# POTENTIAL KERNEL FOR TWO-DIMENSIONAL RANDOM WALK 

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#### Abstract

It is proved that the potential kernel of a recurrent, aperiodic random walk on the integer lattice $\mathbb{Z}^{2}$ admits an asymptotic expansion of the form


$$
(2 \pi \sqrt{|Q|})^{-1} \ln Q\left(x_{2},-x_{1}\right)+\mathrm{const}+|x|^{-1} U_{1}\left(\omega^{x}\right)+|x|^{-2} U_{2}\left(\omega^{x}\right)+\cdots,
$$

where $|Q|$ and $Q(\theta)$ are, respectively, the determinant and the quadratic form of the covariance matrix of the increment $X$ of the random walk, $\omega^{x}=x /|x|$ and the $U_{k}(\omega)$ are smooth functions of $\omega,|\omega|=1$, provided that all the moments of $X$ are finite. Explicit forms of $U_{1}$ and $U_{2}$ are given in terms of the moments of $X$.

1. Introduction and statements of results. Let $X^{(1)}, X^{(2)}, \ldots$ be a sequence of $\mathbb{Z}^{2}$-valued i.i.d. mean-zero random variables with finite variance and $\left\{S_{n}\right\}_{n=0}^{\infty}$ the associated random walk on the integer lattice $\mathbb{Z}^{2}$ starting at the origin; that is, $S_{0}=0, S_{n}=\sum_{i=1}^{n} X^{(i)}$. We write $X$ for $X^{(1)}$ for brevity. We assume that the random walk $\left\{S_{n}\right\}_{n=0}^{\infty}$ is aperiodic (i.e., the smallest additive subgroup containing $\left\{x \in \mathbb{Z}^{2}: P\{X=x\}>0\right\}$ agrees with $\mathbb{Z}^{2}$ ). As in [3], we define the potential function (Green function) $a(x)$ by

$$
a(x)=\sum_{n=0}^{\infty}\left(P\left\{S_{n}=0\right\}-P\left\{S_{n}=-x\right\}\right), \quad x \in \mathbb{Z}^{2}
$$

Let $Q(\theta)$ be the moment quadratic form of $X$. That is, $Q(\theta)=E\left\{(\theta \cdot X)^{2}\right\}$, $\theta=\left(\theta_{1}, \theta_{2}\right) \in \mathbb{R}^{2}$. We sometimes write $Q\left(\theta_{1}, \theta_{2}\right)$ for $Q(\theta)$. Let $Q$ also denote the covariance matrix of $X$, and $Q^{-1}$ its inverse matrix, and define

$$
\|x\|:=\sqrt{x \cdot Q^{-1} x}=\sqrt{Q\left(x_{2},-x_{1}\right) / \operatorname{det} Q}
$$

[here $x=\left(x_{1}, x_{2}\right)$ is thought to be a column vector when the matrix is operated from the left]. The square root of $Q$ that is symmetric and positive definite is denoted by $\sqrt{Q}$. We need the moment conditions

$$
(\mathrm{MC}: k+\delta) \quad E\left\{|X|^{k+\delta}\right\}<\infty \quad \text { for some } \delta>0
$$

where $k$ will take the values $2,3, \ldots$. The following result is due to Spitzer [3].

[^0]Theorem 1. Suppose the moment condition (MC: $2+\delta$ ) holds. Then

$$
\begin{equation*}
\lim _{|x| \rightarrow \infty}\left[a(x)-\left(\pi \sigma_{1} \sigma_{2}\right)^{-1} \ln \|x\|-C_{0}\right]=0 \tag{1.1}
\end{equation*}
$$

where $C_{0}$ is a certain constant that depends on the distribution of $X$, and $\sigma_{1}$ and $\sigma_{2}$ are the eigenvalues of the matrix $\sqrt{Q}$.

To state the main result of this paper, we put

$$
\psi_{i, 1}(\theta):=-2 E\left\{(\theta \cdot X)^{3}\right\} / 3 Q^{2}(\theta)
$$

and

$$
\psi_{r, 1}(\theta):=\left[\frac{1}{2} Q(\theta) E\left\{(\theta \cdot X)^{4}\right\}-\frac{2}{3}\left(E\left\{(\theta \cdot X)^{3}\right\}\right)^{2}\right] / 3 Q^{3}(\theta) .
$$

These functions are the principal parts of the real and imaginary parts of the function $\psi(\theta):=\left(1-E\left\{e^{i X \cdot \theta}\right\}\right)^{-1}-2 / Q(\theta)$ [cf. (3.4) and (3.5)]. We also define

$$
\begin{equation*}
g^{\sharp}(\theta):=\text { p.v. } \int_{-\infty}^{\infty} g\left(\theta_{1}-t \theta_{2}, \theta_{2}+t \theta_{1}\right) d t \tag{1.2}
\end{equation*}
$$

for $\theta \in \mathbb{R} \backslash\{0\}$ and a function $g$ for which the principal value on the right-hand side exists. [The principal value here is, of course, the limit of the integral on the symmetrical interval $(-L, L)$.]

Theorem 2. If the moment condition (MC: $2+m+\delta$ ) holds ( $m \geq 1$ ), then

$$
\begin{equation*}
a(x)-\frac{1}{\pi \sigma_{1} \sigma_{2}} \ln \|x\|-C_{0}=\frac{U_{1}\left(\omega^{x}\right)}{|x|}+\cdots+\frac{U_{m}\left(\omega^{x}\right)}{|x|^{m}}+o\left(\frac{1}{|x|^{m}}\right) \tag{1.3}
\end{equation*}
$$

as $|x| \rightarrow \infty$ in $\mathbb{Z}^{2}$, where $\omega^{x}=x /|x|, \sigma_{1}, \sigma_{2}$ and $C_{0}$ are the same constants as in (1.1) and $U_{k}, k=1,2, \ldots$, are smooth functions of $\omega=\left(\omega_{1}, \omega_{2}\right),|\omega|=1$; moreover, the first and the second of them are given by

$$
U_{1}(\omega)=\frac{1}{2 \pi^{2}} \psi_{i, 1}^{\sharp}(\omega) \quad \text { and } \quad U_{2}(\omega)=\frac{1}{2 \pi^{2}}\left(\omega \cdot \nabla \psi_{r, 1}\right)^{\sharp}(\omega) \text {. }
$$

Remark 1. The function $U_{1}(\omega)$ is identically 0 if and only if all the third moments $E\left\{\left(X_{1}\right)^{k}\left(X_{2}\right)^{3-k}\right\}, k=0,1,2,3$, vanish. If $X$ is symmetric, that is, $X$ has the same distribution as $-X$, then $U_{k}$ vanishes for every odd number $k$.

