# On some combinatorial relations concerning the symmetric random walk 

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Dedicated to the three inseparable friends P. Erdös, T. Gallai, and P. Turán. at the occasion of all being close to 50

## § 1. Introduction and notations

1. In this paper we shall consider sequences $\vartheta=\left(\vartheta_{1}, \vartheta_{2}, \ldots, \vartheta_{2 n}\right)$ of $n(+1)$-s and $n(-1)$-s, each possible array being of equal probability $1 /\binom{2 n}{n}$. Thus the partial sums $s_{0}=0, s_{i}=\vartheta_{1}+\vartheta_{2}+\ldots+\vartheta_{i}(i=1,2, \ldots, 2 n)$ generate a simple symmetric random walk, returning after $2 n$ steps to the origin.

We use the following notations:

$$
x=\max _{0 \leqq i \leqq 2 n} s_{i} ; \quad \varrho=\min \left\{i: s_{i}=x\right\} \text { (index of the first maximum). }
$$

$\lambda-1$ is the number of the intersections, i. e. the number of $i$-s with $s_{i}=0$, $s_{i-1} s_{i+1}=-1$ (thus $\lambda$ is the number of half-waves).
$\gamma$ is the Galton-statistics (i. e. $2 \gamma$ is the number of indices $i$ for which either $s_{i}>0$, or $\left.s_{i}=0, s_{i-1}=+1\right)$.

The authors have found the following asymptotic relation [6], [2]:

$$
\begin{aligned}
& \lim _{n \rightarrow \infty} P(x<y \sqrt{2 n}, \varrho<2 n z)=\lim _{n \rightarrow \infty} P(\lambda<y \sqrt{2 n}, \gamma<n z)= \\
&=\sqrt{\frac{2}{\pi}} \int_{0}^{y} \int_{0}^{z} \frac{u^{2}}{[v(1-v)]^{3 / 2}} e^{-\frac{u^{2}}{2 v(1-v)}} d u d v \quad(y \geqq 0 ; 1 \geqq z \geqq 0) .
\end{aligned}
$$

In connection with this relation, E. Sparre Andersen raised the question ${ }^{1}$ ), whether there exists some equivalence principle for the finite case too.

In the following we give some equivalence theorems and prove among others that

$$
\begin{align*}
& P(x=l)=\frac{1}{2}(P(\lambda=l)+P(\lambda=l+1)) \quad(l=0,1, \ldots, n),  \tag{1.1}\\
& P(x=l, \varrho=r)=P\left(\lambda^{\prime}=l, \pi=r\right) \tag{1.2}
\end{align*}
$$

[^0]where $\pi$ denotes the number of positve terms in $\left(s_{0}, s_{1}, \ldots, s_{2 n}\right)$, while $\lambda^{\prime}$ is the number of indices $i$ for which $s_{i-1}=0, s_{i}=+1$. (1.2) implies
$$
P(\varrho=r)=P(\pi=r),
$$
which is a special case of a well-known result of Sparre Andersen [1]; it implies also the following result of Mihalevič [5]:
$$
P(\varkappa=l)=P\left(\lambda^{\prime}=l\right) .
$$

Thus we have a joint equivalence between $(\varkappa, \varrho)$ and $\left(\lambda^{\prime}, \pi\right)$. We would like to point out furthermore that each of our theorems is proved by means of one-to-one correspondences between the sets of paths considered. This indicates a combinatorial and geometrical background of these equivalences.

We also remark that in our constructions $\chi$ appears virtually more as the number of ladder indices (see Feller [3]), than as the maximum, both coinciding for the special variables $\vartheta_{i}= \pm 1$.
2. We shall make use of the following further terminology and notations:

The polygonal line whose subsequent vertices have the coordinat es ( $i, s_{i}$ ) ( $i=0,1,2, \ldots, j$ ) is called the path $\left(s_{0}, s_{1}, \ldots, s_{j}\right)$.
$E_{2 n}$ is a path $\left(s_{0}, s_{1}, \ldots, s_{2 n}\right)$ with $s_{2 n}=0$. A point $\left(2 i, s_{2 i}\right)$ of the path $E_{2 n}$, for which $s_{2 i}=0$ and $s_{2 i-1} s_{2 i+1}=-1$, as well as the points $(0,0)$ and $(2 n, 0)$ of $E_{2 n}$, are called intersection points or $T$-points. As defined above, $\lambda+1$ is the number of $T$-points.

