# Strong Approximation of Very Weak Bernoulli Processes 

Ernst Eberlein<br>Institut f. Mathem. Stochastik, Universität Freiburg, Hebelstraße 27, D-7800 Freiburg, Federal Republic of Germany

Summary. Very weak Bernoulli processes with values in a separable metric space are introduced. An estimate for the Prohorov distance in the central limit theorem is obtained. This estimate is used to establish a strong (almost sure) approximation of the partial sums of a very weak Bernoulli process by a Brownian motion where the error term is of the order $O\left(t^{1 / 2-\gamma}\right)$. The proofs are based on a new version of the Berkes-Philipp approximation theorem.

## 1. Introduction and Statement of Results

Functional central limit theorems as well as the functional law of the iterated logarithm and other asymptotic fluctuation results for processes $X=\left(x_{k}\right)_{k \in Z}$ can be derived from a strong (or almost sure) approximation of the partial sums of the process by a Brownian motion with a sufficiently small error term.

In a remarkable paper Kuelbs and Philipp [14] recently established a number of strong approximation results for Banach space valued sequences satisfying various mixing conditions. Their proofs rest on an approximation theorem given first by Berkes and Philipp [2] and in more general form by Philipp [16]. Dehling and Philipp [7] showed that various approximation results stated in [14] for $\phi$-mixing sequences continue to hold for absolutely regular sequences. Dehling [4] succeeded to improve the error term in the approximation of the last mentioned sequences. By constructing a counterexample he pointed out the limits of the argumentation used ([4], [5]). For processes with a finite state space absolute regularity is equivalent to being weak Bernoulli.

In this paper we will get an approximation of the order

$$
\begin{equation*}
\sum_{v \leqq t} x_{v}-X(t) \ll t^{1 / 2-\gamma} \tag{1.1}
\end{equation*}
$$

for very weak Bernoulli processes (Definition 1). Very weak Bernoulli processes with finite state spaces were introduced by Ornstein (see e.g. [15]) in connection with the solution of the isomorphism problem of ergodic theory.

Our first result (Theorem 1) is a new version of the Berkes-Philipp approximation theorem. It is used in the proof of Theorem 2 to estimate the Prohorov distance in the central limit theorem for very weak Bernoulli processes. Both Theorem 1 as well as Theorem 2 are needed to prove (1.1) which is stated as Theorem 3. The last result finally is a further application of our basic approximation theorem.

Some remarks concerning the notation: $\mathscr{L}(X)$ denotes distribution of a random variable $X$ defined on a probability space $(\Omega, \mathfrak{M}, P)$. If $A \in \mathfrak{N}, P(A)>0$ we denote by $\mathfrak{L}(X \mid A)$ the conditional distribution of $X$ given $A$. $\mathfrak{L}\left(\left(X_{1}, \ldots, X_{n}\right) \mid A\right)$ is used if $X=\left(X_{1}, \ldots, X_{n}\right)$ is a vector. The variance of a realvalued random variable $X$ is denoted by $\operatorname{Var}(X) . N\left(0, \tau^{2}\right)$ stands for a normal distribution with mean zero and variance $\tau^{2}$. Stationarity of a process will always mean strict stationarity. Given a Banach space $B, B^{*}$ is the topological dual space. For any subset $E$ of a metric space ( $S, \sigma$ ) and $\varepsilon>0, E^{\varepsilon}$ denotes the open $\varepsilon$-neighborhood of $E . E^{\subset}$ means complement of a set $E . Z$ are the integers. "«" and " $O($.$) " are used with the same meaning.$
Theorem 1. Let $\left\{\left(S_{k}, \sigma_{k}\right) \mid k \geqq 1\right\}$ be a sequence of complete separable metric spaces and let $\left(X_{k}\right)_{k \geqq 1}$ be a sequence of random variables with values in $S_{k}$. Let $\left(\mathfrak{F}_{k}^{\prime}\right)_{k \geqq 1}$ be a sequence of $\sigma$-algebras such that $X_{k}$ is $\mathfrak{F}_{k}^{\prime}$-measurable. Suppose that for every $k \geqq 1$ there exist $D_{k} \in \mathfrak{V}_{k}^{\prime}=\bigvee_{j \leq k} \mathscr{F}_{j}^{\prime}, \varepsilon_{k} \geqq 0$ and $\eta_{k} \geqq 0$ such that $P\left(D_{k}^{\subset}\right) \leqq \eta_{k}$ and if $A \in \mathfrak{Y}_{k-1}^{\prime}, A \subset D_{k-1}, P(A)>0 \begin{gathered}j \leq k \\ \text { then }\end{gathered}$

$$
\begin{equation*}
\inf _{\lambda \in \mathfrak{P}_{k}} \int \sigma_{k}(u, v) d \lambda(u, v) \leqq \varepsilon_{k}^{2} / 2 \tag{1.2}
\end{equation*}
$$

where $\mathfrak{P}_{k}$ is the class of all Borel probability measures $\lambda$ on $S_{k} \times S_{k}$ with marginals $\mathfrak{L}\left(X_{k} \mid A\right)$ and $\mathcal{L}\left(X_{k}\right)$. Let $\left(G_{k}\right)_{k \geqq 1}$ be a sequence of Borel probability measures on $S_{k}$ and $\rho_{k}, \delta_{k}$ nonnegative numbers such that

$$
\begin{equation*}
\mathfrak{L}\left(X_{k}\right)(E) \leqq G_{k}\left(E^{\rho_{k}}\right)+\delta_{k} \tag{1.3}
\end{equation*}
$$

for all Borel sets $E \subset S_{k}$.
Then without loss of generality there exists a sequence $\left(Y_{k}\right)_{k \geqq 1}$ of independent random variables such that $\mathscr{Q}\left(Y_{k}\right)=G_{k}$ and

$$
\begin{equation*}
P\left[\sigma_{k}\left(X_{k}, Y_{k}\right) \geqq 2 \rho_{k}+\varepsilon_{k}\right] \leqq \varepsilon_{k}+\eta_{k-1}+\delta_{k} \tag{1.4}
\end{equation*}
$$

Here and in Theorem 3 the phrase "without loss of generality ..." is to be understood in the sense of Strassen [20]: without changing its distribution we can redefine the sequence on a new probability space on which there exists a sequence $\left(Y_{k}\right)_{k \geqq 1}$ satisfying (1.4) (resp. a Brownian motion $(X(t))_{t \geqq 0}$ satisfying (1.1)).
(1.2) means that the distributions $\mathfrak{L}\left(X_{k} \mid A\right)$ and $\mathfrak{L}\left(X_{k}\right)$ have distance less than $\varepsilon_{k}^{2} / 2$ in the Wasserstein metric which is defined in general as follows. Let $\mu, v$ be two Borel probability measures on a complete separable metric space $(S, \sigma)$ then we define their Wasserstein distance by

$$
\begin{equation*}
\rho(\mu, v)=\inf _{\lambda \in \mathfrak{P}} \int \sigma(u, v) d \lambda(u, v) \tag{1.5}
\end{equation*}
$$

where $\mathfrak{P}$ is the class of Borel probability measures on $S \times S$ with marginals $\mu$ and $v$. Since $S$ is complete and separable both of the measures $\mu$ and $v$ are tight. This implies that the infimum in (1.5) is actually a minimum.

Recall that the Prohorov distance is defined by

$$
\begin{equation*}
\pi(\mu, \nu)=\inf \left\{\varepsilon>0 \mid \mu(E) \leqq v\left(E^{\varepsilon}\right)+\varepsilon \text { for all closed } E \subset S\right\} . \tag{1.6}
\end{equation*}
$$

We will make use of the following relation

$$
\begin{equation*}
\pi(\mu, v) \leqq(\rho(\mu, v))^{1 / 2} \tag{1.7}
\end{equation*}
$$

This inequality can easily be derived from the Čebyšev-Markov inequality and the following alternative definition for $\pi$

$$
\pi(\mu, v)=\inf _{\lambda \in \mathfrak{P}} \inf \{\varepsilon>0 \mid \lambda(\{(u, v) \mid \sigma(u, v)>\varepsilon\}) \leqq \varepsilon\}
$$

which is a consequence of the Strassen-Dudley theorem [9].
We will now define the Wasserstein distance for distributions of processes. For each $n$ we consider a metric $\sigma_{n}$ on the product space $S^{n}=S \times \ldots \times S$

$$
\sigma_{n}(u, v)=n^{-1} \sum_{i=1}^{n} \sigma\left(u_{i}, v_{i}\right)
$$

if $u=\left(u_{1}, \ldots, u_{n}\right), v=\left(v_{1}, \ldots, v_{n}\right) \in S^{n}$. Given two $S$-valued processes $X=\left(x_{k}\right)_{k \in Z}$ and $Y=\left(y_{k}\right)_{k \in Z}$ we define for $n \geqq 1, k \in Z$

$$
\begin{equation*}
\bar{\rho}_{n}\left(\left(x_{k+1}, \ldots, x_{k+n}\right),\left(y_{k+1}, \ldots, y_{k+n}\right)\right)=\inf _{\lambda \in \mathfrak{P}_{n}} \int_{n} \sigma_{n} d \lambda \tag{1.8}
\end{equation*}
$$

where $\mathfrak{P}_{n}$ is the class of all Borel probability measures on $S^{n} \times S^{n}$ with marginals $\mathfrak{L}\left(\left(x_{k+1}, \ldots, x_{k+n}\right)\right)$ and $\mathfrak{L}\left(\left(y_{k+1}, \ldots, y_{k+n}\right)\right)$. Since $S^{n}$ is separable and complete again the infimum in (1.8) is a minimum. If $k=0$ then we write for short $\bar{\rho}_{n}(X, Y)$ in (1.8). The processes $Y$ we consider in the following are processes derived from a fixed process $X$ by conditioning $X$ on certain sets $A$ of positive measure, i.e. we consider measures $\lambda$ with one of the marginals being $\mathscr{L}\left(\left(x_{k+1}, \ldots, x_{k+n}\right) \mid A\right)$. Again we shall write $X \mid A$ for short if $k=0$.