Remark 2. For the simple random walk in particular, Theorem 2 gives the asymptotic expansion

$$
a(x)=\frac{2}{\pi} \ln |x|+\frac{\ln 8+2 \gamma}{\pi}+\frac{1}{6 \pi} \frac{8\left(\omega_{1}^{x} \omega_{2}^{x}\right)^{2}-1}{|x|^{2}}+\frac{U_{4}\left(\omega^{x}\right)}{|x|^{4}}+\cdots
$$

( $\gamma$ is Euler's constant), which is an improvement of a result of Stöhr [4], where $a(x)$ is computed up to $O\left(|x|^{-2}\right)$.

Remark 3. We shall see in Section 5 that $U_{k}(\theta /|\theta|) /|\theta|^{k}$ is a rational fraction of the form $\left\{\theta^{2 \nu-k}\right\} /\|\theta\|^{2 \nu}$, where $\left\{\theta^{k}\right\}$ represents a homogeneous polynomial of degree $k$. Accordingly, (1.3) may be rewritten as

$$
\begin{equation*}
a(x)-\frac{1}{\pi \sigma_{1} \sigma_{2}} \ln \|x\|-C_{0}=\frac{\tilde{U}_{1}\left(\tilde{\omega}^{x}\right)}{\|x\|}+\cdots+\frac{\tilde{U}_{m}\left(\tilde{\omega}^{x}\right)}{\|x\|^{m}}+o\left(\frac{1}{|x|^{m}}\right) . \tag{1.4}
\end{equation*}
$$

Here $\tilde{\omega}^{x}:=\sqrt{Q^{-1}} x /\|x\| ; \tilde{U}_{k}(\omega)$ is a polynomial of $\omega=\left(\omega_{1}, \omega_{2}\right)$ of degree (at most) $3 k$ for $k=1,2, \ldots$; in particular, $\tilde{U}_{1}(\omega)=\left(\psi_{i, 1}{ }^{\circ} \sqrt{Q^{-1}}\right)^{\sharp}(\omega) /\left(2 \pi^{2} \sigma_{1} \sigma_{2}\right)$. [See (1.8) below.]

In the case of a simple random walk, $a(x)$ can be neatly expressed by a contour integral on the complex plane (as given and applied, e.g., in [4] and [5]) and the complex function theory accordingly provides us machinery for computation, though the proof given in [4] is still quite involved.

In our approach, we employ only real analytic arguments as in Spitzer [3] and it is a key step to establish an asymptotic expansion of an integral of the form

$$
\begin{equation*}
\int_{[-\pi, \pi] \times[-\pi, \pi]} \frac{p(\theta)}{q(\theta)} \sin (x \cdot \theta) d \theta \tag{1.5}
\end{equation*}
$$

as $x \rightarrow \infty$, where $p(\theta)$ and $q(\theta)$ are homogeneous polynomials of degree $2 \nu-1$ and $2 \nu, \nu \geq 1$, respectively, and $q(\theta)$ is supposed to be positive for $\theta \neq 0$, so that

$$
q(\theta) \geq c|\theta|^{2 \nu}, \quad \theta \in \mathbb{R}^{2}
$$

for a constant $c>0$. We formulate the result on the integral (1.5) in the following theorem.

Theorem 3. Let $p$ and $q$ be as above. Let $D$ be a two-dimensional bounded domain containing the origin and having piecewise smooth boundary. Let $m$ be a positive integer and $\xi(\theta)$ a function on the closure $\bar{D}$ such that $\xi$ has partial derivatives up to order $m$ that are continuous on $\bar{D} \backslash\{0\}$ and integrable on $D$. Then for the function

$$
g(\theta)=\frac{p(\theta)}{q(\theta)}+\xi(\theta)
$$

it holds that

$$
\begin{equation*}
\int_{D} g(\theta) e^{i r \omega \cdot \theta} d \theta=-\frac{2}{i r}\left(\frac{p}{q}\right)^{\#}(\omega)+\sum_{l=1}^{m} \frac{1}{(i r)^{l}} B_{l}(r, \omega)+o\left(\frac{1}{r^{m}}\right) \tag{1.6}
\end{equation*}
$$

as $r \uparrow \infty$ uniformly for $\omega,|\omega|=1$, where $\#$ is defined by (1.2) and

$$
\begin{equation*}
B_{l}(r, \omega):=\int_{\partial D}(-\omega \cdot \nabla)^{l-1} g(\theta) e^{i r \omega \cdot \theta} \omega \cdot n d s . \tag{1.7}
\end{equation*}
$$

[ Here, $n=n(\theta)$ is the outward unit normal vector to $\partial D$ and $d s=d s_{\theta}$ is a line element of $\partial D$.]

REMARK 4. As a consequence of Theorem 3, we obtain a useful property of the transform $p / q \rightarrow(p / q)^{\sharp}$. Let $A$ be a regular $2 \times 2$ matrix. Then

$$
\begin{equation*}
\left(\frac{p}{q}\right)^{\#}(\theta)=|\operatorname{det} A|\left(\left(\frac{p}{q}\right) \circ A\right)^{\#}\left(A^{t} \theta\right) . \tag{1.8}
\end{equation*}
$$

( $A^{t}$ denotes the transpose of the matrix $A$.) The verification of (1.8) is immediate from (1.6) [take $\xi$ so that both $g$ and $\nabla \cdot g$ vanish on $\partial D$, change the variable according to $\theta=A \theta^{\prime}$ on the left-hand side of (1.6) and let $r \rightarrow \infty$ ], while it is not so simple a matter to establish (1.8) if one only looks at the defining expression (1.2). (See Section 5 for further properties of the transform \#.)

The following version of Theorem 3 is convenient for application.
Remark 5. Theorem 3 may be extended to a more complete form. Let $D$ and $\xi$ be as in Theorem 3. Let $\psi(\omega)$ be a smooth function on the unit circle. Then for a function $g$ of the form

$$
g(\theta)=h(\theta)+\xi(\theta) \text { with } h(\theta)=\frac{\psi(\theta /|\theta|)(\theta)}{|\theta|}
$$

it holds that

$$
\int_{D} g(\theta) e^{i r \omega \cdot \theta} d \theta=2 \pi h^{b}(\theta)+i 2 h^{\sharp}(\theta)+\sum_{k=1}^{m} \frac{1}{(i r)^{k}} B_{k}(r, \omega)+o\left(\frac{1}{r^{m}}\right)
$$

as $r \uparrow \infty$, uniformly in $\omega,|\omega|=1$. Here $B_{k}$ is the same as in (1.7) and

$$
h^{b}(\theta)=\frac{1}{2}\left[h\left(\theta_{2},-\theta_{1}\right)+h\left(-\theta_{2}, \theta_{1}\right)\right]
$$

The method developed in this paper can be adapted for deriving the asymptotic expansion of the potential kernel for the higher-dimensional random walk, which will be studied in a separate paper. Theorem 3, in particular, has a $d$-dimensional version $(d \geq 3)$ in which $g^{\sharp}$ takes an analogous or different form according as $d$ is even or odd.

