# Vidyadhar S. Mandrekar WEAK CONVERGENCE OF STOCHASTIC PROCESSES 

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Weak Convergence of Stochastic Processes
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# Vidyadhar S. Mandrekar Weak Convergence of Stochastic Processes 

With Applications to Statistical Limit Theorems

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## Contents

1 Weak convergence of stochastic processes - 1 ..... Introduction - 1
2 Weak convergence in metric spaces - 4
2.1 Cylindrical measures - 4
2.2 Kolmogorov consistency theorem ..... 6
2.3 The finite dimensional family for Brownian motion - 9
$2.4 \quad$ Properties of Brownian motion - 11
2.5 Kolmogorov continuity theorem - 14
2.6 Exit time for Brownian motion and Skorokhod theorem ..... 16
2.6.1 Skorokhod theorem - 19
2.7 Embedding of sums of i.i.d. random variable in Brownian motion ..... 21
2.8 Donsker's theorem - ..... 23
2.9 Empirical distribution function ..... 25
2.1 Weak convergence of probability measure on Polish space ..... 29
2.10 .1 Prokhorov theorem - ..... 34
2.11 Tightness and compactness in weak convergence ..... 36
3 Weak convergence on $C[0,1]$ and $D[0, \infty)$ - 41
3.1 Structure of compact sets in $C[0,1]-41$3.1.1
41
3.2 Invariance principle of sums of i.i.d. random variables ..... 45
3.3 Invariance principle for sums of stationary sequences - 48
3.4 Weak convergence on the Skorokhod space - 50
3.4.13.4.2Skorokhod topology -52
3.5 Metric of $D[0,1]$ to make it complete - 53
3.6 Separability of the Skorokhod space - 56
3.7 Tightness in the Skorokhod space - 58
$3.8 \quad$ The space $D[0, \infty)$ ..... 60
3.8.1 Separability and completeness ..... 64
3.8.2 Compactness - ..... 64
3.8.3 Tightness ..... 66
3.8.4 Aldous's tightness criterion ..... 67
4 Central limit theorem for semi-martingalesand applications - 70
4.1 Local characteristics of semi-martingale - 70
4.2
4.3 Central limit theorem for semi-martingale - 76
4.4 Application to survival analysis ..... 81
4.5 Asymptotic distribution of $\hat{\beta}(t)$ and Kaplan-Meier estimate - $\mathbf{8}$
$5 \quad$ Central limit theorems for dependent random variables - 89
6 Empirical process - 110
6.1 Spaces of bounded functions - 114
6.2 Maximal inequalities and covering numbers ..... 119
6.3 Sub-Gaussian inequalities - 125
6.4 Symmetrization ..... 126
6.4.1 Glivenko-Cantelli theorems ..... 129
6.4.2 Donsker theorems ..... 131
6.5 Lindberg-type theorem and its applications ..... 133
Bibliography ..... 142

## 1 Weak convergence of stochastic processes

## Introduction

The study of limit theorems in probability has been important in inference in statistics. The classical limit theorems involved methods of characteristic functions, in other words, Fourier transform. One can see this in the recent book of Durrett [8]. To study the weak convergence of stochastic processes, one may be tempted to create concepts of Fourier's theory in infinite dimensional cases. However, such an attempt fails, as can be seen from the work of Sazanov and Gross [14]. The problem was also solved by Donsker [5] by choosing techniques of weak convergence of measures on the space of continuous functions by interpolating the partial sum in the central limit theorem. If we denote for $X_{1}, \ldots, X_{n}$ i.i.d. (independent identically distributed random variable with zero expectation and finite variance)

$$
S_{n, m}=X_{1}+\ldots+X_{m / \sqrt{n}} \text {, where } m=1,2, \ldots, n,
$$

and $S_{n, 0}=0$. One can define

$$
S_{n, u}=\left\{\begin{array}{l}
S_{n, m} \text { if } u=0,1, \ldots n \\
\text { linear if } u \epsilon[m-1, m]
\end{array}\right.
$$

Then, the continuous process $S_{[n, t]}, t \in[0,1]$ generates a sequence of measures on $C[0,1]$ if we can prove, with appropriate definition $P_{n}$ converges weakly to the measure $P$ given by Brownian motion on $C[0,1]$, we can conclude that

$$
\max _{0 \leq t \leq 1}\left|S_{n, t}\right| \text { converges to } \max _{0 \leq t \leq 1}|W(t)|
$$

and

$$
\frac{R_{n}}{\sqrt{n}}=\frac{1}{\sqrt{n}}+\max _{1 \leq m \leq n} S_{n, m}-\min _{0 \leq t \leq 1} S_{n, m}
$$

converges to

$$
\max _{1 \leq t \leq 1} W(t)-\min _{0 \leq t \leq 1} W(t)
$$

giving more information than the classical central theorem. One can also prove empirical distribution function, if properly interpolated,

$$
\hat{F}_{n}(x)=\frac{1}{n} \sum_{1}^{n} 1\left(x_{i} \leq x\right), x \in \mathbb{R},
$$

satisfies

$$
\sqrt{n} \sup _{x}\left|\hat{F}_{n}(x)-F(x)\right|,
$$

which converges to the maximum of the Brownian bridge $\max _{0 \leq t \leq 1}|W(t)-W(1)|$, justifying the Kolmogorov-Smirnov statistics for F continuous distribution function of $X_{1}$. We prove these results using convergence in probability in $C[0,1]$ space of $S[n$. to $W[$.] in the supremum distance in $C[0,1]$. For this, we used the embedding theorem of Skorokhod [23]. We follow [8] for the proof of this theorem of Skorokhod.

As we are interested in the convergence of semi-martingales to prove the statistical limit theorems for censored data that arise in clinical trials, we introduce the following Billingsley [2] convergence in a separable metric (henceforth, Polish space). Here we show the compactness of the sequences of probability measures on Polish space using the so-called tightness condition of Prokhorov [22], which connects compactness of measures with a compact set having large measures for the whole sequence of probability measures. We then consider the form of compact sets in $C[0,1]$ using the Arzela-Ascoli theorem. We also consider the question if one can define on the space of functions with jump discontinuity $D[0, T]$ a metric to make it a Polish space. Here we follow Billingsley [2] to present the so-called Skorokhod topology on $D[0, T$ ] and $D[0, \infty)$. We then characterize the compact sets in these space to study the tightness of a sequence of probability measures. Then we use a remarkable result of Aldous [1] to consider compactness in terms of stopping times. This result is then exploited to study the weak convergence of semi-martingales in $D[0, \infty)$.

Following Durrett and Resnick [9], we then generalize the Skorokhod embedding theorem for sums of dependent random variables. This allows us to extend weak convergence result, as in Chapter 2, to martingale differences [3]. Using the work of Gordin [13], one can reduce a similar theorem for stationary sequences to that for martingale differences. If one observes that empirical measures take values in $D[0,1]$ and we try to use the maximum norm, we get a nonseparable space. Thus, we can ask the question, "can one study weak convergence of stochastic processes taking values in non-separable metric spaces? This creates a problem with measurability as one can see from Dudley-Phillip's [7] work presented in [18]. To handle these problems, we study the so-called empirical processes following Van der Vaart and Wellner [25]. We introduce covering numbers, symmetrization, and sub-Gaussian inequalities to find conditions for weak convergence of measures.

We begin chapter 2 by constructing a measure on $\mathbb{R}^{T}$ to get a stochastic process. We then construct a Gaussian process with given covariance. We obtain sufficient conditions on the moments of a process indexed by $T=\mathbb{R}$ to have a continuous version. Using this, we construct Brownian motion. Then we prove the Skorokhod embedding theorem for sums of independent random variables. It is then exploited to obtain a convergence in $C[0,1]$ of the functions of the sum of independent random variables as described earlier. As a consequence of the central limit theorem in $C[0,1]$,
we prove using [20] the results of [11] and the convergence of symmetric statistics. The chapter ends by giving a weak convergence of probability measures and Prokhorov's tightness result on a Polish space. In chapter 3, we study compact sets in $C[0,1]$ and use them to get alternate proof of theorems similar to the one in chapter 2 using weak convergence in Polish spaces. This is followed by studying the topology of Skorokhod on $D[0, T]$ and $D[0, \infty)$ and proving that these are Polish spaces. Again, we study the compact sets and prove the result of Aldous. Chapter 4 studies the weak convergence of semi-martingales, which requires Lenglart [16] inequality to prove compactness using the result of Aldous. As a consequence of this, we derive weak convergence of Nelson and Kaplan-Meier estimates by simplifying the proof of Gill [12]. We do not present convergence of the Susarla-Van Ryzin Bayes estimate [24], but it can be obtained by similar methods as shown in [4]. For convergence of Linden-Bell estimates arising in astronomy, see [21] where similar techniques are used. Chapter 5 considers limit theorems as in chapter 2 using generalization of Skorokhod theorem from [9]. Limit theorems in chapters 2 and 5 use techniques given in Durrett's book [8]. Chapter 3 follows the presentation in the book of Billingsley [2], and chapter 4 uses the simplified version of techniques in [17] (see also [15]). We present in the last chapter the convergence of empirical processes using the techniques mentioned above taken from [25].

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## 2 Weak convergence in metric spaces

We begin in this chapter the process of associating a probability measure on the function space $\mathbb{R}^{T}$ for any set $T$ given a family of probability measures on $\mathbb{R}^{S}, S \subseteq T$ finite set with certain conditions. This allows us to define a family of real random variables $\left\{X_{t}, t \epsilon T\right\}$ on a probability space $(\Omega, \mathcal{F}, P)$. Such a family is referred to as a stochastic process. This idea originated from the work of Wiener (see [19]) in the special case of Brownian motion. Kolmogorov generalized it for the construction of any stochastic process and gave conditions under which one can find a continuous version (cf. section 2.5), that is, a stochastic process with continuous paths. We use his approach to construct Brownian motion. We explain it in the next section.

### 2.1 Cylindrical measures

Let $\left\{X_{t}, t \in T\right\}$ be family of random variable on a probability space ( $\Omega, \mathcal{F}, P$ ). Assume that $X_{t}$ takes values in $\left(\mathcal{X}_{t}, \mathcal{A}_{t}\right)$. For any finite set $S \subset T$,

$$
\mathcal{X}_{S}=\prod_{t \in S} \mathcal{X}_{t}, \mathcal{A}_{S}=\bigotimes_{t \in S} \mathcal{A}_{t}, Q_{S}=P \circ\left(X_{t}, t \in S\right)^{-1},
$$

where $Q_{S}$ is the induced measure. Check, if $\Pi_{S^{\prime} S}: \mathcal{X}_{S^{\prime}} \rightarrow \mathcal{X}_{S}$ for $S \subset S^{\prime}$, then

$$
\begin{equation*}
Q_{S}=Q_{S^{\prime}} \circ \Pi_{S^{\prime} S}^{-1} \tag{2.1}
\end{equation*}
$$

Suppose we are given a family $\left\{Q_{S}, S \subseteq T\right.$ finite-dimensional $\}$, a probability measure where $Q_{S}$ on $\left(\mathcal{X}_{S}, \mathcal{A}_{S}\right)$. Assume they satisfy (2.1). Then, there exists $Q$ on $\left(\mathcal{X}_{T}, \mathcal{A}_{T}\right)$ such that

$$
Q \circ \Pi_{S}^{-1}=Q_{S},
$$

where $\mathcal{X}_{T}=\prod_{t \epsilon T} \mathcal{X}_{t}, \mathcal{A}_{T}=\sigma\left(\bigcup_{S \subset T} \mathcal{C}_{S}\right), \mathcal{C}_{S}=\Pi_{S}^{-1}\left(\mathcal{A}_{S}\right)$.

## Remark: For

$$
\begin{aligned}
S & \subset T \text {, finite, } C \in \mathcal{C}_{S} \\
\text { define } Q_{0}(C) & =Q_{S}(A) \text {, where } C=\Pi_{S}^{-1}(A) \\
\text { as } \mathcal{C}_{S} & =\Pi_{S}^{-1}\left(\mathcal{A}_{S}\right)
\end{aligned}
$$

We can define $Q_{0}$ on $\bigcup_{S \subset T} \mathcal{C}_{S}$. Then, for $C \in \mathcal{C}_{S}$ and $\mathcal{C}_{S^{\prime}}$, then by (2.1),

$$
Q_{0}(C)=Q_{S}(A)=Q_{S^{\prime}}(A),
$$

and hence $Q_{0}$ is well-defined.

Note that as

$$
\mathcal{C}_{S_{1}} \cup \mathcal{C}_{S_{2}} \cup \cdots \cup \mathcal{C}_{S_{k}} \subset \mathcal{C}_{S_{1} \cup \cdots \cup S_{k}}
$$

$Q_{0}$ is finitely additive on $\bigcup_{S \subset T} \mathcal{C}_{S}$. We have to show the countable additivity.

Definition 2.1: A collection of subsets $\mathcal{K} \subset \mathcal{X}$ is called a compact class if for every sequence $\left\{C_{k}, k=1,2, \cdots n\right\}, n$ finite,

$$
\bigcap_{k=1}^{n} C_{k} \neq \emptyset \Longrightarrow \bigcap_{k=1}^{\infty} C_{k} \neq \emptyset
$$

Exercise 1: Every subcollection of compact class is compact.

Exercise 2: If $\mathcal{X}$ and $\mathcal{Y}$ are two spaces and $T: \mathcal{X} \rightarrow \mathcal{Y}$ and $\mathcal{K}$ is a compact class in $\mathcal{Y}$, then $T^{-1}(\mathcal{K})$ is a compact class in $\mathcal{X}$.

Definition 2.2: A finitely additive measure $\mu$ on $\left(\mathcal{X}, \mathcal{A}_{0}\right)$ is called compact if there exists a compact class $\mathcal{K}$ such that for every $A \in \mathcal{A}_{0}$ and $\epsilon>0$, there exists $C_{\epsilon} \in \mathcal{K}$, and $A_{\epsilon} \in \mathcal{A}_{0}$ such that

$$
A_{\epsilon} \subset C_{\epsilon} \subset A \quad \text { and } \quad \mu\left(A-A_{\epsilon}\right)<\epsilon
$$

We call $\mathcal{K}$ is $\mu$-approximates $\mathcal{A}_{0}$.
Lemma 2.1: Every compact finitely additive measure is countably additive.

Proof: Suppose $\left(\mathcal{X}, \mathcal{A}_{0}, \mu\right)$ is given. There exists a compact class $\mathcal{K}$ that is $\mu$-approximates $\mathcal{A}_{0}$. Let $\left\{A_{n}\right\} \subset \mathcal{A}_{0}$ such that $A_{n} \searrow \emptyset$. We need to show that $\mu\left(A_{n}\right) \searrow 0$. For given $\epsilon>0$, let $B_{n} \in \mathcal{A}_{0}$ and $C_{n} \in \mathcal{K}$, such that

$$
B_{n} \subset C_{n} \subset A_{n} \quad \text { and } \quad \mu\left(A_{n}-B_{n}\right)<\frac{\epsilon}{2^{n}} .
$$

Suppose $\mu\left(A_{n}\right)$ does not go to 0 , i.e., for all $n, \mu\left(A_{n}\right)>\epsilon$. Since we know that

$$
\mu\left(A_{n}-\bigcap_{k=1}^{n} B_{k}\right)=\mu\left(\bigcap_{k=1}^{n} A_{k}\right)-\mu\left(\bigcap_{k=1}^{n} B_{k}\right)<\frac{\epsilon}{2},
$$

we conclude that for all $n$,

$$
\mu\left(\bigcap_{k=1}^{n} B_{k}\right)>\frac{\epsilon}{2} .
$$

Next, for all $n$

$$
\bigcap_{k=1}^{n} B_{k} \neq \emptyset
$$

and hence, we have for all $n$

$$
\bigcap_{k=1}^{n} C_{k} \neq \emptyset
$$

which implies

$$
\bigcap_{k=1}^{\infty} C_{k} \neq \emptyset
$$

since $C_{k} \in \mathcal{K}$ and $\mathcal{K}$ is compact. Therefore, it follows that

$$
\bigcap_{k=1}^{\infty} A_{k} \supset \bigcap_{k=1}^{\infty} C_{k} \neq \emptyset
$$

implies

$$
\lim _{n \rightarrow \infty} A_{n} \neq \emptyset,
$$

which is a contradiction.

### 2.2 Kolmogorov consistency theorem

Suppose $S \subset T$ is finite subset and $Q_{S}$ is measure on $\left(\mathcal{X}_{S}, \mathcal{A}_{S}\right)$ satisfying consistency condition (2.1). Let $\left(\mathcal{X}_{\{t\}}, \mathcal{A}_{\{t\}}, Q_{\{t\}}\right)$ be a compact probability measure space. For each $t \in T$, there exists a compact class $\mathcal{K}_{t} \subseteq \mathcal{A}_{\{t\}}$ and $\mathcal{K}_{t}, Q_{t}$ approximates $\mathcal{A}_{\{t\}}$. Then, there exists a unique probability measure $Q_{0}$ on $\left(\mathcal{X}_{T}, \mathcal{A}_{T}\right)$ such that

$$
\Pi_{S}: \mathcal{X}_{T} \rightarrow \mathcal{X}_{S}, \text { and } Q_{0} \circ \Pi_{S}^{-1}=Q_{S}
$$

Proof: Define

$$
\mathcal{D}=\left\{C: C=\Pi_{t}^{-1}(K), K \in \mathcal{K}_{t}, t \in T\right\}
$$

Let

$$
\left\{\Pi_{t_{i}}^{-1}\left(C_{t_{i}}\right), i=1,2 \ldots\right\}
$$

be a countable family of sets and

$$
B_{t}=\bigcup_{t_{i}=t} \Pi_{t_{i}}^{-1}\left(C_{t_{i}}\right)
$$

If the countable intersection of $\left\{B_{t}, t \in T\right\}$ is empty, then $B_{t_{0}}$ is empty for some $t_{0}$. Since $\mathcal{K}_{t_{0}}$ is a compact class and all $C_{t_{i}} \in \mathcal{K}_{t_{0}}$, we get a finite set of $t_{\mathrm{i}}$ 's $\left(t_{i}=t_{0}\right)$. Let us call it $J$ for which

$$
\bigcup_{t_{i} \in J} C_{t_{i}}=\emptyset \Longrightarrow \bigcup_{t_{i} \in J} \Pi_{t_{i}}^{-1}\left(C_{t_{i}}\right)=\emptyset
$$

Since $\mathcal{D}$ is a compact class, $\mathcal{K}$ as a countable intersections of sets in $\mathcal{D}$ is a compact finitely additive class. We shall show that $Q_{0}$ is a compact measure, i.e., $\mathcal{K}$ $Q_{0}$-approximates $\mathcal{C}_{0}$. Take $C \in \mathcal{C}_{0}$ and $\epsilon>0$. For some $S \subset T$,

$$
C=\Pi_{S}^{-1}(B) .
$$

Choose a rectangle

$$
\prod_{t \in S}\left(A_{t}\right) \subset B
$$

so that for $A_{t} \in \mathcal{A}_{t}$

$$
\begin{aligned}
& Q_{S}\left(B-\prod_{t \in S} A_{t}\right)<\frac{\epsilon}{2} \\
& Q_{0}\left(\Pi_{S}^{-1}(B)-\Pi_{S}^{-1}\left(\prod_{t \in S} A_{t}\right)\right)<\frac{\epsilon}{2} .
\end{aligned}
$$

For each $t$, choose $K_{t} \in \mathcal{K}_{t}$ such that $K_{t} \subset A_{t}$ and

$$
Q_{t}\left(A_{t}\right)<Q_{t}\left(K_{t}\right)+\frac{\epsilon}{\text { cardinality }(S)} .
$$

Let

$$
K=\Pi_{t \in S} K_{t} \text { for } K_{t} \in \mathcal{K}_{t} .
$$

Then, $K \subset C$ and

$$
\begin{aligned}
Q_{0}\left(\Pi_{S}^{-1}(B)-\Pi_{S}^{-1}\left(\prod_{t \in S} K_{t}\right)\right)= & Q_{0}\left(\Pi_{S}^{-1}\left(B-\prod_{t \in S} A_{t}\right)\right) \\
& +Q_{0}\left(\Pi_{S}^{-1}\left(\prod_{t \in S} A_{t}\right)-\Pi_{S}^{-1}\left(\prod_{t \in S} K_{t}\right)\right) \\
< & \epsilon .
\end{aligned}
$$

$Q_{0}$ extends to a countable additive measure on $\sigma(\mathcal{C})$.

$$
\text { ( } \mathcal{X}_{t}=R, R^{d}, \text { or a complete separate metric space.) }
$$

Example: $T=N, \mathcal{X}_{t}=R$.
Suppose

$$
Q_{n}=\bigotimes_{t \in\{1,2, \ldots, n\}} Q_{\{t\}}
$$

Then, there exists $\left\{X_{n}, n \in N\right\}$ of random variables defined on $R^{\infty}$.

Example: $T=[0,1]$.
Let $\{C(t, s), t, s \in T\}$ be a set of real valued function with $C(t, s)=C(s, t)$ and

$$
\sum_{t, s \in S} a_{t} a_{s} C(t, s) \geq 0
$$

for $S$ finite. ( $\left\{a_{t}, t \in S\right\} \subset R$ ) Let $Q_{S}$ be a probability measure with characteristic function for $\mathbf{t} \in R^{d}$

$$
\begin{equation*}
\phi_{Q_{S}}(\mathbf{t})=\exp \left(-\frac{1}{2} \mathbf{t}^{\prime} \sum_{S} \mathbf{t}\right), \tag{2.2}
\end{equation*}
$$

where

$$
\sum_{S}=(C(u, v))_{u, v \in S}
$$

$Q_{S}$ satisfies the condition of Kolmogorov consistency theorem. Therefore, there exists $\left\{X_{t}, t \in T\right\}$, a family of random variables such that joint distribution of $\left\{X_{t}, t \in S\right\}$ is $Q_{S}$.

Example: Take $t, s \in[0,1]$ and $C(t, s)=\min (t, s)$.

$$
\begin{aligned}
C(t, s) & =\int_{0}^{1} 1_{[0, t]}(u) \cdot 1_{[0, s]}(u) d u \\
& =\min (t, s)
\end{aligned}
$$

Let $S=\left\{t_{1}, \ldots, t_{n}\right\},\left\{a_{1}, \ldots, a_{n}\right\} \subset R$. Then,

$$
\begin{aligned}
\sum_{i, j} a_{i} a_{j} C\left(t_{i}, t_{j}\right) & =\sum_{i, j} a_{i} a_{j} \int_{0}^{1} 1_{\left[0, t_{i}\right]}(u) \cdot 1_{\left[0, t_{j}\right]}(u) d u \\
& =\int_{0}^{1}\left(\sum_{i=1}^{n} a_{i} 1_{\left[0, t_{i}\right]}(u)\right)^{2} d u \\
& \geq 0
\end{aligned}
$$

since $\sum \sum a_{i} a_{j} C\left(t_{i}, t_{j}\right)$ is non-negative definite. Therefore, there exists a Gaussian process with covariance $C(t, s)=\min (t, s)$.

### 2.3 The finite dimensional family for Brownian motion

Given covariance function $C(t, s)=C(s, t)$ with $s, t \in T$, and for all $\left\{a_{1}, \ldots, a_{k}\right\}$ and $\left\{t_{1}, \ldots, t_{k}\right\} \subset T$, such that

$$
\sum_{k, j} a_{k} a_{j} C\left(t_{k}, t_{j}\right) \geq 0
$$

then there exists a Gaussian process $\left\{X_{t}, t \in T\right\}$ with $E X_{t}=0$ for all $t$ and $C(t, s)=$ $E\left(X_{t} X_{s}\right)$.

Example: $C(t, s)=\min (t, s)$, and $T=[0,1]$.

$$
\begin{aligned}
\min (t, s) & =\int_{0}^{1} 1_{[o, t]}(u) \cdot 1_{[o, s]}(u) d u \\
& =E X_{t} X_{s}
\end{aligned}
$$

There exists $\left\{X_{t}, t \in T\right\}$ such that $C(t, s)=E\left(X_{t} X_{s}\right)=\min (t, s)$.
Since $X_{t}$ is Gaussian, we know that

$$
X_{t} \in L^{2}(\Omega, \mathcal{F}, P)
$$

Let, with $\overline{\mathrm{SP}}^{L^{2}}$ denoting closure in $L_{2}$ of the subspace generated by

$$
M(X)=\overline{\mathrm{SP}}^{L^{2}}\left(X_{t}, t \in T\right),
$$

Consider the map

$$
I\left(X_{t}\right)=1_{[0, t]}(u) \in L_{2}([0,1]) \text { (Lebesgue measure). }
$$

$I$ is an isometry. Therefore, for $\left(t_{1}, \ldots, t_{n}\right)$ with $t_{1} \leq \ldots \leq t_{n}$

$$
\begin{aligned}
I\left(X_{t_{k}}-X_{t_{k-1}}\right) & =I\left(X_{t_{k}}\right)-I\left(X_{t_{k-1}}\right)(\text { because } I \text { is a linear map }) \\
& =1_{\left[0, t_{k}\right]}(u)-1_{\left[0, t_{k-1}\right]}(u) \\
& =1_{\left(t_{k-1}, t_{k}\right]}(u)
\end{aligned}
$$

For $k \neq j$,

$$
\begin{aligned}
E\left(X_{t_{k}}-X_{t_{k-1}}\right)\left(X_{t_{j}}-X_{t_{j-1}}\right) & =\int_{0}^{1} 1_{\left(t_{k-1}, t_{k}\right]}(u) \cdot 1_{\left(t_{j-1}, t_{j}\right]}(u) d u \\
& =0
\end{aligned}
$$

$X_{t_{k}}-X_{t_{k-1}}$ is independent of $X_{t_{j}}-X_{t_{j-1}}$ if $\left(t_{k-1}, t_{k}\right] \cap\left(t_{j-1}, t_{j}\right]=\emptyset$ because

$$
\left(t_{k-1}, t_{k}\right] \cap\left(t_{j-1}, t_{j}\right]=\emptyset \Longrightarrow E\left(X_{t_{k}}-X_{t_{k-1}}\right)\left(X_{t_{j}}-X_{t_{j-1}}\right)=0
$$

$\left\{X_{t}, t \in T\right\}$ is an independent increment process.

Given $t_{0}=0, X_{0}=0, t_{0} \leq t_{1} \leq \ldots \leq t_{n}$, we have

$$
P\left(X_{t_{k}}-X_{t_{k-1}} \leq x_{k}, k=1,2 \ldots, n\right)=\prod_{k=1}^{n} \int_{-\infty}^{\infty} \frac{1}{\sqrt{2 \pi\left(t_{k}-t_{k-1}\right)}} e^{-\frac{\nu_{k}^{2}}{2\left(t_{k}-t_{k-1}\right)}} d y_{k} .
$$

Using transformation

$$
\begin{aligned}
Y_{t_{1}} & =X_{t_{1}} \\
Y_{t_{2}} & =X_{t_{2}}-X_{t_{1}} \\
\vdots & =\vdots \\
Y_{t_{n}} & =X_{t_{n}}-X_{t_{n-1}}
\end{aligned}
$$

we can compute the joint density of $\left(X_{t_{1}}, \ldots, X_{t_{n}}\right)$.

$$
f_{\left(X_{t_{1}}, \ldots, X_{t_{n}}\right)}\left(x_{1}, \ldots, x_{n}\right)=\frac{1}{\prod_{k=1}^{n} \sqrt{2 \pi\left(t_{k}-t_{k-1}\right)}} \exp \left[-\frac{1}{2} \sum_{k=1}^{n} \frac{\left(x_{k}-x_{k-1}\right)^{2}}{\left(t_{k}-t_{k-1}\right)}\right] .
$$

Define

$$
p(t, x, B)=\frac{1}{\sqrt{2 \pi t}} \int_{B} e^{-\frac{(v-x)^{2}}{t}} d y, \quad t \geq 0
$$

and

$$
\tilde{p}(t, s, x, B)=p(t-s, x, B) .
$$

Exercise: Prove for $0=t_{0} \leq t_{1} \leq \ldots \leq t_{n}$

$$
\begin{aligned}
Q_{X_{t_{1}}, \ldots, X_{t_{n}}}\left(B_{1} \times \cdots \times B_{n}\right)= & \int_{B_{n}} \cdots \int_{B_{1}} p\left(t_{1}-t_{0}, 0, d y_{1}\right) \\
& p\left(t_{2}-t_{1}, y_{1}, d y_{2}\right) \cdots p\left(p\left(t_{n}-t_{n-1}, y_{n-1}, d y_{n}\right)\right) .
\end{aligned}
$$

Suppose we are given transition function $p(t, s, x, B)$ with $x \leq t$. Assume that $s \leq t \leq u$

$$
p(u, s, x, B)=\int p(t, s, x, d y) p(s, t, y, B) \text { (C-Kolmogorov condition). }
$$

Then, for $0=t_{0} \leq t_{1} \leq \ldots \leq t_{n}$

$$
\begin{aligned}
Q_{t_{1}, \ldots, t_{n}}^{x}\left(B_{1} \times \cdots \times B_{n}\right)= & \int_{B_{n}} \cdots \int_{B_{1}} p\left(t_{1}, t_{0}, x, d y_{1}\right) \\
& p\left(t_{2}, t_{1}, y_{1}, d y_{2}\right) \cdots p\left(p\left(t_{n}, t_{n-1}, y_{n-1}, d y_{n}\right)\right) .
\end{aligned}
$$

(Use Fubini's theorem): For this consistent family, there exists a stochastic process with $Q$ as finite dimensional distributions.

Exercise 2: Check that $\tilde{p}(t, s, x, B)$ satisfies the above condition.

$$
Q^{x}\left(X_{s} \in B_{1}, X_{t} \in B_{2}\right)=\int_{B_{1}} Q^{x}\left(X_{t} \in B_{2} \mid X_{s}=y\right) Q^{x} \circ X_{s}^{-1}(d y),
$$

where

$$
Q^{x}\left(X_{t_{n}} \in B_{n} \mid X_{t_{1}}, \ldots, X_{t_{n-1}}\right)=p\left(t_{n}, t_{n-1}, X_{t_{n-1}}, B_{n}\right) \text { (Markov). }
$$

The Gaussian process with covariance

$$
\min (t, s), \quad t, s, \in[0,1]
$$

has independent increments and is Markov.

Remarks: Consider

- $\quad X(\omega) \in R^{[0,1]}$
- $\quad C[0,1] \notin \sigma\left(C\left(R^{[0,1]}\right)\right)$
$C[0,1]$ is not measurable, but $C[0,1]$ has $Q_{0}^{*}(C[0,1])=1$ (outer measure).


### 2.4 Properties of Brownian motion

Definition 2.3: A Gaussian process is a stochastic process $X_{t}, t \in T$ for which any finite linear combination of $\left\{X_{t}\right\}$ has a Gaussian distribution.

Notation-wise, one can write $X \sim G P(m, K)$, meaning the random function $X$ is distributed as a GP with mean function $m$ and covariance function $K$.

## Remark:

$$
\begin{aligned}
X & \sim N\left(\mu_{X}, \sigma_{X}^{2}\right) \\
Y & \sim N\left(\mu_{Y}, \sigma_{Y}^{2}\right) \\
Z & =X+Y
\end{aligned}
$$

Then, $Z \sim N\left(\mu_{Z}, \sigma_{Z}^{2}\right)$, where

$$
\mu_{Z}=\mu_{X}+\mu_{Y} \text { and } \sigma_{Z}^{2}=\sigma_{X}^{2}+\sigma_{Y}^{2}+2 \rho_{X, Y} \sigma_{X} \sigma_{Y} .
$$

Proposition 2.1: Given covariance $C(t, s)=C(s, t)$ with $s, t \in T$ and for all $\left\{a_{1}, \ldots, a_{k}\right\} \subset R$ and $\left\{t_{1}, \ldots, t_{k}\right\} \subset T$

$$
\sum_{k, j} a_{k} a_{j} C\left(t_{k}, t_{j}\right) \geq 0
$$

then there exists a Gaussian process such that

$$
\left\{X_{t}, t \in T\right\} \text { with for all } t, E X_{t}=0, C(t, s)=E\left(X_{t}, X_{s}\right)
$$

Example: For $C(t, s)=\min (t, s)$ and $T=[0,1]$, we recall the properties of the associated Gaussian process.

In this example,

$$
\begin{aligned}
\min (t, s) & =\int_{0}^{1} 1_{[0, t]}(u) 1_{[0, s]}(u) d u \\
& =E X_{t} X_{s}
\end{aligned}
$$

Thus, there exists $\left\{X_{t}, t \in T\right\}$, which is a Gaussian process, such that for all $t$

$$
E X_{t}=0, \text { and } C(t, s)=E\left(X_{t}, X_{s}\right)=\min (t, s)
$$

Since $X_{t}$ is Gaussian, we know that

$$
X_{t} \in L^{2}(\Omega, \mathcal{F}, P)
$$

Let

$$
M(X)=\overline{S P}^{L^{2}}\left(X_{t}, t \in T\right) \text { (SP means "span"). }
$$

Consider the map $I: M(X) \rightarrow \overline{S P}\left\{1_{[0, t](u)}, t \in[0,1]\right\}$ such that

$$
I\left(X_{t}\right)=1_{[0, t]}(u)
$$

and

$$
I\left(\sum a_{k} X_{t_{k}}\right)=\sum a_{k} I\left(X_{t_{k}}\right)
$$

Proposition 2.2: $I$ is a linear map.
Exercise: Prove Proposition 2.2.
Proposition 2.3: $I$ is an isometry

Proof: Since $\left\{X_{t}, t \in T\right\}$ is Gaussian,

$$
\begin{aligned}
\operatorname{Var}\left(X_{t_{k}}-X_{t_{k-1}}\right) & =\operatorname{Var}\left(X_{t_{k}}\right)+\operatorname{Var}\left(X_{t_{k}}\right)-2 \operatorname{Cov}\left(X_{t_{k}}, X_{t_{k-1}}\right) \\
& =C\left(t_{k}, t_{k}\right)+C\left(t_{k-1}, t_{k-1}\right)-2 C\left(t_{k-1}, t_{k}\right) \\
& =t_{k}+t_{k-1}-2 t_{k-1} \\
& =t_{k}-t_{k-1} .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\left\|X_{t_{k}}-X_{t_{k-1}}\right\|_{L_{2}}^{2} & =\int_{[0,1]}\left(X_{t_{k}}-X_{t_{k-1}}\right)^{2} d P \\
& =E\left(X_{t_{k}}-X_{t_{k-1}}\right)^{2} \\
& =\operatorname{Var}\left(X_{t_{k}}-X_{t_{k-1}}\right) \\
& =t_{k}-t_{k-1} .
\end{aligned}
$$

Also,

$$
\begin{aligned}
\left\|I\left(X_{t_{k}}\right)-I\left(X_{t_{k-1}}\right)\right\|_{L_{2}}^{2} & =\left\|1_{\left[0, t_{k}\right]}(u)-1_{\left[0, t_{k-1}\right]}(u)\right\|_{L_{2}}^{2} \\
& =\left\|1_{\left(t_{k-1}, t_{k}\right]}(u)\right\|_{L_{2}}^{2} \\
& =t_{k}-t_{k-1} \\
& =\left\|X_{t_{k}}-X_{t_{k-1}}\right\|_{L_{2}}^{2} .
\end{aligned}
$$

This completes the proof.
Exercise: $I$ can be extended by continuity on $M(X)$ onto $L^{2}([0,1], \lambda)$
We call $I^{-1}(f)$ the stochastic integral of $f$ with respect to Brownian motion for $f \in L^{2}([0,1], \lambda)$.

Suppose that $t_{2} \leq \ldots \leq t_{k}$. Then,

$$
\begin{aligned}
I\left(X_{t_{k}}-X_{t_{k-1}}\right) & =I\left(X_{t_{k}}\right)-I\left(X_{t_{k-1}}\right) \text { (since } I \text { is an linear map) } \\
& =1_{\left[0, t_{k}\right]}(u)-1_{\left[0, t_{k-1}\right]}(u) \\
& =1_{\left(t_{k-1}, t_{k}\right]}(u) .
\end{aligned}
$$

$X_{t_{k}}-X_{t_{k-1}}$ is independent of $X_{t_{j}}-X_{t_{j-1}}$ if

$$
\left(t_{k-1}, t_{k}\right] \cap\left(t_{j-1}, t_{j}\right]=\emptyset .
$$

Proposition 2.4: If $X_{t_{k}}-X_{t_{k-1}}$ is independent of $X_{t_{j}}-X_{t_{j-1}}$, then

$$
E\left(X_{t_{k}}-X_{t_{k-1}}\right)\left(X_{t_{j}}-X_{t_{j-1}}\right)=0 .
$$

Proof: For $k \neq j$,

$$
\begin{aligned}
& E\left(X_{t_{k}}-X_{t_{k-1}}\right)\left(X_{t_{j}}-X_{t_{j-1}}\right) \\
= & \int_{0}^{1} 1_{\left(t_{k-1}, t_{k}\right]}(u) 1_{\left(t_{j-1}, t_{j]}\right]}(u) d u \\
= & 0 .
\end{aligned}
$$

Suppose $\left\{X_{t}, t \in T\right\}$ is an independent increment process such that $t_{0}=0, X_{0}=0$, and $t_{0} \leq t_{1} \leq \ldots \leq t_{n}$. Then, $X_{t_{k}}-X_{t_{k-1}}$ is Gaussian with mean 0 and variance $t_{k}-t_{k-1}$.

$$
\begin{aligned}
P\left(X_{t_{k}}-X_{t_{k-1}} \leq x_{k}, k=1,2 \ldots, n\right) & =\prod_{k=1}^{n} P\left(X_{t_{k}}-X_{t_{k-1}} \leq x_{k}\right) \\
& =\prod_{k=1}^{n} \int_{-\infty}^{x_{k}} \frac{1}{\sqrt{2 \pi\left(t_{k}-t_{k-1}\right)}} e^{-\frac{v_{k}^{2}}{2\left(t_{k}-t_{k-1}\right)}} d y_{k} .
\end{aligned}
$$

Let

$$
\begin{aligned}
Y_{t_{1}} & =X_{t_{1}} \\
Y_{t_{2}} & =X_{t_{2}}-X_{t_{1}} \\
\vdots & \vdots \\
Y_{t_{n}} & =X_{t_{n}}-X_{t_{n-1}} .
\end{aligned}
$$

Then,

$$
f_{X_{t_{1}}, \ldots, X_{t_{n}}}\left(x_{1}, \ldots, x_{n}\right)=\frac{1}{\prod_{k=1}^{n} \sqrt{2 \pi\left(t_{k}-t_{k-1}\right)}} \exp \left[-\frac{1}{2} \sum_{k=1}^{n} \frac{\left(x_{k}-x_{k-1}\right)^{2}}{\left(t_{k}-t_{k-1}\right)}\right] .
$$

### 2.5 Kolmogorov continuity theorem

For each $t$, if $P\left(\tilde{X}_{t}=X_{t}\right)=1$, then we say the finite dimensional distributions of $\tilde{X}_{t}$ and $X_{t}$ are the same. We call $\tilde{X}_{t}$ a version of $X_{t}$.
Proposition 2.5: Let $\left\{X_{t}, t \in[0,1]\right\}$ be a stochastic process with

$$
E\left|X_{t}-X_{s}\right|^{\beta} \leq C|t-s|^{1+\alpha} \text { with } C, \alpha, \beta>0 .
$$

Then, there exists a version of $\left\{X_{t}, t \in[0,1]\right\}$, which has a continuous sample paths.
Corollary 2.1: The Gaussian process with covariance function

$$
C(t, s)=\min (t, s), t, s \in[0,1]
$$

(has independent increment and is Markovian) has a version that is continuous.

$$
\begin{aligned}
E\left(X_{t}-X_{s}\right)^{2} & =E X_{t}^{2}-2 E X_{t} X_{s}+E X_{s}^{2} \\
& =t-2 s+s \\
& =|t-s| \\
E\left(X_{t}-X_{s}\right)^{4} & =3\left[E\left(X_{t}-X_{s}\right)^{2}\right]^{2} \\
& =3|t-s|^{2}
\end{aligned}
$$

We shall denote the continuous version by $W_{t}$.
Proof of Proposition 2.5: Take $0<\gamma<\frac{\alpha}{\beta}$ and $\delta>0$ such that

$$
(1-\delta)(1+\alpha-\beta \gamma)>1+\delta
$$

For $0 \leq i \leq j \leq 2^{n}$ and $|j-i| \leq 2^{n \delta}$,

$$
\begin{aligned}
\sum_{i, j} P\left(\left|X_{j 2^{-n}}-X_{i 2^{-n}}\right|>\left[(j-i) 2^{-n}\right]^{\gamma}\right) & \leq C \sum_{i, j}\left[(j-i) 2^{-n}\right]^{-\beta \gamma+(1+\delta)} \text { (by Chevyshev) } \\
& =C_{2} 2^{n[(1+\delta)-(1-\delta)(1+\alpha-\beta \gamma)]} \\
& <\infty
\end{aligned}
$$

where $(1+\delta)-(1-\delta)(1+\alpha-\beta \gamma)=-\mu$.
Then, by the Borell-Cantelli lemma,

$$
P\left(\left|X_{j 2^{-n}}-X_{i 2^{-n}}\right|>\left[(j-i) 2^{-n}\right]^{\gamma}\right)=0,
$$

i.e., there exists $n_{0}(\omega)$ such that for all $n \geq n_{0}(\omega)$

$$
\left|X_{j 2^{-n}}-X_{i 2^{-n}}\right| \leq\left[(j-i) 2^{-n}\right]^{\gamma} .
$$

Let $t_{1}<t_{2}$ be rational numbers in $[0,1]$ such that

$$
\begin{aligned}
t_{2}-t_{1} & \leq 2^{-n_{0}(1-\delta)} \\
t_{1} & =i 2^{-n}-2^{-p_{1}}-\ldots-2^{-p_{k}}\left(n<p_{1}<\ldots<p_{k}\right) \\
t_{2} & =j 2^{-n}-2^{-q_{1}}-\ldots-2^{-q_{k}}\left(n<q_{1}<\ldots<q_{k}\right) \\
t_{1} & \leq i 2^{-n} \leq j 2^{-n} \leq t_{2} .
\end{aligned}
$$

Let

$$
h(t)=t^{\gamma} \text { for } 2^{-(n+1)(1-\delta)} \leq t \leq 2^{-n(1-\delta)} .
$$

Then,

$$
\begin{aligned}
\left|X_{t_{1}}-X_{i 2^{-n}}\right| & \leq C_{1} h\left(2^{-n}\right) \\
\left|X_{t_{2}}-X_{j 2^{-n}}\right| & \leq C_{2} h\left(2^{-n}\right) \\
\left|X_{t_{2}}-X_{t_{1}}\right| & \leq C_{3} h\left(t_{2}-t_{1}\right)
\end{aligned}
$$

and

$$
\left|X_{i 2^{-n}-2^{-p_{1}} \ldots \ldots 2^{-p_{k}}}-X_{i 2^{-n}-2^{-p_{1}} \ldots-2^{-p_{k-1}}}\right| \leq 2^{-p_{k} \gamma} .
$$

Under this condition, the process $\left\{X_{t}, t \in[0,1]\right\}$ is uniformly continuous on rational numbers in $[0,1]$.

Let

$$
\psi:\left(\Omega_{Q}, \mathcal{C}_{\Omega_{Q}}\right) \rightarrow(C[0,1], \sigma(C[0,1])),
$$

which extends uniformly continuous functions in rational to continuous function in $[0,1]$.

Let $P$ is a measure generated by the same finite dimensional on rationals. Then

$$
\tilde{P}=P \circ \psi^{-1}
$$

is the measure of $\tilde{X}_{t}$, version of $X_{t}$.
For $\left\{X_{t}, t \in[0,1]\right\}$, there exists a version of continuous sample path. In case of Brownian motion, there exists a continuous version. We call it $\left\{W_{t}\right\}$.
$\left\{W_{t+t_{0}}-W_{t_{0}}, t \in[0, \infty]\right\}$ is a Weiner process.

### 2.6 Exit time for Brownian motion and Skorokhod theorem

It is well known that if $X_{1}, \ldots, X_{n}$ are independent identically distributed (iid) random variables with $E X_{1}=0$ and $E X_{1}^{2}<\infty$, then

$$
Y_{n}=\frac{X_{1}+\ldots+X_{n}}{\sqrt{n}}
$$

converges in distribution to standard normal random variable $Z$. This is equivalent to weak convergence of $Y_{n}$ to $Z$, i.e., $E f\left(Y_{n}\right) \rightarrow E f(Z)$ for all bounded continuous functions for $\mathbb{R}$. If we denote by

$$
Y_{n, m}=\frac{X_{1}+\ldots+X_{m}}{\sqrt{n}} \text { and } Y_{n, 0}=0
$$

and define interpolated continuous process

$$
Y_{n, t}=\left\{\begin{array}{l}
Y_{n, m} t=m \epsilon\{0,1, \ldots n\} \\
\text { linear if } t \in[m-1, m]
\end{array}\right.
$$

for $m=0,1,2, \ldots n$ and we shall show that $E f\left(Y_{n,[n .]} \rightarrow E f(W\right.$.$) for all bounded conti-$ nuous functions for $C[0,1]$ with $\{W(t), t \epsilon[0,1]\}$ Brownian motion. In fact, we shall use Skorokhod embedding to prove for any $\epsilon>0$

$$
P\left(\sup _{0 \leq t \leq 1}\left|Y_{n,[n t]}-W(t)\right|>\epsilon\right) \rightarrow 0
$$

as $n \rightarrow \infty$, which implies the above weak convergences (see Theorem 2.3).
Let

$$
\mathcal{F}_{t}^{W}=\sigma\left(W_{s}, s \leq t\right)
$$

$\tau$ is called "stopping time" if

$$
\{\tau \leq t\} \in \mathcal{F}_{t}^{W}
$$

Define

$$
\mathcal{F}_{\tau}=\left\{A: A \cap\{\tau \leq t\} \in \mathcal{F}_{t}^{W}\right\} .
$$

Then $\mathcal{F}_{\tau}$ is a $\sigma$-field. $\frac{W_{t}}{\sqrt{t}}$ is a standard normal variable.
Define

$$
T_{a}=\inf \left\{t: W_{t}=a\right\}
$$

Then, $T_{a}<\infty$ a.e. and is a stopping time.
Theorem 2.1: Let $a<x<b$. Then

$$
P_{x}\left(T_{a}<T_{b}\right)=\frac{b-x}{b-a} .
$$

Remark: $W_{t}$ is Gaussian and has independent increment. Also, for $s \leq t$

$$
E\left(W_{t}-W_{s} \mid \mathcal{F}_{s}\right)=0
$$

and hence $\left\{W_{t}\right\}$ is martingale.