In the following $\left(\mathscr{F}_{k}\right)_{k \in Z}$ will denote a family of sub- $\sigma$-algebras, where we will always assume without further mention that $\mathscr{F}_{k}$ is countably generated for each $k$.
Definition 1. Let $X=\left(x_{k}\right)_{k \in Z}$ be a strictly stationary $S$-valued process defined on a probability space $(\Omega, \mathfrak{U}, P)$ and let $\left(\mathfrak{F}_{k}\right)_{k \in Z}$ be a family of sub- $\sigma$-algebras of $\mathfrak{H}$ such that $x_{k}$ is $\mathfrak{F}_{k}$-measurable. The process is called very weak Bernoulli (with respect to $\left.\left(\mathcal{F}_{k}\right)_{k \in Z}\right)$ if for each $\varepsilon>0$ there is an $n$ such that for all integers $m \geqq 0$, $k \in Z$ there exists a set

$$
D=D(n, m, k) \in \mathfrak{A}_{k-m}^{k}=\bigvee_{j=k-m}^{k} \mathfrak{F}_{j}
$$

such that $P\left(D^{\subset}\right)<\varepsilon$ and if $A \in \mathfrak{A}_{k-m}^{k}, A \subset D, P(A)>0$ then

$$
\begin{equation*}
\bar{\rho}_{n}\left(\left(x_{k+1}, \ldots, x_{k+n}\right),\left(x_{k+1}, \ldots, x_{k+n}\right) \mid A\right)<\varepsilon . \tag{1.9}
\end{equation*}
$$

Given a sequence $(\varepsilon(n))_{n \geqq 1}$ of positive numbers decreasing to 0 we say $X$ $=\left(x_{k}\right)_{k \in Z}$ is very weak Bernoulli (with respect to $\left(\mathscr{F}_{k}\right)_{k \in Z}$ ) at rate $(\varepsilon(n))_{n \geqq 1}$ if for all integers $m \geqq 0, n \geqq 1, k \in Z$ there exists a set $D=D(n, m, k) \in \mathfrak{S}_{k-m}^{k}$ such that $P\left(D^{\subset}\right) \leqq \varepsilon(n)$ and if $A \in \mathfrak{M}_{k-m}^{k}, A \subset D, P(A)>0$ then

$$
\begin{equation*}
\bar{\rho}_{n}\left(\left(x_{k+1}, \ldots, x_{k+n}\right),\left(x_{k+1}, \ldots, x_{k+n}\right) \mid A\right) \leqq \varepsilon(n) . \tag{1.10}
\end{equation*}
$$

The process is called strictly very weak Bernoulli resp. strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ if (1.9) resp. (1.10) hold for any set $A \in \mathfrak{M}_{k-m}^{k}, P(A)>0$.

The dependence structure given by this definition differs from the mixing conditions in the sense that the rate $(\varepsilon(\mathrm{n}))_{n \geq 1}$ has implications on the qualitative dependence structure. More precisely, a process which is very weak Bernoulli at rate $\varepsilon(n)=o\left(n^{-1}\right)$ is a Bernoulli process, i.e. is independent. It is easy to see that this rate implies that the process is strongly mixing. By a simple extension of the argument it was shown recently [6] that this rate actually implies independence. ("Strongly mixing" is used here in the probabilistic sense, not in the terminology of ergodic theory).

If the state space $S$ is a finite set endowed with the discrete metric and the $\mathfrak{F}_{k}$ are the $\sigma$-algebras generated by the $x_{k}$, Definition 1 coincides with the one due to Ornstein. In the following two theorems we consider real-valued processes. Thus $\bar{\rho}_{n}$ will be computed with respect to
for $u, v \in \mathbb{R}^{n}$.

$$
\begin{equation*}
\sigma_{n}(u, v)=n^{-1} \sum_{i=1}^{n}\left|u_{i}-v_{i}\right| \tag{1.11}
\end{equation*}
$$

Theorem 2. Let $X=\left(x_{k}\right)_{k=Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate

$$
\begin{equation*}
\varepsilon(n) \ll n^{-1} \tag{1.12}
\end{equation*}
$$

Suppose $x_{0}$ is centered and bounded with probability 1 and $\lim _{n \rightarrow \infty} \operatorname{Var}\left(\sum_{k=1}^{n} x_{k}\right)=\infty$ then there exists a finite positive value $\sigma^{2}$ such that $\lim _{n \rightarrow \infty} n^{-1} \operatorname{Var}\left(\sum_{k=1}^{n} x_{k}\right)=\sigma^{2}$ and
for some $\kappa>0$

$$
\begin{equation*}
\pi\left(Q\left(n^{-1 / 2} \sum_{v=1}^{n} x_{v}\right), N\left(0, \sigma^{2}\right)\right) \ll n^{-\kappa} . \tag{1.13}
\end{equation*}
$$

The explicit value of the exponent is $\kappa=1 / 64$.
Theorem 3. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate (1.12). Suppose $x_{0}$ is centered and bounded with probability 1 and $\lim _{n \rightarrow \infty} \operatorname{Var}\left(\sum_{k=1}^{n} x_{k}\right)=\infty$ then there exists a finite positive value $\sigma^{2}$ such that $\lim _{n \rightarrow \infty} n^{-1} \operatorname{Var}\left(\sum_{k=1}^{n} x_{k}\right)=\sigma^{2}$ and without loss of generality there exists a Brownian motion $(X(t))_{t \geqq 0}$ with variance $\sigma^{2}$ such that

$$
\begin{equation*}
\sum_{v \leqq t} x_{v}-X(t) \ll t^{1 / 2-\gamma} \tag{1.14}
\end{equation*}
$$

for some $\gamma>0$.

The next definition is a straightforward extension of a notion given in [19]. Let $X$ and $Y$ again be $S$-valued processes. Recall that $\bar{\rho}_{n}(X, Y)$ denotes the Wasserstein distance of $\mathfrak{Q}\left(\left(x_{1}, \ldots, x_{n}\right)\right)$ and $\mathscr{L}\left(\left(y_{1}, \ldots, y_{n}\right)\right)$. Set

$$
\begin{equation*}
\bar{\rho}(X, Y)=\sup _{n \geqq 1} \bar{\rho}_{n}(X, Y) \tag{1.15}
\end{equation*}
$$

Definition 2. Let $X=\left(x_{k}\right)_{k \in \mathcal{Z}}$ be a stationary $S$-valued process. $X$ is called almost block independent if given $\varepsilon>0$ there is an $N$ such that if $n \geqq N$ and $Y$ is the process defined by the two conditions

$$
\begin{align*}
& \quad \mathcal{L}\left(\left(y_{n k+1}, \ldots, y_{n(k+1)}\right)\right)=\mathscr{Q}\left(\left(x_{1}, \ldots, x_{n}\right)\right) \quad \text { for all } k \in Z,  \tag{1.16}\\
& \text { for each } k,\left(y_{n k+1}, \ldots, y_{n(k+1)}\right) \text { is independent of }\left\{\left(y_{j}\right)_{j \leqq n k}\right\} \tag{1.17}
\end{align*}
$$

then

$$
\bar{\rho}(X, Y) \leqq \varepsilon
$$

Theorem 4. Let $X=\left(x_{k}\right)_{k \in Z}$ be a very weak Bernoulli $S$-valued process and suppose that $\sigma$ is bounded, then $X$ is almost block independent.

Let us finally point out another aspect of this paper. The central limit theorem and the strong approximation result established here are of a rather different nature than the aims and statements one is looking towards in the isomorphism theory. But it became obvious that some of the underlying approximation ideas are exactly the same.

## 2. Proof of Theorem 1

We consider the case of discrete random variables $X_{k}$ and probability measures $G_{k}=\mathscr{L}\left(X_{k}\right)$ first. This means that the assumptions hold for $\rho_{k}=\delta_{k}=0(k \geqq 1)$.

Define $Y_{1}=X_{1}$ and suppose $Y_{1}, \ldots, Y_{k-1}$ have been constructed. Suppose furthermore $Y_{j}$ is $\mathfrak{H}_{j}^{\prime}$-measurable $(1 \leqq j \leqq k-1)$. Define $D=D\left(b_{1}, \ldots, b_{k-1}\right)=\left\{Y_{1}\right.$ $\left.=b_{1}, \ldots, Y_{k-1}=b_{k-1}\right\} \cap D_{k-1}$ for values $b_{1}, \ldots, b_{k-1}$ in the range of $Y_{1}, \ldots, Y_{k-1}$ respectively. Then $D \in \mathfrak{V}_{k-1}^{\prime}$ and $D \subset D_{k-1}$.

As we mentioned already the infimum in (1.2) is a minimum. Therefore if $P(D)>0$ there is a Borel probability measure $\lambda$ on $S_{k} \times S_{k}$ with marginals $\mathscr{L}\left(X_{k} \backslash D\right), \mathscr{L}\left(X_{k}\right)$ and

$$
\begin{equation*}
\lambda\left(\left\{(u, v) \mid \sigma_{k}(u, v) \geqq \varepsilon_{k}\right\}\right) \leqq \varepsilon_{k}^{-1} \int \sigma_{k} d \lambda \leqq \varepsilon_{k} \tag{2.1}
\end{equation*}
$$

Note that for this part of the proof we use only the bound $\varepsilon_{k}^{2}$ in (1.2). Let $\left\{a_{i} \mid i \geqq 1\right\}$ be the range of $X_{k}$. Then $\lambda$ is concentrated on $\left\{\left(a_{i}, a_{j}\right) \mid i, j \geqq 1\right\}$ and

$$
P\left[X_{k}=a_{i} \mid D\right]=\sum_{j \geqq 1} \lambda\left(a_{i}, a_{j}\right)
$$

for each $i \geqq 1$. All the $\mathfrak{F}_{k}^{\prime}$ can be assumed to be atomless since $\left(X_{k}\right)_{k \geqq 1}$ can be redefined on a richer probability space if necessary. Therefore $\left\{X_{k}=a_{i}\right\} \cap D \in \mathfrak{G}_{k}^{\prime}$ can be partitioned into $\mathscr{Y}_{k}^{\prime}$-measurable sets $D_{i j}$ such that $P\left(D_{i j} \mid D\right)=\lambda\left(a_{i}, a_{j}\right)$. We
define $Y_{k}=a_{j}$ on $D_{i j}(i, j \geqq 1)$. Then the joint distribution of $\left(X_{k}, Y_{k}\right)$ on $D$ is $\lambda$. By (2.1) this means

$$
\begin{equation*}
P\left(\left\{\sigma_{k}\left(X_{k}, Y_{k}\right) \geqq \varepsilon_{k}\right\} \cap D\right) \leqq P(D) \varepsilon_{k} \tag{2.2}
\end{equation*}
$$

On sets $D$ with $P(D)=0$ and on $D^{\prime}=\left\{Y_{1}=b_{1}, \ldots, Y_{k-1}=b_{k-1}\right\} \cap D_{k-1}^{\subset}$ we define $Y_{k}$ such that $\mathcal{Q}\left(Y_{k} \mid D\right)=\mathfrak{Q}\left(Y_{k} \mid D^{\prime}\right)=\mathscr{L}\left(X_{k}\right)$. Now $Y_{k}$ is defined on the whole space and by (2.2)

$$
\begin{equation*}
P\left[\sigma_{k}\left(X_{k}, Y_{k}\right) \geqq \varepsilon_{k}\right] \leqq \sum_{D \subset D_{k-1}} P(D) \varepsilon_{k}+\sum_{D^{\prime} \subset D_{k-1}^{\in}} P\left(D^{\prime}\right) \leqq \varepsilon_{k}+\eta_{k-1} \tag{2.3}
\end{equation*}
$$

By the very construction $Y_{k}$ is $\mathfrak{X}_{k}^{\prime}$-measurable, $\mathfrak{L}\left(Y_{k}\right)=\mathfrak{L}\left(X_{k}\right)$ and $Y_{k}$ is independent of $Y_{1}, \ldots, Y_{k-1}$.

