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# A. J. Stam <br> Expectation and variance of the volume covered by a large number of independent random sets 

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# EXPECTATION AND VARIANCE OF THE VOLUME COVERED BY A LARGE NUMBER OF INDEPENDENT RANDOM SETS 

A.J. Stam

## 1. Introduction

Throughout this paper the following assumptions and notations apply. The set $A \in \mathbb{R}^{m}$ is bounded and measurable with Lebesgue measure $|A|>0$, but sometimes more stringent assumptions on $A$ will be made. The random $m$-vectors $u_{1}, u_{2}, \ldots$ are independent with common probability density $p(u)=f(\|u\|)$, where $f$ is nonincreasing and positive. Let $N$ be either a nonrandom positive integer or a random variable having Poisson distribution with parameter $\lambda$ and independent of the $u_{i}$. The volume in $\mathbb{R}^{m}$ covered by at least $k$ of the random sets $A+u_{i}, i=1, \ldots, N$, is denoted $V_{k}$. Here $A+y=\{x+y \mid x \in A\}$. The volume covered by exactly $k$ of the random sets is denoted $W_{k}$. We take $k \geqq 1$. If $N=0$ we define $V_{k}=W_{k}=0$. Replacing $A$ by $A+a$ does not change $W_{k}$ or $V_{k}$, so it is no restriction to assume $0 \in A$ or $0 \in \operatorname{Int}(A)$ if $A$ has nonempty interior.

Under the above conditions we have $V_{k} \rightarrow \infty$, a.s. if $N \rightarrow \infty$ and in probability if $\lambda \rightarrow \infty$, so that $E V_{k} \rightarrow \infty$. Moran in his 1974 paper [6] took $m=3$, assumed that the $u$, have standard normal distribution and derived the principal term in the asymptotic expansion of $E V_{1}$ as $\lambda \rightarrow \infty$ or $N \rightarrow \infty$. This term does not depend on $A$. He also proved the curious result that if $A$ is a ball the variance $\sigma^{2}\left(V_{1}\right)$ tends to zero as $\lambda \rightarrow \infty$, or $N \rightarrow \infty$, viz.

$$
\begin{equation*}
0<K_{1} \leqq(\log E N)^{1 / 2} \sigma^{2}\left(V_{1}\right) \leqq K_{2}<\infty, \quad E N \geqq 2 \tag{1.1}
\end{equation*}
$$

and he showed that $V_{1}$ when centered and normalized in asymptotically normal. We will extend the results on expectation and variance, with special attention to the dependence on $A$, to densities of the form

$$
\begin{align*}
& p(u)=f(\|u\|)=c \exp \left(-\alpha^{-2 \theta}\|u\|^{2 \theta}\right),  \tag{1.2}\\
& p(u)=f(\|u\|)=c \exp \{-g(\|u\|)\},  \tag{1.3}\\
& p(u)=f(\|u\|)=1 / h(\|u\|) \tag{1.4}
\end{align*}
$$

where $\theta>0$ and $g$ and $h$ vary regularly at $\infty$. We need some results on regularly varying functions, to be found in De Haan [1], Feller [3] and Seneta [7]. The function $f:(0, \infty) \rightarrow(0, \infty)$ varies regularly of order $\rho$ at $\infty$ or at 0 , respectively, if

$$
\begin{equation*}
f(t x) / f(t) \rightarrow x^{\rho}, \quad t \rightarrow \infty(t \rightarrow 0), \quad x>0 \tag{1.5}
\end{equation*}
$$

If $\rho=0$ we say that $f$ varies slowly. If (1.5) holds and $\rho>-1$, then

$$
\begin{equation*}
\int_{0}^{x} f(s) \mathrm{d} s \sim(\rho+1)^{-1} x f(x), \quad x \rightarrow \infty(x \rightarrow 0) \tag{1.6}
\end{equation*}
$$

If $f$ varies regularly of order $\rho$ at $\infty$, it has the Karamata representation

$$
\begin{equation*}
f(x)=c(x) x^{\rho}\left\{\exp \int_{1}^{x} s^{-1} \epsilon(s) \mathrm{d} s\right\}, \quad x \geqq 1 \tag{1.7}
\end{equation*}
$$

with $c(x)>0, x \geqq 1, c(x) \rightarrow c>0, \epsilon(x) \rightarrow 0, x \rightarrow \infty$. If $f$ is nondecreasing, it follows from (1.7) with $f(\lambda x) \leqq f(\lambda), x \leqq 1$, that

$$
\begin{equation*}
f(\lambda x) / f(\lambda) \leqq \max \left\{1, C x^{\rho+\delta}\right\} \tag{1.8}
\end{equation*}
$$

with $C=C(\delta)$, for $x>0, \lambda \geqq \lambda(\delta)$, for any $\delta>0$. Let $g:(0, \infty) \rightarrow(0, \infty)$ be nondecreasing with $g(x) \rightarrow \infty, x \rightarrow \infty$. The generalized inverse $g^{-1}$ of $g$ is defined by

$$
\begin{equation*}
g^{-1}(y)=\sup \{x \mid g(x) \leqq y\}, \quad g\left(0^{+}\right)<y<\infty \tag{1.9}
\end{equation*}
$$

If $\rho>0$, then $g$ varies regularly at $\infty$ of order $\rho$ if and only if $g^{-1}$ varies regularly at $\infty$ of order $\rho^{-1}$. But $g^{-1}$ may vary slowly without $g$ varying regularly and then $x^{n} / g(x) \rightarrow 0, x \rightarrow \infty, n>0$.

In Section 3 the asymptotic behaviour of $E V_{k}$ and $E W_{k}$ as $\lambda \rightarrow \infty$ for Poissonian $N$ will be studied. Under (1.4) they tend to $\infty$ of the same order depending on $h$ and $m$ with coefficients involving $k$ and $|A|$. Under (1.3) there is no dependence on $A$ in the principal term and we may have $E W_{k} \rightarrow 0$.

Estimating the variance is far more difficult and we have to restrict our attention to $V_{1}$ since in contrast to (2.12) below the expressions for the other variances contain terms with opposite signs. In Section 4 we consider Poissonian $N$. Our results suggest that $\sigma^{2}\left(V_{1}\right)$ may tend to zero only if $p(u)$ decreases faster than exponentially. If the $u_{,}$have probability density (1.2) with $\theta>\frac{1}{2}$ and $A$ is convex and bounded by a regular surface $\partial A$ with curvatures bounded below and above,

$$
\begin{align*}
& 0<K_{1} \leqq(\log E N)^{-b_{1}} \sigma^{2}\left(V_{1}\right) \leqq K_{1}<\infty, \quad E N \geqq 2  \tag{1.10}\\
& b_{1}=(2 \theta)^{-1}\left\{m-1+(m+3)\left(\frac{1}{2}-\theta\right)\right\}, \tag{1.11}
\end{align*}
$$

in accordance with (1.1). So $\sigma^{2}\left(V_{1}\right) \rightarrow 0, \sigma^{2}\left(V_{1}\right)$ is bounded and $\sigma^{2}\left(V_{1}\right) \rightarrow$ $\infty$, if $m-1+(m+3)\left(\frac{1}{2}-\theta\right)$ is negative, zero and positive, respectively. However, the order of decrease or increase of $\sigma^{2}\left(V_{1}\right)$ does not only depend on $m$ and $\theta$, but, strongly, on the shape of $A$. It was shown in Stam [8] that if $\theta=1$ and $A$ is a parallelopipedum, we have for Poissonian $N$

$$
\sigma^{2}\left(V_{1}\right) \sim B(\log \log \lambda)^{m-1} / \log \lambda, \quad \lambda \rightarrow \infty .
$$

For subexponentially decreasing $p$ of the form (1.3) or (1.4) the exact first-order term for $\sigma^{2}\left(V_{1}\right)$ is derived. We see that then $\sigma^{2}\left(V_{1}\right) \rightarrow \infty$.

In Section 5 the results for Poissonian $N$ are extended to nonrandom $N$ by comparing the corresponding expectations and variances. Section 2 contains the general relations from which our proofs will start. We will take $m \geqq 2$. Results for $m=1$ will be stated in remarks. The proofs that may different from those for $m \geqq 2$ are given in Stam [9], as are the derivations of the second-order terms that in some cases may be found.

We write $\int g(x) \mathrm{d} x$ for the Lebesgue integral of $g$ over $\mathbb{R}^{m}$,

$$
\begin{equation*}
d_{A}=\operatorname{diam}(A)=\max \{\|x-y\|: x, y \in A\} \tag{1.12}
\end{equation*}
$$

$A-B=\{x-y \mid x \in A, y \in B\},-A=\{-x \mid x \in A\}$ and $I_{A}$ for the indicator function of $A$. The inner product of $m$-vectors is written $(x, y)$. The notation $(r, \omega)$ will always stand for the spherical coordinates of $x \in \mathbb{R}^{m}$, so that $r=\|x\|, \omega=x / r$ and $\omega$ is the generic point of the unit sphere $\Omega$. The area element of $\Omega$ will be written $\operatorname{d} \sigma(\omega)$ or $\mathrm{d} \sigma$, so that

$$
\begin{equation*}
S_{m}=\int_{\Omega} \mathrm{d} \sigma(\omega)=2 \pi^{m / 2} / \Gamma(m / 2), \quad m \geqq 2 \tag{1.13}
\end{equation*}
$$

is the area of the unit sphere in $\mathbb{R}^{m}$. If $m=1$ a good definition turns out to be $S_{1}=2$. By the identification $x=(r, \omega)$ we will sometimes write $g(r, \omega)$ for $g(x)$. The argument $\omega$ often will be suppressed.

In estimates and remainder terms we will often have to do with functions of $r, \omega, \lambda$ or $N$ that are bounded by some constant. These functions all will be denoted by the same symbol $\eta$ or by $\eta_{+}$if they are nonnegative. The argument usually will be suppressed.

