rightarrow \infty$,

$$
1-F(u) \sim u^{-\alpha \hat{L}(u)}
$$

where $\hat{L}(s) \geq 0$ is of slow variation. Then, using the Corollary, p. 272 of Feller (1966), and provided $E\left(S^{2}\right)<\infty,(4.3 .2)$ reduces to

$$
\int_{e}^{\infty} \hat{L}(u) /\left(u^{\alpha}(\log \log u)^{\frac{1}{2}}\right) d u<\infty
$$

so that even if $\hat{L}(u)$ is bounded, we still require $\alpha>1$, a rather strong condition. However, $R(u)=\sum_{j=0}^{\infty} \operatorname{Pr}\{S=j\}\left(1-F^{j^{*}}(u)\right)$, so that (4.2.5) will require $\alpha>2$, a much stronger condition.

The assumption that $X_{n}$ and $Y_{n}$ jointly satisfy a FLIL warrants some attention. Certainly this is true if the centre process is a stationary renewal process (or the superposition of stationary renewal processes) whose second moment exists, and the clusters are i.i.d.: for, if $W_{1}$ and $W_{2}$ are independent Brownian motions on a common space $\Omega$, and we redefine $X_{n}$ and $Y_{n}$ also to be on $\Omega$, we can show $\rho_{s}\left(X_{n}, W_{n}^{I}\right) \rightarrow 0$ a.s. (Iglehart
(197la), Section 2) and $\rho_{s}\left(Y_{n}, W_{n}^{2}\right) \rightarrow 0$ a.s. (Strassen (1964)), where $W_{n}^{i}(t)=W_{i}(n t) / \phi(n), \quad i=1,2$. Hence if $\rho_{s}^{2}$ is the product supremum metric for $C^{2}[0, s]$ or $D^{2}[0, s]$,

$$
\begin{equation*}
\rho_{s}^{2}\left(\left(X_{n}, Y_{n}\right),\left(W_{n}^{1}, W_{n}^{2}\right)\right) \rightarrow 0 \quad \text { a.s. } \tag{4.3.3}
\end{equation*}
$$

The result then follows from Strassen (1964), Theorem l, who employed the Skorokhod representation theorem approach. Alternatively, if the centre process $\eta$ and process of subsidiaries $\eta$ possess some form of asymptotic independence, say $\phi$-mixing, and given this, satisfy the conditions of Corollary 3 of Heyde and Scott (1973), then again the above Skorokhod type approach, as used by Heyde and Scott, will prove that the joint FLIL follows merely from the individual FLIL's and the independence of $X_{n}$ and $Y_{n}$. Alternatively, for $\eta$ and $\eta \quad \phi$-mixing or strongly mixing, we may possibly use Chover's (1967) approach as applied by Oodaira and Yoshihara (197la, b) to demonstrate this. It is intuitively clear, however, that the joint FLIL will not necessarily follow from merely the individual FLIL's and independenc we need some statement such as (4.3.3) on the 'density' of subsequences converging to particular points in the limit set. Indeed, in order to satisfactorily achieve this, we may need to define the LIL for processes via an appropriate generalization of (4.3.3).

### 3.5. Related topics

### 3.5.1. FUNCTIONAL STRONG LAWS FOR CLUSTER POINT PROCESSES

The strong law of large numbers for cluster processes has been investigated by Daley (1972). Here we strengthen his theorem, as well as generalize it to the functional case.

Let $T_{y}: N \rightarrow N, S_{k}: N^{Z} \rightarrow N^{Z}$ and $T_{y}^{\prime}: N \times N^{Z} \rightarrow N \times N^{Z}$ be defined analogously to (3.1), (3.2) and (3.4) of Chapter 2, and let $T_{1}, T_{2}$ and $T$
denote their respective invariant $\sigma$-fields. We will assume $P_{1}$ and $P_{2}$ stationary (Theorem 3.2.3).

THEOREM 3.5.1. If the centre process $P_{1}$ is stationary with $E_{1} N(0,1]<\infty$, and the process of subsidiaries $P_{2}$ is stationary with $E_{2} N_{1}(R)<\infty$, then

$$
\lim _{n \rightarrow \infty} \eta_{c}(0, n] / n=E_{1} \times E_{2}\left\{\eta_{c}(0,1] \mid T\right\}, P_{1} \times P_{2}-a . s .,
$$

where $\left(E_{1} \times E_{2}\right) \eta_{c}(0,1]=E_{1}\{N(0,1]\} E_{2}\left\{N_{1}(R)\right\}<\infty$, and

$$
\begin{equation*}
E_{1} \times E_{2}\left\{\eta_{c}(0,1] \mid T\right\}=E_{1}\left\{N(0,1] \mid T_{1}\right\} E_{2}\left\{N_{1}(R) \mid T_{2}\right\} \tag{5.1}
\end{equation*}
$$

Remark. The LHS and RHS of (5.1) are random variables on $N \times N^{Z}$ (see Problem 6.4.6) .

Proof. Omitted. Similar to the proofs of Theorem 5, Daley (1972) and Theorem 2.4.1 (a).

COROLLARY 3.5.2. If the conditions of Theorem 3.5.1 are satisfied, and $\{N(0, n]: n \geq 1\}$ and $\left\{N_{i}(R): i \geq 1\right\}$ both satisfy the strong law of large numbers $P_{1}-a . s$. and $P_{2}-a . s$. respectively, then as $n \rightarrow \infty$,

$$
\eta_{c}(0, n] / n \rightarrow E_{1}\{N(0,1]\} E_{2}\left\{N_{1}(R)\right\}<\infty, P_{1} \times P_{2}-a . s
$$

COROLLARY 3.5.3. Let $L_{s}(t)=E_{1}\left\{N(0,1] \mid T_{1}\right\} E_{2}\left\{N_{1}(R) \mid T_{2}\right\} t$, $t \in[0, s]$. Then if $\rho_{s}$ is the supremum metric on $D[0, s]$, and the conditions of Theorem 3.5.1 are satisfied,

$$
\rho_{s}\left(\eta_{c}(0, n \cdot] / n, L_{s}\right) \rightarrow 0, P_{I} \times P_{2}-a . s
$$

Proof. Repeat the argument in Iglehart (1971b), Theorem 3.1, for each $(N, N) \in N \times N^{Z}$ for which convergence of $\eta_{c}(0, n] / n$ holds.
3.5.2. THE LAW OF THE ITERATED LOGARITHM FOR THE $G / G / \infty$ QUEUE

The LIL for the "number of servers process" of the $G / G / \infty$ queue may be proved using similar techniques to those for cluster processes (note that
here we only have right-hand 'clusters'). Hence we will give no proofs. Let the arrival process $\eta$ and the process of service times $\left\{V_{j}\right\}$ be define on a common probability space ( $\Omega, F, P r$ ) (we need this formulation in Corollary 3.5.5). Let

$$
\begin{aligned}
& A_{n}(t) \equiv(n(0, n t]-m n t) / \sigma_{1} \phi(n) \\
& S_{n}(t) \equiv \sum_{i=1}^{[n t]}\left(V_{i}-E\left(V_{1}\right)\right) / \sigma_{i} \phi(n)
\end{aligned}
$$

for some appropriate constants $\sigma_{1}, \sigma_{2}, m=E n(0,1]<\infty$, where
$\phi(n)=(2 n \log \log n)^{\frac{1}{2}}$.
THEOREM 3.5.4. Suppose $\eta$ and $\left\{V_{j}\right\}$ are stationary with finite second moments, and $\left(A_{n}, S_{n}\right)$ jointly obey a FLIL; then the process

$$
\left(\int_{0}^{n t} m(s) d s-m E\left(V_{1}\right) n t\right) / \sigma \phi(n)
$$

obeys a FLIL, where $\sigma^{2}=\left(E V_{1}\right)^{2} \sigma_{1}^{2}+m \sigma_{2}^{2}$.
Once again we require $m<1$. A similar theorem holds if $\eta$ is replaced by $\eta^{0}$, the arrival process corresponding to $P_{1}^{0}$ (see Section
2.4). The proof again consists of approximating by a coherent mapping, and showing various remainder terms converge a.s. to zero. Geometric subsequence arguments suffice for this latter half of the proof, indeed, seem necessary for the $\eta^{0}$ case, but in the $\eta$ case neater proofs can be devised via Kronecker's Lemma and monotonicity.

We will say that a process $X(t), t \geq 0$, obeys the ordinary law of the iterated logarithm (OLIL) if, with $W(t), t \geq 0$, a Brownian motion defined on the same space $(\Omega, F, P r)$,

$$
\left|X(t)-\phi(t)^{-1} W(t)\right| \rightarrow 0 \quad \text { a.s. }
$$

The techniques of the proof of Theorem 3.5.4 reveal
COROLLARY 3.5.5 (see Problem 6.4.7). If the armival process is
stationary with $m<\infty$, and $A_{n}(1)$ obeys an OLIL, and the service times are i.i.d. with finite second moment, then

$$
\left(\int_{0}^{t} m(s) d s-m E\left(V_{1}\right) t\right) / \sigma \phi(t)
$$

obeys the OLIL.
Note that the FLIL has been proved under weaker moment conditions than the functional central limit theorem. Also, if a non-stationary arrival process $\eta$ (with appropriate centring and norming) satisfies the OLIL, and $E_{\eta}(I)=O(|I|),|I| \rightarrow \infty$, then so will the "accumulated number of servers" process.

### 3.5.3. LIMIT LAWS FOR THE DOUBLY STOCHASTIC POISSON PROCESS

Kingman (1964) has shown that a doubly stochastic Poisson process may be represented as a random time transformation of a stationary Poisson process of unit parameter. Hence it is unnecessary to prove functional limit laws for these processes, as they will follow from e.g. Section 17 of Billingsley (1968), or Lemma 3.4.7 here. Ordinary central limit theorems can be deduced from Lemma 2.5.1, rendering unnecessary the characteristic function techniques of Grandell (1971), at least in his cases $0<k<\infty$ (see Problem 6.4.8).

## CHAPTER 4

## ASYMPTOTIC INDEPENDENCE OF POINT PROCESSES, PARTICULARLY CLUSTER PROCESSES

### 4.1. Introduction

The concept of asymptotic independence arose simultaneously with ergodic theory, and its suitability, in the form of strong mixing and $\phi$-mixing, as a sufficient condition for central limit theorems has been known since the papers of Rosenblatt (1956), Billingsley ((1956), (1962)) and Ibragimov (1962). Important recent developments are due to Oodaira and Yoshihara ((197la), (1971b), (1972)) and Heyde (1974). Stronger mixing conditions have also been studied e.g. Philipp (1969).

Our original motivation for considering $\phi$-mixing of point processes was the hope that it would be preserved under the clustering operation (independent subsidiaries), and hence give a simple avenue to functional limit theorems for cluster point processes (Daley and Vere-Jones (1972), Theorem 8.6). In fact, Westcott (1973), (Concluding remarks), suggests that limit laws obtained by more direct methods (as in Chapter 3) should be weakened in the presence of a mixing condition. However, our investigations suggest that, unlike weaker forms of asymptotic independence (Westcott (1971), (1972)), $\phi$-mixing and strong mixing are only maintained under extra conditions, which are quite severe (bounded clusters) in the case of $\phi$-mixing.

Our basic definitions conclude this section, and in Section 4.2 we give two examples of point processes which are not $\phi$-mixing; the first of these examples illustrates, in Section 4.3, the problem of characterizing complete $\phi-$ and strong mixing, and that complete mixing (or $\phi$-mixing) on a determining class is not always sufficient to establish complete mixing (respectively $\phi$-mixing). Section 4.4 investigates the preservation of various modes of mixing under clustering, and Section 4.5 contains concludir
remarks.
We will need the following definitions: let

$$
\sigma(N(B)) \equiv \sigma_{\text {gen }}\left\{\{N(A)=k\}, A \in B(R) \cap B, k \in Z_{+}\right\}, B \in B(R),
$$

i.e., the smallest $\sigma$-field such that the maps $N \mapsto N(A)$ (for any $A \in B(R) \cap B$ ) from $N \rightarrow R$ are measurable.

We shall introduce the term 'complete mixing' for a stationary point process $P$ satisfying, for a given function $\gamma(\cdot)$,

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \gamma(\tau) P(C) P(D) \tag{1.1}
\end{equation*}
$$

whenever $C \in \sigma(N(-\infty, t]), D \in \sigma(N(t+\tau, \infty)), t \in R, \tau>0$ 。Here $\gamma:[0, \infty) \rightarrow[0,1]$ is a monotone decreasing function satisfying $\lim \gamma(\tau)=0$. This type of mixing appears in at least one paper of Philipp $\tau \rightarrow \infty$
(1969). As a non-trivial example (not a point process), consider a process that is defined as a real-valued function on the state space of a discrete time aperiodic irreducible stationary Markov chain on some finite state space. Such a process is completely mixing (for proof, see Billingsley (1968), pp. 167-8, where the example is used to illustrate the weaker concept of $\phi$-mixing).

A stationary point process $P$ is $\phi$-mixing for a given function $\phi(\cdot)$ if

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \phi(\tau) P(C), \tag{1.2}
\end{equation*}
$$

and strong mixing for a given function $\alpha(\cdot)$ if

$$
\begin{equation*}
|P(C \cap D)-P(C) P(D)| \leq \alpha(\tau), \tag{1.3}
\end{equation*}
$$

where $C$ and $D$ are as before, and $\phi(\tau)$ and $\alpha(\tau)$ have the same properties as $\gamma(\tau)$ (see Problem 6.5.13).

In the following we will assume, as usual, that the centre process is stationary, and that the process of subsidiaries is stationary as well as being independent of the centre process. Also, we will take it for granted that our cluster processes exist (in the sense of Theorem 3.2.2), and the subsidiaries are a.s. finite.

### 4.2. Two examples of point processes which are not $\phi$-mixing

It is well-known that the output of a stationary $M / G / \infty$ queue is a stationary Poisson process, being equivalent to the random translations of points of the original (input) process (see, e.g., Daley (1975) for references).

Let $G_{\delta}(\delta>0)$ denote a distribution function satisfying

$$
\begin{align*}
G_{\delta}(x) & =0, \quad x \leq \delta, \\
& =G(x-\delta), \quad x>\delta, \tag{2.1}
\end{align*}
$$

where $G(x)$ is any arbitrary distribution function on $[0, \infty)$. The following is true:

THEOREM 4.2.1. Superposing the input and output of a $M / G / \infty$ queue for which $G(x)<1$ for all finite $x$ results in a stationary point process which is not $\phi$-mixing.

Remark. In this example, and the next, we will take our processes on an arbitrary probability space $(\Omega, F, P r)$, although we could equally as well work from $N \times R_{+}^{Z}$.

Proof (see also Problem 6.5.2). We will only consider the case where $G(x)=F_{\delta}(x)$, some $\delta>0$, and some $F$. The stronger result can be proved using the techniques of Example 4.2.2, but the $F_{\delta}$ case is neater.

Let $\lambda$ be the common parameter of the input $\left(n_{1}\right)$ and output $\left(\eta_{2}\right)$ processes. The superposed input and output process is denoted by $n\left(=n_{1}+n_{2}\right)$.

Let $h$ satisfy $\exp (-2 \lambda h)>\frac{3}{2}, 0<h \leq \delta$. Then $\eta_{1}(-h, 0]$ and $\eta_{2}(-h, 0]$ are independent Poisson variables, and hence, for all $k \in Z_{+}$,

$$
\begin{equation*}
\operatorname{Pr}\{n(-h, 0]=k\}=e^{-2 \lambda h}(2 \lambda h)^{k} / k!. \tag{2.2}
\end{equation*}
$$

Now using Milne (1970), Theorem 2, or from first principles,
$\operatorname{Pr}\{n(-h, 0]=k, n(\tau, \tau+h]=0\}=\exp (-4 \lambda h+2 \lambda h p(\tau))(2 \lambda h(1-p(\tau)))^{k} / k$ ! , (2.3) where $p(\tau)=\int_{-h}^{0}\left[F_{\delta}(\tau+h-v)-F_{\delta}(\tau-v)\right] d v / 2 h$.

Thus from (2.2) and (2.3), we have that
$\operatorname{Pr}\{n(\tau, \tau+h]=0 \mid n(-h, 0]=k\}-\operatorname{Pr}\{n(\tau, \tau+h]=0\}$
$=\exp (-2 \lambda \hbar) \cdot\left[\exp (2 \lambda h p(\tau))(1-p(\tau))^{k}-1\right]$.
For all finite $x, F_{\delta}(x)<1$, so for any given finite $\tau$ we can find $\tau^{*}>\tau$ such that $p\left(\tau^{*}\right)>0$, and hence $1-p\left(\tau^{*}\right)<1$ 。 Thus $\sup _{k}\left|\operatorname{Pr}\left\{n\left(\tau^{*}, \tau^{*}+h\right]=0 \mid n(-h, 0]=k\right\}-\operatorname{Pr}\left\{n\left(\tau^{*}, \tau^{*}+h\right]=0\right\}\right|$

$$
=\exp (-2 \lambda h)>\frac{1}{2} .
$$

Let $P_{\eta} \equiv P_{P} n^{-1}$ be the probability measure on $(N, \sigma(N))$ corresponding to n. Clearly,

$$
\phi(\tau) \equiv \sup _{C, D}\left|P_{\eta}(D / C)-P_{\eta}(D)\right| \geq \phi\left(\tau^{*}\right)>\frac{3}{2},
$$

where the supremum is taken over $C \in \sigma(N(-\infty, t]), D \in \sigma(N(t+\tau, \infty))$. Hence $\phi(\tau) \nmid 0(\tau \rightarrow \infty)$, which shows that the process is not $\phi$-mixing.