The result of Stöhr [4] mentioned above is used for estimating a certain hitting distribution by Kesten [1] (cf. also [2]). As another example of application of our expansion (1.3), we shall compute in Section 6 the distributions of hitting places of lines $x_{2}=N$ up to $O\left(|x|^{-3}\right)$.

The proof of Theorem 1, which is essentially the same as in [3], prepares that of Theorem 2, and our task for the latter is to get the estimate of the remainder term, which will be reduced to Theorem 3 with not much difficulty.

We shall proceed in logical order, namely, we first prove Theorem 1 in Section 2, secondly Theorem 3 in Section 3 and then Theorem 2 in Section 4.
2. Proof of Theorem 1. Here we outline the proof of Theorem 1. It is identical to that given in Proposition 12.3 of Spitzer [3] [where $Q(\theta)$ is assumed to be a constant multiple of $|\theta|^{2}$ ] except for a simple modification by
a change of variable, but we need the content of it, since our proof of Theorem 2 being a continuation of it.

Let $\phi(\theta)$ be the characteristic function of $X$, that is, $\phi(\theta)=E\left\{e^{i X \cdot \theta}\right\}$. The function $a(x)$ is expressed as follows:

$$
\begin{equation*}
a(x)=\frac{1}{(2 \pi)^{2}} \int_{T} \frac{1-e^{i x \cdot \theta}}{1-\phi(\theta)} d \theta, \quad x \in \mathbb{Z}^{2} \tag{2.1}
\end{equation*}
$$

where $T=[-\pi, \pi] \times[-\pi, \pi]$. Introducing

$$
\psi(\theta):=\frac{1}{1-\phi(\theta)}-\frac{2}{Q(\theta)},
$$

which is integrable on $T$ since the condition (MC: $2+\delta$ ) implies $1-$ $\phi(\theta)=\frac{1}{2} Q(\theta)+O\left(|\theta|^{2+\delta}\right)$ as $|\theta| \rightarrow 0$ (cf. [3], Proposition 12.3). We make the decomposition

$$
\begin{equation*}
4 \pi^{2} a(x)=\int_{T} \frac{2}{Q(\theta)}(1-\cos x \cdot \theta) d \theta+\operatorname{Re} \int_{T} \psi(\theta)\left(1-e^{i x \cdot \theta}\right) d \theta \tag{2.2}
\end{equation*}
$$

( $\operatorname{Re} z$ indicates the real part of a complex number $z$ ). In view of the Rie-mann-Lebesgue lemma, the second term converges, as $|x| \rightarrow \infty$, to $\operatorname{Re} \int_{T} \psi(\theta) d \theta$, contributing to the constant $C_{0}$ and leaving the $o(1)$ term

$$
\begin{equation*}
\Psi(x):=-\operatorname{Re} \int_{T} \psi(\theta) e^{i x \cdot \theta} d \theta \tag{2.3}
\end{equation*}
$$

For the evaluation of the first term, we consider the mapping

$$
\begin{equation*}
\theta \mapsto \theta^{\prime}=\sqrt{Q} \theta, \tag{2.4}
\end{equation*}
$$

which, entailing the identity $Q(\theta)=\left|\theta^{\prime}\right|^{2}$, transforms the ellipse $Q(\theta)=r^{2}$ into the circle $\left|\theta^{\prime}\right|=r$. Let $(r, \alpha)$ be the polar coordinates of $\theta^{\prime}$, namely,

$$
\begin{equation*}
\theta^{\prime}=\left(r \cos \left(\alpha-\alpha_{0}\right), r \sin \left(\alpha-\alpha_{0}\right)\right), \tag{2.5}
\end{equation*}
$$

where $\alpha_{0}$ is a constant chosen arbitrarily. Since $x \cdot \theta=\sqrt{Q^{-1}} x \cdot \theta^{\prime}$ and $\|x\|=\mid \sqrt{Q^{-1}} x$, we can choose the constant $\alpha_{0}=\alpha_{0}(x)$ so as to get $x \cdot \theta=$ $\|x\| r \sin \alpha$.

Now, putting

$$
B=\left\{\theta: \sqrt{Q(\theta)} \leq\left(\sigma_{1} \wedge \sigma_{2}\right) \pi\right\}
$$

( $B$ is an elliptic region inscribed in $T$ ), we decompose the first integral on the right-hand side of (2.2) into two parts, one the integral over $B$ and the other that over $T \backslash B$. The latter converges to $\int_{T \backslash B}[2 / Q(\theta)] d \theta$, leaving the second $o$ (1) term

$$
\begin{equation*}
-\int_{T \backslash B} \frac{2}{Q(\theta)} \cos (x \cdot \theta) d \theta . \tag{2.6}
\end{equation*}
$$

The former equals

$$
\begin{align*}
& \int_{0}^{2 \pi} d \alpha \int_{0}^{\left(\sigma_{1} \wedge \sigma_{2}\right) \pi} \frac{2}{\sigma_{1} \sigma_{2} r}[1-\cos (\|x\| r \sin \alpha)] d r \\
& \quad=4 \int_{0}^{\pi / 2} d \alpha \int_{0}^{c(x) \sin \alpha} \frac{2}{\sigma_{1} \sigma_{2} u}(1-\cos u) d u \\
& \quad=\frac{8}{\sigma_{1} \sigma_{2}}\left[\int_{0}^{\pi / 2} d \alpha \int_{0}^{1} \frac{1-\cos u}{u} d u-\int_{0}^{\pi / 2} d \alpha \int_{1}^{\infty} \frac{\cos u}{u} d u\right.  \tag{2.7}\\
& \left.\quad+\int_{0}^{\pi / 2} d \alpha \int_{1}^{c(x) \sin \alpha} \frac{1}{u} d u+\int_{0}^{\pi / 2} d \alpha \int_{c(x) \sin \alpha}^{\infty} \frac{\cos u}{u} d u\right]
\end{align*}
$$