## 2. General formulae

The probability that the point $x \in \mathbb{R}^{m}$ is covered by the set $A+u_{t}$ is, since $p(u)=f(\|u\|)$,

$$
\begin{align*}
P(x) & =P(r, \omega)=P\left(u_{t} \in-A+x\right)=\int_{-A+x} p(u) \mathrm{d} u \\
& =\int_{A} p(x-v) \mathrm{d} v=\int_{A} f\left(\left\{r^{2}-2 r(v, \omega)+\|v\|^{2}\right\}^{1 / 2}\right) \mathrm{d} v \tag{2.1}
\end{align*}
$$

The probability that $x$ and $y$ are covered by $A+u$, is

$$
\begin{aligned}
P(x, y) & =P\left(u_{1} \in(-A+x) \cap(-A+y)\right) \\
& =\int I_{-A+x}(u) I_{-A+y}(u) p(u) \mathrm{d} u=\int_{A(y-x)} p(x-v) \mathrm{d} v,
\end{aligned}
$$

$$
\begin{equation*}
A(z)=A_{z}=A \cap(A-z) \tag{2.2}
\end{equation*}
$$

By interchanging intergrations and using the fifth term in (2.1) and the third term in (2.2) we find

$$
\begin{equation*}
\int P(x) \mathrm{d} x=|A|, \quad \int P(x, y) \mathrm{d} y=|A| P(x) \tag{2.4}
\end{equation*}
$$

From (2.2) and (2.3)

$$
\begin{equation*}
P(x, y)=0, \quad y-x \notin A-A . \tag{2.5}
\end{equation*}
$$

Let $\xi(x)$ be the number of random sets that cover $x$. Then

$$
\begin{align*}
& E W_{k}=\int P(\xi(x)=k) \mathrm{d} x  \tag{2.6}\\
& \sigma^{2}\left(V_{1}\right)= \iint\{P(\xi(x) \geqq 1, \xi(y) \geqq 1) \\
&-P(\xi(x) \geqq 1) P(\xi(y) \geqq 1)\} \mathrm{d} x \mathrm{~d} y \\
&= \iint\{P(\xi(x)=\xi(y)=0) \\
&-P(\xi(x)=0) P(\xi(y)=0)\} \mathrm{d} x \mathrm{~d} y \tag{2.7}
\end{align*}
$$

From (2.6), the independence of the $u_{t}$ and (2.1) we find for nonrandom $N$

$$
\begin{equation*}
E W_{k}=\binom{N}{k} \int P^{k}(x)(1-P(x))^{N-k} \mathrm{~d} x \tag{2.8}
\end{equation*}
$$

and for Poissonian $N$ by conditioning with respect to $N$

$$
\begin{equation*}
E W_{k}=\int(\lambda P(x))^{k} \exp (-\lambda P(x)) \mathrm{d} x / k!. \tag{2.9}
\end{equation*}
$$

In a similar way, for nonrandom $N$

$$
\begin{align*}
& P(\xi(x)=0)=(1-P(x))^{N}  \tag{2.10}\\
& P(\xi(x)=\xi(y)=0)=(1-P(x)-P(y)+P(x, y))^{N} \tag{2.11}
\end{align*}
$$

and for Poissonian $N$

$$
\begin{equation*}
\sigma^{2}\left(V_{1}\right)=\iint \exp (-\lambda P(x)-\lambda P(y))\{\exp (\lambda P(x, y))-1\} \mathrm{d} x \mathrm{~d} y \tag{2.12}
\end{equation*}
$$

Since $V_{k}=W_{k}+W_{k+1}+\ldots$, we have for Poissonian $N$ by (2.9) and the relation

$$
\begin{align*}
& \sum_{j=k}^{\infty} \alpha^{\prime} \mathrm{e}^{-\alpha} / j!=\int_{0}^{\alpha} y^{k-1} \mathrm{e}^{-y} /(k-1)!, \quad k \geqq 1, \\
& E V_{k}=\int\left\{\int_{0}^{P(x)} y^{k-1} \mathrm{e}^{-y} \mathrm{~d} y\right\} \mathrm{d} x /(k-1)!, \quad k \geqq 1 . \tag{2.13}
\end{align*}
$$

Let $w_{k}(R)$ and $v_{k}(R)$ be the contribution to (2.9) and (2.13) from the domain $D_{R}=\{x \mid r \leqq R\}$.

Since $p(u) \geqq c(R)>0$ on $D_{R}$ we have for some $\gamma=\gamma(R)>0$

$$
\begin{equation*}
w_{k}(R)=\mathcal{O}\left(\mathrm{e}^{-\gamma \lambda}\right), \quad v_{k}(R)-\left|D_{R}\right|=\mathcal{O}\left(\mathrm{e}^{-\gamma \lambda}\right), \quad \lambda \rightarrow \infty . \tag{2.14}
\end{equation*}
$$

Now let $\mu$ be the measure on $\mathbb{R}$ induced from Lebesgue measure on $D_{R}^{c}$ by the restriction to $D_{R}^{c}$ of $P$ given by (2.1), i.e.

$$
\begin{equation*}
\mu(A)=\left|\left\{x \in D_{R}^{c}: P(x) \in A\right\}\right| \tag{2.15}
\end{equation*}
$$

Then from (2.13), with $H=\sup \{P(x) \mid r>R\}$,

$$
\begin{align*}
E V_{k} & =v_{k}(R)+\int_{0}^{H}\left(\int_{0}^{\lambda \xi} y^{k-1} \mathrm{e}^{-y} \mathrm{~d} y\right) \mathrm{d} \mu(\xi) /(k-1)! \\
& =v_{k}(R)+\int_{0}^{\lambda H} y^{k-1} \mathrm{e}^{-y} \mu\left[\lambda^{-1} y, H\right] \mathrm{d} y /(k-1)! \tag{2.16}
\end{align*}
$$

We will always take $R$ so large that $P(r, \omega)$ is nonincreasing in $r$ for $r \geqq R$. Let

$$
\begin{equation*}
r_{u}=r_{u}(\omega)=r(u, \omega)=\sup \{r \mid P(r, \omega) \geqq u\} \tag{2.17}
\end{equation*}
$$

Then, since for $0<u \leqq H$ and $m \geqq 2$

$$
\begin{gather*}
\mu[u, H]=m^{-1} \int_{\Omega} r^{m}(u, \omega) \mathrm{d} \sigma(\omega)-\left|D_{R}\right|,  \tag{2.18}\\
E V_{k}=v_{k}^{\prime}(R)+m^{-1} \int_{\Omega} \mathrm{d} \sigma(\omega) \int_{0}^{\lambda H} y^{k-1} \\
\times \mathrm{e}^{-y} r^{m}\left(\lambda^{-1} y, \omega\right) \mathrm{d} y /(k-1)! \tag{2.19}
\end{gather*}
$$

for $m \geqq 2$ where we have by (2.14) for some $\gamma=\gamma(R)>0$

$$
\begin{equation*}
v_{k}^{\prime}(R)=\mathcal{O}(\exp (-\gamma \lambda)), \quad \lambda \rightarrow \infty \tag{2.20}
\end{equation*}
$$

If $P(r, \omega)$ for $r \geqq R$ has a strictly negative continuous partial derivative

$$
\begin{equation*}
P^{\prime}(r)=P^{\prime}(r, \omega)=\partial / \partial r P(r, \omega) \tag{2.21}
\end{equation*}
$$

so that $r(\cdot, \omega)$ defined by $(2.17)$ is the inverse of $P(\cdot, \omega)$, we have from (2.9) for Poissonian $N$ and $m \geqq 2$

$$
\begin{align*}
E W_{k}= & w_{k}(R)+\int_{\Omega} \mathrm{d} \sigma(\omega) \int_{0}^{P(R, \omega)}(\lambda u)^{k} \mathrm{e}^{-\lambda u} r_{u}^{m-1}(\omega) \\
& \times\left|P^{\prime}\left(r_{u}\right)\right|^{-1} \mathrm{~d} u / k! \tag{2.22}
\end{align*}
$$

## 3. Expectation for Poissonian $N$

Theorem 3.1.: Let the $u_{1}$ have probability density (1.4), with $h$ nondecreasing and $h^{-1}$, defined by (1.9), varying regularly at $\infty$ of order $\rho \in\left[0, m^{-1}\right)$. Then as $\lambda \rightarrow \infty$

$$
\begin{align*}
& E V_{k} \sim m^{-1} S_{m}|A|^{m \rho} \Gamma(k-m \rho)\left(h^{-1}(\lambda)\right)^{m} /(k-1)!  \tag{3.1}\\
& E W_{k} \sim \rho S_{m}|A|^{m \rho} \Gamma(k-m \rho)\left(h^{-1}(\lambda)\right)^{m} / k! \tag{3.2}
\end{align*}
$$

Remark: if $m=1$ these relations hold with $S_{1}=2$, see Stam [9]. That $\rho \geqq 0$ follows from the fact that $h$ is nondecreasing. From (1.7) we see that $\rho m \leqq 1$ is necessary in order that $\int p(x) \mathrm{d} x<\infty$, since $h$ varies regularly of order $\rho^{-1}$ if $\rho>0$.

Proof: To prove (3.1) we start from (2.19) where we take $R>d_{A}=$ $\operatorname{diam}(A)$. From (2.1), for $r>R$,

$$
|A| / h\left(r+d_{A}\right) \leqq P(r, \omega) \leqq|A| / h\left(r-d_{A}\right)
$$

Since $P(r, \omega)$ is nonincreasing in $r$ for $r>R$ we have by (1.9) and (2.17), for $0<u<P(R, \omega)$

$$
\begin{equation*}
h^{-1}\left(u^{-1}|A|\right)-d_{A} \leqq r(u, \omega) \leqq h^{-1}\left(u^{-1}|A|\right)+d_{A} . \tag{3.3}
\end{equation*}
$$

By the regular variation of $h^{-1}$ and by (1.8), since $(a+b)^{m}<2^{m}\left(a^{m}+\right.$ $\mathrm{b}^{m}$ ), $a \geqq 0, b \geqq 0$,

$$
\begin{aligned}
& \left\{h^{-1}\left(\lambda y^{-1}|A|\right)+d_{A}\right\}^{m} /\left\{h^{-1}(\lambda)\right\}^{m} \rightarrow y^{-m \rho}|A|^{m \rho}, \quad \lambda \rightarrow \infty \\
& \left\{h^{-1}\left(\lambda y^{-1}|A|\right)+d_{A}\right\}^{m} /\left\{h^{-1}(\lambda)\right\}^{m} \leqq \max \left\{1, C y^{-m \rho-\delta}\right\}
\end{aligned}
$$

with $m \rho+\delta<1$ for $\lambda>\lambda(\delta)$. It follows from (2.19), (2.20) and (3.3) by dominated convergence that

$$
\limsup _{\lambda \rightarrow \infty} E V_{k} /\left\{h^{-1}(\lambda)\right\}^{m} \leqq m^{-1} S_{m}|A|^{m \rho} \Gamma(k-m \rho) /(k-1)!
$$

From the other inequality in (3.3) the reversed inequality for liminf follows and this proves (3.1).