Unbounded translations of a $\phi$-mixing point process do not necessarily result in a $\phi$-mixing point process. Counter-examples can be difficult to establish however, since, as remarked before, it is well-known that Poisson processes are invariant under translation, and the "counting" behaviour of non-Poisson centre processes after translation is in general algebraically intractable (or well-nigh so), with the exception of the translation of compound Poisson processes, as in the following example.

EXAMPLE 4.2.2. Consider a centre process of Poisson doublets (i.e., a stationary Poisson process with each point doubled) with rate 1 , and i.i.d. translations $X$ such that $\operatorname{Pr}\{X>t\}=e^{-t}$. Hence if $\eta$ is the cluster process (see the remark at the end of this section), then by standard techniques we can prove

$$
\begin{equation*}
E\left\{z^{\eta(-h, 0]}\right\}=\exp \left(-\sigma+z \rho+z^{2} \nu\right) \tag{2.5}
\end{equation*}
$$

where $\rho=2\left(1-e^{-h}\right), \quad \nu=h-\left(1-e^{-h}\right)$ and $\sigma=\rho+\nu$. (This is the probability generating function of the so-called Hermite distribution (Kemp and Kemp (1965)).)

Similarly we can prove

$$
\begin{equation*}
\left.E\left\{z^{n(-h, 0]}\right]_{\{n(\tau, \tau+h]=0\}}\right\}=\exp \left(-\gamma+z \delta+z^{2} \nu\right) \tag{2.6}
\end{equation*}
$$

where $\delta=\rho-\left(1-e^{-h}\right)^{2} e^{-\tau}$, and $\gamma=2 \sigma-\left(1-e^{-h}\right)^{2} e^{-\tau}$ and $1_{A}(\omega), \omega \in \Omega$, is the indicator function of the set $A$.

Let $a_{n}(\alpha, \beta)=\sum_{j=0}^{n} \alpha^{n-j} \beta^{2 j} /((n-j)!(2 j)!)$ for any $\alpha, \beta \geq 0$. From
(2.5) and (2.6),
$\operatorname{Pr}\{n(\tau, \tau+h]=0 \mid n(-h, 0]=2 k\}=\exp \left(\sigma-\left(1-e^{-h}\right)^{2} e^{-\tau}\right) \cdot d_{k}(\nu, \delta) / d_{k}(\nu, \rho)$.
We assert that $\lim _{k \rightarrow \infty} d_{k}(\nu, \delta) / d_{k}(\nu, \rho)=0$. To prove this, choose
$\varepsilon \geq 0$ and $J$ such that $\delta^{2 j} \leq \varepsilon \rho^{2 j}$ for $j>J$. Then for $k>J$,

$$
\begin{equation*}
d_{k}(\nu, \delta) / d_{k}(\nu, \rho) \leq \varepsilon+\left[\sum_{j=0}^{J} \nu^{k-j} \delta^{2 j} /((k-j)!(2 j)!)\right] / d_{k}(\nu, \rho) \tag{2.7}
\end{equation*}
$$

Let $\tau$ be any fixed integer, $0 \leq \tau \leq J$. Then

$$
\left(\nu^{k-\tau} \delta^{2 l}\right) /\left((k-\imath)!(2 \imath)!d_{k}(\nu, \rho)\right)=\left(\delta^{2} / \nu\right)^{2} /\left((2 \tau)!D_{k}\right)
$$

where $D_{k}=\sum_{j=0}^{k}(k-Z)!\left(\delta^{2} / \nu\right)^{j} /((k-j)!(2 j)!)$.
But the $(\tau+1)$ th term of $D_{k}$ is $(k-\eta)\left(\delta^{2} / \nu\right)^{\tau+l} /(2 \tau+2)$ ! . Hence $D_{k} \rightarrow \infty$ as $k \rightarrow \infty$, and all the terms in the fixed sum of (2.7) have limit 0 . Hence for $C \in \sigma(N(-\infty, t]), D \in \sigma(N(t+\tau, \infty))$,

$$
\sup _{C, D}\left|P_{\eta}(D \mid C)-P_{\eta}(D)\right| \geq \exp (-\sigma),
$$

where $P_{\eta}$ is the probability measure corresponding to $\eta$, so that $\eta$ is not $\phi$-mixing.

Note that the point process in this last example is the same as an elementary Neyman-Scott cluster process with Poisson centres, and in which each subsidiary consists of 2 points whose positive distances $D_{1}, D_{2}$ from the centre have independent exponential distributions with parameter 1 .
4.3. On characterisations of complete, $\phi$ - and strong mixing

Westcott (1971) has given a neat characterisation of mixing in terms of the probability generating functional (p.g.fl.) (see also Westcott (1972)). This is very useful for cluster point processes, since, if $G_{1}[\xi]$ is the p.g.fl. of the centre process, and $G_{2}(\xi \mid t)$ is the p.g.fl. of a subsidiary given its centre is at $t$ (we are considering independent subsidiaries here) then the p.g.fl. $G[\xi]$ of the cluster point process is given by

$$
G[\xi]=G_{1}\left[G_{2}(\xi \mid t)\right] .
$$

(Our definition of p.g.fls. is over a suitable class $V$ of real-valued measurable functions $\xi$ for which $I-\xi$ has bounded support, i.e., $\xi(t)=1$ for $t$ outside some bounded set.) Let $V_{I} \subset V$ denote those functions $\xi$ in $V$ for which $1-\xi$ vanishes outside $I$. Then we can prove the following

THEOREM 4.3.1. Let $\xi_{1} \in V_{(-\infty, t]}, \xi_{2} \in V_{(t, \infty)}$, and thus $S_{\tau} \xi(u) \equiv \xi(u-\tau) \in V_{(t+\tau, \infty)}, \tau>0$. If a stationary point process with p.g.fl. G is
(a) completely mixing with rate $\gamma(\tau)$, then

$$
\begin{equation*}
\left|G\left[\xi_{1} S_{\tau} \xi_{2}\right]-G\left[\xi_{1}\right] G\left[\xi_{2}\right]\right| \leq \gamma(\tau) G\left[\xi_{1}\right] G\left[\xi_{2}\right] ; \tag{3.1}
\end{equation*}
$$

(b) $\phi$-mixing, then

$$
\begin{equation*}
\left|G\left[\xi_{1} S_{\tau} \xi_{2}\right]-G\left[\xi_{1}\right] G\left[\xi_{2}\right]\right| \leq 2 \phi(\tau) G\left[\xi_{1}\right] . \tag{3.2}
\end{equation*}
$$

The proof of ( $a$ ) is almost immediate. For a proof of (b), see Theorem 4.3.3, in which similar, but slightly more complex techniques are used (see

Problem 6.5.7).
The converses to Theorem 4.3.1 do not hold. It is easy to demonstrate that the p.g.fl. of the superposed input and output of the $M / M / \infty$ queue satisfies (3.1) (and thus also (3.2)), with

$$
G[\xi]=\exp \left\{-\int_{-\infty}^{+\infty}\left[1-\xi(t) \int_{0}^{\infty} \xi(t+x) \mu d x\right] \lambda d t\right\}
$$

and

$$
\begin{equation*}
\gamma(\tau)=\min \left(1, \exp \left[\lambda \mu^{-1} e^{-\mu \tau}\right]-1\right) \tag{3.3}
\end{equation*}
$$

where $\lambda$ is the arrival rate, and $\mu$ the service time parameter (see Problem 6.5.4).

For orderly point processes $P$, mixing can be characterized using the zero probability function

$$
\begin{equation*}
\phi(B)=P\{N(B)=0\}, B \in B(R) \tag{3.4}
\end{equation*}
$$

Kurtz (1974) has stated without proof similar characterisations of stationarity and ergodicity.

THEOREM 4.3.2. If, for an orderly point process $P$,

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} \phi\left(B_{1} \cap\left(B_{2}+\tau\right)\right)=\phi\left(B_{1}\right) \phi\left(B_{2}\right) \tag{3.5}
\end{equation*}
$$

for any $B_{1}, B_{2} \in R(I)$, the ring of finite unions of intervals, then $P$ is mixing.

Proof. Firstly note that $S \equiv\{\{N(B)=0\}, B \in R(I)\}$ is closed under intersections, and that $\sigma_{\text {gen. }}(S)=\sigma(N)$ (Kallenberg (1973)). Then, for an fixed $A_{2} \in S$,

$$
D_{1} \equiv\left\{A_{1}: \lim _{\tau \rightarrow \infty} P\left(A_{1} \cap T_{\tau} A_{2}\right)=P\left(A_{1}\right) P\left(A_{2}\right)\right\}
$$

is a Dynkin system (Ash (1972), p. 168) and hence by the Dynkin system theorem, $D_{1} \supset \sigma_{\text {gen }}(S)$. Similarly, for any fixed $A_{1} \in \sigma(N)$,

$$
D_{2} \equiv\left\{A_{2}: \lim _{\tau \rightarrow \infty} P\left(A_{1} \cap T_{\tau} A_{2}\right)=P\left(A_{1}\right) P\left(A_{2}\right)\right\}
$$

is a Dynkin system so that $D_{2} \supset \sigma_{\text {gen }}$. (S) .
The equations for complete mixing and $\phi$-mixing analogous to (3.5)
follow from (3.1) and (3.2) by taking $\xi^{\prime} s$ of the form $1-1_{A}(t)$. Hence, similar characterisations do not hold for complete mixing and. $\phi$-mixing. If they did, the converses to Theorem 4.3.1 would be true for orderly point processes.

We turn now to characterisations of strong mixing, aiming at an analogue of Theorem 4.3.1. Firstly we note that we can extend the domain of the probability generating functionals to include complex-valued, measurable functions $\phi \in \Phi$ with $\phi(x)=1$ on the complement of a bounded set, and satisfying $\sup _{x \in R}|\phi(x)|=1$ (see, e.g., Fisher (1972)). This, because we anticipate that the RHS of (3.2) will no longer contain $G\left[\xi_{1}\right]$, and the modulus on the LHS can be interpreted more liberally. If, however, we introduce the characteristic functional (ch.fl.)

$$
\begin{equation*}
C[\theta] \equiv G[\exp i \theta(t)] \tag{3.6}
\end{equation*}
$$

where $\theta(t)$ has bounded support (i.e., $\theta \in V$ ), and define

$$
\Phi_{(t, t+x]} \equiv\{\phi \in \Phi: 1-\phi \text { has support within }(t, t+x]\}
$$

we can easily show, via the multidimensional maximum modulus principle (e.g., Gunning and Rossi (1965), p. 7) that

$$
\begin{equation*}
\left|G\left[\phi_{1} S_{\tau} \phi_{2}\right]-G\left[\phi_{1}\right] G\left[\phi_{2}\right]\right| \leq\left|C\left[\theta_{1}+S_{\tau} \theta_{2}\right]-C\left[\theta_{1}\right] C\left[\theta_{2}\right]\right| \tag{3.7}
\end{equation*}
$$

where $\phi_{1} \in \Phi_{(-\infty, t]}, \phi_{2} \in \Phi_{(t, \infty)} ; \theta_{1} \in V_{(-\infty, t]}, \theta_{2} \in V_{(t, \infty)}$. Hence the ch.fl. is (potentially) more useful than the p.g.fl. in characterising strong mixing.

THEOREM 4.3.3. If a stationary point process $P$ is strongly mixing with rate $\alpha(\tau)$, then its ch.fl. satisfies

$$
\begin{equation*}
\left|C\left[\theta_{1}+S_{\tau} \theta_{2}\right]-C\left[\theta_{1}\right] C\left[\theta_{2}\right]\right| \leq 4 \alpha(\tau) \tag{3.8}
\end{equation*}
$$

where $\theta_{1} \in V_{(-\infty, t]}, \theta_{2} \in V_{(t, \infty)}$.
Proof. Let $I_{1}, \ldots, I_{k} \in B(R) \cap(t+\tau, \infty)$, and
$J_{1}, \ldots, J_{Z} \in B(R) \cap(-\infty, t]$. Denote $\left(I_{1}, \ldots, I_{k}\right)$ by $I$ and
$\left(J_{1}, \ldots, J_{2}\right)$ by $J$. Also, put

$$
\begin{aligned}
\mathrm{P} & \equiv P(\mathrm{I}, \mathrm{n} ; \mathrm{J}, \mathrm{~m}) \\
\equiv P\left\{N\left(I_{1}\right)\right. & \left.=n_{1}, \ldots, N\left(I_{k}\right)=n_{k} ; N\left(J_{1}\right)=m_{1}, \ldots, N\left(J_{\imath}\right)=m_{l}\right\} \\
& -P\left\{N\left(I_{1}\right)=n_{1}, \ldots, N\left(I_{k}\right)=n_{k}\right\} P\left\{N\left(J_{1}\right)=m_{1}, \ldots, N\left(J_{\imath}\right)=m_{l}\right\} .
\end{aligned}
$$

Then

$$
\begin{align*}
\mid E\left(\Pi_{r} \Pi_{s} \exp i\left[\theta_{r} N\left(I_{r}\right)+\psi_{s} N\left(J_{s}\right)\right]\right)- & E\left(\Pi_{r} \exp \left[i \theta_{r} N\left(I_{r}\right)\right]\right) E\left(\Pi_{s} \exp \left[i \theta_{s} N\left(J_{s}\right)\right]\right) \mid \\
& =\left|\Sigma_{n} \Sigma_{m} P(I, n ; J, m) \exp [i(\theta \cdot n+\psi \cdot m)]\right| \tag{3.9}
\end{align*}
$$

where $\theta \equiv\left(\theta_{1}, \ldots, \theta_{k}\right), \psi \equiv\left(\psi_{1}, \ldots, \psi_{Z}\right)$. Now set $\gamma=\theta \cdot n+\psi \cdot m$ and define $D_{1}=\{n, m: P \cos \gamma \geq 0\}, D_{3}=\{n, m: P \sin \gamma \geq 0\}$, and let $D_{2}, D_{4}$ be their respective complements. Clearly (3.9) $\leq \max \left(\left|\Sigma_{D_{1}} P \cos \gamma\right|,\left|\Sigma_{D_{2}} P \cos \gamma\right|\right)$

$$
\begin{equation*}
+\max \left(\left|\Sigma_{D_{3}} P \sin \gamma\right|,\left|\Sigma_{D_{4}} P \sin \gamma\right|\right) . \tag{3.10}
\end{equation*}
$$

Now let $E_{1}=\{n, m: \cos \gamma \geq 0\}, E_{3}=\{n, m: \sin \gamma \geq 0\}$ and $E_{2}, E_{4}$ be their respective complements. Then

$$
\begin{equation*}
\left|\Sigma_{D_{1}} P \cos \gamma\right| \leq\left|\Sigma_{D_{1} \cap E_{1}} P \cos \gamma\right|+\left|\Sigma_{D_{1} \cap E_{2}} P \cos \gamma\right| \leq 2 \alpha(\tau), \tag{3.11}
\end{equation*}
$$

and similarly for the other expressions in the right-hand side of (3.10). Hence (3.9) $\leq 4 \alpha(\tau)$. If we define

$$
\theta_{1}(u)=\sum_{i=1}^{l} \psi_{i} 1_{J_{i}}(u), \quad \theta_{2}(u)=\sum_{i=1}^{k} \theta_{i} 1_{I_{i}-\tau}(u)
$$

then we have

$$
\begin{equation*}
\left|C\left[\theta_{1}+S_{\tau} \theta_{2}\right]-C\left[\theta_{1}\right] C\left[\theta_{2}\right]\right| \leq 4 \alpha(\tau) . \tag{3.12}
\end{equation*}
$$

A measurable function $\theta \in V$ can be approximated pointwise by a sequence of simple functions within the same class, so, since $C[\theta]$ is continuous for $\theta \in V$ (cf. Westcott (1972)), (3.8) follows.

Examples of disjoint distributions with uniformly close characteristic functions are well known in Fourier series theory (for probabilistic
examples, see Chung (1968), Ex. 6.3.12, or Dudley (1968), Section 4), so that a converse to Theorem 4.3.3 would clearly require extra conditions. According to Prohorov and Rozanov ((1969), P. 162), convergence "in variation' (which we require) cannot adequately be expressed in terms of characteristic functions, so we abandon this approach in favour of more direct techniques.

Finally, for reference, we quote a known theorem which characterizes complete, $\phi$ - and strong mixing in terms of rings of events $R(N(I))$ generating the $\sigma$-algebras $\sigma(N(I))$ used in the definitions.

THEOREM 4.3.4. A stationary point process $P$ is
(a) complete mixing if and only if

$$
|P(C \cap D)-P(C) P(D)| \leq \gamma(\tau) P(C) P(D)
$$

(b) $\phi$-mixing if and only if

$$
|P(C \cap D)-P(C) P(D)| \leq \phi(\tau) P(C)
$$

(c) strong mixing if and only if

$$
|P(C \cap D)-P(C) P(D)| \leq \alpha(\tau)
$$

for all $C \in R(N(-\infty, t]), D \in R(N(t+\tau, \infty)), \quad t \in R, \tau>0$, where $\gamma, \phi, \alpha:[0, \infty) \rightarrow[0,1]$ are monotone decreasing functions satisfying $\lim _{\tau \rightarrow \infty} \gamma(\tau)=\lim _{\tau \rightarrow \infty} \phi(\tau)=\lim _{\tau \rightarrow \infty} \alpha(\tau)=0$.

Proof. (b) is proved on p. 167 of Billingsley (1968), and (a) and (c) are proved similarly.

As far as we know, Theorem 4.3 .4 cannot usefully be weakened. For example, $R(N(-\infty, t])$ may, in (c), be replaced by a class $S_{1}$ such that $R(N(-\infty, t])=S_{1} \cup\left\{C^{c}: C \in S_{1}\right\}$, and similarly for $R(N(t+\tau, \infty))$, but we know of no stronger reduction.

The characterization Theorem 4.3.4 is as it stands not particularly useful to us. It needs to be supplemented in particular cases by characterizations of the sets in the classes $R(N(-\infty, t])$ and $R(N(t+\tau, \infty))$ so that complete, $\phi^{-}$, or strong mixing can be proved on these classes.