Proof: Let $T=T_{a} \wedge T_{b}$. We know $T_{a}, T_{b}<\infty$ a.e., and

$$
W_{T_{a}}=a, \text { and } W_{T_{b}}=b
$$

Since $\left\{W_{t}\right\}$ is MG,

$$
\begin{aligned}
E_{x} W_{T} & =E_{\chi} W_{0} \\
& =x \\
& =a P\left(T_{a}<T_{b}\right)+b\left(1-P\left(T_{a}<T_{b}\right)\right) .
\end{aligned}
$$

Therefore,

$$
P_{x}\left(T_{a}<T_{b}\right)=\frac{b-x}{b-a} .
$$

$\left\{W_{t}, t \in[0, \infty]\right\}$ is a Weiner Process, starting from $x$ with $a<x<b$. We know that

$$
\begin{aligned}
E\left(\left(W_{t}-W_{s}\right)^{2} \mid \mathcal{F}_{s}^{W}\right) & =E\left(W_{t}-W_{s}\right)^{2} \\
& =(t-s)
\end{aligned}
$$

Also,

$$
\begin{aligned}
E\left(\left(W_{t}-W_{s}\right)^{2} \mid \mathcal{F}_{s}^{W}\right) & =E\left(W_{t}^{2} \mid \mathcal{F}_{s}^{W}\right)-2 E\left(W_{t} W_{s} \mid \mathcal{F}_{s}^{W}\right)+E\left(W_{s}^{2} \mid \mathcal{F}_{s}^{W}\right) \\
& =E\left(W_{t}^{2} \mid \mathcal{F}_{s}^{W}\right)-W_{s}^{2} \\
& =E\left(W_{t}^{2}-W_{s}^{2} \mid \mathcal{F}_{s}^{W}\right) \\
& =(t-s) .
\end{aligned}
$$

Therefore,

$$
E\left(W_{t}^{2}-t \mid \mathcal{F}_{s}^{W}\right)=W_{s}^{2}-s
$$

and hence $\left\{\left(W_{t}^{2}-t\right), t \in[0, \infty]\right\}$ is a martingale.
Suppose that $x=0$ and $a<0<b$. Then $T=T_{a} \wedge T_{b}$ is a finite stopping time. Therefore,

## $T \wedge t$

is also stopping time.

$$
\begin{aligned}
& \quad E\left(W_{T \wedge t}^{2}-T \wedge t\right)=0 \\
& E_{0}\left(W_{T}^{2}\right)=E_{0} T \\
& E W_{T}^{2}=E T \\
& =a^{2} P\left(T_{a}<T_{b}\right)+b^{2}\left(1-P\left(T_{a}<T_{b}\right)\right) \\
& =-a b
\end{aligned}
$$

Suppose $X$ has two values $a, b$ with $a<0<b$ and

$$
\begin{aligned}
E X & =a P(X=a)+b P(X=b) \\
& =0
\end{aligned}
$$

## Remark:

$$
P(X=a)=\frac{b}{b-a} \text { and } P(X=b)=-\frac{a}{b-a}
$$

Let $T=T_{a} \wedge T_{b}$. Then, $W_{T}$ has the same distribution as $X$. We denote

$$
\mathcal{L}\left(W_{T}\right)=\mathcal{L}(X)
$$

or

$$
W_{T}={ }_{\mathcal{D}} X
$$

### 2.6.1 Skorokhod theorem

Let $X$ be random variable with $E X=0$ and $E X^{2}<\infty$. Then, there exists a $\mathcal{F}_{t}^{W}$-stopping time $T$ such that

$$
\mathcal{L}\left(W_{T}\right)=\mathcal{L}(X) \text { and } E T=E X^{2} .
$$

Proof: Let $F(x)=P(X \leq x)$.

$$
\begin{aligned}
E X=0 & \Rightarrow \int_{-\infty}^{0} u d F(u)+\int_{0}^{\infty} v d F(v)=0 \\
& \Rightarrow-\int_{-\infty}^{0} u d F(u)=\int_{0}^{\infty} v d F(v)=C .
\end{aligned}
$$

Let $\psi$ be a bounded function with $\psi(0)=0$. Then,

$$
\begin{aligned}
& C \int_{R} \psi(x) d F(x) \\
= & C\left(\int_{0}^{\infty} \psi(v) d F(v)+\int_{-\infty}^{0} \psi(u) d F(u)\right) \\
= & \int_{0}^{\infty} \psi(v) d F(v) \int_{-\infty}^{0}-u d F(u)+\int_{-\infty}^{0} \psi(u) d F(u) \int_{0}^{\infty} v d F(v) \\
= & \int_{0}^{\infty} d F(v) \int_{-\infty}^{0} d F(u)(v \psi(u)-u \psi(v)) .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\int_{R} \psi(x) d F(x) & =C^{-1} \int_{0}^{\infty} d F(v) \int_{-\infty}^{0} d F(u)(v \psi(u)-u \psi(v)) \\
& =C^{-1} \int_{0}^{\infty} d F(v) \int_{-\infty}^{0} d F(u)(v-u)\left[\frac{v}{v-u} \psi(u)-\frac{u}{v-u} \psi(v)\right]
\end{aligned}
$$

Consider $(U, V)$ be a random vector in $R^{2}$ such that

$$
P[(U, V)=(0,0)]=F(\{0\})
$$

and for $A \subset(-\infty, 0) \times(0, \infty)$

$$
P((U, V) \in A)=C^{-1} \iint_{A} d F(u) d F(v)(v-u)
$$

If $\psi=1$,

$$
\begin{aligned}
& P((U, V) \in(-\infty, 0) \times(0, \infty)) \\
& \quad=C^{-1} \int_{0}^{\infty} d F(v) \int_{-\infty}^{0} d F(u)(v-u) \\
& \quad=C^{-1} \int_{0}^{\infty} d F(v) \int_{-\infty}^{0} d F(u)(v-u)\left[\frac{v}{v-u} \psi(u)-\frac{u}{v-u} \psi(v)\right] \\
& \quad=\int_{R} \psi(x) d F(x) \\
& \quad=\int_{R} d F(x) \\
& \quad=1
\end{aligned}
$$

and hence, $P$ is a probability measure.
Let $u<0<v$ such that

$$
\mu_{U, V}(\{u\})=\frac{v}{v-u} \text { and } \mu_{U, V}(\{v\})=-\frac{u}{v-u} .
$$

Then, by Fubini,

$$
\int \psi(x) d F(x)=E \int \psi(x) \mu_{u, v}(d x)
$$

On product space $\Omega \times \Omega^{\prime}$, let

$$
\begin{aligned}
W_{t}\left(\omega, \omega^{\prime}\right) & =W_{t}(\omega) \\
(U, V)\left(\omega, \omega^{\prime}\right) & =(U, V)\left(\omega^{\prime}\right) .
\end{aligned}
$$

$T_{U, V}$ is not a stopping time on $\mathcal{F}_{t}^{W}$.

We know that if $U=u$ and $V=v$

$$
\mathcal{L}\left(T_{U, V}\right)=\mathcal{L}(X) .
$$

Then,

$$
\begin{aligned}
E T_{U, V} & =E_{U, V} E\left(T_{U, V} \mid U, V\right) \\
& =-E U V \\
& =C^{-1} \int_{-\infty}^{0} d F(u)(-u) \int_{0}^{\infty} d F(v) v(v-u) \\
& =-\int_{0}^{\infty} d F(v)(-u)\left[C^{-1} \int_{0}^{\infty} v d F(v)-u\right] \\
& =E X^{2} .
\end{aligned}
$$

### 2.7 Embedding of sums of i.i.d. random variable in Brownian motion

Let $t_{0} \in[0, \infty)$. Then,

$$
\left\{W\left(t+t_{0}\right)-W\left(t_{0}\right), t \geq 0\right\}
$$

is a Brownian motion and independent of $\mathcal{F}_{t_{0}}$.
Let $\tau$ be a stopping time w.r.t. $\mathcal{F}_{t}^{W}$. Then, one can see that

$$
W_{t}^{*}(\omega)=W_{\tau(\omega)+t}(\omega)-W_{\tau(\omega)}(\omega)
$$

is a Brownian motion w.r.t. $\mathcal{F}_{\tau}^{W}$ where

$$
\left.\mathcal{F}_{\tau}^{W}=\left\{B \in \mathcal{F}: B \cap\{\tau \leq t\} \in \mathcal{F}_{t}^{W}\right\} \text { for all } t \leq 0\right\} .
$$

Let $V_{0}$ be countable. Then,

$$
\left\{\omega: W_{*}^{t}(\omega) \in B\right\}=\bigcup_{t_{0} \in V_{0}}\left\{\omega: W\left(t+t_{0}\right)-W\left(t_{0}\right) \in B, \tau=t_{0}\right\}
$$

For $A \in \mathcal{F}_{\tau}^{W}$,

$$
\begin{aligned}
P\left[\left\{\left(W_{t_{1}}^{*}, \ldots, W_{t_{k}}^{*}\right) \in B_{k}\right\} \cap A\right]= & \sum_{t_{0} \in V_{0}} P\left[\left\{\left(W_{t_{1}}, \ldots, W_{t_{k}}\right) \in B_{k}\right\} \cap A \cap\left\{\tau=t_{0}\right\}\right] \\
= & \sum_{t_{0} \in V_{0}} P\left(\left(W_{t_{1}}^{*}, \ldots, W_{t_{k}}^{*}\right) \in B_{k}\right) \cdot P\left(A \cap\left\{\tau=t_{0}\right\}\right) \\
& \text { (because of independence) }
\end{aligned}
$$

$$
\begin{aligned}
& =P\left(\left(W_{t_{1}}^{*}, \ldots, W_{t_{k}}^{*}\right) \in B_{k}\right) \sum_{t_{0} \in V_{0}} P\left(A \cap\left\{\tau=t_{0}\right\}\right) \\
& =P\left(\left(W_{t_{1}}^{*}, \ldots, W_{t_{k}}^{*}\right) \in B_{k}\right) P(A)
\end{aligned}
$$

Let $\tau$ be any stopping time such that

$$
\tau_{n}= \begin{cases}0, & \text { if } \tau=0 \\ \frac{k}{2^{n}}, & \text { if } \frac{k-1}{2^{n}} \leq \tau<\frac{k}{2^{n}} .\end{cases}
$$

If $\frac{k}{2^{n}} \leq t<\frac{k+1}{2^{n}}$,

$$
\left\{\tau_{n} \leq t\right\}=\left\{\tau \leq \frac{k}{2^{n}}\right\} \in \mathcal{F}_{2^{n}} \subset \mathcal{F}_{t} .
$$

Claim: $\mathcal{F}_{\tau} \subset \mathcal{F}_{\tau_{n}}$.
Proof: Suppose $C \in \mathcal{F}_{\tau}=\left\{B: B \cap\{\tau \leq t\} \in \mathcal{F}_{t}\right\}$. Then, since $\frac{k}{2^{n}} \leq t$,

$$
C \cap\{\tau \leq t\}=C \cap\left\{\tau \leq \frac{k}{2^{n}}\right\} \in \mathcal{F}_{\frac{k}{2^{n}}} \in \mathcal{F}_{t} .
$$

This completes the proof.

$$
W_{t_{1}}^{n}=W_{\tau_{n}+t}-W_{\tau_{n}} \text { is a Brownian Motion for each } n \text {, independent of } \mathcal{F}_{\tau} .
$$

Theorem 2.2: Let $X_{1}, \ldots, X_{n}$ be i.i.d. with $E X_{i}=0, E X_{i}^{2}<\infty$ for all $i$. Then, there exists a sequence of stopping time $T_{0}=0, T_{1}, \ldots, T_{n}$ such that

$$
\mathcal{L}\left(S_{n}\right)=\mathcal{L}\left(W_{T_{n}}\right),
$$

where ( $T_{i}-T_{i-1}$ ) are i.i.d.
Proof: $\left(U_{1}, V_{1}\right), \ldots,\left(U_{n}, V_{n}\right)$ i.i.d. as $(U, V)$ and independent of $W_{t}$.

$$
T_{0}=0, T_{k}=\inf \left\{t \geq T_{k-1}, W_{t+T_{k-1}}-W_{T_{k-1}} \in\left(U_{k}, V_{k}\right)\right\}
$$

$T_{k}-T_{k-1}$ are i.i.d.

$$
\begin{aligned}
\mathcal{L}\left(X_{1}\right) & =\mathcal{L}\left(W_{T_{1}}\right) \\
\mathcal{L}\left(X_{2}\right) & =\mathcal{L}\left(W_{T_{2}}-W_{T_{1}}\right) \\
\vdots & \vdots \\
\mathcal{L}\left(X_{n}\right) & =\mathcal{L}\left(W_{T_{n}}-W_{T_{n-1}}\right) .
\end{aligned}
$$

Then,

$$
\begin{aligned}
\mathcal{L}\left(\frac{S_{n}}{\sqrt{n}}\right) & =\mathcal{L}\left(\frac{W\left(T_{n}\right)}{\sqrt{n}}\right) \\
& =\mathcal{L}\left(\frac{W\left(T_{n} / n \cdot n\right)}{\sqrt{n}}\right) \\
& =\mathcal{L}\left(W\left(\frac{T_{n}}{n}\right)\right) \quad\left(\text { since } \frac{W(n t)}{\sqrt{n}} \approx W(t)\right) .
\end{aligned}
$$

Assume $E X_{1}^{2}=1$. Then

$$
\frac{T_{n}}{n} \rightarrow_{a . s .} E\left(T_{1}\right)=E X_{1}^{2}=1,
$$

and hence

$$
\frac{S_{n}}{\sqrt{n}} \rightarrow W_{1} .
$$

### 2.8 Donsker's theorem

$X_{n, 1}, \ldots, X_{n, n}$ are i.i.d. for each $n$ with $E X_{n, m}=0, E X_{n, m}^{2}<\infty$, and $S_{n, m}=X_{n, 1}+\cdots+$ $X_{n, m}=W_{\tau_{m}^{n}}$, where $\tau_{m}^{n}$ is stopping time and $W$ is Brownian motion. Define

$$
S_{n, u}=\left\{\begin{array}{l}
S_{n, m}, \quad \text { if } u=m \in\{0,1,2, \ldots, n\} ; \\
\text { linear, if } u \in[m-1, m] .
\end{array}\right.
$$

Lemma 2.2: If $\tau_{[n s]}^{n} \rightarrow s$ for $s \in[0,1]$, then with sup norm $\left\|\|_{\infty}\right.$ on $C[0,1]$,

$$
\left\|S_{n,[n \cdot]}-W(\cdot)\right\|_{\infty} \rightarrow 0 \text { in probability. }
$$

Proof: For given $\epsilon>0$, there exists $\delta>0(1 / \delta$ is an integer) such that

$$
\begin{equation*}
P\left(\left|W_{t}-W_{s}\right|<\epsilon, \text { for all } t, s \in[0,1],|t-s|<2 \delta\right)>1-\epsilon \tag{2.3}
\end{equation*}
$$

$\tau_{m}^{n}$ is increasing in $m$. For $n \geq N \delta$,

$$
P\left(\left|\tau_{n k \delta}^{n}-k \delta\right|<\delta, k=1,2, \ldots, \frac{1}{\delta}\right) \geq 1-\epsilon
$$

since $\tau_{[n s]}^{n} \rightarrow s$. For $s \in((k-1) \delta, k \delta)$, we have

$$
\tau_{[n s]}^{n}-s \geq \tau_{[n(k-1) \delta]}^{n}-k \delta
$$

$$
\tau_{[n s]}^{n}-s \leq \tau_{[n k \delta]}^{n}-(k-1) \delta
$$

Combining these, we have for $n \geq N \delta$

$$
\begin{equation*}
P\left(\sup _{0 \leq s \leq 1}\left|\tau_{[n s]}^{n}-s\right|<2 \delta\right)>1-\epsilon . \tag{2.4}
\end{equation*}
$$

For $\omega$ in event in (2.3) and (2.4), we get for $m \leq n$, as $W_{t_{m}^{n}}=S_{n, m}$,

$$
\left|W_{\tau_{m}^{n}}-W_{\frac{m}{n}}\right|<\epsilon .
$$

For $t=\frac{m+\theta}{n}$ with $0<\theta<1$,

$$
\begin{aligned}
\left|S_{n,[n t]}-W_{t}\right| \leq & (1-\theta)\left|S_{n, m}-W_{\frac{m}{n}}\right|+\theta\left|S_{n, m+1}-\frac{W_{m+1}}{n}\right| \\
& +(1-\theta)\left|W_{\frac{m}{n}}-W_{t}\right|+\theta\left|\frac{W_{n+1}}{n}-W_{t}\right|
\end{aligned}
$$

For $n \geq N_{\delta}$ with $\frac{1}{n}<2 \delta$,

$$
P\left(\left\|S_{n,[n s]}-W_{s}\right\|_{\infty} \geq 2 \epsilon\right)<2 \epsilon .
$$

Theorem 2.3: Let $f$ be bounded and continuous function on $[0,1]$. Then

$$
E f\left(S_{n,[n \cdot]}\right) \rightarrow E f(W(\cdot))
$$

Proof: For fixed $\epsilon>0$, define

$$
G_{\delta}=\left\{W, W^{\prime} \in C[0,1]:\left\|W-W^{\prime}\right\|_{\infty}<\delta \text { implies }\left|f(W)-f\left(W^{\prime}\right)\right|<\epsilon\right\}
$$

Observe that $G_{\delta} \uparrow C[0,1]$ as $\delta \downarrow 0$. Then,

$$
\left|E f\left(S_{n,[n \cdot]}\right)-E f(W(\cdot))\right| \leq \epsilon+2 M\left(P\left(G_{\delta}^{c}\right)+P\left(\left\|S_{n,[n \cdot]}-W(\cdot)\right\|>\delta\right)\right) .
$$

Since $P\left(G_{\delta}^{c}\right) \rightarrow 0$ and $P\left(\left\|S_{n,[n \cdot]}-W(\cdot)\right\|>\delta\right) \rightarrow 0$ by Lemma 2.2.
For

$$
f(x)=\max _{t}|x(t)|,
$$

we have

$$
\max _{t}\left|\frac{S_{[n t]}}{\sqrt{n}}\right| \rightarrow \max _{t}|W(t)| \text { in distribution }
$$

and

$$
\max _{1 \leq m \leq n}\left|\frac{S_{m}}{\sqrt{n}}\right| \rightarrow \max _{t}|W(t)| \text { in distribution. }
$$

Let

$$
R_{n}=1+\max _{1 \leq m \leq n} S_{m}-\min _{1 \leq m \leq n} S_{m}
$$

Then

$$
\frac{R_{n}}{\sqrt{n}} \Rightarrow \text { weakly } \max _{0 \leq t \leq 1} W(t)-\min _{0 \leq t \leq 1} W(t)
$$

We now derive from Donsker theorem invariance principle for $U$-statistics

$$
\prod_{i=1}^{[n t]}\left(1+\frac{\theta X_{i}}{\sqrt{n}}\right)=\sum_{k=1}^{[n t]} n^{-\frac{k}{2}} \sum_{1 \leq i_{1} \leq \cdots \leq i_{k} \leq n} X_{i_{1}} \cdots X_{i_{k}}
$$

where $X_{i}$ are i.i.d. and $E X_{i}^{2}<\infty$. Next, using CLT, SLLN, and the fact $P\left(\max \left|\frac{X_{i}}{\sqrt{n}}\right|>\epsilon\right)$ $\rightarrow 0$,

$$
\begin{aligned}
\log \left[\prod_{i=1}^{[n t]}\left(1+\frac{\theta X_{i}}{\sqrt{n}}\right)\right] & =\sum_{i=1}^{[n t]} \log \left(1+\frac{\theta X_{i}}{\sqrt{n}}\right) \\
& =\theta \sum_{i=1}^{[n t]} \frac{X_{i}}{\sqrt{n}}-\frac{\theta^{2}}{2} \sum_{i=1}^{[n t]} \frac{X_{i}^{2}}{n}+\frac{\theta^{3}}{3} \sum_{i=1}^{[n t]} \frac{X_{i}^{3}}{n \sqrt{n}}-\cdots \\
& \Rightarrow \theta W(t)-\frac{\theta^{2}}{2} t
\end{aligned}
$$

and hence

$$
\prod_{i=1}^{[n t]}\left(1+\frac{\theta X_{i}}{\sqrt{n}}\right) \Rightarrow e^{\theta W(t)-\frac{\theta^{2}}{2} t}
$$

### 2.9 Empirical distribution function

Let us define empirical distribution function

$$
\hat{F}_{n}(x)=\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}\left(X_{i} \leq x\right), \quad x \in \mathbb{R} .
$$

Then Glivenko-Cantelli lemma says

$$
\sup _{x}\left|\hat{F}_{n}(x)-F(x)\right| \rightarrow_{\text {a.s. }} 0 .
$$

Assume $F$ is continuous. Let $U_{i}=F\left(X_{i}\right)$. For $y \in[0,1]$, define

$$
\hat{G}(y)=\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}\left(F\left(X_{i}\right) \leq y\right) .
$$

Then by 1-1 transformation,

$$
\sqrt{n} \sup _{x}\left|\hat{F}_{n}(x)-F(x)\right|=\sqrt{n} \sup _{y \in[0,1]}\left|\hat{G}_{n}(y)-y\right| .
$$

Let $U_{1}, U_{2}, \ldots, U_{n}$ be uniform distribution and let $U_{(i)}$ be order statistic such that

$$
U_{(1)}(\omega) \leq \cdots \leq U_{(n)}(\omega) .
$$

Next,

$$
\begin{aligned}
& f\left(U_{(1)}, \ldots, U_{(n)}\right)=f\left(U_{\pi(1)}, \ldots, U_{\pi(n)}\right) \\
& f_{U_{\pi(1)}, \ldots, U_{\pi(n)}}\left(u_{1}, \ldots, u_{n}\right) \\
& =f_{U_{1}, \ldots, U_{n}}\left(u_{1}, \ldots, u_{n}\right) \text { if } u_{1}<u_{2}<\ldots<u_{n} . \\
& \\
& =\left\{\begin{array}{l}
1, \text { if } \mathbf{u} \in[0,1]^{n} ; \\
0, \text { if } \mathbf{u} \notin[0,1]^{n} .
\end{array}\right.
\end{aligned}
$$

For bounded $g$,

$$
\begin{aligned}
\operatorname{Eg}\left(U_{1}, \ldots, U_{n}\right)=\sum_{\pi \in \Pi} & \int_{u_{1}<u_{2}<\ldots<u_{n}} g_{U_{\pi(1)}, \ldots, \pi(2), \ldots, U_{\pi(n)}}^{\left(u_{1}, u_{n}\right)}\left(u_{1}, \ldots, u_{n}\right) \\
& \times f\left(u_{\pi(1)}, \ldots, u_{\pi(n)}\right) d u_{1} \cdots d u_{n}
\end{aligned}
$$

So we get

$$
f_{U_{1}, \ldots, U_{n}}\left(u_{1}, \ldots, u_{n}\right)=\left\{\begin{array}{l}
n!, \text { if } u_{1}<u_{2}<\cdots<u_{n} \\
0, \text { otherwise }
\end{array}\right.
$$

Theorem 2.4: Let $e_{j}$ be i.i.d. exponential distribution with failure rate $\lambda$. Then

$$
\mathcal{L}\left(U_{1}, \ldots, U_{n}\right)=\mathcal{L}\left(\frac{Z_{1}}{Z_{n+1}}, \ldots, \frac{Z_{n}}{Z_{n+1}}\right),
$$

where

$$
Z_{i}=\sum_{j=1}^{i} e_{j} .
$$

Proof: First we have

$$
f_{e_{1}, \ldots, e_{n+1}}\left(u_{1}, \ldots, u_{n+1}\right)= \begin{cases}\lambda^{n+1} e^{-\sum_{i=1}^{n+1} u_{i} \lambda}, & \text { if } u_{i} \geq 0 \\ 0, & \text { otherwise }\end{cases}
$$

Let $s_{i}-s_{i-1}=u_{i}$ for $i=1,2, \ldots, n+1$. Then

$$
f_{Z_{1}, \ldots, Z_{n+1}}\left(s_{1}, \ldots, s_{n+1}\right)=\prod_{i=1}^{n+1} \lambda e^{-\left(s_{i}-s_{i-1}\right) \lambda}
$$

Let

$$
\begin{aligned}
v_{i} & =\frac{s_{i}}{s_{n+1}} \text { for } i \leq n \\
v_{n+1} & =s_{n+1} .
\end{aligned}
$$

Then

$$
\begin{aligned}
f_{V_{1}, \ldots, V_{n+1}}\left(v_{1}, \ldots, v_{n+1}\right) & =\prod_{i=1}^{n}\left(\lambda e^{-\lambda v_{n+1}\left(v_{i}-v_{i-1}\right)}\right) \lambda e^{-\lambda v_{n+1}\left(1-v_{n}\right)} \\
& =\lambda^{n+1} e^{-\lambda v_{n} v_{n+1}-\lambda v_{n+1}+\lambda v_{n} v_{n+1}} v_{n+1}^{n}
\end{aligned}
$$

Exercise: Integrate with respect to $v_{n+1}$ to complete the proof.

$$
\begin{aligned}
D_{n} & =\sqrt{n} \max _{1 \leq m \leq n}\left|\frac{Z_{m}}{Z_{n+1}}-\frac{m}{n}\right| \\
& =\frac{n}{Z_{n+1}} \max _{1 \leq m \leq n}\left|\frac{Z_{m}}{\sqrt{n}}-\frac{m}{n} \frac{Z_{n+1}}{\sqrt{n}}\right| \\
& =\frac{n}{Z_{n+1}} \max _{1 \leq m \leq n}\left|\frac{Z_{m}-m}{\sqrt{n}}-\frac{m}{n} \frac{Z_{n+1}-n}{\sqrt{n}}\right| \\
& =\frac{n}{Z_{n+1}} \max _{1 \leq m \leq n}\left|W_{n}(t)-t\left(W_{n}(1)+\frac{Z_{n+1}-Z_{n}}{\sqrt{n}}\right)\right| .
\end{aligned}
$$

where

$$
W_{n}(t)= \begin{cases}\frac{Z_{m}-m}{\sqrt{n}}, \quad \text { if } t=\frac{m}{n} \\ \text { linear, } & \text { between } .\end{cases}
$$

We know that $n / Z_{n+1} \rightarrow$ a.s. $\lambda$ and

$$
E\left(\frac{Z_{n+1}-Z_{n}}{\sqrt{n}}\right)^{2}=\frac{1}{n} E e_{n}^{2} \rightarrow 0,
$$

and hence by Chevyshev's inequality,

$$
\frac{Z_{n+1}-Z_{n}}{\sqrt{n}} \rightarrow_{p} 0
$$

Since $\max (\cdot)$ is a continuous function and

$$
\left(W_{n}(\cdot)-\cdot W_{n}(1)\right) \Rightarrow_{D}(W(\cdot)-\cdot W(1)),
$$

we have for $\lambda=1$

$$
\begin{aligned}
D_{n} & =\underbrace{\frac{n}{Z_{n+1}} \max _{1 \leq m \leq n}|W_{n}(t)-t(W_{n}(1)+\underbrace{\frac{Z_{n+1}-Z_{n}}{\sqrt{n}}}_{\rightarrow p^{0}})|}_{\rightarrow a, s^{1}} \\
& \Rightarrow{ }_{D} \max _{0 \leq t \leq 1}|W(t)-t W(1)| .
\end{aligned}
$$

The process $\{W(t)-t W(1), 0 \leq t \leq 1\}$ is called Brownian Bridge.
Observe,

$$
\begin{aligned}
& P\left(W_{t_{1}} \leq x_{1}, W_{t_{2}} \leq x_{2}, \ldots, W_{t_{k}} \leq x_{k}, W_{1}=0\right) \\
& P\left(W_{t_{1}} \leq x_{1}, W_{t_{2}} \leq x_{2}, \ldots, W_{t_{k}} \leq x_{k}\right) \cdot P\left(W_{1}=0\right) \\
& P\left(W_{t_{1}} \leq x_{1}, W_{t_{2}} \leq x_{2}, \ldots, W_{t_{k}} \leq x_{k} \mid W(1)=0\right) \\
& \quad=\frac{P\left(W_{t_{1}} \leq x_{1}, W_{t_{2}} \leq x_{2}, \ldots, W_{t_{k}} \leq x_{k}, W(1)=0\right)}{P\left(W_{1}=0\right)} \\
& \quad P\left(W_{t_{1}} \leq x_{1}, W_{t_{2}} \leq x_{2}, \ldots, W_{t_{k}} \leq x_{k}\right)
\end{aligned}
$$

$\left\{W_{t}^{0}\right\}$ is called Brownian Bridge if

$$
\begin{aligned}
E W_{t}^{0} W_{s}^{0} & =E\left(W_{t}-t W(1)\right)\left(W_{s}-s W(1)\right) \\
& =\min (t, s)-s t-t s+t s \\
& =s(1-t)
\end{aligned}
$$

for $s \leq t$.
The above calculations show the conditional distribution of Brownian Motion given $\{W(1)=0\}$ is the distribution of $\left\{W_{t}^{0}, 0 \leq t \leq 1\right\}$.

### 2.10 Weak convergence of probability measure on Polish space

In chapter 3, we shall discuss the weak convergence in the space of functions with jump discontinuities with a topology that makes it a Polish space. This will be applied in chapter 4 for the weak convergence of semi-martingales. To make this book selfcontained, we present to the readers the weak convergence of probability measures on Polish space.

Let $(S, \rho)$ be a complete separable metric space. $\left\{P_{n}\right\}$, a sequence of probability measure on $\mathcal{B}(S)$, converges weakly to $P$ if for all bounded continuous function on $\mathcal{X}$

$$
\int f d P_{n} \longrightarrow \int f d P
$$

and we write $P_{n} \Rightarrow P$. Let $\mathcal{B}(S)=\mathcal{S}$.
Theorem 2.5: Every probability measure $P$ on $(S, \mathcal{S})$ is regular, that is, for every $\mathcal{S}$-set $A$ and every $\epsilon$ there exist a closed set $F$ and an open set $G$ such that $F \subset A \subset G$ and $P(G-F)<\epsilon$.

Proof: Denote the metric on $S$ by $\rho(x, y)$ and the distance from $x$ to $A$ by $\rho(x, A)=$ $\inf \{\rho(x, y): y \in A\}$. If $A$ is closed, we can take $F=A$ and $G=A^{\delta}=\{x: \rho(x, A)<\delta\}$ for some $\delta$, since the latter sets decrease to $A$ as $\delta \downarrow 0$. Hence, we need only show that the class $\mathcal{G}$ of $\mathcal{S}$-sets with the asserted property is a $\sigma$-field. Given sets $A_{n}$ in $\mathcal{G}$, choose closed sets $F_{n}$ and open sets $G_{n}$ such that $F_{n} \subset A_{n} \subset G_{n}$ and $P\left(G_{n}-F_{n}\right)<\epsilon / 2^{n+1}$. If $G=\bigcup_{n} G_{n}$, and if $F=\bigcup_{n \leq n_{0}} F_{n}$, with $n_{0}$ so chosen that

$$
P\left(\bigcup_{n} F_{n}-F\right)<\frac{\epsilon}{2}
$$

then $F \subset \bigcup_{n} A_{n} \subset G$ and $P(G-F)<\epsilon$. Thus, $\mathcal{G}$ is closed under the formation of countable unions; since it is obviously closed under complementation, $\mathcal{G}$ is a $\sigma$-field.

Equation (2.5) implies that $P$ is completely determined by the values of $P(F)$ for closed sets $F$. The next theorem shows that $P$ is also determined by the values of $\int f d P$ for bounded, continuous $f$. The proof depends on the approximation of indicator $I_{F}$ by such an $f$, and the function $f(x)=(1-\rho(x, F) / \epsilon)^{+}$works. It is bounded, and it is continuous, even uniformly continuous, because $|f(x)-f(y)| \leq \rho(x, y) / \epsilon . x \in F$ implies $f(x)=1$, while $x \notin F^{\epsilon}$ implies $\rho(x, F) \geq \epsilon$ and hence $f(x)=0$. Therefore,

$$
\begin{equation*}
I_{F}(x) \leq f(x)=(1-\rho(x, F) / \epsilon)^{+} \leq I_{F e}(x) . \tag{2.5}
\end{equation*}
$$

Theorem 2.6: Probability measures $P$ and $Q$ on $\mathcal{S}$ coincide if and only if $\int f d P=\int f d Q$ for all bounded, uniformly continuous real functions $f$.

Proof: $(\Rightarrow)$ Trivial.
$(\Leftarrow)$ For the bounded, uniformly continuous $f$ of (2.5), $P(F) \leq \int f d P=\int f d Q \leq$ $Q\left(F^{\epsilon}\right)$. Letting $\epsilon \downarrow 0$ gives $P(F)=Q(F)$, provided $F$ is closed. By symmetry and (Theorem 2.5), $P=Q$.

The following notion of tightness plays a fundamental role both in the theory of weak convergence and in its applications. A probability measure $P$ on $(S, \mathcal{S})$ is tight if for each $\epsilon$ there exists a compact set $K$ such that $P(K) \geq 1-\epsilon$. By (2.5), $P$ is tight if and only if for each $A \in \mathcal{S}$

$$
P(A)=\sup \{P(K): K \subset A, \quad K \text { is compact. }\}
$$

Theorem 2.7: If $S$ is separable and complete, then each probability measure on $(S, \mathcal{S})$ is tight.

Proof: Since $S$ is separable, there is, for each $k$, a sequence $A_{k 1}, A_{k 2}, \ldots$ of open $1 / k$-balls covering $S$. Choose $n_{k}$ large enough that

$$
P\left(\bigcup_{i \leq n_{k}} A_{k i}\right)>1-\frac{\epsilon}{2^{k}} .
$$

By the completeness hypothesis, the totally bounded set

$$
\bigcap_{k \geq 1} \bigcup_{i \leq n_{k}} A_{k i}
$$

has compact closure $K$. However, clearly, $P(K)>1-\epsilon$. This completes the proof.
The following theorem provides useful conditions equivalent to weak convergence; any of them could serve as the definition. A set $A$ in $\mathcal{S}$ whose boundary $\partial A$ satisfies $P(\partial A)=0$ is called $P$-continuity set. Let $P_{n}, P$ be probability measures on ( $\mathcal{X}, \mathcal{B}(\mathcal{X}))$.

Theorem 2.8 (The Portmanteau theorem): The following are equivalent.

1. For bounded and continuous $f$

$$
\lim _{n \rightarrow \infty} \int f d P_{n}=\int f d P
$$

2. For closed set $F$

$$
\limsup _{n \rightarrow \infty} P_{n}(F) \leq P(F) .
$$

3. For open set $G$

$$
\liminf _{n \rightarrow \infty} P_{n}(G) \geq P(G)
$$

4. For all set $A$ with $P(\partial A)=0$

$$
\lim _{n \rightarrow \infty} P_{n}(A)=P(A)
$$



Proof: (1) $\rightarrow$ (2) : Let $f_{k}(x)=\gamma_{k}(\rho(x, F))$, with $\gamma_{k}(x)$ given by the above graph. First of all, we know that

$$
f_{k}(x) \searrow 1_{F}(x) .
$$

Then, for any $\delta>0$, there exists $K$ such that for all $k \geq K$

$$
\begin{aligned}
\limsup _{n \rightarrow \infty} P_{n}(F) & =\limsup _{n \rightarrow \infty} \int 1_{F} d P_{n} \\
& \leq \limsup _{n \rightarrow \infty} \int f_{k} d P_{n} \\
& =\lim _{n \rightarrow \infty} \int f_{k} d P_{n} \\
& =\int f_{k} d P \\
& \leq P(F)+\delta
\end{aligned}
$$

The last inequality follows from the fact that

$$
\int_{F} f_{n} d P \searrow P(F)
$$

As a result, for all $\delta>0$, we have

$$
\limsup _{n \rightarrow \infty} P_{n}(F) \leq P(F)+\delta
$$

(2) $\rightarrow(3)$

Let $G=F^{c}$. Then, it follows directly.
$(2)+(3) \rightarrow(4)$ trivial.
$(4) \rightarrow(1)$ Approximate $f$ let simple functions $f_{n}$ of sets $A$ with $P(\bar{A})=P(A)$.

Theorem 2.9: A necessary and sufficient condition for $P_{n} \Rightarrow P$ is that each subsequence $\left\{P_{n_{i}}\right\}$ contain a further subsequence $\left\{P_{n_{i}(m)}\right\}$ converging weakly to $P$.

Proof: The necessary is easy. As for sufficiency, if $P_{n}$ does not converge weakly to $P$, then there exists some bounded and continuous $f$ such that $\int f d P_{n}$ does not converge to $\int f d P$. However, for some positive $\epsilon$ and some subsequence $P_{n_{i}}$,

$$
\left|\int f d P_{n_{i}}-\int f d P\right|>\epsilon
$$

for all $i$, and no further subsequence can converge weakly to $P$.
Suppose that $h$ maps $\mathcal{X}$ into another metric space $\mathcal{X}^{\prime}$, with metric $\rho^{\prime}$ and Borel $\sigma$-field $\mathcal{B}\left(\mathcal{X}^{\prime}\right)$. If $h$ is measurable $\mathcal{X} / \mathcal{X}^{\prime}$, then each probability $P$ on $(\mathcal{X}, \mathcal{B}(\mathcal{X}))$ induces on ( $\left.\mathcal{X}^{\prime}, \mathcal{B}\left(\mathcal{X}^{\prime}\right)\right)$ a probability $P \circ h^{-1}$ defined as usual by $P \circ h^{-1}(A)=P\left(h^{-1}(A)\right)$. We need conditions under which $P_{n} \Rightarrow P$ implies $P_{n} \circ h^{-1} \Rightarrow P \circ h^{-1}$. One such condition is that $h$ is continuous: If $f$ is bounded and continuous on $\mathcal{X}^{\prime}$, then $f h$ is bounded and continuous on $\mathcal{X}$, and by change of variable, $P_{n} \Rightarrow P$ implies

$$
\begin{equation*}
\int_{\mathcal{X}^{\prime}} f(y) P_{n} \circ h^{-1}(d y)=\int_{\mathcal{X}} f(h(x)) P_{n}(d x) \rightarrow \int_{\mathcal{X}} f(h(x)) P(d x)=\int_{\mathcal{X}^{\prime}} f(y) P \circ h^{-1}(d y) \tag{2.6}
\end{equation*}
$$

Let $D_{h}$ be discontinuity set of $h$.
Theorem 2.10: Let $(\mathcal{X}, \rho)$ and $\left(\mathcal{X}^{\prime}, \rho^{\prime}\right)$ be two Polish space and

$$
h: \mathcal{X} \rightarrow \mathcal{X}^{\prime}
$$

with $P\left(D_{h}\right)=0$. Then, $P_{n} \Rightarrow P$ implies

$$
P_{n} \circ h^{-1} \Rightarrow P \circ h^{-1} .
$$

Proof: Since

$$
h^{-1}(F) \subset \overline{h^{-1}(F)} \subset D_{h} \cup h^{-1}(F),
$$

$$
\begin{aligned}
\limsup _{n \rightarrow \infty} P_{n}\left(h^{-1}(F)\right) & \leq \limsup _{n \rightarrow \infty} P_{n}\left(\overline{h^{-1}(F)}\right) \\
& \leq P\left(D_{h} \cup h^{-1}(F)\right) \\
& \leq P\left(h^{-1}(F)\right) \text { (since } D_{h} \text { is a set of } P \text {-measure zero). }
\end{aligned}
$$

Therefore, for all closed set $F$,

$$
\limsup _{n \rightarrow \infty} P_{n} \circ h^{-1}(F) \leq P \circ h^{-1}(F)
$$

and hence, by Theorem 2.8, the proof is completed.

Let $X_{n}$ and $X$ be random variables( $\mathcal{X}$-valued). Then, we say $X_{n} \rightarrow_{\mathcal{D}} X$ if $P \circ X_{n}^{-1} \Rightarrow$ $P \circ X^{-1}$.

Observation: If $X_{n} \rightarrow_{\mathcal{D}} X$ and $\rho\left(X_{n}, Y_{n}\right) \rightarrow_{p} 0$, then

$$
Y_{n} \rightarrow_{\mathcal{D}} X .
$$

Remark: We use the following property of limsup and liminf.

$$
\limsup _{n \rightarrow \infty}\left(a_{n}+b_{n}\right) \leq \limsup _{n \rightarrow \infty} a_{n}+\limsup _{n \rightarrow \infty} b_{n}
$$

and

$$
\liminf _{n \rightarrow \infty}\left(a_{n}+b_{n}\right) \geq \liminf _{n \rightarrow \infty} a_{n}+\liminf _{n \rightarrow \infty} b_{n} .
$$

Proof: Consider closed set $F$. Let $F^{\epsilon}=\{x: \rho(x, F) \leq \epsilon\}$. Then, $F^{\epsilon} \searrow F$ as $\epsilon \rightarrow 0$ and

$$
\begin{aligned}
\left\{X_{n} \notin F^{\epsilon}\right\} & =\left\{\omega: X_{n}(\omega) \notin F^{\epsilon}\right\} \\
& =\left\{\omega: \rho\left(X_{n}(\omega), F\right)>\epsilon\right\}
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\omega \in\left\{X_{n} \notin F^{\epsilon}\right\} \cap\left\{\rho\left(X_{n}, Y_{n}\right)<\epsilon\right\} & \Rightarrow \rho\left(X_{n}(\omega), F\right)>\epsilon \text { and } \rho\left(X_{n}(\omega), Y_{n}(\omega)\right)<\epsilon \\
& \Rightarrow \rho\left(Y_{n}(\omega), F\right)>0 \text { (draw graph.) } \\
& \Rightarrow Y_{n}(\omega) \notin F \\
& \Rightarrow \omega \in\left\{Y_{n} \notin F\right\}
\end{aligned}
$$

Thus,

$$
\left\{X_{n} \notin F^{\epsilon}\right\} \cap\left\{\rho\left(X_{n}, Y_{n}\right)<\epsilon\right\} \subset\left\{Y_{n} \notin F\right\} .
$$

Therefore,

$$
P\left(Y_{n} \in F\right) \leq P\left(\rho\left(X_{n}, Y_{n}\right)>\epsilon\right)+P\left(X_{n} \in F^{\epsilon}\right) .
$$

Let $P_{n}^{Y}=P \circ Y_{n}^{-1}$ and $P^{X}=P \circ X^{-1}$. Then, for all $\epsilon>0$,

$$
\begin{aligned}
\limsup _{n \rightarrow \infty} P_{n}(F) & =\limsup _{n \rightarrow \infty} P\left(Y_{n} \in F\right) \\
& \leq \limsup _{n \rightarrow \infty} P(\underbrace{\rho\left(X_{n}, Y_{n}\right)}_{\rightarrow_{p} 0}>\epsilon)+\limsup _{n \rightarrow \infty} P\left(X_{n} \in F^{\epsilon}\right)
\end{aligned}
$$

$$
\begin{aligned}
& =\limsup _{n \rightarrow \infty} P\left(X_{n} \in F^{\epsilon}\right) \\
& =P\left(X \in F^{\epsilon}\right)\left(\text { since } X_{n} \Rightarrow_{\mathcal{D}} X\right) .
\end{aligned}
$$

Therefore, for all closed set $F$, we have

$$
\limsup _{n \rightarrow \infty} P_{n}^{Y}(F) \leq P^{X}(F),
$$

and hence, by Theorem 2.5,

$$
P_{n}^{Y} \Longrightarrow P^{X}
$$

which implies $Y_{n} \Rightarrow_{\mathcal{D}} X$.
We say that a family of probability measure $\Pi \subset \mathcal{P}(\mathcal{X})$ is tight if, given $\epsilon>0$, there exists compact $K_{\epsilon}$ such that

$$
P\left(K_{\epsilon}\right)>1-\epsilon \text { for all } P \in \Pi .
$$

### 2.10.1 Prokhorov theorem

Definition 2.4: $\Pi$ is relatively compact if, for $\left\{P_{n}\right\} \subset \Pi$, there exists a subsequence $\left\{P_{n_{i}}\right\} \subset \Pi$ and probability measure $P$ (not necessarily an element of $\Pi$ ), such that

$$
P_{n_{i}} \Rightarrow P .
$$

Even though $P_{n_{i}} \Rightarrow P$ makes no sense if $P(\mathcal{X})<1$, it is to be emphasized that we do require $P(\mathcal{X})=1$ and we disallow any escape of mass, as discussed below. For the most part, we are concerned with the relative compactness of sequences $\left\{P_{n}\right\}$; this means that every subsequence $\left\{P_{n_{i}}\right\}$ contains a further subsequence $\left\{P_{n_{i}(m)}\right\}$, such that $P_{n_{i}(m)} \Rightarrow P$ for some probability measure $P$.

Example: Suppose we know of probability measures $P_{n}$ and $P$ on $(C, \mathcal{C})$ that the finite-dimensional distributions of $P_{n}$ converges weakly to those of $P: P_{n} \pi_{t_{1}, \ldots, t_{k}}^{-1} \Rightarrow$ $P \pi_{t_{1}, \ldots, t_{k}}^{-1}$ for all $k$ and all $t_{1}, \ldots, t_{k}$. Notice that $P_{n}$ need not converge weakly to $P$. Suppose, however, that we also know that $\left\{P_{n}\right\}$ is relatively compact. Then each $\left\{P_{n}\right\}$ contains some $\left\{P_{n_{i}(m)}\right\}$ converging weakly to some $Q$. Since the mapping theorem then gives $P_{n_{i}(m)} \pi_{t_{1}, \ldots, t_{k}}^{-1} \Rightarrow Q \pi_{t_{1}, \ldots, t_{k}}^{-1}$ and since $P_{n} \pi_{t_{1}, \ldots, t_{k}}^{-1} \Rightarrow P \pi_{t_{1}, \ldots, t_{k}}^{-1}$ by assumption, we have $P \pi_{t_{1}, \ldots, t_{k}}^{-1}=Q \pi_{t_{1}, \ldots, t_{k}}^{-1}$ for all $t_{1}, \ldots, t_{k}$. Thus, the finite-dimensional distributions of $P$ and $Q$ are identical, and since the class $\mathcal{C}_{f}$ of finite-dimensional sets is a separating class, $P=Q$. Therefore, each subsequence contains a further subsequence converging weakly to $P$, not to some fortuitous limit, but specifically to $P$. It follows by Theorem 2.9 that the entire sequence $\left\{P_{n}\right\}$ converges weakly to $P$.

Therefore, $\left\{P_{n}\right\}$ is relatively compact and the finite-dimensional distributions of $P_{n}$ converge weakly to those of $P$, then $P_{n} \Rightarrow P$. This idea provides a powerful method for proving weak convergence in $C$ and other function spaces. Note that if $\left\{P_{n}\right\}$ does converge weakly to $P$, then it is relatively compact, so that this is not too strong a condition.

Theorem 2.11: Suppose $(\mathcal{X}, \rho)$ is a Polish space and $\Pi \subset \mathcal{P}(\mathcal{X})$ is relatively compact, then it is tight.

This is the converse half of Prohorov's theorem.

Proof: Consider open sets, $G_{n} \nearrow \mathcal{X}$. For each $\epsilon>0$, there exists $n$, such that for all $P \in \Pi$

$$
P\left(G_{n}\right)>1-\epsilon .
$$

Otherwise, for each $n$, we can find $P_{n}$ such that $P_{n}\left(G_{n}\right)<1-\epsilon$. Then by by relative compactness, there exists $\left\{P_{n_{i}}\right\} \subset \Pi$ and probability measure $Q \in \Pi$ such that $P_{n_{i}} \Rightarrow Q$. Thus,

$$
\begin{aligned}
Q\left(G_{n}\right) & \leq \liminf _{i \rightarrow \infty} P_{n_{i}}\left(G_{n}\right) \\
& \leq \liminf _{i \rightarrow \infty} P_{n_{i}}\left(G_{n_{i}}\right)\left(\text { since } n_{i} \geq n \text { and hence } G_{n} \subset G_{n_{i}}\right) \\
& <1-\epsilon
\end{aligned}
$$

Since $G_{n} \nearrow \mathcal{X}$,

$$
\begin{aligned}
1 & =Q(\mathcal{X}) \\
& =\lim _{n \rightarrow \infty} Q\left(G_{n}\right) \\
& <1-\epsilon,
\end{aligned}
$$

which is contradiction. Let $A_{k_{m}}, m=1,2, \ldots$ be open ball with radius $\frac{1}{k_{m}}$, covering $\mathcal{X}$ (separability). Then, there exists $n_{k}$ such that for all $P \in \Pi$,

$$
P\left(\bigcup_{i \leq n_{k}} A_{k_{i}}\right)>1-\frac{\epsilon}{2^{k}}
$$

Then, let

$$
K_{\epsilon}=\overline{\bigcap_{k \geq 1} \bigcup_{i \leq n_{k}} A_{k_{i}}},
$$

where $\bigcap_{k \geq 1} \bigcup_{i \leq n_{k}} A_{k_{i}}$ is totally bounded set. Then, $K_{\epsilon}$ is compact(completeness), and $P\left(K_{\epsilon}\right)>1-\epsilon$.