Now we consider the general case: Fix $k \geqq 1$. Given $\delta>0$ we can construct a partition $\left(E_{i}\right)_{i \geqq 1}$ of the separable space $S_{k}$ into Borel sets such that $\operatorname{diam}\left(E_{i}\right)<\delta(i \geqq 1)$. We do this for $\delta=\min \left(\varepsilon_{k}^{2} / 4, \rho_{k} / 4\right)$ and choose a point $u_{i} \in E_{i}$ $(i \geqq 1)$. The discrete random variable $X_{k}^{\prime}$ defined by $X_{k}^{\prime}(\omega)=u_{i}$ for $\omega \in X_{k}^{-1}\left(E_{i}\right)$ is $\mathfrak{F}_{k}^{\prime}$-measurable and satisfies $\sigma_{k}\left(X_{k}, X_{k}^{\prime}\right)<\delta$. Now let be given a set $A \subset D_{k-1}$, $A \in \mathfrak{A}_{k-1}, P(A)>0$ and a Borel probability measure $\lambda$ on $S_{k} \times S_{k}$ with marginals $\mathscr{U}\left(X_{k} \mid A\right)$ and $\mathscr{L}\left(X_{k}\right)$ such that

$$
\begin{equation*}
\int \sigma_{k}(u, v) d \lambda(u, v) \leqq \varepsilon_{k}^{2} / 2 \tag{2.4}
\end{equation*}
$$

We get a probability measure $\lambda^{\prime}$ on $S_{k} \times S_{k}$ concentrated on $\left\{\left(u_{i}, u_{j}\right) \mid i, j \geqq 1\right\}$ if we define $\lambda^{\prime}\left(u_{i}, u_{j}\right)=\lambda\left(E_{i} \times E_{j}\right)$. $\lambda^{\prime}$ has marginals $\mathcal{L}\left(X_{k}^{\prime} \mid A\right)$ and $\mathscr{L}\left(X_{k}^{\prime}\right)$. Furthermore using (2.4)

$$
\begin{aligned}
& \int \sigma_{k}(u, v) d \lambda^{\prime}(u, v)=\sum_{i, j} \int_{E_{i} \times E_{j}} \sigma_{k}\left(u_{i}, u_{j}\right) d \lambda(u, v) \\
& \leqq \sum_{i, j} \int_{E_{i} \times E_{j}}\left(\sigma_{k}(u, v)+\varepsilon_{k}^{2} / 2\right) d \lambda(u, v) \\
& =\int \sigma_{k} d \lambda+\varepsilon_{k}^{2} / 2 \leqq \varepsilon_{k}^{2} .
\end{aligned}
$$

Thus we have constructed a sequence $\left(X_{k}^{\prime}\right)_{k \geq 1}$ of discrete, $\mathfrak{F}_{k}^{\prime}$-measurable random variables satisfying $\sigma_{k}\left(X_{k}, X_{k}^{\prime}\right)<\min \left(\varepsilon_{k}^{2} / 4, \rho_{k} / 4\right)$ and

$$
\inf _{\lambda^{\prime}} \int \sigma_{k}(u, v) d \lambda^{\prime}(u, v) \leqq \varepsilon_{k}^{2}
$$

where infimum is taken over all Borel probability measures $\lambda^{\prime}$ on $S_{k} \times S_{k}$ with marginals $\mathfrak{L}\left(X_{k}^{\prime} \mid A\right)$ and $\mathscr{L}\left(X_{k}^{\prime}\right)$ for sets $A$ as considered above. By the first part of the proof there is a sequence $\left(Y_{k}^{\prime}\right)_{k \geqq 1}$ of independent random variables such that $\mathscr{L}\left(Y_{k}^{\prime}\right)=\mathscr{L}\left(X_{k}^{\prime}\right)$ and

$$
P\left[\sigma_{k}\left(X_{k}^{\prime}, Y_{k}^{\prime}\right) \geqq \varepsilon_{k}\right] \leqq \varepsilon_{k}+\eta_{k-1}
$$

The proof is finished following exactly Philipp's arguments in [16] p. 176. We just indicate the changes in the quantities $\varepsilon_{k}, \eta_{k}, \rho_{k}, \delta_{k}$. By (1.3) we get for any Borel set $E \subset S_{k}$

$$
\mathfrak{L}\left(Y_{k}^{\prime}\right)(E) \leqq G_{k}\left(E^{5 \rho_{k} / 4}\right)+\delta_{k} .
$$

Using the Strassen-Dudley theorem one constructs a probability measure $\lambda^{\prime \prime}$ on $S_{k} \times S_{k}$ with marginals $\mathscr{E}\left(Y_{k}^{\prime}\right)$ and $G_{k}$ such that

$$
\lambda^{\prime \prime}\left(\left\{(u, v) \mid \sigma_{k}(u, v) \geqq 3 \rho_{k} / 2\right\}\right) \leqq \delta_{k} .
$$

This implies that there is a sequence $\left(Y_{k}\right)_{k \geqq 1}$ of independent random variables, $\mathcal{L}\left(Y_{k}\right)=G_{k}$ and

$$
\begin{aligned}
& P\left[\sigma_{k}\left(X_{k}, Y_{k}\right) \geqq 2 \rho_{k}+\varepsilon_{k}\right] \\
& \leqq P\left[\sigma_{k}\left(X_{k}, X_{k}^{\prime}\right) \geqq \rho_{k} / 4\right]+P\left[\sigma_{k}\left(X_{k}^{\prime}, Y_{k}^{\prime}\right) \geqq \varepsilon_{k}\right]+P\left[\sigma_{k}\left(Y_{k}^{\prime}, Y_{k}\right) \geqq 3 \rho_{k} / 2\right] \\
& \leqq \varepsilon_{k}+\eta_{k-1}+\delta_{k} .
\end{aligned}
$$

## 3. Inequalities for Moments

In this section if not specified more precisely $B$ will mean a separable Banach space with norm $\|\cdot\|, X=\left(x_{k}\right)_{k \in Z}$ a strictly stationary $B$-valued stochastic process defined on a probability space $(\Omega, \mathfrak{U}, P)$ and $\left(\mathcal{F}_{k}\right)_{k \in Z}$ a fixed family of countably generated sub- $\sigma$-algebras of $\mathfrak{U}$ such that $x_{k}$ is $\mathfrak{F}_{k}$-measurable ( $k \in Z$ ). For $m \leqq n$ we shall write $\mathfrak{A}_{m}^{n}$ for the $\sigma$-algebra generated by $\left\{\mathfrak{F}_{k} \mid m \leqq k \leqq n\right\}$. Given a set $A \in \mathfrak{A}, P(A)>0, P(. \mid A)$ denotes conditional probability given $A$.

Besides of the Banach space norm $\|\cdot\|$ various other norms will be used. $\|\cdot\|_{*}$ denotes the norm in $B^{*}$ and $\|\cdot\|_{r}$ stands for $L^{r}$-norm, i.e. if $\int\|x\|^{r} d P<\infty$ for some $B$-valued random variable $x$ we write

$$
\|x\|_{r}=\left(\int\|x\|^{r} d P\right)^{1 / r} \quad(r \geqq 1) .
$$

$\|\cdot\|_{r}$ will also be used for the usual $L^{r}$-norm, i.e. if $\int|h|^{r} d P<\infty$ for some realvalued random variable $h$ then

$$
\|h\|_{r}=\left(\oint|h|^{r} d P\right)^{1 / r} \quad(r \geqq 1) .
$$

Lemma 3.1. Let $h$ be a real-valued random variable measurable with respect to $\mathfrak{H}_{-m}^{0}$ for some $m \geqq 0$ and having finite $r$-th absolute moment for some $r>1$. Let $g$ be a uniformly continuous, bounded function $g: B^{n} \rightarrow \mathbb{R}$ with modulus of continuity $c(\tau)$ and bound $|g| \leqq C$. Then, if $X=\left(x_{k}\right)_{k \in Z}$ is very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$,

$$
\begin{align*}
& \left|E\left[h g\left(x_{1}, \ldots, x_{n}\right)\right]-E[h] E\left[g\left(x_{1}, \ldots, x_{n}\right)\right]\right| \\
& \leqq\|h\|_{r}\left(\left|c(\tau)+2 C \varepsilon(n) \tau^{-1}\right|^{s}+|2 C|^{s} \varepsilon(n)\right)^{1 / s} . \tag{3.1}
\end{align*}
$$

Here $s$ is defined by $r^{-1}+s^{-1}=1$ and $\tau>0$ is any real number.
Proof. We may assume that $h$ is of the form $h=\sum_{i=1}^{l} a_{i} 1_{A_{i}}$ for a partition of sets $A_{i} \in \mathfrak{A l}_{-m}^{0}$ satisfying $P\left(A_{i}\right)>0$. Then using Hölder's inequality

$$
\begin{aligned}
& \left|E\left[h g\left(x_{1}, \ldots, x_{n}\right)\right]-E[h] E\left[g\left(x_{1}, \ldots, x_{n}\right)\right]\right| \\
& \left.=\mid \sum_{i=1}^{l} a_{i} \iint g\left(x_{1}, \ldots, x_{n}\right) d P\left(. \mid A_{i}\right)-\int g\left(x_{1}, \ldots, x_{n}\right) d P\right\} P\left(A_{i}\right) \mid \\
& \leqq\|h\|_{r}\left(\sum_{i=1}^{l}\left|\int g\left(x_{1}, \ldots, x_{n}\right) d P\left(. \mid A_{i}\right)-\int g\left(x_{1}, \ldots, x_{n}\right) d P\right|^{s} P\left(A_{i}\right)\right)^{1 / s} .
\end{aligned}
$$

Since $D$ chosen according to Definition 1 is an element of $\mathfrak{R}_{-m}^{0}$ we can assume - taking intersections if necessary - that $A_{i} \subset D$ or $A_{i} \subset D^{\subset}$ for each $i$.