The relation (3.2) follows from (3.1) since $W_{k}=V_{k}-V_{k+1}$.
Theorem 3.1 applies to densities like (1.2) and (1.3) with $g(x) / \log x$ $\rightarrow \infty$. Then we have $\rho=0$ so that (3.2) gives no more than $E W_{k} /\left(H^{-1}(\lambda)\right)^{m} \rightarrow 0$. The principal term for $W_{k}$ then is given by

Theorem 3.2: Let the $u$, have probability density (1.3) where $g$ has $a$ continuous positive derivative $g^{\prime}$ that varies regularly at $\infty$ of order $\tau \geqq-1$, where if $\tau=-1$ we assume that $g^{-1}(\log x)$ varies slowly at $\infty$. Then, as $\lambda \rightarrow \infty$,

$$
\begin{equation*}
E W_{k} \sim(m k)^{-1} S_{m} \xi_{0}(\log \lambda), \quad \xi_{0}(x)=\frac{\mathrm{d}}{\mathrm{~d} x}\left(g^{-1}(x)\right)^{m} \tag{3.4}
\end{equation*}
$$

Proof: We start from (2.22) with $R$ sufficiently large. From (2.1) by applying (1.7) to $g^{\prime}$

$$
\begin{align*}
\left|P^{\prime}(r, \omega)\right| & =\int_{A} p(x-v) g^{\prime}(\|x-v\|)\|x-v\|^{-1}(r-(v, \omega)) \mathrm{d} v \\
& =\int_{A} p(x-v) g^{\prime}(r)(1+\epsilon(r, v, \omega)) \mathrm{d} v \\
& =g^{\prime}(r) P(r, \omega)(1+\epsilon(r, \omega)) \tag{3.5}
\end{align*}
$$

where $\epsilon(r, v, \omega) \rightarrow 0, \epsilon(r, \omega) \rightarrow 0, r \rightarrow \infty$, uniformly in $(v, \omega)$ and $\omega$. For
$r>d_{A}$ with (2.17), noting that $P(r, \omega)$ decrease with $r$ for $r>d_{A}$,

$$
\begin{align*}
& c|A| \exp \left(-g\left(r+d_{A}\right)\right) \leqq P(r, \omega) \leqq c|A| \exp \left(-g\left(r-d_{A}\right)\right) \\
& -d_{A}+g^{-1}\left(\log c|A| u^{-1}\right) \leqq r_{u}(\omega) \leqq g^{-1}\left(\log c|A| u^{-1}\right)+d_{A} \tag{3.6}
\end{align*}
$$

Since $P\left(r_{u}, \omega\right)=u$ we have with (2.14) and with (3.5), noting that as $u \rightarrow \infty r(u, \omega) \rightarrow \infty$ uniformly in $\omega$ by (3.5), for $R=R_{1}(\delta)$,

$$
\begin{equation*}
E W_{k} \leqq(1+\delta) \int_{\Omega} \mathrm{d} \sigma(\omega) \int_{0}^{P(R, \omega)} \lambda^{k} u^{k-1} r_{u}^{m-1} \mathrm{e}^{-\lambda u}\left\{k!g^{\prime}\left(r_{u}\right)\right\}^{-1} \mathrm{~d} u \tag{3.7}
\end{equation*}
$$

If $\tau>-1$ we see from (1.6) that $g$ varies regularly at $\infty$ of order $\tau+1>0$, so that $g^{-1}$ varies regularly of order $(\tau+1)^{-1}$. By (3.6) and (1.7) for $g$

$$
\begin{equation*}
r(u, \omega) \sim \psi(u) \stackrel{\text { def }}{=} g^{-1}\left(\log u^{-1}\right), \quad u \downarrow 0 \tag{3.8}
\end{equation*}
$$

uniformly in $\omega$. If $\tau=1$, this follows from the slow variation of $g^{-1}(\log x)$. From (3.8) and (1.7) for $g^{\prime}$, noting that $\psi(u) \rightarrow \infty$ as $u \downarrow 0$,

$$
\begin{equation*}
g^{\prime}(r(u, \omega)) \sim g^{\prime}(\psi(u)), \quad u \downarrow 0 \tag{3.9}
\end{equation*}
$$

uniformly in $\omega$. From (3.7), (3.8) and (3.9) for $R=R_{2}(\delta)$ and some $a>0$

$$
\begin{equation*}
E W_{k} \leqq(1+\delta)^{2} S_{m} \lambda^{k} \int_{0}^{a} \psi^{m-1}(u) u^{k-1} \mathrm{e}^{-\lambda u}\left\{k!g^{\prime} \circ \psi(u)\right\}^{-1} \mathrm{~d} u \tag{3.10}
\end{equation*}
$$

From (3.8) and the regular variation of $g^{-1}$ if $\tau<1$, or by assumption if $\tau=1$, it follows that $\psi$ varies slowly at 0 and then, e.g. by (1.7), we see that $g^{\prime} o \psi$ varies slowly at 0 . So the integrand in (3.10) varies regularly of order $k$ at 0 . By applying the Abelian theorem given in Feller [3], Ch. XIII.5, we see that the right-hand side of (3.10) is asymptotically equal to

$$
(1+\delta)^{2} k^{-1} \psi^{m-1}\left(\lambda^{-1}\right) / g^{\prime} \circ \psi\left(\lambda^{-1}\right)=(1+\delta)^{2}(m k)^{-1} S_{m} \xi_{0}(\log \lambda)
$$

So, since this holds for all $\delta>0$,

$$
\underset{\lambda \rightarrow \infty}{\limsup } E W_{k} / \xi_{0}(\log \lambda) \leqq(m k)^{-1} S_{m}
$$

From (3.5) and (3.6) we find in the same way the reversed inequality for liminf and this proves (3.4).

If the density $p$ has the special form (1.2) a second-order term in the expansion of $E V_{k}$ and $E W_{k}$ may be found. The proofs start from (2.19)
and (2.22) and proceed by expanding (2.1) and (3.5) into a suitable number of powers of $r$ and a remainder term. From the expansion of $P(r, \omega)$ or from an inequality like (3.6) an expansion for $r(u, \omega)$ with remainder term uniform in $\omega$ is derived. Substitution into (2.19) and (2.22) then gives the asymptotic relations. The proofs are given in Stam [9]. Writing

$$
\begin{aligned}
& R_{h}(\lambda)=E V_{h}-m^{-1} S_{m} \alpha^{m}(\log \lambda)^{p}, \quad p=m / 2 \theta \\
& T_{h}(\lambda)=E W_{h}-S_{m} \alpha^{m}(2 \theta k)^{-1}(\log \lambda)^{p-1}
\end{aligned}
$$

we find for $m \geqq 2$, if $p(u)$ is given by (1.2): If $\theta>\frac{1}{2}$ and $A$ is a convex body and $0 \in \operatorname{Int}(A)$,

$$
\begin{aligned}
& R_{h}(\lambda) \sim M_{A} \alpha^{m-1}(\log \lambda)^{q}, \quad q=(m-1) / 2 \theta \\
& T_{k}(\lambda) \sim(m-1) M_{A} \alpha^{m-1}(2 \theta k)^{-1}(\log \lambda)^{q-1} \\
& M_{A}=\int_{\Omega} \max \{(\omega, x) \mid x \in A\} \operatorname{d} \sigma(\omega)
\end{aligned}
$$

If $m=2$, then $M_{A}$ is equal to the length of the perimeter of $A$ and if $m=3$ it is equal to the integral of the mean curvature of $\partial A$. See Valentine [10], Ch. XII, and Hadwiger [5], §18. If $\theta=\frac{1}{2}$,

$$
\begin{aligned}
& R_{h}(\lambda)=b_{2}(\log \lambda)^{m-1}+\mathcal{O}(\log \lambda)^{m-2} \\
& T_{h}(\lambda)=(m-1) k^{-1} b_{2}(\log \lambda)^{m-2}+\mathcal{O}(\log \lambda)^{m-3} \\
& b_{2}=\alpha^{m}\left\{\int_{\Omega} \log \varphi(\omega) \mathrm{d} \sigma(\omega)-S_{m} \Gamma^{\prime}(k) / \Gamma(k)\right\} \\
& \varphi(\omega)=c \int_{A} \exp \left(-\alpha^{-1}(v, \omega)\right) \mathrm{d} v
\end{aligned}
$$

If $0<\theta<\frac{1}{2}$, with $a=\min \left\{2,(2 \theta)^{-1}\right\}$,

$$
\begin{aligned}
& R_{k}(\lambda)=b_{3}(\log \lambda)^{P-1}+\mathcal{O}(\log \lambda)^{p-a} \\
& T_{k}(\lambda)=(p-1) k^{-1} b_{3}(\log \lambda)^{P-2}+\mathcal{O}(\log \lambda)^{p-1-a} \\
& b_{3}=\alpha^{m} S_{m}(2 \theta)^{-1}\left\{\log c|A|-\Gamma^{\prime}(k) / \Gamma(k)\right\}
\end{aligned}
$$

## 4. Variance for Poissonian $\mathbf{N}$

Theorem 4.1: Let $m \geqq 2$ and let the $u$, have probability density (1.2) with $2 \theta>1$. If $A$ is convex and its boundary $\partial A$ is a twice continuously
differentiable surface in $\mathbb{R}^{m}$ whose principal curvatures $k_{1, \ldots, k_{m-1} \text { satisfy }}$ $0<c_{1} \leqq k_{1}(x) \leqq c_{2}<\infty, i=1, \ldots, m-1, x \in \partial A$, then, with $b_{1}$ given by (1.11),

$$
\begin{align*}
& \frac{1}{2} B_{1} \leqq \limsup _{\lambda \rightarrow \infty}(\log \lambda)^{-b_{1}} \sigma^{2}\left(V_{1}\right) \leqq \limsup _{\lambda \rightarrow \infty}(\log \lambda)^{-b_{1}} \sigma^{2}\left(V_{1}\right) \leqq B_{1}  \tag{4.1}\\
& B_{1}= \alpha^{(3 m+1) / 2} \theta^{-(m+3) / 2} S_{m-1} m!\left(m^{2}-1\right)^{-1}\left(\Gamma\left(\frac{1}{2} m+\frac{1}{2}\right)\right)^{-1} \\
& \times \int_{\partial A} K^{1 / 2}(\xi) \mathrm{d} \mathcal{O}(\xi) \tag{4.2}
\end{align*}
$$

Here $S_{m}$ is given by (1.13), $\mathrm{d} \mathcal{O}(\xi)$ is the surface area element of $\partial A$ at $\xi$ and $K(\xi)$ the gaussian curvature of $\partial A$ at $\xi$, i.e.

$$
\begin{equation*}
K(\xi)=k_{1}(\xi) \ldots k_{m-1}(\xi) \tag{4.3}
\end{equation*}
$$

Remark: If $m=2$ we will have to define $S_{m-1}=2$, as will be seen from the proof. It is conjectured that the inequality (4.1) may be replaced by an actual limit. If $m=1$ and $A=[-a, a]$ we have, see Stam [9],

$$
(\log \lambda)^{2-\theta^{-1}} \sigma^{2}\left(V_{1}\right) \rightarrow\left(\pi \alpha \theta^{-1}\right)^{2} / 6, \quad \lambda \rightarrow \infty
$$

If $p(u)$ decreases with $\|u\|$ faster than exponentially, estimates for $\sigma^{2}\left(V_{1}\right)$ are difficult. This makes the proof of Theorem 4.1 very long. It is given in the appendix.