[^1]Examples occur in the next section.

### 4.4. Strong mixing of cluster processes

In this section we develop a technique which allows us decide whether a cluster process is mixing in some sense if its centre process and the process of subsidiaries are mixing in the same sense. The most interesting results occur in the case in which the centre process is strongly mixing, and hence we will introduce the ideas in that setting.

Let the centre process of a cluster process be a measure $P_{1}$ on $(N, \sigma(N))$ and the process of clusters be a measure $P_{2}$ on $\left(N^{R}, \sigma(N)^{R}\right)$. Here $\sigma(N)^{A}$ for any given $A \in B(R)$ is a product $\sigma$-field on $N^{R}$, i.e., the smallest $\sigma$-field such that the maps $N^{R} \rightarrow R^{k}\left(k \in Z_{+}\right)$defined by $\left\{N_{v}\right\} \rightarrow\left(N_{v_{1}}\left(B_{1}\right), \ldots, N_{v_{k}}\left(B_{k}\right)\right), B_{i} \in B(R), v_{i} \in A, 1 \leq i \leq k$, are measurable. Also we will denote members of $N^{R}$ by $N$.

Let $N_{\infty}$ be the set of non-negative integer or infinite valued measures on $R$ which may be infinite on bounded Borel sets. Define $\eta_{c}: N \times N^{R} \rightarrow N_{\infty}$ by

$$
\begin{equation*}
\eta_{c}\left(N,\left\{N_{v}\right\}\right)=\int N_{v}(\cdot-v) d N(v) \tag{4,1}
\end{equation*}
$$

Throughout the following, we will assume that $\eta_{c} \in N, P_{1} \times P_{2}-$ a.s., and also that $\eta_{c}$ is measurable with respect to the product $\sigma$-field $\sigma(N) \times \sigma(N)^{R}$ (see Problem 6.5.10). We define the cluster process via (4.1) (i.e., as $\left.P_{c} \equiv\left(P_{1} \times P_{2}\right) \eta_{c}^{-1}\right)$, rather than as in Chapter 3, to avoid substantial indexing problems. Note that our process therefore differs from the usual model, in that if the centre process has a multiple event of size $n$ at $x$, then the contribution to the RHS of (4.1) is $n N(\cdot-x)$, rather
than the sum of $n$ different subsidiaries all having the same origin. Our model and the usual correspond if the centre process is orderly. We point out that if the subsidiaries are i.i.d., the theorems of this section may be proved for non-orderly centre processes in the manner of Kerstan, Matthes and Mecke (1974), p. 326. We prefer the simplicity of our formulation (4.1). As usual, we require that $P_{1}$ be stationary with respect to $T_{y}: N \rightarrow N$ defined by $T_{y} N(\cdot)=N(\cdot+y)$, and also we will take $P_{2}$ stationary with respect to $S_{x}: N^{R} \rightarrow N^{R}$ defined by $S_{x}\left\{N_{v}\right\}=\left\{N_{v+x}\right\}$. Hence let $T_{y} \times S_{y}: N \times N^{R} \rightarrow N \times N^{R}$ by specified by $T_{y} \times S_{y}(N, N)=\left(T_{y} N, S_{y} N\right)$. Then clearly $P_{1} \times P_{2}$ is stationary with respect to $T_{y} \times S_{y}$, and, for $C \in \sigma(N)$,

$$
\begin{aligned}
P_{c}\left(T_{y} C\right) & =P_{1} \times P_{2}\left\{(N, N): \int N_{v}(\cdot-v) d N(v) \in T_{y} C\right\} \\
& =P_{1} \times P_{2}\left\{(N, N): \int N_{v}(\cdot-v+y) d N(v) \in C\right\} \\
& =P_{1} \times P_{2}\left\{(N, N): \int N_{v+y}(\cdot-v) d N(v+y) \in C\right\} \\
& =P_{1} \times P_{2}\left\{T_{y} \times S_{y}\left(n_{c} \in C\right)\right\} \\
& =P_{c}(C)
\end{aligned}
$$

so that $P_{c}$ is also stationary.

Our most relevant results will occur in the case of $P_{2}$ having
independent increments, i.e., if $C \in \sigma(N)^{I}, D \in \sigma(N)^{J}, I, J$ intervals, $I \cap J=\emptyset$, then $P_{2}(C \cap D)=P_{2}(C) P_{2}(D)$.

Finally, before proving our theorem, we define an event $G_{\tau}$ in $\sigma(N) \times \sigma(N)^{R}$ as follows: for an arbitrary fixed $t \in R$, $G_{\tau}=\left\{(N, N): \int_{-\infty}^{t+2 \tau} N_{v}(t+3 \tau-v, \infty) d N(v)+\int_{t+\tau}^{\infty} N_{v}(-\infty, t-v] d N(v)=0\right\}$. (4.2)

This event imposes restrictions on the extent of interactions of the subsidiary processes. Its occurrence is sufficient to allow the properties of the centre process to dominate, since it effectively means that events of the cluster process within ( $-\infty, t$ ) and $(t+3 \tau, \infty)$ only interact via the centre process if $P_{2}$ has independent increments (see Problem 6.5.1).

THEOREM 4.4.1. If $P_{1}$ is strongly mixing with rate $\alpha(\tau), P_{2}$ has independent increments, and $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1(\tau \rightarrow \infty)$, then $P_{c}$ is strongly mixing with rate $\alpha\left(\frac{1}{3} \tau\right)+6 P_{1} \times P_{2}\left(G_{\frac{1}{3} \tau}^{c}\right)$.

Proof (see Problem 6.5.13). Let $t$ be an arbitrary but fixed number in $R$. We will denote by $\sigma_{-\infty}^{t}$ and $\sigma_{t}^{\infty}$ the product $\sigma$-field
$\sigma(N(I)) \times \sigma(N)^{I}$ for $I=(-\infty, t]$ and $I=(t, \infty)$ respectively. We will say that $P_{1} \times P_{2}$ is strongly mixing if there exists $\psi:[0, \infty) \rightarrow[0,1]$ such that $\lim _{\tau \rightarrow \infty} \psi(\tau)=0$, and

$$
\begin{equation*}
\left|P_{1} \times P_{2}(C \cap D)-P_{1} \times P_{2}(C) P_{1} \times P_{2}(D)\right| \leq \psi(\tau) \tag{4.3}
\end{equation*}
$$

if $C \in \sigma_{-\infty}^{t}$ and $D \in \sigma_{t+\tau}^{\infty}$.
We first prove that $P_{1} \times P_{2}$ is strongly mixing with rate $\alpha(\tau)$.

Consider the class of sets in $\sigma_{-\infty}^{t}$ of the form $C=\sum_{i=1}^{n} C_{1 i} \times C_{2 i}$;
$C_{1 i} \in \sigma(N(-\infty, t]), C_{2 i} \in \sigma\left(N^{(-\infty, t]}\right)$ and $n \in Z_{+}$, where $\left\{C_{2 i}\right\}$ are
disjoint. Then this class forms a ring, for if $F=\sum_{j=1}^{m} F_{1 j} \times F_{2 j}$, then
$C \cup F=\sum_{i=1}^{n}\left[C_{1 i} \times\left(C_{2 i}-\bigcup_{j=1}^{m} F_{2 j}\right)\right]+\sum_{j=1}^{m}\left[F_{1 j} \times\left(F_{2 j}-\bigcup_{i=1}^{n} C_{2 i}\right)\right]$

$$
+\sum_{i=1}^{n} \sum_{j=1}^{m}\left[\left(C_{1 i} \cup F_{1 j}\right) \times\left(C_{2 i} \cap F_{2 j}\right)\right],
$$

$$
C-F=\sum_{i=1}^{n} \sum_{j=1}^{m}\left[C_{1 i} \times\left(C_{2 i}-F_{2 j}\right)\right]+\sum_{i=1}^{n} \sum_{j=1}^{m}\left[\left(C_{1 i}-F_{1 j}\right) \times\left(C_{2 i} \cap F_{2 j}\right)\right]
$$

are in the same class. The same of course holds for sets in $\sigma_{t+\tau}^{\infty}$. By Theorem 4.3.4 (which holds more generally than for point processes), we only need prove $P_{1} \times P_{2}$ strongly mixing for sets in these classes. Hence, with $C$ as above, and $D=\sum_{j=1}^{m} D_{1 j} \times D_{2 j} \in \sigma_{t+\tau}^{\infty}$,

$$
\begin{align*}
&\left|P_{1} \times P_{2}(C \cap D)-P_{1} \times P_{2}(C) P_{1} \times P_{2}(D)\right| \\
&=\left|\sum_{i=1}^{n} \sum_{j=1}^{m}\left[P_{1}\left(C_{1 i} \cap D_{1 j}\right)-P_{1}\left(C_{1 i}\right) P_{1}\left(D_{1 j}\right)\right] P_{2}\left(C_{2 i}\right) P_{2}\left(D_{2 j}\right)\right| \\
& \leq \alpha(\tau) . \tag{4.4}
\end{align*}
$$

For any set $F \in \sigma(N)$, let $F_{\eta} \equiv\left\{\eta_{c} \in F\right\} \in \sigma_{-\infty}^{+\infty}$. Hence if $C \in \sigma((-\infty, t])$, then

$$
\begin{aligned}
C_{\eta} \cap G_{\tau} & =\left\{(N, N): \int_{-\infty}^{+\infty} N_{v}(\cdot-v) d N(v) \in C\right\} \cap G_{\tau} \\
& =\left\{(N, N): \int_{-\infty}^{t+\tau} N_{v}(\cdot-v) d N(v) \in C\right\} \cap G_{\tau} \\
& =C_{\eta}^{\prime} \cap G_{\tau},
\end{aligned}
$$

say, where $C_{\eta}^{\prime} \in \sigma_{-\infty}^{t+\tau}$. Correspondingly, for $D \in \sigma(N(t+3 \tau, \infty))$, $D_{\eta} \cap G_{\tau}=D_{\eta}^{\prime} \cap G_{\tau}$, where $D_{\eta}^{\prime} \in \sigma_{t+2 \tau}^{\infty}$. Thus

$$
\begin{aligned}
& \mid P_{c}(C \cap D)-P_{c}^{(C) P_{c}(D) \mid} \\
& =\left|P_{1} \times P_{2}\left(C_{\eta} \cap D_{\eta}\right)-P_{1} \times P_{2}\left(C_{\eta}\right) P_{1} \times P_{2}\left(D_{\eta}\right)\right| \\
& \leq\left|P_{1} \times P_{2}\left(C_{\eta} \cap D_{\eta} \cap G_{\tau}\right)-P_{1} \times P_{2}\left(C_{\eta} \cap G_{\tau}\right) P_{1} \times P_{2}\left(D_{\eta} \cap G_{\tau}\right)\right|+3 P_{1} \times P_{2}\left(G_{\tau}^{e}\right) \\
& =\left|P_{1} \times P_{2}\left(C_{\eta}^{\prime} \cap D_{\eta}^{\prime} \cap G_{\tau}\right)-P_{1} \times P_{2}\left(C_{\eta}^{\prime} \cap G_{\tau}\right) P_{1} \times P_{2}\left(D_{\eta}^{\prime} \cap G_{\tau}\right)\right|+3 P_{1} \times P_{2}\left(G_{\tau}^{c}\right) \\
& \leq\left|P_{1} \times P_{2}\left(C_{\eta}^{\prime} \cap D_{n}^{\prime}\right)-P_{1} \times P_{2}\left(C_{\eta}^{\prime}\right) P_{1} \times P_{2}\left(D_{n}^{\prime}\right)\right| \\
& +P_{1} \times P_{2}\left[\left(C_{\eta}^{\prime} \cap D_{\eta}^{\prime}\right) \cup G_{\tau}\right]+P_{1} \times P_{2}\left(C_{\eta}^{\prime} \cup G_{\tau}\right)+P_{1} \times P_{2}\left(D_{\eta}^{\prime} \cup G_{\tau}\right) \\
& -3 P_{1} \times P_{2}\left(G_{\tau}\right)+3 P_{1} \times P_{2}\left(G_{\tau}^{c}\right) \\
& \leq \alpha(\tau)+6 P_{1} \times P_{2}\left(G_{\tau}^{c}\right) .
\end{aligned}
$$

It is clear that the same techniques will yield
COROLLARY 4.4.2. (i) If $P_{1}$ is ergodic and $P_{2}$ weakly mixing, then $P_{c}$ is ergodic. Similarly, if $P_{1}$ is weakly mixing and $P_{2}$ ergodic, $P_{c}$ is ergodic.
(ii) $P_{c}$ is weakly mixing if both $P_{1}$ and $P_{2}$ are.
(iii) $P_{c}$ is mixing if both $P_{1}$ and $P_{2}$ are.
(iv) If $P_{1}$ is strongly mixing with rate $\alpha(\tau)$ and $P_{2}$ is completely mixing with rate $\gamma(\tau)$ and $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1(\tau \rightarrow \infty)$, then $P_{c}$ is strongly mixing with rate

$$
\alpha\left(\frac{1}{3} \tau\right)+\gamma\left(\frac{1}{3} \tau\right)+6 P_{1} \times P_{2}\left(G_{\frac{1}{3} \tau}^{c}\right) .
$$

The same conclusion holds if $P_{1}$ is completely mixing with rate $\alpha(\tau)$ and $P_{2}$ is strongly mixing with rate $\gamma(\tau)$.
(v) If $P_{1}$ is $\phi$-mixing and $P_{2}$ completely mixing with rate $\gamma(\tau)$, and the subsidiary processes are $P_{2}-a . s$. bounded (i.e., $P_{2}\left\{N: N_{1}\left(K^{C}\right)=0\right\}=1$ for some bounded interval $\left.K:|K|=\delta\right\}$, then $P_{c}$
is
$\phi_{c}$-mixing, where, for $\tau \geq 0$,

$$
\begin{aligned}
\phi_{c}(\tau) & =\phi(\tau-\delta)+\gamma(\tau-\delta), & \tau \geq \delta, \\
& =1 & , \quad \tau<\delta .
\end{aligned}
$$

Alternatively, if $P_{1}$ is completely mixing (rate $\gamma(\tau)$ ) and $P_{2}$ is $\phi$-mixing and the subsidiaries are bounded as above, then $P_{c}$ is $\phi_{c}$-mixing. (vi) If $P_{1}$ is completely mixing with rate $\gamma_{1}(\tau)$ and the subsidiary processes are completely mixing with rate $\gamma_{2}(\tau)$ and are $P_{2}$-a.s. bounded as in $(v)$, then $P_{c}$ is completely mixing with rate $\gamma_{c}(\tau)$, where

$$
\begin{array}{rlr}
\gamma_{c}(\tau) & =\gamma_{1}(\tau-\delta)+\gamma_{2}(\tau-\delta), \quad \tau \geq \delta, \\
& =1 & , \quad \tau<\delta .
\end{array}
$$

Proof (see Problem 6.5.8). We prove only the first half of (iv), the rest being similar. Clearly all that is needed is an extension of the argument at (4.4). So

$$
\begin{aligned}
\mid P_{1} \times & P_{2}(C \cap D)-P_{1} \times P_{2}(C) P_{1} \times P_{2}(D) \mid \\
= & \mid \sum_{i, j}\left[P_{1}\left(C_{1 i} \cap D_{1 j}\right)-P_{1}\left(C_{1 i}\right) P_{1}\left(D_{1 j}\right)\right] P_{2}\left(C_{2 i} \cap D_{2 j}\right) \\
& +\sum_{i, j}\left[P_{2}\left(C_{2 i} \cap D_{2 j}\right)-P_{2}\left(C_{2 i}\right) P_{2}\left(D_{2 j}\right)\right] P_{1}\left(C_{1 i}\right) P_{1}\left(D_{1 j}\right) \mid \\
\leq & \alpha(\tau)+\gamma(\tau) .
\end{aligned}
$$

Unfortunately the strong mixing rate in Theorem 4.4.1 involves $P_{1} \times P_{2}\left(G_{\tau}\right)$, a quantity which is very difficult to calculate without further knowledge of the structure of the process. Of course, Chebyshev's inequality yields the crude upper bound for $P_{1} \times P_{2}\left(G G_{\tau}^{c}\right)$ of (with $\left.m \equiv E_{1} N(0,1]<\infty\right)$

$$
m \int_{\tau}^{\infty}\left\{E_{2} N_{1}(-\infty,-u]+E_{2} N_{1}(u, \infty)\right\} d u
$$

but by arguing more precisely we can discover a finer condition involving only the component processes which guarantees $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1(\tau \rightarrow \infty)$. Define

$$
\begin{equation*}
L(x)=P_{2}\left\{N: N_{1}(-\infty, x]>0\right\}, R(x)=P_{2}\left\{N: N_{1}(x, \infty)>0\right\} . \tag{4.5}
\end{equation*}
$$

Note that $N_{1}$ could be replaced by any $N_{2}$ because of the stationarity of $P_{2}$.

LEMMA 4.4.3. A sufficient stochastic condition for $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1$ $(\tau \rightarrow \infty) \quad i s$

$$
\begin{equation*}
\int_{a}^{\infty} L(\alpha-v) d N(v)+\int_{-\infty}^{a} R(\alpha-v) d N(v)<\infty, P_{1}-a . s . \tag{4.6}
\end{equation*}
$$

for at least one $a \in R$. If $P_{2}$ has independent increments, and there is a uniform bound on the multiplicity of events of the centre process, then this condition is also necessary.

Remark. This result overlaps with Lemma 3.3.1, but for clarity, we give all details.