Remark: The last inequality is from the following. Let $B_{i}$ be such that $P\left(B_{i}\right)>1-\frac{\epsilon}{2^{i}}$. Then,

$$
\begin{aligned}
P\left(B_{i}\right)>1-\frac{\epsilon}{2^{i}} & \Rightarrow P\left(B_{i}^{c}\right) \leq \frac{\epsilon}{2^{i}} \\
& \Rightarrow P\left(\cup_{i=1}^{\infty} B_{i}^{c}\right) \leq \epsilon \\
& \Rightarrow P\left(\cap_{i=1}^{\infty} B_{i}\right)>1-\epsilon .
\end{aligned}
$$

### 2.11 Tightness and compactness in weak convergence

Theorem 2.12: If $\Pi$ is tight, then for $\left\{P_{n}\right\} \subset \Pi$, there exists a subsequence $\left\{P_{n_{i}}\right\} \subset\left\{P_{n}\right\}$ and probability measure $P$ such that

$$
P_{n_{i}} \Rightarrow P .
$$

Proof: Choose compact $K_{1} \subset K_{2} \subset \ldots$, such that for all $n$

$$
P_{n}\left(K_{u}\right)>1-\frac{1}{u}
$$

from tightness condition. Look at $\bigcup_{u} K_{u}$. We know that there exists a countable family of open sets, $\mathcal{A}$, such that if $x \in \cup_{u} K_{u}$ and $G$ is open, then

$$
x \in A \subset \bar{A} \subset G
$$

for some $A \in \mathcal{A}$. Let

$$
\mathcal{H}=\{\emptyset\} \cup\left\{\text { finite union of sets of the form } \bar{A} \cap K_{u} u \geq 1, A \in \mathcal{A}\right\} .
$$

Then, $\mathcal{H}$ is a countable family. Using the Cantor diagonalization method, there exists $\left\{n_{i}\right\}$ such that for all $H \in \mathcal{H}$,

$$
\alpha(H)=\lim _{i \rightarrow \infty} P_{n_{i}}(H) .
$$

Our aim is to construct a probability measure $P$ such that for all open set $G$,

$$
\begin{equation*}
P(G)=\sup _{H \subset G} \alpha(H) . \tag{2.7}
\end{equation*}
$$

Suppose we showed (2.7) above. Consider an open set $G$. Then, for $\epsilon>0$, there exists $H_{\epsilon} \subset G$ such that

$$
\begin{aligned}
P(G) & =\sup _{H \subset G} \alpha(H) \\
& <\alpha\left(H_{\epsilon}\right)+\epsilon \\
& =\lim _{i} P_{n_{i}}\left(H_{\epsilon}\right)+\epsilon \\
& =\liminf _{i} P_{n_{i}}\left(H_{\epsilon}\right)+\epsilon \\
& \leq \liminf _{i} P_{n_{i}}(G)+\epsilon
\end{aligned}
$$

and hence, for all open set $G$,

$$
P(G) \leq \liminf _{i} P_{n_{i}}(G),
$$

which is equivalent to $P_{n_{i}} \Rightarrow P$.
Observe $\mathcal{H}$ is closed under finite union and

1. $\alpha\left(H_{1}\right) \leq \alpha\left(H_{2}\right)$ if $H_{1} \subset H_{2}$
2. $\alpha\left(H_{1} \cup H_{2}\right)=\alpha\left(H_{1}\right)+\alpha\left(H_{2}\right)$ if $H_{1} \cap H_{2}=\emptyset$
3. $\alpha\left(H_{1} \cup H_{2}\right) \leq \alpha\left(H_{1}\right)+\alpha\left(H_{2}\right)$.

Define for open set $G$

$$
\begin{equation*}
\beta(G)=\sup _{H \subset G} \alpha(H) . \tag{2.8}
\end{equation*}
$$

Then, $\alpha(\emptyset)=\beta(\emptyset)=0$ and $\beta$ is monotone.
Define for $M \subset \mathcal{X}$

$$
\gamma(M)=\inf _{M \subset G} \beta(G)
$$

Then,

$$
\begin{aligned}
\gamma(M) & =\inf _{M \subset G} \beta(G) \\
& =\inf _{M \subset G}\left(\sup _{H \subset G} \alpha(H)\right) \\
\gamma(G) & =\inf _{G \subset G^{\prime}} \beta\left(G^{\prime}\right) \\
& =\beta(G)
\end{aligned}
$$

$M$ is $\gamma$-measurable if for all $L \subset \mathcal{X}$

$$
\gamma(L) \geq \gamma(M \cap L)+\gamma\left(M^{c} \cap L\right) .
$$

We shall prove that $\gamma$ is outer measure, and hence open and closed sets are $\beta$-measurable.
$\gamma$-measurable sets $M$ form a $\sigma$-field, $\mathcal{M}$, and

$$
\left.\gamma\right|_{\mathcal{M}}
$$

is a measure.

Claim: Each closed set is in $\mathcal{M}$ and

$$
P=\left.\gamma\right|_{\mathcal{B}(\mathcal{X})}
$$

open set $G$

$$
P(G)=\gamma(G)=\beta(G) .
$$

Note that $P$ is a probability measure. $K_{u}$ has finite covering of sets in $\mathcal{A}$ when $K_{u} \in \mathcal{H}$.

$$
\begin{aligned}
1 & \geq P(\mathcal{X}) \\
& =\beta(\mathcal{X}) \\
& =\sup _{u} \alpha\left(K_{u}\right) \\
& =\sup _{u}\left(1-\frac{1}{u}\right) \\
& =1 .
\end{aligned}
$$

Step 1: If $F \subset G$ ( $F$ is closed and $G$ is open), and if $F \subset H$ for some $H \in \mathcal{H}$, then there exists some $H_{0} \in \mathcal{H}$ such that

$$
H \subset H_{0} \subset G
$$

Proof: Consider $x \in F$ and $A_{x} \in \mathcal{A}$ such that

$$
x \in A_{x} \subset \bar{A}_{x} \subset G .
$$

Since $F$ is closed subset of compact, $F$ is compact. Since $A_{x}$ covers $F$, there exists finite subcovers $A_{x_{1}}, A_{x_{2}}, \ldots, A_{x_{k}}$. Take

$$
H_{0}=\bigcup_{i=1}^{k}\left(\overline{A_{x_{i}}} \cap K_{u}\right) .
$$

Step 2: $\beta$ is finitely sub-additive on open set. Suppose $H \subset G_{1} \cup G_{2}$ and $H \in \mathcal{H}$. Let

$$
\begin{aligned}
& F_{1}=\left\{x \in H: \rho\left(x, G_{1}^{c}\right) \geq \rho\left(x, G_{2}^{c}\right)\right\} \\
& F_{2}=\left\{x \in H: \rho\left(x, G_{2}^{c}\right) \geq \rho\left(x, G_{1}^{c}\right)\right\} .
\end{aligned}
$$

If $x \in F_{1}$ but not in $G_{1}$, then $x \in G_{2}$, and hence $x \in H$. Suppose $x$ is not in $G_{2}$. Then $x \in G_{2}^{c}$, and hence, $\rho\left(x, G_{2}^{c}\right)>0$. Therefore,

$$
\begin{aligned}
0 & =\rho\left(x, G_{1}^{c}\right)\left(\text { since } x \in G_{1}^{c}\right) \\
& <\underbrace{\rho\left(x, G_{2}^{c}\right)}_{>0},
\end{aligned}
$$

which contradicts $x \in F_{1}$, and hence contradicts $\rho\left(x, G_{1}^{c}\right) \geq \rho\left(x, G_{2}^{c}\right)$. Similarly, if $x \in$ $F_{2}$ but not in $G_{2}$, then $x \in G_{1}$. Therefore, $F_{1} \subset G_{1}$ and $F_{2} \subset G_{2}$. Since $F_{i}$ 's are closed, by step 1, there exist $H_{1}$ and $H_{2}$ such that

$$
F_{1} \subset H_{1} \subset G_{1}
$$

and

$$
F_{2} \subset H_{2} \subset G_{2} .
$$

Therefore,

$$
\begin{aligned}
& \alpha(H) \leq \alpha\left(H_{1}\right)+\alpha\left(H_{2}\right) \\
& \beta(G) \leq \beta\left(G_{1}\right)+\beta\left(G_{2}\right) .
\end{aligned}
$$

Step 3: $\beta$ is countably sub-additive on open-set $H \subset \bigcup_{n} G_{n}$, where $G_{n}$ is an open set. Since $H$ is compact (union of compacts), there exist a finite subcovers, i.e., there exists $n_{0}$ such that

$$
H \subset \bigcup_{n \leq n_{0}} G_{n}
$$

and

$$
\begin{aligned}
\alpha(H) & \leq \beta(H) \\
& \leq \beta\left(\bigcup_{n \leq n_{0}} G_{n}\right) \\
& =\sum_{n \leq n_{0}} \beta\left(G_{n}\right) \\
& =\sum_{n} \beta\left(G_{n}\right) .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\beta\left(\bigcup_{n} G_{n}\right) & =\sup _{H \subset \cup_{n} G_{n}} \alpha(H) \\
& \leq \sup _{H \subset \cup_{n} G_{n}} \sum_{n} \beta\left(G_{n}\right) \\
& =\sum_{n} \beta\left(G_{n}\right) .
\end{aligned}
$$

Step 4: $\gamma$ is an outer measure. We know $\gamma$ is monotonic by definition and is countably sub-additive. Given $\epsilon>0$ and subsets $\left\{M_{n}\right\} \subset \mathcal{X}$, choose open sets $G_{n}, M_{n} \subset G_{n}$ such that

$$
\begin{aligned}
\beta\left(G_{n}\right) & \leq \gamma\left(M_{n}\right)+\frac{\epsilon}{2^{n}} \\
\gamma\left(\bigcup_{n} M_{n}\right) & \leq \beta\left(\bigcup_{n} G_{n}\right) \\
& =\sum_{n} \beta\left(G_{n}\right) \\
& =\sum_{n} \gamma\left(M_{n}\right)+\epsilon .
\end{aligned}
$$

Step 5: $F$ is closed $G$ is open.

$$
\beta(G) \geq \gamma(F \cap G)+\gamma\left(F^{c} \cap G\right) .
$$

Choose $\epsilon>0$ and $H_{1} \in \mathcal{H}, H_{1} \subset F^{c} \cap G$ such that

$$
\alpha\left(H_{1}\right)>\beta\left(G \cap F^{c}\right)-\epsilon .
$$

Chose $H_{0}$ such that

$$
\alpha\left(H_{0}\right)>\beta\left(H_{1}^{c} \cap G\right)-\epsilon .
$$

Then, $H_{0}, H_{1} \subset G$, and $H_{0} \cap H_{1}=\emptyset$,

$$
\begin{aligned}
\beta(G) & \geq \alpha\left(H_{0} \cup H_{1}\right) \\
& =\alpha\left(H_{0}\right)+\alpha\left(H_{1}\right) \\
& >\beta\left(H_{1}^{c} \cap G\right)+\beta\left(F^{c} \cap G\right)-2 \epsilon \\
& \geq \gamma(F \cap G)+\gamma\left(F^{c} \cap G\right)-2 \epsilon .
\end{aligned}
$$

Step 6: If $F \in \mathcal{M}$, then $F$ are all closed. If $G$ is open and $L \subset G$, then

$$
\beta(G) \geq \gamma(F \cap L)+\gamma\left(F^{c} \cap L\right) .
$$

Then,

$$
\begin{aligned}
& \inf \beta(G) \geq \inf \left(\gamma(F \cap L)+\gamma\left(F^{c} \cap L\right)\right) \\
\Longrightarrow & \gamma(L) \geq \gamma(F \cap L)+\gamma\left(F^{c} \cap L\right) .
\end{aligned}
$$

## 3 Weak convergence on $C[0,1]$ and $D[0, \infty)$

We now use the techniques developed in sections 2.10 and 2.11 to the case of space $C[0,1]$, which is a complete separable metric space with sup norm. In later sections, we consider the space of functions with discontinuities of the first kind. Clearly, one has to define distance on this space denoted by $D[0,1]$, such that the space is complete separable metric space. This was done by Skorokhod with the introduction of the topology called Skorokhod topology. In view of the fact that convergence of finite dimensional distributions determines the limiting measure and the Prokhorov theorem from chapter 2 , we need to characterize compact sets in $C[0,1]$ with sup norm and $D[0,1]$ with Skorokhod topology, which is described in section 3.4. The tightness in this case is described in section 3.7.

### 3.1 Structure of compact sets in $C[0,1]$

Let $\mathcal{X}$ be a complete separable metric space. We showed that $\Pi$ is tight if and only if $\Pi$ is relatively compact. Consider $P_{n}$ is measure on $C[0,1]$ and let

$$
\pi_{t_{1}, \ldots, t_{k}}(x)=\left(x\left(t_{1}\right), \ldots, x\left(t_{k}\right)\right)
$$

and suppose that

$$
P_{n} \circ \pi_{t_{1}, \ldots, t_{k}}^{-1} \Longrightarrow P \circ \pi_{t_{1}, \ldots, t_{k}}^{-1}
$$

does not imply

$$
P_{n} \Longrightarrow P
$$

on $C[0,1]$. However, $P_{n} \circ \pi_{t_{1}, \ldots, t_{k}}^{-1} \Rightarrow P \circ \pi_{t_{1}, \ldots, t_{k}}^{-1}$ and $\left\{P_{n}\right\}$ is tight. Then, $P_{n} \Rightarrow P$.
Proof: Since tightness is implied (as we proved), there exists a subsequence $\left\{P_{n_{i}}\right\}$ of $\left\{P_{n}\right\}$ such that the sequence converges weakly to a probability measure $Q$, i.e $P_{n_{i}} \Longrightarrow Q$.

Hence, by Theorem 2.9, $P_{n i} \circ \pi_{t_{1} \cdots t_{k}}^{-1} \longrightarrow Q \circ \pi_{t_{1} \ldots t_{k}}^{-1}$, Giving $P=Q$.
However, all subsequences have the same limit. Hence, $P_{n}(\rightarrow) P$.
What is a compact set in $C[0,1]$ ?

### 3.1.1 Arzela-Ascoli theorem

Definition 3.1: The uniform norm (or sup norm) assigns to real- or complex-valued bounded functions $f$ defined on a set $S$ the non-negative number

$$
\|f\|_{\infty}=\|f\|_{\infty, S}=\sup \{|f(x)|: x \in S\} .
$$

This norm is also called the supremum norm, the Chebyshev norm, or the infinity norm. The name "uniform norm" derives from the fact that a sequence of functions $\left\{f_{n}\right\}$ converges to $f$ under the metric derived from the uniform norm if and only if $f_{n}$ converges to $f$ uniformly.

Theorem 3.1: The set $A \subset C[0,1]$ is relatively compact in sup topology if and only if

$$
\begin{aligned}
& \text { (i) } \sup _{x \in A}|x(0)|<\infty . \\
& \text { (ii) } \lim _{\delta}\left(\sup _{x \in A} w_{x}(\delta)\right)=0 .
\end{aligned}
$$

Remark (modulus of continuity): Here

$$
w_{x}(\delta)=\sup _{|s-t| \leq \delta}|x(t)-x(s)| .
$$

Proof: Consider function

$$
f: C[0,1] \rightarrow R,
$$

such that $f(x)=x(0)$.
Claim: $f$ is continuous.
Proof of claim: We want to show that for $\epsilon>0$, there exists $\delta$ such that

$$
\|x-y\|_{\infty}=\sup _{t \in[0,1]}\{|x(t)-y(t)|\}<\delta \text { implies }|f(x)-f(y)|=|x(0)-y(0)|<\delta .
$$

Given $\epsilon>0$, let $\delta=\epsilon$. Then, we are done.
Since $A$ is compact, continuous mapping $x \mapsto x(0)$ is bounded. Therefore,

$$
\sup _{x \in A}|x(0)|<\infty .
$$

$w_{x}\left(\frac{1}{n}\right)$ is continuous in $x$ uniformly on $A$ and hence

$$
\lim _{n \rightarrow \infty} w_{x}\left(\frac{1}{n}\right)=0
$$

Suppose (i) and (ii) hold. Choose $k$ large enough so that

$$
\sup _{x \in A} w_{x}\left(\frac{1}{k}\right)=\sup _{x \in A}\left(\sup _{|s-t| \leq \frac{1}{k}}|x(s)-x(t)|\right)
$$

is finite. Since

$$
|x(t)|<|x(0)|+\sum_{i=1}^{k}\left|x\left(\frac{i t}{k}\right)-x\left(\frac{(i-1) t}{k}\right)\right|,
$$

we have

$$
\alpha=\sup _{0 \leq t \leq 1}\left(\sup _{x \in A}|x(t)|\right)<\infty .
$$

Choose $\epsilon>0$ and finite $\epsilon$-covering $H$ of $[-\alpha, \alpha]$. Choose $k$ large enough so that

$$
w_{x}\left(\frac{1}{k}\right)<\epsilon .
$$

Take $B$ to be finite set of polygonal functions on $C[0,1]$ that are linear on $\left[\frac{i-1}{k}, \frac{i}{k}\right]$ and takes the values in $H$ at end points.

If $x \in A$ and $\left|x\left(\frac{1}{k}\right)\right| \leq \alpha$ so that there exists a point $y \in B$ such that

$$
\left|x\left(\frac{i}{k}\right)-y\left(\frac{i}{k}\right)\right|<\epsilon, \quad i=1,2, \ldots, k
$$

then

$$
\left|y\left(\frac{i}{k}\right)-x(t)\right|<2 \epsilon \quad \text { for } t \in\left[\frac{i-1}{k}, \frac{i}{k}\right] .
$$

$y(t)$ is convex combination of $y\left(\frac{i}{k}\right), y\left(\frac{i-1}{k}\right)$, so it is within $2 \epsilon$ of $x(t),\|x-y\|_{\infty}<2 \epsilon, B$ is finite, $B$ is $2 \epsilon$-covering of $A$. This implies $A$ is compact.

Theorem 3.2: $\left\{P_{n}\right\}$ is tight on $C[0,1]$ if and only if

1. For each $\eta>0$, there exists $a$ and $n_{0}$ such that for $n \geq n_{0}$

$$
P_{n}(\{x:|x(0)|>a\})<\eta
$$

2. For each $\epsilon, \eta>0$, there exists $0<\delta<1$ and $n_{0}$ such that for $n \geq n_{0}$

$$
P_{n}\left(\left\{x: w_{x}(\delta) \geq \epsilon\right\}\right)<\eta
$$

Proof: Since $\left\{P_{n}\right\}$ is tight, given $\delta>0$, choose $K$ compact such that

$$
P_{n}(K)>1-\eta .
$$

Note that by the Arzela-Ascoli theorem, for large $a$

$$
K \subset\{x:|x(0)| \leq a\}
$$

and for small $\delta$

$$
K \subset\left\{x: w_{x}(\delta) \leq \epsilon\right\} .
$$

Now, $C[0,1]$ is a complete separable metric space. So for each $n, P_{n}$ is tight, and hence, we get the necessity condition. Given $\eta>0$, there exists $a$ such that

$$
P(\{x:|x(0)|>a\})<\eta
$$

and $\epsilon, \eta>0$, there exists $\delta>0$ such that

$$
P_{n}\left(\left\{x: w_{x}(\delta) \geq \epsilon\right\}\right)<\eta .
$$

This happens for $P_{n}$, where $n \leq n_{0}$ with $n_{0}$ is finite. Assume (i) and (ii) holds for all $n$. Given $\eta$, choose $a$ so that

$$
B=\{x:|x(0)| \leq a\},
$$

$a$ satisfies for all $n$

$$
P_{n}(B)>1-\eta,
$$

and choose $\delta_{k}$ such that

$$
B_{k}=\left\{x: w_{x}\left(\delta_{k}\right) \leq \frac{1}{k}\right\}
$$

with

$$
P_{n}\left(B_{k}\right)>1-\frac{\eta}{2^{k}}
$$

Let $K=\bar{A}$ where

$$
A=B \cap\left(\bigcap_{k} B_{k}\right)
$$

$K$ is compact by Arzela-Ascoli Theorem,

$$
P_{n}(K)>1-2 \eta .
$$

### 3.2 Invariance principle of sums of i.i.d. random variables

Let $X_{i}$ 's be i.i.d. with $E X_{i}=0$ and $E X_{i}^{2}=\sigma^{2}$. Define

$$
W_{t}^{n}(\omega)=\frac{1}{\sigma \sqrt{n}} S_{[n t]}(\omega)+(n t-[n t]) \frac{1}{\sigma^{2} \sqrt{n}} X_{[n t]+1} .
$$

Consider linear interpolation

$$
W_{\frac{k}{n}}=\frac{S_{k}}{\sqrt{n}} \quad \text { where } W^{n} \in C[0,1] \text { a.e. } P_{n} .
$$

Let

$$
\psi_{n t}=(n t-[n t]) \cdot \frac{X_{[n t+1]}}{\sigma \sqrt{n}} .
$$

Claim: For fixed $t$, by Chevyshev's inequality, as $n \rightarrow \infty$

$$
\psi_{n t} \rightarrow 0
$$

## Proof of Claim:

$$
\begin{align*}
P\left(\left|\psi_{n t}\right| \epsilon\right) & =P\left(\left|X_{[n t+1]}\right|>\frac{\sigma \sqrt{n} \epsilon}{(n t-[n t])}\right) \\
& \leq \frac{E\left|X_{[n t+1]}\right|^{2}}{\frac{\sigma^{2} \epsilon^{2}}{(n t-[n t])^{2}}} \\
& =\frac{(n t-[n t])^{2}}{n \epsilon^{2}} \\
& \leq \frac{1}{n \epsilon^{2}} \rightarrow 0 . \tag{3.1}
\end{align*}
$$

By CLT,

$$
\frac{S_{[n t]}}{\sigma \sqrt{[n t]}} \Rightarrow_{\mathcal{D}} N(0,1) .
$$

Since $\frac{[n t]}{n} \rightarrow t$, by CLT,

$$
\frac{S_{[n t]}}{\sigma \sqrt{n}}=\frac{S_{[n t]}}{\sigma \sqrt{[n t]}} \times \frac{\sqrt{[n t]}}{\sqrt{n}} \Rightarrow_{\mathcal{D}} \sqrt{t} Z
$$

by Slutsky's equation. Therefore,

$$
W_{t}^{n} \Rightarrow_{\mathcal{D}} \sqrt{t} Z
$$

Then,

$$
\begin{aligned}
\left(W_{s}^{n}, W_{t}^{n}-W_{s}^{n}\right)= & \frac{1}{\sigma \sqrt{n}}\left(S_{[n s]}, S_{[n t]}-S_{[n s]}\right)+\left(\psi_{n s}, \psi_{n t}-\psi_{n s}\right) \\
& \Longrightarrow{ }_{\mathcal{D}}\left(N_{1}, N_{2}\right)
\end{aligned}
$$

Since $S_{[n s]}$ and $S_{[n t]}-S_{[n s]}$ are independent, $N_{1}$ and $N_{2}$ are independent normal with variance $s$ and $t-s$. Thus,

$$
\begin{aligned}
\left(W_{s}^{n}, W_{t}^{n}\right)= & \left(W_{s}^{n},\left(W_{t}^{n}-W_{s}^{n}\right)+W_{s}^{n}\right) \\
& \Rightarrow{ }_{\mathcal{D}}\left(N_{1}, N_{1}+N_{2}\right) .
\end{aligned}
$$

The two-dimensional distributions of $\left(W_{t}^{n}\right)_{t \in[0,1]}$ converges to two dimensional distributions of Brownian Motion. We consider similarly $k$-dimensional distributions of $\left(W_{t}^{n}\right)_{t \in[0,1]}$ converge to those of Brownian Motion.

We considered two-dimensional. We can take $k$-dimensional. Similar argument shows that

$$
\left(W_{1}^{n}, \ldots, W_{k}^{n}\right) \Rightarrow_{\mathcal{D}} \text { finite dimensional distribution of Brownian motion. }
$$

Now, we have to show that $P_{n}$ is tight. Recall the Arzela-Ascoli theorem.
Theorem 3.3: $\left\{P_{n}\right\}$ is tight on $C[0,1]$ if and only if

1. For each $\eta>0$, there exists $a$ and $n_{0}$ such that for $n \geq n_{0}$

$$
P_{n}(\{x:|x(0)|>a\})>\eta .
$$

2. For each $\epsilon, \eta>0$, there exists $0<\delta<1$ and $n_{0}$ such that for $n \geq n_{0}$

$$
P_{n}\left(\left\{x: w_{x}(\delta) \geq \epsilon\right\}\right)<\eta .
$$

Theorem 3.4: Suppose $0=t_{0}<t_{1}<\cdots<t_{v}=1$ and

$$
\begin{equation*}
\min _{1<i<v}\left(t_{i}-t_{i-1}\right) \geq \delta . \tag{3.2}
\end{equation*}
$$

Then for arbitrary $x$,

$$
\begin{equation*}
w_{x}(\delta) \leq 3 \max _{1 \leq i \leq v}\left(\sup _{t_{i-1} \leq \leq \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right|\right) \tag{3.3}
\end{equation*}
$$

and for any $P$ on $C[0,1]$

$$
\begin{equation*}
P\left(x: w_{x}(\delta) \geq 3 \epsilon\right) \leq \sum_{i=1}^{v} P\left(x: \sup _{t_{i-1} \leq s \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right| \geq \epsilon\right) . \tag{3.4}
\end{equation*}
$$

Proof: Let $m$ denote the maximum in (3.3), i.e.

$$
m=\max _{1 \leq i \leq v}\left(\sup _{t_{i-1} \leq s \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right|\right)
$$

If $s, t$ lie in $I_{i}=\left[t_{i-1}, t_{i}\right]$. Then

$$
\begin{aligned}
|x(s)-x(t)| & \leq\left|x(s)-x\left(t_{i-1}\right)\right|+\left|x(t)-x\left(t_{i-1}\right)\right| \\
& \leq 2 m .
\end{aligned}
$$

Suppose $s, t$ lie in adjoining intervals $I_{i-1}$ and $I_{i}$. Then,

$$
\begin{aligned}
|x(s)-x(t)| & \leq\left|x(s)-x\left(t_{i-1}\right)\right|+\left|x\left(t_{i}\right)-x\left(t_{i-1}\right)\right|+\left|x(t)-x\left(t_{i}\right)\right| \\
& \leq 3 m .
\end{aligned}
$$

Since

$$
\min _{1<i<v}\left(t_{i}-t_{i-1}\right) \geq \delta .
$$

for $s$ and $t$ to be such that $|s-t|<\delta, s$ and $t$ should lie in the same interval or adjoining intervals. Therefore,

$$
\begin{aligned}
w_{x}(\delta) & =\sup _{|s-t| \leq \delta}|x(t)-x(s)| \\
& \leq \max \left\{\sup _{s, t \in \text { same interval }}|x(t)-x(s)|, \sup _{s, t \in a d j o i n i n g ~ i n t e r v a l ~}|x(t)-x(s)|\right\} \\
& \leq 3 m .
\end{aligned}
$$

This proves (3.3). Note that if $X \geq Y$, then

$$
P(X>a) \geq P(Y>a) .
$$

Therefore,

$$
\begin{aligned}
P\left(x: w_{x}(\delta)>3 \epsilon\right) & \leq P\left(3 \max _{1 \leq i \leq v}\left(\sup _{t_{i-1} \leq s \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right|\right)>3 \epsilon\right) \\
& =P\left(x: \max _{1 \leq i \leq v}\left(\sup _{t_{i-1} \leq \leq \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right|\right)>\epsilon\right) \\
& =\sum_{i=1}^{v} P\left(x: \sup _{t_{i-1} \leq \leq \leq t_{i}}\left|x(s)-x\left(t_{i-1}\right)\right|>\epsilon\right) .
\end{aligned}
$$

This proves the theorem.
Condition (ii) of the Arzela-Ascoli theorem holds if for each $\epsilon, \eta$, there exists $\delta \in(0,1)$ and $n_{0}$ such that for all $n \geq n_{0}$

$$
\frac{1}{\delta} P_{n}\left(x: \sup _{t \leq s \leq t+\delta}|x(s)-x(t)| \geq \epsilon\right)>\eta .
$$

Now apply Theorem 3.4 with $t_{i}=i \delta$ for $i<v=[1 / \delta]$. Then by (3.4), condition (ii) of Theorem 3.3 holds.

### 3.3 Invariance principle for sums of stationary sequences

Definition 3.2: $\left\{X_{n}\right\}$ is stationary if for any $m$,

$$
\left(X_{i_{i}}, \ldots, X_{i_{k}}\right)={ }^{\mathcal{D}}\left(X_{i_{i}+m}, \ldots, X_{i_{k}+m}\right)
$$

Lemma 3.1: Suppose $\left\{X_{n}\right\}$ is stationary and $W^{n}$ is defined as above. If

$$
\lim _{\lambda \rightarrow \infty} \limsup _{n \rightarrow \infty} \lambda^{2} P\left(\max _{k \leq n}\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right)=0
$$

then, $W^{n}$ is tight.
Proof: Since $W_{0}^{n}=0$, the condition (i) of Theorem 3.3 is satisfied. Let $P_{n}$ is induced measure of $W^{n}$, i.e. consider now $P_{n}(w(\delta)>\epsilon)$, we shall show that for all $\epsilon>0$

$$
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} P_{n}\left(w^{(\delta)} \geq \epsilon\right)=0
$$

If

$$
\min \left(t_{t}-t_{i-1}\right) \geq \delta
$$

then by Theorem 3.4,

$$
P\left(w\left(W^{n}, \delta\right) \geq 3 \epsilon\right) \leq \sum_{i=1}^{v} P\left(\sup _{t_{i-1} \leq s \leq t_{i}}\left|W_{s}^{n}-W_{t}^{n}\right| \geq \epsilon\right)
$$

Take $t_{i}=\frac{m_{i}}{n}, 0=m_{0}<m_{1}<\cdots<m_{v}=n . W_{t}^{n}$ is polygonal and hence,

$$
\sup _{t_{i-1} \leq s \leq t_{i}}\left|W_{s}^{n}-W_{t}^{n}\right|=\max _{m_{i-1} \leq k \leq m_{i}} \frac{\left|S_{k}-S_{m_{i-1}}\right|}{\sigma \sqrt{n}}
$$

Therefore,

$$
\begin{aligned}
P\left(w\left(W^{n}, \delta\right) \geq 3 \epsilon\right) & \leq \sum_{i=1}^{v} P\left(\max _{m_{i-1} \leq k \leq m_{i}} \frac{\left|S_{k}-S_{m_{i-1}}\right|}{\sigma \sqrt{n}} \geq \epsilon\right) \\
& \leq \sum_{i=1}^{v} P\left(\max _{m_{i-1} \leq k \leq m_{i}}\left|S_{k}-S_{m_{i-1}}\right| \geq \sigma \sqrt{n} \epsilon\right) \\
& =\sum_{i=1}^{v} P\left(\max _{k \leq m_{i}-m_{i-1}}\left|S_{k}\right| \geq \sigma \sqrt{n} \epsilon\right) \text { (by stationarity). }
\end{aligned}
$$

This inequality holds if

$$
\frac{m_{i}}{n}-\frac{m_{i-1}}{n} \geq \delta \text { for } 1<i<v
$$

Take $m_{i}=i m$ for $0 \leq i<v$ and $m_{v}=n$. For $i<v$ choose $\delta$ such that

$$
m_{i}-m_{i-1}=m \geq n \delta .
$$

Let $m=[n \delta], v=\left[\frac{n}{m}\right]$. Then,

$$
m_{v}-m_{v-1} \leq m
$$

and

$$
v=\left[\frac{n}{m}\right] \longrightarrow \frac{1}{\delta} \text { where } \frac{1}{2 \delta}<\frac{1}{\delta}<\frac{2}{\delta}
$$

Therefore, for large $n$

$$
\begin{aligned}
P\left(w\left(W^{n}, \delta\right) \geq 3 \epsilon\right) & \leq \sum_{i=1}^{v} P\left(\max _{k \leq m_{i}-m_{i-1}}\left|S_{k}\right| \geq \sigma \sqrt{n} \epsilon\right) \\
& \leq v P\left(\max _{k \leq m}\left|S_{k}\right| \geq \sigma \sqrt{n} \epsilon\right) \\
& \leq \frac{2}{\delta} P\left(\max _{k \leq m}\left|S_{k}\right| \geq \sigma \sqrt{n} \epsilon\right) .
\end{aligned}
$$

Take $\lambda=\frac{\epsilon}{\sqrt{28}}$. Then,

$$
P\left(w\left(W^{n}, \delta\right) \geq 3 \epsilon\right) \leq \frac{4 \lambda^{2}}{\epsilon^{2}} P\left(\max _{k \leq m}\left|S_{k}\right| \geq \lambda \sigma \sqrt{n}\right) .
$$

By the condition of the Lemma, given $\epsilon, \eta>0$, there exists $\lambda>0$ such that

$$
\frac{4 \lambda^{2}}{\epsilon^{2}} \limsup _{n \rightarrow \infty} P\left(\max _{k \leq n}\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right)<\eta .
$$

Now, for fixed $\lambda, \delta$, let $m \rightarrow \infty$ with $n \rightarrow \infty$.
Look at $X_{k}$ i.i.d. Then,

$$
\lim _{\lambda \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left(\max _{k \leq n}\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right)=0 .
$$

We know that

$$
P\left(\max _{u \leq m}\left|S_{u}\right|>\alpha\right) \leq 3 \max _{u \leq m} P\left(\left|S_{u}\right|>\frac{\alpha}{3}\right) .
$$

To show

$$
\lim _{\lambda \rightarrow \infty} \limsup _{n \rightarrow \infty} \lambda^{2} P\left(\max _{k \leq n}\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right)=0
$$

we assume that $X_{i}$ i.i.d. normal, and hence, $S_{k} / \sqrt{k}$ is asymptotically normal, $N$. Since we know

$$
P(|N|>\lambda) \leq \frac{E N^{4}}{\lambda^{4}}=\frac{3 \sigma^{4}}{\lambda^{4}}
$$

we have for $k \leq n,\left(\frac{\sqrt{n}}{\sqrt{k}}>1\right)$

$$
\begin{aligned}
P\left(\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right) & =P(\sqrt{k}|N|>\lambda \sigma \sqrt{n}) \\
& \leq \frac{3}{\lambda^{4} T^{4}}
\end{aligned}
$$

$k_{\lambda}$ is large and $k_{\lambda} \leq k \leq n$. Then,

$$
\begin{aligned}
P\left(\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right) & \leq P\left(\left|S_{k}\right|>\lambda \sigma \sqrt{k}\right) \\
& \leq \frac{3}{\lambda^{4}}
\end{aligned}
$$

Also,

$$
\begin{aligned}
P\left(\left|S_{k}\right|>\lambda \sigma \sqrt{n}\right) & \leq \frac{E\left|S_{k}\right|^{2} / \sigma^{2}}{\lambda^{2} n} \\
& \leq \frac{k_{\lambda}}{\lambda^{2} n}
\end{aligned}
$$

and hence, we get the above convergence $(A)$ is true.

### 3.4 Weak convergence on the Skorokhod space

### 3.4.1 The space $D[0,1]$

Let

$$
x:[0,1] \rightarrow R
$$

be right-continuous with left limit such that

1. for $0 \leq t<1$

$$
\lim _{s \searrow t} x(s)=x(t+)=x(t)
$$

2. for $0<t \leq 1$

$$
\lim _{s>t} x(s)=x(t-) .
$$

We say that $x(t)$ has discontinuity of the first kind at t if left and right limit exist.
For $x \in D$ and $T \subset[0,1]$,

$$
w_{x}(T)=w(x, T)=\sup _{s, t \in T}|x(s)-x(t)| .
$$

We define the modulus of continuity

$$
w_{x}(\delta)=\sup _{0 \leq t \leq 1-\delta} w_{x}([t, t+\delta))
$$

Lemma 3.2 (D1): For each $x \in D$ and $\epsilon>0$, there exist points $0=t_{0}<t_{1}<\cdots$ $<t_{v}=1$ and $w_{x}\left(\left[t_{i-1}, t_{i}\right)\right)<\epsilon$.

Proof: Call $t_{i}$ 's above $\delta$-sparse. If $\min _{i}\left\{\left(t_{i}-t_{i-1}\right)\right\} \geq \delta$, define for $0<\delta<1$

$$
w_{x}^{\prime}(\delta)=w^{\prime}(x, \delta)=\inf _{\left\{t_{i}\right\}} \max _{1 \leq i \leq v} w_{x}\left(\left[t_{i-1}, t_{i}\right)\right) .
$$

(If we prove the above lemma, we get $x \in D, \lim _{\delta \rightarrow 0} w_{x}^{\prime}(\delta)=0$.)
If $\delta<\frac{1}{2}$, we can split $[0,1)$ into subintervals $\left[t_{i-1}, t_{i}\right)$ such that

$$
\delta<\left(t_{i}-t_{i-1}\right) \leq 2 \delta
$$

and hence,

$$
w_{x}^{\prime}(\delta) \leq w_{x}(2 \delta) .
$$

Let us define jump function

$$
j(x)=\sup _{0 \leq t \leq 1}|x(t)-x(t-)| .
$$

We shall prove that

$$
w_{x}(\delta) \leq 2 w_{x}^{\prime}(\delta)+j(x) .
$$

Choose $\delta$-sparse sequence $\left\{t_{i}\right\}$ such that

$$
w_{x}\left(\left[t_{i-1}, t_{i}\right)\right)<w_{x}^{\prime}(\delta)+\epsilon .
$$

We can do this from the definition

$$
w_{x}^{\prime}(\delta)=w^{\prime}(x, \delta)=\inf _{\left\{t_{i}\right\}} \max _{1 \leq i \leq v} w_{x}\left(\left[t_{i-1}, t_{i}\right)\right) .
$$

If $|s-t|<\delta$, then $s, t \in\left[t_{i-1}, t_{i}\right)$ or belongs to adjoining intervals. Then,

$$
|x(s)-x(t)| \begin{cases}w_{x}^{\prime}(\delta)+\epsilon, & \text { if } s, t \text { belong to the same interval; } \\ 2 w_{x}^{\prime}(\delta)+\epsilon+j(x), & \text { if } s, t \text { belong to adjoining intervals. }\end{cases}
$$

If $x$ is continuous, $j(x)=0$ and hence,

$$
w_{x}(\delta) \leq 2 w_{x}^{\prime}(\delta) .
$$

### 3.4.2 Skorokhod topology

Let $\Lambda$ be the class of strictly increasing functions on $[0,1]$ and $\lambda(0)=0, \lambda(1)=1$. Define

$$
d(x, y)=\inf \left\{\epsilon: \exists \lambda \in \Lambda \text { such that } \sup _{t}|\lambda(t)-t|<\epsilon \text { and } \sup _{t}|x(\lambda(t))-y(t)|<\epsilon\right\} .
$$

$d(x, y)=0$ implies there exists $\lambda_{n} \in \Lambda$ such that $\lambda_{n}(t) \rightarrow t$ uniformly and $x\left(\lambda_{n}(t)\right) \rightarrow$ $y(t)$ uniformly. Therefore, with

$$
\begin{aligned}
\|\lambda-I\| & =\sup _{t \in[0,1]}|\lambda(t)-t| \\
\|x-y \circ \lambda\| & =\sup _{t \in[0,1]}|x(t)-y(\lambda(t))| \\
d(x, y) & =\inf _{\lambda}(\|\lambda-I\| \vee\|x-y \circ \lambda\|) .
\end{aligned}
$$

If $\lambda(t)=t$, then

1. $d(x, y)=\sup |x(t)-y(t)|<\infty$ since we showed $|x(s)-x(t)| \leq w_{x}^{\prime}(\delta)<\infty$.
2. $d(x, y)=d(y, x)$.
3. $d(x, y)=0$ only if $x(t)=y(t)$ or $x(t)=y(t-)$.

If $\lambda_{1}, \lambda_{2} \in \Lambda$ and $\lambda_{1} \circ \lambda_{2} \in \Lambda$

$$
\left\|\lambda_{1} \circ \lambda_{2}-I\right\| \leq\left\|\lambda_{1}-I\right\|+\left\|\lambda_{2}-I\right\| .
$$

If $\lambda_{1}, \lambda_{2} \in \Lambda$, then the following holds:

1. $\lambda_{1} \circ \lambda_{2} \in \Lambda$.
2. $\left\|\lambda_{1} \circ \lambda_{2}-I\right\| \leq\left\|\lambda_{1}-I\right\|+\left\|\lambda_{2}-I\right\|$.
3. $\left\|x-z \circ\left(\lambda_{1} \circ \lambda_{2}\right)\right\| \leq\left\|x-y \circ \lambda_{2}\right\|+\left\|y-z \circ \lambda_{1}\right\|$.
4. $d(x, z) \leq d(x, y)+d(y, z)$.

Therefore, Skorokhod topology is given by $d$.
( $\mathrm{D}[0,1], \mathrm{d}$ ) is not complete. To see this, choose $x_{n}=1_{\left[0, \frac{1}{\left.2^{n}\right]}\right.}(t)$ and $\lambda_{n}$ be such that $\lambda_{n}\left(1 / 2^{n}\right)=1 / 2^{n+1}$ and linear in $\left[0,1 / 2^{n}\right]$ and $\left[1 / 2^{n+1}, 1\right]$ then $\left\|x_{n+1} 0 \lambda_{n}-x_{n}\right\|=0$ and \| $\lambda_{n}-I \|=\frac{1}{2^{n+1}}$. Meanwhile, $\lambda_{n}\left(1 / 2^{n}\right) \neq 1 / 2^{n+1}$, then $\left\|x_{n}-x_{n+1}\right\|=1$ and $d\left(x_{n}, x_{n+1}\right)=$ $1 / 2^{n+1}$, i.e. $x_{n}$ is d-Cauchy and $d\left(x_{n}, 0\right)=1$.

Choose $\lambda \in \Lambda$ near identity. Then for $t, s$ close, $\frac{\lambda(t)-\lambda(s)}{t-s}$ is close to 1 . Therefore,

$$
\|\lambda\|^{0}=\sup _{s<t}\left|\log \frac{\lambda(t)-\lambda(s)}{t-s}\right| \in(0, \infty) .
$$

### 3.5 Metric of $D[0,1]$ to make it complete

Let $\lambda \in \Lambda(\lambda$ is non-decreasing, $\lambda(0)=0$, and $\lambda(1)=1)$. Recall

$$
\|\lambda\|^{0}=\sup _{s<t}\left|\log \frac{\lambda(t)-\lambda(s)}{t-s}\right| \in(0, \infty) .
$$

Consider $d^{0}$

$$
\begin{aligned}
d^{0}(x, y) & =\inf \left\{\epsilon>0: \exists \lambda \in \Lambda \text { with }\|\lambda\|^{0}<\epsilon \text { and } \sup _{t}|x(t)-y(\lambda(t))|<\epsilon\right\} \\
& =\inf _{\lambda \in \Lambda}\left\{\|\lambda\|^{0} \vee\|x-y \circ \lambda\|\right\} .
\end{aligned}
$$

since, for $u>0$,

$$
|u-1| \leq e^{|\log u|}-1
$$

we have

$$
\begin{aligned}
\sup _{0 \leq t \leq 1}|\lambda(t)-t| & =\sup _{0 \leq t \leq 1} t\left|\frac{\lambda(t)-\lambda(0)}{t-0}-1\right| \\
& =e^{\|\lambda\|^{0}}-1 .
\end{aligned}
$$

For any $v, v \leq e^{v}-1$, and hence

$$
d(x, y) \leq e^{d^{0}(x, y)}-1 .
$$

Thus, $d^{0}\left(x_{n}, y\right) \rightarrow 0$ implies $d\left(x_{n}, y\right) \rightarrow 0$.
Lemma 3.3 (D2): If $x, y \in D[0,1]$ and $d(x, y)<\delta^{2}$, then $d^{0}(x, y) \leq 4 \delta+w_{x}^{\prime}(\delta)$.

Proof: Take $\epsilon<\delta$ and $\left\{t_{i}\right\} \delta$-sparse with

$$
w_{x}\left(\left[t_{i-1}, t_{i}\right)\right)<w_{x}^{\prime}(\delta)+\epsilon \quad \forall i .
$$

We can do this from definition of $w_{x}^{\prime}(\delta)$. Choose $\mu \in \Lambda$ such that

$$
\begin{equation*}
\sup _{t}|x(t)-y(\mu(t))|=\sup _{t}\left|x\left(\mu^{-1}(t)\right)-y(t)\right|<\delta^{2} \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{t}|\mu(t)-t|<\delta^{2} \tag{3.6}
\end{equation*}
$$

This follows from $d(x, y)<\delta^{2}$. Take $\lambda$ to agree with $\mu$ at points $t_{i}$ and linear between. $\mu^{-1} \circ \lambda$ fixes $t_{i}$ and is increasing in $t$. Also, $\left(\mu^{-1} \circ \lambda\right)(t)$ lies in the same interval $\left[t_{i-1}, t_{i}\right)$. Thus, from (3.5) and (3.6),

$$
\begin{aligned}
|x(t)-y(\lambda(t))| & \leq\left|x(t)-x\left(\left(\mu^{-1} \circ \lambda\right)(t)\right)\right|+\left|x\left(\left(\mu^{-1} \circ \lambda\right)(t)\right)-y(\lambda(t))\right| \\
& =w_{x}^{\prime}(\delta)+\epsilon+\delta^{2} .
\end{aligned}
$$

$\delta<\frac{1}{2}<4 \delta+w_{x}^{\prime}(\delta) . \lambda$ agrees with $\mu$ at $t_{i}$ 's. Then by (3.5), (3.6), and $\left(t_{i}-t_{i-1}\right)>\delta$ ( $\delta$-sparse),

$$
\begin{aligned}
\left|\left(\lambda\left(t_{i}\right)-\lambda\left(t_{i-1}\right)\right)-\left(t_{i}-t_{i-1}\right)\right| & <2 \delta^{2} \\
& <2 \delta\left(t_{i}-t_{i-1}\right)
\end{aligned}
$$

and

$$
|(\lambda(t)-\lambda(s))-(t-s)| \leq 2 \delta|t-s|
$$

for $t, s \in\left[t_{i-1}, t_{i}\right)$ by polygonal property. Now, we take care of adjoining interval. For $u_{1}, u_{2}, u_{3}$
$\left|\left(\lambda\left(u_{3}\right)-\lambda\left(u_{1}\right)\right)-\left(u_{3}-u_{1}\right)\right| \leq\left|\left(\lambda\left(u_{3}\right)-\lambda\left(u_{2}\right)\right)-\left(u_{3}-u_{2}\right)\right|+\left|\left(\lambda\left(u_{2}\right)-\lambda\left(u_{1}\right)\right)-\left(u_{2}-u_{1}\right)\right|$.
If $t$ and $s$ are in adjoining intervals, we get the same bound. Since for $u<\frac{1}{2}$

$$
|\log (1 \pm u)| \leq 2 u
$$

we have

$$
\log (1-2 \delta) \leq \log \frac{\lambda(t)-\lambda(s)}{t-s} \leq \log (1+2 \delta)
$$

Therefore,

$$
\|\lambda\|^{0}=\sup _{s<t}\left|\log \frac{\lambda(t)-\lambda(s)}{t-s}\right|<4 \delta
$$

and hence, $d^{0}$ and $d$ are equivalent. Now, we shall show that $D^{0}$ is separable and is complete.