If $p(u)$ decreases with $\|u\|$ exponentially or slower, we may obtain first-order asymptotic terms for $\sigma^{2}\left(V_{1}\right)$ with the methods of Section 3. We use the notation $\omega=x / r, r=\|x\|$ defined in Section 1 and start from (2.12). Let $U_{1}$ be the contribution to (2.12) for $r<\gamma(\lambda)$ to be chosen later on. Putting $y-x=z$ and noting (2.4) and (2.5) we have, with $P(y) \leqq$ $P(x, y)$,

$$
\begin{align*}
\sigma^{2}\left(V_{1}\right)= & U_{1}+\int_{A-A} \mathrm{~d} z \int_{\Omega} \mathrm{d} \sigma(\omega) \int_{\gamma(\lambda)}^{\infty} \mathrm{d} r r^{m-1} \\
& \cdot \exp \{-\lambda P(x)-\lambda P(x+z)\}\{\exp (\lambda P(x, x+z))-1\}, \tag{4.4}
\end{align*}
$$

$$
\begin{align*}
U_{1} & \leqq \int_{r<\gamma(\lambda)} \mathrm{d} x \exp (-\lambda P(x)) \int\{1-\exp (-\lambda P(y)) \mathrm{d} y \\
& \leqq \lambda \int_{r<\gamma(\lambda)} \mathrm{d} x \exp (-\lambda P(x)) \int P(y) \mathrm{d} y \\
& =\lambda|A| \int_{r<\gamma(\lambda)} \exp (-\lambda P(x)) \mathrm{d} x . \tag{4.5}
\end{align*}
$$

Theorem 4.2: Let $m \geqq 2$ and let the $u$, have probability density (1.2) with $\theta=\frac{1}{2}$. Then as $\lambda \rightarrow \infty$

$$
\begin{align*}
& \sigma^{2}\left(V_{1}\right) \sim \alpha^{m}(\log \lambda)^{m-1} \int_{A-A} \mathrm{~d} z \int_{\Omega} \mathrm{d} \sigma(\omega) \log W(z, \omega)  \tag{4.6}\\
& W(z, \omega)=\left\{\exp \left(-\alpha^{-1}(z, \omega)\right)+1\right\}\left\{\exp \left(-\alpha^{-1}(z, \omega)\right)\right. \\
&\left.+1-\chi\left(A_{z}\right) / \chi(A)\right\}^{-1} \tag{4.7}
\end{align*}
$$

where $A_{z}$ is given by (2.3) and

$$
\begin{equation*}
\chi(E)=\chi(E, \omega)=\int_{E} \exp \left(\alpha^{-1}(v, \omega)\right) \mathrm{d} v, \quad E \in \mathscr{B}\left(\mathbb{R}^{m}\right) \tag{4.8}
\end{equation*}
$$

Remark 1: The integral in (4.6) is finite since $0 \leqq \chi\left(A_{z}\right) \leqq \chi(A)$ and $\exp \left(-\alpha^{-1}(z, \omega)\right) \geqq c_{0}>0$.

Remark 2: If $A$ is a ball with radius $a$, then $\chi(A)$ does not depend on $\omega$ and, see Gröbner und Hofreiter [4], §313, $21^{b}$, and (1.13),

$$
\begin{aligned}
\chi(A) & =(m-1)^{-1} S_{m-1} \int_{-a}^{a} \exp (-x / \alpha)\left(a^{2}-x^{2}\right)^{(m-1) / 2} \mathrm{~d} x \\
& =a^{m}(m-1)^{-1} S_{m-1} \Gamma\left(\frac{1}{2}\right) \Gamma\left(\frac{1}{2} m+\frac{1}{2}\right)(2 \alpha / i a)^{m / 2} J_{m / 2}\left(i a \alpha^{-1}\right) \\
& =a^{m} \pi^{m / 2}(2 \alpha / i a)^{m / 2} J_{m / 2}\left(i a \alpha^{-1}\right)
\end{aligned}
$$

where $J_{\beta}$ is the Bessel function of index $\beta$. Also $\chi\left(A_{z}, \omega\right)=\psi(\rho, \varphi)$ depends only on $\rho=\|z\|$ and the angle $\varphi$ between $z$ and $\omega$. So, since $A-A$ is a ball with radius 2 a and $\varphi$ is distributed uniformly on $[0,2 \pi]$,

$$
\begin{aligned}
\sigma^{2}\left(V_{1}\right) \sim & \alpha^{m}(\log \lambda)^{m-1}(2 \pi)^{-1} S_{m-1}^{2} \int_{0}^{2 a} \rho^{m-1} \mathrm{~d} \rho \\
& \times \int_{0}^{2 \pi} \log W^{*}(\rho, \varphi) \mathrm{d} \varphi
\end{aligned}
$$

$$
\begin{aligned}
W^{*}(\rho, \varphi)= & \left\{\exp \left(-\alpha^{-1} \rho \cos \varphi\right)+1\right\} \\
& \times\left\{\exp \left(-\alpha^{-1} \rho \cos \varphi\right)+1-\psi(\rho, \varphi) / \chi(A)\right\} .
\end{aligned}
$$

Proof of Theorem 4.2: From (2.1), (2.2) and

$$
\begin{align*}
& \|x-v\|=r\left(1-2 r^{-1}(v, \omega)+r^{-2}\|v\|^{2}\right)^{1 / 2}=r-(v, \omega)+\eta r^{-1} \\
& P(x)=c \exp (-r / \alpha) \chi(A, \omega)\left(1+\eta r^{-1}\right)  \tag{4.9}\\
& P(x+z)=c \exp \left(-r \alpha^{-1}-(z, \omega) \alpha^{-1}\right) \chi(A, \omega)\left(1+\eta r^{-1}\right)  \tag{4.10}\\
& P(x, x+z)=c \exp (-r / \alpha) \chi\left(A_{z}, \omega\right)\left(1-\eta r^{-1}\right) \tag{4.11}
\end{align*}
$$

In (4.4) and (4.5) we take $\gamma(\lambda)=\frac{1}{2} \alpha \log \lambda$. From (4.9) for some $c_{1}>0$

$$
\begin{equation*}
U_{1}=\mathcal{O}\left(\lambda(\log \lambda)^{m}\right) \exp \left(-c_{1} \lambda^{1 / 2}\right), \quad \lambda \rightarrow \infty \tag{4.12}
\end{equation*}
$$

Substituting (4.9)-(4.11) into (4.4) and putting $y=\lambda \exp (-r / \alpha)$ we find

$$
\begin{aligned}
\sigma^{2}\left(V_{1}\right)= & U_{1}+\alpha^{m} \int_{A-A} \mathrm{~d} z \int_{\Omega} \mathrm{d} \sigma(\omega) \int_{0}^{\lambda^{1 / 2}} y^{-1} \mathrm{~d} y\left(\log \lambda y^{-1}\right)^{m-1} \\
& \cdot \exp \left\{-\operatorname{cy\chi }(A, \omega)\left(1+\mathrm{e}^{-(z, \omega) / \alpha}\right)\left(1+\eta / \log \lambda y^{-1}\right)\right\} \\
& \cdot\left[-1+\exp \left\{\operatorname{cy\chi }\left(A_{z}, \omega\right)\left(1+\eta / \log \lambda y^{-1}\right)\right\}\right]
\end{aligned}
$$

After dividing by $(\log \lambda)^{m-1}$ we may interchange limit and integral, since $\chi(A, \omega) \geqq b_{0}>0$ and $\lambda y^{-1}>\lambda^{1 / 2}$ in the integral. With (4.12) and Gröbner und Hofreiter [4], $\S 313,3^{b}$, we find (4.6).

Theorem 4.3: Let the $u_{i}$ have probability density (1.3), where $g$ has a positive continuous derivative $g^{\prime}$ that varies regularly at $\infty$ of order $\tau \in[-$ $1,0)$. Let $g^{-1}(\log x)$ vary slowly at $\infty$. Then as $\lambda \rightarrow \infty$, with $A_{z}$ given by (2.3) and $\xi_{0}$ by (3.4),

$$
\begin{equation*}
\sigma^{2}\left(V_{1}\right) \sim m^{-1} S_{m} \xi_{0}(\log \lambda) \int_{A-A} \mathrm{~d} z \log \left(1-\frac{1}{2}\left|A_{z}\right| /|A|\right)^{-1} \tag{4.13}
\end{equation*}
$$

Remark 1: If $m=1$ we have to take $S_{m}=2$. If $\tau>-1$, the regular variation of $g^{\prime}$ implies the slow variation of $g^{-1}(\log x)$, see the proof of (3.8).

Remark 2: If $g(x)=(x / \alpha)^{2 \theta}$, with $0<2 \theta<1$, we have $\xi_{0}(\log \lambda)=$ $m \alpha^{m}(\log \lambda)^{p-1} / 2 \theta, p=m / 2 \theta$. Comparison with (4.1) and (4.6) shows
that the order of the asymptotic behaviour of $\sigma^{2}\left(V_{1}\right)$ depends continuously on $\theta$ if $p(u)$ is given by (1.2), but the coefficient of the leading term does not.