Consider the case $A_{i} \subset D$. Denote $\mathfrak{L}\left(\left(x_{1}, \ldots, x_{n}\right)\right)$ by $\mu$ and $\mathfrak{L}\left(\left(x_{1}, \ldots, x_{n}\right) \mid A_{i}\right)$ by $\bar{\mu}$. Furthermore let $\lambda$ be a measure such that

$$
\begin{equation*}
\bar{\rho}_{n}\left(X, X \mid A_{i}\right)=\int \sigma_{n} d \lambda \tag{3.2}
\end{equation*}
$$

then

$$
\begin{aligned}
& \left|\int g\left(x_{1}, \ldots, x_{n}\right) d P\left(\cdot \mid A_{i}\right)-\int g\left(x_{1}, \ldots, x_{n}\right) d P\right| \\
& =\left|\int_{B^{n}} g(v) d \bar{\mu}(v)-\int_{B^{n}} g(u) d \mu(u)\right| \\
& \leqq \int_{B^{n} \times B^{n}}|g(u)-g(v)| d \lambda(u, v) .
\end{aligned}
$$

Splitting this integral in the two parts where $\left\{\sigma_{n} \leqq \tau\right\}$ and $\left\{\sigma_{n}>\tau\right\}$ we get the bound $c(\tau)+2 C \lambda\left(\left\{\sigma_{n}>\tau\right\}\right)$. By the Čebyšev-Markov inequality and (1.10)

$$
\begin{equation*}
\lambda\left(\left\{\sigma_{n}>\tau\right\}\right) \leqq \tau^{-1} \int \sigma_{n} d \lambda \leqq \tau^{-1} \varepsilon(n) . \tag{3.3}
\end{equation*}
$$

Therefore

$$
\begin{aligned}
& \left|E\left[h g\left(x_{1}, \ldots, x_{n}\right)\right]-E[h] E\left[g\left(x_{1}, \ldots, x_{n}\right)\right]\right| \\
& \leqq\|h\|_{r}\left(\sum_{A_{i} \in D}\left|c(\tau)+2 C \tau^{-1} \varepsilon(n)\right|^{s} P\left(A_{i}\right)+\sum_{A_{i} \subset D C}|2 C|^{s} P\left(A_{i}\right)\right)^{1 / s} .
\end{aligned}
$$

(3.1) follows since $P\left(D^{\complement}\right) \leqq \varepsilon(n)$.

Considering less general functions $g$ will allow us to deduce sharper bounds. In the limit theorems to be proved $g$ enters always in the form $g\left(\sum_{i=1}^{n} x_{i}\right)$.
Write Write

$$
S_{n}=\sum_{i=1}^{n} x_{i}
$$

Lemma 3.2. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary B-valued process which is very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$. Suppose $\left\|x_{0}\right\| \leqq C$ with probability 1. Let $h$ be a realvalued random variable measurable with respect to $\mathfrak{H}_{-m}^{0}$ for some $m \geqq 0$ and having finite $r$-th absolute moment for some $r>1$. Then for $g \in B^{*}$ and $n \geqq 1$

$$
\begin{align*}
& \left|E\left[h g\left(S_{n}\right)\right]-E[h] E\left[g\left(S_{n}\right)\right]\right| \\
& \leqq\|h\|_{r}\|g\|_{*}\left((n \varepsilon(n))^{s}+(2 n C)^{s} \varepsilon(n)\right)^{1 / s} \tag{3.4}
\end{align*}
$$

where $s$ is defined by $r^{-1}+s^{-1}=1$. If $X$ is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ and $x_{0}$ is integrable then for $h$ and $g$ as above

$$
\begin{equation*}
\left|E\left[h g\left(S_{n}\right)\right]-E[h] E\left[g\left(S_{n}\right)\right]\right| \leqq\|h\|_{r}\|g\|_{*} n \varepsilon(n) \tag{3.5}
\end{equation*}
$$

Proof. Let $D=D(n, m, 0) \in \mathfrak{H}_{-m}^{0}$ be a set chosen according to Definition 1 then for any $A \in \mathfrak{A}_{-m}^{0}, A \subset D, P(A)>0$

$$
\begin{aligned}
& \left|\int g\left(S_{n}\right) d P(\cdot \mid A)-\int g\left(S_{n}\right) d P\right| \\
& =\left|\int_{B^{n} \times B^{n}}\left(g\left(\sum_{i=1}^{n} u_{i}\right)-g\left(\sum_{i=1}^{n} v_{i}\right)\right) d \lambda\left(\left(u_{1}, \ldots, u_{n}\right),\left(v_{1}, \ldots, v_{n}\right)\right)\right| \\
& \leqq\|g\|_{*} \int_{B^{n} \times B^{B^{n}}} \sum_{i=1}^{n}\left\|u_{i}-v_{i}\right\| d \lambda\left(\left(u_{1}, \ldots, u_{n}\right),\left(v_{1}, \ldots, v_{n}\right)\right) \\
& =\|g\|_{*} \int_{B^{n} \times B^{n}} n \sigma_{n}(u, v) d \lambda(u, v) \leqq\|g\|_{*} n \varepsilon(n) .
\end{aligned}
$$

Here again $\lambda$ is a measure satisfying (3.2). Now we take the same representation for $h$ as in the proof before and apply Hölder's inequality, then

$$
\begin{aligned}
& \left|E\left[h g\left(S_{n}\right)\right]-E[h] E\left[g\left(S_{n}\right)\right]\right| \\
& \leqq\|h\|_{r}\left(\sum_{i=1}^{l}\left|\int g\left(S_{n}\right) d P\left(. \mid A_{i}\right)-\int g\left(S_{n}\right) d P\right|^{s} P\left(A_{i}\right)\right)^{1 / s} .
\end{aligned}
$$

(3.5) follows immediately since if $X$ is strictly very weak Bernoulli, the set $D$ has probability 1 , i.e. the estimate above applies to each of the $A_{i}$. In order to derive (3.4) we split the sum according to $A_{i} \subset D$ or $A_{i} \subset D^{\subset}$ and use the trivial estimate $\left|g\left(S_{n}\right)\right| \leqq\|g\|_{*} n C$ in the latter case. Thus we get the bound

$$
\|h\|_{r}\left\{\sum_{A_{i} \in D}\left(\|g\|_{*} n \varepsilon(n)\right)^{s} P\left(A_{i}\right)+\sum_{A_{i} \in D}\left(2\|g\|_{*} n C\right)^{s} P\left(A_{i}\right)\right\}^{1 / s}
$$

(3.4) follows since $P\left(D^{C}\right) \leqq \varepsilon(n)$.

Lemma 3.3. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ has mean zero and finite variance then for any $N \geqq 1$ and $M=2 N$ or $M=2 N+1$

$$
\begin{equation*}
\left(2\left(1-c_{M}\right)\right)^{1 / 2}\left\|S_{N}\right\|_{2} \leqq\left\|S_{M}\right\|_{2} \leqq\left\|S_{N}\right\|_{2}\left(2\left(1+c_{M}\right)\right)^{1 / 2} \tag{3.6}
\end{equation*}
$$

where

$$
c_{2 N}=C_{0} /\left\|S_{N}\right\|_{2}, \quad c_{2 N+1}=2 C_{0} /\left\|S_{N}\right\|_{2}+\left\|x_{0}\right\|_{2}^{2} / 2\left\|S_{N}\right\|_{2}^{2}
$$

and $C_{0}$ is the constant implied by $<$ in (1.12).
Proof. Consider the case $M=2 N$ first. We use Lemma 3.2 with $r=2$. Then (3.5) yields

$$
\begin{equation*}
\left|E\left[\left(\sum_{i=-N+1}^{0} x_{i}\right)\left(\sum_{i=1}^{N} x_{i}\right)\right]\right| \leqq\left\|S_{N}\right\|_{2} C_{0} \tag{3.7}
\end{equation*}
$$

Therefore by stationarity

$$
\begin{aligned}
\left\|S_{2 N}\right\|_{2}^{2} & =2\left\|S_{N}\right\|_{2}^{2}+2 E\left[\left(\sum_{i=1}^{N} x_{i}\right)\left(\sum_{i=N+1}^{2 N} x_{i}\right)\right] \\
& \leqq 2\left\|S_{N}\right\|_{2}^{2}+2\left\|S_{N}\right\|_{2} C_{0}
\end{aligned}
$$

This gives the upper inequality. The lower inequality follows in the same way if we use the representation $S_{N}=S_{2 N}-\left(S_{2 N}-S_{N}\right)$.

Now let $M=2 N+1$. Applying again (3.5) to $E\left[\left(\sum_{i=-N+1}^{0} x_{i}\right)\left(\sum_{i=1}^{N+1} x_{i}\right)\right]$ and to $E\left[\left(\sum_{i=-N+1}^{0} x_{i}\right) x_{1}\right]$ we get

$$
\begin{aligned}
\left\|S_{2 N+1}\right\|_{2}^{2}= & 2\left\|S_{N}\right\|_{2}^{2}+\left\|x_{2 N+1}\right\|_{2}^{2}+2 E\left[\left(\sum_{i=1}^{N} x_{i}\right)\left(\sum_{i=N+1}^{2 N+1} x_{i}\right)\right] \\
& +2 E\left[\left(\sum_{i=N+1}^{2 N} x_{i}\right) x_{2 N+1}\right] \\
& \leqq 2\left\|S_{N}\right\|_{2}^{2}+\left\|x_{0}\right\|_{2}^{2}+4\left\|S_{N}\right\|_{2} C_{0} \\
= & \left\|S_{N}\right\|_{2}^{2}\left(2\left(1+2 C_{0} /\left\|S_{N}\right\|_{2}+\left\|x_{0}\right\|_{2}^{2} / 2\left\|S_{N}\right\|_{2}^{2}\right)\right)
\end{aligned}
$$

The lower bound in (3.6) can be derived similarly.
Lemma 3.4. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ has mean zero, finite variance and $\lim _{n \rightarrow \infty} \operatorname{Var}\left(S_{n}\right)=\infty$ then

$$
\begin{equation*}
\operatorname{Var}\left(S_{n}\right)=n h(n) \tag{3.8}
\end{equation*}
$$

where $h(n)$ is a slowly varying function of the integral variable $n$.
Proof. We have to show that for every $k \geqq 1$

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \operatorname{Var}\left(S_{n k}\right) / \operatorname{Var}\left(S_{n}\right)=k \tag{3.9}
\end{equation*}
$$

If we define $y_{j}=\sum_{i=1}^{n} x_{(j-1) n+i}$ for $j=1, \ldots, k$ then by stationarity

$$
\operatorname{Var}\left(S_{n k}\right)=k \operatorname{Var}\left(S_{n}\right)+2 \sum_{t=1}^{k-1} E\left[y_{1}\left(y_{2}+\ldots+y_{l+1}\right)\right]
$$

By (3.5)

$$
\left|E\left[y_{1}\left(y_{2}+\ldots+y_{l+1}\right)\right]\right| \leqq \sum_{k=1}^{n} x_{k} \|_{2} \ln \varepsilon(\ln ) \leqq\left(\operatorname{Var}\left(S_{n}\right)\right)^{1 / 2} C_{0}
$$

Therefore

$$
\operatorname{Var}\left(S_{n k}\right)=k \operatorname{Var}\left(S_{n}\right)+o\left(\operatorname{Var}\left(S_{n}\right)\right)
$$