Consider $\sigma=\left\{s_{u}\right\}$ with $0=s_{0}<\cdots<s_{k}=1$ and define $A_{\sigma}: D \rightarrow D$ by

$$
\left(A_{\sigma} x\right)(t)=x\left(s_{u-1}\right)
$$

for $t \in\left[s_{u-1}, s_{u}\right)$ with $1 \leq u \leq k$ with $\left(A_{\sigma} x\right)\left(s_{k}\right)=x(1)$.
Lemma 3.4 (D3): If $\max \left(s_{u}-s_{u-1}\right) \leq \delta$, then

$$
d\left(A_{\sigma} x, x\right) \leq \delta \vee w_{x}^{\prime}(\delta)
$$

Proof: Let $A_{\sigma} x \equiv \hat{x}$. Let $\zeta(t)=s_{u-1}$ if $t \in\left[s_{u-1}, s_{u}\right)$ with $\zeta(1)=s_{k}=1$. Then, $\hat{x}(t)=$ $x(\zeta(t))$. Given $\epsilon>0$, find $\delta$-sparse set $\left\{t_{i}\right\}$ such that

$$
w_{x}\left(\left[t_{i-1}, t_{i}\right)\right)<w_{x}^{\prime}(\delta)+\epsilon
$$

for all i. Let $\lambda\left(t_{i}\right)$ be defined by

1. $\lambda\left(t_{0}\right)=s_{0}$.
2. $\lambda\left(t_{i}\right)=s_{v}$ if $t_{i} \in\left[s_{v-1}, s_{v}\right)$, where

$$
t_{i}-t_{i-1}>\delta \geq s_{v}-s_{v-1}
$$

Then, $\lambda\left(t_{i}\right)$ is increasing. Now, extend it to $\lambda \in \Lambda$ by linear interpolation.

## Claim:

$$
\begin{aligned}
\| \hat{x}(t)-x\left(\lambda^{-1}(t)\right) & =\left|x(\zeta(t))-x\left(\lambda^{-1}(t)\right)\right| \\
& <w_{x}^{\prime}(\delta)+\epsilon .
\end{aligned}
$$

Clearly, if $t=0$, or $t=1$, it is true. Let us look at $0<t<1$. First we observe that $\zeta(t), \lambda^{-1}(t)$ lie in the same interval $\left[t_{i-1}, t_{i}\right)$.(We will prove it.) This follows if we shows

$$
t_{j} \leq \zeta(t) \text { iff } t_{j} \leq \lambda^{-1}(t)
$$

or equivalently,

$$
t_{j}>\zeta(t) \text { iff } t_{j}>\lambda^{-1}(t)
$$

This is true for $t_{j}=0$. Suppose $t_{j} \in\left(s_{v-1}, s_{v}\right]$ and $\zeta(t)=s_{i}$ for some i. By definition, $\zeta(t)$ $t \leq \zeta(t)$ is equivalent to $s_{v} \leq t$. Since $t_{j} \in\left(s_{v-1}, s_{v}\right], \lambda\left(t_{j}\right)=s_{v}$. This completes the proof.

### 3.6 Separability of the Skorokhod space

$d^{0}$-convergence is stronger than $d$-convergence.
Theorem 3.5 (D1): The space $(D, d)$ is separable, and hence, so is $\left(D, d^{0}\right)$.
Proof: Let $B_{k}$ be the set of functions taking constant rational value on $\left[\frac{u-1}{k}, \frac{u}{k}\right]$ and taking rational value at 1 . Then, $B=\cup_{k} B_{k}$ is countable. Given $x \in D, \epsilon>0$, choose $k$ such that $\frac{1}{k}<\epsilon$ and $w_{x}\left(\frac{1}{k}\right)$. Apply Lemma D3 with $\sigma=\left\{\frac{u}{k}\right\}$. Note that $A_{\sigma} x$ has finite many values and

$$
d\left(x, A_{\sigma} x\right)<\epsilon .
$$

Since $A_{\sigma} x$ has finitely many real values, we can find $y \in B$ such that given $d(x, y)<\epsilon$,

$$
d\left(A_{\sigma} x, y\right)<\epsilon .
$$

Now, we shall prove the completeness.
Proof: We take $d^{0}$-Cauchy sequence. Then it contains a $d^{0}$-convergent subsequence. If $\left\{x_{k}\right\}$ is Cauchy, then there exists $\left\{y_{n}\right\}=\left\{x_{k_{n}}\right\}$ such that

$$
d^{0}\left(y_{n}, y_{n+1}\right)<\frac{1}{2^{n}} .
$$

There exists $\mu_{n} \in \Lambda$ such that

1. $\left\|\mu_{n}\right\|^{0}<\frac{1}{2^{n}}$.
2. 

$$
\begin{aligned}
\sup _{t}\left|y_{n}(t)-y_{n+1}\left(\mu_{n}(t)\right)\right| & =\sup _{t}\left|y_{n}\left(\mu_{n}^{-1}(t)\right)-y_{n+1}(t)\right| \\
& <\frac{1}{2^{n}}
\end{aligned}
$$

We have to find $y \in D$ and $\lambda_{n} \in \Lambda$ such that

$$
\left\|\lambda_{n}\right\|^{0} \rightarrow 0
$$

and

$$
y_{n}\left(\lambda_{n}^{-1}(t)\right) \rightarrow y(t)
$$

uniformly.

Heuristic (not a proof): Suppose $y_{n}\left(\lambda_{n}^{-1}(t)\right) \rightarrow y(t)$. Then, by (2), $y_{n}\left(\mu_{n}^{-1}\left(\lambda_{n+1}^{-1}(t)\right)\right)$ is within $\frac{1}{2^{n}}$ of $y_{n+1}\left(\lambda_{n+1}^{-1}(t)\right)$. Thus, $y_{n}\left(\lambda_{n}^{-1}(t)\right) \rightarrow y(t)$ uniformly.

Find $\lambda_{n}$ such that

$$
y_{n}\left(\mu_{n}^{-1}\left(\lambda_{n+1}^{-1}(t)\right)\right)=y_{n}\left(\lambda_{n}^{-1}(t)\right),
$$

i.e.

$$
\mu_{n}^{-1} \circ \lambda_{n+1}^{-1}=\lambda_{n}^{-1} .
$$

Thus,

$$
\begin{aligned}
\lambda_{n} & =\lambda_{n+1} \mu_{n} \\
& =\lambda_{n+2} \mu_{n+1} \mu_{n} \\
& \vdots \\
& =\cdots \mu_{n+2} \mu_{n+1} \mu_{n} .
\end{aligned}
$$

Proof: Since

$$
e^{u}-1 \leq 2 u
$$

for $0 \leq u \leq \frac{1}{2}$, we have

$$
\sup _{t}|\lambda(t)-t| \leq e^{\|\lambda\|^{0}}-1
$$

Therefore,

$$
\begin{aligned}
\sup _{t}\left|\left(\mu_{n+m+1} \mu_{n+m} \cdots \mu_{n}\right)(t)-\left(\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right)(t)\right| & \leq \sup _{s}\left|\mu_{n+m+1}(s)-s\right| \\
& \leq 2\left\|\mu_{n+m+1}\right\|^{0} \\
& =\frac{1}{2^{n+m}} .
\end{aligned}
$$

For fixed $n$,

$$
\left(\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right)(t)
$$

converges uniformly in $t$ as $n$ goes to $\infty$. Let

$$
\lambda_{n}(t)=\lim _{m \rightarrow \infty}\left(\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right)(t) .
$$

Then, $\lambda_{n}$ is continuous and non-decreasing with $\lambda_{n}(0)=0$ and $\lambda_{n}(1)=1$. We have to prove $\left\|\lambda_{n}\right\|^{0}$ is finite. Then, $\lambda_{n}$ is strictly increasing.

$$
\begin{aligned}
& \left|\log \frac{\left(\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right)(t)-\left(\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right)(s)}{t-s}\right| \\
& \leq\left\|\mu_{n+m} \mu_{n+m-1} \cdots \mu_{n}\right\|^{0} \\
& \left(\text { since } \lambda_{n} \in \Lambda,\left\|\lambda_{n}\right\|^{0}<\infty\right) \\
& \leq\left\|\mu_{n+m}\right\|^{0}+\cdots+\left\|\mu_{n}\right\|^{0} \\
& \left(\text { since }\left\|\lambda_{1} \lambda_{2}\right\|^{0} \leq\left\|\lambda_{1}\right\|^{0}+\left\|\lambda_{2}\right\|^{0}\right) \\
& <\frac{1}{2^{n-1}} .
\end{aligned}
$$

Let $m \rightarrow \infty$. Then, $\left\|\lambda_{n}\right\|^{0}<\frac{1}{2^{n-1}}$ is finite, and hence, $\lambda_{n}$ is strictly increasing. Now, by (2),

$$
\begin{aligned}
\sup _{t}\left|y_{n}\left(\lambda_{n}^{-1}(t)\right)-y_{n}\left(\lambda_{n+1}^{-1}(t)\right)\right| & \leq \sup _{s}\left|y_{n}(s)-y_{n+1}\left(\mu_{n}(s)\right)\right| \\
& <\frac{1}{2^{n}}
\end{aligned}
$$

Therefore, $\left\{y_{n}\left(\lambda_{n}^{-1}(t)\right)\right\}$ is Cauchy under supnorm and

$$
y_{n}\left(\lambda_{n}^{-1}(t)\right) \rightarrow y(t) \in D
$$

and hence converges in $d^{0}$.

### 3.7 Tightness in the Skorokhod space

We turn now to the problem of characterizing compact sets in $D$. We will prove an analogue of the Arzelà-Ascoli theorem.

Theorem 3.6: A necessary and sufficient condition for a set $A$ to be relatively compact in the Skorohod topology is that

$$
\begin{equation*}
\sup _{x \in A}\|x\|<\infty \tag{3.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \sup _{x \in A} w_{x}^{\prime}(\delta)=0 \tag{3.8}
\end{equation*}
$$

Proof of sufficiency: Let

$$
\alpha=\sup _{x \in A}\|x\| .
$$

Given $\epsilon>0$, choose a finite $\epsilon$-net $H$ in $[-\alpha, \alpha]$ and choose $\delta$ so that $\delta<\epsilon$ and $w_{x}^{\prime}(\delta)<\epsilon$ for all $x$ in $A$. Apply Lemma 3.4 for any $\sigma=\left\{s_{u}\right\}$ satisfying $\max \left(s_{u}-s_{u-1}\right)<\delta: x \in A$ implies $d\left(x, A_{\sigma} x\right)<2 \epsilon$. Take $B$ to be the finite set of $y$ that assume on each $\left[s_{u-1}, s_{u}\right)$ a constant value from $H$ and satisfy $y(1) \in H$. Since $B$ contains a $y$ for which $d\left(x, A_{\sigma} x\right)$, it is a finite $2 \epsilon$-net for $A$ in the sense of $d$. Thus, $A$ is totally bounded in the sense of $d$. However, we must show that $A$ is totally bounded in the sense of $d^{0}$, since this is the metric under which $D$ is complete. Given (a new) $\epsilon$, choose a new $\delta$ so that $0<\delta \leq 1 / 2$ and so that $4 \delta+w_{x}^{\prime}(\delta)<\epsilon$ holds for all $x$ in $A$. We have already seen that $A$ is $d$-totally bounded, and thus, there exists a finite set $B^{\prime}$ that is a $\delta^{2}$-net for $A$ in the sense of $d$. However, by Lemma 2, $B^{\prime}$ is an $\epsilon$-net for $A$ in the sense of $d^{0}$.

The proof of necessity requires a Lemma 3.3 and a definition.
Definition 3.3: In any metric space, $f$ is upper semi-continuous at $x$, if for all $\epsilon>0$, there exists $\delta>0$ such that

$$
\rho(x, y)<\delta \Rightarrow f(y)<f(x)+\epsilon .
$$

Lemma 3.5: For fixed $\delta, w^{\prime}(x, \delta)$ is upper-semicontinuous in $x$.
Proof: Let $x, \delta$, and $\epsilon$ be given. Let $\left\{t_{i}\right\}$ be a $\delta$-spars set such that $w_{x}\left[t_{i-1}, t_{i}\right)<w_{x}^{\prime}(\delta)+\epsilon$ for each $i$. Now choose $\eta$ small enough that $\delta+2 \eta<\min \left(t_{i}-t_{i-1}\right)$ and $\eta<\epsilon$. Suppose that $d(x, y)<\eta$. Then, for some $\lambda$ in $\Lambda$, we have

$$
\begin{aligned}
& \sup _{t}|y(t)-x(\lambda t)|<\eta, \\
& \sup _{t}\left|\lambda^{-1} t-t\right|<\eta
\end{aligned}
$$

Let $s_{i}=\lambda^{-1} t_{i}$. Then $s_{i}-s_{i-1}>t_{i}-t_{i-1}-2 \eta>\delta$. Moreover, if $s$ and $t$ both lies in $\left[s_{i-1}, s_{i}\right)$, then $\lambda s$ and $\lambda t$ both lie in $\left[t_{i-1}, t_{i}\right.$ ), and hence $|y(s)-y(t)|<|x(\lambda s)-x(\lambda t)|+2 \eta \leq$ $w_{x}^{\prime}(\delta)+\epsilon+2 \eta$. Thus, $d(x, y)<\eta$ implies $w_{y}^{\prime}(\delta)<w_{x}^{\prime}(\delta)+3 \epsilon$.

Definition 3.4 ( $d$-bounded): $A$ is $d$-bounded if diameter is bounded, i.e.

$$
\operatorname{diameter}(A)=\sup _{x, y \in A} d(x, y)<\infty .
$$

Proof of necessity in Theorem 3.6: If $A^{-}$is compact, then it is $d$-bounded, and since $\sup _{t}|x(t)|$ is the $d$-distance from $x$ to the 0 -function, (3.7) follows. By Lemma 3.1, $w^{\prime}(x, \delta)$ goes to 0 with $\delta$ for each $x$. However, since $w^{\prime}(\cdot, \delta)$ is upper-semicontinuous the convergence is uniform on compact sets.

Theorem 3.6, which characterizes compactness in $D$, gives the following result. Let $\left\{P_{n}\right\}$ be a sequence of probability measure on $(D, \mathcal{D})$.

Theorem 3.7: The sequence $\left\{P_{n}\right\}$ is tight if and only if these two conditions hold: We have

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \lim _{n} \sup _{n} P_{n}(\{x:\|x\| \geq a\})=0 \tag{3.9}
\end{equation*}
$$

(ii) for each $\epsilon$,

$$
\begin{equation*}
\lim _{\delta \rightarrow \infty} \lim _{n} \sup _{n} P_{n}\left(\left\{x: w_{x}^{\prime}(\delta) \geq \epsilon\right\}\right)=0 . \tag{3.10}
\end{equation*}
$$

Proof: Conditions (i) and (ii) here are exactly conditions (i) and (ii) of Azela-Ascoli theorem with $\|x\|$ in place of $|x(0)|$ and $w^{\prime}$ in place of $w$. Since $D$ is separable and complete, a single probability measure on $D$ is tight, and so the previous proof goes through.

### 3.8 The space $D[0, \infty)$

Here we extend the Skorohod theory to the space $D_{\infty}=D[0, \infty)$ of cadlag functions on $[0, \infty)$, a space more natural than $D=D[0,1]$ for certain problems.

In addition to $D_{\infty}$, consider for each $t>0$ the space $D_{t}=D[0, t]$ of cadlag functions on $[0, t]$. All the definitions for $D_{1}$ have obvious analogues for $D_{t}: \sup _{s \leq t}|x(s)|$, $\Lambda_{t},\|\lambda\|_{t}^{0}, d_{t}^{0}, d_{t}$. All the theorems carry over from $D_{1}$ to $D_{t}$ in an obvious way. If $x$ is an element of $D_{\infty}$, or if $x$ is an element of $D_{u}$ and $t<u$, then $x$ can also be regarded as an element of $D_{t}$ by restricting its domain of definition. This new cadlag function will be denoted by the same symbol; it will always be clear what domain is intended.

One might try to define the Skorohod convergence $x_{n} \rightarrow x$ in $D_{\infty}$ by requiring that $d_{t}^{0}\left(x_{n}, x\right) \rightarrow 0$ for each finite, positive $t$. However, in a natural theory, $x_{n}=I_{[0,1-1 / n]}$ will converge to $x=I_{[0,1]}$ in $D_{\infty}$, while $d_{1}^{0}\left(x_{n}, x\right)=1$. The problem here is that $x$ is discontinuous at 1, and the definition must accommodate discontinuities.

Lemma 3.6: Let $x_{n}$ and $x$ be elements of $D_{u}$. If $d_{u}^{0}\left(x_{n}, x\right) \rightarrow 0$ and $m<u$, and if $x$ is continuous at $m$, then $d_{m}^{0}\left(x_{n}, x\right) \rightarrow 0$.

Proof: We can work with the metrics $d_{u}$ and $d_{m}$. By hypothesis, there are elements $\lambda_{n}$ of $\Lambda_{u}$ such that

$$
\left\|\lambda_{n}-I\right\|_{u} \rightarrow 0
$$

and

$$
\left\|x_{n}-x \lambda_{n}\right\|_{u} \rightarrow 0
$$

Given $\epsilon$, choose $\delta$ so that $|t-m| \leq 2 \delta$ implies $|x(t)-x(m)|<\epsilon / 2$. Now choose $n_{0}$ so that, if $n \geq n_{0}$ and $t \leq u$, then $\left|\lambda_{n} t-t\right|<\delta$ and $\left|x_{n}(t)-x\left(\lambda_{n} t\right)\right|<\epsilon / 2$. Then, if $n \geq n_{0}$ and $|t-m| \leq \delta$, we have $\left|\lambda_{n} t-m\right| \leq\left|\lambda_{n} t-t\right|+|t-m|<2 \delta$ and hence $\left|x_{n}(t)-x(m)\right| \leq$ $\left|x_{n}(t)-x\left(\lambda_{n} t\right)\right|+\left|x\left(\lambda_{n} t\right)-x(m)\right|<\epsilon$. Thus

$$
\begin{equation*}
\sup _{|t-m| \leq \delta}|x(t)-x(m)|<\epsilon, \quad \sup _{|t-m| \leq \delta}\left|x_{n}(t)-x(m)\right|<\epsilon, \quad \text { for } n \geq n_{0} . \tag{3.11}
\end{equation*}
$$

If
(i) $\lambda_{n} m<m$, let $p_{n}=m-\frac{1}{n}$;
(ii) $\lambda_{n} m>m$, let $p_{n}=\lambda^{-1}\left(m-\frac{1}{n}\right)$;
(iii) $\lambda_{n} m=m$, let $p_{n}=m$.

Then,
(i) $\left|p_{n}-m\right|=\frac{1}{n}$;
(ii) $\left|p_{n}-m\right|=\leq\left|\lambda_{n}^{-1}\left(m-n^{-1}\right)-\left(m-n^{-1}\right)\right|+\frac{1}{n}$;
(iii) $\left|p_{n}-m\right|=m$.

Therefore, $p_{n} \rightarrow m$. Since

$$
\left|\lambda_{n} p_{n}-m\right| \leq\left|\lambda_{n} p_{n}-p_{n}\right|+\left|p_{n}-m\right|,
$$

we also have $\lambda_{n} p_{n} \rightarrow m$. Define $\mu_{n} \in \Lambda_{n}$ so that $\mu_{n} t=\lambda_{n} t$ on [ $0, p_{n}$ ] and $\mu_{n} m=m$, and interpolate linearly on $\left[p_{n}, m\right.$ ]. Since $\mu_{n} m=m$ and $\mu_{n}$ is linear over [ $p_{n}, m$ ], we have $\left|\mu_{n} t-t\right| \leq\left|\lambda_{n} p_{n}-p_{m}\right|$ there, and therefore, $\mu_{n} t \rightarrow t$ uniformly on $[0, m]$. Increase the $n_{0}$ of (3.11) so that $p_{n}>m-\delta$ and $\lambda_{n} p_{n}>m-\delta$ for $n \geq n_{0}$. If $t \leq p_{n}$, then $\left|x_{n}(t)-x\left(\mu_{n} t\right)\right|=\left|x_{n}(t)-x\left(\lambda_{n} t\right)\right| \leq\left\|x_{n}-x \lambda_{n}\right\|_{u}$. Meanwhile, if $p_{n} \leq t \leq m$ and $n \geq n_{0}$, then $m \geq t \geq p_{n}>m-\delta$ and $m \geq \mu_{n} t \geq \mu_{n} p_{n}=\lambda_{n} p_{n}>m-\delta$, and therefore, by (3.11), $\left|x_{n}(t)-x\left(\mu_{n} t\right)\right| \leq\left|x_{n}(t)-x(m)\right|+\left|x(m)-x\left(\mu_{n} t\right)\right|<2 \epsilon$. Thus, $\left|x_{n}(t)-x\left(\mu_{n} t\right)\right| \rightarrow 0$ uniformly on $[0, m]$.

The metric on $D_{\infty}$ will be defined in terms of the metrics $d_{m}^{0}(x, y)$ for integral $m$, but before restricting $x$ and $y$ to $[0, m]$, we transform them in such a way that they are continuous at $m$. Define

$$
g_{n}(t)= \begin{cases}1, & \text { if } t \leq m-1  \tag{3.12}\\ m-t, & \text { if } m-1 \leq t \leq m ; \\ 0, & t \geq m\end{cases}
$$

For $x \in D_{\infty}$, let $x^{m}$ be the element of $D_{\infty}$ defined by

$$
\begin{equation*}
x^{m}(t)=g_{m}(t) x(t), \quad t \geq 0 \tag{3.13}
\end{equation*}
$$

Now take

$$
\begin{equation*}
d_{\infty}^{0}(x, y)=\sum_{m=1}^{\infty} 2^{-m}\left(1 \wedge d_{m}^{0}\left(x^{m}, y^{m}\right)\right) \tag{3.14}
\end{equation*}
$$

If $d_{\infty}^{0}(x, y)=0$, then $d_{m}^{0}(x, y)=0$ and $x^{m}=y^{m}$ for all $m$, and this implies $x=y$. The other properties being easy to establish, $d_{\infty}^{0}$ is a metric on $D_{\infty}$; it defines the Skorohod topology there. If we replace $d_{m}^{0}$ by $d_{m}$ in (3.14), we have a metric $d_{\infty}$ equivalent to $d_{\infty}^{0}$.

Let $\Lambda_{\infty}$ be the set of continuous, increasing maps of $[0, \infty)$ onto itself.
Theorem 3.8: There is convergence $d_{\infty}^{0}\left(x_{n}, x\right) \rightarrow 0$ in $D_{\infty}$ if and only if there exist elements $\lambda_{n}$ of $\Lambda_{\infty}$ such that

$$
\begin{equation*}
\sup _{t<\infty}\left|\lambda_{n} t-t\right| \rightarrow 0 \tag{3.15}
\end{equation*}
$$

and for each $m$,

$$
\begin{equation*}
\sup _{t \leq m}\left|x_{n}\left(\lambda_{n} t\right)-x(t)\right| \rightarrow 0 . \tag{3.16}
\end{equation*}
$$

Proof: Suppose that $d_{\infty}^{0}\left(x_{n}, x\right)$ and $d_{\infty}\left(x_{n}, x\right)$, go to 0 . Then there exist elements $\lambda_{n}^{m}$ of $\Lambda_{m}$ such that

$$
\epsilon_{n}^{m}=\left\|I-\lambda_{n}^{m}\right\|_{m} \vee\left\|x_{n}^{m} \lambda_{n}^{m}-x^{m}\right\|_{m} \rightarrow 0
$$

for each $m$. Choose $l_{m}$ so that $n \geq l_{m}$ implies $\epsilon_{n}^{m}<1 / m$. Arrange that $l_{m}<l_{m+1}$, and for $l_{m} \leq n<l_{m+1}$, let $m_{n}=m$. Since $l_{m} \leq n<l_{m+1}$, we have $m_{n} \rightarrow n$ and $\epsilon_{n}^{m_{n}}<1 / m_{n}$. Define

$$
\lambda_{n} t= \begin{cases}\lambda_{n}^{m_{n}} t, & \text { if } t \leq m_{n} ; \\ t+\lambda_{n}^{m_{n}}\left(m_{n}\right)-m_{n}, & \text { if } t \geq m_{n} .\end{cases}
$$

Then, $\left|\lambda_{n} t-t\right|<1 / m_{n}$ for $t \geq m_{n}$ as well as for $t \leq m_{n}$, and therefore,

$$
\sup _{t}\left|\lambda_{n} t-t\right| \leq \frac{1}{m_{n}} \rightarrow 0 .
$$

Hence, (24). Fix $c$. If $n$ is large enough, then $c<m_{n}-1$, and so

$$
\left\|x_{n} \lambda_{n}-x\right\|_{c}=\left\|x_{n}^{m_{n}} \lambda_{n}^{m_{n}}-x^{m_{n}}\right\|_{c} \leq \frac{1}{m_{n}} \rightarrow 0
$$

which is equivalent to (3.16).

Now suppose that (3.15) and (3.16) hold. Fix $m$. First,

$$
\begin{equation*}
x_{n}^{m}\left(\lambda_{n} t\right)=g_{m}\left(\lambda_{n} t\right) x_{n}\left(\lambda_{n} t\right) \rightarrow g_{m}(t) x(t)=x^{m}(t) \tag{3.17}
\end{equation*}
$$

holds uniformly on $[0, m]$. Define $p_{n}$ and $\mu_{n}$ as in the proof of Lemma 1. As before, $\mu_{n} t \rightarrow t$ uniformly on $[0, m]$. For $t \leq p_{n},\left|x^{m}(t)-x_{n}^{m}\left(\mu_{n} t\right)\right|=\left|x^{m}(t)-x_{n}^{m}\left(\lambda_{n} t\right)\right|$, and this goes to 0 uniformly by (3.17). For the case $p_{n} \leq t \leq m$, first note that $\left|x^{m}(u)\right| \leq$ $g_{m}(u)\|x\|_{m}$ for all $u \geq 0$ and hence,

$$
\begin{equation*}
\left|x^{m}(t)-x_{n}^{m}\left(\mu_{n} t\right)\right| \leq g_{m}(t)\|x\|_{m}+g_{m}\left(\mu_{n} t\right)\left\|x_{n}\right\|_{m} . \tag{3.18}
\end{equation*}
$$

By (3.15), for large $n$, we have $\lambda_{n}(2 m)>m$ and hence $\left\|x_{n}\right\|_{m} \leq\left\|x_{n} \lambda_{n}\right\|_{2 m}$; and $\left\|x_{n} \lambda_{n}\right\|_{2 m} \rightarrow\left\|x_{n}\right\|_{2 m}$ by (3.16). This means that $\left\|x_{n}\right\|_{m}$ is bounded( $m$ is fixed). Given $\epsilon$, choose $n_{0}$ so that $n \geq n_{0}$ implies that $p_{n}$ and $\mu_{n} p_{n}$ both lies in ( $m-\epsilon, m$ ], an interval on which $g_{m}$ is bounded by $\epsilon$. If $n \geq n_{0}$ and $p_{n} \leq t \leq m$, then $t$ and $\mu_{n} t$ both lie in ( $m-\epsilon, m$ ], and it follows by (3.18) that $\left|x^{m}(t)-x_{n}^{m}\left(\mu_{n} t\right)\right| \leq \epsilon\left(\|x\|_{m}+\left\|x_{n}\right\|_{m}\right)$. Since $\left\|x_{n}\right\|_{m}$ is bounded, this implies that $\left|x^{m}(t)-x_{n}^{m}\left(\mu_{n} t\right)\right| \rightarrow 0$ holds uniformly on $\left[p_{n}, m\right]$ as well as on $\left[0, p_{n}\right]$. Therefore, $d_{m}^{0}\left(x_{n}^{m}, x^{m}\right) \rightarrow 0$ for each $m$ and hence $d_{\infty}^{0}\left(x_{n}, x\right)$ and $d_{\infty}\left(x_{n}, x\right)$ go to 0 . This completes the proof.

Theorem 3.9: There is convergence $d_{\infty}^{0}\left(x_{n}, x\right) \rightarrow 0$ in $D_{\infty}$ if and only if $d_{t}^{0}\left(x_{n}, x\right) \rightarrow 0$ for each continuity point $t$ of $x$.

Proof: If $d_{\infty}^{0}\left(x_{n}, x\right) \rightarrow 0$, then $d_{\infty}^{0}\left(x_{n}^{m}, x^{m}\right) \rightarrow 0$ for each $m$. Given a continuity point $t$ of $x$, fix an integer $m$ for which $t<m-1$. By Lemma 1 (with $t$ and $m$ in the roles of $m$ and $u$ ) and the fact that $y$ and $y^{m}$ agree on $[0, t], d_{t}^{0}\left(x_{n}, x\right)=d_{t}^{0}\left(x_{n}^{m}, x^{m}\right) \rightarrow 0$.

To prove the reverse implication, choose continuity points $t_{m}$ of $x$ in such a way that $t_{m} \uparrow \infty$. The argument now follows the first part of the proof of (3.8). Choose elements $\lambda_{n}^{m}$ of $\Lambda_{t_{m}}$ in such a way that

$$
\epsilon_{n}^{m}=\left\|\lambda_{n}^{m}-I\right\|_{t_{m}} \vee\left\|x_{n} \lambda_{n}^{m}-x\right\|_{t_{m}} \rightarrow 0
$$

for each $m$. As before, define integers $m_{n}$ in such a way that $m_{n} \rightarrow \infty$ and $\epsilon_{n}^{m_{n}}<1 / m_{n}$, and this time define $\lambda_{n} \in \Lambda_{\infty}$ by

$$
\lambda_{n} t= \begin{cases}\lambda_{n}^{m_{n}} t, & \text { if } t \leq t_{m_{n}} \\ t, & \text { if } t \geq t_{m_{n}}\end{cases}
$$

The $\left|\lambda_{n} t-t\right| \leq 1 / m_{n}$ for all $t$, and if $c<t_{m_{n}}$, then $\left\|x_{n} \lambda_{n}-x\right\|_{c}=\left\|x_{n} \lambda_{n}^{m_{n}}-x\right\|_{c} \leq$ $1 / m_{n} \rightarrow 0$. This implies that (3.15) and (3.16) hold, which in turn implies that $d_{\infty}^{0}\left(x_{n}, x\right) \rightarrow 0$. This completes the proof.

### 3.8.1 Separability and completeness

For $x \in D_{\infty}$, define $\psi_{m} x$ as $x^{m}$ restricted to $[0, m]$. Then, since $d_{m}^{0}\left(\psi_{m} x_{n}, \psi_{m} x\right)=$ $d_{m}^{0}\left(x_{n}^{m}, x^{m}\right), \psi_{m}$ is a continuous map of $D_{\infty}$ into $D_{m}$. In the product space $\Pi=D_{1} \times$ $D_{2} \times \cdots$, the metric

$$
\rho(\alpha, \beta)=\sum_{m=1}^{\infty} 2^{-m}\left(1 \wedge d_{m}^{0}\left(\alpha_{m}, \beta_{m}\right)\right)
$$

defines the product topology, that of coordinatewise convergence. Now define $\psi$ : $D_{\infty} \rightarrow \Pi$ by $\psi x=\left(\psi_{1} x, \psi_{2} x, \ldots\right)$ :

$$
\psi_{m}: D_{\infty} \rightarrow D_{m}, \quad \psi: D_{\infty} \rightarrow \Pi .
$$

Then $d_{\infty}^{0}(x, y)=\rho(\psi x, \psi y): \psi$ is an isometry of $D_{\infty}$ into $\Pi$.
Lemma 3.7: The image $\psi D_{\infty}$ is closed in $\Pi$.

Proof: Suppose that $x_{n} \in D_{\infty}$ and $\alpha \in \Pi$ and $\rho\left(\psi x_{n}, \alpha\right) \rightarrow 0$; then $d_{m}^{0}\left(x_{n}^{m}, \alpha_{m}\right) \rightarrow 0$ for each $m$. We must find an $x$ in $D_{\infty}$ such that $\alpha=\psi x$-that is, $\alpha_{m}=\psi_{m} x$ for each $m$. Let $T$ be the dense set of $t$ such that for every $m \geq t, \alpha_{m}$ is continuous at $t$. Since $d_{m}^{0}\left(x_{n}^{m}, \alpha_{m}\right) \rightarrow 0, t \in T \cap[0, m]$ implies $x_{n}^{m}(t)=g_{n}(t) x_{n}(t) \rightarrow \alpha_{m}(t)$. This means that for every $t$ in $T$, the limit $x(t)=\lim x_{n}(t)$ exists (consider an $m>t+1$, so that $g_{n}(t)=1$ ). Now $g_{m}(t) x(t)=\alpha_{m}(t)$ on $T \cap[0, m]$. It follows that $x(t)=\alpha_{m}(t)$ on $T \cap[0, m-1]$, so that $x$ can be extended to a cadlag function on each $[0, m-1]$ and then to a cadlag function on $[0, \infty]$. Now, by right continuity, $g_{m}(t) x(t)=\alpha_{m}(t)$ on $[0, m]$, or $\psi_{m} x=x^{m}=\alpha_{m}$. This completes the proof.

Theorem 3.10: The space $D_{\infty}$ is separable and complete.
Proof: Since $\Pi$ is separable and complete, so are the closed subspace $\psi D_{\infty}$ and its isometric copy $D_{\infty}$. This completes the proof.

### 3.8.2 Compactness

Theorem 3.11: Set $A$ is relatively compact in $D_{\infty}$ if and only if, for each $m, \psi_{m} A$ is relatively compact in $D_{m}$.

Proof: If $A$ is relatively compact, then $\bar{A}$ is compact and hence the continuous image $\psi_{m} \bar{A}$ is also compact. But then, $\psi_{m} A$, as a subset of $\psi_{m} \bar{A}$, is relatively compact.

Conversely, if each $\psi_{m} A$ is relatively compact, then each $\overline{\psi_{m} A}$ is compact, and therefore, $B=\overline{\psi_{1} A} \times \overline{\psi_{2} A} \times \cdots$ and $E=\psi D_{\infty} \cap B$ are both compact in $\Pi$. But $x \in A$
implies $\psi x \in \overline{\psi_{m} A}$ for each $m$, so that $\psi x \in B$. Hence $\psi A \subset E$, which implies that $\psi A$ is totally bounded and so is its isometric image $A$. This completes the proof.

For an explicit analytical characterization of relative compactness, analogous to the Arzela-Ascoli theorem, we need to adapt the $w^{\prime}(x, \delta)$ to $D_{\infty}$. For an $x \in D_{m}$, define

$$
\begin{equation*}
w_{m}^{\prime}(x, \delta)=\inf \max _{1 \leq i \leq v} w\left(x,\left[t_{i-1}, t_{i}\right)\right), \tag{3.19}
\end{equation*}
$$

where the infimum extends over all decompositions $\left[t_{i-1}, t_{i}\right), 1 \leq i \leq v$, of $[0, m)$ such that $t_{t}-t_{i-1}>\delta$ for $1 \leq i \leq v$. Note that the definition does not require $t_{v}-t_{v-1}>\delta$; Although 1 plays a special role in the theory of $D_{1}$, the integers $m$ should play no special role in the theory of $D_{\infty}$.

The exact analogue of $w^{\prime}(x, \delta)$ is (3.19), but with the infimum extending only over the decompositions satisfying $t_{t}-t_{i-1}>\delta$ for $i=v$ as well as for $i<v$. Call this $\bar{w}_{m}(x, \delta)$. By an obvious extension, a set $B$ in $D_{m}$ is relatively compact if and only if $\sup _{x}\|x\|_{m}<\infty$ and $\lim _{\delta} \sup _{x} \bar{w}(x, \delta)=0$. Suppose that $A \subset D_{\infty}$ and transform the two conditions by giving $\psi_{m} A$ the role of $B$. By (Theorem 3.11), $A$ is relatively compact if and only if, for every $m$

$$
\begin{equation*}
\sup _{x \in A}\left\|x^{m}\right\|_{m}<\infty \tag{3.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \sup _{x \in A} \bar{w}_{m}\left(x^{m}, \delta\right)=0 \tag{3.21}
\end{equation*}
$$

The next step is to show that (3.20) and (3.21) are together equivalent to the condition that, for every $m$,

$$
\begin{equation*}
\sup _{x \in A}\|x\|_{m}<\infty \tag{3.22}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \sup _{x \in A} w_{m}^{\prime}(x, \delta)=0 \tag{3.23}
\end{equation*}
$$

The equivalence of (3.20) and (3.22) follows easily because $\left\|x^{m}\right\|_{m} \leq\|x\|_{m} \leq$ $\left\|x^{m+1}\right\|_{m+1}$. Suppose (3.22) and (3.23) both hold, and let $K_{m}$ be the supremum in (3.22). If $x \in A$ and $\delta<1$, then we have $\left|x^{m}(t)\right| \leq K_{m} \delta$ for $m-\delta \leq t<m$. Given $\epsilon$, choose $\delta$ so that $K_{m} \delta<\epsilon / 4$ and the supremum in (3.23) is less than $\epsilon / 2$. If $x \in A$ and $m-\delta$ lies in the interval $\left[t_{j-1}, t_{j}\right.$ ) of the corresponding partition, replace the intervals $\left[t_{i-1}, t_{i}\right.$ ) for $i \geq j$ by the single interval $\left[t_{j-1}, m\right)$. This new partition shows that $\bar{w}_{m}(x, \delta)$. Hence, (3.21).

That (3.21) implies (3.23) is clear because $w_{m}^{\prime}(x, \delta) \leq \bar{w}_{m}(x, \delta)$ : An infimum increases if its range is reduced. This gives us the following criterion.

Theorem 3.12: A set $A \in D_{\infty}$ is relatively compact if and only if (3.22) and (3.23) hold for all $m$.

### 3.8.3 Tightness

Theorem 3.13: The sequence $\left\{P_{n}\right\}$ is tight if and only if there two conditions hold:
(i) For each $m$

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \lim _{n} \sup _{n} P_{n}\left(\left\{x:\|x\|_{m} \geq a\right\}\right)=0 \tag{3.24}
\end{equation*}
$$

(ii) For each $m$ and $\epsilon$,

$$
\begin{equation*}
\lim _{\delta} \lim _{n} \sup _{n}\left(\left\{x: w_{m}^{\prime}(x, \delta) \geq \epsilon\right\}\right)=0 . \tag{3.25}
\end{equation*}
$$

There is the corresponding corollary. Let

$$
\begin{equation*}
j_{m}(x)=\sup _{t \leq m}|x(t)-x(t-)| . \tag{3.26}
\end{equation*}
$$

Corollary 3.1: Either of the following two conditions can be substituted for (i) in Theorem 3.13:
(i') For each $t$ in a set $T$ that is dense in $[0, \infty)$,

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \lim _{n} \sup _{n}(\{x:|x(t)| \geq a\})=0 . \tag{3.27}
\end{equation*}
$$

(ii') The relation (3.27) holds for $t=0$, and for each $m$,

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \lim _{n} \sup _{n} P_{n}\left(\left\{x: j_{m}(x) \geq a\right\}\right)=0 \tag{3.28}
\end{equation*}
$$

Proof: The proof is almost the same as that for the corollary to Theorem 3.7.
Assume (ii) and (i'). Choose points $t_{i}$ such that $0=t_{0}<t_{1}<\cdots<t_{v}=m, t_{i}-$ $t_{i-1}>\delta$ for $1 \leq i \leq v-1$, and $w_{x}\left[t_{i-1}, t_{i}\right)<w_{m}^{\prime}(x, \delta)+1$ for $1 \leq i \leq v$. Choose from $T$ points $s_{j}$ such that $0=s_{0}<s_{1}<\cdots<s_{k}=m$ and $s_{j}-s_{j-1}<\delta$ for $1 \leq j \leq k$. Let $m(x)=$ $\max _{0 \leq j \leq k}\left|x\left(s_{j}\right)\right|$. If $t_{v}-t_{v-1}>\delta$, then $\|x\|_{m} \leq m(x)+w_{m}^{\prime}(x, \delta)+1$, just as before. If $t_{v}-t_{v-1} \leq \delta\left(\right.$ and $\delta<1$, so that $\left.t_{v-1}>m-1\right)$, then $\|x\|_{m-1} \leq m(x)+w_{m}^{\prime}(x, \delta)+1$. The old argument now gives (3.24), but with $\|x\|_{m}$ replaced by $\|x\|_{m-1}$, which is just as good.

In the proof that (ii) and (i') imply (i), we have ( $v-1$ ) $\delta \leq m$ instead of $v \delta<$ !. However, $v \leq m \delta^{-1}+1$, and the old argument goes through. This completes the proof.

### 3.8.4 Aldous's tightness criterion

Consider two conditions, assuming (3.24).
Condition $\mathbf{1}^{\mathbf{0}}$. For each $\epsilon, \eta, m$, there exist a $\delta_{0}$ and an $n_{0}$ such that, if $\delta \leq \delta_{0}$ and $n \geq n_{0}$, and if $\tau$ is a discrete $X^{n}$-stopping time satisfying $\tau \leq m$, then

$$
\begin{equation*}
P\left(\left|X_{\tau+\delta}^{n}-X_{\tau}^{n}\right| \geq \epsilon\right) \leq \eta . \tag{3.29}
\end{equation*}
$$

Condition $2^{\mathbf{0}}$. For each $\epsilon, \eta, m$, there exist a $\delta$ and an $n_{0}$ such that, if $n \geq n_{0}$, and if $\tau_{1}$ and $\tau_{2}$ are a discrete $X^{n}$-stopping time satisfying $0 \leq \tau_{1} \leq \tau_{2} \leq m$, then

$$
\begin{equation*}
P\left(\left|X_{\tau_{2}}^{n}-X_{\tau_{1}}^{n}\right| \geq \epsilon, \tau_{2}-\tau_{1} \leq \delta\right) \leq \eta . \tag{3.30}
\end{equation*}
$$

Theorem 3.14: Conditions $1^{0}$ and $2^{0}$ are equivalent.
Proof: Note that $\tau+\delta$ is a stopping time since

$$
\{\tau+\delta \leq t\}=\{\tau \leq t-\delta\} \in \mathcal{F}_{t}^{X_{n}}
$$

In Condition $2^{0}$, put $\tau_{2}=\tau, \tau_{1}=\tau$. Then it gives Condition $1^{0}$. For the converse, suppose that $\tau \leq m$ and choose $\delta_{0}$ so that $\delta \leq 2 \delta_{0}$ and $n \geq n_{0}$ together imply (3.29). Fix an $n \geq n_{0}$ and a $\delta \leq \delta_{0}$, and let (enlarge the probability space for $X^{n}$ ) $\theta$ be a random variable independent of $\mathcal{F}^{n}=\sigma\left(X_{n}^{s}: s \geq 0\right)$ and uniformly distributed over $J=[0,2 \delta]$. For the moment, fix an $x$ in $D_{\infty}$ and points $t_{1}$ and $t_{2}$ satisfying $0 \leq t_{1} \leq t_{2}$. Let $\mu$ be the uniform distribution over $J$, and let $I=[0, \delta], M_{i}=\left\{s \in J:\left|x\left(t_{i}+s\right)-x\left(t_{i}\right)\right|<\epsilon\right\}$, and $d=t_{2}-t_{1}$.

Suppose that

$$
\begin{equation*}
t_{2}-t_{1} \leq \delta \tag{3.31}
\end{equation*}
$$

and

$$
\begin{equation*}
\mu\left(M_{i}\right)=P\left(\theta \in M_{i}\right)>\frac{3}{4}, \quad \text { for } i=1,2 \tag{3.32}
\end{equation*}
$$

If $\mu\left(M_{2} \cap I\right) \leq \frac{1}{4}$, then $\mu\left(M_{2}\right) \leq \frac{3}{4}$, which is a contradiction. Hence, $\mu\left(M_{2} \cap I\right)>\frac{1}{4}$, and for $d(0 \leq d \leq \delta), \mu\left(\left(M_{2}+d\right) \cap J\right) \leq \mu\left(\left(M_{2} \cap I\right)+d\right)=\mu\left(\left(M_{2} \cap I\right)\right) \frac{1}{4}$. Thus $\mu\left(M_{1}\right)+\mu\left(\left(M_{2}+d\right) \cap J\right)>$ 1 , which implies $\mu\left(M_{1} \cap\left(M_{2}+d\right)\right)>0$. There is therefore an $s$ such that $s \in M_{1}$ and $s-d \in M_{2}$, from which follows

$$
\begin{equation*}
\left|x\left(t_{1}\right)-x\left(t_{2}\right)\right|<2 \epsilon \tag{3.33}
\end{equation*}
$$

Thus, (3.31) and (3.32) together implies (3.33). To put it another way, if (3.31) holds but (3.33) does not, then either $P\left(\theta \in M_{1}^{c}\right) \geq \frac{1}{4}$ or $P\left(\theta \in M_{2}^{c}\right) \geq \frac{1}{4}$. Therefore,

$$
\begin{aligned}
P\left(\left|X_{\tau_{2}}^{n}-X_{\tau_{1}}^{n}\right| \geq 2 \epsilon, \tau_{2}-\tau_{1} \leq \delta\right) & \leq \sum_{i=1}^{2} P\left[P\left(\left|X_{\tau_{i}+\theta}^{n}-X_{\tau_{i}}^{n}\right| \geq \epsilon \mid \mathcal{F}^{n}\right) \geq \frac{1}{4}\right] \\
& \leq 4 \sum_{i=1}^{2} P\left(\left|X_{\tau_{i}+\theta}^{n}-X_{\tau_{i}}^{n}\right| \geq \epsilon\right) .
\end{aligned}
$$

Since $0 \leq \theta \leq 2 \delta \leq 2 \delta_{0}$, and since $\theta$ and $\mathcal{F}^{n}$ are independent, it follows by (3.29) that the final term here is at most $8 \eta$. Therefore, Condition $1^{0}$ implies Condition $2^{0}$.

This is Aldous's theorem:
Theorem 3.15 (Aldous): If (3.24) and Condition $1^{\circ}$ hold, then $\left\{X^{n}\right\}$ is tight.

Proof: By Theorem 3.13, it is enough to prove that

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \limsup _{n} P\left(w_{m}^{\prime}\left(X^{n}, \delta\right) \geq \epsilon\right)=0 . \tag{3.34}
\end{equation*}
$$

Let $\Delta_{k}$ be the set of nonnegative dyadic rationals ${ }^{1} j / 2^{k}$ of order $k$. Define random variables $\tau_{0}^{n}, \tau_{1}^{n}, \ldots$ by $\tau_{0}^{n}=0$ and

$$
\tau_{i}^{n}=\min \left\{t \in \Delta_{k}: \tau_{i-1}^{n}<t \leq m,\left|X_{t}^{n}-X_{\tau_{i-1}^{n}}^{n}\right| \geq \epsilon\right\},
$$

with $\tau_{i}^{n}=m$ if there is no such $t$. The $\tau_{i}^{n}$ depend on $\epsilon, m$, and $k$ as well as on $i$ and $n$, although the notation does not show this. It is easy to prove by induction that the $\tau_{i}^{n}$ are all stopping times.