Proof of Theorem 4.3: Since $g^{\prime}(x) \rightarrow 0, x \rightarrow \infty$, we have $f(x+b) / f(x)$ $\rightarrow 1, x \rightarrow \infty$, uniformly with respect to $b$ in compacta. So from (2.1)-(2.3)

$$
\begin{align*}
& P(x)=|A| f(r)\left(1+\epsilon_{1}(x)\right), \quad P(x+z)=|A| f(r)\left(1+\epsilon_{2}(x, z)\right) \\
& P(x, x+z)=\left|A_{z}\right| f(r)\left(1+\epsilon_{3}(x, z)\right. \tag{4.14}
\end{align*}
$$

where $\epsilon_{l} \rightarrow 0, r \rightarrow \infty, i=1,2,3$, uniformly with respect to $\omega$ and $z \in A-A$. In (4.4) we take $\gamma(\lambda)=\gamma$ fixed and so large that $\epsilon_{i} \geqq-\delta, i=1,2, \epsilon_{3} \leqq \delta$ for $r \geqq \gamma$. Then

$$
\begin{align*}
& \sigma^{2}\left(V_{1}\right) \leqq \\
& \quad U_{1}+\int_{A-A} \mathrm{~d} z \int_{\Omega} \mathrm{d} \sigma(\omega) \int_{\gamma}^{\infty} r^{m-1} \mathrm{~d} r  \tag{4.15}\\
&  \tag{4.16}\\
& \cdot \exp (-\lambda \nu f(r))\{-1+\exp (\lambda \beta(z) f(r))\} \\
& \nu=2|A|(1-\delta), \quad \beta(z)=\left|A_{z}\right|(1+\delta)
\end{align*}
$$

The integrand in (4.15) does not depend on $\omega$. We replace the integration variable $r$ by $y=\lambda f(r)$ and find

$$
\begin{align*}
\sigma^{2}\left(V_{1}\right) \leqq & U_{1}+m^{-1} S_{m} \int_{A-A} \mathrm{~d} z \int_{0}^{\lambda f(\gamma)} \xi_{0}\left(\log \lambda c y^{-1}\right) \\
& \times \mathrm{e}^{-\nu y}\left(\mathrm{e}^{\nu \beta(z)}-1\right) y^{-1} \mathrm{~d} y \tag{4.17}
\end{align*}
$$

with $\xi_{0}$ given by (3.4). Here $U_{1} \rightarrow 0$, exponentially as $\lambda \rightarrow \infty$, by (4.5). Since $g^{-1}(\log x)$ varies slowly and $g^{\prime}(x)$ varies regularly at $\infty$, also $\xi_{0}(\log x)$ varies slowly. We want to divide (4.17) by $\xi_{0}(\log \lambda)$ and interchange limit and integral in the right-hand side. In the domain of integration we have $c \lambda y^{-1}>\exp (g(\gamma))$. Since $g(x) \rightarrow \infty, x \rightarrow \infty$, we may take $\gamma$ so large that by applying (1.7) with $\rho=0$ to $\xi_{0}(\log x)$ we find, with $|\epsilon(s)|<\delta_{1}, s \geqq s_{1}$,

$$
\xi_{0}\left(\log \lambda c y^{-1}\right) / \xi_{0}(\log \lambda) \leqq c_{1} \max \left\{y^{-\delta_{1}}, y^{\delta_{1}}\right\}
$$

So, taking into account (4.16), we may apply dominated convergence.

This gives, with Gröbner und Hofreiter [4], §313,3b,

$$
\begin{aligned}
& \limsup _{\lambda \rightarrow \infty} \sigma^{2}\left(V_{1}\right) / \xi_{0}(\log \lambda) \\
& \quad \leqq m^{-1} S_{m} \int_{A-A} \mathrm{~d} z \int_{0}^{\infty} \exp (-\nu y)\{\exp (y \beta(z))-1\} y^{-1} \mathrm{~d} y \\
& \quad=-m^{-1} S_{m} \int_{A-A} \mathrm{~d} z \log \left\{1-\frac{1}{2}(1+\delta)\left|A_{z}\right|(1-\delta)^{-1}|A|^{-1}\right\}
\end{aligned}
$$

A reversed inequality for liminf may be obtained from (4.4) with Fatou's lemma in a similar way by taking $\gamma(\lambda)=\gamma$ so large that in (4.14) we have $\epsilon_{t} \leqq \delta, i=1,2, \epsilon_{3} \geqq-\delta$. Since these inequalities hold for all $\delta>0$, the relation (4.13) follows.

Theorem 4.4: Let the $u$, have probability density (1.4) where $h$ is nondecreasing and varies regularly at $\infty$ of order $\rho^{-1}>m$. Then, with $h^{-1}$ defined by (1.9) and $S_{m}$ by (1.13) and with $S_{l}=2$,

$$
\begin{align*}
& \sigma^{2}\left(V_{1}\right) \sim M_{4}\left(h^{-1}(\lambda)\right)^{m}, \quad \lambda \rightarrow \infty  \tag{4.18}\\
& M_{4}=m^{-1} S_{m} \Gamma(1-m \rho) \int_{A-A}\left\{2|A|^{m \rho}-\left(2|A|-\left|A_{z}\right|\right)^{m \rho}\right\} \mathrm{d} z \tag{4.19}
\end{align*}
$$

Proof: With $\gamma=\gamma(\lambda)$ sufficiently large in (4.4), the relations (4.14)-(4.16) continue to hold since $f(x+b) / f(x) \rightarrow 1, x \rightarrow \infty$, uniformly in $b$, by (1.7) applied to $h$. Define the measure $\mu$ on the Borel sets of $\mathbb{R}_{+}$by

$$
\mu(E)=\int_{B} r^{m-1} \mathrm{~d} r, \quad B=\{x \mid r \geqq \gamma, \quad f(x) \in E\} .
$$

Then (4.15) becomes

$$
\begin{align*}
\sigma^{2}\left(V_{1}\right) & \leqq U_{1}+S_{m} \int_{A-A} G(z, \lambda) \mathrm{d} z  \tag{4.19}\\
G(z, \lambda) & =\int_{0}^{\infty} \mathrm{d} \mu(u) \mathrm{e}^{-\lambda \nu u}\{\exp (\lambda u \beta(z))-1\} \\
& =\int_{0}^{\infty} \mathrm{d} \mu(u) \int_{\nu u-\beta(z) u}^{\nu u} \lambda \exp (-\lambda t) \mathrm{d} t \\
& =\int_{0}^{\infty} \lambda \mathrm{e}^{-\lambda t}\left\{\mu\left[\nu^{-1} t,(\nu-\beta(z))^{-1} t\right)\right\} \mathrm{d} t \tag{4.20}
\end{align*}
$$

For $0<\xi<f(\gamma)$ we have

$$
\mu[\xi, \infty)=m^{-1}\left(h^{-1}\left(\xi^{-1}\right)\right)^{m}-m^{-1} \gamma^{m}
$$

where $h^{-1}\left(\xi^{-1}\right)$ varies regularly at 0 of order $-\rho$. So as $t \downarrow 0$

$$
\mu\left[\nu^{-1} t,(\nu-\beta(z))^{-1} t\right) \sim m^{-1}\left\{\nu^{m \rho}-(\nu-\beta(z))^{m \rho}\right\}\left(h^{-1}\left(t^{-1}\right)\right)^{m}
$$

where the right-hand side varies regularly at 0 of order $-m \rho$ with $-1<-m \rho<0$. By applying the Abelian theorem in Feller [3], Ch. XIII. 5 to (4.20) we see that as $\lambda \rightarrow \infty$,

$$
G(z, \lambda)\left(h^{-1}(\lambda)\right)^{-m} \rightarrow m^{-1} \Gamma(1-m \rho)\left\{\nu^{m \rho}-(\nu-\beta(z))^{m \rho}\right\}
$$

This limit may be interchanged with the integration in (4.19) since with (4.16) and $\left|A_{z}\right| \leqq|A|$

$$
\mu\left[\nu^{-1} t,(\nu-\beta(z))^{-1} t\right) \leqq \mu\left[\nu^{-1} t,|A|^{-1}(1-3 \delta)^{-1} t\right)
$$

where the right-hand side in (4.20) would make applicable the same Abelian theorem, showing that $G(z, \lambda)\left(h^{-1}(\lambda)\right)^{m}$ is bounded by a constant for $\lambda \geqq \lambda_{0}$, so that there is dominated convergence. So from (4.19)

$$
\begin{aligned}
& \limsup _{\lambda \rightarrow \infty} \sigma^{2}\left(V_{1}\right)\left(h^{-1}(\lambda)\right)^{-m} \\
& \quad \leqq m^{-1} S_{m} \Gamma(1-m \rho) \int_{A-A}\left(\nu^{m \rho}-(\nu-\beta(z))^{m \rho}\right) \mathrm{d} z
\end{aligned}
$$

Since this inequality, and a similar reversed one for liminf to be derived by Fatou's lemma, hold for any $\delta>0$, the relation (4.18) follows with (4.16).

## 5. A nonrandom number of sets

We compare corresponding expectations and corresponding variances for a nonrandom number $N=M$ of sets and for a number of sets that is Poissonian with parameter $\lambda=M$. Expectations and variances for the Poissonian number of sets will be denoted by $E_{1}, \sigma_{1}^{2}$ and for the nonrandom number of sets by $E_{2}, \sigma_{2}{ }^{2}$.

Theorem 5.1: We have, for $k=1,2, \ldots$,

$$
\begin{align*}
\left|E_{1} W_{k}-E_{2} W_{k}\right| \leqq & \frac{1}{2} k(k-1) M^{-1} E_{1} W_{k} \\
& +\frac{1}{2} \mathrm{e}(k+1)(k+2) M^{-1} E_{1} W_{k+2} \\
& +k(k+1) \mathrm{e}^{k} M^{-1} E_{1} W_{k+1} \tag{5.1}
\end{align*}
$$

Proof: We write, also for $k=0$,

$$
\begin{align*}
& Q_{k}(x)=(M P(x))^{k} \exp (-M P(x)) / k!,  \tag{5.2}\\
& R_{k}(x)=\binom{M}{k} P^{k}(x)(1-P(x))^{M-k},  \tag{5.3}\\
& Q_{k}(x)-R_{k}(x)=T_{1}(x)+T_{2}(x)+T_{3}(x),  \tag{5.4}\\
& T_{1}(x)=Q_{k}(x)-\binom{M}{k} P^{k}(x) \exp (-M P(x)), \\
& T_{2}(x)=\binom{M}{k} P^{k}(x)\left\{\exp (-M P(x))-(1-P(x))^{M}\right\}, \\
& T_{3}(x)=\binom{M}{k} P^{k}(x)(1-P(x))^{M}-R_{k}(x) .
\end{align*}
$$

We use the inequalities for $x \geqq 0, y \geqq 0$ :

$$
\begin{align*}
1-y \leqq \mathrm{e}^{-y}, \quad 0 & \leqq \mathrm{e}^{-y}-1+y \leqq \frac{1}{2} y^{2}  \tag{5.5}\\
\mathrm{e}^{-M y}-(1-y)^{M} & =\left(\mathrm{e}^{-y}-1+y\right) \sum_{j=0}^{M-1} \mathrm{e}^{-\jmath y}(1-y)^{M-1-\jmath} \\
& \leqq \frac{1}{2} M y^{2} \exp \{-(M-1) y\}  \tag{5.6}\\
(1-x)(1-y) \geqq & 1-x-y . \tag{5.7}
\end{align*}
$$