Proposition 3.5. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ has mean zero, finite variance and $\lim \operatorname{Var}\left(S_{n}\right)=\infty$ then there exists a finite positive value $\sigma^{2}$ such that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} n^{-1} \operatorname{Var}\left(S_{n}\right)=\sigma^{2} \tag{3.10}
\end{equation*}
$$

Before proving this Proposition we restate a lemma from Bradley [3] (Lemma 2). Define under the assumptions of the Proposition above

$$
\begin{equation*}
g(n)=n^{-1 / 2}\left\|S_{n}\right\|_{2} \tag{3.11}
\end{equation*}
$$

then the following holds

Lemma 3.6. Given any $\varepsilon>0$ and any positive integer $L_{0}$, there are positive integers $N$ and $L$ with $L \geqq L_{0}$ such that for all $l, L \leqq l \leqq 2 L$ one has

$$
(1-\varepsilon) g(N) \leqq g(l) \leqq(1+\varepsilon) g(N) .
$$

The proof of this lemma is the same as in [3], since the basic ingredient used there, namely

$$
\begin{equation*}
\lim _{N \rightarrow \infty} g(m N) / g(N)=1 \tag{3.12}
\end{equation*}
$$

for all $m \geqq 1$ holds under our assumptions by Lemma 3.4.
Again as in [3] we choose $0<A<1$ sufficiently small that if $\left(a_{n}\right)_{n \geqq 1}$ is any sequence of real numbers such that $\sum_{n \geqq 1}\left|a_{n}\right| \leqq A$ then

$$
\left|1-\prod_{n \geqq 1}\left(1+a_{n}\right)\right| \leqq 2 \sum_{n \geqq 1}\left|a_{n}\right| \quad \text { and } \quad\left|1-\prod_{n \geqq 1}\left(1+a_{n}\right)^{-1}\right| \leqq 2 \sum_{n \geqq 1}\left|a_{n}\right| .
$$

Proof of Proposition 3.5. Let $0<\varepsilon<A$ be given and denote by $C_{0}$ the constant implied by $\ll$ in (1.12). Let $r \geqq 2$ be a fixed integer and $C>0$ a constant to be determined later. Since $\lim _{n \rightarrow \infty} \operatorname{Var}\left(S_{n}\right)=\infty$ we can choose $L_{0} \geqq 0$ sufficiently large such that $L_{0}^{-1}<\varepsilon / 4$ and for all $N \geqq L_{0}$

$$
\begin{gather*}
1-2 C_{0} /\left\|S_{N}\right\|_{2}-\left\|x_{0}\right\|_{2}^{2} / 2\left\|S_{N}\right\|_{2}^{2}>2^{-1 / r},  \tag{3.13}\\
\left\|S_{N}\right\|_{2}>C \tag{3.14}
\end{gather*}
$$

and

$$
\begin{equation*}
\left\|x_{0}\right\|_{2}^{2} / 2\left\|S_{N}\right\|_{2} \leqq C_{0} \tag{3.15}
\end{equation*}
$$

By Lemma 3.6 there exist positive integers $H$ and $L$ with $L \geqq L_{0}$ such that for all $l, L \leqq l \leqq 2 L$ one has

$$
(1-\varepsilon) g(H) \leqq g(l) \leqq(1+\varepsilon) g(H)
$$

Now let $m$ be any integer $m \geqq 2 L$ then $2^{M} L \leqq m \leqq 2^{M+1} L$ for some $M$. There exist integers $J_{0}, J_{1}, \ldots, J_{M}$ such that $m=J_{M}, L \leqq J_{0} \leqq 2 L$ and for all $n=0$, $1, \ldots, M-1, J_{n+1}=2 J_{n}$ or $J_{n+1}=2 J_{n}+1$. By (3.6)

$$
g\left(J_{n+1}\right) \leqq g\left(J_{n}\right)\left(1+3 C_{0} /\left\|S_{J_{n}}\right\|_{2}\right)^{1 / 2}
$$

where we made use of (3.15). Introducing (3.13) in the lower half of (3.6) and making use of (3.14) we see that $\left\|S_{J_{n}}\right\|_{2} \geqq 2^{n(r-1) / 2 r} C$. Therefore

$$
g\left(J_{n+1}\right) \leqq g\left(J_{n}\right)\left(1+3 C_{0} / C 2^{n(r-1) / 2 r}\right)^{1 / 2}
$$

and

$$
g(m) \leqq g\left(J_{0}\right)\left(\prod_{n=0}^{\infty}\left(1+3 C_{0} / C 2^{n(r-1) / 2 r}\right)\right)^{1 / 2}
$$

Now we choose $C$ such that $\sum_{n=0}^{\infty} 3 C_{0} / C 2^{n(r-1) / 2 r}<\varepsilon / 4$ then

$$
\begin{equation*}
g(m) \leqq g\left(J_{0}\right)(1+\varepsilon)^{1 / 2} \leqq g(H)(1+\varepsilon)^{2} \tag{3.16}
\end{equation*}
$$

Using the lower half of (3.6) we get

$$
\begin{aligned}
g\left(J_{n+1}\right) & \geqq g\left(J_{n}\right)\left(2 J_{n} / J_{n+1}\right)^{1 / 2}\left(1-3 C_{0} /\left\|S_{J_{n}}\right\|_{2}\right)^{1 / 2} \\
& \geqq g\left(J_{n}\right)\left(1-2^{-(n+1)} L_{0}^{-1}\right)\left(1-3 C_{0} / C 2^{n(r-1) / 2 r}\right)^{1 / 2}
\end{aligned}
$$

or

$$
g(m) \geqq g\left(J_{0}\right)\left(\prod_{n=0}^{\infty}\left(1-2^{-(n+1)} L_{0}^{-1}\right)\left(1-3 C_{0} / C 2^{n(r-1) / 2 r}\right)\right)^{1 / 2}
$$

By the choice of $L_{0}$ and C

$$
\begin{equation*}
g(m) \geqq g\left(J_{0}\right)(1-\varepsilon) \geqq g(H)(1-\varepsilon)^{2} \tag{3.17}
\end{equation*}
$$

This together with (3.16) shows that $\lim _{m \rightarrow \infty} g(m)$ exists, which proves (3.10).
Proposition 3.7. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ has mean zero, finite variance and $\lim _{n \rightarrow \infty} \operatorname{Var}\left(S_{n}\right)=\infty$ then

$$
\begin{equation*}
m^{-1} \operatorname{Var}\left(S_{m}\right)-\sigma^{2} \ll m^{-1 / 2} \tag{3.18}
\end{equation*}
$$

where $\sigma^{2}=\lim _{n \rightarrow \infty} n^{-1} \operatorname{Var}\left(S_{n}\right)$.
Proof. Since $\lim _{m \rightarrow \infty} g(m)$ exists and is positive by (3.10) we can choose constants $0<C_{1}<C_{2}$ such that $C_{1} \leqq g(m) \leqq C_{2}$ for all $m$. Let $C_{0}$ be the same constant as before and let $L$ be a positive integer such that $C_{0}\left(\sum_{n=0}^{\infty} 2^{-n / 2}\right) / C_{1} L<A$. Given an integer $m$ we write $2^{M} L^{2} \leqq m \leqq 2^{M+1} L^{2}$ and for each $n=0,1, \ldots$ we set $J_{n}$ $=2^{n} m$. By (3.6)

$$
g\left(J_{n+1}\right) \leqq g\left(J_{n}\right)\left(1+C_{0} / C_{1} L 2^{(n+M) / 2}\right)^{1 / 2}
$$

Since $\lim _{m \rightarrow \infty} g(m)^{2}=\sigma^{2}$ this implies

$$
g(m)^{2} \geqq \sigma^{2} \prod_{n=0}^{\infty}\left(1+C_{0} / C_{1} L 2^{(n+M) / 2}\right)^{-1}
$$

Using the lower half of (3.6) we get

$$
g(m)^{2} \leqq \sigma^{2} \prod_{n=0}^{\infty}\left(1-C_{0} / C_{1} L 2^{(n+M) / 2}\right)^{-1}
$$

Both estimates together yield

$$
\left|g(m)^{2}-\sigma^{2}\right| \leqq 2 \sigma^{2} \sum_{n=0}^{\infty} C_{0} / C_{1} L 2^{(n+M) / 2}=O\left(m^{-1 / 2}\right)
$$

Lemma 3.8. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ is centered and
bounded with probability 1 and $\lim _{n \rightarrow \infty} \operatorname{Var}\left(S_{n}\right)=\infty$ then for any $0 \leqq \delta<1$

$$
\begin{equation*}
E\left[\left|\sum_{k=1}^{n} x_{k}\right|^{2+\delta}\right] \ll n^{1+\delta / 2} \tag{3.19}
\end{equation*}
$$

Proof. Write for short $a_{n}=E\left[\left|S_{n}\right|^{2+\delta}\right]$ and $\sigma_{n}=E\left[\left|S_{n}\right|^{2}\right]^{1 / 2}$ for any $n \geqq 1$. From Proposition 3.5 we know that

$$
\begin{equation*}
\sigma_{n}^{2}=\sigma^{2} n(1+o(1)) \tag{3.20}
\end{equation*}
$$

for some $\sigma^{2}>0$. Besides of $S_{n}=\sum_{i=1}^{n} x_{i}$ we use $\tilde{S}_{n}=\sum_{i=n+1}^{2 n} x_{i}$ then by stationarity we get as in Doob [8], p. 226

$$
\begin{align*}
a_{2 n} \leqq & 2 a_{n}+E\left[S_{n}^{2}\left|\tilde{S}_{n}\right|^{\delta}\right]+E\left[\mid S_{n} \tilde{\delta}^{\delta} \tilde{S}_{n}^{2}\right]  \tag{3.21}\\
& +2 E\left[\left|S_{n}\right|^{1+\delta}\left|\tilde{S}_{n}\right|\right]+2 E\left[\left|S_{n}\right|\left|\tilde{S}_{n}\right|^{1+\delta}\right] .
\end{align*}
$$

Let $C_{0}$ again be the constant implied by $\ll$ in (1.12). If $A \in \mathfrak{M d}_{-n+1}^{0}, P(A)>0$ and $\lambda$ is a probability measure satisfying (3.2) then

$$
\begin{align*}
& \left|\int\right| S_{n}\left|d P(\cdot \mid A)-\int\right| S_{n}|d P| \\
& =\left|\int_{\mathbb{R}^{n} \times \mathbb{R}^{n}}\left(\left|\sum_{k=1}^{n} u_{k}\right|-\left|\sum_{k=1}^{n} v_{k}\right|\right) d \lambda(u, v)\right|  \tag{3.22}\\
& \leqq \int n \sigma_{n}(u, v) d \lambda(u, v) \leqq n \varepsilon(n) \leqq C_{0} .
\end{align*}
$$