Because of Theorem 3.14, we can assume that condition 2 holds. For given $\epsilon, \eta, m$, choose $\delta^{\prime}$ and $n_{0}$ so that

$$
P\left(\left|X_{\tau_{i}}^{n}-X_{\tau_{i-1}}^{n}\right| \geq \epsilon, \tau_{i}^{n}-\tau_{i-1}^{n} \leq \delta^{\prime}\right) \leq \eta
$$

for $i \geq 1$ and $n \geq n_{0}$. Since $\tau_{i}^{n}<m$ implies that $\left|X_{\tau_{i}}^{n}-X_{\tau_{i-1}}^{n}\right| \geq \epsilon$, we have

$$
\begin{equation*}
P\left(\tau_{i}^{n}<m, \tau_{i}^{n}-\tau_{i-1}^{n} \leq \delta^{\prime}\right) \leq \eta, \quad i \geq 1, n \geq n_{0} . \tag{3.35}
\end{equation*}
$$

Now choose an integer $q$ such that $q \delta \geq 2 m$. There is also a $\delta$ such that

$$
\begin{equation*}
P\left(\tau_{i}^{n}<m, \tau_{i}^{n}-\tau_{i-1}^{n} \leq \delta\right) \leq \frac{\eta}{q}, \quad i \geq 1, n \geq n_{0} . \tag{3.36}
\end{equation*}
$$

1 Dyadic rational is a rational number whose denominator is a power of 2, i.e. a number of the form $a / 2 b$, where $a$ is an integer and $b$ is a natural number; for example, $1 / 2$ or $3 / 8$, but not $1 / 3$. These are precisely the numbers whose binary expansion is finite.

However,

$$
\begin{equation*}
P\left(\bigcup_{i-1}^{q}\left\{\tau_{i}^{n}<m, \tau_{i}^{n}-\tau_{i-1}^{n} \leq \delta\right\}\right) \leq \eta, \quad n \geq n_{0} . \tag{3.37}
\end{equation*}
$$

Although $\tau_{i}^{n}$ depends on $k$, (3.35) and (3.37) hold for all $k$ simultaneously. By (3.35),

$$
\begin{aligned}
E\left(\tau_{i}^{n}-\tau_{i-1}^{n} \mid \tau_{q}^{n}<m\right) & \geq \delta^{\prime} P\left(\tau_{i}^{n}-\tau_{i-1}^{n} \geq \delta^{\prime} \mid \tau_{q}^{n}<m\right) \\
& \geq \delta^{\prime}\left(1-\frac{\eta}{P\left(\tau_{q}^{n}<m\right)}\right)
\end{aligned}
$$

and therefore,

$$
\begin{aligned}
m & \geq E\left(\tau_{q}^{n} \mid \tau_{q}^{n}<m\right) \\
& =\sum_{i=1}^{q} E\left(\tau_{i}^{n}-\tau_{i-1}^{n} \mid \tau_{q}^{n}<m\right) \\
& \geq q \delta^{\prime}\left(1-\frac{\eta}{P\left(\tau_{q}^{n}<m\right)}\right) .
\end{aligned}
$$

Since $q \delta^{\prime} \geq 2 m$ by the choice of $q$, this leads to $P\left(\tau_{q}^{n}<m\right) \geq 2 \eta$. By this and (3.37),

$$
\begin{equation*}
P\left(\left\{\tau_{q}^{n}<m\right\} \cup \bigcup_{i-1}^{q}\left\{\tau_{i}^{n}<m, \tau_{i}^{n}-\tau_{i-1}^{n} \leq \delta\right\}\right) \leq 3 \eta, \quad k \geq 1, n \geq n_{0} \tag{3.38}
\end{equation*}
$$

Let $A_{n_{k}}$ be the complement of the set in (3.38). On this set, let $v$ be the first index for which $\tau_{v}^{n}=m$. Fix an $n$ beyond $n_{0}$. There are points $t_{i}^{k}\left(\tau_{i}^{n}\right)$ such that $0=t_{0}^{k}<\cdots$ $<t_{v}^{k}=m$ and $t_{i}^{k}-t_{i-1}^{k}>\delta$ for $1 \leq i<v .\left|X_{t}^{n}-X_{s}^{n}\right|<\epsilon$ if $s, t$ lie in the same $\left[t_{i-1}^{k}, t_{i}^{k}\right)$ as well as in $\Delta_{k}$. If $A_{n}=\limsup \sup _{k} A_{n_{k}}$, then $P\left(A_{n}\right) \geq 1-3 \eta$, and on $A_{n}$ there is a sequence of values of $k$ along which $v$ is constant $(v \leq q)$, and for each $i \leq v, t_{i}^{k}$ converges to some $t_{i}$. However, $0=t_{0}<\cdots<t_{v}=m, t_{i}-t_{i-1} \geq \delta$ for $i<v$, and by right continuity, $\left|X_{t}^{n}-X_{s}^{n}\right| \leq \epsilon$ if $s, t$ lie in the same $\left[t_{i-1}, t_{i}\right)$. It follows that $w^{\prime}\left(X^{n}, \delta\right) \leq \epsilon$ on a set of probability at least $1-3 \eta$ and hence (3.34).

As a corollary to the theorem we get.
Corollary 3.2: If for each $m$, the sequences $\left\{X_{0}^{n}\right\}$ and $\left\{j_{m}\left(X^{n}\right)\right\}$ are tight on the line, and if Condition $1^{\circ}$ holds, then $\left\{X^{n}\right\}$ is tight.

Skorokhod introduced the topology on $D[0, \infty)$ to study the weak convergence of processes with independent increment. Semi-martingales are generalizations of these processes. In the next chapter, we shall discuss the weak convergence of semimartingales.

# 4 Central limit theorem for semi-martingales and applications 

As stated at the end of chapter 3, Skorokhod developed the theory of weak convergence of stochastic processes with values in $D[0, T]$ to consider the limit theorems (or invariance principle) with convergence to process with independent increments. Since these are special classes of semi-martingales that have sample paths in $D[0, \infty)$, we now study the work of Liptser and Shiryaev [17] (see [15]) on weak convergence of a sequence of semi-martingales to a limit. We begin with the definition of semimartingales and their structure, including semi-martingale characteristics [15]. Based on this, we obtain conditions for the weak convergence of semi-martingale sequence to a limit. We begin with some preliminary lemmas, which will be needed in the proofs. We end the chapter by giving applications to statistics of censored data that arises in survival analysis in clinical trials.

### 4.1 Local characteristics of semi-martingale

In this chapter, we study the central limit theorem by Lipster and Shiryayev. We begin by giving some preliminaries. We consider ( $\Omega,\left\{\mathcal{F}_{t}\right\}, \mathcal{F}, P$ ) a filtered probability space, where $\mathbf{F}=\left\{\mathcal{F}_{t}, t \geq 0\right\}$ is a non-decreasing family of sub $\sigma$-field of $\mathcal{F}$, satisfying $\bigcap_{t \geq s} \mathcal{F}_{t}=\mathcal{F}_{s}$. We say that $\{\underline{X}, \mathbf{F}\}$ is a martingale if for each $t, X_{t} \in \underline{X}=\left\{X_{t}\right\} \subset$ $L_{1}\left(\Omega, \mathcal{F}_{t}, P\right)$ and $E\left(X(t) \mid \mathcal{F}_{s}\right)=X_{s}$ a.e. P. WLOG, we assume $\left\{X_{t}, t \geq 0\right\}$ is $D[0, \infty)$ valued (or a.s. it is cadlag) as we can always find a version. A martingale $X$ is said to be square-integrable if $\sup _{t} E X_{t}^{2}<\infty$. We say that $\left\{X_{t}, t \geq 0\right\}$ is locally square integrable martingale if there exists an increasing sequence $\sigma_{n}$ of $\left(\mathcal{F}_{t}\right)$-stopping times such that $0 \leq \sigma_{n}<\infty$ a.e. $\lim _{n} \sigma_{n}=\infty$, and $\left\{X\left(t \wedge \sigma_{n}\right) 1_{\left\{\sigma_{n}>0\right\}}\right\}$ is a square integrable martingale. A process ( $\underline{X}, \mathbf{F}$ ) is called a semi-martingale if it has the decomposition

$$
X_{t}=X_{0}+M_{t}+A_{t}
$$

where $\left\{M_{t}\right\}$ is local martingale, $M_{0}=0, A$ is right continuous process with $A_{0}, A_{t}, \mathcal{F}_{t^{-}}$ measurable and has sample paths of finite variation. We now state condition of $A$ to make this decomposition unique (called canonical decomposition). For this we need the following. We say that a sub $\sigma$-field of $[0, \infty) \times \Omega$ generated by sets of the form $\left\{(s, t] \times A, 0<s \leq t<\infty\right.$ and $\left.A \in \mathcal{F}_{s}\right\}$ and $\{0\} \times B\left(B \in \mathcal{F}_{0}\right)$ is the $\sigma$-field of predictable sets from now on called predictable $\sigma$-algebra $\mathcal{P}$. In the above decomposition, $A$ is $\mathcal{P}$-measurable then it is canonical.

We remark that if the jumps of the semi-martingale are bounded, then the decomposition is canonical.

We now introduce the concept of local characteristics.

Let $(\underline{X}, F)$ be a semi-martingale. Set with $\Delta X(s)$ as jump at $s$,

$$
\tilde{X}(t)=\sum_{s \leq t} \Delta X(s) 1(|\Delta X(s)| \geq \epsilon) .
$$

The $\underline{X}(t)=X(t)-\tilde{X}(t)$ is a semi-martingale with unique canonical decomposition

$$
\underline{X}(t)=X(0)+M(t)+A(t),
$$

where $(M, \mathbf{F})$ is a local martingale and $A$ is predictable process of finite variation. Thus,

$$
X(t)=X(0)+M(t)+A(t)+\tilde{X}(t)
$$

Let

$$
\mu((0, t] ; A \cap\{|x| \geq \epsilon\})=\sum 1(\Delta X(s) \in A,|\Delta X(s)|>\epsilon)
$$

and $v((0, t] ; \cdot \cap\{|x| \geq \epsilon\})$ its predictable projection. Then for each $\epsilon>0$, we can write with $v$ as a mesaure generated on $v(* x 0)$

$$
X(t)=X(0)+M^{\prime}(t)+A(t)+\int_{0}^{1} \int_{|x|>\epsilon} x v(d s, d x),
$$

where ( $M^{\prime}, \mathbf{F}$ ) is a local martingale. Now the last two terms are predictable. Thus, the semi-martingale is described by $M^{\prime}, A$, and $v$. We thus have the following.

Definition 4.1: The local characteristic of a semi-martingale $X$ is defined by the triplet ( $A, C, v$ ), where

1. $A$ is the predictable process of finite variation appearing in the above decomposition of $\underline{X}(t)$.
2. $C$ is a continuous process defined by $C_{t}$

$$
C_{t}=[\underline{X}, \underline{X}]_{t}^{c}=\left\langle M^{c}\right\rangle_{t} .
$$

3. $v$ is the predictable random measure on $R_{+} \times R$, the dual predictable projection of the measure $\mu$ associated to the jumps of $X$ given on $\left(\{0\}^{c}\right)$ by

$$
\mu(w, d t, d x)=\sum_{s>0} 1(\Delta X(s, w) \neq 0) \delta_{(s, \Delta X(s, w))}(d t, d x)
$$

with $\delta$ being the dirac delta measure.

### 4.2 Lenglart inequality

Lemma 4.1 Mcleish lemma: Let $F_{n}(t), n=1,2, \ldots$ and $F(t)$ be in $D[0, \infty)$ such that $F_{n}(t)$ is increasing in $t$ for each $n$, and $F(t)$ is a.s. continuous and for each $t>0$, there exists $t_{n} \rightarrow t$ such that $F_{n}\left(t_{n}\right) \rightarrow F(t)$. Then,

$$
\sup _{t}\left|F_{n}(t)-F(t)\right| \rightarrow 0,
$$

where sup is taken over a compact set of $[0, \infty)$.
Proof: WLOG, choose compact set $[0,1]$ : For $\epsilon>0$ choose $\left\{t_{n_{i}}, i=0,1, \ldots, k\right\}$ for fixed $k \geq 1 / \epsilon$ such that $t_{n_{i}} \rightarrow i \epsilon$ and $F_{n}\left(t_{n_{i}}\right) \rightarrow_{p} F(i \epsilon)$ as $n \rightarrow \infty$. Then

$$
\begin{aligned}
\sup _{t}\left|F_{n}(t)-F(t)\right| \leq & \sup _{i}\left|F_{n}\left(t_{n_{i+1}}\right)-F_{n}\left(t_{n_{i}}\right)\right|+\sup _{i}\left|F_{n}\left(t_{n_{i}}\right)-F\left(t_{n_{i}}\right)\right| \\
& +\sup _{i}\left|F_{n}\left(t_{n_{i+1}}\right)-F\left(t_{n_{i}}\right)\right|+\epsilon
\end{aligned}
$$

As $n \rightarrow \infty$, choose $\epsilon$ such that $|F((i+1) \epsilon)-F(i \epsilon)|$ is small.
We assume that $A_{t}$ is an increasing process for each $t$ and is $\mathcal{F}_{t}$-measurable.
Definition 4.2: An adapted positive right continuous process $X_{t}$ is said to be dominated by an increasing predictable process $A$ if for all finite stopping times $T$ we have

$$
E X_{T} \leq E A_{T}
$$

Example: Let $M_{t}^{2}$ is square martingale. Consider $X_{t}=M_{t}^{2}$. Then, we know $X_{t}-\langle M\rangle_{t}$ is a martingale, and hence, $X_{T^{-}}\langle M\rangle_{T}$ is a martingale. Thus,

$$
\begin{aligned}
& E\left(X_{T^{-}}<M>_{T}\right)=E X_{0}=0 \\
\Rightarrow & E X_{T}=E<M>_{T}
\end{aligned}
$$

Let

$$
X_{t}^{*}=\sup _{s \leq t}\left|X_{s}\right|
$$

Lemma 4.2: Let $T$ be a stopping time and $X$ be dominated by increasing process $A$ (as above). Then,

$$
P\left(X_{T}^{*} \geq c\right) \leq \frac{E\left(A_{T}\right)}{c}
$$

for any positive $c$.

Proof: Let $S=\inf \left\{s \leq T \wedge n: X_{s} \geq c\right\}$, clearly $S \leq T \wedge n$. Thus,

$$
\begin{aligned}
E A_{T} & \geq E A_{S}(\text { since } A \text { is an increasing process }) \\
& \geq E X_{S}(\text { since } X \text { is dominated by } A) \\
& \geq \int_{\left\{X_{T \wedge n}^{*}>c\right\}} X_{S} d P\left(\text { since } X_{S}>0 \text { on }\left\{X_{T \wedge n}^{*}>c\right\}\right) \\
& \geq c \cdot P\left(X_{T \wedge n}^{*}>c\right)
\end{aligned}
$$

Therefore, we let $n$ go to $\infty$, then we get

$$
E A_{T} \geq c \cdot P\left(X_{T}^{*}>c\right)
$$

Theorem 4.1 (Lenglart Inequality): If $X$ is dominated by a predictable increasing process, then for every positive $c$ and $d$

$$
P\left(X_{T}^{*}>c\right) \leq \frac{E\left(A_{T} \wedge d\right)}{c}+P\left(A_{T}>d\right) .
$$

Proof: It is enough to prove for predictable stopping time $T>0$,

$$
\begin{equation*}
P\left(X_{T_{-}}^{*} \geq c\right) \leq \frac{1}{c} E\left(A_{T_{-}} \wedge d\right)+P\left(A_{T_{-}} \geq d\right) \tag{4.1}
\end{equation*}
$$

We choose $T^{\prime}=\infty$. Then $T^{\prime}$ is predictable, $\sigma_{n}=n$ and apply to $X_{t}^{T}=X_{t \wedge T}$ for $T$ finite stopping time $X_{T-}^{T^{*}}=X_{T}^{*}$.

To prove (4.1)

$$
\begin{align*}
P\left(X_{T-}^{*} \geq c\right) & =P\left(X_{T_{-}}^{*} \geq c, A_{T_{-}}<d\right)+P\left(X_{T-}^{*} \geq c, A_{T_{-}} \geq d\right) \\
& \leq P\left(X_{T-}^{*} \geq c\right)+P\left(A_{T_{-}} \geq d\right) \tag{4.2}
\end{align*}
$$

Let $S=\inf \left\{t: A_{t} \geq d\right\}$. It is easy to show that $S$ is a stopping time. Also, $S$ is predictable.
On $\left\{\omega: A_{T-}<d\right\}, S(\omega) \geq T(\omega)$, and hence

$$
1\left(A_{t-}<d\right) X_{T-}^{*} \leq X_{(T \wedge S)-}^{*}
$$

By (4.2), we have

$$
\begin{aligned}
P\left(X_{T-}^{*} \geq c\right) & \leq P\left(X_{T-}^{*} \geq c\right)+P\left(A_{T-} \geq d\right) \\
& \leq P\left(X_{T \wedge S}^{*} \geq c\right)+P\left(A_{T-} \geq d\right)
\end{aligned}
$$

Let $\epsilon>0, \epsilon<c$ and $S_{n} \nearrow S \wedge T$. Then,

$$
\begin{aligned}
P\left(X_{(T \wedge S)-}^{*} \geq c\right) & \leq \liminf _{n} P\left(X_{S_{n}}^{*} \geq c-\epsilon\right)(\text { by Fatou’s lemma }) \\
& \leq \frac{1}{c-\epsilon} \lim _{n \rightarrow \infty} E A_{S_{n}}(\text { by Lemma 4.2 }) \\
& =\frac{1}{c-\epsilon} E A_{(S \wedge T)}(\text { by Monotone convergence theorem })
\end{aligned}
$$

Since $\epsilon$ is arbitrary,

$$
\begin{aligned}
P\left(X_{(T \wedge S)-}^{*} \geq c\right) & \leq \frac{1}{c} E A_{S_{n}} \leq \frac{1}{c} E\left(A_{(S \wedge T)-}\right) \\
& \leq \frac{1}{c} E\left(A_{(T-\wedge d)}\right)
\end{aligned}
$$

This completes the proof.
Corollary 4.1: Let $M \in \mathcal{M}_{L O C}^{2}\left(\left(\mathcal{F}_{t}\right), P\right)$ (class of locally square integrable martingale). Then,

$$
P\left(\sup _{t \leq T}\left|M_{t}\right|>a\right) \leq \frac{1}{a^{2}} E\left(<M>_{T} \wedge b\right)+P\left(<M>_{T} \geq b\right)
$$

Proof: Use $X_{t}=\left|M_{t}\right|^{2}, c=a^{2}, b=d, A_{t}=\langle M\rangle_{t}$.
Lemma 4.3: Consider $\left\{\mathcal{F}_{t}^{n}, P\right\}$. Let $\left\{M^{n}\right\}$ be locally square martingale. Assume that

$$
<M^{n}>_{t} \longrightarrow_{p} f(t)
$$

where $f$ is continuous and deterministic function (hence, $f$ will be an increasing function). Then, $\left\{P \circ\left(M^{n}\right)^{-1}\right\}$ on $D[0, \infty)$ is relatively compact on $D[0, \infty)$.

Proof: It suffices to show for constant $T<\infty$, and for any $\eta>0$ there exists $a>0$ such that

$$
\begin{equation*}
\sup _{n} P\left(\sup _{t \leq T}\left|M_{t}^{n}\right|>a\right)<\eta \tag{4.3}
\end{equation*}
$$

For each $T<\infty$ and $\eta, \epsilon>0$, there exist $n_{0}, \delta$ such that for any stopping time $\tau^{n}$ (w.r.t $\left.\left(\mathcal{F}_{t}^{n}\right), \tau^{n} \leq T, \tau^{n}+\delta<T\right)$.

$$
\begin{equation*}
\sup _{n \geq n_{0}} P\left(\sup _{0 \leq t \leq \delta}\left|M_{\tau^{n}+t}^{n}-M_{\tau_{n}}^{n}\right| \geq \epsilon\right)<\eta . \tag{4.4}
\end{equation*}
$$

Observe that by corollary to Lenglart inequality,

$$
P\left(\sup _{t \leq T}\left|M_{t}^{n}\right|>a\right) \leq \frac{1}{a^{2}} E\left(<M^{n}>_{T} \wedge b\right)+P\left(<M^{n}>_{T} \geq b\right)
$$

Let $b=f(T)+1$, then under the hypothesis, there exists $n_{1}$ such that for all $n \geq n_{1}$

$$
P\left(<M^{n}>_{T} \geq b\right)<\frac{\eta}{2} .
$$

Thus, for $n \geq n_{1}$

$$
\sup _{n} P\left(\sup _{t \leq T}\left|M_{t}^{n}\right|>a\right) \leq \frac{b}{a^{2}}+\frac{\eta}{2}+\sum_{k=1}^{n_{1}} P\left(\sup _{t \leq T}\left|M_{t}^{k}\right|>a\right) .
$$

Choose $a$ large to obtain (4.3).
We again note that $M_{\tau+t}^{n}-M_{\tau}^{n}$ is a locally square integrable martingale. Hence, by Corollary 4.1 and triangle inequality

$$
\begin{aligned}
P\left(\sup _{0 \leq t \leq \delta}\left|M_{\tau+t}^{n}-M_{t}^{n}\right| \geq \epsilon\right) \leq & \frac{1}{\epsilon^{2}} E\left(\left(\left\langle M^{n}\right\rangle_{\tau+\delta}-\left\langle M^{n}\right\rangle_{\tau}\right) \wedge b\right) \\
& +P\left(\left\langleM^{n}>_{\tau+\delta}-\left\langle M^{n}>_{\tau} \geq b\right)\right.\right. \\
\leq & \frac{1}{\epsilon^{2}} E\left(\sup _{t \leq T}\left(\left\langle M^{n}\right\rangle_{t+\delta}-\left\langle M^{n}\right\rangle_{t}\right) \wedge b\right) \\
& +P\left(\sup _{t \leq T}\left|\left\langle M^{n}\right\rangle_{t+\delta}-\left\langle M^{n}\right\rangle_{t}\right| \geq b\right) \\
\leq & \frac{1}{\epsilon^{2}} E\left(\sup _{t \leq T}\left|M_{t+\delta}^{n}-f(t+\delta)\right| \wedge b\right) \\
& +\frac{1}{\epsilon^{2}} E\left(\sup _{t \leq T}\left|M_{t}^{n}-f(t+\delta)\right| \wedge b\right) \\
& +\frac{1}{\epsilon^{2}} \sup _{t \leq T}|f(t+\delta)-f(t)| \\
& +P\left(\sup _{t \leq T}\left|\left\langle M^{n}\right\rangle_{t+\delta}-f(t+\delta)\right| \geq \frac{b}{3} \geq b\right) \\
& +P\left(\sup _{t \leq T}\left|\left\langle M^{n}\right\rangle_{t}-f(t)\right| \geq \frac{b}{3} \geq b\right) \\
& +1\left(\sup _{t \leq T}|f(t+\delta)-f(t)| \geq \frac{b}{3}\right) \\
& +P\left(\sup _{t \leq T}\left|\left\langle M^{n}\right\rangle_{t+\delta}-<M^{n}\right\rangle_{t} \mid \geq b\right) .
\end{aligned}
$$

Using the McLeish lemma, each term goes to 0 . This completes the proof.

If our conditions guarantee that for a locally square integrable martingale the associated increasing process converges, then the problem reduces to the convergence of finite-dimensional distributions for weak convergence of $\left\{M_{t}^{n}\right\}$ to the convergence of finite-dimensional distributions.

### 4.3 Central limit theorem for semi-martingale

Theorem 4.2: Let $\left\{X^{n}\right\}$ be a sequence of semi-martingale with characteristics ( $B^{n}\left[X^{n}, X^{n}\right], v^{n}$ ) and $M$ be a continuous Gaussian martingale with increasing process $<M>$.
(i) For any $t>0$ and $\epsilon \in(0,1)$, let the following conditions be satisfied:
(A)

$$
\int_{0}^{t} \int_{|x|>\varepsilon} v^{n}(d s, d x) \rightarrow_{p} 0
$$

(B)

$$
B_{t}^{n c}+\sum_{0 \leq s \leq t} \int_{|x| \leq \epsilon} x v^{n}(\{s\}, d x) \rightarrow_{p} 0
$$

(C)

$$
<X^{n c}>_{t}+\int_{0}^{t} \int_{|x| \leq \epsilon} x^{2} v^{n}(d s, d x)-\sum_{0 \leq s \leq t}\left(\int_{|x| \leq \epsilon} x v^{n}(\{s\}, d x)\right)^{2} \rightarrow_{p}<M>_{t}
$$

Then $X^{n} \Rightarrow M$ for finite dimension.
(ii) If (A) and (C) are satisfied as well as the condition

$$
\begin{equation*}
\sup _{0<s \leq t}\left|B_{s}^{n c}+\sum_{0 \leq u \leq s} \int_{|x| \leq \epsilon} x v^{n}(\{u\}, d x)\right| \rightarrow_{P} 0 \tag{4.5}
\end{equation*}
$$

for any $t$ and $\epsilon \in(0,1]$, then $X \Rightarrow M$ in $D[0, T]$.

Proof: Let $\mu^{n}$ be as defined in Section 4.1 and $v^{n}$ be associated predictable projection. For $\epsilon \in(0,1]$,

$$
\begin{aligned}
X_{t}^{n}= & \left(\sum_{0 \leq s \leq t} \int_{\epsilon<|x| \leq 1} x v^{n}(\{s\}, d x)\right)+\left(B_{t}^{n c}+\sum_{0 \leq s \leq t} \int_{|x| \leq \epsilon} x v^{n}(\{s\}, d x)\right) \\
& +\left(\int_{0}^{t} \int_{|x|>1} x \mu^{n}(d s, d x)+\int_{0}^{t} \int_{\epsilon<|x| \leq 1}\left(\mu^{n}-v^{n}\right)(d s, d x)\right) \\
& +\left(X_{t}^{n c}+\int_{0}^{t} \int_{|x| \leq \epsilon}\left(\mu^{n}-v^{n}\right)(d s, d x)\right) \\
= & \alpha_{t}^{n}(\epsilon)+\beta_{t}^{n}(\epsilon)+\gamma_{t}^{n}(\epsilon)+\Delta_{t}^{n}(\epsilon) .
\end{aligned}
$$

where

$$
\begin{aligned}
& \alpha_{t}^{n}(\epsilon)=\sum_{0 \leq s \leq t} \int_{\epsilon<|x| \leq 1} x v^{n}(\{s\}, d x) \\
& \beta_{t}^{n}(\epsilon)=B_{t}^{n c}+\sum_{0 \leq s \leq t} \int_{|x| \leq \epsilon} x v^{n}(\{s\}, d x) \\
& \gamma_{t}^{n}(\epsilon)=\int_{0}^{t} \int_{|x|>1} x \mu^{n}(d s, d x)+\int_{0}^{t} \int_{\epsilon<|x| \leq 1}\left(\mu^{n}-v^{n}\right)(d s, d x) \\
& \Delta_{t}^{n}(\epsilon)=X_{t}^{n c}+\int_{0}^{t} \int_{|x| \leq \epsilon}\left(\mu^{n}-v^{n}\right)(d s, d x) .
\end{aligned}
$$

By (A) we have

$$
\sup _{s \leq t} \alpha_{s}^{n}(\epsilon) \rightarrow_{p} 0
$$

By (B) we have

$$
\beta_{t}^{n}(\epsilon) \rightarrow_{p} 0
$$

By (4.3) we have

$$
\sup _{s \leq t} \beta_{s}^{n}(\epsilon) \rightarrow_{p} 0
$$

Let

$$
Y_{t}^{n}=\gamma_{t}^{n}(\epsilon)+\Delta_{t}^{n}(\epsilon) .
$$

It suffices to prove $Y^{n} \rightarrow M$ on $D[0, T]$ for each $T$ (by the decomposition $Y^{n}$ does not depend on $\epsilon$ ). Next, we have

$$
\begin{equation*}
\sup _{0<t \leq T}\left|\gamma_{t}^{n}(\epsilon)\right| \leq \int_{0}^{T} \int_{|x|>1}|x| \mu^{n}(d s, d x)+\int_{0}^{T} \int_{|x|>\epsilon} \mu^{n}(d s, d x)+\int_{0}^{T} \int_{|x|>\epsilon} v^{n}(d s, d x) \tag{4.6}
\end{equation*}
$$

Note that the second term on RHS $\rightarrow 0$ by Lenglart inequality and the third by (A).
Therefore, if we can show that the first term of RHS goes to 0 , then $\sup _{0<t \leq T}$ $\left|\gamma_{t}^{n}(\epsilon)\right| \rightarrow 0$. We have

$$
\int_{0}^{T} \int_{|x|>1}|x| \mu^{n}(d s, d x)=\sum_{0<s \leq T}\left|\Delta X_{s}^{n}\right| 1_{\left(\left|\Delta X_{s}^{n}\right|>1\right)}
$$

For $\delta \in(0,1)$,

$$
\left\{\sum_{0<s \leq T}\left|\Delta X_{s}^{n}\right| 1_{\left(\left|\Delta X_{s}^{n}\right|>1\right)}>\delta\right\} \subset\left\{\sum_{0<s \leq T} 1_{\left(\left|\Delta X_{s}^{n}\right|>1\right)}>\delta\right\} .
$$

and

$$
\sum_{0<s \leq T} 1_{\left(\left|\Delta X_{s}^{n}\right|>1\right)}=\int_{0}^{T} \int_{|x|>1} \mu^{n}(d s, d x) \rightarrow_{p} 0 .
$$

by Lenglart inequality. Therefore, by (4.6), we have

$$
\sup _{0<t \leq T}\left|\gamma_{t}^{n}(\epsilon)\right| \rightarrow 0
$$

Now, only thing left is to show that

$$
\begin{gathered}
\Delta_{t}^{n}(\epsilon) \rightarrow 0 \\
\Delta_{t}^{n}(\epsilon)=X_{t}^{n c}+\int_{0}^{T} \int_{|x| \leq \epsilon} x\left(\mu^{n}-v^{n}\right)(d s, d x) .
\end{gathered}
$$

Since $\left(\mu^{n}-v^{n}\right)$ is martingale and $X^{n c}$ is martingale, $\Delta^{n}$ is martingale. Since $\Delta^{n}(\epsilon) \epsilon$ $\mathcal{M}_{L O C}\left(\left(\mathcal{F}^{n}\right)_{t}, P\right)$,

$$
\begin{aligned}
<\Delta^{n}(\epsilon)>_{t}= & <X^{n c}>_{t}+\int_{0}^{t} \int_{|x| \leq \varepsilon} x^{n} v^{n}(d s, d x) \\
& -\sum_{0<s \leq t}\left(\int_{|x| \leq \epsilon} x v^{n}(\{s\}, d x)\right)^{2} \\
& \longrightarrow<M>_{t}
\end{aligned}
$$

by condition (C).
By McLeish lemma,

$$
\begin{equation*}
\sup _{t \leq T}\left|<\Delta^{n}(\epsilon)>_{t}-<M>_{t}\right| \rightarrow_{p} 0 \tag{4.7}
\end{equation*}
$$

We showed $\sup _{t \leq T}\left|\gamma_{t}^{n}(\epsilon)\right| \rightarrow 0$. Combining this with (4.7), we have

$$
\max \left(\sup _{t \leq T}\left|<\Delta^{n}(\epsilon)\right\rangle_{t}-<M>_{t}\left|, \sup _{t \leq T}\right| \gamma_{t}^{n}(\epsilon) \mid\right) \rightarrow 0
$$

Then, there exists $\left\{\epsilon_{n}\right\}$ such that

$$
\begin{gather*}
\sup _{t \leq T}\left|<\Delta^{n}\left(\epsilon_{n}\right)>_{t}-<M>_{t}\right| \rightarrow 0, \quad \sup _{t \leq T}\left|\gamma_{t}^{n}\left(\epsilon_{n}\right)\right| \rightarrow 0  \tag{4.8}\\
M_{t}^{n}=\Delta_{t}^{n}\left(\epsilon_{n}\right), \quad Y_{t}^{n}=\Delta_{t}^{n}\left(\epsilon_{n}\right)+\gamma_{t}^{n}\left(\epsilon_{n}\right) .
\end{gather*}
$$

It suffices to prove that $M^{n} \Rightarrow M .\left\{M_{t}^{n}=\Delta_{t}^{n}\left(\epsilon_{n}\right)\right\}$ is compact by Lemma 4.3 and (4.7). It suffices to prove finite-dimensional convergence.

Let $H(t), 0 \leq t \leq T$ be a piecewise constant left-continuous function assuming finitely many values. Let

$$
N_{t}^{n}=\int_{0}^{t} H(s) d M_{s}^{n}, \quad N_{t}=\int_{0}^{t} H(s) d M_{s} .
$$

Since $M$ is Gaussian, $N$ is also Gaussian.
Remark: Cramer-Wold criterion for finite-dimensional convergence

$$
E e^{i N_{T}^{n}} \rightarrow E e^{i N_{T}}=e^{\left.-\frac{1}{2} \int_{0}^{T} H^{2}(s) d<M\right\rangle_{s}} .
$$

Let $A$ be predictable, $A \in \mathcal{A}_{\text {LOC }}(\mathcal{F}, P)$

$$
e(A)_{t}=e^{A_{t}} \prod_{0 \leq s \leq t}\left(1+\Delta A_{s}\right) e^{-A_{s}} .
$$

Then $e^{A_{t}}$ will be a solution of $d Z_{t}=Z_{t-} d A_{t}$ by Dolean-Dade's work. If $m \in \mathcal{M}_{L O C}$, then

$$
A_{t}=-\frac{1}{2}\left\langle m^{c}\right\rangle_{t}+\int_{0}^{t} \int_{R-\{0\}}\left(e^{i s x}-1-i x\right) v_{m}(d s, d x)
$$

Lemma 4.4: For some $a>0$ and $c \in(0,1)$, let $\left\langle m>_{\infty} \leq a, \sup _{t}\right| \Delta m_{t} \mid \leq c$. Then $\left(e\left(A_{t}\right), \mathcal{F}_{t}\right)$ is such that

$$
\left|e(A)_{t}\right| \geq \exp \left(-\frac{2 a}{1-c^{2}}\right)
$$

and the process $\left(Z_{t}, \mathcal{F}_{t}\right)$ with $Z_{t}=e^{i m t}\left(e(A)_{t}\right)^{-1}$ is a uniformly integrable martingale. We will use the lemma to prove

$$
E e^{i N_{T}^{n}} \rightarrow E e^{i N_{T}}=e^{-\frac{1}{2} \int_{0}^{T} H^{2}(s) d<M>_{s}} .
$$

Case 1: Let us first assume $\left\langle M^{n}>_{T} \leq a\right.$ and $a>1+\langle M\rangle_{T}$. Observe that

1. $\operatorname{By}(4.7),\left\langle N^{n}\right\rangle_{t} \rightarrow\langle N\rangle_{t}$.
2. $\left|\Delta_{t}^{n}\right| \leq 2 \epsilon_{n}$.
3. $\left|\Delta N_{t}^{n}\right| \leq 2 \lambda \epsilon_{n}=d_{n}$, where $\lambda=\max _{t \leq T}|H(s)|$.

We want to prove

$$
\begin{equation*}
E \exp \left(i N_{T}^{n}+\frac{1}{2}<N>_{T}\right) \rightarrow 1 \tag{4.9}
\end{equation*}
$$

Let $A_{t}^{n}$ be increasing process associated with $N_{t}^{n}$. Let $Z_{t}=e^{i N_{t}^{n}}\left(e\left(A^{n}\right)_{t}\right)^{-1}$. Choose $n_{0}$ such that $d_{n_{0}}=2 \lambda \epsilon_{n_{0}} \leq 1 / 2$. By Lemma 4.4, $Z^{n}$ is a martingale with $E Z_{T}^{n}=1$. To prove (57) is equivalent to proving

$$
\lim _{n \rightarrow \infty}\left(E \exp \left(i N_{T}^{n}+\frac{1}{2}<N>_{T}\right)-E e^{i N_{T}^{n}}\left(e\left(A^{n}\right)_{T}\right)^{-1}\right)=0 E \exp \left(i N_{T}^{n}+\frac{1}{2}<N>_{T}\right) \rightarrow 1 .
$$

Thus, it is sufficient to prove

$$
e\left(A^{n}\right)_{T} \rightarrow e^{\left.-\frac{1}{2}<N\right\rangle_{T}} .
$$

Recall

$$
A_{t}^{n}=-\frac{1}{2}<N>_{t}+\int_{0}^{t} \int_{|x| \leq d_{n}}\left(e^{i x}-1-i x\right) \tilde{v}^{n}(d s, d x)
$$

Let

$$
\alpha_{t}^{n}=\int_{|x| \leq d_{n}}\left(e^{i x}-1-i x\right) \tilde{v}^{n}(\{t\}, d x)
$$

Since $\left(e^{i x}-1-i x\right) \leq x^{2} / 2$, we have $\alpha_{t}^{n} \leq d_{n}^{2} / 2$. Therefore,

$$
\begin{aligned}
\sum_{0 \leq t \leq T}\left|\alpha_{t}^{n}\right| & =\frac{1}{2} \int_{0}^{T} \int_{|x| \leq d_{n}} x^{2} \tilde{v}^{n}(d t, d x) \\
& =\frac{1}{2}<N^{n}>_{T} \\
& =\frac{1}{2} \frac{\lambda^{2} a}{2} .
\end{aligned}
$$

Then,

$$
\prod_{0<t \leq T}\left(1+\alpha_{t}^{n}\right) e^{-\alpha_{t}^{n}} \rightarrow 1
$$

By definition of $e(A)_{t}$, it remains to prove

$$
\frac{1}{2}<N^{n c}>_{T}-\int_{0}^{T} \int_{|x| \leq d_{n}}\left(e^{i s x}-1-i s x\right) \tilde{v}^{n}(d s, d x) \rightarrow_{p} \frac{1}{2}<N>_{T} .
$$

By observation (a) and the form of $\left\langle N^{n}\right\rangle_{T}$, it suffices to prove

$$
\int_{0}^{T} \int_{|x| \leq d_{n}}\left(e^{i s x}-1-i s x\right) \tilde{v}^{n}(d s, d x) \rightarrow_{p} 0
$$

We have, since $\left|\left(e^{i s x}-1-i s x+x^{2} / 2\right)\right| \leq c_{s} \frac{x^{3}}{3}$,

$$
\begin{aligned}
\int_{0}^{T} \int_{|x| \leq d_{n}} \underbrace{\left(\left(e^{i s x}-1-i s x\right)+\frac{x^{2}}{2}\right)}_{\leq \frac{|x|^{3}}{6}} \tilde{v}^{n}(d s, d x) & \leq \frac{d_{n}}{6} \int_{0}^{T} \int_{|x| \leq d_{n}} x^{2} \tilde{v}^{n}(d s, d x) \\
& \leq \frac{d_{n}}{6}<N^{n}>_{T} \\
& \leq \frac{d_{n}}{6} \lambda^{2} a
\end{aligned}
$$

To dispose of assumption, define

$$
\left.\tau_{n}=\min \left\{t \leq T:<M^{n}\right\rangle_{t} \geq<M>_{T}+1\right\} .
$$

Then $\tau_{n}$ is stopping time. We have $\tau_{n}=T$ if $\left\langle M^{n}\right\rangle_{t} \ll M>_{T}+1$. Let $\tilde{M}^{n}=M_{t \wedge \tau_{n}}^{n}$. Then

$$
<\tilde{M}^{n}>_{T} \leq 1+\left\langle M>_{T}+\epsilon_{n}^{2} \leq 1+\left\langle M>_{T}+\epsilon_{1}^{2}\right.\right.
$$

and

$$
\lim _{n} P\left(\left|<\tilde{M}^{n}>_{t}-<\tilde{M}>_{t}\right|>\epsilon\right) \leq \lim _{n} P\left(\tau_{n}>T\right)=0 .
$$

Next,

$$
\begin{aligned}
\lim _{n \rightarrow \infty} E e^{i N_{T}^{n}} & =\lim _{n \rightarrow \infty} E\left(e^{i N_{t}^{n}}-e^{i N_{\Lambda \Lambda \tau_{n}}^{n}}\right)+\lim _{n \rightarrow \infty} E e^{i N_{\Lambda \wedge \tau_{n}}^{n}} \\
& =\lim _{n \rightarrow \infty} E\left(e^{i N_{T}^{n}}-e^{i N_{T \Lambda \tau_{n}}^{n}}\right)+E e^{i N_{T}} \\
& =E e^{i N_{T}} .
\end{aligned}
$$

The last equality follows from

$$
\lim _{n \rightarrow \infty}\left|E\left(e^{i N_{T}^{n}}-e^{i N_{T \Lambda \tau_{n}}^{n}}\right)\right| \leq 2 \lim _{n \rightarrow \infty} P\left(\tau_{n}>T\right)=0
$$

This completes the proof.

### 4.4 Application to survival analysis

Let $X$ be a positive random variable. Let $F$ and $f$ be cumulative distribution function and probability density function of $X$. Then, survival function $\bar{F}$ is defined as

$$
\bar{F}(t)=P(X>t)=1-F(t)
$$

Then, we have

$$
\begin{aligned}
P(t<X \leq t+\Delta t \mid X>t) & =\frac{P(t<X \leq t+\Delta t)}{\bar{F}(t)} \\
& =\frac{\int_{t}^{t+\Delta t} d F(s)}{\bar{F}(t)} .
\end{aligned}
$$

Since we know

$$
\frac{1}{\Delta t} \int_{t}^{t+\Delta t} f(x) d s \longrightarrow f(t) .
$$

as $\Delta t \rightarrow 0$, hazard rate, is defined as

$$
\begin{aligned}
h(t) & =\frac{f(t)}{\bar{F}(t)} \\
& =-\frac{d}{d t} \log \bar{F}(t) .
\end{aligned}
$$

Therefore, survival function can be written as

$$
\bar{F}(t)=\exp \left(-\int_{0}^{t} h(s) d s\right) .
$$

If the integrated hazard rate is given, then it determines uniquely life distribution. For example, think about the following:

$$
\tau_{F}=\sup \{s: F(s)<1\} .
$$

Consider now the following problem arising in clinical trials.
Let

1. $\quad X_{1}, \ldots ., X_{n}$ be i.i.d $F$ (life time distribution)
2. $U_{1}, \ldots ., U_{n}$ be i.i.d measurable function with distribution function $G$ with $G(\infty)<1$, which means $U_{i}$ are not random variable in some sense.

Now, consider

1. indicator for "alive or not at time $s$ ": $1\left(X_{i} \leq U_{i}, X_{i} \wedge U_{i} \leq s\right)$.
2. indicator for "alive and leave or not at time $s$ ": $U_{i} 1\left(X_{i} \wedge U_{i} \leq s\right)$.
3. indicator for "leave or not at time $s$ ": $1\left(X_{i} \wedge U_{i} \geq s\right)$.
and $\sigma$-field
$\mathcal{F}_{t}^{n}=\sigma\left(\left\{1\left(X_{i} \leq U_{i}, X_{i} \wedge U_{i} \leq s\right), U_{i} 1\left(X_{i} \wedge U_{i} \leq s\right), 1\left(X_{i} \wedge U_{i} \geq s\right), s \leq t, i=1,2, \ldots, n\right\}\right)$.
$\mathcal{F}_{t}^{n}$ is called information contained in censored data.

Let

$$
\beta(t)=\int_{0}^{t} h(s) d s
$$

If $\hat{\beta}(t)$ is an estimate of $\beta(t)$, then we can estimate survival function. The estimator of survival function will be

$$
\hat{\bar{F}}(t)=e^{-\hat{\beta}(t)}
$$

which will be approximately be

$$
\prod_{s \leq t}(1-d(\hat{\beta}(s))) .
$$

This is alternate estimate of survival function, which is called, Nelson estimate.
Let

$$
\begin{aligned}
& N_{n}(t)=\sum_{i=1}^{n} 1\left(X_{i} \leq U_{i}, X_{i} \wedge U_{i} \leq t\right) \\
& Y_{n}(t)=\sum_{i=1}^{n} 1\left(X_{i} \wedge U_{i} \geq t\right)
\end{aligned}
$$

Then, $\hat{\beta}(t)$, which is called Breslow estimator, will be

$$
\hat{\beta}(t)=\int_{0}^{t} \frac{d N_{n}(s)}{Y_{n}(s)} \approx \frac{\int_{t}^{t+\Delta t} d F(s)}{\bar{F}(t)} .
$$

Now, we consider another estimator of survival function, which is called the KaplanMeier estimator. It will be

$$
\prod_{s \leq t}\left(1-\frac{\Delta N_{n}(s)}{Y_{n}(s)}\right)
$$

Gill [12] showed the asymptomatic properties of the Kaplan-Meier estimator.
We can show that

$$
\begin{equation*}
\left|e^{-\hat{\beta}(t)}-\prod_{s \leq t}\left(1-\frac{\Delta N_{n}(s)}{Y_{n}(s)}\right)\right|=O\left(\frac{1}{n}\right) \tag{4.11}
\end{equation*}
$$

using following lemma.

Lemma 4.5: Let $\left\{\alpha^{n}(s), 0 \leq s \leq T, n \geq 1\right\}$ be real-valued function such that

1. $\left\{s \in(0, u]: \alpha^{n}(s) \neq 0\right\}$ is $P$-a.e. at most countable for each $n$.
2. $\quad \sum_{0<s \leq u}\left|\alpha^{n}(s)\right| \leq C$ with $C$ constant.
3. $\sup _{s \leq u}\left\{\left|\alpha^{n}(s)\right|\right\}=O\left(a_{n}\right)$, where $a_{n} \searrow 0$ as $n$ goes $\infty$.

Then,

$$
\sup _{t \leq u}\left|\prod_{0<s \leq t}\left(1-\alpha^{n}(s)\right)-\prod_{0<s \leq t} e^{-\alpha^{n}(s)}\right|=O\left(a_{n}\right) .
$$

Proof: We choose $n_{0}$ large such that for $n \geq n_{0} O\left(a_{n}\right)<\frac{1}{2}$. Since

$$
\begin{aligned}
\prod_{0<s \leq t}\left(1-\alpha^{n}(s)\right) e^{\alpha^{n}(s)} & =\exp \left(\sum_{0<s \leq t} \log \left(1-\alpha^{n}(s)\right)+\alpha^{n}(s)\right) \\
& =\exp \left(\sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)\right)^{j}\right) \text { (by Taylor expansion.), }
\end{aligned}
$$

for $n \geq n_{0}$, we have

$$
\begin{aligned}
\left|\prod_{0<s \leq t}\left(1-\alpha^{n}(s)\right) e^{\alpha^{n}(s)}-1\right| & \leq\left|\exp \left(\sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)^{j}-1\right)\right)\right| \\
& \leq e^{\eta}\left|\sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)\right)^{j}\right| .
\end{aligned}
$$

where

$$
0 \wedge \sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)\right)^{j}<\eta<0 \vee \sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)\right)^{j} .
$$

For large $n$,

$$
\begin{aligned}
\left|\sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{(-1)^{j+2}}{j}\left(\alpha^{n}(s)\right)^{j}\right| \leq & \sum_{0<s \leq t} \sum_{j=2}^{\infty} \frac{\left|\alpha^{n}(s)\right|^{j}}{j} \\
\leq & \underbrace{\sup _{s \leq u}\left|\alpha^{n}(s)\right|}_{=O\left(a_{n}\right)} \underbrace{\sum_{0<s \leq t}\left|\alpha^{n}(s)\right|}_{\leq t \cdot M} \underbrace{\sum_{j=1}^{\infty}\left(\frac{1}{2}\right)^{j-2} \frac{1}{j}}_{<\infty} \\
& \longrightarrow 0 .
\end{aligned}
$$

$\sum_{0<s \leq t}\left|\alpha^{n}(s)\right| \leq t \cdot M$ holds since $\left|\alpha^{n}(s)\right|$ will be bounded by $M$.
To prove (4.11), we let

$$
\alpha^{n}(s)=\frac{1}{Y_{n}(s)}, \quad s \leq T
$$

and

$$
\triangle N(s)=0, \quad s>T .
$$

We get $a_{n}=1 / n$ using the Glivenko-Cantelli theorem.