Proof. In (4.2), we may write $G_{\tau}=G_{1 \tau} \cap G_{2 \tau}$, where $G_{1 \tau}=\left\{\int_{-\infty}^{t+2 \tau} N_{v}(t+3 \tau-v, \infty) d N(v)=0\right\}, \quad G_{2 \tau}=\left\{\int_{t+\tau}^{\infty} N_{v}(-\infty, t-v] d N(v)=0\right\}$. It is necessary and sufficient for $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1 \quad(\tau \rightarrow \infty)$ that $P_{1} \times P_{2}\left(G_{1 \tau}\right) \rightarrow 1$ and $P_{1} \times P_{2}\left(G_{2 \tau}\right) \rightarrow 1(\tau \rightarrow \infty)$. Let

$$
G_{1 \tau}^{\prime}=\left\{\int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)=0\right\}
$$

By stationarity, $P_{1} \times P_{2}\left(G_{1 \tau}^{\prime}\right)=P_{1} \times P_{2}\left(G_{1 \tau}\right)$. So let $t \in R, \tau_{1}>0$ be arbitrary and fixed. Then note that, for given $N \in N$,

$$
\begin{align*}
\left\{N: \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)<\infty\right\} & =\bigcup_{\tau>\tau_{1}}\left\{N: \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)<\infty\right\} \\
& =\lim _{\tau \rightarrow \infty}\left\{N: \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)=0\right\} \tag{4.7}
\end{align*}
$$

and this limit is monotone $\uparrow$. Hence, writing $F(v)=\left\{N: N_{v}(\alpha-v, \infty)>0\right\}$

$$
\begin{align*}
& \int_{-\infty}^{a} R(a-v) N(\ddot{d v})<\infty \\
& \\
& \quad \Leftrightarrow \int_{-\infty}^{t-\tau_{1}} R(t-v) N(d v)<\infty \quad P_{1} \text {-a.s. for any fixed } a \in R \\
&  \tag{4.8}\\
& \Leftrightarrow \sum_{t_{i}(N) \leq t-\tau_{1}} P_{2}\left\{F\left(t_{i}\right)\right\}<\infty \quad P_{1}-\text { a.s. } \\
& \quad \Rightarrow P_{2}\left\{F\left(t_{i}\right) \text { i.o. }, t_{i} \leq t-\tau_{1}\right\}=0 \quad P_{1}-a . s .
\end{align*}
$$

by the Borel-Cantelli Lemma,

$$
\begin{array}{r}
\Leftrightarrow P_{2}\left\{N: \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)<\infty\right\}=1 \quad P_{1}-\text { a.s., }  \tag{4.9}\\
\text { since the } N_{v}^{\prime} \text { s are a.s. finite, }
\end{array}
$$

$$
\begin{equation*}
\Leftrightarrow P_{2}\left\{N: \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)=0\right\} \rightarrow 1 \quad(\tau \rightarrow \infty) \quad P_{1}-\text { prob., } \tag{4.10}
\end{equation*}
$$

by (4.7),

$$
\Leftrightarrow P_{1} \times P_{2}\left\{(N, N): \int_{-\infty}^{t-\tau} N_{v}(t-v, \infty) N(d v)=0\right\} \rightarrow 1 \quad(\tau \rightarrow \infty)
$$

by dominated convergence,
$\Leftrightarrow P_{1} \times P_{2}\left\{G_{1 \tau}\right\} \rightarrow 1 \quad(\tau \rightarrow \infty)$.
In the forward direction, we have $(4.9) \Rightarrow(4.10)$ by (4.7), but in the reverse direction we are using the fact that if random variables $X_{n} \rightarrow X$ a.s., and $X_{n} \rightarrow Y$ in probability, then $X=Y$ a.s.. Also (4.11) $\Rightarrow$ (4.10) because in (4.11) we have convergence in first mean. Finally, observe that if the centre process is orderly, and $P_{2}$ has independent increments, then the reverse implication in (4.8) is immediate, since $\left\{F\left(t_{i}\right)\right\}$ are then independent events. If the centre process is non-orderly (see Problem 6.5.1: then we can find a subsequence $\left\{t_{i}^{\prime}\right\}$ of $\left\{t_{i}\right\}$ such that $t_{i}^{\prime} \neq t_{j}^{\prime}, i \neq j$ by counting multiple events as one. If there is an upper bound $M$ on the
multiplicity of centre points, then

$$
M \sum_{t_{i}^{\prime \leq t-\tau_{1}}} P_{2}\left\{F\left(t_{i}^{!}\right)\right\} \geq \sum_{t_{i} \leq t-\tau_{1}} P_{2}\left\{F\left(t_{i}\right)\right\}=\infty
$$

so that the Borel-Cantelli Lemma still applies to give $P_{2}\left\{F\left(t_{i}^{\prime}\right)\right.$ i.o. $\}=1$ and hence $P_{2}\left\{F\left(t_{i}\right)\right.$ i.O. $\}=1$.

In a similar manner we can handle the other half of (4.6).
COROLLARY 4.4.4. A sufficient condition for $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1$ is

$$
\begin{equation*}
\int_{0}^{\infty}\{L(-t)+R(t)\} d t<\infty \tag{4.12}
\end{equation*}
$$

if the centre process has $m=E_{1} N(0,1]<\infty$.
Proof. One proof is to take expectations in (4.6), but a more instructive proof is as follows: clearly

$$
\begin{aligned}
P_{1} \times P_{2}\left(G_{\tau}^{c}\right) & \leq P_{1} \times P_{2}\left(G_{1 \tau}^{\prime c}\right)+P_{1} \times P_{2}\left(G_{2 \tau}^{c}\right) \\
P_{1} \times P_{2}\left(G_{1 \tau}^{\prime c}\right) & =\int_{N} P_{2}\left\{{ }_{t_{i} \leq t-\tau} F^{c}\left(t_{i}(N)\right)\right\} d P_{1}(N) \\
& \leq \int_{N} \int_{-\infty}^{t-\tau} P_{2}\left\{F^{c}(v)\right\} d N(v) d P_{1}(N) \\
& =m \int_{\tau}^{\infty} R(v) d v
\end{aligned}
$$

Similarly for $P_{1} \times P_{2}\left(G_{2 \tau}^{c}\right)$.
It would appear from Lemma 4.4.3 and Corollary 4.4.4 that (4.12) is a fairly fine condition for $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1$. This is confirmed in the next result, in which we interpret (4.12) in special cases (see Lawrance (1972)).

LEMMA 4.4.5. (a) If the centre process. $P_{1}$ is Poisson, and $P_{2}$ has independent increments, (4.12) is also necessary.
(b) If the subsidiary structure is Neyman-Scott with $F$ denoting the distribution of each point from the centre and $S$ the number of points per
subsidiary, (4.12) becomes

$$
\begin{equation*}
\sum_{k=0}^{\infty} \operatorname{Pr}\{S=k\} \int_{0}^{\infty}\left[1-F(t)^{k}\right] d t<\infty . \tag{4.13}
\end{equation*}
$$

(c) If the subsidiary structure is Bartlett-Lewis with inter-epoch distribution $F$ and $S$ points per subsidiary, (4.12) becomes

$$
\begin{equation*}
\sum_{k=0}^{\infty} \operatorname{Pr}\{S=k\} \int_{0}^{\infty}\left[1-F^{k^{*}}(t)\right] d t<\infty \tag{4.14}
\end{equation*}
$$

Proof. We prove only (a). The random process $\left\{V_{x}\right\}$ in $R^{R}$ specified by $V_{x}=\inf \left\{t: N_{x}(t, \infty)=0\right\}$ clearly has independent increments. Note that

$$
\begin{align*}
G_{1 \tau}^{\prime} & =\left\{\int_{-\infty}^{t-\tau} N_{x}(t-x, \infty) d N(x)=0\right\} \\
& \left.=\left\{\int_{-\infty}^{t-\tau} 1_{\left\{N_{x}\right.}(t-x, \infty)>0\right\} d N(x)=0\right\} \\
& =\left\{\int_{-\infty}^{t-\tau} 1_{\left\{V_{x}>t-x\right\}} d N(x)=0\right\} \tag{4.15}
\end{align*}
$$

so that calculating $P_{1} \times P_{2}\left(G_{1 \tau}\right)$ becomes equivalent to determining the probability of (4.15) for a Poisson process (of rate $\lambda$, say) subjected to i.i.d. translations $V$. This is easily calculated using the techniques of Milne (1970) to be

$$
\begin{equation*}
P_{1} \times P_{2}\left(G_{1 \tau}\right)=\exp \left(-\lambda \int_{\tau}^{\infty} R(t) d t\right) \tag{4.16}
\end{equation*}
$$

Similarly, we may calculate

$$
\begin{equation*}
P_{1} \times P_{2}\left(G_{2 \tau}\right)=\exp \left(-\lambda \int_{\tau}^{\infty} L(-t) d t\right) . \tag{4.17}
\end{equation*}
$$

### 4.5. Conclusions

Since a Poisson process subjected to i.i.d. translations is also Poisson, the conditions of bounded subsidiaries for $\phi$-mixing and complete
mixing and of $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1$ for strong mixing are not also necessary for preservation of the mixing property. However, the counter-examples contained in Section 4.2 suggest that the clustering operation may not preserve $\phi$-mixing under fairly wide circumstances. The conditioning event $\{N(-h, 0]=n\}$ implies for large $n$ that there are a large number of centres in the neighbourhood of $(-h, 0]$ and hence that $\{N(\tau, \tau+h]=0\}$ will have a reduced probability of occurrence. Hence we may conjecture that the conditions $R(x)>0$ for all finite $x$, and $P_{2}\left\{N_{1}[0, \infty)>1\right\}>0$ may be sufficient for a cluster process not to be $\phi$-mixing.

We should also remark on the applicability of Theorem 4.4.1 to limit laws. All such theorems (Oodaira and Yoshihara (1971a, b), (1972), Heyde and Scott (1973), Heyde (1974)) seem to impose on the strong mixing rate $\alpha^{\prime}(\tau)$ the condition

$$
\begin{equation*}
\int_{0}^{\infty} \alpha^{\prime}(\tau) d \tau<\infty \tag{5.1}
\end{equation*}
$$

if not the stronger condition in which the integrand is raised to the power $(\delta /(2+\delta))$, for some $\delta>0$. As in the proof of Corollary 4.4.4, we have the following upper bound for the strong mixing rate $\alpha_{c}(\tau)$ of the cluster process:

$$
\begin{equation*}
\alpha_{c}(3 \tau) \leq \alpha(\tau)+m \int_{\tau}^{\infty}[L(-t)+R(t)] d t . \tag{5.2}
\end{equation*}
$$

(Write $a(\tau)$ for the second term in the RHS of (5.2). If the centre process is Poisson, $\alpha(\tau)=0$, and $\alpha_{c}(3 \tau)=O(\alpha(\tau))$, but $\neq o(\alpha(\tau))$, when the subsidiaries are i.i.d., by (4.16) and (4.17).) Thus, (5.1) becomes

$$
\begin{equation*}
\int_{0}^{\infty}\{\alpha(t)+t[L(-t)+R(t)]\} d t<\infty \tag{5.3}
\end{equation*}
$$

which is stronger than (3.6) or (for processes with right-hand clusters) (4.2.5) of Chapter 3, so that it does not appear promising to prove limit
laws for cluster processes via strong mixing theorems. However, we point out that strong mixing theorems are potentially able to provide laws of the iterated logarithm for unbounded double-sided clusters, which Theorem 3.4.8 was unable to do.

## CHAPTER 5

## RENEWAL CHARACTERIZATIONS OF POINT PROCESSES

### 5.1. Introduction

Problems of characterization of point process systems arise, amongst other reasons, from the fact that very few classes of point processes can be handled reasonably; renewal processes, after Poisson processes, are prime examples of those which are tractable. Consequently it is pertinent to study characterizations via renewal processes. It should be pointed out, though, that such characterizations are not particularly useful unless complemented by some idea of "robustness". For example, it is known that the superposition of $n$ i.i.d. stationary point processes is approximately Poisson for large $n$, and it is believed (Conjecture 5.3.1 below) that the superposition of $n$ i.i.d. stationary point processes is a renewal process if and only if all processes are Poisson. Thus, although characterization problems are of mathematical interest, they are not necessarily of such practical importance. A suitable definition of "robustness" (see Problem 6.6.1) may be difficult to come by, however. A possibility is to replace the renewal characterizations here by processes with uncorrelated inter-epoch times (Problems 6.6.2, 6.6.4). Alternatively, we may ask for the Prohorov distance $((5.4)$ of Chapter 1$)$ between the point process in question and a renewal (or Poisson) process of the same intensity. The context will probably decide the appropriateness or otherwise of any particular concept.

In this chapter we tackle two characterization problems: the first in a queueing situation, and the second an example concerning superpositions of point processes.

### 5.2. Characterizing the finite capacity $G I / M / 1$ queue with renewal output

### 5.2.1. INTRODUCTION AND SUMMARY

In tandem queueing systems in which the output of one queue becomes the input of the next, it is computationally convenient to know when this output is a renewal process. The problem of characterizing outputs of queueing systems has been of interest since Burke (1956) proved that the stationary $M / M / s$ queue has Poisson output; a summary of some of the work in the problem is available in Daley (1975). The object of this section is to complete the proof of the following results for certain single-server queueing systems with waiting room capacity $N$, i.e., for certain $G I / G / 1 / N$ systems.

THEOREM 5.2.1. The only stationary $G I / M / 1 / N(0 \leq N \leq \infty)$ queueing systems with renewal output are the $M / M / 1 / 0$ and stable (traffic intensity < 1 ) $M / M / 1 / \infty$ systems.

THEOREM 5.2.2. The only stationary $G I / M / 1 / N(0 \leq N \leq \infty)$ queueing systems with adjacent departure intervals independent are the $M / M / 1 / 0$ and stable $M / M / 1 / \infty$ systems.

Theorem 5.2.2 is, of course, stronger than Theorem 5.2.1, which, however, is of more interest to us, in that it demonstrates the nonconservation of a local property (that of being renewal) of a point process when subjected to a sufficiently rough transformation.

Daley ((1968), (1974)) has completely characterized the stationary $G I / G / 1 / 0$ and $G I / M / 1 / \infty$ systems with renewal output; our proof of Theorem 5.2.1 handles only the intermediate situations ( $1 \leq N<\infty$ ). Complete renewal characterization results have now been obtained for the stationary $M / G / 1 / N(0 \leq N \leq \infty)$ (Disney et al. (1973)) and $G I / M / I / N$ systems; the solution techniques involve using the imbedded Markov chain of queue lengths at service completions or at arrivals. The general problem of characterizing the $G I / G / 1 / N$ queue with renewal departure process (Problem
6.6.3) seems very difficult (but see Daley (1975)).

The point of Theorem 5.2.2 is the contrast with the $M / G / 1 / N$
( $1 \leq N<\infty$ ) situation, in which Daley and Shanbhag (1975) show that there are service time distribution functions $G$ for which adjacent departure intervals are independent, although the output is not renewal. We will only prove Theorem 5.2.2 for $1 \leq N \leq \infty$; Daley (1974) handles the $G I / G / 1 / 0$ situation (see the analysis prior to his Theorem 1): adjacent departure intervals are independent if and only if the output is renewal.

Throughout this section we will assume without further comment that our random variables are all defined on a common probability space ( $\Omega, F, P r$ ) .

Consider a single server queueing system with independent service times $\left\{S_{n}\right\},(n=\ldots-1,0,1, \ldots)$, where $\operatorname{Pr}\left\{S_{n} \leq x\right\}=B(x)=1-e^{-\mu x}$, $x \geq 0$ and $\mu>0$. Potential customers amrive at successive epochs $\left\{t_{i}\right\}$ $\left(\ldots \leq t_{0}=0<t_{1} \leq \ldots\right)$ of a renewal process, with inter-epoch times $\left\{\Pi_{i}=t_{i}-t_{i-1}\right\}, \operatorname{Pr}\left\{\Pi_{i} \leq x\right\}=A(x), A(0+)=0$. An arrival finding a queue with $N$ customers waiting does not enter the system, and is not considered here as part of the output (compare Boes, (1969)). The $n$th served arrival finds $Q_{n}$ customers in the system (queue and service), waits for a time $W_{n}$ until service, and is served for a time $S_{n}$.

Let $q_{i}$ denote the number of customers in the system the instant before arrival epoch $t_{i}$, and $W^{(i)}$ the waiting time of the arrival at $t_{i}$ if $q_{i} \leq N$. If $v(n)$ is defined to be the index of the $n$th served arrival, the

$$
\begin{aligned}
& v(0)=\sup \left\{i: t_{i}+W^{(i)} \leq 0, q_{i} \leq N\right\} \\
& v(n)=\inf \left\{i>v(n-1): q_{i} \leq N\right\}, n>0, \\
& v(n)=\sup \left\{i<v(n+1): q_{i} \leq N\right\}, n<0 .
\end{aligned}
$$

Observe that $v(n+1)=v(n)+1$ if $Q_{n}<N$, while when $Q_{n}=N$,

$$
v(n+1)=\inf \left\{i: \sum_{j=v(n)+1}^{i} \pi_{j}>S_{n-N}^{\prime}\right\}
$$

where $S_{n-N}^{\prime}$ is the time elapsing from the $\nu(n)$ th arrival epoch to the completion of servicing of the customer being serviced. Since the service time distribution function $B(\cdot)$ is exponential, $B(\cdot)$ is also the distribution function of $S_{n-N}^{\prime}$.