where $c(x)=\left(\sigma_{1} \wedge \sigma_{2}\right) \pi\|x\|$. Within the brackets of the last expression in (2.7), the first and the second terms are constants; the third is equal to

$$
\frac{1}{2} \pi \log \|x\|+\frac{1}{2} \pi \log \left[\left(\sigma_{1} \wedge \sigma_{2}\right) \pi\right]+\int_{0}^{\pi / 2} \log (\sin \alpha) d \alpha
$$

and the fourth vanishes as $|x| \rightarrow \infty$ [this gives the last $o(1)$ term]. The proof of Theorem 1 is complete.
3. Proof of Theorem 3. Put

$$
\omega=\frac{x}{|x|} \quad \text { and } \quad R_{\omega}=\left(\begin{array}{cc}
\omega_{1} & -\omega_{2} \\
\omega_{2} & \omega_{1}
\end{array}\right)
$$

and introduce the new variables

$$
(u, v)=R_{\omega}^{-1} \theta
$$

so that $\omega \cdot \theta=u$. Let

$$
D^{\omega}=R_{\omega}^{-1} D \quad \text { and } \quad g^{\omega}(u, v)=g\left(R_{\omega}(u, v)\right) .
$$

Then

$$
\begin{equation*}
\int_{D} g(\theta) e^{i r \omega \cdot \theta} d \theta=\int_{D^{\omega}} g^{\omega}(u, v) e^{i r u} d u d v \tag{3.1}
\end{equation*}
$$

We are to carry out the integration by parts for the integral with respect to $u$ on the right-hand side above. This amounts to applying the divergence theorem to the integral on the left-hand side, into which we substitute

$$
\begin{equation*}
e^{i r \omega \cdot \theta}=\nabla \cdot \mathbf{A}(\theta), \quad \text { where } \mathbf{A}(\theta)=(\text { ir })^{-1} e^{i r \omega \cdot \theta} \omega . \tag{3.2}
\end{equation*}
$$

We can apply the divergence theorem repeatedly $m$ times to the integral of $\xi(\theta) e^{i r \omega \cdot \theta}=\nabla \cdot(\xi \mathbf{A})-\nabla \xi \cdot \mathbf{A}$, which results in the boundary integrals given in (1.7) with $\xi$ in place of $g$ plus the remainder term of the order $o\left(r^{-m}\right)$. It therefore suffices to prove (1.6) in the case when $\xi=0$. [Although $\xi$ may be singular at the origin, it can be approximated in Sobolev norm by a smooth function (under the assumption on $\xi$ in Theorem 3) so that the divergence theorem is applicable at least $m$ times.] Let $g=p / q$ with polynomials $p$ and $q$ as described in Theorem 3.

We suppose for simplicity that $D^{\omega}$ contains the square

$$
K=\{(u, v):|u| \leq 1,|v| \leq 1\}
$$

and we decompose the integral on the right-hand side of (3.1) into that over $K$ and the rest. We formulate the result of the computation as follows.

LEMMA 1. Let $g=p / q$, where $p$ and $q$ are homogeneous polynomials as given in Theorem 3. Then

$$
\frac{d}{d u} \int_{-1}^{1} g^{\omega}(u, v) d v=\int_{-1}^{1} \frac{\partial g^{\omega}}{\partial u}(u, v) d v=\frac{f(u, \omega)}{q^{\omega}(u, 1) q^{\omega}(u,-1)}, \quad u \neq 0
$$

where $q^{\omega}=q \circ R_{\omega}$ and $f(u, \omega)$ is a polynomial of $\left(u, \omega_{1}, \omega_{2}\right)$, and

$$
\begin{aligned}
\int_{D} g(\theta) e^{i r \omega \cdot \theta} d \theta= & -\frac{2}{i r} g^{\sharp}(\omega)+\frac{1}{i r} \int_{\partial D} g(\theta) e^{i r \omega \cdot \theta} \omega \cdot n d s \\
& -\frac{1}{i r} \int_{D^{\omega} \backslash K} \frac{\partial g^{\omega}}{\partial u}(u, v) e^{i r u} d u d v \\
& -\frac{1}{i r} \int_{-1}^{1} e^{i r u} d u \int_{-1}^{1} \frac{\partial g^{\omega}}{\partial u}(u, v) d v
\end{aligned}
$$

Theorem 3 readily follows from Lemma 1 . In fact, if we apply the integra-tion-by-parts formula to the integral relative to $u$ in the last two integrals on the right-hand side above, the contributions of the boundary terms that thereby come up are reduced to

$$
\begin{equation*}
\frac{1}{(i r)^{2}} \int_{\partial D}(-\omega \cdot \nabla) g(\theta) e^{i r \omega \cdot \theta} \omega \cdot n d s \tag{3.3}
\end{equation*}
$$

because of cancellation between those from $\partial K$ [recall the remark made when $\mathbf{A}$ is introduced in (3.2)]. We can repeat the integration by parts in the same way in view of the first half of Lemma 1 to arrive at (1.6). Now it remains to prove Lemma 1.

Proof of Lemma 1. By the divergence theorem

$$
\begin{align*}
\int_{D^{\omega} \backslash K} g^{\omega}(u, v) e^{i r u} d u d v= & \frac{1}{i r} \int_{\partial\left(D^{\omega} \backslash K\right)} g^{\omega}(u, v) e^{i r u} n d s \\
& -\frac{1}{i r} \int_{D^{\omega} \backslash K} \frac{\partial g^{\omega}}{\partial u}(u, v) e^{i r u} d u d v . \tag{3.4}
\end{align*}
$$

We cannot apply the divergence theorem to the integral over $K$ directly. We consider the function

$$
F(u):=\int_{-1}^{1} g^{\omega}(u, v) d v=\int_{-1 / u}^{1 / u} g^{\omega}(1, t) d t
$$

Here, to obtain the second expression, we have applied the assumption that $g$ is in the special form $p / q$ [which the function $g^{\omega}=g \circ R_{\omega}$ clearly inherits, so that $\operatorname{vg}^{\omega}(u, v)$ is the ratio of two homogeneous polynomials of degree $2 \nu$ ]. Clearly, $F(-u)=-F(u)$. Although $g^{\omega}(1, t)$ is not integrable on $\mathbf{R}$, there exists the principal value

$$
F(0+)=\text { p.v. } \int_{-\infty}^{\infty} g^{\omega}(1, t) d t:=\lim _{L \rightarrow \infty} \int_{-L}^{L} g^{\omega}(1, t) d t .
$$