With (5.7) and (2.9)

$$
\begin{align*}
& 0 \leqq T_{1}(x) \leqq \frac{1}{2} k(k-1) M^{k-1} P^{k}(x) \exp (-M P(x)) / k!, \\
& 0 \leqq \int T_{1}(x) \mathrm{d} x \leqq \frac{1}{2} k(k-1) M^{-1} E_{1} W_{k}, \quad k \leqq 1 . \tag{5.8}
\end{align*}
$$

With (5.6) and (2.9), since $P(x) \leqq 1$,

$$
\begin{align*}
& 0 \leqq T_{2}(x) \leqq \frac{1}{2} \mathrm{e} M^{k+1} P^{k+2}(x) \exp (-M P(x)) / k!, \\
& 0 \leqq \int T_{2}(x) \mathrm{d} x \leqq \frac{1}{2} \mathrm{e}(k+1)(k+2) M^{-1} E_{1} W_{k+2}, \quad k \leqq 0 . \tag{5.9}
\end{align*}
$$

With (5.5), (5.7) and (2.9), since $P(x) \leqq 1$,

$$
\begin{align*}
& 0 \leqq-T_{3}(x) \leqq k \mathrm{e}^{k} M^{k} P^{k+1}(x) \exp (-M P(x)) / k! \\
& 0 \leqq-\int T_{3}(x) \mathrm{d} x \leqq k(k+1) \mathrm{e}^{k} M^{-1} E_{1} W_{k+1}, \quad k \leqq 0 . \tag{5.10}
\end{align*}
$$

The relation (5.1) follows from (2.8), (2.9), (5.2), (5.3), (5.4), (5.8), (5.9), and (5.10). For $k=0$ we have $T_{1}(x)=T_{3}(x)=0$, so with (5.9)

$$
\begin{equation*}
0 \leqq \int\left(Q_{0}(x)-R_{0}(x)\right) \mathrm{d} x \leqq \mathrm{e} M^{-1} E_{1} W_{2} \tag{5.11}
\end{equation*}
$$

Corollary 5.1: For nonrandom $N=M$ theorems 3.1 and 3.2 with $\lambda=M$ continue to hold.

Proof: Since

$$
\sum_{j=0}^{\infty} Q_{k}(x)=\sum_{j=1}^{\infty} R_{h}(x)=1
$$

we have from (2.8), (2.9) and $V_{k}=\sum_{j \geqq k} W_{k}$,

$$
\begin{align*}
E_{1} V_{k}-E_{2} V_{k} & =\sum_{j=0}^{k-1} \int\left(R_{j}(x)-Q_{j}(x)\right) \mathrm{d} x \\
& =\int\left(R_{0}(x)-Q_{0}(x)\right) \mathrm{d} x+\sum_{j=1}^{k-1}\left(E_{2} W_{j}-E_{1} W_{j}\right) . \tag{5.12}
\end{align*}
$$

The corollary follows from (5.1), (5.11) and (5.12) since in theorem 3.1 we have $E W_{k}=\mathcal{O}\left(E V_{j}\right)$ and in Theorem 3.1 and 3.2 all $E W_{k}$ are of the same order.

Remark: The asymptotic relations stated at the end of Section 3 also hold for nonrandom $N$, since they only contain logarithmic terms.

Theorem 5.2: with $\sigma_{l}^{2}\left(V_{l}\right), i=1,2$ as defined above,

$$
\begin{equation*}
\left|\sigma_{1}^{2}\left(V_{1}\right)-\sigma_{2}^{2}\left(V_{1}\right)\right| \leqq e^{2} M^{-1}\left\{4|A-A| E_{1} W_{2}+\left(E_{1} W_{1}\right)^{2}\right\} . \tag{5.13}
\end{equation*}
$$

Proof: From (2.7), (2.10), (2.11), (2.12) and (2.5) we have, writing $U(x, y)=P(x)+P(y)-P(x, y)$ and $L=\{(x, y) \mid y-x \in A-A\}$,

$$
\begin{align*}
& \sigma_{1}^{2}\left(V_{1}\right)-\sigma_{2}^{2}\left(V_{1}\right)=K_{1}+K_{2}+K_{3}  \tag{5.14}\\
& K_{1}=\iint_{L}\left\{\exp (-M U(x, y))-(1-U(x, y))^{M}\right\} \mathrm{d} x \mathrm{~d} y \\
& K_{2}=\iint_{L}\left\{(1-P(x))^{M}(1-P(y))^{M}\right. \\
& \quad-\exp (-M P(x)-M P(y))\} \mathrm{d} x \mathrm{~d} y
\end{align*}
$$

$$
\begin{aligned}
K_{3}= & \iint_{L^{\prime}}\left\{(1-P(x))^{M}(1-P(y))^{M}\right. \\
& \left.-(1-P(x)-P(y))^{M}\right\} \mathrm{d} x \mathrm{~d} y .
\end{aligned}
$$

From (5.6) by symmetry, with $U(x, y) \geqq P(x)$ and by (2.9)

$$
\begin{align*}
0 & \leqq K_{1} \leqq \frac{1}{2} M \iint_{L} U^{2}(x, y) \exp (-(M-1) U(x, y)) \mathrm{d} x \mathrm{~d} y \\
& \leqq M \iint_{L \cap\{P(y)<P(x)\}}(2 P(x))^{2} \exp (-(M-1) P(x)) \mathrm{d} x \mathrm{~d} y \\
& \leqq 4 M|A-A| e \int P^{2}(x) \exp (-M P(x)) \mathrm{d} x=8 M^{-1} \mathrm{e}|A-A| E_{1} W_{2} \tag{5.15}
\end{align*}
$$

In a similar way

$$
\begin{align*}
& 0 \leqq-K_{2} \leqq M \iint_{L}\{\exp (-P(x)-P(y))-1+P(x)+P(y)\} \\
& \cdot \exp \{-(M-1)(P(x)+P(y))\} \mathrm{d} x \mathrm{~d} y \\
& \leqq \frac{1}{2} M e^{2} \iint_{L}(P(x)+P(y))^{2} \exp (-M P(x)-M P(y)) \mathrm{d} x \mathrm{~d} y \\
& \leqq 4 M^{-1} \mathrm{e}^{2}|A-A| E_{1} W_{2}  \tag{5.16}\\
& 0 \leqq K_{3} \leqq M \iint P(x) P(y)(1-P(x))^{M-1}(1-P(y))^{M-1} \mathrm{~d} x \mathrm{~d} y \\
& \leqq M e^{2}\left\{\int P(x) \exp (-M P(x))\right\}^{2}=M^{-1} \mathrm{e}^{2}\left(E_{1} W_{2}\right)^{2} . \tag{5.17}
\end{align*}
$$

The relation (5.13) follows from (5.14)-(5.17).
Corollary 5.2: For nonrandom $N=M$ Theorems 4.1, 4.2, 4.3 and 4.4 with $\lambda=M$ continue to hold.

Proof: From (5.13). In Theorem 4.1 and 4.2 the behaviour of $\sigma_{1}^{2}\left(V_{1}\right)$ is logarithmic, whereas $E_{1} W_{k} \rightarrow 0, \lambda \rightarrow \infty$, by Theorem 3.1. In Theorem 4.3 there is slow variation of $\sigma_{1}^{2}\left(V_{1}\right)$ as $\lambda \rightarrow \infty$ and the same holds for $E_{1} W_{k}$ by Theorem 3.2. In Theorem 4.4 the order of $\sigma_{1}^{2}\left(V_{1}\right)$ and $E_{1} W_{k}$ is the
same by Theorem 3.1 and we have $M^{-1} E_{1} W_{1} \rightarrow 0$ since $\left(h^{-1}(x)\right)^{m}$ varies regularly at $\infty$ of order $m \rho<1$.

## Appendix

Proof of Theorem 4.1: Define $d=d(x)$ and $T=T(x)$ by $d(x)=\|x\|$, $x \in A$, and

$$
\begin{align*}
& d(x)=\inf \{\|x-u\|: u \in A\}, \quad x \notin A,  \tag{1}\\
& T(x) \in \partial A, \quad\|x-T(x)\|=d(x), \quad x \notin A . \tag{2}
\end{align*}
$$

So $d(x)$, if $x \notin A$, is the distance of $x$ to the foot $T(x)$ of the unique normal to $\partial A$, that passes through $x$. For fixed $x \notin A$ we define a Cartesian coordinate system with coordinates $\left(w_{1}, \ldots, w_{m}\right)=w$ and origin $T(x)$. The $w_{m}$-axis is the inner normal to $\partial A$ at $T(x)$ and the other axes are along the principal directions in the tangent plane to $\partial A$ at $T(x)$. See Do Carmo [2], §3.2 and §3.3. In a neighbourhood of $T(x)$ the surface $\partial A$ is described by

$$
\begin{align*}
& 2 w_{m}=\sum_{i=1}^{m-1} k_{i} w_{t}^{2}+\varphi\left(w_{1}, \ldots, w_{m-1}\right)  \tag{3}\\
& t^{-2} \varphi\left(w_{1} \ldots w_{m-1}\right) \rightarrow 0, \quad t \rightarrow 0, \quad t=\left(w_{1}^{2}+\ldots w_{m-1}^{2}\right)^{1 / 2} \tag{4}
\end{align*}
$$

where $k_{t}=k_{t}(T(x))$.
In (2.1) we take $w$ as new integration variable. Note that the $v$ - and $w$-systems have different origins. We have

$$
\begin{align*}
& \|x-v\|^{2}=\left(d(x)+w_{m}\right)^{2}+t^{2}  \tag{5}\\
& P(x)=c \int_{B} \exp \left\{-\alpha^{-2 \theta}\left(\left(d(x)+w_{m}\right)^{2}+t^{2}\right)^{\theta}\right\} \mathrm{d} w \tag{6}
\end{align*}
$$

where $B=B(x)$ is the set $A$ "as seen from $T(x)$ in the $w$-coordinates". For $d=d(x) \geqq d_{1}$ sufficiently large, with the convention about bounded functions stated in section 1, we have

$$
\begin{align*}
& \left(\left(d+w_{m}\right)^{2}+t^{2}\right)^{\theta}=d^{2 \theta}+2 \theta w_{m}\left(1+\eta d^{-1}\right) d^{2 \theta-1}+\eta_{+} t^{2} d^{2 \theta-2}  \tag{7}\\
& P(x)=c \exp \left(-\alpha^{-2 \theta} d^{2 \theta}(x)\right) \int_{0}^{b(x)} \mathrm{d} w_{m} \\
& \quad \times \int_{S\left(w_{m}\right)} \mathrm{d} w_{1} \ldots \mathrm{~d} w_{m-1} \psi_{1}(x, w)  \tag{8}\\
& \psi_{1}(x, w)=\exp \left\{-2 \theta \alpha^{-2 \theta} w_{m}\left(1+\eta d^{-1}\right) d^{2 \theta-1}-\eta_{+} t^{2} d^{2 \theta-2}\right\} \tag{9}
\end{align*}
$$

where $b(x)=\max \left\{w_{m} \mid w \in B(x)\right\}$ and $S(y)$ is the intersection of $B(x)$ with the plane $w_{m}=y$. Let $P_{\epsilon}(x)$ be the contribution to (8) for $0 \leqq w_{m} \leqq \epsilon$. Then $P_{\epsilon}(x) / P(x) \rightarrow 1$ as $d(x) \rightarrow \infty$, uniformly in $T(x)$. In a neighbourhood of $T(x)$ the boundary $\partial A$ lies between the paraboloids

$$
\begin{equation*}
2 w_{m}=(1 \pm \delta) \sum_{i=1}^{m-1} k_{l}(x) w_{i}^{2} \tag{10}
\end{equation*}
$$

so we may take $\epsilon$ so small that for $0 \leqq w_{m} \leqq \epsilon$

$$
\begin{equation*}
S\left(w_{m}\right) \subset\left\{\left(w_{1}, \ldots, w_{m-1}\right) \mid \sum_{i=1}^{m-1} k_{t} w_{t}^{2}<2(1-\delta)^{-1} w_{m}\right\} . \tag{11}
\end{equation*}
$$