Denote the stationary state probabilities of the number of customers in the system the instant before an arrival occurs by $\left\{\pi_{j}\right\}, 0 \leq j \leq N+1$, i.e., $\pi_{j}=\operatorname{Pr}\left\{q_{i}=j\right\}$. These probabilities are known to exist (e.g., Keilson (1966)). Clearly

$$
\begin{equation*}
\operatorname{Pr}\left\{Q_{n}=j\right\}=\pi_{j} /\left(\pi_{0}+\ldots+\pi_{N}\right)=\pi_{j} /\left(1-\pi_{N+1}\right), 0 \leq j \leq N, \tag{2.1.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{Pr}\left\{W_{n} \leq x\right\}=B^{Q} n^{*}(x), \quad x \geq 0 \tag{2.1.2}
\end{equation*}
$$

where

$$
\begin{aligned}
B^{0^{*}}(x) & =1 \text { if } x \geq 0, \\
& =0 \quad \text { if } x<0,
\end{aligned}
$$

and $B^{j^{*}}(x)$ is the $j$-fold convolution of $B$ with itself.
Always, then, the time between arrivals of the $n$th and ( $n+1$ )th customers to be served is

$$
\begin{equation*}
\Pi_{n}^{\prime} \equiv \Pi_{v(n)+1}+\ldots+\Pi_{v(n+1)} \tag{2.1.3}
\end{equation*}
$$

Let

$$
\begin{equation*}
I_{n+1} \equiv \max \left(0, \Pi_{n}^{\prime}-W_{n}-S_{n}\right) \tag{2.1.4}
\end{equation*}
$$

denote the idle time for the server between the completion of the $n$th service and the beginning of the $(n+1)$ th service. We define the output process of the queueing system via the sequence $\left\{D_{n}\right\}$ of inter-departure times, these being almost surely positive random variables given by

$$
\begin{equation*}
D_{n}=S_{n}+I_{n} \tag{2.1.5}
\end{equation*}
$$

where $S_{n}$ and $I_{n}$ are independent, since the arrival process and service
times are independent.

### 5.2.2. THE DISTRIBUTION OF INTER-DEPARTURE TIMES

From the definitions, it is easily seen that

$$
\begin{equation*}
\left(1-\pi_{N+1}\right) \operatorname{Pr}\left\{I_{n+1} \leq y\right\}=\sum_{j=0}^{N} \operatorname{Pr}\left\{I_{n+1} \leq y \mid Q_{n}=j\right\} \pi_{j} \tag{2.2.1}
\end{equation*}
$$

After some algebra, we find

$$
\begin{align*}
& \operatorname{Pr}\left\{I_{n+1} \leq y \mid Q_{n}=j\right\}=\int_{0}^{\infty} A(u+y) d B^{(j+1) *}(u), 0 \leq j \leq N-1  \tag{2.2.2}\\
& \operatorname{Pr}\left\{I_{n+1} \leq y \mid Q_{n}=N\right\}
\end{align*}
$$

$$
\begin{align*}
& =\int_{0}^{\infty} d B(u) \int_{0}^{\infty} \operatorname{Pr}\left\{\Pi_{n}^{\prime} \leq y+u+v \mid S_{n-N}^{\prime}=u\right\} d B^{I V^{*}}(v) \\
& =\int_{0}^{\infty} d B(u) \int_{0}^{\infty} d B^{N^{*}}(v) \int_{u}^{y+u+v}[1-A(y+u+v-w)] d H(w) \tag{2.2.3}
\end{align*}
$$

where

$$
\begin{equation*}
H(x)=\sum_{i=1}^{\infty} A^{i *}(x)=A(x)+\int_{0}^{x} A(x-u) d H(u) \tag{2.2.4}
\end{equation*}
$$

is the renewal function of the arrival process. By interchanging the order of integration in (2.2.3), we can replace the renewal function by the Laplace-Stieltjes transform $\alpha(\mu)$ where

$$
\begin{equation*}
\alpha(s)=\int_{0}^{\infty} e^{-s t} d A(t), \quad \operatorname{Re}(s) \geq 0 \tag{2.2.5}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{n+1} \leq y \mid Q_{n}=N\right\}=[1-\alpha(\mu)]^{-1}\left[\int_{0}^{\infty} A(u+y) d B^{(N+1)^{*}}(u)-\alpha(\mu)\right] \tag{2.2.6}
\end{equation*}
$$

Putting

$$
\begin{array}{ll}
\pi_{j}^{\prime}=\pi_{j}, & 0 \leq j \leq N-1 \\
\pi_{N}^{\prime}=\pi_{N} /[1-\alpha(\mu)]=\pi_{N}+\pi_{N+1} & \tag{2.2.7}
\end{array}
$$

we obtain, for $y \geq 0$,

$$
\begin{equation*}
\left(1-\pi_{N+1}\right) \operatorname{Pr}\left\{I_{n+1} \leq y\right\}=\sum_{j=0}^{N} \pi_{j}^{\prime} \int_{0}^{\infty} A(u+y) d B^{(j+1)^{*}}(u)-\pi_{N}^{\prime} \alpha(\mu) . \tag{2.2.8}
\end{equation*}
$$

Using similar techniques, we can obtain the joint distribution of $\left(D_{n}, I_{n+1}\right)$. If $Q_{n}>0$, then $I_{n}=0$, since $S_{n-1}$ has not yet been completed, and by (2.1.2), (2.1.3), (2.1.4) and (2.1.5) we can discuss the joint distribution of $\left(D_{n}, I_{n+1}\right)$ in terms of the arrival of the $n$th served customer. Otherwise we have to consider the ( $n-1$ )th served arrival. For $x \geq 0, y \geq 0$, $\left(1-\pi_{N+1}\right) \operatorname{Pr}\left\{D_{n} \leq x, I_{n+1} \leq y\right\}$

$$
=\sum_{j=1}^{N} \pi_{j} \operatorname{Pr}\left\{D_{n} \leq x, I_{n+1} \leq y \mid Q_{n}=j\right\}
$$

$$
+\sum_{j=0}^{N} \pi_{j} \operatorname{Pr}\left\{D_{n} \leq x, I_{n+1} \leq y, Q_{n}=0 \mid Q_{n-1}=j\right\}
$$

$$
=\sum_{j=1}^{N} \pi_{j}^{\prime} \int_{0}^{\infty} d B^{j^{*}}(u) \int_{0}^{x} A(u+v+y) d B(v)-\pi_{N}^{\prime} \alpha(\mu) B(x)
$$

$$
\begin{equation*}
+\sum_{j=0}^{N} \pi_{j}^{\prime} \int_{0}^{\infty} d B^{(j+1) *}(u) \int_{u}^{u+x} d A(v) \int_{0}^{u+x-v} A(w+y) d B(w) . \tag{2.2.9}
\end{equation*}
$$

### 5.2.3. PROOFS OF THE THEOREMS

A necessary condition for the output to be renewal is that
$\operatorname{Pr}\left\{D_{n} \leq x, D_{n+1} \leq y\right\}=\operatorname{Pr}\left\{D_{n} \leq x\right\} \operatorname{Pr}\left\{D_{n+1} \leq y\right\}$, for all $x, y$. By recognising that $\left(D_{n}, I_{n+1}\right)$ is jointly independent of $S_{n+1}$, and taking Laplace-Stieltjes transforms, this necessary condition reduces to

$$
\begin{equation*}
\operatorname{Pr}\left\{D_{n} \leq x\right\} \operatorname{Pr}\left\{I_{n+1} \leq y\right\}=\operatorname{Pr}\left\{D_{n} \leq x, I_{n+1} \leq y\right\} . \tag{2.3.1}
\end{equation*}
$$

In particular, for all $\gamma \geq 0, \phi \geq-\mu$,

$$
\begin{align*}
\int_{0}^{\infty} e^{-\gamma x} d P r\left\{D_{n} \leq x\right\} & \int_{0+}^{\infty} e^{-(\phi+\mu) y} d P r\left\{I_{n+1} \leq y\right\} \\
& =\int_{0}^{\infty} e^{-\gamma x} \int_{0+}^{\infty} e^{-(\phi+\mu)} y_{x} d_{y} \operatorname{Pr}\left\{D_{n} \leq x, I_{n+1} \leq y\right\} . \tag{2.3.2}
\end{align*}
$$

Write (2.3.2) as $L(\phi, \gamma)=R(\phi, \gamma)$ say. Let $P_{n}(\phi, \gamma)$
$(n=1,2,3,4)$ denote functions of $\phi$ and $\gamma$ that are in fact polynomials of degree $N$ in $\phi$. Then (we omit some algebraic detail; see Section 5.2.4) $L$ and $R$ are expressible in the form

$$
\begin{equation*}
\left(1-\pi_{N+1}\right)(-\phi)^{N+1} L(\phi, \gamma)=\alpha(\phi+\mu) P_{1}(\phi, \gamma)+P_{2}(\phi, \gamma) \tag{2.3.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(1-\pi_{N+1}\right)(\phi-\gamma)(-\phi)^{N} R(\phi, \gamma)=\alpha(\phi+\mu) P_{3}(\phi, \gamma)+P_{4}(\phi, \gamma) \tag{2.3.4}
\end{equation*}
$$

The coefficients of $\phi^{i} \quad(0 \leq i \leq N)$ in each of the polynomials $P_{n}(\phi, \gamma) \quad(1 \leq n \leq 4)$ simplify considerably using the properties of the imbedded Markov chain of queue lengths just prior to arrival. Solving the equation $L(\phi, \gamma)=R(\phi, \gamma)$ then leads to the following expression for $\alpha(\phi+\mu)$, valid for $\phi \geq-\mu, \gamma>0, \phi \neq \gamma$,
$\alpha(\phi+\mu)=\frac{-\gamma \psi \pi_{N+1} \mu^{N+1}+\sum_{j=0}^{N-1}(-\phi)^{N-j} \mu^{j+1}\left[\mu \pi_{j+2}(1-\psi)-\gamma \psi \pi_{j+1}\right]+(-\phi)^{N+1} C(\gamma)}{-\gamma \psi \pi_{N^{\prime}} \mu^{N+1}+\sum_{j=0}^{N-1}(-\phi)^{N-j} \mu^{j+1}\left[\mu \pi_{j+1}^{\prime}(1-\psi)-\gamma \psi \pi_{j}^{\prime}\right]+(-\phi)^{N+1} D(\gamma)}$
where

$$
\begin{equation*}
\psi \equiv \psi(\gamma)=E\left(e^{-\gamma D} n\right) \tag{2.3.6}
\end{equation*}
$$

and $C(\cdot)$ and $D(\cdot)$ are independent of $\phi$; we need note only the form of $D(\cdot)$, namely

$$
\begin{equation*}
D(\gamma)=-\mu \pi_{0} \psi+\left(1-\pi_{N+1}\right)\left[(\mu+\gamma) \psi-\mu \operatorname{Pr}\left\{I_{n}=0\right\}\right] . \tag{2.3.7}
\end{equation*}
$$

We see immediately that $C(\gamma) \equiv 0$, since $\alpha(\phi+\mu) \rightarrow 0$ as $\phi \rightarrow \infty$, but since this yields a complicated functional equation for $\alpha$, we pass it by. Instead, observe that the term independent of $\phi$ in the denominator of (2.3.5) cannot vanish for all $Y$, and that therefore the denominator cannot vanish identically in $\phi$.

From (2.2.8) we observe that

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{n}=0\right\}=1-\pi_{0} /\left(1-\pi_{N+1}\right) \tag{2,3.8}
\end{equation*}
$$

We require $\alpha(0)=1$, which from (2.3.5) is possible only if

$$
\begin{equation*}
\psi(\gamma)=m /(m+\gamma) \tag{2.3.9}
\end{equation*}
$$

where (using (2.3.8))

$$
\begin{equation*}
m=\mu\left[1-\pi_{1} /\left(1-\pi_{0}-\pi_{N+1}\right)\right] . \tag{2.3.10}
\end{equation*}
$$

Taking Laplace-Stieltjes transforms in (2.1.5), and employing (2.3.9), we find

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{n}=0\right\}=m / \mu \tag{2.3.11}
\end{equation*}
$$

and (2.3.7) vanishes identically.
Thus, substituting (2.3.9) into (2.3.5), we obtain, for $\phi \geq-\mu$,

$$
\begin{equation*}
\alpha(\phi+\mu)=\frac{-m \pi_{N+1} \mu^{N+1}+\sum_{j=0}^{N-1}(-\phi)^{N-j} \mu^{j+1}\left[\mu \pi_{j+2}-m \pi_{j+1}\right]}{-m \pi_{N^{\prime} \mu^{N+1}}+\sum_{j=0}^{N-1}(-\phi)^{N-j} \mu^{j+1}\left[\mu \pi_{j+1}^{\prime}-m \pi_{j}^{\prime}\right]} . \tag{2.3.12}
\end{equation*}
$$

Again we set $\alpha(0)=1$, and (2.3.12) yields

$$
\begin{equation*}
m / \mu=\pi_{1} / \pi_{0} \tag{2.3.13}
\end{equation*}
$$

Substituting (2.3.13) into (2.3.12), we deduce that

$$
\begin{equation*}
m / \mu=\pi_{i+1} / \pi_{i}, \quad 0 \leq i \leq N \tag{2.3.14}
\end{equation*}
$$

It follows that, for $\phi \geq-\mu$,

$$
\begin{equation*}
\alpha(\phi+\mu)=m \pi_{N+1} /\left(m \pi_{N}^{\prime}+\phi \pi_{N+1}\right) . \tag{2.3.15}
\end{equation*}
$$

Thus $A(x)=1-e^{-m x}, x \geq 0$, i.e., the system has a Poisson arrival process.

Finally, we appeal to a known result (Finch (1959)) that the $M / M / 1 / N$ queueing system ( $1 \leq N<\infty$ ) never has renewal output, which proves Theorem 5.2.1.

If $N<\infty$, then Theorem 5.2.2 follows immediately from (2.3.15) and Daley and Shanbhag (1975), but Daley (1968) asks for $D_{n}, D_{n+2}$ to be independent in the $G I / M / 1 / \infty$ situation. However, a similar analysis to the above yields instead of (2.3.5),

$$
\begin{equation*}
\alpha(\phi+\mu)=\frac{\sum_{j=0}^{\infty}(-\mu / \phi)^{j+1}\left[\mu \pi_{j+2}(1-\psi)-\gamma \psi \pi_{j+1}\right]+C(\gamma)}{\sum_{j=0}^{\infty}(-\mu / \phi)^{j+1}\left[\mu \pi_{j+1}(1-\psi)-\gamma \psi \pi_{j}\right]+D(\gamma)} \tag{2.3.16}
\end{equation*}
$$

valid for $\phi=-\mu, \phi>\mu, \gamma>\mu$, where

$$
\begin{equation*}
D(\gamma)=-\mu \pi_{0} \psi+\left[(\mu+\gamma) \psi-\mu \operatorname{Pr}\left\{I_{n}=0\right\}\right] \tag{2.3.17}
\end{equation*}
$$

The same procedure gives rise to $\psi(\gamma)=m /(m+\gamma), \quad(\gamma>\mu)$, and also $m / \mu=\pi_{i+1} / \pi_{i}, i \geq 0$. However, (2.3.16) is now unhelpful, in that the RHS gives $0 / 0$; similarly for the identity $C(\gamma) \equiv 0$. However, it is not difficult to establish (Daley (1968), Theorem 3) that

$$
\psi(\gamma)=\mu(\mu+\gamma)^{-1}[\delta \gamma-\mu(1-\delta) \alpha(\gamma)][\gamma-\mu(1-\delta)]^{-1}
$$

where $\delta$ is the unique root in $0<\delta<1$ of $\delta=\alpha(\mu(1-\delta))$. Solving for $\alpha(\gamma)$, we obtain

$$
\alpha(\gamma)=\left[\gamma^{2}(m-\mu \delta)-m \mu^{2}(1-\delta)\right] \cdot\left[-\mu^{2}(m+\gamma)(1-\delta)\right]^{-1}
$$

and hence $m=\mu \delta$, since $\alpha(\gamma) \rightarrow 0$ as $\gamma \rightarrow \infty$. Thus $\alpha(\gamma)=m /(m+\gamma)$, $(\gamma>\mu)$ which establishes $A(\cdot)$ uniquely (Feller (1966), p. 4l0) and thus the input process must be Poisson.