### 4.5 Asymptotic distribution of $\hat{\beta}(t)$ and Kaplan-Meier estimate

$X_{i}$ and $U_{i}$ are defined as previously. Again, define $N_{n}(t), Y_{n}(t)$ as

$$
\begin{aligned}
& N_{n}(t)=\sum_{i=1}^{n} 1\left(X_{i} \leq U_{i}, X_{i} \wedge U_{i} \leq t\right) \\
& Y_{n}(t)=\sum_{i=1}^{n} 1\left(X_{i} \wedge U_{i} \geq t\right)
\end{aligned}
$$

Then,

$$
\begin{aligned}
\beta(t) & =\int_{0}^{t} h(s) d s \\
& =\int_{0}^{t} \frac{f(s)}{\bar{F}(s)} d s \\
& =\int_{0}^{t} \frac{d F(s)}{1-F(s)} \\
\hat{\beta}(t) & =\int_{0}^{t} \frac{d N_{n}(s)}{Y_{n}(s)} .
\end{aligned}
$$

Using the above lemma, we can show that the difference between Kaplan-Meier and Nelson estimates is asymptotically of order $\frac{1}{n}$.

$$
\begin{aligned}
& \sup _{t \leq u}\left|\prod_{s \leq t}\left(1-\frac{\Delta N_{n}(s)}{Y_{n}(s)}\right)-e^{-\hat{\beta}(t)}\right| \\
= & \sup _{t \leq u}\left|\prod_{s \leq t}\left(1-\frac{\Delta N_{n}(s)}{Y_{n}(s)}\right)-\exp \left(-\int_{0}^{t} \frac{d N_{n}(s)}{Y_{n}(s)}\right)\right| \\
= & O_{P}\left(\frac{1}{n}\right) .
\end{aligned}
$$

Let

$$
\begin{aligned}
Q_{1}(s) & =P\left(X_{1} \wedge U_{1} \leq s, X_{1} \leq U_{1}\right) \\
H(s) & =P\left(X_{1} \wedge U_{1} \leq s\right) \\
\beta_{1}(t) & =\int_{0}^{t} \frac{d Q_{1}(s)}{(1-H(s-))}
\end{aligned}
$$

Assume that $X$ (with $F$ ) and $U$ (with $G$ ) are independent. Then,

$$
\begin{aligned}
\Delta Q_{1}(s) & =P\left(s \leq X_{1} \wedge U_{1} \leq s+\Delta s, X_{1} \leq U_{1}\right) \\
& =P\left(s \leq X_{1} \leq s+\Delta s, X_{1} \leq U_{1}\right) \\
& =P\left(s \leq X_{1} \leq s+\Delta s, U_{1} \geq s+\Delta s\right) \\
& =P\left(s \leq X_{1} \leq s+\Delta s\right) \cdot P\left(U_{1} \geq s+\Delta s\right) \\
d Q_{1}(s) & =(1-G(s-)) d F(s) .
\end{aligned}
$$

Similarly,

$$
\begin{equation*}
(1-H(s-))=(1-G(s-))(1-F(s)) . \tag{4.12}
\end{equation*}
$$

Then,

$$
\begin{aligned}
\beta_{1}(t) & =\int_{0}^{t} \frac{d Q_{1}(s)}{(1-H(s-))} \\
& =\int_{0}^{t} \frac{(1-G(s-)) d F(s)}{(1-G(s-))(1-F(s))} \\
& =\int_{0}^{t} \frac{d F(s)}{(1-F(s))}
\end{aligned}
$$

## Lemma 4.6 (PL 1): Let

$\mathcal{F}_{t}^{n}=\sigma\left(\left\{1\left(X_{i} \leq U_{i}, X_{i} \wedge U_{i} \leq s\right), U_{i} 1\left(X_{i} \wedge U_{i} \leq s\right), 1\left(X_{i} \wedge U_{i} \geq s\right), s \leq t, i=1,2, \ldots, n\right\}\right)$.
Suppose we have

1. $\left(N_{n}(t), \mathcal{F}_{t}^{n}\right)$ is a point process and

$$
\left\{N_{n}(t)-\int_{0}^{t} Y_{n}(s) \frac{d Q_{1}(s)}{(1-H(s-))}\right\}
$$

is martingale.
and
2. $\quad \hat{\beta}_{n}(t)=\int_{0}^{t} \frac{d N_{n}(t)}{Y_{n}(t)}$.

Then, $m_{n}(t)=\hat{\beta}_{n}(t)-\beta_{1}(t)$ is a locally square integrable martingale, and increasing process $\left\langle m_{n}\right\rangle_{t}$ will be

$$
<m_{n}>_{t}=\int_{0}^{t} \frac{1}{Y_{n}(s)} \frac{d Q_{1}(s)}{(1-H(s-))}
$$

Remark: If $X$ and $U$ are independent, from (4.12), we have

$$
\frac{d Q_{1}(s)}{(1-H(s-))}=\frac{d F(s)}{1-F(s)}
$$

Let us assume that $\left\{X_{i}\right\}$ and $\left\{U_{i}\right\}$ are independent and

$$
A_{n}(t)=\int_{0}^{t} Y_{n}(s) \frac{d Q_{1}(s)}{(1-H(s-))} .
$$

From the previous theorem we know that $m_{n}(t)=\left(\hat{\beta}_{n}(t)-\beta_{1}(t)\right)$ is a locally square integtrable martingale with

$$
<m_{n}>_{t}=\int_{0}^{t} \frac{1}{Y_{n}(s)} \frac{d Q_{1}(s)}{(1-H(s-))}
$$

Hence

$$
\sqrt{n}\left(\hat{\beta}_{n}(t)-\beta_{1}(t)\right) \Rightarrow_{\text {a.s. }} \gamma_{t}
$$

where $Y_{t}$ is a Gaussian martingale

$$
\begin{aligned}
\left\langle\sqrt{n} m_{n}\right\rangle_{t} & =\int_{0}^{t} \frac{n}{Y_{n}(s)} \frac{d Q_{1}(s)}{(1-H(s-))} \rightarrow C_{1}(t) \\
C_{1}(t) & =\int_{0}^{t} \frac{d Q_{1}(s)}{(1-H(s-))^{2}} \\
r_{t} & =C_{1}(t) .
\end{aligned}
$$

Also $\left\langle m_{n}\right\rangle_{t}=A_{n}(t)$ which gives by Glivenko-Cantelli Lemma $\left\langle m_{n}\right\rangle_{t}=0_{p}\left(\frac{1}{n}\right)$. Using Lenglart inequality we get

$$
\begin{equation*}
\sup _{s \leq t}\left|\hat{\beta}_{n}(s)-\beta_{1}(s)\right| \rightarrow_{P} 0 . \tag{4.13}
\end{equation*}
$$

Hence $\hat{\beta}_{n}(s)$ is consistent estimate of integrated hazard rate under the independence assumption above. We note that under the assumption

$$
\beta_{1}(t)=\int_{0}^{t} \frac{d F(s)}{(1-F(s))}
$$

and

$$
C_{1}(t)=\int_{0}^{1} \frac{d F(s)}{(1-F(s))^{2}(1-G(s-))} .
$$

With $\tau_{H}=\inf \{s: H(s-)<1\}$, the above results hold for all $t<\tau_{H}$ only.
Lemma 4.7: 1 . For $t<\tau_{H}$

$$
\sup _{s \leq t}\left|\exp \left(-\int_{0}^{t} \frac{d F(s)}{(1-F(s))}\right)-\exp \left(-\int_{0}^{t} \frac{d N_{n}(s)}{Y_{n}(s)}\right)\right| \rightarrow_{p} 0
$$

2. 

$$
\sqrt{n}\left[\exp \left(-\int_{0} \frac{d F(s)}{(1-F(s))}\right)-\exp \left(-\int_{0} \frac{d N_{n}(s)}{Y_{n}(s)}\right)\right] \rightarrow_{\infty} r \cdot \exp \left(-\int_{0} \frac{d F(s)}{(1-F(s))}\right)
$$

in $D(0, t]$ for $t<\tau_{H}$ with $\gamma$ as above.

Proof: Using Taylor expansion we get

$$
\begin{aligned}
& \exp \left(-\int_{0}^{t} \frac{d F(s)}{(1-F(s))}\right)-\exp \left(-\hat{\beta}_{n}(t)\right) \\
& =\exp \left(-\int_{0}^{t} \frac{d F(s)}{(1-F(s))}\right)\left\{\left(\beta_{1}(t)-\hat{\beta}_{n}(t)\right)+\frac{\left(\beta_{1}(t)-\hat{\beta}_{n}(t)\right)^{2}}{2} \exp \left(-h_{n}\right)\right\}
\end{aligned}
$$

with $h_{n}$ is a random variable satisfying $\beta_{1}(t) \wedge \hat{\beta}_{n}(t) \leq h_{n} \leq \beta_{1}(t) \vee \hat{\beta}_{n}(t)$. Since for $t<\tau_{H}, \exp \left(-h_{n}\right)$ is bounded by convergence of $\sup _{s \leq t}\left(\hat{\beta}_{n}(s)-\beta_{1}(s)\right) \rightarrow_{p} 0$, the result follows. To prove the second part, note that

$$
\sqrt{n}\left(\beta_{1}(\cdot)-\hat{\beta}_{n}(\cdot)\right)^{2}=\sqrt{n}\left(\beta_{1}(\cdot)-\hat{\beta}_{n}(\cdot)\right)\left(\beta_{1}(\cdot)-\hat{\beta}_{n}(\cdot)\right) \Rightarrow_{\mathcal{D}} \gamma \cdot 0=0
$$

by Slutsky theorem and the first term converges in distribution to $\gamma$.
Theorem 4.3 (R. Gill): Let $\hat{F}_{n}(t)=\prod_{s \leq t}\left(1-\frac{\Delta N_{n}(s)}{Y_{n}(s)}\right)$. Then under independence of $\left\{X_{i}\right\}$, $\left\{U_{i}\right\}$, we get that

$$
\frac{n\left(\hat{F}_{n}(\cdot)-F(\cdot)\right)}{1-F(\cdot)} \Rightarrow \gamma_{o}
$$

in $D[0, t]$ for $t<\tau_{H}$.
Proof: We have $\exp \left(-\beta_{1}(t)\right)=1-F(t)$ for $t<\tau_{H}$. Hence by previous Lemma

$$
\frac{\sqrt{n}}{(1-F(\cdot))}\left(\exp \left(-\hat{\beta}_{n}(\cdot)\right)-(1-F(\cdot))\right) \Rightarrow_{\mathcal{D}} \gamma .
$$

Although we do not present it here, similar techniques work to study so-called LindenBell estimator which arises in the study of truncated data (e.g. in randomly starting clinical trials). The latter has applications in astronomy. Interested readers are referred to ([26], [21]).

## 5 Central limit theorems for dependent random variables

When one collects data, the observed sample could produce dependent random variables. A simplest example of this could be

$$
S_{n, m}=\sum_{1}^{m} X_{n, m},
$$

which could be a martingale for each $n$. These kind of theorems were first considered by Billingsley [3]. A major breakthrough for convergence of interpolated $S_{n,[n, t]}$ was obtained by Dvoretsky [10]. Gordin [13] showed that for second-order stationary processes with "mixing" condition, the central limit problem can be reduced to that for martingales. We show, following Durrett and Resnick [9], that one can derive the convergence of the interpolated sequence associated with the martingale central limit theorem, which can be obtained by their extension of the Skorokhod embedding theorem and using the work in [8] to obtain convergence in the of the interpolated sequence associated with the stationary sequence central limit problem of the Brownian motion. This proves weak convergence of such sequence in $C[0,1]$.

In this chapter, we study the central limit theorems for dependent random variables using the Skorokhod embedding theorem.

Theorem 5.1 (Martingale central limit theorem (discrete): Let $\left\{S_{n}\right\}$ be a martingale. Let $S_{0}=0$ and $\{W(t), 0 \leq t<\infty\}$ be Brownian motion. Then there exists a sequence of stopping time, $0=T_{0} \leq T_{1} \leq T_{2} \cdots \leq T_{n}$, with respect to $\mathcal{F}_{t}^{W}$ such that

$$
\left(S_{0}, \ldots, S_{n}\right)={ }_{d}\left(W\left(T_{0}\right), \ldots, W\left(T_{n}\right)\right) .
$$

Proof: We use induction.

$$
T_{0}=0 .
$$

Assume there exists $\left(T_{0}, \ldots, T_{k-1}\right)$ such that

$$
\left(S_{0}, \ldots, S_{k-1}\right)={ }_{d}\left(W\left(T_{0}\right), \ldots, W\left(T_{k-1}\right)\right) .
$$

Note that the strong Markov property implies that $\left\{W\left(T_{k-1}+t\right)-W\left(T_{k-1}\right), t \geq 0\right\}$ is a Brownian motion, independent of $\mathcal{F}_{t}^{W}$. Look at the regular conditional distribution of $S_{k}-S_{k-1}$ given $S_{0}=s_{0}, \ldots, S_{k-1}=s_{k-1}$. Denote it by

$$
\begin{gathered}
\mu\left(S_{0}, \ldots, S_{k-1} ; B\right)=P\left(S_{k}-S_{k-1} \in B \mid S_{0}=S_{0}, \ldots, S_{k-1}=S_{k-1}\right) \text { for } B \in(B(\mathbb{R})) . \\
\text { So } \mu\left(S_{0}, S_{1}, \ldots, S_{k-1} ; B\right)=P\left(S_{k}-S_{k-1} \in B \mid S_{0} \ldots, S_{k-1}\right) .
\end{gathered}
$$

Since $S_{k}$ is a martingale, we have

$$
0=E\left(S_{k}-S_{k-1} \mid S_{0}, \ldots, S_{k-1}\right)=\int x \mu_{k}\left(S_{0}, \ldots, S_{k-1} ; d x\right)
$$

By Skorokhod's representation theorem, we see that for a.e. $S \equiv\left(S_{0}, \ldots, S_{k-1}\right)$, there exists a stopping time $\tilde{\tau}_{S}\left(\right.$ exist time from $\left.\left(U_{k}, V_{k}\right)\right)$ such that

$$
W\left(T_{k-1}+\tilde{\tau}_{S}\right)-W\left(T_{k-1}\right)=\tilde{W}\left(\tau_{k}\right)={ }_{d} \mu_{k}\left(S_{0}, \ldots, S_{k-1} ; \cdot\right) .
$$

We let $T_{k}=T_{k-1}+\tilde{\tau}_{S}$, then

$$
\left(S_{0}, S_{1}, \ldots, S_{k}\right)={ }_{d}\left(W\left(T_{0}\right), \ldots, W\left(T_{k}\right)\right),
$$

and the result follows by induction.
Remark: If $E\left(S_{k}-S_{k-1}\right)^{2}<\infty$, then

$$
E\left(\tilde{\tau}_{S} \mid S_{0}, \ldots, S_{k-1}\right)=\int x^{2} \mu_{k}\left(S_{0}, \ldots, S_{k-1} ; d x\right)
$$

since $W_{t}^{2}-t$ is a martingale and $\tilde{\tau}_{S}$ is the exit time from a randomly chosen interval $\left(S_{k-1}+U_{k}, S_{k-1}+V_{k}\right)$.

Definition 5.1: We say that $X_{n, m}, \mathcal{F}_{n, m}, 1 \leq m \leq n$, is a martingale difference array if $X_{n, m}$ is $\mathcal{F}_{n, m}$ measurable and $E\left(X_{n, m} \mid \mathcal{F}_{n, m-1}\right)=0$ for $1 \leq m \leq n$, where $\mathcal{F}_{n, 0}=\{\emptyset, \Omega\}$.

Notation: Let

$$
S_{(u)}= \begin{cases}S_{k}, & \text { if } u=k \in \mathbf{N} ; \\ \text { linear on } u, & \text { if } u \in[k, k+1] \text { for } k \in \mathbf{N} .\end{cases}
$$

and

$$
S_{n,(u)}= \begin{cases}S_{n, k}, & \text { if } u=k \in N ; \\ \text { linear on } u, & \text { if } u \in[k, k+1] .\end{cases}
$$

Consider $X_{n, m}(1 \leq m \leq n)$ be triangular arrays of random variables with

$$
\begin{aligned}
E X_{n, m} & =0 \\
S_{n, m} & =X_{n, 1}+\cdots+X_{n, m} .
\end{aligned}
$$

We shall use the following fact:

$$
S_{n, m}=W\left(\tau_{m}^{n}\right)
$$

and

$$
\tau_{[n s]}^{n} \rightarrow_{P} s \text { for } s \in[0,1],
$$

then $\left\|S_{n,(n \cdot)}-W(\cdot)\right\|_{\infty} \rightarrow_{P} 0$.
Theorem 5.2: Let $\left\{X_{n, m}, \mathcal{F}_{n, m}\right\}$ be a martingale difference array and $S_{n, m}=X_{n, 1}+\cdots+$ $X_{n, m}$. Assume that

1. $\left|X_{n, m}\right| \leq \epsilon_{n}$ for all $m$, and $\epsilon_{n} \rightarrow 0$ as $n \rightarrow \infty$
2. with $V_{n, m}=\sum_{k=1}^{m} E\left(X_{n, k}^{2} \mid \mathcal{F}_{n, k-1}\right), V_{n,[n t]} \rightarrow t$ for all $t$.

Then $S_{n,(n \cdot)} \Rightarrow W(\cdot)$.
Proof: We stop $V_{n, k}$ first time if it exceeds 2 (call it $k_{0}$ ) and set $X_{n, m}=0, m>k_{0}$. We can assume without loss of generality

$$
V_{n, n} \leq 2+\epsilon_{n}^{2}
$$

for all $n$. Using Theorem 5.1, we can find stopping times $T_{n, 1}, \ldots, T_{n, n}$ so that

$$
\left(0, S_{n, 1}, \ldots, S_{n, n}\right)=_{\mathscr{D}}\left(W(0), W\left(T_{n, 1}\right), \ldots, W\left(T_{n, n}\right)\right)
$$

Using Lemma 5.1 and the above equality, it suffices to show that $T_{n,[n t]} \rightarrow_{P} t$ for each $t$. Let

$$
t_{n, m}=T_{n, m}-T_{n, m-1} \text { with }\left(T_{n, 0}=0\right) .
$$

Using the Skohorod embedding theorem, we have

$$
E\left(t_{n, m} \mid \mathcal{F}_{n, m-1}\right)=E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)
$$

The last observation with Hypothesis 2 implies

$$
\sum_{m=0}^{[n t]} E\left(t_{n, m} \mid \mathcal{F}_{n, m-1}\right) \rightarrow_{p} t
$$

Observe that

$$
\begin{aligned}
E\left(T_{n,[n t]}-V_{n,[n t]}\right)^{2} & =E(\sum_{m=1}^{[n t]}(\underbrace{t_{n, m}-E\left(t_{n, m} \mid \mathcal{F}_{n, m-1}\right)}_{\text {any two terms are orthogonal }}))^{2} \\
& =\sum_{m=1}^{[n t]} E\left(t_{n, m}-E\left(t_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right)^{2}
\end{aligned}
$$

$$
\begin{align*}
& \leq \sum_{m=1}^{[n t]} E\left(t_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& \leq \sum_{m=1}^{[n t]} C \cdot E\left(X_{n, m}^{4} \mid \mathcal{F}_{n, m-1}\right) \text { ( we will show } \\
& \\
& \text { that } C=4 .) \\
& \leq \sum_{m=1}^{[n t]} C \epsilon_{n}^{2} E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \text { (by Assumption 1) } \\
& =C \epsilon_{n}^{2} V_{n, n}  \tag{5.1}\\
& \leq C \epsilon_{n}^{2}\left(2+\epsilon_{n}^{2}\right) \rightarrow 0 .
\end{align*}
$$

Since $L^{2}$ convergence implies convergence in probability,

$$
E\left(T_{n,[n t]}-V_{n,[n t]}\right)^{2} \longrightarrow 0
$$

and

$$
V_{n,[n t]} \longrightarrow{ }_{P} t
$$

together implies

$$
T_{n,[n t]} \longrightarrow{ }_{p} t .
$$

Proof of (5.1): If $\theta$ is real, then

$$
E\left(\left.\exp \left(\theta(W(t)-W(s))-\frac{1}{2} \theta^{2}(t-s)\right) \right\rvert\, \mathcal{F}_{S}^{W}\right)=1
$$

Since

$$
E\left(\left.\exp \left(\theta W(t)-\frac{1}{2} \theta^{2} t\right) \right\rvert\, \mathcal{F}_{s}^{W}\right)=\exp \left(\theta W(s)-\frac{1}{2} \theta^{2} s\right),
$$

we know that $\left\{\exp \left(\theta W(t)-\frac{1}{2} \theta^{2} t\right), \mathcal{F}_{t}^{W}\right\}$ is a martingale. Then, for all $A \in \mathcal{F}_{s}^{W}$,

$$
\begin{aligned}
& E 1_{A}\left(\exp \left(\theta W(s)-\frac{1}{2} \theta^{2} s\right)\right)= \int_{A}\left(\exp \left(\theta W(s)-\frac{1}{2} \theta^{2} s\right)\right) d P \\
&= \int_{A} \exp \left(\theta W(t)-\frac{1}{2} \theta^{2} t\right) d P \\
& \quad \text { (by definition of conditional expectation) } \\
&= E 1_{A}\left(\exp \left(\theta W(t)-\frac{1}{2} \theta^{2} t\right)\right)
\end{aligned}
$$

Take a derivative in $\theta$ and find a value at $\theta=0$.

Number of derivative

| 1 | $W(t)$ is MG |
| :--- | :--- |
| 2 | $W^{2}(t)-t$ is MG |
| 3 | $W^{3}(t)-3 t W(t)$ is MG |
| 4 | $W^{4}(t)-6 t W^{2}(t)+3 t^{2}$ is MG |

For any stopping time $\tau$,

$$
E\left(W^{4}(\tau)-6 \tau W^{2}(\tau)+3 \tau^{2}\right)=0
$$

Therefore,

$$
\begin{aligned}
& E W_{\tau}^{4}-6 E\left(\tau W_{\tau}^{2}\right)=-3 E W_{\tau}^{2} \\
\Rightarrow & E W_{\tau}^{2} \leq 2 E\left(\tau W_{\tau}^{2}\right) .
\end{aligned}
$$

Since

$$
E\left(\tau W_{\tau}^{2}\right) \leq\left(E \tau^{2}\right)^{1 / 2} \cdot\left(E W_{\tau}^{4}\right)^{1 / 2}
$$

by Schwartz Inequality, we have

$$
\left(E \tau^{2}\right)^{1 / 2} \leq 2\left(E W_{\tau}^{4}\right)^{1 / 2}
$$

Therefore,

$$
E\left(t_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \leq 4 E\left(X_{n, m}^{4} \mid \mathcal{F}_{n, m-1}\right)
$$

Theorem 5.3 Generalization of the Lindberg-Feller theorem: Let $\left\{X_{n, m}, \mathcal{F}_{n, m}\right\}$ be a martingale difference array and $S_{n, m}=X_{n, 1}+\cdots+X_{n, m}$. Assume that

1. $\quad V_{n,[n t]}=\sum_{k=1}^{[n t]} E\left(X_{n, k}^{2} \mid \mathcal{F}_{n, k-1}\right) \rightarrow_{P} t$.
2. $\widehat{V}(\epsilon)=\sum_{m \leq n} E\left(X_{n, m}^{2} 1\left(\left|X_{n, m}\right|>\epsilon\right) \mid \mathcal{F}_{n, m-1}\right) \rightarrow_{P} 0$, for all $\epsilon>0$.

Then $S_{n,(n \cdot)} \Rightarrow W(\cdot)$.
For i.i.d. $X_{n, m}$ and $t=1$, we get the Lindberg-Feller theorem.
Lemma 5.2: There exists $\epsilon_{n} \rightarrow 0$ such that $\epsilon_{n}^{2} \hat{V}\left(\epsilon_{n}\right) \rightarrow_{P} 0$.
Proof: Since $\hat{V}(\epsilon) \rightarrow_{P} 0$, we choose large $N_{m}$ such that

$$
P\left(\hat{V}\left(\frac{1}{m}\right)>\frac{1}{m^{3}}\right)<\frac{1}{m}
$$

for $m \geq N_{m}$. Here we choose $\epsilon_{n}=\frac{1}{m}$ with $n \in\left[N_{m}, N_{m+1}\right)$. For $\delta>\frac{1}{m}$, we have

$$
P\left(\epsilon_{n}^{-2} \widehat{V}\left(\epsilon_{n}\right)>\delta\right) \leq P\left(m^{2} \widehat{V}\left(\frac{1}{m}\right)>\frac{1}{m}\right)<\frac{1}{m} .
$$

This completes the proof of lemma.
Let

$$
\begin{aligned}
& \bar{X}_{n, m}=X_{n, m} 1\left(\left|X_{n, m}\right| \leq \epsilon_{n}\right) \\
& \widehat{X}_{n, m}=X_{n, m} 1\left(\left|X_{n, m}\right|>\epsilon_{n}\right) \\
& \widetilde{X}_{n, m}=\bar{X}_{n, m}-E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right) .
\end{aligned}
$$

Lemma 5.3: $\widetilde{S}_{n,[n \cdot]} \Rightarrow W(\cdot)$
Proof: We will show that $\widetilde{X}_{n, m}$ satisfies Theorem 5.2.

$$
\begin{aligned}
\left|\widetilde{X}_{n, m}\right| & =\left|\bar{X}_{n, m}-E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| \\
& \leq\left|\bar{X}_{n, m}\right|+\left|E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| \\
& \leq 2 \epsilon_{n} \rightarrow 0 .
\end{aligned}
$$

and hence, the first condition is satisfied. Since

$$
X_{n, m}=\bar{X}_{n, m}+\widetilde{X}_{n, m},
$$

we have

$$
\begin{align*}
E\left(\bar{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) & =E\left(\left(X_{n, m}-\widehat{X}_{n, m}\right)^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& \left.=E\left(X_{n, m}^{2}-2 X_{n, m} \widehat{X}_{n, m}+\widehat{X}_{n, m}\right)^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& =E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)-E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) . \tag{5.3}
\end{align*}
$$

The last equality follows from $E\left(X_{n, m} \widehat{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)=E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)$. Since $X_{n, m}$ is a martingale difference array, and hence, $E\left(X_{n, m} \mid \mathcal{F}_{n, m-1}\right)=0$. The last observation implies $E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)=-E\left(\widehat{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)$, and hence,

$$
\begin{aligned}
{\left[E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right]^{2} } & =\left[E\left(\widehat{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right]^{2} \\
& \leq E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \text { (by Jensen's inequality). }
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\sum_{m=1}^{n}\left[E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right]^{2} & \leq \sum_{m=1}^{n} E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& =\widehat{V}\left(\epsilon_{n}\right) \\
& \rightarrow_{P} 0
\end{aligned}
$$

by given condition. Finally,

$$
\begin{aligned}
\sum_{m=1}^{n} E\left(\widetilde{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)= & \sum_{m=1}^{n} E\left(\bar{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& -\sum_{m=1}^{n} E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)^{2} \\
& \text { (by the conditional variance formula) } \\
= & \sum_{m=1}^{n}\left(E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)-E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)\right) \\
& -\sum_{m=1}^{n} E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)^{2} \\
& (\text { from equation (5.3)) } \\
= & \sum_{m=1}^{n} E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)-\sum_{m=1}^{n} E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& -\sum_{m=1}^{n}\left[E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right]^{2} .
\end{aligned}
$$

As $\hat{V}\left(\epsilon_{n}\right) \xrightarrow{p} 0$ and the last term converges to zero in probability, we get

$$
\lim _{n \rightarrow \infty} \sum_{m=1}^{n} E\left(\widetilde{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)=\lim _{n \rightarrow \infty} \sum_{m=1}^{n} E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right)
$$

Since

$$
V_{n,[n t]}=\sum_{m=1}^{[n t]} E\left(X_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \rightarrow_{P} t,
$$

we conclude that

$$
\sum_{m=1}^{[n t]} E\left(\widetilde{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \rightarrow_{P} t
$$

This show that the second condition is satisfied, and this completes the proof.

## Lemma 5.4:

$$
\left\|S_{n,(n \cdot)}-\widetilde{S}_{n,(n \cdot)}\right\|_{\infty} \leq \sum_{m=1}^{n}\left|E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| .
$$

Proof: Note that if we prove this lemma, then, since $\sum_{m=1}^{n}\left|E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| \rightarrow_{P} 0$ (we will show this), and we construct a Brownian motion with $\left\|\widetilde{S}_{n,(n \cdot)}-W(\cdot)\right\|_{\infty} \rightarrow 0$, the desired result follows from the triangle inequality.

Since $X_{n, m}$ is martingale difference array, we know that $E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)=$ $-E\left(\widehat{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)$, and hence,

$$
\begin{aligned}
\sum_{m=1}^{n}\left|E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right|= & \sum_{m=1}^{n}\left|E\left(\widehat{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| \\
\leq & \sum_{m=1}^{n} E\left(\left|\widehat{X}_{n, m}\right| \mid \mathcal{F}_{n, m-1}\right) \text { by Jensen } \\
\leq & \frac{1}{\epsilon_{n}} \sum_{m=1}^{n} E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \\
& \left(\text { if }\left|X_{n, m}\right|>\epsilon_{n}, \widehat{X}_{n, m} \leq \frac{X_{n, m}^{2}}{\epsilon_{n}}=\frac{\widehat{X}_{n, m}^{2}}{\epsilon_{n}}\right) \\
= & \frac{\widehat{V}\left(\epsilon_{n}\right)}{\epsilon_{n}} \rightarrow_{p} 0(\text { by Lemma } 5.2) .
\end{aligned}
$$

On $\left\{\left|X_{n, m}\right| \leq \epsilon_{n}, 1 \leq m \leq n\right\}$, we have $\bar{X}_{n, m}=X_{n, m}$, and hence, $S_{n,(n \cdot)}=\bar{S}_{n,(n \cdot)}$. Thus,

$$
\begin{aligned}
\left\|S_{n,(n \cdot)}-\widetilde{S}_{n,(n \cdot)}\right\|_{\infty} & =\left\|S_{n,(n \cdot)}-\bar{S}_{n,(n \cdot)}+\sum_{m=1}^{[n \cdot]} E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right\|_{\infty} \\
& \leq \sum_{m=1}^{[n \cdot]}\left|E\left(\bar{X}_{n, m} \mid \mathcal{F}_{n, m-1}\right)\right| \rightarrow_{P} 0 .
\end{aligned}
$$

Now, to complete the proof, we have to show that Lemma 5.3 holds on $\left\{\left|X_{n, m}\right|>\epsilon_{n}, 1 \leq\right.$ $m \leq n\}$. It suffices to show that

## Lemma 5.5:

$$
P\left(\left|X_{n, m}\right|>\epsilon_{n}, \text { for some } m, 1 \leq m \leq n\right) \rightarrow 0 .
$$

To prove Lemma 5.5, we use Dvoretsky's proposition.

Proposition 5.1 (Dvoretsky): Let $\left\{\mathcal{G}_{n}\right\}$ be a sequence of $\sigma$-fields with $\mathcal{G}_{n} \subset \mathcal{G}_{n+1}$. If $A_{n} \in \mathcal{G}_{n}$ for each $n$, then for each $\delta \geq 0$, measurable with respect to $\mathcal{G}_{0}$,

$$
\begin{equation*}
P\left(\bigcup_{m=1}^{n} A_{m} \mid \mathcal{G}_{0}\right) \leq \delta+P\left(\sum_{m=1}^{n} P\left(A_{m} \mid \mathcal{G}_{m-1}\right)>\delta \mid \mathcal{G}_{0}\right) \tag{5.4}
\end{equation*}
$$

Proof: We use induction.
(i) $n=1$

We want to show

$$
\begin{equation*}
P\left(A_{1} \mid \mathcal{G}_{0}\right) \leq \delta+P\left(P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta \mid \mathcal{G}_{0}\right) . \tag{5.5}
\end{equation*}
$$

Consider $\Omega_{\ominus}=\left\{\omega: P\left(A_{1} \mid \mathcal{G}_{0}\right) \leq \delta\right\}$. Then (5.5) holds. Also, consider $\Omega_{\oplus}=\{\omega$ : $\left.P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta\right\}$. Then

$$
\begin{aligned}
P\left(P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta \mid \mathcal{G}_{0}\right) & =E\left(\mathbf{1}\left(P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta\right) \mid \mathcal{G}_{0}\right) \\
& =\mathbf{1}\left(P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta\right) \\
& =1
\end{aligned}
$$

and hence (5.5) also holds.
(ii) $n>1$

Consider $\omega \in \Omega_{\oplus}$. Then

$$
\begin{aligned}
P\left(\sum_{m=1}^{n} P\left(A_{m} \mid \mathcal{G}_{m-1}\right)>\delta \mid \mathcal{G}_{0}\right) & \geq P\left(P\left(A_{1} \mid \mathcal{G}_{0}\right)>\delta \mid \mathcal{G}_{0}\right) \\
& =1_{\Omega_{\oplus}}(\omega) \\
& =1
\end{aligned}
$$

Then, (5.4) holds. Consider $\omega \in \Omega_{\ominus}$. Let $B_{m}=A_{m} \cap \Omega_{\ominus}$. Then, for $m \geq 1$,

$$
\begin{aligned}
P\left(B_{m} \mid \mathcal{G}_{m-1}\right) & =P\left(A_{m} \cap \Omega_{\ominus} \mid \mathcal{G}_{m-1}\right) \\
& =P\left(A_{m} \mid \mathcal{G}_{m-1}\right) \cdot P\left(\Omega_{\ominus} \mid \mathcal{G}_{m-1}\right) \\
& =P\left(A_{m} \mid \mathcal{G}_{m-1}\right) \cdot 1_{\Omega_{\ominus}}(\omega) \\
& =P\left(A_{m} \mid \mathcal{G}_{m-1}\right) .
\end{aligned}
$$

Now suppose $\gamma=\delta-P\left(B_{1} \mid \mathcal{G}_{0}\right) \geq 0$ and apply the last result for $n-1$ sets (induction hypothesis).

$$
P\left(\bigcup_{m=2}^{n} B_{m} \mid \mathcal{G}_{1}\right) \leq \gamma+P\left(\sum_{m=2}^{n} P\left(B_{m} \mid \mathcal{G}_{m-1}\right)>\gamma \mid \mathcal{G}_{1}\right) .
$$

Recall $E\left(E\left(X \mid \mathcal{G}_{0}\right) \mid \mathcal{G}_{1}\right)=E\left(X \mid \mathcal{G}_{0}\right)$ if $\mathcal{G}_{0} \subset \mathcal{G}_{1}$. Taking conditional expectation w.r.t. $\mathcal{G}_{0}$ and noting $\gamma \in \mathcal{G}_{0}$,

$$
P\left(\bigcup_{m=2}^{n} B_{m} \mid \mathcal{G}_{0}\right) \leq P\left(\gamma+P\left(\sum_{m=2}^{n} P\left(B_{m} \mid \mathcal{G}_{m-1}\right)>\gamma \mid \mathcal{G}_{1}\right) \mid \mathcal{G}_{0}\right)
$$

$$
\begin{aligned}
& =\gamma+P\left(\sum_{m=2}^{n} P\left(B_{m} \mid \mathcal{G}_{m-1}\right)>\gamma \mid \mathcal{G}_{0}\right) \\
& =\gamma+P\left(\sum_{m=1}^{n} P\left(B_{m} \mid \mathcal{G}_{m-1}\right)>\delta \mid \mathcal{G}_{0}\right) .
\end{aligned}
$$

Since $\cup B_{m}=\left(\cup A_{m}\right) \cap \Omega_{\ominus}$, on $\Omega_{\ominus}$ we have

$$
\sum_{m=1}^{n} P\left(B_{m} \mid \mathcal{G}_{m-1}\right)=\sum_{m=1}^{n} P\left(A_{m} \mid \mathcal{G}_{m-1}\right) .
$$

Thus, on $\Omega_{\ominus}$,

$$
P\left(\bigcup_{m=2}^{n} A_{m} \mid \mathcal{G}_{0}\right) \leq \delta-P\left(A_{1} \mid \mathcal{G}_{0}\right)+P\left(\sum_{m=1}^{n} P\left(A_{m} \mid \mathcal{G}_{m-1}\right)>\delta \mid \mathcal{G}_{0}\right) .
$$

The result follows from

$$
P\left(\bigcup_{m=1}^{n} A_{m} \mid \mathcal{G}_{0}\right) \leq P\left(A_{1} \mid \mathcal{G}_{0}\right)+P\left(\bigcup_{m=2}^{n} A_{m} \mid \mathcal{G}_{0}\right)
$$

using monotonicity of conditional expectation and $1_{A \cup B} \leq 1_{A}+1_{B}$.
Proof of Lemma 5.5: Let $A_{m}=\left\{\left|X_{n, m}\right|>\epsilon_{n}\right\}, \mathcal{G}_{m}=\mathcal{F}_{n, m}$, and let $\delta$ be a positive number. Then, Proposition 5.1 implies

$$
P\left(\left|X_{n, m}\right|>\epsilon_{n} \text { for some } m \leq n\right) \leq \delta+P\left(\sum_{m=1}^{n} P\left(\left|X_{n, m}\right|>\epsilon_{n} \mid \mathcal{F}_{n, m-1}\right)>\delta\right) .
$$

To estimate the right-hand side, we observe that "Chebyshev's inequality" implies

$$
\sum_{m=1}^{n} P\left(\left|X_{n, m}\right|>\epsilon_{n} \mid \mathcal{F}_{n, m-1}\right) \leq \epsilon_{n}^{-2} \sum_{m=1}^{n} E\left(\widehat{X}_{n, m}^{2} \mid \mathcal{F}_{n, m-1}\right) \rightarrow 0
$$

so $\lim \sup P\left(\left|X_{n, m}\right|>\epsilon_{n}\right.$ for some $\left.m \leq n\right) \leq \delta$. Since $\delta$ is arbitrary, the proof of lemma and theorem is complete.

Theorem 5.4 (Martingale cental limit theorem): Let $\left\{X_{n}, \mathcal{F}_{n}\right\}$ be a martingale difference sequence, $X_{n, m}=X_{m} / \sqrt{n}$, and $V_{k}=\sum_{n=1}^{k} E\left(X_{n}^{2} \mid \mathcal{F}_{n-1}\right)$. Assume that

1. $V_{k} / k \rightarrow_{P} \sigma^{2}$,
2. $n^{-1} \sum_{m \leq n} E\left(X_{m}^{2} 1\left(\left|X_{m}\right|>\epsilon \sqrt{n}\right)\right) \rightarrow 0$.

Then,

$$
\frac{S_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot)
$$

Definition 5.2: A process $\left\{X_{n}, n \geq 0\right\}$ is called stationary if

$$
\left\{X_{m}, X_{m+1}, \ldots, X_{m+k}\right\}==_{\mathscr{D}}\left\{X_{0}, \ldots, X_{k}\right\}
$$

for any $m, k$.

Definition 5.3: Let $(\Omega, \mathcal{F}, P)$ be a probability space. A measurable map $\varphi: \Omega \rightarrow \Omega$ is said to be measure preserving if $P\left(\varphi^{-1} A\right)=P(A)$ for all $A \in \mathcal{F}$.

Theorem 5.5: If $X_{0}, X_{1}, \ldots$ is a stationary sequence and $g: \mathbf{R}^{\mathbf{N}} \rightarrow \mathbf{R}$ is measurable, then $Y_{k}=g\left(X_{k}, X_{k+1}, \ldots\right)$ is a stationary sequence.

Proof: If $\mathbf{x} \in \mathbf{R}^{\mathbf{N}}$, let $g_{k}(\mathbf{x})=g\left(x_{k}, x_{k+1}, \ldots\right)$, and if $B \in \mathcal{R}^{\mathbf{N}}$ let

$$
A=\left\{\mathbf{x}:\left(g_{0}(\mathbf{x}), g_{1}(\mathbf{x}), \ldots\right) \in B\right\} .
$$

To check stationarity now, we observe:

$$
\begin{aligned}
P\left(\left\{\omega:\left(Y_{0}, Y_{1}, \ldots\right) \in B\right\}\right)= & P\left(\left\{\omega:\left(g\left(X_{0}, X_{1}, \ldots\right), g\left(X_{1}, X_{2}, \ldots\right), \ldots\right) \in B\right\}\right) \\
= & P\left(\left\{\omega:\left(g_{0}(\mathbf{X}), g_{1}(\mathbf{X}), \ldots\right) \in B\right\}\right) \\
= & P\left(\left\{\omega:\left(X_{0}, X_{1}, \ldots\right) \in A\right\}\right) \\
= & P\left(\left\{\omega:\left(X_{k}, X_{k+1}, \ldots\right) \in A\right\}\right) \\
& \left(\text { since } X_{0}, X_{1}, \ldots \text { is a stationary sequence }\right) \\
= & P\left(\left\{\omega:\left(Y_{k}, Y_{k+1}, \ldots\right) \in B\right\}\right) .
\end{aligned}
$$

Exercise: Show the above equality.

Definition 5.4: Assume that $\theta$ is measure preserving. A set $A \in \mathcal{F}$ is said to be invariant if $\theta^{-1} A=A$. Denote by $\mathcal{I}=\left\{A\right.$ in $\left.\mathcal{F}, A=\theta^{-1} A\right\}$ in measure.

Definition 5.5: A measure preserving transformation on $(\Omega, \mathcal{F}, P)$ is said to be ergodic if $\mathcal{I}$ is trivial, i.e., for every $A \in \mathcal{I}, P(A) \in\{0,1\}$.

Example: Let $X_{0}, X_{1}, \ldots$ be the i.i.d. sequence. If $\Omega=\mathbf{R}^{\mathbf{N}}$ and $\theta$ is the shift operator, then an invariant set $A$ has $\{\omega: \omega \in A\}=\{\omega: \theta \omega \in A\} \in \sigma\left(X_{1}, X_{2}, \ldots\right)$. Iterating gives

$$
A \in \bigcap_{n=1}^{\infty} \sigma\left(X_{n}, X_{n+1}, \ldots\right)=\mathcal{T} \text {, the tail } \sigma \text {-field }
$$

so $\mathcal{I} \subset \mathcal{T}$. For an i.i.d. sequence, Kolmogorov's 0-1 law implies $\mathcal{T}$ is trivial, so $\mathcal{I}$ is trivial and the sequence is ergodic. We call $\theta$ the ergodic transformation.

Theorem 5.6: Let $g: \mathbf{R}^{\mathbf{N}} \rightarrow \mathbf{R}$ be measurable. If $X_{0}, X_{1}, \ldots$ is an ergodic stationary sequence, then $Y_{k}=g\left(X_{k}, X_{k+1}, \ldots\right)$ is ergodic.

Proof: Suppose $X_{0}, X_{1}, \ldots$ is defined on sequence space with $X_{n}(\omega)=\omega_{n}$. If $B$ has $\left\{\omega:\left(Y_{0}, Y_{1}, \ldots\right) \in B\right\}=\left\{\omega:\left(Y_{1}, Y_{2}, \ldots\right) \in B\right\}$, then $A=\left\{\omega:\left(Y_{0}, Y_{1}, \ldots\right) \in B\right\}$ is shift invariant.

Theorem 5.7 (Birkhoff's ergodic theorem): For any $f \in L_{1}(P)$,

$$
\frac{1}{n} \sum_{m=0}^{n-1} f\left(\theta^{m} \omega\right) \rightarrow E(f \mid \mathcal{G}) \text { a.s. and in } L_{1}(P)
$$

where $\theta$ is measure preserving transformation on $(\Omega, \mathcal{F}, P)$ and $\mathcal{G}=\left\{A \in \mathcal{F}: \theta^{-1}\right.$ $A=A\}$.

For the proof, see [8].

Theorem 5.8: Suppose $\left\{X_{n}, n \in \mathbf{Z}\right\}$ is an ergodic stationary sequence of martingale differences, i.e., $\sigma^{2}=E X_{n}^{2}<\infty$ and $E\left(X_{n} \mid \mathcal{F}_{n-1}\right)=0$ with respect to $\mathcal{F}_{n}=\sigma\left(X_{m}\right.$, $m \leq n)$. Let $S_{n}=X_{1}+\cdots+X_{n}$. Then,

$$
\frac{S_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot)
$$

Proof: Let $u_{n}=E\left(X_{n}^{2} \mid \mathcal{F}_{n-1}\right)$. Then $u_{n}$ can be written as $\theta\left(X_{n-1}, X_{n-2}, \ldots\right)$, and hence, by Theorem $5.6 u_{n}$ is stationary and ergodic. By Birkhoff's ergodic theorem $(\mathcal{G}=\{\emptyset, \Omega\})$,

$$
\frac{1}{n} \sum_{m=1}^{n} u_{m} \rightarrow E u_{0}=E X_{0}^{2}=\sigma^{2} \text { a.s. }
$$

The last conclusion shows that (i) of Theorem 5.4 holds. To show (ii), we observe

$$
\begin{aligned}
\frac{1}{n} \sum_{m=1}^{n} E\left(X_{m}^{2} 1\left(\left|X_{m}\right|>\epsilon \sqrt{n}\right)\right)= & \frac{1}{n} \sum_{m=1}^{n} E\left(X_{0}^{2} 1\left(\left|X_{0}\right|>\epsilon \sqrt{n}\right)\right) \\
& \text { (because of stationarity) } \\
= & E\left(X_{0}^{2} 1\left(\left|X_{0}\right|>\epsilon \sqrt{n}\right)\right) \rightarrow 0
\end{aligned}
$$

by the dominated convergence theorem. This completes the proof.
Let's consider stationary process, (with $\xi_{i} i . i . d ., E \xi_{i}=0$ and $E \xi_{i}^{2}<\infty$ )

$$
X_{n}=\sum_{k=0}^{\infty} c_{k} \xi_{n-k}
$$

with

$$
\sum_{k=0}^{\infty} c_{k}^{2}<\infty
$$

If $\xi_{i}$ are i.i.d., $X_{n}$ is definitely stationary, but it is not martingale difference process. This is called moving average process. What we will do is we start with stationary ergodic process, and then we will show that limit of this process is the limit of martingale difference sequence. Then this satisfies the conditions of martingale central limit theorem. Note that in our example, $E X_{n}^{2}<\infty$ for all $n$.

If $X_{n}$ is stationary second-order, then $r(n)=E X_{n} \bar{X}_{o}$ is positive definite.
We can separate phenomenon into two parts: new information(nondeterministic) and non-new information (deterministic).

$$
E X_{n} X_{0}=\int_{0}^{2 \pi} e^{i n \lambda} d F(\lambda)
$$

where $F$ is the spectral measure. In case $X_{n}=\sum_{k=0}^{\infty} c_{k} \xi_{n-k}$, then $F \ll$ Lebesgue measure.

In fact, we have the following theorem:
Theorem 5.9: There exist $\bar{c}_{k}$ and $\phi$ such that
$-f\left(e^{i \lambda}\right)=\left|\phi\left(e^{i \lambda}\right)\right|^{2}$
$-\quad \phi\left(e^{i \lambda}\right)=\sum_{k=0}^{\infty} \bar{c}_{k} e^{i k \lambda}$
if and only if

$$
\int_{0}^{2 \pi} \log f(\lambda) d F(\lambda)>-\infty
$$

Now we start with $\left\{X_{n}: n \in \mathbf{Z}\right\}$ ergodic stationary sequence such that

- $E X_{n}=0, E X_{n}^{2}<\infty$.
- $\quad \sum_{n=1}^{\infty}\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2}<\infty$.