We may assume that

$$
\begin{equation*}
\left|\sum_{k=-n+1}^{0} x_{k}\right|=\sum_{i=1}^{l} a_{i} 1_{A_{i}} \tag{3.23}
\end{equation*}
$$

for a partition of sets $A_{i} \in \mathfrak{H}_{-n+1}^{0}$ satisfying $P\left(A_{i}\right)>0$. Then by stationarity and Hölder's inequality

$$
\begin{aligned}
& \left|E\left[\left|S_{n}\right|^{1+\delta}\left|\tilde{S}_{n}\right|\right]-E\left[\left|S_{n}\right|^{1+\delta}\right] E\left[\left|\tilde{S}_{n}\right|\right]\right| \\
& =\left|\sum_{i=1}^{l} a_{i}^{1+\delta}\left\{\int\left|S_{n}\right| d P\left(. \mid A_{i}\right)-\int\left|S_{n}\right| d P\right\} P\left(A_{i}\right)\right| \\
& \leqq \sigma_{n}^{1+\delta}\left\{\sum_{i=1}^{l}\left|\int\right| S_{n}\left|d P\left(. \mid A_{i}\right)-\int\right| S_{n}|d P|^{2 /(1-\delta)} P\left(A_{i}\right)\right\}^{(1-\delta) / 2} \leqq \sigma_{n}^{2+\delta} C_{0} .
\end{aligned}
$$

In the last inequality we assumed $\sigma_{n} \geqq 1$ which is true for sufficiently large $n$. Since $E\left[\left|S_{n}\right|^{1+\delta}\right] \leqq \sigma_{n}^{1+\delta}$ and $E\left[\left|\tilde{S_{n}}\right|\right] \leqq \sigma_{n}$ we get

$$
\begin{equation*}
E\left[\left|S_{n}\right|^{1+\delta}\left|\tilde{S}_{n}\right|\right] \leqq\left(C_{0}+1\right) \sigma_{n}^{2+\delta} \tag{3.24}
\end{equation*}
$$

Let $C$ be the bound for $x_{0}$ then $\left|\sum_{k=1}^{n} x_{k}\right| \leqq n C$ and we can apply the elementary inequality $\left|y^{2}-z^{2}\right| \leqq 2 n C|y-z|$ which holds for $0 \leqq y, z \leqq n C$. Then similar to (3.22)

$$
\begin{aligned}
& \left.\left|\int\right| S_{n}\right|^{2} d P(. \mid A)-\int\left|S_{n}\right|^{2} d P \mid \\
& \leqq 2 n C \int_{\mathbb{R}^{n} \times \mathbb{R}^{n}}\left\|\sum_{k=1}^{n} u_{k}|-| \sum_{k=1}^{n} v_{k}\right\| d \lambda(u, v) \leqq 2 n C C_{0} .
\end{aligned}
$$

We introduce this in the following estimate where (3.23) and Hölder's inequality are used

$$
\begin{aligned}
& \left|E\left[\left|S_{n}\right|^{\delta}\left|\tilde{S}_{n}\right|^{2}\right]-E\left[\left|S_{n}\right|^{\delta}\right] E\left[\left|\tilde{S}_{n}\right|^{2}\right]\right| \\
& \leqq\left.\sum_{i=1}^{l} a_{i}^{\delta}\left|\int\right| S_{n}\right|^{2} d P\left(. \mid A_{i}\right)-\int\left|S_{n}\right|^{2} d P \mid P\left(A_{i}\right) \leqq \sigma_{n}^{\delta} 2 n C C_{0}
\end{aligned}
$$

By (3.20) the last term is bounded by $\sigma_{n}^{2+\delta} 4 C C_{0} \sigma^{-2}$ for large $n$. Since $E\left[\left|S_{n}\right|^{\delta}\right]$ $\leqq \sigma_{n}^{\delta}$ and $E\left[\left|\tilde{S}_{n}\right|^{2}\right] \leqq \sigma_{n}^{2}$ we deduce

$$
\begin{equation*}
E\left[\left|S_{n}\right|^{\delta}\left|\tilde{S}_{n}\right|^{2}\right] \leqq\left(1+4 C C_{0} \sigma^{-2}\right) \sigma_{n}^{2+\delta} . \tag{3.25}
\end{equation*}
$$

The following inequality $\left||y|^{1+\delta}-|z|^{1+\delta}\right| \leqq(1+\delta)(n C)^{\delta}| | x|-|y||$ which holds for $0 \leqq y, z \leqq n C$ implies together with (3.22)

$$
\left.\left|\int\right| S_{n}\right|^{1+\delta} d P(. \mid A)-\int\left|S_{n}\right|^{1+\delta} d P \mid \leqq(1+\delta)(n C)^{\delta} C_{0} .
$$

Therefore by Hölder's inequality

$$
\begin{aligned}
& \left|E\left[\left|S_{n}\right|\left|\tilde{S}_{n}\right|^{1+\delta}\right]-E\left[\left|S_{n}\right|\right] E\left[\left|\tilde{S}_{n}\right|^{1+\delta}\right]\right| \\
& \leqq\left.\sum_{i=1}^{l}\left|a_{i}\right|\left|\int\right| S_{n}\right|^{1+\delta} d P\left(. \mid A_{i}\right)-\int\left|S_{n}\right|^{1+\delta} d P \mid P\left(A_{i}\right) \\
& \leqq \sigma_{n}(1+\delta)(n C)^{\delta} C_{0} \leqq \sigma_{n}^{1+2 \delta}(1+\delta) C^{\delta} C_{0} 2^{\delta} \sigma^{-2 \delta} .
\end{aligned}
$$

In the last line (3.20) is again used. Our assumptions $\lim _{n \rightarrow \infty} \sigma_{n}^{2}=\infty$ and $\delta<1$ imply $\lim _{n \rightarrow \infty} \sigma_{n}^{1-\delta}=\infty$. Therefore the last term above is less than $\sigma_{n}^{2+\delta}$ for $n$ large enough. We conclude

$$
\begin{equation*}
E\left[\left|S_{n}\right|\left|\tilde{S}_{n}\right|^{1+\delta}\right] \leqq 2 \sigma_{n}^{2+\delta} . \tag{3.26}
\end{equation*}
$$

Now we consider the remaining summand in (3.21).

$$
E\left[S_{n}^{2}\left|\tilde{S}_{n}\right|^{\delta}\right] \leqq \sigma_{n}^{2(1-\delta)}\left(E\left[\left|S_{n}\right|^{2}\left|\tilde{S}_{n}\right|\right]\right)^{\delta} .
$$

Using (3.22) we easily derive

$$
\left|E\left[S_{n}^{2}\left|\tilde{S}_{n}\right|\right]-E\left[S_{n}^{2}\right] E\left[\left|\tilde{S}_{n}\right|\right]\right| \leqq C_{0} \sigma_{n}^{2}
$$

Therefore

$$
E\left[S_{n}^{2} \mid \tilde{S}_{n} \|^{\delta}\right] \leqq \sigma_{n}^{2(1-\delta)}\left(\left(C_{0} \sigma_{n}^{2}\right)^{\delta}+\sigma_{n}^{3 \delta}\right)=C_{0}^{\delta} \sigma_{n}^{2}+\sigma_{n}^{2+\delta}
$$

The last term is bounded by $2 \sigma_{n}^{2+\delta}$ for large $n$. Thus

$$
\begin{equation*}
E\left[\left|S_{n}\right|^{2}\left|\tilde{S}_{n}\right|^{\delta}\right] \leqq 2 \sigma_{n}^{2+\delta} . \tag{3.27}
\end{equation*}
$$

Introducing (3.24)-(3.27) in (3.21) we get for some constant $C^{\prime}>0$

$$
a_{2 n} \leqq 2 a_{n}+C^{\prime} \sigma_{n}^{2+\delta}
$$

Now define $b_{n}=a_{n} / \sigma_{n}^{2+\delta}$ then since $\lim _{n \rightarrow \infty} \sigma_{2 n}^{2} / \sigma_{n}^{2}=2$ the last inequality implies that there exist constants $\lambda, 0<\lambda<1$ and $C^{\prime \prime}>0$ such that for all $n$

$$
b_{2 n} \leqq \lambda b_{n}+C^{\prime \prime}
$$

Hence

$$
\begin{equation*}
\sup _{r} b_{2^{r}} \leqq b_{1}+C^{\prime \prime}(1-\lambda)^{-1}<\infty \tag{3.28}
\end{equation*}
$$

The deduction of $a_{n} \leqq a \sigma_{n}^{2+\delta}$ for some constant $a>0$ from (3.28) is routine (see Doob [8] p. 227 or Ibragimov-Linnik [13] p.343). Proposition 3.5 has to be used here. This proves the lemma.

Let us mention that essentially by using (3.5) it is easy to prove the following inequality for higher moments
Remark. Let $X=\left(x_{k}\right)_{k \in Z}$ be a stationary real-valued process which is strictly very weak Bernoulli at rate $(\varepsilon(n))_{n \geqq 1}$ given by (1.12). Suppose $x_{0}$ is centered and bounded then for any integer $l \geqq 2$

$$
E\left[\left|\sum_{k=1}^{n} x_{k}\right|^{l}\right] \lll n^{l-1} .
$$

## 4. Proof of Theorem 2

The existence of $\sigma^{2}$ was shown in Proposition 3.5. Define for every $n \geqq 1, N$ $=N(n)=\left[n^{13 / 16}\right], l=l(n)=\left[n^{3 / 16}\right]$ and

$$
X_{j}=\sum_{v=(j-1) N+1}^{j N} x_{v} \quad(1 \leqq j \leqq l)
$$

Furthermore we consider the following sub- $\sigma$-algebras

$$
\mathfrak{F}_{j}^{\prime}=\bigvee_{v=(j-1) N+1}^{j N} \mathfrak{F}_{v} \quad(1 \leqq j \leqq l)
$$

and as in Theorem 1

$$
\mathfrak{A}_{k}^{\prime}=\bigvee_{j \leqq k} \mathfrak{F}_{j}^{\prime}=\bigvee_{v=1}^{k N} \mathfrak{F}_{v} \quad(1 \leqq k \leqq l)
$$