### 5.2.4. SUPPLEMENTS TO SECTION 5.2.3

For reasons of clarity, we outline the transition from (2.3.2) to
(2.3.5) in this section, since the steps are purely manipulative. Straightforwardly, for $\phi \geq-\mu, \gamma>0, \phi \neq \gamma$,

$$
\begin{align*}
P_{3}(\phi, \gamma)=-\mu(-\phi)^{N}\{ & \left.\int_{0}^{\infty} e^{-\gamma u_{F}(u, \gamma) d A(u)}+\sum_{j=1}^{N}(-\mu / \gamma)^{j} \pi_{j}^{\prime}\right\} \\
& -(\phi-\gamma) \sum_{j=1}^{N}(-\mu / \gamma)^{j+1} \pi_{j}^{\prime} \sum_{k=0}^{j-1}(-\gamma)^{k+1}(-\phi)^{N-k-1}, \tag{2.4.1}
\end{align*}
$$

where we have used only

$$
\begin{equation*}
F(u, \gamma) \equiv \sum_{j=0}^{N} \pi_{j}^{\prime} \int_{0}^{u} e^{\gamma w_{d B}(j+1) *}(w) \tag{2.4.2}
\end{equation*}
$$

and the trivial identity

$$
\begin{equation*}
\int_{0}^{t} e^{-\mu x} \mu^{j+1} x^{j} / j!d x=1-e^{-\mu t} \sum_{k=0}^{j}(\mu t)^{k} / k!. \tag{2.4.3}
\end{equation*}
$$

Expanding the final term in (2.4.1) gives

$$
\begin{equation*}
-\mu \sum_{j=1}^{N}(-\mu / \gamma)^{j} \pi_{j}^{\prime}(-\phi)^{N}+\sum_{j=1}^{N} \mu^{j+1} \pi_{j+1}^{\prime}(-\phi)^{N-j} \tag{2.4.4}
\end{equation*}
$$

so that (2.4.1) simplifies to

$$
\begin{equation*}
\left.P_{3}(\phi, \gamma)=-\mu(-\phi)^{N} \int_{0}^{\infty} e^{-\gamma u} u^{(u}, \gamma\right) d A(u)-\sum_{j=1}^{N} \mu^{j+1} \pi_{j}^{\prime}(-\phi)^{N-j} . \tag{2.4.5}
\end{equation*}
$$

A similar surprising simplification occurs in $P_{4}(\phi, \gamma)$,
$P_{4}(\phi, \gamma)=-\alpha(\gamma+\mu) \quad\{$ first term of (2.4.1)\}
$+(\phi-\gamma) \sum_{j=1}^{N}(-\mu / \gamma)^{j+1} \pi_{j}^{\prime} \sum_{k=0}^{j-1}(-\gamma)^{k+1} \int_{0}^{\infty} e^{-\mu t} \sum_{\imath=0}^{k}(-\phi)^{N+\eta-k-1}\left(t^{\eta} / \imath!\right) d A(t)$
and expanding the second term of (2.4.6) we obtain

$$
-(\phi / \gamma)^{N} h(\mu, \gamma)+h(\mu, \phi)
$$

where

$$
\begin{equation*}
h(\mu, \phi)=\sum_{j=1}^{N} \sum_{k=0}^{j-1} \mu^{j+1} \pi_{j}^{\prime}(-\phi)^{N+k-j} \int_{0}^{\infty} e^{-\mu t}\left(t^{k} / k!\right) d A(t) \tag{2.4.7}
\end{equation*}
$$

Denote the transition probabilities of the imbedded Markov chain of the number of people in the queue the instant before arrivals occur by $p_{i j}$. Clearly

$$
\begin{array}{ll}
p_{i j}=\int_{0}^{\infty} e^{-\mu t}(\mu t)^{i-j+1} /(i-j+1)!d A(t) & , 0 \leq i \leq N, 1 \leq j \leq i+1 \\
p_{i 0}=\int_{0}^{\infty}\{1-A(t)\} d B^{(i+1)^{*}}(t) & , 0 \leq i \leq N, \\
p_{N+1, j}=p_{N, j} & , 0 \leq j \leq N+1 . \\
\quad \text { Looking at the coefficient of }(-\phi)^{N-i}, 1 \leq i \leq N, \text { in (2.4.7) we }
\end{array}
$$

$$
\begin{align*}
(-\phi)^{N-i} \sum_{j=i}^{N} \pi_{j}^{\prime} j^{j+1} \int_{0}^{\infty} e^{-\mu t} t^{j-i} /(j-i)!d A(t) & =(-\phi)^{N-i} \mu^{i+1} \sum_{j=i}^{N} \pi_{j}^{\prime} p p_{j, i+1} \\
& =(-\phi)^{N-i} \mu^{i+1} \pi_{i+1} \tag{2.4.8}
\end{align*}
$$

Hence (2.4.6) becomes

$$
P_{4}(\phi, \gamma)=-\alpha(\gamma+\mu) \quad\{\text { first term of (2.4.1) }\}
$$

$$
\begin{equation*}
-(-\phi)^{N} \sum_{i=1}^{N} \mu^{i+1} \pi_{i+1}\left\{(-\gamma)^{-i}-(-\phi)^{-i}\right\} . \tag{2.4.9}
\end{equation*}
$$

With reference to (2.4.5) and (2.4.9) (compare (2.3.7)) note that

$$
\begin{align*}
\mu \int_{0}^{\infty} e^{-\gamma u} F(u, \gamma) d A(u) & =\sum_{j=0}^{N} \pi_{j}^{\prime} \int_{0}^{\infty} e^{-\gamma u} d A(u) \int_{0}^{u} e^{\gamma \omega} d B(j+1) *(w) \\
& =\left(1-\pi_{N+1}\right) \int_{0+}^{\infty} e^{-\gamma u} d \operatorname{Pr}\left\{I_{n+1} \leq u\right\} \\
& =\left(1-\pi_{N+1}\right)\left[(\mu+\gamma) E\left(e^{-\gamma D} n\right)-\mu \operatorname{Pr}\left\{I_{n}=0\right\}\right] . \tag{2.4.10}
\end{align*}
$$

Straightforward manipulation also yields (put $\psi(\gamma) \equiv E\left(e^{-\gamma D} n\right)$ )

$$
\begin{aligned}
& P_{1}(\phi, \gamma)=\psi(\gamma) \sum_{j=0}^{N} \pi_{j}^{\prime} \mu^{j+1}(-\phi)^{N-j} \\
& P_{2}(\phi, \gamma)=-\psi(\gamma) \sum_{j=0}^{N} \pi_{j+1} \mu^{j+1}(-\phi)^{N-j}
\end{aligned}
$$

arguing as in (2.4.8) . Substituting (2.4.10) and the above expressions for $P_{n}(\phi, \gamma)$ into (2.3.3) and (2.3.4) then gives (2.3.5), with

$$
C(\gamma)=-\alpha(\mu+\gamma)\{\text { first term of }(2.4 .1)\}-\mu \sum_{j=1}^{N}(-\mu / \gamma)^{j} \pi_{j+1}-\mu \psi(\gamma) \pi_{1} .
$$

The expression $C(\gamma)=0$ only readily yields information on $\alpha(\mu+\gamma)$ for $\gamma>0$.

It is interesting to note that equation (2.3.5) has its LHS independent of $\gamma$, but not its RHS. We have exploited this fact only peripherally in our analysis. A more direct argument is as follows:

There are two possibilities to consider:
(a) the coefficient of $\phi^{N}$ in the numerator of (2.3.5) vanishes identically, from which we easily deduce

$$
\psi(\gamma)=m /(m+\gamma), \quad m=\mu \pi_{2} / \pi_{1}
$$

and the analysis proceeds as from (2.3.12);
(b) the coefficient of $\phi^{N}$ in the numerator does not vanish identically, and neither therefore can $D(\gamma)$. In this case, since the LHS of (2.3.5) is independent of $\gamma$, so is the RHS, and so the patio of the coefficient of $\phi^{N}$ in the numerator to $D(\gamma)$ must be independent of $\gamma$, equal to some constant $K$, say. This leads to

$$
\begin{equation*}
\psi(\gamma)=m^{\prime} /\left(m^{\prime}+\gamma\right) \tag{2.4.11}
\end{equation*}
$$

where $m^{\prime}=\left[\mu \pi_{2}+\mu K\left(1-\pi_{N+1}\right) \operatorname{Pr}\left\{I_{n}=0\right\}\right] /\left[\pi_{1}+\left(1-\pi_{N+1}\right) K\right]$.
But if (2.4.11) is true, then from (2.1.5) and (2.2.8) we find

$$
\operatorname{Pr}\left\{I_{n}=0\right\}=m^{\prime} / \mu=1-\pi_{0} /\left(1-\pi_{N+1}\right)
$$

and $D(\gamma)$ vanishes identically. Thus case (b) is impossible.

### 5.3. Point processes whose superposition is a renewal process

In this section we solve a problem in support of the following
CONJECTURE 5.3.1. $n$ independent and identically distributed stationary point processes $\eta_{i}(1 \leq i \leq n)$ superpose to an orderly stationary renewal process $\eta$ if and only if all processes are Poisson.

Daley (1973a, b) found counter-examples to the corresponding conjecture in the non-identically distributed case, but now conjectures (private communication) only that at least one of the superposed processes must be Poisson (see Problem 6.6.5). The above conjecture, which he proves (1973a) when the summands are renewal (see also Störmer (1969)) is therefore a special case of his. We prove here only that his counter-example (1973a) does not apply to the above conjecture, and we suspect this non-extension anyway since the argument for his Theorem 1 does not reproduce. Since the
proof of his Theorem 2 and our theorem are virtually identical, we sketch only enough to point out the differences, and assume complete familiarity with Daley's paper.

THEOREM 5.3.2. $n$ i.i.d. stationary alternating renewal processes $\eta_{i}$ superpose to an orderly stationary renewal process $\eta$ if and only if all processes are Poisson.

Proof. As in Daley (1973a), consider

$$
\begin{aligned}
G(x) & \equiv E[\eta(0, x]-m x \mid \eta(\{0\})>0] \\
G(x, y) & \equiv E[\{n[-x, 0)-m x\}\{\eta(0, y]-m y\} \mid \eta(\{0\})>0]
\end{aligned}
$$

where $m=\operatorname{En}(0,1]$. Define $G_{i}(\cdot), G_{i}(\cdot, \cdot)$ for $\eta_{i}(1 \leq i \leq n)$
similarly. It is easily deduced from (9) of Daley (ibid.) that

$$
\begin{equation*}
G(x, y)=G_{1}(x, y)+m\left(1-n^{-1}\right) \int_{0}^{y}\left[G_{1}(x+u)-G_{1}(u)\right] d u \tag{3.1}
\end{equation*}
$$

and immediately it follows that

$$
m\left(1-n^{-1}\right) \int_{0}^{y}\left[G_{1}(x+u)-G_{1}(u)\right] d u=K(x) K(y)
$$

for some function $K(\cdot)$. Hence Daley's Sections $1,2,3$ and 5 immediately carry over. Thus if the lifetime distributions of the alternating renewal processes are $F_{1}, F_{2}$ respectively, then with $Q_{i}(x) \equiv 1-F_{i}(x)$, $i=1,2$,
(i) $Q_{i}(x)=A_{i} e^{-c_{i}^{\prime} x}+\left(1-A_{i}\right) e^{-c_{i}^{\prime \prime} x}, \quad 0<c_{i}^{\prime}<c_{i}^{\prime \prime}$,

$$
0 \leq A_{i} \leq c_{i}^{\prime \prime} /\left(c_{i}^{\prime \prime}-c_{i}^{\prime}\right), A_{i} \neq 1
$$

(ii) $Q_{i}(x)=\left(1+A_{i} x\right) e^{-c_{i} x}, 0<A_{i} \leq c_{i}$,
where $Q_{1}, Q_{2}$ cannot both be of the form (ii), nor both of the form (i) if $A_{i}>0, c_{1}^{\prime}=c_{2}^{\prime}<c_{1}^{\prime \prime}=c_{2}^{\prime \prime}$. If we define

$$
\begin{align*}
R(x) & =\operatorname{Pr}\{\eta(0, x]=0 \mid \eta(\{0\})>0\}  \tag{3.3}\\
R(x, y) & =\operatorname{Pr}\{\eta[-x, 0)=0, \eta(0, y]=0 \mid \eta(\{0\})>0\}, \tag{3.4}
\end{align*}
$$

then since $\eta$ is stationary renewal,

$$
\begin{equation*}
R(x, y)=R(x) R(y) \tag{3.5}
\end{equation*}
$$

Hence, if $R_{i}(\bullet), R_{i}(\bullet, \cdot)$ are defined similarly,
$\int_{x}^{\infty} \int_{y}^{\infty} R_{1}(u, v)\left[\int_{u+v}^{\infty} R_{1}(w) d w\right]^{n-1} d u d v$

$$
\begin{equation*}
=\left(m^{n-1} / n^{n+1}\right)\left[\int_{x}^{\infty} R_{1}(u) d u \int_{y}^{\infty} R_{I}(v) d v\right]^{n} \tag{3.6}
\end{equation*}
$$

where $2 R_{1}(x)=Q_{1}(x)+Q_{2}(x)$, and $2 R_{1}(x, y)=Q_{1}(x) Q_{2}(y)+Q_{2}(x) Q_{1}(y)$. If $Q_{2}$ is of the form (ii), and $Q_{1}$ of (i), then the RHS of (3.6) contains a term in $(x y)^{n}$, whereas the LHS does not. If $Q_{1}$ and $Q_{2}$ are both of the form (i), then the LHS consists of sums of terms of the form $\exp \left(-\sum_{\eta=1}^{n-1} c_{k_{i}} \cdot(x+y)-c_{i} x-c_{j} y\right)$, where $c_{i}$ is $c_{1}^{\prime}$ or $c_{1}^{\prime \prime}$, and $c_{j}$ is $c_{2}^{\prime}$ or $c_{2}^{\prime \prime}$, and $c_{k_{i}}$ is any of $c_{i}^{\prime}, c_{i}^{\prime \prime}$, and the RHS is a sum of terms of the form $\exp \left(-\sum_{z=1}^{n} c_{k_{\eta}} x-\sum_{s=1}^{n} c_{r_{s}} y\right)$. But there exists a smallest (or largest) exponential constant $c$, say, with coefficient $A$, say. Then the coefficient of $\exp (-n c(x+y))$ on the RHS is $(A / 2 c)^{2 n} \cdot\left(m^{n-1} / n^{n+1}\right)$, but zero on the LHS. Since $c$ is the largest or smallest exponential constant, no other sum of $n$ exponential constants can contribute to $\exp (-n c(x+y))$.

## CHAPTER 6

## UNSOLVED PROBLEMS AND GENERALIZATIONS

### 6.1. Introduction

In this final chapter we collect together, for easy reference (rather than having them scattered randomiy throughout the thesis) some unsolved problems and possible generalizations which relate to the ideas presented in the body of the thesis. Of course, there are many more unsolved problems than there were solved, but only few of those we suggest will fall into that elusive class "interesting and (possibly) resolvable". Some of the rest (probably) defy currently available techniques, while others may be of only borderline interest. The problems from each chapter will be listed approximately in order of occurrence with miscellaneous problems last. Reference to its source will be given at the beginning of any problem whose origin is not specified in its statement. If the formulation is due to another person, then he will be appropriately accredited. We will not attempt to impose our own bias by indicating those we regard as most important. Some partial solutions will be provided.

Before embarking on this project, however, we will present in the next section a problem which is solved, but not in the sense of this thesis, and hence is included in this chapter.

### 6.2. Identifiability of the cluster structure of a stationary Poisson cluster process

### 6.2.1. INTRODUCTION

Suppose we know that the centre process of a cluster process is stationary Poisson with finite intensity; can we deduce from a complete record of the cluster process alone the structure of the subsidiary processes? The answer to this question is in the negative when each subsidiary is of unit
size, because Poisson processes subjected to independent random translations are themselves Poisson. Milne (1970), Brown (1970) and Ross (1970) have all shown that the service time distribution of the $M / G / \infty$ queue can be identified from a complete input-output record. We answer affirmatively the corresponding question for Poisson cluster processes here, that is, we prove that the structure of a subsidiary process can be identified given realizations of a cluster process and of the Poisson process from which it is derived, without any information being given on the linkage between the two processes. Bendrath (1974) uses completely different techniques (no probability generating function(al)s) in a generalization of Milne's (1970) result. I thank Dr R.K. Milne for pointing out the existence of this alternative generalization and Dr D.J. Daley for finding the reference.

### 6.2.2. IDENTIFIABILITY WITH GENERAL CLUSTER STRUCTURE

Our solution will be in the spirit of Milne (1970). We shall assume throughout that the subsidiaries are i.i.d., and employ the notation of Section 4.4.

To identify the cluster structure, it is sufficient to specify

$$
\begin{equation*}
P_{2}\left\{N: N_{1}\left(A_{1}\right)=m_{1}, \ldots, N_{1}\left(A_{k}\right)=m_{k}\right\} \tag{2.2.1}
\end{equation*}
$$

for $A_{1}, \ldots, A_{k}$ left-open, right-closed adjacent intervals with rational endpoints, and $k, m_{1}, \ldots, m_{k}$ each within $Z_{+}$(Theorem 1.2.2). We need to define a mapping $\left(i, \eta_{c}\right): N \times N^{R} \rightarrow N \times N_{\infty}$ (compare (4.1) of Chapter 4) by

$$
\begin{equation*}
\left(i, \eta_{c}\right)(N, N)=\left(N, \int_{R} N_{v}(\cdot-v) d N(v)\right) \tag{2.2.2}
\end{equation*}
$$

and define a measure $P$ (with expectation denoted by $E$ ) on $\left(N \times N_{\infty}, \sigma(N) \times \sigma\left(N_{\infty}\right)\right)$ by $P=\left(P_{1} \times P_{2}\right)\left(i, \eta_{c}\right)^{-1}$. We will let $N_{0}$ stand for a typical member of $N_{\infty}$. Note that here we are assuming $\left(i, \eta_{c}\right)$ is
measurable with respect to $\sigma(N) \times \sigma(N)^{R}$ (which will follow from the measurability of $\eta_{c}$ defined at (4.1) of Chapter 4). We will also assume that $\left(i, \eta_{c}\right) \in N^{2}, P_{1} \times P_{2}$-a.s. If $T_{y}: N \rightarrow N$ is defined as usual by $T y^{N(\cdot)}=N(\cdot+y)$, then it is easy to verify that $P$ is stationary with respect to $T_{y} \times T_{y}$, and that the ergodicity of $P$ will follow immediately from the ergodicity of $P_{1} \times P_{2}$, which in turn follows from $P_{1}$ being ergodic, and $P_{2}$ weakly mixing (of. Theorem 4.4.1).

Let $A \equiv(0, r]$ for a given rational $r>0$, and $A_{i} \equiv\left(r_{i-1}, r_{i}\right]$, $r_{0} \equiv 0, r_{i}$ rational, $1 \leq i \leq k$. Setting $\hat{R}=\max \left(r, r_{k}\right)$, we define the sets $D_{i} \in \sigma(N)^{2}$ for $i \in Z_{+}$by

$$
D_{i}=\left\{\left(N, N_{c}\right): N(A+i \hat{R})=m, N_{c}\left(A_{\eta}+i \hat{R}\right)=m_{\eta}(1 \leq \eta \leq k)\right\}=\left(T_{\hat{R}} \times T_{\hat{R}}\right)^{i} D_{0},
$$

where $m, m_{\ell}(l \leq \ell \leq k)$ are each within $Z_{+}$. Clearly
$1_{D_{i}}: N^{2} \rightarrow\{0,1\}$ satisfies $E\left|1_{D_{i}}\right|<\infty$, and hence by the pointwise ergodic theorem,

$$
\begin{equation*}
k^{-1} \sum_{i=0}^{k-1} 1_{D_{i}}+E\left(1_{D_{0}}\right)=P\left(D_{0}\right), P-\text { a.s. } \tag{2.2.3}
\end{equation*}
$$

Hence for $z, z_{1}, \ldots, z_{k}$ rational, we have identified

$$
\begin{equation*}
\left.M\left(z, z_{1}, \ldots, z_{k}\right) \equiv E\left[z^{N(A)} z_{1}{ }_{c}{ }^{( } A_{1}\right) \quad \ldots z_{k}{ }_{c}{ }_{c}\left(A_{k}\right)\right] . \tag{2.2.4}
\end{equation*}
$$

By continuity, we have identified $M\left(z, z_{1}, \ldots, z_{k}\right)$ for all $\left(z, z_{1}, \ldots, z_{k}\right) \in[0,1]^{k+1}$. In particular, we have also identified $\log M\left(z, z_{1}, \ldots, z_{k}\right)$.