Furthermore, for $u \neq 0$,

$$
\begin{aligned}
F^{\prime}(u) & =-u^{-2}\left[g^{\omega}\left(1, \frac{1}{u}\right)+g^{\omega}\left(1, \frac{-1}{u}\right)\right] \\
& =-u^{-1}\left[g^{\omega}(u, 1)+g^{\omega}(u,-1)\right] \\
& =-\frac{1}{u} \frac{p^{\omega}(u, 1) q^{\omega}(u,-1)+p^{\omega}(u,-1) q^{\omega}(u, 1)}{q^{\omega}(u, 1) q^{\omega}(u,-1)} .
\end{aligned}
$$

Because of cancellation of the constant terms in the numerator of the quotient above, we conclude that, for $u \neq 0$,

$$
F^{\prime}(u)=f(u, \omega) / q^{\omega}(u, 1) q^{\omega}(u,-1) \quad \text { with } q^{\omega}(u, \pm 1) \geq c\left(1+u^{2}\right)^{\nu}
$$

where $f(u, \omega)$ is a polynomial. This proves the first half of Lemma 1 . Now

$$
\begin{aligned}
\int_{K} g^{\omega}(u, v) e^{i r u} d u d v & =\int_{-1}^{1} F(u) e^{i r u} d u \\
& =\left.\frac{1}{i r} F(u) e^{i r u}\right|_{-1} ^{1}-\frac{1}{i r} \int_{-1}^{1} e^{i r u} d F(u) \\
& =\left.\frac{1}{i r} F(u) e^{i r u}\right|_{-1} ^{1}-\frac{2}{i r} F(0+)-\frac{1}{i r} \int_{-1}^{1} e^{i r u} F^{\prime}(u) d u
\end{aligned}
$$

The boundary term appearing above cancels out the contribution from $\partial K$ to the boundary integral in (3.4). The contribution from $\partial D^{\omega}$ to the latter integral equals $(i r)^{-1} \int_{\partial D} g(\theta) e^{r \omega \cdot \theta} \omega \cdot n d s$. Finally, $F(0+)=g^{\sharp}(\omega)$. Thus we obtain the second relation of Lemma 1.
4. Proof of Theorem 2. We collect all the error terms that we neglected as $o(1)$ terms in the proof of Theorem 1 and write $\Psi(x)+\Lambda(x)$ for their sum, where

$$
\Psi(x)=-\operatorname{Re} \int_{T} \psi(\theta) e^{i x \cdot \theta} d \theta
$$

[as already introduced in (2.3)] and

$$
\Lambda(x):=\frac{8}{\sigma_{1} \sigma_{2}} \int_{0}^{\pi / 2} d \alpha \int_{c(x) \sin \alpha}^{\infty} \frac{\cos u}{u} d u-\int_{T \backslash B} \frac{2}{Q(\theta)} \cos (x \cdot \theta) d \theta .
$$

[The first term of $\Lambda(x)$ comes from the last double integral in (2.7) and the second term is (2.6).] As in the previous section, we put $\omega=x /|x|$ and make the change of variables $(u, v)=R_{\omega}^{-1} \theta$ so that $\omega \cdot \theta=u$. Let

$$
T^{\omega}=R_{\omega}^{-1} T \quad \text { and } \quad Q^{\omega}(u, v)=Q\left(R_{\omega}(u, v)\right) .
$$

Lemma 2.

$$
\Lambda(x)=\int_{-\infty}^{\infty} d v \int_{\left\{u:(u, v) \in \mathbb{R}^{2} \backslash T^{\omega}\right\}} \frac{2}{Q^{\omega}(u, v)} \cos (|x| u) d u .
$$

Proof. Recall that we obtained the first term in (2.7), an expression for the integral $\int_{B}[2 / Q(\theta)](1-\cos x \cdot \theta) d \theta$, via a change of variables according to (2.4) and (2.5). By formally reversing the procedure, we see that the first term of the expression defining $\Lambda(x)$ is equal to

$$
\int_{\mathbb{R}^{2} \backslash B} \frac{2}{Q(\theta)} \cos (x \cdot \theta) d \theta
$$

or, by changing the variables of integration, to the iterated integral

$$
\begin{equation*}
\int_{-\infty}^{\infty} d v \int_{\left\{u:(u, v) \in \mathbb{R}^{2} \backslash B^{\omega}\right\}} \frac{2}{Q^{\omega}(u, v)} \cos (|x| u) d u, \tag{4.1}
\end{equation*}
$$

where $B^{\omega}=R_{\omega}^{-1} B$. Hence we obtain an expression for $\Lambda(x)$ in Lemma 2. This argument, however, must be justified because the function $[1 / Q(\theta)] \cos (x \cdot \theta)$, not being Lebesgue integrable on $\mathbb{R}^{2} \backslash B$, does not admit the application of Fubini's theorem.

For justification we consider the integral

$$
\begin{equation*}
I(L):=\frac{8}{\sigma_{1} \sigma_{2}} \int_{0}^{\pi / 2} d \alpha \int_{c(x) \sin \alpha}^{L c(x) \sin \alpha} \frac{\cos u}{u} d u . \tag{4.2}
\end{equation*}
$$

Since the inner integral is bounded by $1+\log ^{+}(1 /(c(x) \sin \alpha)), I(L)$ converges to the first term of the expression defining $\Lambda$ as $L \rightarrow \infty$. Since the function $u^{-1} \cos u$ is integrable on $\{(\alpha, u): c(x) \sin \alpha<u<L c(x) \sin \alpha, 0<$ $\alpha \leq \pi / 2\}$, we may follow the recipe discussed at the beginning of this proof to get $I(L)=\int_{-\infty}^{\infty} f_{L}(v) d v$, where

$$
f_{L}(v):=\int_{\left\{u: \lambda<Q^{\omega}(u, v)<L \lambda\right\}} \frac{2}{Q^{\omega}(u, v)} \cos (|x| u) d u, \quad \lambda:=\left(\left(\sigma_{1} \wedge \sigma_{2}\right) \pi\right)^{2} .
$$

We have only to show that $f_{L}(v)$ is dominated by an integrable function that is independent of $L$ since we can then apply Lebesgue's convergence theorem to see that $I(L)$ converges to (4.1) as $L \rightarrow \infty$. Clearly, $f_{L}$ is bounded uniformly for $L \geq 1$. It therefore suffices to show that, for $a<b$,

$$
\begin{equation*}
\left|\int_{a}^{b} \frac{1}{Q^{\omega}(u, v)} \cos (|x| u) d u\right| \leq \frac{M}{|x| v^{2}}, \tag{4.3}
\end{equation*}
$$

where $M$ is a constant depending on ( $\sigma_{1} \wedge \sigma_{2}$ ) only. However, since the function $1 / Q^{\omega}(u, v)$ with $v$ fixed does not fluctuate at all, that is, it has only one peak for $u \in \mathbb{R}$, and is bounded above by $\left[\left(\sigma_{1} \wedge \sigma_{2}\right) v^{2}\right]^{-1}$, the integral of $\cos (|x| u) / Q^{\omega}(u, v)$ over $u \in(a, b)$ is dominated by $2\left[\left(\sigma_{1} \wedge \sigma_{2}\right)|x| v^{2}\right]^{-1}$ in absolute value; hence (4.3). The proof of Lemma 2 is complete.