The idea is if we go back, then $X_{n}$ will be independent of $X_{0}$.
Let

$$
\begin{aligned}
H_{n} & =\left\{Y \in \mathcal{F}_{n} \text { with } E Y^{2}<\infty\right\}=L^{2}\left(\Omega, \mathcal{F}_{n}, P\right) \\
K_{n} & =\left\{Y \in H_{n} \text { with } E Y Z=0 \text { for all } Z \in H_{n-1}\right\}=H_{n} \ominus H_{n-1} .
\end{aligned}
$$

Geometrically, $H_{0} \supset H_{-1} \supset H_{-2} \ldots$ is a sequence of subspaces of $L^{2}$ and $K_{n}$ is the orthogonal complement of $H_{n-1}$. If $Y$ is a random variable, let

$$
\left(\theta^{n} Y\right)(\omega)=Y\left(\theta^{n} \omega\right)
$$

i.e., $\theta$ is isometry(measure-preserving) on $L^{2}$. Generalizing from the example $Y=f\left(X_{-j}, \ldots, X_{k}\right)$, which has $\theta^{n} Y=f\left(X_{n-j}, \ldots, X_{n+k}\right)$, it is easy to see that if $Y \in H_{k}$, then $\theta^{n} Y \in H_{k+n}$, and hence $Y \in K_{j}$ then $\theta^{n} Y \in K_{n+j}$.

Lemma 5.6: Let $P$ be a projection such that $X_{j} \in H_{-j}$ implies $P_{H_{-j}} X_{j}=X_{j}$. Then,

$$
\begin{aligned}
\theta P_{H_{-1}} X_{j} & =P_{H_{0}} X_{j+1} \\
& =P_{H_{0}} \theta X_{j} .
\end{aligned}
$$

Proof: For $j \leq-1$,

$$
\theta \underbrace{P_{H_{-j}} X_{j}}_{X_{j}}=\theta X_{j}=X_{j+1} .
$$

We will use this property. For $Y \in H_{-1}$,

$$
\begin{aligned}
& X_{j}-P_{H_{-1}} X_{j} \perp Y \\
\Rightarrow & \left(X_{j}-P_{H_{-1}} X_{j}, Y\right)_{2}=0 \\
\Rightarrow & \left(\theta\left(X_{j}-P_{H_{-1}} X_{j}\right), \theta Y\right)_{2}=0\left(\text { since } \theta \text { is isometry on } L^{2}\right)
\end{aligned}
$$

Since $Y \in H_{-1}, \theta Y$ generates $H_{0}$. Therefore, for all $Z \in H_{0}$, we have

$$
\begin{aligned}
& \left(\left(\theta X_{j}-\theta P_{H_{-1}} X_{j}\right), Z\right)_{2}=0 \\
\Rightarrow & \left(\theta X_{j}-\theta P_{H_{-1}} X_{j}\right) \perp Z \\
\Rightarrow & \theta P_{H_{-1}} X_{j}=P_{H_{0}} \theta X_{j}=P_{H_{0}} X_{j+1} .
\end{aligned}
$$

We come now to the central limit theorem for stationary sequences.
Theorem 5.10: Suppose $\left\{X_{n}, n \in \mathbf{Z}\right\}$ is an ergodic stationary sequence with $E X_{n}=0$ and $E X_{n}^{2}<\infty$. Assume

$$
\sum_{n=1}^{\infty}\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2}<\infty
$$

Let $S_{n}=X_{1}+\ldots+X_{n}$. Then

$$
\frac{S_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot)
$$

where we do not know what $\sigma$ is.

Proof: If $X_{0}$ happened to be in $K_{0}$ since $X_{n}=\theta^{n} X_{0} \in K_{n}$ for all $n$, and taking $Z=1_{A} \in$ $H_{n-1}$, we would have $E\left(X_{n} 1_{A}\right)=0$ for all $A \in \mathcal{F}_{n-1}$ and hence $E\left(X_{n} \mid \mathcal{F}_{n-1}\right)=0$. The next best thing to having $X_{n} \in K_{0}$ is to have

$$
\begin{equation*}
X_{0}=Y_{0}+Z_{0}-\theta Z_{0} \tag{*}
\end{equation*}
$$

with $Y_{0} \in K_{0}$ and $Z_{0} \in L^{2}$. Let

$$
\begin{aligned}
Z_{0} & =\sum_{j=0}^{\infty} E\left(X_{j} \mid \mathcal{F}_{-1}\right) \\
\theta Z_{0} & =\sum_{j=0}^{\infty} E\left(X_{j+1} \mid \mathcal{F}_{0}\right) \\
Y_{0} & =\sum_{j=0}^{\infty}\left(E\left(X_{j} \mid \mathcal{F}_{0}\right)-E\left(X_{j} \mid \mathcal{F}_{-1}\right)\right) .
\end{aligned}
$$

Then we can solve (*) formally

$$
\begin{equation*}
Y_{0}+Z_{0}-\theta Z_{0}=E\left(X_{0} \mid \mathcal{F}_{0}\right)=X_{0} \tag{5.6}
\end{equation*}
$$

We let

$$
S_{n}=\sum_{m=1}^{n} X_{m}=\sum_{m=1}^{n} \theta^{m} X_{0} \quad \text { and } T_{n}=\sum_{m=1}^{n} \theta^{m} Y_{0} .
$$

We want to show that $T_{n}$ is martingale difference sequence. We have $S_{n}=T_{n}+\theta Z_{0}-$ $\theta^{n+1} Z_{0}$. The $\theta^{m} Y_{0}$ are a stationary ergodic martingale difference sequence (ergodicity follows from Theorem 5.6), so Theorem 5.8 implies

$$
\frac{T_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot), \text { where } \sigma^{2}=E Y_{0}^{2}
$$

To get rid of the other term, we observe

$$
\frac{\theta Z_{0}}{\sqrt{n}} \rightarrow 0 \text { a.s. }
$$

and

$$
\begin{aligned}
P\left(\max _{0 \leq m \leq n-1}\left|\theta^{m+1} Z_{0}\right|>\epsilon \sqrt{n}\right) & \leq \sum_{m=0}^{n-1} P\left(\left|\theta^{m+1} Z_{0}\right|>\epsilon \sqrt{n}\right) \\
& =n P\left(\left|Z_{0}\right|>\epsilon \sqrt{n}\right) \\
& \leq \epsilon^{-2} E\left(Z_{0}^{2} \mathbf{1}_{\left\{\left|Z_{0}\right|>\epsilon \sqrt{n}\right\}}\right) \rightarrow 0
\end{aligned}
$$

The last inequality follows from the stronger form of Chevyshev,

$$
E\left(Z_{0}^{2} \mathbf{1}_{\left\{\left|Z_{0}\right|>\epsilon \sqrt{n}\right\}}\right) \geq \epsilon^{2} n P\left(\left|Z_{0}\right|>\epsilon \sqrt{n}\right) .
$$

Therefore, in view of the above comments,

$$
\begin{aligned}
& \frac{S_{n}}{\sqrt{n}}=\frac{T_{n}}{\sqrt{n}}+\frac{\theta Z_{0}}{\sqrt{n}}-\frac{\theta^{n+1} Z_{0}}{\sqrt{n}} \\
\Rightarrow & \frac{S_{n}}{\sqrt{n}}-\frac{T_{n}}{\sqrt{n}} \xrightarrow{p} 0 \\
\Rightarrow & \lim _{n \rightarrow \infty} \frac{S_{(n \cdot)}}{\sqrt{n}}=\lim _{n \rightarrow \infty} \frac{T_{(n \cdot)}}{\sqrt{n}}=\sigma W(\cdot) .
\end{aligned}
$$

Theorem 5.11: Suppose $\left\{X_{n}, n \in \mathbf{Z}\right\}$ is an ergodic stationary sequence with $E X_{n}=0$ and $E X_{n}^{2}<\infty$. Assume

$$
\sum_{n=1}^{\infty}\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2}<\infty
$$

Let $S_{n}=X_{1}+\ldots+X_{n}$. Then

$$
\frac{S_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot)
$$

where

$$
\sigma^{2}=E X_{0}^{2}+2 \sum_{n=1}^{\infty} E X_{0} X_{n}
$$

If $\sum_{n=1}^{\infty} E X_{0} X_{n}$ diverges, theorem will not be true. We will show that $\sum_{n=1}^{\infty}\left|E X_{0} X_{n}\right|<\infty$. This theorem is different from previous theorem since we now specify $\sigma^{2}$.

Proof: First,

$$
\begin{aligned}
\left|E X_{0} X_{m}\right| & =\left|E\left(E\left(X_{0} X_{m} \mid \mathcal{F}_{0}\right)\right)\right| \\
& \leq E\left|X_{0} E\left(X_{m} \mid \mathcal{F}_{0}\right)\right| \\
& \leq\left\|X_{0}\right\|_{2} \cdot\left\|E\left(X_{m} \mid \mathcal{F}_{0}\right)\right\|_{2}(\text { by Cauchy Schwarz inequality }) \\
& =\left\|X_{0}\right\|_{2} \cdot\left\|E\left(X_{0} \mid \mathcal{F}_{-m}\right)\right\|_{2}(\text { by shift invariance }) .
\end{aligned}
$$

Therefore, by assumption,

$$
\sum_{n=1}^{\infty}\left|E X_{0} X_{n}\right| \leq\left\|X_{0}\right\|_{2} \sum_{n=1}^{\infty}\left\|E\left(X_{0} \mid \mathcal{F}_{-m}\right)\right\|_{2}<\infty .
$$

Next,

$$
\begin{aligned}
E S_{n}^{2} & =\sum_{j=1}^{n} \sum_{k=1}^{n} E X_{j} X_{k} \\
& =n E X_{0}^{2}+2 \sum_{m=1}^{n-1}(n-m) E X_{0} X_{m}
\end{aligned}
$$

From this, it follows easily that

$$
\frac{E S_{n}^{2}}{n} \rightarrow E X_{0}^{2}+2 \sum_{m=1}^{\infty} E X_{0} X_{m} .
$$

To finish the proof, let $T_{n}=\sum_{m=1}^{n} \theta^{m} Y_{0}$, observe $\sigma^{2}=E Y_{0}^{2}$, and

$$
\begin{aligned}
n^{-1} E\left(S_{n}-T_{n}\right)^{2} & =n^{-1} E\left(\theta Z_{0}-\theta^{n+1} Z_{0}\right)^{2} \\
& \leq \frac{3 E Z_{0}^{2}}{n} \rightarrow 0
\end{aligned}
$$

since $(a-b)^{2} \leq(2 a)^{2}+(2 b)^{2}$.
We proved central limit theorem of ergodic stationary process. We will discuss examples: M-dependence and moving average.

Example 1. M-dependent sequences: Let $X_{n}, n \in \mathbb{Z}$ be a stationary sequence with $E X_{n}=0, E X_{n}^{2}<\infty$. Assume that $\sigma\left(\left\{X_{j}, j \leq 0\right\}\right)$ and $\sigma\left(\left\{X_{j}, j \geq M\right\}\right)$ are independent. In this case, $E\left(X_{0} \mid \mathcal{F}_{-n}\right)=0$ for $n>M$, and $\sum_{n=0}^{\infty}\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2}<\infty$. Let $\mathcal{F}_{-\infty}=\cap_{m} \sigma\left(\left\{X_{j}, j \geq M\right\}\right)$ and $\mathcal{F}_{k}=\sigma\left(\left\{X_{j}, j \leq k\right\}\right)$. If $m-k>M$, then $\mathcal{F}_{-\infty} \perp \mathcal{F}_{k}$. Recall Kolmogorov 0-1 law. If $A \in \mathcal{F}_{k}$ and $B \in \mathcal{F}_{-\infty}$, then $P(A \cap B)=P(A) \cdot P(B)$. For all $A \in$ $\cup_{k} \mathcal{F}_{k}, A \in \sigma\left(\cup_{k} \mathcal{F}_{k}\right)$. Also, $A \in \mathcal{F}_{-\infty}$, where $\mathcal{F}_{-\infty} \subset \cup_{k} \mathcal{F}_{-k}$. Therefore, by Kolmogorov 0-1 law, $P(A \cap A)=P(A) \cdot P(A)$, and hence, $\left\{X_{n}\right\}$ is stationary and ergodic. Thus, Theorem 5.8 implies

$$
\frac{S_{n,(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot)
$$

where

$$
\sigma^{2}=E_{0}^{2}+2 \sum_{m=1}^{M} E X_{0} X_{m}
$$

Example 2. Moving average: Suppose

$$
X_{m}=\sum_{k \geq 0} c_{k} \xi_{m-k}, \text { where } \sum_{k \geq 0} c_{k}^{2}<\infty,
$$

and $\xi_{i}, i \in \mathbf{Z}$ are i.i.d. with $E \xi_{i}=0$ and $E \xi_{i}^{2}=1$. Clearly, $\left\{X_{n}\right\}$ is a stationary sequence since series converges. Check whether $\left\{X_{n}\right\}$ is ergodic. We have

$$
\bigcap_{n} \sigma\left(\left\{X_{m}, m \leq n\right\}\right) \subset \bigcap_{n} \sigma\left(\left\{\xi_{k}, k \leq n\right\}\right) ;
$$

therefore, by Kolmogorov 0-1 law, $\left\{X_{n}\right\}$ is ergodic. Next, if $\mathcal{F}_{-n}=\sigma\left(\left\{\xi_{m}, m \leq-n\right\}\right)$, then

$$
\begin{aligned}
\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2} & =\left\|\sum_{k \geq n} c_{k} \xi_{k}\right\|_{2} \\
& =\left(\sum_{k \geq n} c_{k}^{2}\right)^{1 / 2} .
\end{aligned}
$$

If, for example, $c_{k}=(1+k)^{-p},\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2} \sim n^{(1 / 2-p)}$, and Theorem 5.11 applies if $p>3 / 2$.

Let $\mathcal{G}, \mathcal{H} \subset \mathcal{F}$, and

$$
\alpha(\mathcal{G}, \mathcal{H})=\sup _{A \in \mathcal{G}, B \in \mathcal{H}}\{|P(A \cap B)-P(A) P(B)|\} .
$$

If $\alpha=0, \mathcal{G}$, and $\mathcal{H}$ are independent, $\alpha$ measures the dependence of two $\sigma$-algebras.
Lemma 5.7: Let $p, q, r \in(1, \infty]$ with $1 / p+1 / q+1 / r=1$, and suppose $X \in \mathcal{G}, Y \in \mathcal{H}$ have $E|X|^{p}, E|Y|^{q}<\infty$. Then

$$
|E X Y-E X E Y| \leq 8\|X\|_{p}\|Y\|_{q}(\alpha(\mathcal{G}, \mathcal{H}))^{1 / r}
$$

Here, we interpret $x^{0}=1$ for $x>0$ and $0^{0}=0$.

Proof: If $\alpha=0, X$ and $Y$ are independent and the result is true, so we can suppose $\alpha>0$. We build up to the result in three steps, starting with the case $r=\infty$.
(a). $r=\infty$

$$
|E X Y-E X E Y| \leq 2\|X\|_{p}\|Y\|_{q} .
$$

Proof of (a): Hölder's inequality implies $|E X Y| \leq\|X\|_{p}\|Y\|_{q}$, and Jensen's inequality implies

$$
\|X\|_{p}\|Y\|_{q} \geq|E| X|E| Y| | \geq|E X E Y| ;
$$

thus, the result follows from the triangle inequality.
(b) $X, Y \in L^{\infty}$

$$
|E X Y-E X E Y| \leq 4\|X\|_{\infty}\|Y\|_{\infty} \alpha(\mathcal{G}, \mathcal{H})
$$

Proof of (b): Let $\eta=\operatorname{sgn}(E(Y \mid \mathcal{G})-E Y) \in \mathcal{G} . E X Y=E(X E(Y \mid \mathcal{G}))$, so

$$
\begin{aligned}
|E X Y-E X E Y| & =|E(X(E(Y \mid \mathcal{G})-E Y))| \\
& \leq\|X\|_{\infty} E|E(Y \mid \mathcal{G})-E Y| \\
& =\|X\|_{\infty} E(\eta E(Y \mid \mathcal{G})-E Y) \\
& =\|X\|_{\infty}(E(\eta Y)-E \eta E Y) .
\end{aligned}
$$

Applying the last result with $X=Y$ and $Y=\eta$ gives

$$
|E(Y \eta)-E Y E \eta| \leq\|Y\|_{\infty}|E(\zeta \eta)-E \zeta E \eta|,
$$

where $\zeta=\operatorname{sgn}(E(\eta \mid \mathcal{H})-E \eta)$. Now $\eta=1_{A}-1_{B}$ and $\zeta=1_{C}-1_{D}$, so

$$
\begin{aligned}
|E(\zeta \eta)-E \zeta E \eta|= & \mid P(A \cap C)-P(B \cap C)-P(A \cap D)+P(B \cap D) \\
& -P(A) P(C)+P(B) P(C)+P(A) P(D)-P(B) P(D) \mid \\
\leq & 4 \alpha(\mathcal{G}, \mathcal{H}) .
\end{aligned}
$$

Combining the last three displays gives the desired result.
(c) $q=\infty, 1 / p+1 / r=1$

$$
|E X Y-E X E Y| \leq 6| | X\left\|_{p}\right\| Y \|_{\infty}(\alpha(\mathcal{G}, \mathcal{H}))^{1-1 / p}
$$

Proof of (c): Let $C=\alpha^{-1 / p}\|X\|_{p}, X_{1}=X 1_{(|X| \leq C)}$, and $X_{2}=X-X_{1}$.

$$
\begin{aligned}
|E X Y-E X E Y| & \leq\left|E X_{1} Y-E X_{1} E Y\right|+\left|E X_{2} Y-E X_{2} E Y\right| \\
& \leq 4 \alpha C| | Y\left\|_{\infty}+2\right\| Y \|_{\infty} E\left|X_{2}\right|
\end{aligned}
$$

by (a) and (b). Now

$$
E\left|X_{2}\right| \leq C^{-(p-1)} E\left(|X|^{p} 1_{(|X| \leq C)}\right) \leq C^{-p+1} E|X|^{p} .
$$

Combining the last two inequalities and using the definition of $C$ gives

$$
|E X Y-E X E Y| \leq 4 \alpha^{1-1 / p}\|X\|_{p}\|Y\|_{\infty}+2\|Y\|_{\infty} \alpha^{1-1 / p}\|X\|_{p}^{-p+1+p},
$$

which is the desired result.

Finally, to prove Lemma 5.7, let $C=\alpha^{-1 / q}\|Y\|_{q}, Y_{1}=Y 1_{(|Y| \leq C)}$, and $Y_{2}=Y-Y_{1}$.

$$
\begin{aligned}
|E X Y-E X E Y| & \leq\left|E X Y_{1}-E X E Y_{1}\right|+\left|E X Y_{2}-E X E Y_{2}\right| \\
& \leq 6 C\|X\|_{p} \alpha^{1-1 / p}+2\|X\|_{p}\left\|Y_{2}\right\|_{\theta},
\end{aligned}
$$

where $\theta=(1-1 / p)^{-1}$ by (c) and (a). Now

$$
E|Y|^{\theta} \leq C^{-q+\theta} E\left(|Y|^{q} 1_{(|Y| \leq C)}\right) \leq C^{-1+\theta} E|Y|^{q} .
$$

Taking the $1 / \theta$ root of each side and recalling the definition of $C$

$$
\left\|Y_{2}\right\|_{\theta} \leq C^{-(q-\theta)}\|Y\|_{q}^{q / \theta} \leq \alpha^{(q-\theta) / q \theta}\|Y\|_{q},
$$

so we have

$$
|E X Y-E X E Y| \leq 6 \alpha^{-1 / q}\|Y\|_{q}\|X\|_{p} \alpha^{1-1 / p}+2\|X\|_{p} \alpha^{1 / \theta-1 / q}\|Y\|_{q}^{1 / \theta+1 / q},
$$

proving Lemma 5.7.
Combining Theorem 5.11 and Lemma 5.7 gives:

Theorem 5.12: Suppose $X_{n}, n \in \mathbf{Z}$ is an ergodic stationary sequence with $E X_{n}=0$, $E\left|X_{0}\right|^{2+\delta}<\infty$. Let $\alpha(n)=\alpha\left(\mathcal{F}_{-n}, \sigma\left(X_{0}\right)\right)$, where $\mathcal{F}_{-n}=\sigma\left(\left\{X_{m}, m \leq-n\right\}\right)$, and suppose

$$
\sum_{n=1}^{\infty} \alpha(n)^{\delta / 2(2+\delta)}<\infty
$$

If $S_{n}=X_{1}+\cdots X_{n}$, then

$$
\frac{S_{(n \cdot)}}{\sqrt{n}} \Rightarrow \sigma W(\cdot),
$$

where

$$
\sigma^{2}=E X_{0}^{2}+2 \sum_{n=1}^{\infty} E X_{0} X_{n}
$$

Proof: To use Lemma 5.7 to estimate the quantity in Theorem 5.11 we begin with

$$
\|E(X \mid \mathcal{F})\|_{2}=\sup \left\{E(X Y): Y \in \mathcal{F},\|Y\|_{2}=1\right\} \quad(*) .
$$

Proof of (*): If $Y \in \mathcal{F}$ with $\|Y\|_{2}=1$, then using a by now familiar property of conditional expectation and the Cauchy-Schwarz inequality

$$
E X Y=E(E(X Y \mid \mathcal{F}))=E(Y E(X \mid \mathcal{F})) \leq\|E(X \mid \mathcal{F})\|_{2}\|Y\|_{2}
$$

Equality holds when $Y=E(X \mid \mathcal{F}) /\|E(X \mid \mathcal{F})\|_{2}$.

Letting $p=2+\delta$ and $q=2$ in Lemma 5.7, noticing

$$
\frac{1}{r}=1-\frac{1}{p}-\frac{1}{q}=\frac{\delta}{2(2+\delta)}
$$

and recalling $E X_{0}=0$, showing that if $Y \in \mathcal{F}_{-n}$

$$
\left|E X_{0} Y\right| \leq 8\left\|X_{0}\right\|_{2+\delta}\|Y\|_{2} \alpha(n)^{\delta / 2(2+\delta)}
$$

Combining this with (*) gives

$$
\left\|E\left(X_{0} \mid \mathcal{F}_{-n}\right)\right\|_{2} \leq 8\left\|X_{0}\right\|_{2+\delta} \alpha(n)^{\delta / 2(2+\delta)},
$$

and it follows that the hypotheses of Theorem 5.11 are satisfied.

## 6 Empirical process

We have so far considered the weak convergence of stochastic processes with values in complete separable metric spaces. The main assumption needed to study such convergence in terms of bounded continuous functions is that the measures be defined on Borel subsets of the space. This influences theorems like the Prokhorov theorem. In other words, the random variables with values in separable metric spaces. If we consider the space $D[0,1]$ with sup norm, we get non-separability. We note that if we consider the case of empirical processes

$$
X_{n}(t)=\sqrt{n}\left(\hat{F}_{n}(t)-t\right),
$$

where $\hat{F}_{n}(t)$ is the empirical distribution function of i.i.d. uniform random variable, then the above $X_{n}$ is not a Borel measureable function in $D[0,1]$ with sup norm if we consider the domain of $X_{n}(., w)$ with $w \epsilon[0,1]$ ? Thus, the weak convergence definition above cannot be used in this case. Dudley [6] developed an alternative weak convergence theory. Even this cannot handle the general empirical processes, hence some statistical applications. One therefore has to introduce outer expectations of $X_{n}$ of possibly non-measureable maps as far as limit of $X_{n}$ is Borel measureable. We shall also consider convergence in probability and almost sure convergence for such functions. We therefore follow this idea as in [25]. A similar idea was exploited in the basic paper on invariance principle in non-separable Banach spaces by Dudley and Phillips ([7], see also [18]).

Let $(\Omega, \mathcal{A}, P)$ be an arbitrary probability space and $T: \Omega \rightarrow \overline{\mathbf{R}}$ an arbitrary map. The outer integral of $T$ with respect to $P$ is defined as

$$
E^{*} T=\inf \{E U: U \geq T, U: \Omega \rightarrow \overline{\mathbf{R}} \text { measurable and } E U \text { exists }\} .
$$

Here, as usual, $E U$ is understood to exist if at least one of $E U^{+}$or $E U^{-}$is finite. The outer probability of an arbitrary subset $B$ of $\Omega$ is

$$
P^{*}(B)=\inf \{P(A): A \supset B, A \in \mathcal{A}\}
$$

Note that the functions $U$ in the definition of outer integral are allowed to take the value $\infty$, so that the infimum exists.

The inner integral and inner probability can be defined in a similar fashion. Equivalently, they can be defined by $E_{*} T=-E^{*}(-T)$ and $P_{*}(B)=1-P^{*}(B)$, respectively.

In this section, $\mathbf{D}$ is metric space with a metric $d$. The set of all continuous, bounded functions $f: \mathbf{D} \rightarrow \mathbf{R}$ is denoted by $C_{b}(\mathbf{D})$.

Definition 6.1: Let $\left(\Omega_{\alpha}, \mathcal{A}_{\alpha}, P_{\alpha}\right), \alpha \in I$ be a net of probability spaces, $X_{\alpha}: \Omega_{\alpha} \rightarrow \mathbf{D}$, and $P_{\alpha}=P \circ X_{\alpha}^{-1}$. Then we say that the net $X_{\alpha}$ converges weakly to a Borel measure $L$, i.e., $P_{\alpha} \Rightarrow L$ if

$$
E^{*} f\left(X_{\alpha}\right) \rightarrow \int f d L, \quad \text { for every } f \in C_{b}(\mathbf{D})
$$

Theorem 6.1 (Portmanteau): The following statements are equivalent:
(i) $P_{\alpha} \Rightarrow L$;
(ii) $\liminf P_{*}\left(X_{\alpha} \in G\right) \geq L(G)$ for every open $G$;
(iii) $\lim \sup P^{*}\left(X_{\alpha} \in F\right) \leq L(F)$ for every closed set $F$;
(iv) $\liminf E_{*} f\left(X_{\alpha}\right) \geq \int f d L$ for every lower semicontinuous $f$ that is bounded below;
(v) $\quad \lim \sup E^{*} f\left(X_{\alpha}\right) \leq \int f d L$ for every upper semicontinuous $f$ that is bounded above;
(vi) $\lim P^{*}\left(X_{\alpha} \in B\right)=\lim P_{*}\left(X_{\alpha} \in B\right)=L(B)$ for every Borel set $B$ with $L(\partial B)=0$.
(vii) $\liminf E_{*} f\left(X_{\alpha}\right) \geq \int f d L$ for every bounded, Lipschitz continuous, nonnegative $f$.

Proof: (ii) and (iii) are equivalent by taking complements. Similarly, (ii) and (v) are equivalent by using $f$ by $-f$. The (i) $\Rightarrow$ (vii) is trivial.
(vii) $\Rightarrow$ (ii) Suppose (vii) holds. For every open $G$, consider a sequence of Lipchitz continuous functions $f_{m} \geq 0$ and $f_{m} \uparrow 1_{G}\left(f_{m}(x)=m d(x, D-G)>1\right.$, then

$$
\liminf P_{*}\left(X_{\alpha} \in G\right) \geq \liminf E_{*}\left(f\left(X_{\alpha}\right)\right) \geq \int f_{m} d L
$$

Letting $m \rightarrow \infty$, we get the result.
(ii) $\Rightarrow$ (iv) Let $f$ be lower semicontinuous with $f \geq 0$. Define

$$
f_{m}=\sum_{i=1}^{m^{2}}(1 / m) 1_{G_{i}}, G_{i}=\{x: f(x)>i / m\} .
$$

Then $f_{m} \leq f$ and $\left|f_{m}(x)-f(x)\right| \leq 1 / m$ when $x \leq m$. Fix $m$, use the fact $G_{i}$ is open for all $i$,

$$
\liminf E_{*} f\left(X_{\alpha}\right) \geq \liminf E_{*} f_{m}\left(X_{\alpha}\right) \geq \sum_{i=1}^{m^{2}} \frac{1}{m} P\left(X \in G_{i}\right)=\int f_{m} d L .
$$

Let $n \rightarrow \infty$, we get (iv) for non-negative lower semi continuous $f$. The conclusion follows.

Since a continuous function is both upper and lower continuous (iv) and (v) imply (i). We now prove (vi) is equivalent to others. Consider (ii) $\Rightarrow$ (vi). If (ii) and (iii) hold, then

$$
L(\operatorname{int} B) \leq \liminf P_{*}\left(X_{\alpha} \in \operatorname{int} B\right) \leq \lim \sup P^{*}\left(X_{\alpha} \in B\right) \leq L(\bar{B})
$$

If $L(\delta B)=0$, we get equalities giving (vi).
$(\mathrm{vi}) \Rightarrow$ (iii) Suppose (vi) holds and let $F$ be closed. Write $F^{\epsilon}=\{x: d(X, F)<\epsilon\}$. The sets $\partial F^{\epsilon}$ are disjoint for different $\epsilon>0$, giving at most countably many have non-zero $L$-measure. Choose $\epsilon_{m} \downarrow 0$ with $L\left(\partial F^{\epsilon_{m}}\right)=0$. For fixed $m$,

$$
\lim \sup P^{*}\left(X_{\alpha} \in F\right) \leq \lim \sup P^{*}\left(X_{\alpha} \in \bar{F}^{\varepsilon_{m}}\right)=L\left(\overline{F^{\epsilon_{m}}}\right)
$$

Let $m \rightarrow \infty$ to complete the proof.

Definition 6.2: The net of maps $X_{\alpha}$ is asymptotically measurable if and only if

$$
E^{*} f\left(X_{\alpha}\right)-E_{*} f\left(X_{\alpha}\right) \rightarrow 0, \quad \text { for every } f \in C_{b}(\mathbf{D}) .
$$

The net $X_{\alpha}$ is asymptotically tight if for every $\epsilon>0$ there exists a compact set $K$, such that

$$
\liminf P_{*}\left(X_{\alpha} \in K^{\delta}\right) \geq 1-\epsilon, \quad \text { for every } \delta>0
$$

Here $K^{\delta}=\{y \in \mathbf{D}: d(y, K)<\delta\}$ is the " $\delta$-enlargement" around $K$. A collection of Borel measurable maps $X_{\alpha}$ is called uniformly tight if, for every $\epsilon>0$, there is a compact $K$ with $P\left(X_{\alpha} \in K\right) \geq 1-\epsilon$ for every $\alpha$.

The $\delta$ in the definition of tightness may seem a bit overdone. Its non-asymptotic tightness as defined is essentially weaker than the same condition but with $K$ instead of $K^{\delta}$. This is caused by a second difference with the classical concept of uniform tightness: the enlarged compacts need to contain mass $1-\epsilon$ only in the limit.

Meanwhile, nothing is gained in simple cases: for Borel measurable maps in a Polish space, asymptotic tightness and uniform tightness are the same. It may also be noted that, although $K^{\delta}$ is dependent on the metric, the property of asymptotic tightness depends on the topology only. One nice consequence of the present tightness concept is that weak convergence usually implies asymptotic measurability and tightness.

Lemma 6.1: (i) $X_{\alpha} \Rightarrow X$, then $X_{\alpha}$ is asymptotically measurable.
(ii) If $X_{\alpha} \Rightarrow X$, then $X_{\alpha}$ is asymptotically tight if and only if $X$ is tight.

Proof: (i) This follows upon applying the definition of weak convergence to both $f$ and $-f$.
(ii) Fix $\epsilon>0$. If $X$ is tight, then there is a compact $K$ with $P(X \in K)>1-\epsilon$. By the Portmanteau theorem, $\lim \inf P_{*}\left(X_{\alpha} \in K^{\delta}\right) \geq P\left(X \in K^{\delta}\right)$, which is larger than $1-\epsilon$ for every $\delta>0$. Conversely, if $X_{\alpha}$ is tight, then there is a compact $K$ with $\liminf P_{*}\left(X_{\alpha} \in\right.$ $\left.K^{\delta}\right) \geq 1-\epsilon$. By the Portmanteau theorem, $P\left(X \in \overline{K^{\delta}}\right) \geq 1-\epsilon$. Let $\delta \downarrow 0$.

The next version of Prohorov's theorem may be considered a converse of the previous lemma. It comes in two parts, one for nets and one for sequences, neither one of which implies the other. The sequence case is the deepest of the two.

Theorem 6.2 (Prohorov's theorem): (i) If the net $X_{\alpha}$ is asymptotically tight and asymptotically measurable, then it has a subnet $X_{\alpha(\beta)}$ that converges in law to a tight Borel law.
(ii) If the sequence $X_{n}$ is asymptotically tight and asymptotically measurable, then it has a subsequence $X_{n_{j}}$ that converges weakly to a tight Borel law.

Proof: (i) Consider $\left(E^{*} f\left(X_{\alpha}\right)\right)_{f \in C_{b}(D)}$ as a net in product space

$$
\prod_{f \in C_{b}(D)}\left[-\|f\|_{\infty},\|f\|_{\infty}\right]
$$

which is compact in product topology. Hence, the net has convergent subnet. This implies there exist constants $L(f) \in\left[-\|f\|_{\infty},\|f\|_{\infty}\right]$ such that

$$
E^{*} f\left(X_{\alpha(\beta)}\right) \rightarrow L(f) \text { for } f \in C_{b}(D)
$$

In view of the asymptotic measurability, the numbers are the limits of corresponding $E_{*} f\left(X_{\alpha}\right)$. Now,

$$
\begin{aligned}
L\left(f_{1}+f_{2}\right) & \leq \lim \left(E^{*} f_{1}\left(X_{\alpha(\beta)}\right)+E^{*} f_{2}\left(X_{\alpha(\beta)}\right)\right) \\
& =L\left(f_{1}\right)+L\left(f_{2}\right) \\
& =\lim \left(E_{*} f_{1}\left(X_{\alpha(\beta)}\right)+E_{*} f_{2}\left(X_{\alpha(\beta)}\right)\right) \\
& \leq L\left(f_{1}+f_{2}\right) .
\end{aligned}
$$

Thus, $L: C_{b}(D) \rightarrow \mathbb{R}$ is a additive, and similarly $L(\lambda f)=\lambda L(f)$ for $\alpha \in \mathbb{R}$ and $L$ is positive. $L(f) \geq 0$ for $f \geq 0$. If $f_{m} \downarrow 0, L\left(f_{m}\right) \downarrow 0$. To see this, fix $\epsilon>0$. There is a compact set $K$ such that $\lim \inf P_{*}\left(X_{\alpha} \in K^{\delta}\right)>1-\epsilon$ for all $\delta>0$. Using the Dini theorem, for sufficiently large $m,\left|f_{m}(x)\right| \leq \epsilon$ for all $x \in K$. Using compactness of $K$, there exists $\delta>0$, such that $\left|f_{m}(x)\right| \leq 2 \epsilon$ for every $x \in K^{\delta}$. One gets $A_{\alpha} \subseteq\left\{X_{\alpha} \in K^{\delta}\right\}$ measurable $\left(1_{\left\{X_{\alpha} \in K^{\delta}\right\}}\right)_{*}=1_{A_{\alpha}}$. Hence,

$$
\begin{aligned}
L\left(f_{m}\right)= & \lim E f_{m}\left(X_{\alpha}\right)^{*}\left(\left(1_{\left\{X_{\alpha} \in K^{\delta}\right\}}\right)_{*}\right. \\
& +\left(1_{\left\{X_{\alpha} \notin K_{\delta}\right\}}\right)^{*} \leq 2 \epsilon+\left\|f_{1}\right\|_{\infty} \epsilon .
\end{aligned}
$$

Thus, $L$ is a measure.
(ii) For $m \in \mathbb{N}, K_{m}$ is a compact set with $\lim \inf P_{*}\left(X_{n} \in K_{m}^{\delta}\right) \geq 1-\frac{1}{m}$ for $\delta>0$. Since $K_{m}$ is compact, the space $C_{b}\left(K_{m}\right)$ and $\left\{f: f \in C_{b}\left(K_{m}\right),\|f(x)\| \leq 1\right.$, for $\left.x \in K_{m}\right\}$ are
separable. Using Tietze's theorem, every $f$ in the unit ball of $C_{b}\left(K_{m}\right)$ can be extended to an $f$ in the unit ball of $C_{b}(D)$. Hence, ball of $C_{b}(D)$ the restrictions of which to $K_{m}$ are dense in the unit ball of $C_{b}\left(K_{m}\right)$. Pick such a countable set for every $m$, and let $\Im$ be countably many functions obtained. For a fixed, bounded $f$, there is a subsequence $X_{n_{j}}$ such that $E^{*} f\left(X_{n_{j}}\right)$ converges to some number. Using diagonalization, one obtains a subsequence such that

$$
\bar{E}^{*} f\left(X_{n_{j}}\right) \rightarrow L(f) \text { for } f \in \Im .
$$

with $L(f) \in[-1,1]$.
Let $f \in C_{b}(D)$ taking values in $[-1,1]$. Fix $\epsilon>0$, and $m$. There exists $f_{m} \in \Im$ with $\left|f(x)-f_{m}(x)\right|<\epsilon$ for $x \in K_{m}$. Then as before there exists $\delta>0$ such that $\mid f(x)-$ $f_{m}(x) \mid \leq \epsilon$ for every $x \in K_{m}^{\delta}$. Then,

$$
\begin{aligned}
& \left|E^{*} f\left(X_{n}\right)-E^{*} f_{m}\left(X_{n}\right)\right| \leq E\left|f\left(X_{n}\right)-f_{m}\left(X_{n}\right)\right|^{*}\left(1_{\left\{X_{n} \in K_{m}^{\delta}\right\}}\right)_{*} \\
& +2 P^{*}\left(X_{n} \notin K_{m}^{\delta}\right) \leq 2 \epsilon+\frac{2}{m}
\end{aligned}
$$

for $n$ sufficiently large. Then $E^{*} f\left(X_{n_{j}}\right)$ has the property that, for $\eta>0$, there is a converging subsequence of numbers that is eventually with distance $\eta$. This gives convergence $E^{*} f\left(X_{n_{j}}\right)$ to a limit following as in the proof of (i) we can get the result.

Remark: Let $D_{0} \subseteq D$ and $X$ and $X_{\alpha}$ take values in $D_{0}$. Then $X_{\alpha} \rightarrow X$ as maps in $D_{0}$ iff $X_{\alpha} \rightarrow X$ as maps in $D$ if $D$ and $D_{0}$ are equipped with the same metric. This is easy from Portmaneteau theorem (ii) as each $G_{0} \subseteq D_{0}$ open is of the form $G \cap D_{0}$ with $G$ open in $D$.

### 6.1 Spaces of bounded functions

A vector lattice $\mathcal{F} \subset C_{b}(\mathbf{D})$ is a vector space that is closed under taking positive parts: if $f \in \mathcal{F}$, then $f=f \vee 0 \in \mathcal{F}$. Then automatically $f \vee g \in \mathcal{F}$ and $f \wedge g \in \mathcal{F}$ for every $f, g \in \mathcal{F}$. A set of functions on $\mathbf{D}$ separates points of $\mathbf{D}$ if, for every pair $x \neq y \in \mathbf{D}$, there is $f \in \mathcal{F}$ with $f(x) \neq f(y)$.

Lemma 6.2: Let $L_{1}$ and $L_{2}$ be finite Borel measures on $\mathbf{D}$.
(i) If $\int f d L_{1}=\int f d L_{2}$ for every $f \in C_{b}(\mathbf{D})$, then $L_{1}=L_{2}$.

Let $L_{1}$ and $L_{2}$ be tight Borel probability measures on $\mathbf{D}$.
(ii) If $\int f d L_{1}=\int f d L_{2}$ for every $f$ in a vector lattice $\mathcal{F} \subset C_{b}(\mathbf{D})$ that contains the constant functions and separates points of $\mathbf{D}$, then $L_{1}=L_{2}$.

Proof: (i) For every open set $G$, there exists a sequence of continuous functions with $0 \leq f_{m} \uparrow 1_{G}$. Using the monotone convergence theorem, $L_{1}(G)=L_{2}(G)$ for every open set $G$. Since $L_{1}(D)=L_{2}(D)$, the class of sets $\left\{A: L_{1}(A)=L_{2}(A)\right\}$ is a $\sigma$-field $\supseteq$ open sets.
(ii) Fix $\epsilon>$, take $K$ compact so that $L_{1}(K) L_{2}(K) \geq 1-\epsilon$. Note that by the Stone Weirstvass theorem, $\Im$ (containing constant and separating points of $K$ ) lattice $\subseteq$ $C_{b}(K)$ is uniformly dense in $C_{b}(k)$. Given $g \in C_{b}(D)$ with $0 \leq g \leq 1$, take $f \in \Im$, with $|g(x)-f(x)| \leq \epsilon$ for $x \in K$. Then

$$
\left|\int g d L_{1}-\int g d L_{2}\right| \leq \mid \int(f \wedge 1)^{+} d L_{1}-\inf (f \wedge 1)^{+} d L_{2}+\epsilon .
$$

This equals $4 \epsilon$ as $(f \wedge 1)^{+} \epsilon \Im$. Hence, one can equals $4 \epsilon$ as $(f \wedge 1)^{+} \epsilon \Im$. Hence, one can prove $\int g d L_{1}=\int g d L_{2}$ for all $g \in C_{b}(D)$.

Lemma 6.3: Let the net $X_{\alpha}$ be asymptotically tight, and suppose $E^{*} f\left(X_{\alpha}\right)-E_{*} f\left(X_{\alpha}\right) \rightarrow 0$ for every $f$ in a subalgebra $\mathcal{F}$ of $C_{b}(\mathbf{D})$ that separates points of $\mathbf{D}$. Then the net $X_{\alpha}$ is asymptotically measurable.

Proof: Fix $\epsilon>0$ and $K$ compact such that $\lim \sup P^{*}\left(X_{\alpha} \notin K^{\delta}\right) \leq \epsilon$ for all $\delta>0$. Assume without loss of generality that $\Im$ contains constant functions. Then restrictions of the function $\Im$ into $K$ are uniformly dense in $C_{b}(K)$ as before. Hence, given $f \in C_{b}(D)$, there exists $g \in \Im$ with $|f(x)-g(x)|<\frac{\epsilon}{4}$ for $x \in K$. As $K$ is compact one can show there exists a $\delta>0$ such that $|f(x)-g(x)|<\frac{\epsilon}{3}$ for $x \in K^{\delta}$. Let $\left\{X_{\alpha} \in K^{\delta}\right\}_{*}$ be a measurable set. Then $P\left(\Omega_{\alpha} \backslash\left\{X_{\alpha} \in K^{\delta}\right\}_{*}\right)=P^{*}\left\{X_{\alpha} \notin K^{\delta}\right)$ and for $\alpha$ large

$$
\left.P\left(\mid f\left(X_{\alpha}\right)^{*}-f\left(X_{\alpha}\right)\right)_{*} \mid>\epsilon\right\} .
$$

Let $T$ be an arbitrary set. The space $l^{\infty}(T)$ is defined as the set of all uniformly bounded, real functions on $T$ : all functions $z: T \rightarrow \mathbf{R}$ such that

$$
\|z\|_{T}:=\sup _{t \in T}|z(t)|<\infty
$$

It is a metric space with respect to the uniform distance $d\left(z_{1}, z_{2}\right)=\left\|z_{1}-z_{2}\right\|_{T}$.
The space $l^{\infty}(T)$, or a suitable subspace of it, is a natural space for stochastic processes with bounded sample paths. A stochastic process is simply an indexed collection $\{X(t): t \in T\}$ of random variables defined on the same probability space: every $X(t): \Omega \rightarrow \mathbf{R}$ is a measurable map. If every sample path $t \mapsto X(t, \omega)$ is bounded, then a stochastic process yields a map $X: \Omega \rightarrow l^{\infty}(T)$. Sometimes, the sample paths have additional properties, such as measurability or continuity, and it may be fruitful to consider $X$ as a map into a subspace of $l^{\infty}(T)$. If in either case the uniform metric is used, this does not make a difference for weak convergence of a net, but for measurability it can.

In most cases, a map $X: \Omega \rightarrow l^{\infty}(T)$ is a stochastic process. The small amount of measurability this gives may already be enough for asymptotic measurability. The special role played by the marginals $\left(X\left(t_{1}\right), \ldots, X\left(t_{k}\right)\right)$, which are considered as maps
into $\mathbf{R}^{k}$, is underlined by the following three results. Weak convergence in $l^{\infty}(T)$ can be characterized as asymptotic tightness plus convergence of marginals.

Lemma 6.4: Let $X_{\alpha}: \Omega_{\alpha} \rightarrow l^{\infty}(T)$ be asymptotically tight. Then it is asymptotically measurable if and only if $X_{\alpha}(t)$ is asymptotically measurable for every $t \in T$.

Lemma 6.5: Let $X$ and $Y$ be tight Borel measurable maps into $l^{\infty}(T)$. Then $X$ and $Y$ are equal in Borel law if and only if all corresponding marginals of $X$ and $Y$ are equal in law.

Theorem 6.3: Let $X_{\alpha}: \Omega_{\alpha} \rightarrow l^{\infty}(T)$ be arbitrary. Then $X_{\alpha}$ converges weakly to a tight limit if and only if $X_{\alpha}$ is asymptotically tight and the marginals $\left(X\left(t_{1}\right), \ldots, X\left(t_{k}\right)\right)$ converge weakly to a limit for every finite subset $t_{1}, \ldots, t_{k}$ of $T$. If $X_{\alpha}$ is asymptotically tight and its marginals converge weakly to the marginals $\left(X\left(t_{1}\right), \ldots, X\left(t_{k}\right)\right)$ of a stochastic process $X$, then there is a version of $X$ with uniformly bounded sample paths and $X_{\alpha} \Rightarrow X$.

Proof: For the proof of both lemmas, consider the collection $\mathcal{F}$ of all functions $f$ : $l^{\infty}(T) \rightarrow \mathbf{R}$ of the form

$$
f(z)=g\left(z\left(t_{1}\right), \ldots, z\left(t_{k}\right)\right), \quad g \in C_{b}\left(\mathbf{R}^{k}\right), t_{1}, \ldots, t_{k} \in T, k \in \mathbf{N} .
$$

This forms an algebra and a vector lattice, contains the constant functions, and separates points of $l^{\infty}(T)$. Therefore, the lemmas are corollaries of Lemmas 6.2 and 6.3, respectively. If $X_{\alpha}$ is asymptotically tight and marginals converge, then $X_{\alpha}$ is asymptotically measurable by the first lemma. By Prohorov's theorem, $X_{\alpha}$ is relatively compact $\|z\|_{T}=\sup _{t \epsilon T}|\mathbb{Z}(t)|<\infty$. To prove weak convergence, it suffices to show that all limit points are the same. This follows from marginal convergence and the second lemma.

Marginal convergence can be established by any of the well-known methods for proving weak convergence on Euclidean space. Tightness can be given a more concrete form, either through finite approximation or with the help of the ArzelàAscoli theorem. Finite approximation leads to the simpler of the two characterizations, but the second approach is perhaps of more interest, because it connects tightness to continuity of the sample paths $t \mapsto X_{\alpha}(t)$.

The idea of finite approximation is that for any $\epsilon>0$, the index set $T$ can be partitioned into finitely many subsets $T_{i}$ such that the variation of the sample paths $t \mapsto X_{\alpha}(t)$ is less than $\epsilon$ on every one of the sets $T_{i}$. More precisely, it is assumed that for every $\epsilon, \eta>0$, there exists a partition $T=\cup_{i=1}^{k} T_{i}$ such that

$$
\begin{equation*}
\limsup _{\alpha} P^{*}\left(\sup _{i} \sup _{s, t \in T_{i}}\left|X_{\alpha}(s)-X_{\alpha}(t)\right|>\epsilon\right)<\eta . \tag{6.1}
\end{equation*}
$$

Clearly, under this condition, the asymptotic behavior of the process can be described within error margin $\epsilon, \eta$ by the behavior of the marginal $\left(X_{\alpha}\left(t_{1}\right), \ldots, X_{\alpha}\left(t_{k}\right)\right)$ for arbitrary fixed points $t_{i} \in T_{i}$. If the process can thus be reduced to a finite set of coordinates for any $\epsilon, \eta>0$ and the nets or marginal distributions are tight, then the net $X_{\alpha}$ is asymptotically tight.