Now, it is known for the Poisson cluster process that, for $\xi, \xi_{c}$ within appropriate function classes (see Westcott (1971), Milne (1971),
p. 49), the probability generating functional of $P$ is

$$
\begin{align*}
G\left[\xi, \xi_{c}\right] & \equiv E\left[\exp \int_{-\infty}^{+\infty}\left\{\log \xi(t) d N(t)+\log \xi_{c}(t) d N_{c}(t)\right\}\right] \\
& =\exp \left(-\lambda \int_{-\infty}^{+\infty}\left[1-\xi(t) G_{2}\left(\xi_{c} \mid t\right)\right] d t\right) \tag{2.2.5}
\end{align*}
$$

where $\lambda$ is the rate of the Poisson centre process, and $G_{2}\left(\xi_{c} \mid t\right)$ is the probability generating functional of a subsidiary process whose centre is at $t$, defined by

$$
G_{2}\left(\xi_{c} \mid t\right) \equiv \int_{N^{R}}\left[\exp \int_{-\infty}^{+\infty} \log \xi_{c}(u+t) d N_{1}(u)\right] d P_{2}(N)
$$

If, in particular, we set $\xi(t)=1-(1-z) 1_{A}(t)$, and
$\xi_{c}(u)=1-\sum_{i=1}^{k}\left(1-z_{i}\right) 1_{A_{i}}(u)$, we find that $\log M\left(z, z_{1}, \ldots, z_{k}\right)$ is a
function linear in $\boldsymbol{z}$, so differentiating we have identified

$$
\begin{equation*}
\lambda \int_{-\infty}^{+\infty} 1_{A}(t) E_{2}\left[z_{1}^{N_{1}\left(A_{1}-t\right)} \ldots z_{k}^{N_{1}\left(A_{k}-t\right)}\right] d t \tag{2.2.6}
\end{equation*}
$$

Now $\lambda$ can be identified as $\lim _{n \rightarrow \infty} N(0, n] / n, P_{1}$ - a.s. Since
(2.2.6) is clearly an absolutely convergent power series in $z_{1}, \ldots, z_{k}$, we can equate coefficients to obtain

$$
\begin{equation*}
F(r) \equiv \int_{0}^{r} P_{2}\left\{N_{1}\left(A_{1}-t\right)=m_{1}, \ldots, N_{1}\left(A_{k}-t\right)=m_{k}\right\} d t . \tag{2.2.7}
\end{equation*}
$$

It is not difficult to show that the integrand in $F(r)$ is a continuous function of $t$. To prove this, let us write

$$
P_{2}(A, t, v)=P_{2}\left\{N_{1}\left[\left(A_{i}-t\right) \cap\left(A_{i}-v\right)\right]=m_{i} \quad(1 \leq i \leq k)\right\} .
$$

Then for any $\varepsilon>0$, there exists by consistency condition (v) of Theorem 1.2.1 a $\delta_{1}>0$ such that for $|v-t|<\delta_{1}$

$$
\left|P_{2}(A, t, t)-P_{2}(A, t, v)\right| \leq \varepsilon / 2
$$

and, similarly, a $\delta_{2}>0$ such that for $|v-t|<\delta_{2}$,

$$
\left|P_{2}(A, t, v)-P_{2}(A, v, v)\right| \leq \varepsilon / 2
$$

Hence, for $|v-t|<\min \left(\delta_{1}, \delta_{2}\right)$,

$$
\left|P_{2}(A, t, t)-P_{2}(A, v, v)\right| \leq \varepsilon,
$$

which we require. Thus, evaluating the limit

$$
\left[F\left(r+\Delta_{n}\right)-F(r)\right] / \Delta_{n}
$$

for a sequence of rationals $\Delta_{n} \downarrow 0(n \rightarrow \infty)$, we obtain

$$
P_{2}\left\{N_{1}\left(A_{1}-r\right)=m_{1}, \ldots, N_{1}\left(A_{k}-r\right)=m_{k}\right\}
$$

and thus (2.2.1).
We state this result as
THEOREM 6.2.1. The cluster structure of a stationary Poisson cluster process with i.i.d. subsidiaries is identifiable with probability one from a complete centre process-subsidiary process record.

### 6.2.3. SPECIAL CLUSTER PROCESSES

In both special processes below (see e.g. Lawrance (1972), Section 4, for definitions), we first wish to identify the distribution of $N_{1}(R)$. Note that $\left\{N_{1}(-k, k) \leq \eta\right\} \downarrow\left\{N_{1}(R) \leq \eta\right\}(k \rightarrow \infty), Z \in Z_{+}$, and hence the distribution of $N_{1}(R)$ is identified.
(a) The Bartlett-Lewis process

By standard results

$$
P_{2}\left\{N_{1}(0, t]=0\right\}=1-P_{2}\left\{N_{1}(R)>0\right\} F(t),
$$

so $F$ is identified.
(b) The Neyman-Scott process

Again, looking at the same event,

$$
P_{2}\left\{N_{1}(A)=0\right\}=\sum_{n=0}^{\infty} P_{2}\left\{N_{1}(R)=n\right\}[1-F(A)]^{n}
$$

for any $A \in B(R)$, so that $F$ once again is uniquely specified.

### 6.2.4. CONCLUDING REMARKS

There are several unanswered questions associated with this problem. In particular, we point out that identifiability should be a robust property, and hence we should only need to assume $P$ ergodic. Note that this means that $P_{1}$ must be weakly mixing and $P_{2}$ ergodic or vice versa. One way to tackle the general problem in the case of $P_{2}$ having independent increments would be to expand the p.g.fl. of $P$ (cf. Brillinger (1974)). However, we consider that a more aesthetic technique should be available. Secondly, Milne (1970) has shown that if we are given a realization which is in fact the superposition of a Poisson process and an i.i.d. translation of it (i.e., each point of the Poisson process is independently subjected to a translation with d.f. $G$ say) then, provided $G$ is either symmetric or concentrated on a half-line $[0, \infty)$ or $(-\infty, 0], G$ is still identifiable. We pose the corresponding problem for Poisson centres and right-hand clusters.

Finally, we remark that our formulation gives no idea of the robustness of identifiability of the cluster structure, i.e., given only finite realizations of our processes, can we identify the alternative possible cluster structures and do they differ substantially (in a way that will have to be defined)?

### 6.3. Problems concerning Chapter 2

P.6.3.1 (REMARKS FOLLOWING THEOREM 2.3.1). Define the service process instead as a probability measure $P_{2}$ on $R_{+}^{R}$. Then we can replace $\phi$ by $\phi^{\prime}: N \times R_{+}^{R} \rightarrow \bar{Z}_{+}^{R} \quad$ defined as

$$
\phi^{\prime}(N, x)(t)=\int 1\left[v, v+x_{v}\right)(t) d N(v)
$$

With this new model we can weaken Theorem 2.3.1 (b). But one pays a
price - it is harder to prove measurability of $\phi^{\prime}$, and (more significantly) the model only comesponds with the usual one (i.i.d. service times) if $P_{1}$ is a.s. orderly. We think lack of orderliness more important than the weakening, and so have retained the other model.
P.6.3.2 (THEOREM 2.4.1). Weaken the condition in Theorem 2.4.1 (b) to $E_{2}\left\{x_{1}\right\}<\infty$. Possibly the proof there is insufficient, in that the argument used is not via the ergodic theorem, unlike part ( $\alpha$ ). However, if we disregard the contribution from $t_{i}<0, E_{2}\left\{x_{i}\right\}<\infty$ is sufficient.
P.6.3.3 (LEMMA 2.5.4). Virtually unexplored is the question of replacing $\sqrt{n}$ and $\sqrt{\tau}$ in (5.2) and (5.3) of Chapter 2 by
 which may, however, disappear under closer scrutiny.
P.6.3.4 (END OF SECTION 2.5). Suppose we have a general non-stationary arrival process $P$ with first moment measure $M(A) \equiv E N(A)$. We would like the central limit theorem for $\int_{0}^{\tau} \phi(s) d s$ asking only (in addition to the components satisfying a central limit theorem and possessing suitable moments) that $M[0, x]=O(x) \quad(x \rightarrow \infty) \quad$ (cf. Iglehart and Kennedy (1970)). Then, as in Lemma 2.5.2,

$$
\begin{aligned}
P_{1} \times P_{2}\left\{\tau^{-\delta} \Phi_{1}(\tau) \geq \varepsilon\right\} & \leq E_{1} \times E_{2}\left\{\Phi_{1}(\tau)\right\} / \varepsilon \tau^{\delta} \\
& =\int_{\tau}^{\infty} d F(\tau) \int_{0}^{\tau}(u+\tau-\tau) d M(u) / \varepsilon \tau^{\delta} \\
& +\int_{0}^{\tau} d F(\tau) \int_{\tau-\tau}^{\tau}(u+\tau-\tau) d M(u) / \varepsilon \tau^{\delta}
\end{aligned}
$$

Now the first term is $\leq M(0, \tau] \int_{\tau}^{\infty} \tau d F(\tau) / \varepsilon \tau^{\delta}$, and so is easily
handled. The second term has an upper bound of

$$
\int_{0}^{\tau} \tau M(\tau-\tau, \tau] d F(\tau) / \varepsilon \tau^{\delta} .
$$

If $P$ is stationary, then $M(\tau-\mathcal{L}, \tau]=m \mathcal{L}$, or if $P$ corresponds to a Palm measure then we can use (4.5) of Chapter 2. However, if we have no information such as $M(I)=M(|I|),|I| \rightarrow \infty$, i.e., no information on how nicely segments $I$ (intervals) of our process are behaving (of. functional convergence), then it seems difficult to handle this'term.

The stationarity of the process of service times can also be removed if again one imposes a condition ensuring that no one service time contributes too much to the central limit theorem (of. uniform asymptotic negligibility). For example, suppose that there exists a distribution function $F$ such that

$$
F_{j}(x) \equiv P_{2}\left\{x_{j} \leq x\right\} \geq F(x),
$$

so that, in a sense, the random variable comresponding to $F$ is dominating the service times. Then, after some algebra,

$$
\begin{aligned}
E_{1} \times E_{2}\left\{\Phi_{2}\right\} & =E_{1} \sum_{j=-\infty}^{0}\left\{\tau-\int_{0}^{\tau} F_{j}\left(\tau-t_{j}\right) d \tau\right\} \\
& \leq E_{1} \sum_{j=-\infty}^{0}\left\{\tau-\int_{0}^{\tau} F\left(\tau-t_{j}\right) d l\right\} \\
& =\int_{0}^{\infty} d F(\tau) \int_{-\tau}^{0} \min (\eta+u, \tau) d M(u)
\end{aligned}
$$

and, similarly,

$$
E_{1} \times E_{2}\left\{\Phi_{1}\right\} \leq \int_{0}^{\tau} d M(u) \int_{0}^{\infty} \min (\tau, \tau-u) d F(\tau)
$$

so that the analysis may be carried out as before.
P.6.3.5 (COROLLARY 2.6.12). Our feeling is that this result is suboptimal. It would be interesting to compare this approach to one using characteristic functions and an inequality due to Esseen (Lemma 2, Chung (1968), p. 208) on the stationary $M / G / \infty$ queue (cf. Rao (1966)).
P.6.3.6 (BEGINNING OF SECTION 2.7). The results of Section 2.7 are
therefore not really an extension of those of Section 2.6. An examination of other rates of functional convergence theorems reveals that they too are often expressible in terms of the metric of convergence in probability. This seems strange: one should attempt to obtain rates of functional convergence results asking only that the arrival process and service time process converge to the Wiener process at rates in terms of the Prohorov or dual bounded Lipschitz metrics, although for the $G I / G / \infty$ queue the stronger result is available. Dudley's (1972) comment that rates of functional convergence are better expressed via the Prohorov metric therefore seems open to dispute, although his comparison was with Lipschitz functional formulations. The clue to changing the metric from $\rho$ to $d$ in (7.1) of Chapter 2 is given in Billingsley (1968), p. ll2, line 15.
P.6.3.7 (LEMMA 2.7.1). Rates of convergence results are prime examples for improvement by fine arguments in unexpected directions. It is by no means clear that the $(\log n)^{3 / 4} / n^{\frac{1}{4}}$ term is optimal; in particular, the logarithmic factor.

### 6.4. Problems concerning Chapter 3

P.6.4.1 (THEOREM 3.3.2). The tail sums could hopefully be dealt with as in Barbour (1974), although here we have triangular arrays of convergent series rather than a straight sequence.
P.6.4.2 (LEMMA 3.4.6 (b)). Note that Iglehart (197la), Theorem 3.3 does not need (directly) the properties of $K$ in his proof. In his notation, $\rho\left(S_{n}^{j}, \sigma_{j} \xi_{n}\right) \rightarrow 0$ a.s. has already been proved, where $\xi_{n}(t)=\xi(n t) / \phi(n)$ for $\xi$ a standard Brownian motion. His result thus follows from inequality (5.13) of Chapter 1 , the relative compactness of $\xi_{n}$ (Strassen (1964)) and the Arzelà-Ascoli theorem.
P.6.4.3 (AFTER LEMMA 3.4.6). Of course, these remarks are vacuous if a process $X_{n}$ is defined to obey a FLIL as in Problem 6.4.2 above, i.e., $\rho\left(X_{n}, \xi_{n}\right) \rightarrow 0$ a.s. for $\xi_{n}$ defined there.

In (4.1.7) of Chapter 3 note that if

$$
\begin{aligned}
A_{n} & \equiv\left\{\varepsilon^{2} \phi(n)^{2} \leq X_{1}^{2}<\varepsilon^{2} \phi(n+1)^{2}\right\} \\
E X_{1}^{2} & =\sum_{n=1}^{\infty} \int_{A_{n}} X_{1}^{2} d \operatorname{Pr} \geq 2 \varepsilon^{2} \sum_{n=1}^{\infty} n \operatorname{Pr}\left\{A_{n}\right\}=2 \varepsilon^{2} \sum_{n=1}^{\infty} \operatorname{Pr}\left\{X_{n}^{2} \geq \varepsilon^{2} \phi(n)^{2}\right\}
\end{aligned}
$$

provide the intermediate steps. The ideas succeeding (4.1.7) of Chapter 3 are probably well known - for an instance without details see Heyde and Scott (1973).
P.6.4.4. Removing the condition $m<1$, which is not a logical requirement, but apparently a technical one, in SECTIONS 2.7, 3.3 AND 3.4 seems difficult.
P.6.4.5 (THEOREM 3.4.8). The RHS of (4.2.9) of Chapter 3 is a significant overestimate of the LHS. Using the current arguments, weakening the condition (4.2.4) of Chapter 3 means decreasing in some manageable way the RHS of (4.2.9) of Chapter 3, or estimating tail probabilities in the geometric subsequence argument by something other than Chebyshev's inequality.
P.6.4.6 (THEOREM 3.5.1). A direct proof of the relationship (5.1) of Chapter 3 is not known (to me).
P.6.4.7 (COROLLARY 3.5.5). We seem to require that the service time process obeys a FLIL; the problem is achieving a result parallel to Lemma 3.4.7 for the OLIL.
P.6.4.8 (SECTION 3.5.3). It is by no means clear whether infinite divisibility is a global property. Dr M. Westcott (private communication) has informed me that a student of Matthes is working on the problem: "a doubly stochastic Poisson process is infinitely divisible if and only if its stochastic intensity is infinitely divisible". We pose the question more generally in the light of Kingman's (1964) characterization (not his renewal
characterization - see Section 3.5.3): "if a point process is subjected to a random change of time independent of the original process, then the resultant point process is infinitely divisible if and only if the original and time transformation processes are".