By a $T^{\prime}$-point we mean a point $(2 i+1,1)$ of the path $E_{2 n}$, for which $s_{2 i}=0$, $s_{2 i+1}=+1$ (this kind of points was treated by Mihalevič [5]). $\lambda^{\prime}$ is the number of $T^{\prime}$-points.
$E_{2 n}^{l}$ is a path $E_{2 n}$ with $\lambda=l$,
$\langle i, j\rangle$ is a section of a path lying between the points $\left(i, s_{i}\right)$ and $\left(j, s_{j}\right)$, i. e. the sequence ( $\vartheta_{i+1}, \vartheta_{i+2}, \ldots, \vartheta_{j}$ ).
$k$ is called a strict ladder index (Feller [3]); if $s_{k}>s_{i}$ for $i=0,1, \ldots, k-1 ; k$ is called a strict backward ladder index if $s_{k}>s_{i}$ for $i=k, k+1, \ldots, 2 n$.
$A_{r}^{l}$ is a path $\left(s_{0}, s_{1}, \ldots, s_{r}\right)$, for which $s_{0}=0, s_{1}<l, s_{2}<l, \ldots, s_{r-1}<l, s_{r}=l$, i. e. its $l$-th strict ladder index being $r$.
$N(\cdot)$ is the number of all possible paths whose type is given in the brackets (e. g. $\left.N\left(E_{2 n}\right)=\binom{2 n}{n}\right)$.

## § 2. Equivalence relations

1. The maxium and the number of waves. We shall prove the following

Theorem 2.1. $P(x=l)=\frac{1}{2}[P(\lambda=l)+P(\lambda=l+1)](l=0,1,2, \ldots, \dot{n})$.
Proof. We consider a path $E_{2 n}$ with $\kappa=1$. According to the index $\varrho$ of the first maximum, we distinguish two different cases:
a) $\varrho$ is the only position, for which the maximum takes place;
b) there are more than one maximum places.

In both cases we shall make use of the following
Lemma 2. 1 .

$$
\frac{1}{2} N\left(\dot{E}_{2 n}^{\prime}\right)=N\left(A_{2 n}^{2 l}\right):
$$

This was proved in [2] by-means of a one-to-one correspondence between the sets of paths.

In case a) we consider the sections $\langle 0, \varrho\rangle$ and $\langle\varrho, 2 n\rangle$. Replacing in the second part the steps $\left(\vartheta_{e+1}, \vartheta_{\varrho+2}, \ldots, \vartheta_{2 n}\right)$ by the steps $\left(-\vartheta_{2 n},-\vartheta_{2 n-1}, \ldots,-\vartheta_{\rho+2},-\vartheta_{\rho+1}\right)$, we obtain a path $A_{2 n}^{2 l}$. According to Lemma 2.1 this path can be transformed into a path $E_{2 n}^{l}$ with $s_{1}=+1$.

Obviously this procedure is invertible, by considering the $l$-th strict ladder index of the path $A_{2 n}^{2 l}$.

In case b) let us denote by $\bar{\varrho}$ the index of the last maximum. The path $E_{2 n}$ with $s_{e}=s_{\bar{e}}=l$ consists of the following three sections: $\langle 0, \varrho\rangle,\langle\varrho, \bar{\varrho}\rangle,\langle\bar{\varrho}, 2 n\rangle$. We apply the following transformation: we replace in $\langle\varrho ; \bar{\varrho}\rangle$ the steps $\left(\hat{\vartheta}_{\rho+1}=-1, \vartheta_{\emptyset+2}, \ldots, \vartheta_{\bar{\varphi}}\right)$ by $\left(\vartheta_{\varrho+2}, \ldots, \vartheta_{\bar{\varrho}},+1\right)$ and in $\langle\bar{\varrho}, 2 n\rangle$ the steps $\left(\vartheta_{\bar{p}+1}, \vartheta_{\underline{\rho}+2}, \ldots, \vartheta_{2 n}\right)$ by the steps $\left(-\vartheta_{2 n},-\vartheta_{2 n-1}, \ldots,-\vartheta_{\bar{e}+2},-\vartheta_{\dot{e+1}}\right)$. Thus we obtain a path $A_{2 n}^{2 l+2}$. According to Lemma 2.1 this path can be transformed into a path $E_{2 n}^{l+1}$ with $s_{1}=+1$.