We decompose

$$
\begin{aligned}
\Psi(x) & =-\int_{T} \psi_{r}(\theta) \cos (x \cdot \theta) d \theta+\int_{T} \psi_{i}(\theta) \sin (x \cdot \theta) d \theta \\
& :=\Psi_{c}(x)+\Psi_{s}(x) \quad(\text { say })
\end{aligned}
$$

where $\psi_{r}(\theta)$ is the real part of $\psi(\theta)$ and $\psi_{i}(\theta)$ the imaginary part. Put

$$
c(\theta)=E\{1-\cos (\theta \cdot X)\} \quad \text { and } \quad s(\theta)=E\{\sin (\theta \cdot X)\}
$$

Then

$$
\psi_{r}(\theta)=\frac{c(\theta)}{c^{2}(\theta)+s^{2}(\theta)}-\frac{2}{Q(\theta)} \quad \text { and } \quad \psi_{i}(\theta)=\frac{s(\theta)}{c^{2}(\theta)+s^{2}(\theta)}
$$

Since the random walk is aperiodic, $c^{2}(\theta)+s^{2}(\theta)>0$ for $\theta \in T \backslash\{0\}$. If the moment condition (MC: $4+\delta$ ) holds, then, putting

$$
\begin{equation*}
c_{o}(\theta)=\frac{1}{2} Q(\theta)-\frac{1}{24} E\left\{(\theta \cdot X)^{4}\right\} \quad \text { and } \quad s_{o}(\theta)=-\frac{1}{6} E\left\{(\theta \cdot X)^{3}\right\} \tag{4.4}
\end{equation*}
$$

we obtain

$$
\begin{equation*}
\frac{\partial^{l}}{\partial \theta_{1}^{k} \partial \theta_{2}^{j}}\left(c(\theta)-c_{o}(\theta)\right)=O\left(|\theta|^{4-l+\delta}\right) \quad \text { for } l:=k+j=0,1,2 \tag{4.5}
\end{equation*}
$$

(as $\theta \rightarrow 0$ ) and the same estimate with $s(\theta)-s_{o}(\theta)$ in place of $c(\theta)-c_{o}(\theta)$. Applying these estimates together with the identity $(1+z)^{-1}=(1-z) /(1-$ $z^{2}$ ), we readily deduce

$$
\psi_{r}(\theta)=\frac{1}{3 Q^{3}(\theta)}\left[\frac{1}{2} Q(\theta) E\left\{(\theta \cdot X)^{4}\right\}-\frac{2}{3}\left[E\left\{(\theta \cdot X)^{3}\right\}\right]^{2}\right]+\xi_{r}(\theta)
$$

where

$$
\begin{equation*}
\frac{\partial^{l}}{\partial \theta_{1}^{k} \partial \theta_{2}^{j}} \xi_{r}(\theta)=O\left(|\theta|^{-l+\delta}\right) \quad \text { for } l=k+j=0,1,2 \tag{4.6}
\end{equation*}
$$

similarly,

$$
\begin{equation*}
\psi_{i}(\theta)=-\frac{2 E\left\{(\theta \cdot X)^{3}\right\}}{3 Q^{2}(\theta)}+\xi_{i}(\theta) \tag{4.7}
\end{equation*}
$$

where $\xi_{i}(\theta)$ and its derivatives satisfy (4.6) with $\xi_{i}$ in place of $\xi_{r}$.

Recalling what we noticed just before (3.2) [here $\cos (x \cdot \theta)=\nabla \cdot \mathbf{b}(\theta)$ with $\mathbf{b}(\theta):=|x|^{-1} \sin (x \cdot \theta) \omega$ ], we apply the divergence theorem to see that, in view of Lemma 2,

$$
\begin{align*}
\Psi_{c}(x)+\Lambda(x)= & -\frac{1}{|x|} \int_{\mathbb{R}^{2} \backslash T} \omega \cdot \nabla[2 / Q](\theta) \sin (x \cdot \theta) d \theta \\
& +\frac{1}{|x|} \int_{T} \omega \cdot \nabla \psi_{r}(\theta) \sin (x \cdot \theta) d \theta \tag{4.8}
\end{align*}
$$

which is valid if (MC: $3+\delta$ ) holds. Here the boundary terms cancel out each other since $\psi_{r}(\theta)+2 / Q(\theta)$, as well as $\sin (x \cdot \theta)$, is a doubly periodic function of period $(2 \pi, 2 \pi)$. Noticing $\left(\partial / \partial \theta_{1}\right)^{n}(1 / Q(\theta))=O\left(1 /|\theta|^{2+n}\right)$, we see that the first term on the right-hand side of $(4.8)$ is $O\left(1 /|x|^{2}\right)$. On the other hand, $\left|\nabla \psi_{r}(\theta)\right|$ is integrable on $T$, so that the second term is $o(1 /|x|)$. Consequently, $\Psi_{c}(x)+\Lambda(x)=o(1 /|x|)$ under (MC: $\left.3+\delta\right)$.

If (MC: $4+\delta$ ) holds, we can apply Theorem 3 with $m=2$ to the second integral on the right-hand side of (4.8) in view of (4.6). We can always apply the divergence theorem for the first integral. Again the boundary terms cancel out, resulting in

$$
\Psi_{c}(x)+\Lambda(x)=2|x|^{-2}\left(\omega \cdot \nabla \psi_{r, 1}\right)^{\sharp}(\omega)+o\left(|x|^{-2}\right) .
$$

As for $\Psi_{s}(x)$, we have only to apply Theorem 3 with the help of (4.7) to have

$$
\Psi_{s}(x)=\int_{T} \psi_{i}(\theta) \sin x \cdot \theta d \theta=4 \pi^{2} \frac{U_{1}(\omega)}{|x|}+o\left(\frac{1}{|x|^{m}}\right),
$$

where $m=1$ or 2 according to which moment condition we are assuming. These prove (1.3) for $m=1$ and 2 .