## P.6.4.9. An open question is the generalization of Theorem 4 of Daley

 (1972) to a stationary process of subsidiaries, in the style of Lemma 2.5.4 for the $G / G / \infty$ queue. It is clear that this generalization goes through if we assume that the process of subsidiaries is say, strong mixing, with an appropriate condition on the rate of mixing. We would anticipate that we need only assume$$
\sum_{j=0}^{\infty}\left|E\left\{\left(N_{0}(R)-\mu\right)\left(N_{j}(R)-\mu\right)\right\}\right|<\infty
$$

where $\mu \equiv E\left\{N_{0}(R)\right\}<\infty \quad$ (compare Lemma 3 of Billingsley (1968), p. 172). However, the problem then seems to involve justifying exchanging the order of some limits and integrals.
P.6.4.10 (LEMMA 4, Westcott (1972)). It would be interesting to know whether there exists a small class of functions $\xi$ with nice properties such that the convergence of the p.g.fls.

$$
G_{n}[\xi] \rightarrow G[\xi]
$$

guarantees weak convergence for the corresponding point processes (cf. the dual bounded Lipschitz metric for weak convergence (Section 1.5)).
P.6.4.11 (WEAK CONTINUITY OF CLUSTER PROCESSES). Kennedy (1972a) and Whitt (1974a) have demonstrated weak continuity for queues. Such an idea should extend to cluster processes, but their proofs rely heavily on the queues beginning at $t=0$, and the service times being non-negative. However, in the following, we will demonstrate that weak continuity for cluster processes is a valid concept. The proof is via the p.g.fl., and hence we have not striven for the finest possible conditions via this technique, which demands (unnecessarily, we surmise) that the centre process
and process of subsidiaries be independent, and the subsidiaries themselves be i.i.d. We use the model (2.2) of Chapter 3, but now write
 $P_{2 j}=P^{\prime}$ for some generic point process $P^{\prime}$ for which $P^{\prime}\{N(R)<\infty\}=1$. PROPOSITION 6.4.1. If $P_{1}$ and the sequence $\left\{P_{1 n}\right\}$ are stationary centre processes with finite intensities $m,\left\{m_{n}\right\}$ and satisfy

$$
\begin{equation*}
P_{1 n} \Rightarrow P_{1} \text { and } m<\infty, \quad \underset{n \rightarrow \infty}{\lim \sup } m_{n}<\infty \tag{4.1}
\end{equation*}
$$

and if $P^{\prime}$ and the sequence $\left\{P_{n}^{\prime}\right\}$ of subsidiaries have first moment measures $M^{\prime}(\cdot),\left\{M_{n}^{\prime}(\cdot)\right\}$ satisfying

$$
\begin{equation*}
P_{n}^{\prime} \Rightarrow P^{\prime} \text { and } M^{\prime}(R)<\infty, \quad \underset{n \rightarrow \infty}{\lim \sup _{n}} M_{n}^{\prime}(R)<\infty \tag{4.2}
\end{equation*}
$$

then the sequence of corresponding cluster processes satisfy

$$
P_{c n} \Rightarrow P_{c}
$$

Proof. By Lemma 4 of Westcott (1972) and equation (15) of Vere-Jones (1968), we must prove that the corresponding p.g.fls. satisfy

$$
\begin{equation*}
G_{1 n}\left[G_{n}^{\prime}(\xi \mid t)\right] \rightarrow G_{1}\left[G^{\prime}(\xi \mid t)\right] \tag{4.3}
\end{equation*}
$$

for $\xi \in V$ (recall that weak convergence and convergence of finite dimensional distributions correspond for a sequence of point processes converging to a point process).

We prove (4.3) in two parts. Firstly, by an inequality in the proof of Theorem 2 of Westcott (1972),

$$
\begin{align*}
&\left|G_{1 n}\left[G_{n}^{\prime}(\xi \mid t)\right]-G_{1 n}\left[G^{\prime}(\xi \mid t)\right]\right|  \tag{4.4}\\
& \leq m_{n} \int_{-\infty}^{+\infty}\left|G_{n}^{\prime}(\xi \mid t)-G^{\prime}(\xi \mid t)\right| d t  \tag{4.5}\\
& \leq m_{n}\left\{\int_{-\infty}^{+\infty}\left[1-G_{n}^{\prime}(\xi \mid t)\right] d t+\int_{-\infty}^{+\infty}\left[1-G^{\prime}(\xi \mid t)\right] d t\right\} . \tag{4.6}
\end{align*}
$$

By some easy inequalities, (4.6) has an upper bound of

$$
m_{n}\left\{M_{n}^{\prime}(R)+M^{\prime}(R)\right\}|\operatorname{supp} \cdot[\log \xi]|
$$

and hence by (4.1), (4.2) and the dominated convergence theorem, (4.5) $\rightarrow 0$ $(n \rightarrow \infty)$.

For the second half, let $\xi_{m} \in V$ be a sequence of functions $\not+G_{2}(\xi \mid t)$ (this is always possible since our functions take values in $[0,1]$ ). Then
$\left|G_{1 n}\left[G^{\prime}(\xi \mid t)\right]-G_{1}\left[G^{\prime}(\xi \mid t)\right]\right|$

$$
\leq\left(m_{n}+m\right) ; \int_{-\infty}^{+\infty}\left|\xi_{m^{-}} G^{\prime}(\xi \mid t)\right| d t+\left|G_{1 n}\left[\xi_{m}\right]-G_{1}\left[\xi_{m}\right]\right|
$$

Hence, by Lemma 4 of Westcott (1972),

$$
\begin{equation*}
\lim _{n \rightarrow \infty}(4.7) \leq \kappa \int_{-\infty}^{+\infty}\left|\xi_{m}-G^{\prime}(\xi \mid t)\right| d t \tag{4.8}
\end{equation*}
$$

for some constant $K$. But the RHS of (4.8) has a uniform upper bound of

$$
2 K \int_{-\infty}^{+\infty}\left[1-G^{\prime}(\xi \mid t)\right] d t
$$

which we have already shown to be finite. Hence, by dominated convergence the RHS of $(4.8) \rightarrow 0(m \rightarrow \infty)$, and thus we have proved (4.3).

Hopefully, a more direct analysis will yield a stronger result. We remark, though, that continuity in itself is not really what we are looking for: instead (in an obvious phraseology) we require 'rate of continuity'.

### 6.5. Problems concerning Chapter 4

P.6.5.1 (THEOREM 4.4.1). It was pointed out in Section 4.5 that the condition $P_{1} \times P_{2}\left(G_{\tau}\right) \rightarrow 1(\tau \rightarrow \infty)$ is not a necessary condition for strong mixing of cluster processes to occur. In fact, we know of no example of a cluster process with strong mixing centre which is not strong mixing. But consider the Neyman-Scott process (Poisson centres) with two points per cluster. Then this process always exists (Westcott (1971)), but if $F$, the
distribution of each point of a cluster from its centre, does not have finite mean, then (4.13) of Chapter 4 fails. We have no suggestions, however, as to the relevant events to consider for the potential counter-example.
P.6.5.2 (THEOREM 4.2.1) (D.J. Daley - solved problem). We point out that a continuous time Markov chain need not be $\phi$-mixing. Let $Q(t)$ be the size of a classical immigration-death process; it is well-known that this is a Markov chain. Note that the classical immigration-death process corresponds to the $M / M / \infty$ queue, and $Q(t)$ to the number of servers. Again $\operatorname{Pr}\{Q(\tau)=0 \mid Q(0)=n\}$ can be made arbitrarily small by taking $n$ large enough, so that counter-examples of the style of Section 4.2 follow (compare Bloomfield (1973)).
P.6.5.3. Given a point process which is $\phi$-mixing, and $\phi_{1}:[0, \infty) \rightarrow[0,1]$ satisfying $\phi_{1}(\tau) \downarrow 0(\tau \rightarrow \infty)$ and $\phi_{1}(\tau)>\phi(\tau)$ (all $\tau \in R_{+}$), construct a process which is $\phi_{1}$-mixing but not $\phi$-mixing (P.A.P. Moran).

If we subject the original process to random deletions via an independent process which is $\left(\phi_{1}-\phi\right)$-mixing, then by Corollary 4.4.2(v), the new process is $\phi_{1}$-mixing. It is at least intuitively clear that the derived process will not be $\phi$-mixing.
P.6.5.4 (THEOREMS 4.3.1, 4.3.3). We suspect that the following is true: if the centre process of a cluster process with i.i.d. clusters is complete, $\phi$ - or strong mixing, and its p.g.fl. $G_{1}$ satisfies (3.1), (3.2) or the p.g.fl. equivalent to (3.8) of Chapter 4 respectively, then so does the p.g.fl. $G$ of the cluster process.
P.6.5.5 (D.J. Daley). $\phi$-mixing point processes $P$ may be represented as satisfying, for $C \in \sigma(N(-\infty, t])$ and $D \in \sigma(N(t+\tau, \infty))$,

$$
|P(D \mid C)-P(D)| \leq \phi(\tau) .
$$

Taking $C=\{N(t-h, t]>0\}$ and letting $h \downarrow 0$ suggests that $\phi$-mixing
processes may have nice properties for Palm-Khinchin measures.
P.6.5.6 (D.J. Daley). A more general question is the possibility of a p.g.fl. formulation for the Palm-Khinchin equations, i.e., to relate the p.g.fl. $G(\cdot)$ of a stationary point process to the "conditional p.g.fl. $G_{0}(\cdot)$, the p.g.fl. for the process given a point at zero".
P.6.5.7 (THEOREM 4.3.1) (P.A.P. Moran). Define and study the properties of the "conditional characteristic functional". This remark arose in relation to the inequality (3.2) of Chapter 4. A relevant reference is Bartlett (1938).
P.6.5.8 (COROLLARY 4.4.2). Let $(X, Y) \equiv\left\{\left(X_{j}, Y_{j}\right)\right\}$ be a double sequence of stationary random variables, where $X$ and $Y$ are independent. If $X$ is strong mixing, $Y$ is strong mixing, is $(X, Y)$ strong mixing? Professor E.J. Hannan (private communication) conjectures that this is not so, but knows of no counter-example. Presuming he is correct, then, comparing $(X, Y)$ to $P_{1} \times P_{2}$ and thus $P_{c}$ in Theorem 4.4.1, we conclude (tentatively) that Corollary 4.4.2 (iv) (and (v)) are optimal. If $X$ is ergodic and $Y$ is ergodic then $(X, Y)$ is not necessarily: for a counter-example, see Hannan (1973), p. 163, or Breiman (1968), pp. 100 and 113.
P.6.5.9. Of the same ilk as Problem 8 is the question of superpositions. The techniques of Theorem 4.4 .1 yield: independently superposing a completely mixing point process with a strong mixing point process results in a strong mixing point process. Characteristic functional arguments (Theorem 4.3.3) suggest the stronger conjecture: independently superposing a strong mixing (rate $\alpha_{1}$ ) process with a strong mixing (rate $\alpha_{2}$ ) process yields a strong mixing $\left(\right.$ rate $\alpha_{1}+\alpha_{2}$ ) process. We do not believe the stronger statement.
P.6.5.10 (SECTION 4.4). Find conditions on $\eta_{c}$ which guarantees its measurability. They will probably include the measurability of $v \mapsto N_{v}(A)$ from $R \rightarrow Z_{+}$for any given $\left\{N_{v}\right\} \in N^{R}$ for which $N_{v}(R)<\infty \quad($ all $v \in R)$
and any given $A \in B(R)$.
P.6.5.11. In LEMMA 4.4 .3 we conjecture that if $P_{2}$ has independent increments, the necessity of (4.6) of Chapter 4 holds in general. The problem is to find a path proving

$$
\begin{gathered}
\sum_{t_{i}(N) \leq t-\tau_{1}} P_{2}\left\{F\left(t_{i}\right)\right\}=\infty \quad P_{1}-\text { a.s. } \\
\Rightarrow P_{2}\left\{F\left(t_{i}\right) \text { i.o., } t_{i} \leq t-\tau_{1}\right\}=1 \quad P_{1}-\text { a.s. }
\end{gathered}
$$

by exploiting the independence of subsidiaries located at different points.
P.6.5.12. THEOREM 4.4 .1 should generalize to multivariate point processes. However, a problem of definition is encountered in the generalization to multidimensional point processes: what is meant, for example, by strong mixing of a point process on the plane?
P.6.5.13 (SECTION 4.1) We have not attempted to indicate when a point process is complete, $\phi_{-}$or strong mixing. Many point processes are specified by their inter-epoch times. Therefore a relevant question is: if the inter-epoch times of a point process are complete, $\phi$ - or strong mixing, does the point process itself have the same property?

### 6.6. Problems concerning Chapter 5

P.6.6.1 (SECTION 5.1). Robustness of characterizations of queueing systems is very closely allied to the idea of continuity of queues (Kennedy (1972a), Whitt (1974a)). We discuss this idea further in Problem 6.4.11.

Dr D.J. Daley has suggested the following metric of robustness:

$$
D_{\theta, z}(P, F) \equiv \int_{0}^{\infty} e^{-\theta u}\left|E^{0}\left(z^{N(0, u]}\right)-E_{1}\left(z^{N(0, u]}\right)\right| d u
$$

where $|z|<1, E^{0}$ is the expectation with respect to the Palm measure of a given point process $P$, and $E_{1}$ the expectation of a renewal process $P_{1}$ with lifetime distribution $F$. This metric suffers from the defect
(possibly) that it is zero when $z=1$, but could easily be modified. The point is, though, that $D_{\theta, z}(P, F)$
(a) plays down the importance of what happens for large $u$;
(b) is small when $N(0, u]$ is large,
i.e., it emphasizes local behaviour. (a) and (b) seem to be desirable properties of any 'robustness' metric.
P.6.6.2 (SECTION 5.1) (D.J. Daley). Characterize stationary $G I / M / I / N \quad(1 \leq N<\infty)$ queueing systems with uncorrelated output. Daley (1968) has characterized $G I / M / 1 / \infty$ and $M / G / 1 / \infty$ systems with uncorrelated output, while Daley (1974) and Vlach (1971) have characterized $G I / M / I / O$ systems with such an output. King (1971) has answered the $M / G / 1 / 1$ case, and has given an expression for the covariance in the $M / G / 1 / N$ situation. Preliminary investigations suggest that the techniques of these papers do not assist us here.
P.6.6.3 (SECTION 5.2.1) (D.J. Daley). Characterize stationary $G I / D / I / N \quad(0 \leq N \leq \infty)$ queueing systems with renewal output. It is known (Daley (1974)) that no restrictions on the inter-arrival distribution are required if $N=0$, and that the output is renewal for $N=1$ if the input is Poisson (King (1971)). We conjecture that no other non-degenerate
 to renewal output.
P.6.6.4 (SECTION 5.1) (D.J. Daley). For a pure loss $G I / G / I$ queue (i.e., $G I / G / 1 / 0$ ) find

$$
U(x)=E\{\text { number of departures in }(0, x] \mid \text { departure at } 0\},
$$

and ask when is $U(x)=\lambda x$ for some constant $\lambda$ (i.e., when do departures look like a Poisson process as far as second order properties go)? We could also investigate when $U(x)$ is a renewal function. This may not be a tractable problem, although Laplace transforms may be of assistance.

$$
\text { P.6.6.5 (SECTION 5.3) (D.J. Daley). Let } \eta=\eta_{1}+\eta_{2} \text { be a renewal }
$$

process which is the superposition of independent alternating renewal processes. Prove that either $\eta_{1}$ or $\eta_{2}$ is Poisson. (If Daley's (1973a) technique is to be used, it may be necessary to assume that $\eta_{1}$ is derived from the jump epochs of a Markov chain, i.e., the lifetime distributions are exponential, with rates $\alpha$ and $\beta(\alpha \neq \beta))$.
P.6.6.6 (D.J. Daley). Let $S M$ denote a point process derived from the jump epochs of a stationary irreducible semi-Markov process.

Characterize those SM processes for which

$$
\begin{equation*}
P+S M=R \tag{6.1}
\end{equation*}
$$

( $P$ is Poisson, $R$ renewal ). In particular, it is of interest to know whether any two state semi-Markov process with

$$
p_{i j} \equiv \operatorname{Pr}\{\text { transition } i \rightarrow j \mid \text { lifetime of type } i \text { just ended }\}, i, j=1,2
$$ and with $p_{i i} \neq 0$ can satisfy (6.1). (This is to be compared with Daley's (1973a) counter-example.)

P.6.6.7 (D.J. Daley). Let MC denote a point process derived from the jump epochs of a stationary irreducible continuous time Markov chain on a countable state space. Daley (1973b) has obtained necessary and sufficient conditions on MC to guarantee

$$
P+M C=R .
$$

What is the distribution function $F$ (or its Laplace Transform) of the renewal lifetimes in $R$ ? Daley (private communication) anticipates that

$$
1 / \int_{0}^{\infty} e^{-\theta t} d F(t)=\left(\lambda+\theta+\int_{(0, \infty]}\left(1-e^{-\theta x}\right) \mu(d x)\right) / \lambda,
$$

for some measure $\mu$, because the problem seems to be linked with continuous time regenerative phenomena. Note that this would mean that $R$ is also a doubly stochastic Poisson process (e.g., Kingman, discussion to Bartlett (1963)).
P.6.6.8 (D.J. Daley) Let $\eta(\cdot)$ be a stationary renewal process, and

$$
\begin{aligned}
x_{n}^{h} & =1 \text { if } n(n h,(n+1) h]>0, \\
& =0 \text { otherwise. }
\end{aligned}
$$

CHARACTERIZATION CONJECTURE. $\left\{x_{n}^{h_{n}}\right\}$ is a discrete time stationary renewal process for every $h>0$ if and only if $n(\cdot)$ is Poisson. If 80 , can the conclusion hold if $h$ is restricted to an interval.
P.6.6.9 (D.J. Daley). Let $n(\cdot)$ be a non-arithmetic renewal process with lifetime distribution $F$ having finite mean $\lambda^{-1}$, and set $H(x)=\sum_{n=1}^{\infty} F^{n^{*}}(x)$, where $F^{j *}(x)$ denotes the $j$-fold convolution of $F$. Can skeleton arguments (Kingman (1972), p. 34) be used to show

$$
H(x+h)-H(x) \rightarrow \lambda h \quad(x \rightarrow \infty)
$$

for every $h>0$ ? (That is, can the Blackwell renewal theorem be deduced from the Erdös-Feller-Pollard theorem?)
P.6.6.10 (D.J. Daley). Let a stationary, orderly point process $\eta(\cdot)$
with finite rate $\lambda$ be mixing in some appropriate sense. Does this imply a quasi-Blackwell renewal theorem? That is, does the expectation function $U(x)=E(\eta(0, x] \mid \eta(\{0\})>0)$ satisfy

$$
U(x+h)-U(x) \rightarrow \lambda h \quad(x \rightarrow \infty) ?
$$

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[^0]:    1 The noun is 'weak mixing'; the adjective is 'weak mixing' or 'weakly mixing'; similarly with other types of mixing.

[^1]:    1 I thank Dr C.C. Heyde for this reference.