In order to invert this procedure we have only to find the $l$-th and $l+2$-th ladder indices of the path $A_{2 n}^{2 l+1}$. Cases a) and b) complete the proof of Theorem 2. 1.
2. Two variate equivalences. We shall prove the following

Theorem 2.2. $P(\kappa=l, \varrho=r)=P\left(\lambda^{\prime}=l, \pi=r\right) \quad(l=0, r=0 ; l=1,2, \ldots, n$, $r=l, l+2, \ldots, 2 n-l)$.

Proof ${ }^{2}$ ). For $r=0, l=0$ the paths of both kinds coincide, we have to consider only the case $l \geqq 1$. Then each path with ( $\lambda^{\prime}=l, \pi=r$ ) can be divided by the $T^{\prime}$ points $(2 i+1,1)$ and the points $(2 j, 0)$ with $s_{2 j}=0$ and $s_{2 j-1}=+1$ into $2 l$ or $2 l+1$ sections, some of which are starting from +1 and ending in 0 , all inner points being strictly positive (type $\alpha$ ), while the others are starting from 0 , ending in +1 , all inner points being non-positive (type $\beta$ )).

The first section is always of type $\beta$ ); the last section is either of type $\alpha$ ) or of type $\beta$ ), but in the latter case the last ( $\vartheta_{2 n+1}=+1$ ) step is missing.

There are altogether $l$ sections of type $\alpha$ ) with total length $r$ and $l$ or $(l+1)$ sections of type $\beta$ ).

Let us now consider the sections of type $\alpha$ ). We change all $\theta_{i}$-s occurring in them into $-\vartheta_{i}$ and link together the new sections obtained by this procedure, maintaining their original order of succession. We now link together all sections of type $\beta$ ); denoting the steps of the section thus obtained by $\left(\vartheta_{r+1}^{\prime}, \vartheta_{r+2}^{\prime}, \ldots, \vartheta_{2 n}^{\prime}\right)$ we transform them into ( $-\vartheta_{2 n}^{\prime},-\vartheta_{2 n-1}^{\prime}, \ldots,-\vartheta_{r+2}^{\prime},-\vartheta_{r+1}^{\prime}$ ) and join the respective section to the first section obtained. As a result we obtain a path with $x=l, \varrho=r$.

The reverse procedure transforms each path $\{x=l, \varrho=r\}$ into the corresponding path $\left\{\lambda^{\prime}=l, \pi=r\right\}$; this can be performed by considering the strict ladder indices in section $\langle 0, r\rangle$ and the strict backward ladder indices in section $\langle r, 2 n\rangle$.

[^1]In the following theorem we shall prove two equivalences according to whether the maximum is even or odd:

Theorem 2. 3.

$$
P\left(\dot{\varkappa}=s_{2 r^{\prime}}=2 l\right)=\frac{1}{2} P\left(\lambda=2 l, \gamma=r^{\prime}\right)+P\left(\lambda=2 l+1, \gamma=r^{\prime}\right)+\frac{1}{2} P\left(\lambda=2 l+2, \gamma=r^{\prime}\right) .
$$

and

$$
P\left(x=\dot{s}_{2 r^{\prime}+1}=2 l-1\right)=
$$

$=P\left(\lambda=2 l-1, \gamma=r^{\prime} ; s_{1} \doteq-1\right)+P\left(\lambda=2 l, \gamma=r^{\prime}\right)+P\left(\lambda=2 l+1, \gamma=r^{\prime}, s_{1}=+1\right)$.
Proof. We use the same procedure as in the proof of Lemma 2.1.
The crucial point in the proof of Lemma 2. 1 was the division of a path $\boldsymbol{A}_{2 n}^{2 l}$ by means of its even strict ladder indices. The last step of each section between two consecutive ladder indices is always $(+1)$; omitting this and placing a $(-1)$ before the section, we obtain a negative half wave.