Let $k$ and $A \in \mathfrak{A}_{k-1}^{\prime}, P(A)>0$ be given. Then according to Definition 1

$$
\inf _{\lambda \in \mathfrak{P}_{N}} \int_{\mathbb{R}^{N} \times \mathbb{R}^{N}} \sigma_{N} d \lambda \leqq \varepsilon(N)
$$

where $\mathfrak{P}_{N}$ is the class of Borel probability measures on $\mathbb{R}^{N} \times \mathbb{R}^{N}$ with marginals $\mathscr{L}\left(\left(x_{(k-1) N+1}, \ldots, x_{k N}\right)\right)$ and $\mathscr{L}\left(\left(x_{(k-1) N+1}, \ldots, x_{k N}\right) \mid A\right)$. Let $\lambda \in \mathfrak{B}_{N}$ be a measure such that $\int \sigma_{N} d \lambda \leqq \varepsilon(N)$ and define $T: \mathbb{R}^{N} \times \mathbb{R}^{N} \rightarrow \mathbb{R} \times \mathbb{R}$ by $T\left(\left(u_{1}, \ldots, u_{N}\right)\right.$, $\left.\left(v_{1}, \ldots, v_{N}\right)\right)=\left(n^{-1 / 2} \sum_{i=1}^{N} u_{i}, n^{-1 / 2} \sum_{i=1}^{N} v_{i}\right)$. Then $T(\lambda)$ has marginals $\mathscr{L}\left(n^{-1 / 2} X_{k}\right)$ and $\mathscr{L}\left(n^{-1 / 2} X_{k} \mid A\right)$ and

$$
\begin{aligned}
& \int_{\mathbb{R} \times \mathbb{R}}|u-v| d T(\lambda)(u, v) \\
& =\int_{\mathbb{R}^{v} \times \mathbb{R}^{N}} n^{-1 / 2}\left|\sum_{i=1}^{N}\left(u_{i}-v_{i}\right)\right| d \lambda\left(\left(u_{1}, \ldots, u_{N}\right),\left(v_{1}, \ldots, v_{N}\right)\right) \\
& \leqq n^{-1 / 2} N \varepsilon(N) \leqq C_{0} n^{-1 / 2}
\end{aligned}
$$

where $C_{0}$ means the constant implied by $\ll$ in (1.12). Since constants are of no importance in these estimates we assume $C_{0}=1 / 2$. By Theorem 1 applied to the sequence $\left(n^{-1 / 2} X_{k}\right)_{1 \leqq k \leqq l}$ there exist independent random variables $n^{-1 / 2} Y_{k}$ such that $\mathcal{L}\left(n^{-1 / 2} X_{k}\right)=\mathfrak{Q}\left(n^{-1 / 2} Y_{k}\right)$ and

$$
P\left[\left|n^{-1 / 2} X_{k}-n^{-1 / 2} Y_{k}\right| \geqq n^{-1 / 4}\right] \leqq n^{-1 / 4}
$$

Note that we need here only the special case of Theorem 1 where $G_{k}=\mathscr{I}\left(X_{k}\right)$ and $\rho_{k}=\delta_{k}=\eta_{k}=0$. The last inequality implies

$$
\begin{equation*}
\pi\left(\mathscr{L}\left(n^{-1 / 2} \sum_{k=1}^{l} X_{k}\right), \mathscr{Q}\left(n^{-1 / 2} \sum_{k=1}^{l} Y_{k}\right)\right) \leqq l n^{-1 / 4} \leqq n^{-1 / 16} . \tag{4.1}
\end{equation*}
$$

From the definition of $X_{k}$ and Proposition 3.5 we derive

$$
\begin{aligned}
& E\left[\left|n^{-1 / 2} \sum_{k=1}^{l} X_{k}-n^{-1 / 2} \sum_{v=1}^{n} x_{v}\right|^{2}\right] \\
& =n^{-1} E\left[\left|\sum_{v=1}^{n-l N} x_{v}\right|^{2}\right] \ll n^{-1}(n-l N) \ll n^{-1}(l+N) \ll n^{-3 / 16}
\end{aligned}
$$

which immediately implies

$$
\begin{equation*}
\pi\left(\mathscr{P}\left(n^{-1 / 2} \sum_{k=1}^{l} X_{k}\right), \mathfrak{Q}\left(n^{-1 / 2} \sum_{v=1}^{n} x_{v}\right)\right) \ll n^{-1 / 16} \tag{4.2}
\end{equation*}
$$

The following estimation for the Prohorov distance of certain distributions is due to Yurinskii [22]. We restate it in slightly refined form from [4], Prop. 5.1.
Proposition. Let $X_{1}, \ldots, X_{n}$ be independent $\mathbb{R}^{d}$-valued random variables with mean zero and $E\left[\left\|X_{i}\right\|^{2+\delta^{\prime}}\right]<\infty$. If $\mu_{n}$ denotes the distribution of $n^{-1 / 2} \sum_{i=1}^{n} X_{i}$ and $v_{n}$ the Gaussian measure with mean zero and the same covariance as $\mu_{n}$ then

$$
\pi\left(\mu_{n}, v_{n}\right) \leqq c n^{-\delta^{\prime} / 9} d^{1 / 3}\left(\tilde{\rho}_{2+\delta^{\prime}}\right)^{2 / 9}
$$

where $\tilde{\rho}_{2+\delta^{\prime}}=n^{-1} \sum_{i=1}^{n} E\left[\left\|X_{i}\right\|^{2+\delta^{\prime}}\right]$ and $c$ is an absolute constant.

The bound actually given in [22] (and [4]) is more complicated. We will apply this Proposition to the random variables $(l / n)^{1 / 2} Y_{k}(1 \leqq k \leqq l)$. By Lemma 3.8

$$
E\left[\left|(l / n)^{1 / 2} Y_{k}\right|^{2+\delta^{\prime}}\right] \ll 1
$$

for any $0 \leqq \delta^{\prime}<1$. If $\tilde{\tau}_{N}^{2}$ denotes the variance of $(l / n)^{1 / 2} Y_{1}$ we can conclude using $\delta^{\prime}=3 / 4$

$$
\begin{equation*}
\pi\left(\mathfrak{P}\left(n^{-1 / 2} \sum_{k=1}^{l} Y_{k}\right), N\left(0, \tilde{\tau}_{N}^{2}\right)\right) \ll n^{-1 / 64} \tag{4.3}
\end{equation*}
$$

In the next two steps we use results of Dehling.
Theorem ([4], Th.7). Let $\mu$ and $v$ be two Gaussian measures on $\mathbb{R}^{d}$ with mean zero and covariance operators $S$ and $T$. Then the following estimation holds

$$
\pi(\mu, v) \leqq C\|\mid S-T\|^{1 / 3} d^{1 / 6}\left(1+\left|\log \|S-T\|^{-1} d\right|^{1 / 2}\right)
$$

where $C$ is an absolute constant and $\|\cdot\| \cdot \|$ is defined by

$$
\|R\|=\sup _{\left(e_{i}\right)} \sum_{i=1}^{d}\left|\left(R e_{i}, e_{i}\right)\right|
$$

and the sup is taken over all orthonormal bases $\left(e_{i}\right)$ for $\mathbb{R}^{d}$.
We denote by $\tau_{N}^{2}$ the variance of $N^{-1 / 2} Y_{1}$. Proposition 3.7 implies $\tau_{N}^{2}-\sigma^{2}$ $\ll N^{-1 / 2}$. So if we replace the estimation in the theorem above by $\pi(\mu, v) \leqq C_{0}$ $\|\mid S-T\|^{1 / 4} d^{1 / 4}$ we get

$$
\begin{equation*}
\pi\left(N\left(0, \tau_{N}^{2}\right), N\left(0, \sigma^{2}\right)\right) \ll N^{-1 / 8} \ll n^{-1 / 10} \tag{4.4}
\end{equation*}
$$

The last link in our chain of estimates is a direct consequence of
Lemma ([4], Lemma 2.1). Let $X$ and $Y$ be $\mathbb{R}^{d}$-valued square-integrable random variables with mean zero. Then we have the following estimation for the Prohorov distance of the corresponding Gaussian measures

$$
\pi(N(0, \operatorname{cov} X), N(0, \operatorname{cov} Y)) \leqq\left(E\left[\|X-Y\|^{2}\right]\right)^{1 / 3}
$$

Therefore

$$
\begin{aligned}
& \pi\left(N\left(0, \tilde{\tau}_{N}^{2}\right), N\left(0, \tau_{N}^{2}\right)\right) \leqq\left(E\left[\left|(l / n)^{1 / 2} Y_{1}-N^{-1 / 2} Y_{1}\right|^{2}\right]\right)^{1 / 3} \\
& =\left(n^{-1}\left(n-(N n l)^{1 / 2}\right)\right)^{2 / 3}\left(E\left[\left|N^{-1 / 2} \sum_{v=1}^{N} x_{v}\right|^{2}\right]\right)^{1 / 3}
\end{aligned}
$$

By Proposition 3.5 and the definition of $N$ and $l$ we conclude

$$
\begin{equation*}
\pi\left(N\left(0, \tilde{\tau}_{N}^{2}\right), N\left(0, \tau_{N}^{2}\right)\right) \ll\left(n^{-1}(l+N)\right)^{2 / 3} \ll n^{-1 / 8} . \tag{4.5}
\end{equation*}
$$

Theorem 2 follows from (4.1)-(4.5).

## 5. Proof of Theorem 3

We define $N_{k}=\left[k^{1+\kappa^{-1}}\right](k \geqq 1)$ where $\kappa$ is the exponent from Theorem 2 and set $t_{k}=\sum_{i \leq k-1} N_{i}$ as well as

$$
\begin{equation*}
X_{k}=N_{k}^{-1 / 2} \sum_{i=t_{k}+1}^{t_{k+1}} x_{i} \tag{5.1}
\end{equation*}
$$

Since $\left(x_{k}\right)_{k \in Z}$ is strictly stationary Theorem 2 implies

$$
\pi\left(\mathscr{L}\left(X_{k}\right), N\left(0, \sigma^{2}\right)\right) \leqq C^{\prime} N_{k}^{-\kappa}
$$

for some constant $C^{\prime}>0$ which can be assumed to be 1 for our purposes. Therefore $\mathscr{E}\left(X_{k}\right)$ and $N\left(0, \sigma^{2}\right)$ satisfy assumption (1.3) with $\rho_{k}=\delta_{k}=N_{k}^{-\kappa}$. Furthermore given a set $A \in \mathfrak{N}_{k-1}^{\prime}=\bigvee_{j=1}^{t_{k}} \mathfrak{F}_{j}, P(A)>0$ there is a Borel probability measure $\lambda$ on $\mathbb{R}^{N_{k}} \times \mathbb{R}^{N_{k}} \quad$ with marginals $\mathcal{E}\left(\left(x_{t_{k}+1}, \ldots, x_{t_{k+1}}\right)\right)$ and $\mathcal{L}\left(\left(x_{t_{k}+1}, \ldots, x_{t_{k+1}}\right) \mid A\right)$ such that

$$
\int N_{k}^{-1} \sum_{i=1}^{N_{k}}\left|u_{i}-v_{i}\right| d \lambda(u, v) \leqq \varepsilon\left(N_{k}\right) .
$$