In (8) this gives, by putting $u_{t}=w_{1} d^{\theta-1}(x), i=1, \ldots, m-1, u_{m}=$ $w_{m} d^{2 \theta-1}(x)$,

$$
\begin{align*}
& P_{\epsilon}(x) \leqq P_{\epsilon}^{*}(x)=c \exp \left(-\alpha^{-2 \theta} d^{2 \theta}(x)\right)(d(x))^{m-(m+1) \theta} \\
& \quad \cdot \int_{0}^{\epsilon d^{2 \theta-1}} \mathrm{~d} u_{m} \exp \left\{-2 \theta \alpha^{-2 \theta} u_{m}\left(1+\eta d^{-1}\right)\right\} \varphi_{1}(x, u),  \tag{12}\\
& \varphi_{1}(x, u)=\int_{H} \exp \left\{-\eta_{+}\left(u_{1}^{2}+\ldots u_{m-1}^{2}\right)\right\} \mathrm{d} u_{1} \ldots \mathrm{~d} u_{m-1}  \tag{13}\\
& H=\left\{\left(u_{1}, \ldots, u_{m-1}\right) \mid \sum_{i=1}^{m-1} k_{i} u_{i}^{2} \leqq 2(1-\delta)^{-1} u_{m} / d(x)\right\}
\end{align*}
$$

Since $\partial A$ has uniformly continuous derivatives by the compactness of $\bar{A}$ and $0<c_{1} \leqq k_{l}(x) \leqq c_{2}<\infty, i=1, \ldots, m-1, x \in \partial A$, we have, uniformly in $T(x)$, as $d(x) \rightarrow \infty$,

$$
\begin{aligned}
\varphi_{1}(x, u) \sim & (m-1)^{-1} S_{m-1} K^{-1 / 2}(T(x)) \\
& \times\left(2(1-\delta)^{-1} u_{m} / d(x)\right)^{(m-1) / 2}
\end{aligned}
$$

This gives in (12)

$$
\begin{align*}
& P_{\epsilon}^{*}(x) \sim(1-\delta)^{(1-m) / 2} M_{1} K^{-1 / 2}(T(x))(d(x))^{(m+1)(1 / 2-\theta)} \\
& \times \exp \left(-\alpha^{-2 \theta} d^{2 \theta}(x)\right) \\
& M_{1}=c 2^{(m-1) / 2}(m-1)^{-1} S_{m-1} \Gamma\left(\frac{1}{2} m+\frac{1}{2}\right)\left(\alpha^{2 \theta} / 2 \theta\right)^{(m+1) / 2} \tag{14}
\end{align*}
$$

By considering the reversed inclusion with $1+\delta$, analogous to (11),
following from (10), and noting that $P_{\epsilon}(x) / P(x) \rightarrow 1$ we find that as $d(x) \rightarrow \infty$, uniformly in $T(x)$,

$$
\begin{equation*}
P(x) \sim M_{1} K^{-1 / 2}(T(x))(d(x))^{(m+1)(1 / 2-\theta)} \exp \left(-\alpha^{-2 \theta} d^{2 \theta}(x)\right) . \tag{15}
\end{equation*}
$$

We now return to (2.12). By symmetry and with the inequality $\exp (-a)-\exp (-b) \leqq b-a, 0 \leqq a \leqq b$, noting that $P(x, y) \leqq P(y)$, we have, with

$$
\begin{align*}
J(x)= & \{y \mid d(y)>d(x)\},  \tag{16}\\
\sigma^{2}\left(V_{1}\right) & =2 \int \mathrm{~d} x \int_{J(x)} d y \mathrm{e}^{-\lambda P(x)}\left(\mathrm{e}^{-\lambda P(y)+\lambda P(x, y)}-\mathrm{e}^{-\lambda P(y)}\right) \\
& \leqq 2 \lambda \int d x \exp (-\lambda P(x)) \int_{J(x)} d y P(x, y)=H_{1}+H_{2},
\end{align*}
$$

where $H_{1}$ and $H_{2}$ are the contributions to the integral for $x \in D^{c}$ and $x \in D$, respectively, with

$$
\begin{equation*}
D=D(\lambda)=\{x \mid d(x) \geqq q(\lambda)\} . \tag{17}
\end{equation*}
$$

By (15), since $0<a_{1} \leqq K(\xi) \leqq a_{2}, \xi \in \partial A$, we may choose $q(\lambda)$ so that for some $c_{1}>0$ if $\lambda \geqq \lambda_{1}$,

$$
\begin{equation*}
P(x) \leqq \lambda^{-1 / 2}, \quad x \in D(\lambda), \quad P(x) \geqq c_{1} \lambda^{-1 / 2}, \quad x \in D^{c}(\lambda) . \tag{18}
\end{equation*}
$$

With (2.4)

$$
\begin{align*}
H_{1} & \leqq 2 \lambda|A| \int_{D^{\prime}} P(x) \exp (-\lambda P(x)) \mathrm{d} x \\
& \leqq 2 \lambda|A| \exp \left(-c_{1} \lambda^{1 / 2}\right) \int_{D^{\prime}} P(x) \mathrm{d} x \leqq 2 \lambda|A|^{2} \exp \left(-c_{1} \lambda^{1 / 2}\right) . \tag{19}
\end{align*}
$$

We proceed with

$$
\begin{align*}
& \sigma^{2}\left(V_{1}\right) \leqq H_{1}+2 \lambda \int_{D} \exp (-\lambda P(x)) Q(x) \mathrm{d} x,  \tag{20}\\
& Q(x)=\int_{J(x)} P(x, y) \mathrm{d} y . \tag{21}
\end{align*}
$$

Also from (2.12)

$$
\begin{equation*}
\sigma^{2}\left(V_{1}\right) \geqq 2 \lambda \int_{D} \mathrm{~d} x \int_{J(x)} \mathrm{d} y \exp (-\lambda P(x)-\lambda P(y)) P(x, y) . \tag{22}
\end{equation*}
$$

From (2.5) we see that in (22) we may assume $\|y-x\|<2 d_{A}$ so that $T(x)-T(y) \rightarrow 0$ and $K(T(x))-K(T(y)) \rightarrow 0$ as $d(x) \rightarrow \infty$. Since $q(\lambda)$ $\rightarrow \infty$ as $\lambda \rightarrow \infty$ it follows from (17) that in (22)

$$
\begin{equation*}
P(y) \leqq\left(1+\epsilon_{2}(\lambda)\right) P(x), \quad \epsilon_{2}(\lambda) \rightarrow 0, \quad \lambda \rightarrow \infty \tag{23}
\end{equation*}
$$

So that with (21)

$$
\begin{equation*}
\sigma^{2}\left(V_{1}\right) \geqq 2 \lambda \int_{D} \exp \left\{-\lambda\left(2+\epsilon_{2}(\lambda)\right) P(x)\right\} Q(x) \mathrm{d} x . \tag{24}
\end{equation*}
$$

We now have to estimate $Q(x)$. From (21), (2.2) and (2.3)

$$
\begin{aligned}
Q(x) & =\int_{J(x)} \mathrm{d} y \int_{A} I_{A-y+x}(v) p(x-v) \mathrm{d} v \\
& =\int_{A} p(x-v) \mathrm{d} v \int I_{J(x)}(y) I_{A}(v+y-x) \mathrm{d} y \\
& =\int_{A} p(x-v) \mathrm{d} v \int I_{J(x)}(w+x-v) I_{A}(w) \mathrm{d} w \\
& =\int p(x-v)|A \cap(J(x)-x+v)| \mathrm{d} v
\end{aligned}
$$

With the same coordinates and notations as in (4), (5) and (6)

$$
\begin{equation*}
Q(x)=c \int_{B(x)} \exp \left\{-\alpha^{-2 \theta}\left(\left(d(x)+w_{m}\right)^{2}+t^{2}\right)^{\theta}\right\}|F(x, w)| \mathrm{d} w \tag{25}
\end{equation*}
$$

where $F(x, w)$ for fixed $w$ is the set $A \cap(J(x)-x+v)$ "as seen in the new coordinates" (see Fig. 1), where $F(x, w)=F_{1} \cup F_{2}$, with $F_{1}=$ $F(x, w) \cap\left\{y \mid y_{m} \geq w_{m}\right\}$ and $F_{2}=F(x, w) \cap\left\{y \mid y_{m}<w_{m}\right\}$.