In the case when (MC: $2+m+\delta$ ) is assumed to hold for $m \geq 3$, we can perform the Taylor expansion of $1-\cos (\theta \cdot X)$ and $\sin (\theta \cdot X)$ up to the $m$ th-order terms for defining $c_{o}$ and $s_{o}$ in (4.4). We accordingly obtain the estimates $O\left(|\theta|^{m+2-l+\delta}\right)$ for $l=0,1, \ldots, m$ in (4.5), which in turn yields the following expansion for the real and imaginary parts of $\psi$ :

$$
\begin{equation*}
\psi_{r}(\theta)=\frac{\left\{\theta^{6}\right\}}{Q^{3}(\theta)}+\frac{\left\{\theta^{12}\right\}}{Q^{5}(\theta)}+\cdots+\frac{\left\{\theta^{3 m^{\prime}+6}\right\}}{Q^{m^{\prime}+3}(\theta)}+\xi_{r}(\theta), \tag{4.9}
\end{equation*}
$$

where $m^{\prime}=m$ or $m-1$ according to whether $m$ is even or odd and

$$
\frac{\partial^{l}}{\partial \theta_{1}^{k} \partial \theta_{2}^{j}} \xi_{r}(\theta)=O\left(|\theta|^{m-2-l+\delta}\right) \quad \text { for } l=k+j=0,1, \ldots, m ;
$$

and

$$
\begin{equation*}
\psi_{i}(\theta)=\frac{\left\{\theta^{3}\right\}}{Q^{2}(\theta)}+\frac{\left\{\theta^{9}\right\}}{Q^{4}(\theta)}+\cdots+\frac{\left\{\theta^{3 m^{\prime \prime}+6}\right\}}{Q^{m^{\prime \prime}+3}(\theta)}+\xi_{i}(\theta) \tag{4.10}
\end{equation*}
$$

where $m^{\prime \prime}=m-1$ or $m$ according to whether $m$ is even or odd and

$$
\frac{\partial^{l}}{\partial \theta_{1}^{k} \partial \theta_{2}^{j}} \xi_{i}(\theta)=O\left(|\theta|^{m-2-l+\delta}\right) \quad \text { for } l=k+j=0,1, \ldots, m
$$

Here $\left\{\theta^{k}\right\}$ denotes a certain homogeneous polynomial of degree $k$. For evaluating the integral of $\psi_{i}(\theta) \sin \omega \cdot \theta$, we apply Theorem 3 to the right-hand side of (4.10). All the boundary integrals vanish due to the periodicity of $\psi_{i}$. The resultant is

$$
\frac{1}{4 \pi^{2}} \int_{T} \psi_{i}(\theta) \sin x \cdot \theta d \theta=\frac{U_{1}(\omega)}{|x|}+\frac{U_{3}(\omega)}{|x|^{3}}+\cdots+\frac{U_{m^{\prime \prime}}(\omega)}{|x|^{m^{\prime \prime}}}+o\left(\frac{1}{|x|^{m}}\right)
$$

Similarly, we obtain the analogous expansion for $\Lambda(x)+\Psi_{c}(x)$.
5. Self-reciprocity of $\#$. Let $g$ be a quotient $p / q$ of two homogeneous polynomials $p, q$ of degrees $2 \nu-1$ and $2 \nu$, respectively ( $\nu=1,2, \ldots$ ). Suppose $q>0, \theta \neq 0$, and $p$ is relatively prime to $q$. We prove that

$$
g^{\sharp}(\theta):=\text { p.v. } \int_{-\infty}^{\infty} g\left(\theta_{1}-t \theta_{2}, \theta_{2}+t \theta_{1}\right) d t
$$

is then a function of the same type as $g$ with the same $\nu$ and the transform $g \rightarrow g^{\sharp}$ is self-reciprocal, that is, $\pi^{2} g=\left(g^{\sharp}\right)^{\#}$. The proof is given in (i)-(v) below.

Let $R^{\alpha}$ denote rotation by an angle $\alpha$ (counterclockwise). Then:

$$
\begin{equation*}
g^{\sharp}(\theta)=\text { p.v. } \int_{-\pi / 2}^{\pi / 2} g\left(R^{\alpha} \theta\right) \frac{d \alpha}{\cos \alpha} . \tag{i}
\end{equation*}
$$

This equality is obtained by changing the variable according to $t=$ $\tan \alpha\left(-\frac{1}{2} \pi<\alpha<\frac{1}{2} \pi\right)$ so that

$$
g\left(\omega_{1}-t \omega_{2}, \omega_{2}+t \omega_{1}\right)=g\left(\sqrt{1+t^{2}} R^{\alpha} \omega\right)=g\left(R^{\alpha} \omega\right) / \sqrt{1+t^{2}}
$$

From (i) it follows that $\left(g \circ R^{\alpha}\right)^{\#}=g^{\#} \circ R^{\alpha}$.
(ii) $\quad g^{\#}=|\operatorname{det} A|(g \circ A)^{\#} \circ A^{t} \quad$ if $A$ is a regular $2 \times 2$ matrix.

We have noticed in Remark 4 that (ii) is an easy consequence of Theorem 3. Here we give a direct proof. Once (ii) is proved for diagonal matrices, the general case follows from $\left(g \circ R^{\alpha}\right)^{\sharp}=g^{\sharp} \circ R^{\alpha}$ together with the polar decomposition: $A=S O$ ( $S$, symmetric and $O$, orthogonal). Let $A$ be a diagonal matrix with diagonal elements $\lambda$ and $\mu$. Then if $\theta_{1} \neq 0$ the right-hand side of (ii) equals

$$
\begin{aligned}
& \text { p.v. } \int_{-\infty}^{\infty} g\left(\lambda^{2} \theta_{1}-t \lambda \mu \theta_{2}, \mu^{2} \theta_{2}+t \lambda \mu \theta_{1}\right) d t|\lambda \mu| \\
& \quad=\text { p.v. } \int_{-\infty}^{\infty} g\left(\lambda^{2} \theta_{1}^{2}+\mu^{2} \theta_{2}^{2}-t \theta_{2}, t \theta_{1}\right) d t \\
& \quad=\text { p.v. } \int_{-\infty}^{\infty} g\left(1-t \theta_{2}, t \theta_{1}\right) d t .
\end{aligned}
$$

Thus (ii) follows. Since $q$ is factored into quadratic forms, we can decompose $p / q$ into proper fractions, each of which has for the denominator some power of one of the quadratic forms. Applying (ii) to each fraction of the decomposition, we see that:
(iii) $g^{\sharp}$ is a quotient of two homogeneous polynomials of degrees $2 \nu-1$ and $2 \nu$; the denominator is the product of quadratic forms reciprocal to those making q . [We say here that ( $x \cdot A^{-1} x$ ) is reciprocal to ( $x \cdot A x$ ).]
(iv) $g^{\sharp}(\cos \alpha, \sin \alpha)$ equals the limit value of the allied series for the function $g(-\sin \alpha, \cos \alpha)$; in particular, $\#$ is self-reciprocal: $g=\left(g^{\sharp}\right)^{\sharp}$.