If $T$ is defined on $\mathbb{R}^{N_{k}} \times \mathbb{R}^{N_{k}}$ by

$$
T\left(\left(u_{1}, \ldots, u_{N_{k}}\right),\left(v_{1}, \ldots, v_{N_{k}}\right)\right)=\left(N_{k}^{-1 / 2} \sum_{i=1}^{N_{k}} u_{i}, N_{k}^{-1 / 2} \sum_{i=1}^{N_{k}} v_{i}\right)
$$

then $T(\lambda)$ has marginals $\mathscr{Q}\left(X_{k}\right)$ and $\mathscr{L}\left(X_{k} \mid A\right)$ and

$$
\begin{aligned}
& \int_{\mathbb{R} \times \mathbb{R}}|u-v| d T(\lambda)(u, v) \\
& \leqq N_{k}^{-1 / 2} \int_{\mathbb{R}^{N_{k}} \times \mathbb{R}^{N_{k}}} \sum_{i=1}^{N_{k}}\left|u_{i}-v_{i}\right| d \lambda(u, v) \leqq C_{0} N_{k}^{-1 / 2}
\end{aligned}
$$

Again we can assume that $C_{0}=1 / 2$. This means that we can apply Theorem 1 with $\varepsilon_{k}=N_{k}^{-1 / 4}$ and $\eta_{k}=0$. Without loss of generality there exists a sequence $\left(Y_{k}\right)_{k \geqq 1}$ of independent $N\left(0, \sigma^{2}\right)$-distributed random variables satisfying

$$
P\left[\left|X_{k}-Y_{k}\right| \geqq 2 N_{k}^{-\kappa}+N_{k}^{-1 / 4}\right] \leqq N_{k}^{-1 / 4}+N_{k}^{-\kappa} \ll N_{k}^{-\kappa} \leqq k^{-(1+\kappa)}
$$

which implies by the Borel-Cantelli Lemma

$$
\begin{equation*}
\left|X_{k}-Y_{k}\right| \leqq 2 N_{k}^{-\kappa}+N_{k}^{-1 / 4} \ll N_{k}^{-\kappa} \quad \text { a.s. } \tag{5.2}
\end{equation*}
$$

For any Brownian motion $(X(t))_{t \geqq 0}$ with variance $\sigma^{2}$ the random variables $\left(t_{k+1}-t_{k}\right)^{-1 / 2}\left(X\left(t_{k+1}\right)-X\left(t_{k}\right)\right)$ have the same distribution as $Y_{k}$. Therefore without loss of generality we can assume that there exists a Brownian motion $(X(t))_{t \geq 0}$ with variance $\sigma^{2}$ satisfying

$$
\left(t_{k+1}-t_{k}\right)^{-1 / 2}\left(X\left(t_{k+1}\right)-X\left(t_{k}\right)\right)=Y_{k} \quad(k \geqq 1) .
$$

Now

$$
\left|\sum_{i=t_{k}+1}^{t_{k+1}} x_{i}-\left(X\left(t_{k+1}\right)-X\left(t_{k}\right)\right)\right|=N_{k}^{1 / 2}\left|X_{k}-Y_{k}\right| \ll N_{k}^{1 / 2-\kappa} \quad \text { a.s. }
$$

If we sum over $k$ and make use of the relation

$$
k^{2+\kappa^{-1}} \ll t_{k+1} \ll k^{2+\kappa^{-1}}
$$

we get

$$
\begin{array}{rl}
\mid \sum_{v=1}^{t_{k+1}} x_{v} & X\left(t_{k+1}\right) \mid \\
& \ll \sum_{j=1}^{k} N_{j}^{1 / 2-\kappa} \leqq k^{\left(1+\kappa^{-1}\right)(1 / 2-\kappa)+1}=k^{1 / 2-\kappa+1 / 2 \kappa} \ll t_{k+1}^{1 / 2-\kappa / 2} \quad \text { a.s. }
\end{array}
$$

For any $t>0$ choose $k$ such that $t_{k}<t \leqq t_{k+1}$ then

$$
\begin{aligned}
& \left|\sum_{v \leqq t} x_{v}-X(t)\right| \\
& \leqq\left|\sum_{v=1}^{t_{k}} x_{v}-X\left(t_{k}\right)\right|+\left|\sum_{v=t_{k}+1}^{t} x_{v}\right|+\left|X(t)-X\left(t_{k}\right)\right|
\end{aligned}
$$

Therefore our proof is finished if we show that there is a constant $\gamma>0$ such that with probability 1

$$
\begin{equation*}
\max _{t_{k}<t \leqq t_{k+1}}\left|\sum_{v=t_{k}+1}^{t} x_{v}\right| \ll t_{k}^{1 / 2-\gamma} \tag{5.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{t_{k}<t \leqq t_{k+1}}\left|X(t)-X\left(t_{k}\right)\right| \ll t_{k}^{1 / 2-\gamma} \tag{5.4}
\end{equation*}
$$

With the help of Lemma 3.8 the first of these two statements is shown analogous to the proof of Proposition 2.2 in [14]. (5.4) finally follows directly from Fernique's theorem [11].

## 6. Proof of Theorem 4

Let $X$ be a very weak Bernoulli process. For a fixed integer $N$ and any $k \in Z$ we define $X_{k}=\left(x_{(k-1) N+1}, \ldots, x_{k N}\right)$ and $\mathfrak{F}_{k}^{\prime}=\bigvee_{v=(k-1) N+1}^{k N} \mathfrak{F}_{v}$. Then $X_{k}$ is an $S^{N}$ valued $\mathfrak{F}_{k}^{\prime}$-measurable random variable. Define $\sigma_{N}$ on $S^{N} \times S^{N}$ as before. Given an $\varepsilon>0$ we choose $N$ according to Definition 1 such that for any $k \geqq 1$ there exists a set $D_{k} \in \mathfrak{A}_{k}^{\prime}=\bigvee_{j=1}^{k} \mathfrak{F}_{j}^{\prime}$ such that $P\left(D_{k}^{\subset}\right)<\varepsilon$ and if $A \in \mathfrak{M}_{k-1}^{\prime}, A \subset D_{k-1}$,
$P(A)>0$ then

$$
\inf _{\lambda \in \mathfrak{P}_{N}} \int_{N} \sigma_{N} d \lambda<\varepsilon
$$

where $\mathfrak{P}_{n}$ is the class of all Borel probability measures on $S^{N} \times S^{N}$ with marginals $\mathscr{E}\left(X_{k}\right)$ and $\mathscr{U}\left(X_{k} \mid A\right)$.

By Theorem 1 there exists a sequence $\left(Y_{k}\right)_{k \geqq 1}$ of independent random variables such that $\mathscr{Q}\left(Y_{k}\right)=\mathscr{L}\left(X_{k}\right)$ and for all $k \geqq 1$

$$
\begin{equation*}
P\left[\sigma_{N}\left(X_{k}, Y_{k}\right) \geqq(2 \varepsilon)^{1 / 2}\right] \leqq(2 \varepsilon)^{1 / 2}+\varepsilon . \tag{6.1}
\end{equation*}
$$

Let $Y$ be the process generated by the $Y_{k}$ and let $C$ denote the bound for $\sigma$ then by (1.8) for any $n \geqq 1$

$$
\bar{\rho}_{n N}(X, Y) \leqq n^{-1} \sum_{k=1}^{n} \int_{S^{N} \times S^{N}} \sigma_{N} d \mathscr{L}\left(X_{k}, Y_{k}\right) .
$$

Splitting this integral in the part where $\sigma_{N} \geqq(2 \varepsilon)^{1 / 2}$ and its complement $\left\{\sigma_{N}<(2 \varepsilon)^{1 / 2}\right\}$ we get by (6.1)

$$
\bar{\rho}_{n N}(X, Y) \leqq(C+1)(2 \varepsilon)^{1 / 2}+C \varepsilon .
$$

From this the result follows.

Acknowledgment. I am grateful to E. Berger for valuable comments.

## References

1. Araujo, A., Giné, E.: The central limit theorem for real and Banach valued random variables. New York-Chichester-Brisbane-Toronto: John Wiley \& Sons 1980
2. Berkes, I., Philipp, W.: Approximation theorems for independent and weakly dependent random vectors. Ann. Probability 7, 29-54 (1979)
3. Bradley, R.C.: A sufficient condition for linear growth of variances in a stationary random sequence. Proc. Amer. Math. Soc. 83, 586-589 (1981)
4. Dehling, H.: Grenzwertsätze für Summen schwach abhängiger vektor-wertiger Zufallsvariablen. Dissertation. Göttingen 1981
5. Dehling, H.: A note on a theorem of Berkes and Philipp. Z. Wahrscheinlichkeitstheorie verw. Gebiete 62, 39-42 (1983)
6. Dehling, H., Denker, M., Philipp, W.: Versik processes and very weak Bernoulli processes with summable rates are independent. Preprint 1982 (4 pages)
7. Dehling, H., Philipp, W.: Almost sure invariance principles for weakly dependent vector-valued random variables. Ann. Probability 10, 689-701 (1982)
8. Doob, J.L.: Stochastic processes. New York-London-Sydney: John Wiley \& Sons 1953
9. Dudley, R.M.: Distances of probability measures and random variables. Ann. Math. Statistics 39, 1563-1572 (1968)
10. Dudley, R.M.: Probabilities and metrics. Lecture Notes Series No. 45. Aarhus Universitet 1976
11. Fernique, X.: Integrabilité des vecteurs Gaussiens. C. R. Acad. Sci. Paris Ser. A. 270, 1698-1699 (1970)
12. Gray, R.M., Neuhoff, D.L., Shields, P.C.: A generalization of Ornstein's $\bar{d}$ distance with applications to information theory. Ann. Probability 3, 315-328 (1975)
13. Ibragimov, I.A., Linnik, Yu.V.: Independent and stationary sequences of random variables. Groningen: Wolters-Noordhoff Publishing 1971
14. Kuelbs, J., Philipp, W.: Almost sure invariance principles for partial sums of mixing $B$-valued random variables. Ann. Probability 8, 1003-1036 (1980)
15. Ornstein, D.: Ergodic Theory, Randomness, and Dynamical Systems. New Haven and London: Yale University Press 1974
16. Philipp, W.: Almost sure invariance principles for sums of $B$-valued random variables. In: Probability in Banach spaces IL, Lecture Notes in Math. 709, 171-193. Berlin-Heidelberg-New York: Springer 1979
17. Philipp, W.: Weak and $L^{p}$-invariance principles for sums of $B$-valued random variables. Ann. Probability 8, 68-82 (1980)
18. Serfling, R.J.: Contributions to central limit theory for dependent variables. Ann. Math. Statistics 39, 1158-1175 (1968)
19. Shields, P.C.: Almost block independence. Z. Wahrscheinlichkeitstheorie verw. Gebiete 49, 119-123 (1979)
20. Strassen, V.: An almost sure invariance principle for the law of the iterated logarithm. Z . Wahrscheinlichkeitstheorie verw. Gebiete 3, 211-226 (1964)
21. Strassen, V.: The existence of probability measures with given marginals. Ann. Math. Statistics 36, 423-439 (1965)
22. Yurinskii, V.V.: On the error of the Gaussian approximation for convolutions. Theory Prob. Applications 22, 236-247 (1977)

Received September 30, 1981