Considering a path whose maximum $2 l$ is taken on for the index $2 r^{\prime}$ let us denote by $2 r(2 \bar{r})$ the first (last) index of maximum. The section $\langle 0,2 r\rangle$ is a path $A_{2 r}^{2 l}$, the section $\langle 2 \bar{r}, 2 n\rangle$ is an inverted path $A_{2(n-\bar{r})}^{2 l}$. As described before, both sections can be divided into $l$ parts and each part can be transformed into a negative half wave. The half waves generated by $A_{2 r}^{2 l}$ will be turned into positive half waves by reflection. If $r=r^{\prime}=\bar{r}$ (case ā)) there is no other section; if $r<\bar{r}$ but either $r^{\prime}=r$ or $r^{\prime}=\bar{r}$ (case $\overline{\mathrm{b}}$ ) ) the section $\langle 2 r, 2 \bar{r}\rangle$ is a negative half wave itself. In this case if $r^{\prime}=r$, then this half wave will remain negative, if $r^{\prime}=r$, it will be turned into positive one. If $r<r^{\prime}<\bar{r}$ (case $\overline{\mathrm{c}}$ )) the sections $\left\langle 2 r, 2 r^{\prime}\right\rangle$ and $\left\langle 2 r^{\prime}, 2 \bar{r}\right\rangle$ are half waves themselves. The former will be turned into a positive one, the latter will remain negative. What remains to be done is to connect these half waves, namely a positive after a negative one; in case $\bar{a}$ ) and $\bar{c}$ ) beginning with a positive half wave, in case $\bar{b}$ ) with a negative one if $r^{\prime}=r$ and with a positive one if $r^{\prime}=\bar{r}$.

Each of these procedures determines uniquely the inverse construction, leading to a one-to-one mapping of the sets of corresponding paths. For the second part of this theorem similar construction can be applied.

Summation over $l$ of the relations in Theorem 2:3 results in the following
Corollary 2. 2. $P\left(s_{2 r+1}=\chi\right)=P\left(s_{2 r}=\chi\right)=P(\gamma=r)$ for $r=0,1,2, \ldots, n$.
Another fact proved herewith is expressed in the
Corollary 2. 3.

$$
P(\varkappa=2 l, \varrho=2 r)=\frac{1}{2} P(\lambda=2 l, \gamma=r)+P\left(\lambda=2 l+1, \gamma=r, s_{1}=+1\right) .
$$

and

$$
P(\dot{\varkappa}=2 l-1, \varrho=2 r+1)=
$$

$$
=P\left(\lambda=2 /-1, \gamma=r, s_{1}=-1\right)+\frac{1}{2} P(\lambda=2 /, \gamma=r)
$$

## References

[1] E. S. Andersen, On the fluctuations of sums of random variables. I and II, Math. Scand., 1 (1953); 263-285, and 2 (1954), 195-223.
[2] E. Cșíki and I. Vincze, On some problems connected with the Galton-test, Publ. Math. Inst. Hung. Acad. Sci., 6 (1961), 97-109.
[3] W. Feller, On combinatorial methods in fluctuation theory, Probability and Statistics. The Harald Cramér Volume (1959), 75-91.
[4] Ch. Hobby and R. Pyke, Combinatorial results in fluctuation theory. Technical Report, No. 38 (1962), University of Washington (Seattle).
[5] B. C. Михалевич, О взаимном расположении двух эмпирических функций распределения, Доклады Акад. Наук СССР, 80 (1951), 525-528.
[6] I. Vincze, Einige zweidimensionale Verteilungs- und Grenzverteilungssätze in der Theorie der geordneten Stíchproben, Publ. Math. Inst. Hung. Acad. Sci., 2 (1957), 183-209.
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[^0]:    ${ }^{1}$ ) At the occasion of the Conference on Probability and Statistics held in Oberwolfach, August 20-26, 1961.

[^1]:    ${ }^{2}$ ) Similar construction is used by Ch. HobBy and R. Pyke [4].