If we put, for a function $\phi(\alpha)$,

$$
\begin{equation*}
\phi^{b}(\beta)=\frac{1}{\pi} \text { p.v. } \int_{-\pi / 2}^{\pi / 2} \phi(\beta-\alpha) \frac{d \alpha}{\sin \alpha} . \tag{5.1}
\end{equation*}
$$

Then, making the change of variable $\alpha \rightarrow \pi / 2-\alpha$ for the integral in (i) and using the skew symmetry $g(-\theta)=-g(\theta)$, we have $g^{\sharp}(\omega)=\pi\left(g \circ R^{\pi / 2}\right)^{b}(\omega)$ (where $b$ acts on a restriction of $g \circ R^{\pi / 2}$ to the unit circle). Assertion (iv) follows from the next one:
(v) If $\phi$ is smooth and $\phi(\alpha \pm \pi)=-\phi(\alpha)$, then $\phi^{b}$ agrees with the limit value of the series allied with the Fourier series of the function $\phi$. [Namely, $\phi^{b}(\alpha)=\sum_{n=1}^{\infty} a_{n} \sin n \alpha-b_{n} \cos n \alpha \quad$ if $\quad \phi(\alpha)=\frac{1}{2} a_{0}+\sum_{n=1}^{\infty} a_{n} \cos n \alpha+$ $b_{n} \sin n \alpha$.] In particular, $\left(\phi^{b}\right)^{b}=-\phi$.
To prove (v), substitute the identity $1 / \sin \alpha=\frac{1}{2}\left(\cot \frac{1}{2} \alpha+\tan \frac{1}{2} \alpha\right)$ into the right-hand side of (5.1) and make the change of variable $\alpha \rightarrow \pi-\alpha$ in the integral involving $\tan \frac{1}{2} \alpha$. W then deduce from the assumption $\phi(\alpha \pm \pi)=$ $-\phi(\alpha)$ that

$$
\begin{aligned}
\phi^{b}(\beta) & =\frac{1}{\pi} \text { p.v. } \int_{0}^{\pi} \phi(\beta-\alpha) \frac{d \alpha}{\sin \alpha} \\
& =\frac{1}{2 \pi} \int_{0}^{\pi}[\phi(\beta-\alpha)-\phi(\beta+\alpha)] \cot \frac{1}{2} \alpha d \alpha .
\end{aligned}
$$

This shows the result of (v) since the right-hand side above gives the limit value of the allied series for $\phi$.
6. Hitting distribution of lines. We compute the asymptotic form of the hitting distributions of lines $x_{2}=N$ for large $N$. Suppose that the distribution of $X$ is symmetric with respect to the first coordinate axis $x_{2}=0$ and the random walk $S_{n}$ takes jumps of size at most 1 in the vertical direction, that is, $P\left\{X_{2}=0,1\right.$ or -1$\}=1$. Then the probability that $S_{n}$ enters the line $x_{2}=N$ at a point $x, x_{2}=N$, can be expressed by means of the potential function $a$ as follows:
$H_{N}(x)=\sum_{j=-\infty}^{\infty}\left[a\left(-\left(x_{1}+j, N+1\right)\right)-a\left(-\left(x_{1}+j, N-1\right)\right)\right] P(X=(j, 1))$
(see [1], page 155). We have $Q(\theta)=\sigma_{1}^{2} \theta_{1}^{2}+\sigma_{2}^{2} \theta_{2}^{2}$ and, writing $\tilde{X}_{k}=X_{k} / \sigma_{k}$ and $\quad \tilde{x}_{k}=x_{k} / \sigma_{k}, k=1,2, \psi_{i, 1}{ }^{\circ} \sqrt{Q^{-1}}(-\sin \alpha, \cos \alpha)=b_{1} \sin \alpha+b_{3} \cos \alpha$, where $b_{1}=\frac{1}{2}\left(E\left\{\tilde{X}_{1}^{3}\right\}+E\left\{\tilde{X}_{1}, \tilde{X}_{2}^{2}\right\}\right)$ and $b_{3}=\frac{1}{6}\left(-E\left\{\tilde{X}_{1}^{3}\right\}+3 E\left\{\tilde{X}_{1} \tilde{X}_{2}^{2}\right\}\right)$. With the help of (ii) and (iv) of Section 5, it is obvious that

$$
U_{1}\left(\omega^{x}\right) /|x|=-\left[\left(b_{1}-3 b_{3}\right) \tilde{x}_{1}\|x\|^{-2}+4 b_{3} \tilde{x}_{1}^{3}\|x\|^{-4}\right] /\left(2 \pi \sigma_{1} \sigma_{2}\right)
$$

$\left(\|x\|^{2}=\tilde{x}_{1}^{2}+\tilde{x}_{2}^{2}\right)$. Now, applying (1.3) and making elementary computation, we get

$$
\begin{equation*}
H_{N}(x)=\frac{N}{\pi \sigma_{1} \sigma_{2}\|x\|^{2}}\left(1-2 b_{1} \frac{\tilde{x}_{1}}{\|x\|^{2}}+8 b_{3} \frac{\tilde{x}_{1}^{3}}{\|x\|^{4}}\right)+O\left(|x|^{-3}\right) \tag{6.1}
\end{equation*}
$$

provided that $E\left\{|X|^{5+\delta}\right\}<\infty$. (Here we have applied the smoothness of $U_{2}$.) Relation (6.1) yields, for example, that $\sum_{-m(N) \leq x_{1} \leq m(N)} x_{1} H_{N}(x)$ converges to $\sigma_{1}\left(3 b_{3}-b_{1}\right)=\sigma_{1} E\left\{\tilde{X}_{1}\left(\tilde{X}_{1}^{2}-\tilde{X}_{2}^{2}\right)\right\}$ as $N \rightarrow \infty$ whenever $m(N) \rightarrow \infty$.

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[^0]:    Received July 1995; revised January 1996.
    AMS 1991 subject classifications. 60J15, 60J45, 31C20.
    Key words and phrases. Two-dimensional random walk, potential kernel, Laplace discrete operator, oscillatory integral.