With $S(y)$ defined as following (8) and (9) and with the same


Figure 1. $F(x, w)=F_{1} \cup F_{2}$
inclusions as in (10) and (11), denoting by $\left|\left.\right|_{m-1}\right.$ volume in $\mathbb{R}^{m-1}$, we have

$$
\begin{align*}
\left|F_{2}\right| & =\int_{0}^{w_{m}}|S(y)|_{m-1} \mathrm{~d} y \sim \int_{0}^{w_{m}}(m-1)^{-1} S_{m-1} \prod_{i=1}^{m-1}\left(2 y k_{t}^{-1}\right)^{1 / 2} \mathrm{~d} y \\
& =\left(m^{2}-1\right)^{-1} S_{m-1} K^{-1 / 2}(T(x))\left(2 w_{m}\right)^{(m+1) / 2}, \quad w_{m} \rightarrow 0, \tag{26}
\end{align*}
$$

uniformly in $T(x)$. For $d(x) \rightarrow \infty$ the set $J(x)$ approximates to a ball with radius $d(x)$, uniformly in $T(x)$. The projection of $F_{1}$ on the plane $w_{m}=0$ has diameter $\mathcal{O}\left(w_{m}^{1 / 2}\right)$, uniformly in $x$ and we have

$$
\left|F_{1}\right| \leqq c_{3} \int_{0}^{c_{4} w_{m}^{1 / 2}} t^{m} \mathrm{~d} t / \mathrm{d}(x)=c_{5} w_{m}^{(m+1) / 2} / \mathrm{d}(x), \quad 0 \leqq w_{m} \leqq \epsilon .
$$

With (26) this gives

$$
\begin{align*}
& |F(x, w)|=\left(1+\epsilon\left(w_{m}\right)+\eta d^{-1}\right)\left(m^{2}-1\right)^{-1} S_{m-1} K^{-1 / 2}(T(x)) \\
& \quad \times\left(2 w_{m}\right)^{(m+1) / 2} \tag{27}
\end{align*}
$$

with $\epsilon\left(w_{m}\right) \rightarrow 0, w_{m} \rightarrow 0$, uniformly in $x$. We also may prove (27) by considering the paraboloids approximating $\partial A$ in $T(x)$, see (10), and to $\partial F(x, w)$ in $w$ and integrate the volume between them.

We now have to substitute (7) and (27) into (25), write (25) in the
same form as (8) and apply the same technique of two-sided inclusion that derived (15) from (8). We find

$$
\begin{align*}
& Q(x) \sim M_{2}(d(x))^{(m+1)(1-2 \theta)} \exp \left(-\alpha^{-2 \theta} d^{2 \theta}(x)\right) / K(T(x)), \\
& M_{2}=c(m-1)^{-1}\left(m^{2}-1\right)^{-1} 2^{m} m!S_{m-1}^{2}\left(\alpha^{2 \theta} / 2 \theta\right)^{m+1}, \tag{28}
\end{align*}
$$

as $\mathrm{d}(x) \rightarrow \infty$, uniformly in $T(x)$. With (15) this gives, uniformly in $T(x)$,

$$
\begin{equation*}
Q(x) / P(x) \sim M_{3} K^{-1 / 2}(T(x))(\mathrm{d}(x))^{(m+1)(1 / 2-\theta)}, \quad \mathrm{d}(x) \rightarrow \infty \tag{29}
\end{equation*}
$$

$$
\begin{equation*}
M_{3}=S_{m-1} m!\left(\theta^{-1} \alpha^{2} \theta\right)^{(m+1) / 2} /\left\{\left(m^{2}-1\right) \Gamma\left(\frac{1}{2} m+\frac{1}{2}\right)\right\} . \tag{30}
\end{equation*}
$$

In order to apply (15) and (30) to (20) we define a third, curvilinear, coordinate system on $A^{c}$. The new coordinates $(R, \xi)$ and the Cartesian coordinates $x=\left(x_{1}, \ldots, x_{m}\right)$ of a point in $A^{c}$ are connected by the relations

$$
\begin{align*}
& \xi=T(x), \quad R=\mathrm{d}(x),  \tag{31}\\
& x=\xi-R N(\xi) \tag{32}
\end{align*}
$$

with $T(x)$ and $\mathrm{d}(x)$ defined by (1), (2), so that $\xi \in \partial A$, and $N(\xi)$ the unit inward normal to $\partial A$ at $\xi$. So $x(R, \xi)$ is the endpoint of the outer normal to $\partial A$ at $\xi$ with length $R$. To derive the Jacobian of (32) we replace $\xi$ by a set of coordinates $u_{1}, \ldots, u_{m}$ in $\partial A$ in a neighbourhood of $\xi$, i.e. a diffeomorphism $f: U \subset \mathbb{R}^{m-1} \rightarrow V \cap \partial A$, where $V$ is a neighbourhood in $\mathbb{R}^{m}$ of $\xi$. Then (32) becomes

$$
\begin{equation*}
x=f\left(u_{1}, \ldots, u_{m-1}\right)-R N\left(u_{1}, \ldots, u_{m-1}\right), \tag{33}
\end{equation*}
$$

and the Jacobian $J$ of (33) is the determinant of the matrix consisting of the column vectors $-N, f_{u_{1}}-R N_{u_{1}}, \ldots, f_{u_{m-1}}-R N_{u_{m-1}}$, where $f_{u_{1}}$, etc., denotes partial differentiation. So

$$
J=(-1)^{m} R^{m-1}\left(1+\eta R^{-1}\right) J_{1},
$$

where $J_{1}$ is the determinant of the matrix consisting of the column vectors $N, N_{u_{1}}, \ldots, N_{u_{m-1}}$. Now let $\left\{a_{i j}=a_{i j}(u), i, j=1, \ldots, m-1\right\}$ be the matrix describing the differential of the Gauss map at $\xi$ with respect to the basis of the tangent vectors $f_{u,}, i=1, \ldots, m-1$. Then we have, see Do Carmo
[2], §3.3,

$$
N_{u_{1}}=\sum_{r=1}^{m-1} a_{r_{1}} f_{u_{r}}, \quad i, \ldots, m-1
$$

Applying this to $J_{1}$ we find

$$
J=(-1)^{m} R^{m-1}\left(1+\eta R^{-1}\right) \operatorname{Det}\left\{a_{t \jmath}\right\} J_{2},
$$

where $J_{2}$ is the determinant of the matrix consisting of the column vectors $N, f_{u_{1}}, \ldots, f_{u_{n-1}}$.

Now $N$ is a unit vector orthogonal to the $f_{u_{i}}$. So $J_{2}$ is the $(m-1)$ dimensional volume of the parallelopiped spanned by the $f_{u_{1}}$ and $J_{2}$ $\mathrm{d} u_{1} \ldots \mathrm{~d} u_{m-1}$ is equal to the surface area element $\mathrm{d} \mathcal{O}(\xi)$ at $\xi$ of $\partial A$. We have $\operatorname{Det}\left\{a_{t \jmath}\right\}=K(\xi)$, see Do Carmo [2], §3.3. So

$$
\int_{A^{!}} h(x) \mathrm{d} x=\int_{\partial A} \mathrm{~d} \mathcal{O}(\xi) K(\xi) \int_{0}^{\infty} R^{m-1}\left(1+\eta R^{-1}\right) h(R, \xi) \mathrm{d} R .
$$

Writing $P(R, \xi)$ for $P(x)$ we see from (20), (17) and (29)

$$
\begin{align*}
\sigma^{2}\left(V_{1}\right) \leqq & H_{1}+2 \lambda M_{3} \int_{\partial A} \mathrm{~d} \mathscr{O}(\xi) K^{1 / 2}(\xi) Z(\lambda, \xi)  \tag{34}\\
Z(\lambda, \xi)= & \int_{q(\lambda)}^{\infty} R^{m-1+(m+1)(1 / 2-\theta)}\left(1+\epsilon_{3}\right) P(R, \xi) \\
& \times \exp (-\lambda P(R, \xi)) \mathrm{d} R \tag{35}
\end{align*}
$$

with $\epsilon_{3}=\epsilon_{3}(R, \xi) \rightarrow 0$ as $R \rightarrow \infty$, uniformly in $\xi$. In (35) we take $P(R, \xi)$ as new integration variable. From (6)

$$
|\partial P(R, \xi) / \partial R|=2 \theta \alpha^{-2 \theta} R^{2 \theta-1}\left(1+\eta R^{-1}\right) P(R, \xi)
$$

So with (18) and (1.11)

$$
\begin{equation*}
Z(\lambda, \xi) \leqq \alpha^{2 \theta}(2 \theta)^{-1} \int_{0}^{\lambda^{-1 / 2}} R_{u}^{2 \theta b_{1}}\left(1+\epsilon_{4}\right) \mathrm{e}^{-\lambda u} \mathrm{~d} u \tag{36}
\end{equation*}
$$

where $\epsilon_{4}=\epsilon_{4}(u, \epsilon) \rightarrow 0, u \rightarrow 0$, uniformly in $\xi$ and $R_{u}=R(u, \xi)$ is defined by

$$
\begin{equation*}
P\left(R_{u}, \xi\right)=u \tag{37}
\end{equation*}
$$

Note that $P(R, \xi)$ decreases strictly and continuously with $R$ for $R>2 d_{A}$. From (6)

$$
\begin{aligned}
& c|A| \exp \left\{-\alpha^{-2 \theta}\left(\left(R+d_{A}\right)^{2}+d_{A}^{2}\right)^{\theta}\right\} \leqq P(R, \xi) \\
& \quad \leqq c|A| \exp \left\{-\alpha^{-2 \theta}\left(R-d_{A}\right)^{2 \theta}\right\}
\end{aligned}
$$

So from (37)

$$
\begin{aligned}
& -d_{A}+\alpha\left\{\left(\log c|A| u^{-1}\right)^{\frac{1}{\theta}}-\alpha^{-2} d_{A}^{2}\right\}^{1 / 2} \leqq R(u, \xi) \\
& \quad \leqq d_{A}+\alpha\left(\log c|A| u^{-1}\right)^{1 / 2 \theta}, \\
& R(u, \xi) \sim \alpha\left(\log u^{-1}\right)^{1 / 2 \theta}, \quad u \downarrow 0,
\end{aligned}
$$

uniformly in $\xi$. Substitution into (36) and putting $\lambda u=y$ gives

$$
Z(\lambda, \xi) \leqq \alpha^{2 \theta\left(1+b_{1}\right)}(2 \theta \lambda)^{-1} \int_{0}^{\lambda^{1 / 2}}\left(\log \lambda y^{-1}\right)^{b_{1}}\left(1+\epsilon_{5}\right) \mathrm{e}^{-v} \mathrm{~d} y
$$

where $\epsilon_{5}=\epsilon_{5}\left(\lambda y^{-1}, \xi\right)$ and $\epsilon_{5}(t, \epsilon) \rightarrow 0$ as $t \rightarrow \infty$, uniformly in $\xi$. Substituting into (34) and dividing by $(\log \lambda)^{b_{1}}$ we find the right-hand inequality of (4.1) by (19) and dominated convergence. Note that $\lambda y^{-1} \geqq$ $\lambda^{1 / 2}$ in the domain of integration. This disposes of the case $b_{1}<0$. If $b_{1}>0$ we apply the inequality

$$
\left|\log \lambda y^{-1} / \log \lambda\right| \leqq 1+c_{3}|\log y|, \quad \lambda \geqq \lambda_{2}
$$

From (24) we derive the left-hand inequality of (4.1) in a similar way, using Fatou's lemma, the factor $\frac{1}{2}$ arising from the exponent $-\lambda(2+$ $\left.\epsilon_{2}(\lambda)\right) P(x)$ instead of $-\lambda P(x)$.

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