Theorem 6.4: A net $X_{\alpha}: \Omega_{\alpha} \rightarrow l^{\infty}(T)$ is asymptotically tight if and only if $X_{\alpha}(t)$ is asymptotically tight in $\mathbf{R}$ for every $t$ and, for all $\epsilon, \eta>0$, there exists a finite partition $T=\cup_{i=1}^{k} T_{i}$ such that (6.1) holds.

Proof: The necessity of the conditions follows easily from the next theorem. For instance, take the partition equal to disjointified balls of radius $\delta$ for a semi-metric on $T$ as in the next theorem. We prove sufficiency.

For any partition, as in the condition of the theorem, the norm $\left\|X_{\alpha}\right\|_{T}$ is bounded by $\max _{i}\left|X_{\alpha}\left(t_{i}\right)\right|+\epsilon$, with inner probability at least $1-\eta$, if $t_{i} \in T_{i}$ for each $i$. Since a maximum of finitely many tight nets of real variables is tight, it follows that the net $\left\|X_{\alpha}\right\|_{T}$ is asymptotically tight in $R$.

Fix $\zeta>0$ and a sequence $\epsilon_{m} \downarrow 0$. Take a constant $M$ such that $\lim \sup P^{*}\left(\left\|X_{\alpha}\right\|_{T}>\right.$ $M)<\zeta$, and for each $\epsilon=\epsilon_{m}$ and $\eta=2^{-m} \zeta$, take a partition $T=\cup_{i=1}^{k} T_{i}$ as in (61). For the moment, $m$ is fixed and we do not let it appear in the notation. Let $z_{1}, \ldots, z_{p}$ be the set of all functions in $l^{\infty}(T)$ that are constant on each $T_{i}$ and take on only the values $0, \pm \epsilon_{m}, \ldots, \pm\left[M / \epsilon_{m}\right] \epsilon_{m}$. Let $K_{m}$ be the union of the $p$ closed balls of radius $\epsilon_{m}$ around the $z_{i}$. Then, by construction, the two conditions

$$
\left\|X_{\alpha}\right\|_{T} \leq M \quad \text { and } \quad \sup _{i} \sup _{s, t \in T_{i}}\left|X_{\alpha}(s)-X_{\alpha}(t)\right| \leq \epsilon_{m}
$$

imply that $X_{\alpha} \in K_{m}$. This is true for each fixed $m$.
Let $K=\cap_{m=1}^{\infty} K_{m}$. Then $K$ is closed and totally bounded (by construction of the $K_{m}$ and because $\epsilon \downarrow 0$ ) and hence compact. Furthermore, for every $\delta>0$, there is an $m$ with $K^{\delta} \supset \cap_{m=1}^{m} K_{i}$. If not, then there would be a sequence $z_{m}$ not in $K^{\delta}$, but with $z_{m} \in \cap_{m=1}^{m} K_{i}$ for every $m$. This would have a subsequence contained in one of the balls making up $K_{1}$, a further subsequence eventually contained in one of the balls making up $K_{2}$, and so on. The "diagonal" sequence, formed by taking the first of the first subsequence, the second of the second subsequence, and so on would eventually be contained in a ball of radius $\epsilon_{m}$ for every $m$, hence Cauchy. Its limit would be in $K$, contradicting the fact that $d\left(z_{m}, K\right) \geq \delta$ for every $m$.

Conclude that if $X_{\alpha}$ is not in $K^{\delta}$, then it is not in $\cap_{m=1}^{m} K_{i}$ for some fixed $m$. Then

$$
\lim \sup P^{*}\left(X_{\alpha} \notin K^{\delta}\right) \leq \lim \sup P^{*}\left(X_{\alpha} \notin \cap_{m=1}^{m} K_{i}\right) \leq \zeta+\sum_{i=1}^{m} \zeta 2^{-m}<2 \zeta .
$$

This concludes the proof of the theorem.

The second type of characterization of asymptotic tightness is deeper and relates the concept to asymptotic continuity of the sample paths. Suppose $\rho$ is a semimetric on $T$ A net $X_{\alpha}: \Omega_{\alpha} \rightarrow l^{\infty}(T)$ is asymptotically uniformly $\rho$-equicontinuous in probability if for every $\epsilon, \eta>0$ there exists a $\delta>0$ such that

$$
\limsup _{\alpha} P^{*}\left(\sup _{\rho(s, t)<\delta}\left|X_{\alpha}(s)-X_{\alpha}(t)\right|>\epsilon\right)<\eta .
$$

Theorem 6.5: A net $X_{\alpha}: \Omega_{\alpha} \rightarrow l^{\infty}(T)$ is asymptotically tight if and only if $X_{\alpha}(t)$ is asymptotically tight in $R$ for every $t$ and there exists a semi-metric $\rho$ on $T$ such that ( $T, \rho$ ) is totally bounded and $X_{\alpha}$ is asymptotically uniformly $\rho$-equicontinuous in probability. If, moreover, $X_{\alpha} \Rightarrow X$, then almost all paths $t \mapsto X(t, \omega)$ are uniformly $\rho$ continuous; and the semi-metric $\rho$ can without loss of generality be taken equal to any semi-metric $\rho$ for which this is true and $(T, \rho)$ is totally bounded.

Proof: $(\Leftarrow)$. The sufficiency follows from the previous theorem. First, take $\delta>0$ sufficiently small so that the last displayed inequality is valid. Since $T$ is totally bounded, it can be covered with finitely many balls of radius $\delta$. Construct a partition of $T$ by disjointifying these balls.
$(\Rightarrow)$. If $X_{\alpha}$ is asymptotically tight, then $g\left(X_{\alpha}\right)$ is asymptotically tight for every continuous map $g$; in particular, for each coordinate projection. Let $K_{1} \subset K_{2} \subset \ldots$ be compacts with $\liminf P_{*}\left(X_{\alpha} \in K_{m}^{\epsilon}\right) 1-1 / m$ for every $\epsilon>0$. For every fixed $m$, define a semi-metric $\rho_{m}$ on $T$ by

$$
\rho_{m}(s, t)=\sup _{z \in K_{m}}|z(s)-z(t)|, \quad s, t \in T .
$$

Then ( $T, \rho_{m}$ ) is totally bounded. Indeed, cover $K_{m}$ by finitely many balls of radius $\eta$, centered at $z_{1}, \ldots, z_{k}$. Partition $R^{k}$ into cubes of edge $\eta$, and for every cube pick at most one $t \in T$ such that $\left(z_{1}(t), \ldots, z_{k}(t)\right)$ is in the cube. Since $z_{1}, \ldots, z_{k}$ are uniformly bounded, this gives finitely many points $t_{1}, \ldots, t_{p}$. Now the balls $\left\{t: \rho\left(t, t_{i}\right)<3 \eta\right\}$ cover $T: t$ is in the ball around $t_{i}$ for which $\left(z_{1}(t), \ldots, z_{k}(t)\right)$ and $\left(z_{1}\left(t_{i}\right), \ldots, z_{k}\left(t_{i}\right)\right)$ fall in the same cube. This follows because $\rho_{m}\left(t, t_{i}\right)$ can be bounded by $2 \sup _{z \in K_{m}} \inf _{i}\left\|z-z_{i}\right\|_{T}+$ $\sup _{j}\left|z_{j}\left(t_{i}\right)-z_{j}(t)\right|$. Next set

$$
\rho(s, t)=\sum_{m=1}^{\infty} 2^{-m}\left(\rho_{m}(s, t) \wedge 1\right) .
$$

Fix $\eta>0$. Take a natural number $m$ with $2^{-m}<\eta$. Cover $T$ with finitely many $\rho_{m}$ balls of radius $\eta$. Let $t_{1}, \ldots, t_{p}$ be their centers. Since $\rho_{1} \leq \rho_{2} \leq \ldots$, there is for every $t$ a $t_{i}$ with $\rho\left(t, t_{i}\right) \leq \sum_{k=1}^{m} 2^{-k} \rho_{k}\left(t, t_{i}\right)+2^{-m}<2 \eta$. Thus, $(T, \rho)$ is totally bounded for $\rho$, too. It is clear from the definitions that $|z(s)-z(t)| \leq \rho_{m}(s, t)$ for every $z \in K_{m}$ and
that $\rho_{m}(s, t) \wedge 1 \leq 2^{m} \rho(s, t)$. Also, if $\left\|z_{0}-z\right\|_{T}<\epsilon$ for $z \in K_{m}$ then $\left|z_{0}(s)-z_{0}(t)\right|<$ $2 \epsilon+|z(s)-z(t)|$ for any pair $s, t$. Deduce that

$$
K_{m}^{\epsilon} \subset\left\{z: \sup _{\rho(s, t)<2^{-m_{\epsilon}}}|z(s)-z(t)| \leq 3 \epsilon\right\} .
$$

Thus, for given $\epsilon$ and $m$, and for $\delta<2^{-m} \epsilon$,

$$
\liminf _{\alpha} P_{*}\left(\sup _{\rho(s, t)<\delta}\left|X_{\alpha}(s)-X_{\alpha}(t)\right| \leq 3 \epsilon\right) \geq 1-\frac{1}{m}
$$

Finally, if $X_{\alpha} \Rightarrow X$, then with notation as in the second part of the proof, $P\left(X \in K_{m}\right) \geq$ $1-1 / m$; hence, $X$ concentrates on $\cup_{m=1}^{\infty} K_{m}$. The elements of $K_{m}$ are uniformly $\rho_{m}{ }^{-}$ equicontinuous and hence also uniformly $\rho$-continuous. This yields the first statement. The set of uniformly continuous functions on a totally bounded, semi-metric space is complete and separable, so a map $X$ that takes its values in this set is tight. Next if $X_{\alpha} \Rightarrow X$ and $X$ is tight, the $X_{\alpha}$ is asymptotically tight and the compacts for asymptotical tightness can be chosen equal to the compacts for tightness of $X$. If $X$ has uniformly continuous paths, then the latter compacts can be chosen within the space of uniformly continuous functions. Since a compact is totally bounded, every one of the compacts is necessarily uniformly equicontinuous. The combination of these facts proves the second statement.

### 6.2 Maximal inequalities and covering numbers

We derive a class of maximal inequalities that can be used to establish the asymptotic equicontinuity of the empirical process. Since the inequalities have much wider applicability, we temporarily leave the empirical framework.

Let $\psi$ be a nondecreasing, convex function with $\psi(0)=0$ and $X$ a random variable. Then the Orlicz norm $\|X\|_{\psi}$ is defined as

$$
\|X\|_{\psi}=\inf \left\{C>0: E \psi\left(\frac{|X|}{C}\right) \leq 1\right\} .
$$

Here the infimum over the empty set is $\infty$. Using Jensen's inequality, it is not difficult to check that this indeed defines a norm. The best-known examples of Orlicz norms are those corresponding to the functions $x \mapsto x^{p}$ for $p \geq 1$ : the corresponding Orlicz norm is simply the $L_{p}$-norm

$$
\|X\|_{p}=\left(E|X|^{p}\right)^{1 / p}
$$

For our purpose, Orlicz norms of more interest are the ones given by $\psi_{p}(x)=e^{x^{p}}-1$ for $p \geq 1$, which give much more weight to the tails of $X$. The bound $x^{p} \leq \psi_{p}(x)$ for all nonnegative $x$ implies that $\|X\|_{p} \leq\|X\|_{\psi_{p}}$ for each $p$. It is not true that the exponential Orlicz norms are all bigger than all $L_{p}$-norms. However, we have the inequalities

$$
\begin{aligned}
\|X\|_{\psi_{p}} & \leq\|X\|_{\psi_{p}}(\log 2)^{1 / q-1 / p}, \quad p \leq q \\
\|X\|_{p} & \leq p!\|X\|_{\psi_{1}}
\end{aligned}
$$

Since for the present purposes fixed constants in inequalities are irrelevant, this means that a bound on an exponential Orlicz norm always gives a better result than a bound on an $L_{p}$-norm.

Any Orlicz norm can be used to obtain an estimate of the tail of a distribution. By Markov's inequality,

$$
P(|X|>x) \leq P\left(\psi\left(|X| /\|X\|_{\psi}\right) \geq \psi\left(x /\|X\|_{\psi}\right)\right) \leq \frac{1}{\psi\left(x /\|X\|_{\psi}\right)}
$$

For $\psi_{p}(x)=e^{x^{p}}-1$, this leads to tail estimates $\exp \left(-C x^{p}\right)$ for any random variable with a finite $\psi_{p}$-norm. Conversely, an exponential tail bound of this type shows that $\|X\|_{\psi_{p}}$ is finite.

Lemma 6.6: Let $X$ be a random variable with $P(|X|>x) \leq K e^{-C x^{p}}$ for every $x$, for constants $K$ and $C$, and for $p \geq 1$. Then its Orlicz norm satisfies $\|X\|_{\psi_{p}} \leq((1+K) / C)^{1 / p}$.

Proof: By Fubini's theorem,

$$
E\left(e^{D|X|^{p}}-1\right)=E \int_{0}^{|X|^{p}} D e^{D s} d s=\int_{0}^{\infty} P\left(|X|>s^{1 / p}\right) D e^{D s} d s
$$

Now insert the inequality on the tails of $|X|$ and obtain the explicit upper bound $K D /(C-D)$. This is less than or equal to 1 for $D^{-1 / p}$ greater than or equal to $((1+K) / C)^{1 / p}$. This completes the proof.

Next consider the $\psi$-norm of a maximum of finitely many random variables. Using the fact that max $\left|X_{i}\right|^{p} \leq \sum\left|X_{i}\right|^{p}$, one easily obtains for the $L_{p}$-norms

$$
\left\|\max _{1 \leq i \leq m} X_{i}\right\|_{p}=\left(E \max _{1 \leq i \leq m} X_{i}^{p}\right)^{1 / p} \leq m^{1 / p} \max _{1 \leq i \leq m}\left\|X_{i}\right\|_{p}
$$

A similar inequality is valid for many Orlicz norms, in particular the exponential ones. Here, in the general case, the factor $m^{1 / p}$ becomes $\psi^{-1}(m)$.

Lemma 6.7: Let $\psi$ be a convex, nondecreasing, nonzero function with $\psi(0)=0$ and $\lim \sup _{x, y \rightarrow \infty} \psi(x) \psi(y) / \psi(c x y)<\infty$ for some constant $c$. Then for any random variables $X_{1}, \ldots, X_{n}$,

$$
\left\|\max _{1 \leq i \leq m} X_{i}\right\|_{\psi}=\leq K \psi^{-1}(m) \max _{1 \leq i \leq m}\left\|X_{i}\right\|_{p}
$$

for a constant $K$ depending only on $\psi$.
Proof: For simplicity of notation, assume first that $\psi(x) \psi(y) \leq \psi(c x y)$ for all $x, y \geq 1$. In that case, $\psi(x / y) \leq \psi(c x) / \psi(y)$ for all $x \geq y \geq 1$. Thus, for $y \geq 1$ and any $C$,

$$
\begin{aligned}
\max \psi\left(\frac{\left|X_{i}\right|}{C y}\right) & \leq \max \left[\frac{\psi\left(c\left|X_{i}\right| / C\right)}{\psi(y)}+\psi\left(\frac{\left|X_{i}\right|}{C y}\right) \mathbf{1}\left\{\frac{\left|X_{i}\right|}{C y}<1\right\}\right] \\
& \leq \sum \frac{\psi\left(c\left|X_{i}\right| / C\right)}{\psi(y)}+\psi(1)
\end{aligned}
$$

Set $C=c \max \left\|X_{i}\right\|_{\psi}$, and take expectations to get

$$
E \psi\left(\frac{\max \left|X_{i}\right|}{C y}\right) \leq \frac{m}{\psi(y)}+\psi(1)
$$

When $\psi(1) \leq 1 / 2$, this is less than or equal to 1 for $y=\psi^{-1}(2 m)$, which is greater than 1 under the same condition. Thus,

$$
\left\|\max _{1 \leq i \leq m} X_{i}\right\|_{\psi} \leq \psi^{-1}(2 m) c \max \left\|X_{i}\right\|_{\psi}
$$

By the convexity of $\psi$ and the fact that $\psi(0)=0$, it follows that $\psi^{-1}(2 m) \leq 2 \psi^{-1}(m)$. The proof is complete for every special $\psi$ that meets the conditions made previously. For a general $\psi$, there are constants $\sigma \leq 1$ and $\tau>0$ such that $\phi(x)=\sigma \psi(\tau x)$ satisfies the conditions of the previous paragraph. Apply the inequality to $\phi$, and observe that $\|X\|_{\psi} \leq\|X\|_{\phi} /(\sigma \tau) \leq\|X\|_{\psi} / \sigma$.

For the present purposes, the value of the constant in the previous lemma is irrelevant. The important conclusion is that the inverse of the $\psi$-function determines the size of the $\psi$-norm of a maximum in comparison to the $\psi$-norms of the individual terms. The $\psi$-norms grows slowest for rapidly increasing $\psi$. For $\psi(x)=e^{x^{p}}-1$, the growth is at most logarithmic because

$$
\psi_{p}^{-1}(m)=(\log (1+m))^{1 / p}
$$

The previous lemma is useless in the case of a maximum over infinitely many variables. However, such a case can be handled via repeated application of the lemma via
a method known as chaining. Every random variable in the supremum is written as a sum of "little links," and the bound depends on the number and size of the little links needed. For a stochastic process $\left\{X_{t}: t \in T\right\}$, the number of links depends on the entropy of the index set for the semi-metric

$$
d(s, t)=\left\|X_{s}-X_{t}\right\|_{\psi} .
$$

The general definition of "metric entropy" is as follows.

Definition 6.3 (Covering numbers): Let ( $T, d$ ) be an arbitrary semi-metric space. Then the covering number $N(\epsilon, d)$ is the minimal number of balls of radius $\epsilon$ needed to cover $T$. Call a collection of points $\epsilon$-separated if the distance between each pair of points is strictly larger than $\epsilon$. The packing number $D(\epsilon, d)$ is the maximum number of $\epsilon$-separated points in $T$.

The corresponding entropy numbers are the logarithms of the covering and packing numbers, respectively.

For the present purposes, both covering and packing numbers can be used. In all arguments one can be replaced by the other through the inequalities

$$
N(\epsilon, d) \leq D(\epsilon, d) \leq N\left(\frac{\epsilon}{2}, d\right)
$$

Clearly, covering and packing numbers become bigger as $\epsilon \downarrow 0$. By definition, the semimetric space $T$ is totally bounded if and only if the covering and packing numbers are finite for every $\epsilon>0$. The upper bound in the following maximal inequality depends on the rate at which $D(\epsilon, d)$ grows as $\epsilon \downarrow 0$, as measured through an integral criterion.

Theorem 6.6: Let $\psi$ be a convex, nondecreasing, nonzero function with $\psi(0)=0$ and $\lim \sup _{x, y \rightarrow \infty} \psi(x) \psi(y) / \psi(c x y)<\infty$, for some constant $c$. Let $\left\{X_{t}: t \in T\right\}$ be a separable stochastic process with

$$
\left\|X_{s}-X_{t}\right\|_{\psi} \leq C d(s, t), \quad \text { for every } s, t
$$

for some semi-metric $d$ on $T$ and a constant $C$. Then, for any $\eta, \delta>0$,

$$
\left\|\sup _{d(s, t) \leq \delta}\left|X_{s}-X_{t}\right|\right\|_{\psi} \leq K\left[\int_{0}^{\eta} \psi^{-1}(D(\epsilon, d)) d \epsilon+\delta \psi^{-1}\left(D^{2}(\eta, d)\right)\right],
$$

for a constant $K$ depending on $\psi$ and $C$ only.

Corollary 6.1: The constant $K$ can be chosen such that

$$
\left\|\sup _{s, t}\left|X_{s}-X_{t}\right|\right\|_{\psi} \leq K \int_{0}^{\operatorname{diam} T} \psi^{-1}(D(\epsilon, d)) d \epsilon
$$

where diam $T$ is the diameter of $T$.

Proof: Assume without loss of generality that the packing numbers and the associated "covering integral" are finite. Construct nested sets $T_{0} \subset T_{1} \subset \cdots \subset T$ such that every $T_{j}$ is a maximal set of points such that $d(s, t)>\eta 2^{-j}$ for every $s, t \in T_{j}$, where "maximal" means that no point can be added without destroying the validity of the inequality. By the definition of packing numbers, the number of points in $T_{j}$ is less than or equal to $D\left(\eta 2^{-j}, d\right)$.
"Link" every point $t_{j+1} \in T_{j+1}$ to a unique $t_{j} \in T_{j}$ such that $d\left(t_{j}, t_{j+1}\right) \leq \eta 2^{-j}$. Thus, obtain for every $t_{k+1}$ a chain $t_{k+1}, t_{k}, \ldots, t_{0}$ that connects it to a point in $T_{0}$. For arbitrary points $s_{k+1}, t_{k+1}$ in $T_{k+1}$, the difference in increments along their chains can be bounded by

$$
\begin{aligned}
\left|\left(X_{s_{k+1}}-X_{s_{0}}\right)-\left(X_{t_{k+1}}-X_{t_{0}}\right)\right| & =\left|\sum_{j=0}^{k}\left(X_{s_{j+1}}-X_{s_{j}}\right)-\sum_{j=0}^{k}\left(X_{t_{j+1}}-X_{t_{j}}\right)\right| \\
& \leq 2 \sum_{j=0}^{k} \max \left|X_{u}-X_{v}\right|
\end{aligned}
$$

where for fixed $j$ the maximum is taken over all links $(u, v)$ from $T_{j+1}$ to $T_{j}$. Thus, the $j$ th maximum is taken over at most $\# T_{j+1}$ links, with each link having a $\psi$-norm $\left\|X_{u}-X_{\nu}\right\|_{\psi}$ bounded by $C d(u, v) \leq C \eta 2^{-j}$. It follows with the help of Lemma 6.7 that, for a constant depending only on $\psi$ and $C$,

$$
\begin{align*}
\left\|\max _{s, t \in T_{k+1}} \mid\left(X_{s}-X_{s_{0}}\right)-\left(X_{t}-X_{t_{0}}\right)\right\| \|_{\psi} & \leq K \sum_{j=0}^{k} \psi^{-1}\left(D\left(\eta 2^{-j-1}, d\right)\right) \eta 2^{-j} \\
& \leq 4 K \int_{0}^{\eta} \psi^{-1}(D(\epsilon, d)) d \epsilon . \tag{6.2}
\end{align*}
$$

In this bound, $s_{0}$ and $t_{0}$ are the endpoints of the chains starting at $s$ and $t$, respectively.

The maximum of the increments $\left|X_{s_{k+1}}-X_{t_{k+1}}\right|$ can be bounded by the maximum on the left side of (6.2) plus the maximum of the discrepancies $\left|X_{s_{0}}-X_{t_{0}}\right|$ at the end of the chains. The maximum of the latter discrepancies will be analyzed by a seemingly circular argument. For every pair of endpoints $s_{0}$, $t_{0}$ of chains starting at two points in
$T_{k+1}$ within distance $\delta$ of each other, choose exactly one pair $s_{k+1}, t_{k+1}$ in $T_{k+1}$, with $d\left(s_{k+1}, t_{k+1}\right)<\delta$, whose chains end at $s_{0}, t_{0}$. By definition of $T_{0}$, this gives at most $D^{2}(\eta, d)$ pairs. By the triangle inequality,

$$
\left|X_{s_{0}}-X_{t_{0}}\right| \leq\left|\left(X_{s_{0}}-X_{s_{k+1}}\right)-\left(X_{t_{0}}-X_{t_{k+1}}\right)\right|+\left|X_{s_{k+1}}-X_{t_{k+1}}\right| .
$$

Take the maximum over all pairs of endpoints $s_{0}, t_{0}$ as above. Then the corresponding maximum over the first term on the right in the last display is bounded by the maximum in the left side of (6.2). It $\psi$-norm can be bounded by the right side of this equation. Combine this with (6.2) to find that

$$
\begin{aligned}
& \left\|\max _{s, t \epsilon T_{k+1}, d(s, t)<\delta}\left|\left(X_{s}-X_{s_{0}}\right)-\left(X_{t}-X_{t_{0}}\right)\right|\right\|_{\psi} \\
& \leq 8 K \int_{0}^{\eta} \psi^{-1}(D(\epsilon, d)) d \epsilon+\left\|\max \left|X_{s_{k+1}}-X_{t_{k+1}}\right|\right\|_{\psi}
\end{aligned}
$$

Here the maximum on the right is taken over the pairs $s_{k+1}, t_{k+1}$ in $T_{k+1}$ uniquely attached to the pairs $s_{0}, t_{0}$ as above. Thus the maximum is over at most $D^{2}(\eta, d)$ terms, each of whose $\psi$-norm is bounded by $\delta$. Its $\psi$-norm is bounded by $K \psi^{-1}\left(D^{2}(\eta, d)\right) \delta$.

Thus, the upper bound given by the theorem is a bound for the maximum of increments over $T_{k+1}$. Let $k$ tend to infinity to conclude the proof. The corollary follows immediately from the previous proof, after noting that, for $\eta$ equal to the diameter of $T$, the set $T_{0}$ consists of exactly one point. In that case $s_{0}=t_{0}$ for every pair $s, t$, and the increments at the end of the chains are zero. The corollary also follows from the theorem upon taking $\eta=\delta=\operatorname{diam} T$ and noting that $D(\eta, d)=1$, so that the second term in the maximal inequality can also be written $\delta \psi^{-1}(D(\eta, d))$. Since the function $\epsilon \mapsto \psi^{-1}(D(\epsilon, d))$ is decreasing, this term can be absorbed into the integral, perhaps at the cost of increasing the constant $K$.

Although the theorem gives a bound on the continuity modulus of the process, a bound on the maximum of the process will be needed. Of course, for any $t_{0}$,

$$
\left\|\sup _{t}\left|X_{t}\right|\right\|_{\psi} \leq\left\|X_{t_{0}}\right\|_{\psi}+K \int_{0}^{\operatorname{diam} T} \psi^{-1}(D(\epsilon, d)) d \epsilon .
$$

Nevertheless, to state the maximal inequality in terms of the increments appears natural. The increment bound shows that the process $X$ is continuous in $\psi$-norm, whenever the covering integral $\int_{0}^{\eta} \psi^{-1}(D(\epsilon, d)) d \epsilon$ converges for some $\eta>0$. It is a small step to deduce the continuity of almost all sample paths from this inequality, but this is not needed at this point.

### 6.3 Sub-Gaussian inequalities

A standard normal variable has tails of the order $x^{-1} \exp \left(-\frac{x^{2}}{2}\right)$ and satisfies $P(|X|>x)$ $\leq 2 \exp \left(-\frac{x^{2}}{2}\right)$ for every $x$. In this section, we study random variables satisfying similar tail bounds.

Hoeffding's inequality asserts a "sub-Gaussian" tail bound for random variables of the form $X=\sum X_{i}$ with $X_{1}, \ldots, X_{n}$ i.i.d. with zero means and bounded range. The following special case of Hoeffding's inequality will be needed.

Theorem 6.7 (Hoeffding's inequality): Let $a_{1}, \ldots, a_{n}$ be constants and $\epsilon_{1}, \ldots, \epsilon_{n}$ be independent Rademacher random variables, i.e., with $P\left(\epsilon_{i}=1\right)=P\left(\epsilon_{i}=-1\right)=1 / 2$. Then

$$
P\left(\left|\sum \epsilon_{i} a_{i}\right|>x\right) \leq 2 e^{-\frac{x^{2}}{2\| \| \|^{2}}},
$$

for the Euclidean norm $\|a\|$. Consequently, $\left\|\sum \epsilon_{i} a_{i}\right\|_{\psi_{2}} \leq \sqrt{6}\|a\|$.
Proof: For any $\lambda$ and Rademacher variable $\epsilon$, one has $E e^{\lambda \epsilon}=\left(e^{\lambda}+e^{-\lambda}\right) \leq e^{\lambda^{2} / 2}$, where the last inequality follows after writing out the power series. Thus by Markov's inequality, for any $\lambda>0$,

$$
P\left(\sum \epsilon_{i} a_{i}>x\right) \leq e^{-\lambda x} E e^{\lambda \sum_{i=1}^{n} a_{i} \epsilon_{i}} \leq e^{\left(\frac{\lambda^{2}}{2}\|a\|^{2}-\lambda x\right)}
$$

The best upper bound is obtained for $\lambda=x /\|a\|^{2}$ and is the exponential in the probability bound of the lemma. Combination with a similar bound for the lower tail yields the probability bound. The bound on the $\psi$-norm is a consequence of the probability bound in view of Lemma 6.6.

A stochastic process is called sub-Gaussian with respect to the semi-metric $d$ on its index set if

$$
P\left(\left|X_{s}-X_{t}\right|>x\right) \leq 2 e^{-\frac{x^{2}}{2 d^{2}(s, t)}}, \quad \text { for every } s, t \in T, x>0
$$

any Gaussian process is sub-Gaussian for the standard deviation semi-metric $d(s, t)=$ $\sigma\left(X_{s}-X_{t}\right)$. Another example is Rademacher process

$$
X_{a}=\sum_{i=1}^{n} a_{i} \epsilon_{i}, \quad a \in R^{n}
$$

for Rademacher variables $\epsilon_{1}, \ldots, \epsilon_{n}$. By Hoeffding's inequality, this is sub-Gaussian for the Euclidean distance $d(a, b)=\|a-b\|$.

Sub-Gaussian processes satisfy the increment bound $\left\|X_{s}-X_{t}\right\|_{\psi_{2}} \leq \sqrt{6} d(s, t)$. Since the inverse of the $\psi_{2}$-function is essentially the square root of the logarithm, the general maximal inequality leads for sub-Gaussian processes to a bound in terms of an entropy integral. Furthermore, because of the special properties of the logarithm, the statement can be slightly simplified.

Corollary 6.2: Let $\left\{X_{t}: t \in T\right\}$ be a separable sub-Gaussian process. Then for every $\delta>0$,

$$
E \sup _{d(s, t) \leq \delta}\left|X_{s}-X_{t}\right| \leq K \int_{0}^{\delta} \sqrt{\log D(\epsilon, d)} d \epsilon,
$$

for a universal constant K . In particular, for any $t_{0}$,

$$
E \sup _{t}\left|X_{t}\right| \leq E\left|X_{t_{0}}\right|+K \int_{0}^{\infty} \sqrt{\log D(\epsilon, d)} d \epsilon .
$$

Proof: Apply the general maximal inequality with $\psi_{2}(x)=e^{x^{2}}-1$ and $\eta=\delta$. Since $\psi_{2}^{-1}(m)=\sqrt{\log (1+m)}$, we have $\psi_{2}^{-1}\left(D^{2}(\delta, d)\right) \leq \sqrt{2} \psi_{2}^{-1}(D(\delta, d))$. Thus, the second term in the maximal inequality can first be replaced by $\sqrt{2} \delta \psi^{-1}(D(\eta, d))$ and next be incorporated in the first at the cost of increasing the constant. We obtain

$$
\left\|\sup _{d(s, t) \leq \delta}\left|X_{s}-X_{t}\right|\right\|_{\psi_{2}} \leq K \int_{0}^{\delta} \sqrt{\log (1+D(\epsilon, d))} d \epsilon
$$

Here $D(\epsilon, d) \geq 2$ for every $\epsilon$ that is strictly less than the diameter of $T$. Since $\log (1+m) \leq$ $2 \log m$ for $m \geq 2$, the 1 inside the logarithm can be removed at the cost of increasing $K$.

### 6.4 Symmetrization

Let $\epsilon_{1}, \ldots, \epsilon_{n}$ be i.i.d. Rademacher random variables. Instead of the empirical process

$$
f \mapsto\left(P_{n}-P\right) f=\frac{1}{n} \sum_{i=1}^{n}\left(f\left(X_{i}\right)-P f\right),
$$

consider the symmetrized process

$$
f \mapsto P_{n}^{o} f=\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right),
$$

where $\epsilon_{1}, \ldots, \epsilon_{n}$ are independent of $\left(X_{1}, \ldots, X_{n}\right)$. Both processes have mean function zero. It turns out that the law of large numbers or the central limit theorem for one of these processes holds if and only if the corresponding result is true for the other
process. One main approach to proving empirical limit theorems is to pass from $P_{n}-P$ to $P_{n}^{o}$ and next apply arguments conditionally on the original $X$ 's. The idea is that, for fixed $X_{1}, \ldots, X_{n}$, the symmetrized empirical measure is a Rademacher process, hence a sub-Gaussian process, to which Corollary 6.2 can be applied.

Thus, we need to bound maxima and moduli of the process $P_{n}-P$ by those of the symmetrized process. To formulate such bounds, we must be careful about the possible nonmeasurability of suprema of the type $\left\|P_{n}-P\right\|_{\mathcal{F}}$. The result will be formulated in terms of outer expectation, but it does not hold for every choice of an underlying probability space on which $X_{1}, \ldots, X_{n}$ are defined. Throughout this part, if outer expectations are involved, it is assumed that $X_{1}, \ldots, X_{n}$ are the coordinate projections on the product space ( $\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}$ ), and the outer expectations of functions $\left(X_{1}, \ldots, X_{n}\right) \mapsto h\left(X_{1}, \ldots, X_{n}\right)$ are computed for $P^{n}$. thus "independent" is understood in terms of a product probability space. If auxiliary variables, independent of the $X$ 's, are involved, as in the next lemma, we use a similar convention. In that case, the underlying probability space is assumed to be of the form $\left(\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}\right) \times(\mathcal{Z}, \mathcal{C}, Q)$ with $X_{1}, \ldots, X_{n}$ equal to the coordinate projections on the first $n$ coordinates and the additional variables depending only on the $(n+1)$ st coordinate.

The following lemma will be used mostly with the choice $\Phi(x)=x$.
Lemma 6.8 (symmetrization): For every nondecreasing, convex $\Phi: R \rightarrow R$ and class of measurable functions $\mathcal{F}$,

$$
E^{*} \Phi\left(\left\|P_{n}-P\right\|_{\mathcal{F}}\right) \leq E^{*} \Phi\left(2\left\|P_{n}^{0}\right\|_{\mathcal{F}}\right)
$$

where the outer expectations are computed as indicated in the preceding paragraph.
Proof: Let $Y_{1}, \ldots, Y_{n}$ be independent copies of $X_{1}, \ldots, X_{n}$, defined formally as the coordinate projections on the last $n$ coordinates in the product space $\left(\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}\right) \times$ $(\mathcal{Z}, \mathcal{C}, Q) \times\left(\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}\right)$. The outer expectations in the statement of the lemma are unaffected by this enlargement of the underlying probability space, because coordinate projections are perfect maps. For fixed values $X_{1}, \ldots, X_{n}$,

$$
\begin{aligned}
\left\|P_{n}-P\right\|_{\mathcal{F}} & =\sup _{f \in \mathcal{F}} \frac{1}{n}\left|\sum_{i=1}^{n}\left(f\left(X_{i}\right)-E f\left(Y_{i}\right)\right)\right| \\
& \leq E_{Y}^{*} \sup _{f \in \mathcal{F}} \frac{1}{n}\left|\sum_{i=1}^{n}\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)\right|,
\end{aligned}
$$

where $E_{Y}^{*}$ is the outer expectation with respect to $Y_{1}, \ldots, Y_{n}$ computed for $P^{n}$ for given, fixed values of $X_{1}, \ldots, X_{n}$. Combination with Jensen's inequality yields

$$
\Phi\left(\left\|P_{n}-P\right\|_{\mathcal{F}}\right) \leq E_{Y} \Phi\left(\left\|\frac{1}{n} \sum_{i=1}^{n}\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)\right\|_{\mathcal{F}}^{* Y}\right),
$$

where $* Y$ denotes the minimal measurable majorant of the supremum with respect to $Y_{1}, \ldots, Y_{n}$, still with $X_{1}, \ldots, X_{n}$ fixed. Because $\Phi$ is nondecreasing and continuous, the $* Y$ inside $\Phi$ can be moved to $E_{Y}^{*}$. Next take the expectation with respect to $X_{1}, \ldots, X_{n}$ to get

$$
E^{*} \Phi\left(\left\|P_{n}-P\right\|_{\mathcal{F}}\right) \leq E_{X}^{*} E_{Y}^{*} \Phi\left(\left\|\frac{1}{n} \sum_{i=1}^{n}\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)\right\|_{\mathcal{F}}\right)
$$

Here the repeated outer expectation can be bounded above by the joint outer expectation $E^{*}$ by Fubini's theorem.

Adding a minus sign in front of a term $\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)$ has the effect of exchanging $X_{i}$ and $Y_{i}$. By construction of the underlying probability space as a product space, the outer expectation of any function $f\left(X_{1}, \ldots, X_{n}, Y_{1}, \ldots, Y_{n}\right)$ remains unchanged under permutations of its $2 n$ arguments, hence the expression

$$
E^{*} \Phi\left(\left\|\frac{1}{n} \sum_{i=1}^{n} e_{i}\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)\right\|_{\mathcal{F}}\right)
$$

is the same for any $n$-tuple $\left(e_{1}, \ldots, e_{n}\right) \in\{-1,1\}^{n}$. Deduce that

$$
E^{*} \Phi\left(\left\|P_{n}-P\right\|_{\mathcal{F}}\right) \leq E_{\epsilon} E_{X, Y}^{*} \Phi\left(\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i}\left(f\left(X_{i}\right)-f\left(Y_{i}\right)\right)\right\|_{\mathcal{F}}\right) .
$$

Use the triangle inequality to separate the contributions of the $X$ 's and the $Y$ 's and next use the convexity of $\Phi$ and triangle inequality to bound the previous expression by

$$
\frac{1}{2} E_{\epsilon} E_{X, Y}^{*} \Phi\left(2\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}}\right)+\frac{1}{2} E_{\epsilon} E_{X, Y}^{*} \Phi\left(2\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(Y_{i}\right)\right\|_{\mathcal{F}}\right) .
$$

By perfectness of coordinate projections, the expectation $E_{X, Y}^{*}$ is the same as $E_{X}^{*}$ and $E_{Y}^{*}$ in the two terms, respectively. Finally, replace the repeated outer expectations by a joint outer expectation. This completes the proof.

The symmetrization lemma is valid for any class $\mathcal{F}$. In the proofs of GlivenkoCantelli and Donsker theorems, it will be applied not only to the original set of functions of interest, but also to several classes constructed from such a set $\mathcal{F}$. The next step in these proofs is to apply a maximal inequality to the right side of the lemma, conditionally on $X_{1}, \ldots, X_{n}$. At that point, we need to write the joint outer expectation as the repeated expectation $E_{X}^{*} E_{\epsilon}$, where the indices $X$ and $\epsilon$ mean expectation over $X$ and $\epsilon$ conditionally on remaining variables. Unfortunately, Fubini's theorem is not valid for outer expectations. To overcome this problem, it is assumed that the integrand in the right side of the lemma is jointly measurable in ( $X_{1}, \ldots, X_{n}, \epsilon_{1}, \ldots, \epsilon_{n}$ ).

Since the Rademacher variables are discrete, this is the case if and only if the maps

$$
\begin{equation*}
\left(X_{1}, \ldots, X_{n}\right) \mapsto\left\|\sum_{i=1}^{n} e_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}} \tag{6.3}
\end{equation*}
$$

are measurable for every $n$-tuple $\left(e_{1}, \ldots, e_{n}\right) \in\{-1,1\}^{n}$. For the intended application of Fubini's theorem, it suffices that this is the case for the completion of $\left(\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}\right)$.

Definition 6.4 (measurable class): A class $\mathcal{F}$ of measurable functions $f: \mathcal{X} \rightarrow R$ on a probability space $(\mathcal{X}, \mathcal{A}, P)$ is called a $P$-measurable class if the function (6.3) is measurable on the completion of $\left(\mathcal{X}^{n}, \mathcal{A}^{n}, P^{n}\right)$ for every $n$ and every vector $\left(e_{1}, \ldots, e_{n}\right) \in R^{n}$.

### 6.4.1 Glivenko-Cantelli theorems

In this section, we prove two types of Glivenko-Cantelli theorems. The first theorem is the simplest and is based on entropy with bracketing. Its proof relies on finite approximation and the law of large numbers for real variables. The second theorem uses random $L_{1}$-entropy numbers and is proved through symmetrization followed by a maximal inequality.

Definition 6.5 (Covering numbers): The covering number $N(\epsilon, \mathcal{F},\|\cdot\|)$ is the minimal number of balls $\{g:\|g-f\|<\epsilon\}$ of radius $\epsilon$ needed to cover the set $\mathcal{F}$. The centers of the balls need not belong to $\mathcal{F}$, but they should have finite norms. The entropy (without bracketing) is the logarithm of the covering number.

Definition 6.6 (bracketing numbers): Given two functions $l$ and $u$, the bracket $[l, u]$ is the set of all functions $f$ with $l \leq f \leq u$. An $\epsilon$-bracket is a bracket $[l, u]$ with $\|u-l\|<\epsilon$. The bracketing number $N_{[]}(\epsilon, \mathcal{F},\|\cdot\|)$ is the minimum number of $\epsilon$-brackets needed to cover $\mathcal{F}$. The entropy with bracketing is the logarithm of the bracketing number. In the definition of the bracketing number, the upper and lower bounds $u$ and $l$ of the brackets need not belong to $\mathcal{F}$ themselves but are assumed to have finite norms.

Theorem 6.8: Let $\mathcal{F}$ be a class of measurable functions such that $N_{[\mathrm{l}}\left(\epsilon, \mathcal{F}, L_{1}(P)\right)<\infty$ for every $\epsilon>0$. Then $\mathcal{F}$ is Glivenko-Cantelli.

Proof: Fix $\epsilon>0$. Choose finitely many $\epsilon$-brackets $\left[l_{i}, u_{i}\right]$ whose union contains $\mathcal{F}$ and such that $P\left(u_{i}-l_{i}\right)<\epsilon$ for every $i$. Then for every $f \in \mathcal{F}$, there is a bracket such that

$$
\left(P_{n}-P\right) f \leq\left(P_{n}-P\right) u_{i}+P\left(u_{i}-f\right) \leq\left(P_{n}-P\right) u_{i}+\epsilon
$$

Consequently,

$$
\sup _{f \in \mathcal{F}}\left(P_{n}-P\right) f \leq \max _{i}\left(P_{n}-P\right) u_{i}+\epsilon .
$$

The right side converges almost surely to $\epsilon$ by the strong law of large numbers for real variables. Combination with a similar argument for $\inf _{f \in \mathcal{F}}\left(P_{n}-P\right) f$ yields that $\lim \sup \left\|P_{n}-P\right\|_{\mathcal{F}}^{*} \leq \epsilon$ almost surely, for every $\epsilon>0$. Take a sequence $\epsilon_{m} \downarrow$ 0 to see that the limsup must actually be zero almost surely. This completes the proof.

An envelope function of a class $\mathcal{F}$ is any function $x \mapsto F(x)$ such that $|f(x)| \leq F(x)$, for every $x$ and $f$. The minimal envelope function is $x \mapsto \sup _{f}|f(x)|$. It will usually be assumed that this function is finite for every $x$.

Theorem 6.9: Let $\mathcal{F}$ be a P-measurable class of measurable functions with envelope $F$ such that $P^{*} F<\infty$. Let $\mathcal{F}_{M}$ be the class of functions $f \mathbf{1}\{F \leq M\}$ when $f$ ranges over $\mathcal{F}$. If $\log N\left(\epsilon, \mathcal{F}_{M}, L_{1}\left(P_{n}\right)\right)=o_{P}^{*}(n)$ for every $\epsilon$ and $M>0$, then $\left\|P_{n}-P\right\|_{\mathcal{F}}^{*} \rightarrow 0$ both almost surely and in mean. In particular, $\mathcal{F}$ is Glivenko-Cantelli.

Proof: By the symmetrization lemma, measurability of the class $\mathcal{F}$, and Fubini's theorem,

$$
\begin{aligned}
E^{*}\left\|P_{n}-P\right\|_{\mathcal{F}} & \leq 2 E_{X} E_{\epsilon}\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}} \\
& \leq 2 E_{X} E_{\epsilon}\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}_{M}}+2 P^{*} F\{F>M\}
\end{aligned}
$$

by the triangle inequality, for every $M>0$. For sufficiently large $M$, the last term is arbitrarily small. To prove convergence in mean, it suffices to show that the first term converges to zero for fixed $M$. Fix $X_{1}, \ldots, X_{n}$. If $\mathcal{G}$ is an $\epsilon$-net in $L_{1}\left(P_{n}\right)$ over $\mathcal{F}_{M}$, then

$$
E_{\epsilon}\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}_{M}} \leq E_{\epsilon}\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{G}}+\epsilon .
$$

The cardinality of $\mathcal{G}$ can be chosen equal to $N\left(\epsilon, \mathcal{F}_{M}, L_{1}\left(P_{n}\right)\right)$. Bound the $L_{1}$-norm on the right by the Orlicz-norm for $\psi_{2}(x)=\exp \left(x^{2}\right)-1$, and use the maximal inequality Lemma 6.7 to find that the last expression does not exceed a multiple of

$$
\begin{equation*}
\sqrt{1+\log N\left(\epsilon, \mathcal{F}_{M}, L_{1}\left(P_{n}\right)\right)} \sup _{f \in \mathcal{G}}\left\|\frac{1}{n} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\psi_{2} \mid X}+\epsilon, \tag{6.4}
\end{equation*}
$$

where the Orlicz norms $\|\cdot\|_{\psi_{2} \mid X}$ are taken over $\epsilon_{1}, \ldots, \epsilon_{n}$ with $X_{1}, \ldots, X_{n}$ fixed. By Hoeffding's inequality, they can be bounded by $\sqrt{6 / n}\left(P_{n} f^{2}\right)^{1 / 2}$, which is less than $\sqrt{6 / n} M$. Thus, the last displayed expression is bounded by

$$
\sqrt{1+\log N\left(\epsilon, \mathcal{F}_{M}, L_{1}\left(P_{n}\right)\right)} \sqrt{\frac{6}{n}} M+\epsilon \rightarrow_{P^{*}} \epsilon
$$

It has been shown that the left side of (6.4) converges to zero in probability. Since it is bounded by $M$, its expectation with respect to $X_{1}, \ldots, X_{n}$ converges to zero by the dominated convergence theorem. This concludes the proof that $\left\|P_{n}-P\right\|_{\mathcal{F}}^{*}$ in mean. That it also converges almost surely follows from the fact that the sequence $\left\|P_{n}-P\right\|_{\mathcal{F}}^{*}$ is a reverse martingale with respect to a suitable filtration.

Fact: Let $\mathcal{F}$ be class of measurable functions with envelope $F$ such that $P^{*} F<\infty$. Define $\sum_{n}$ be the $\sigma$-field generated by measurable functions $h: X^{\infty} \rightarrow \mathbb{R}$ that are permutation symmetric in first arguments. Then

$$
E\left(\left\|P_{n}-P\right\|_{\mathcal{F}}^{*} \mid \Sigma_{n+1}\right) \geq\left\|P_{n+1}-P\right\|_{\mathcal{F}}^{*} \quad \text { a.s. }
$$

### 6.4.2 Donsker theorems

Uniform entropy: In this section, the weak convergence of the empirical process will be established under the condition that the envelope function $F$ be square integrable, combined with the uniform entropy bound

$$
\begin{equation*}
\int_{0}^{\infty} \sup _{Q} \sqrt{\log N\left(\epsilon| | F \|_{Q, 2}, \mathcal{F}, L_{2}(Q)\right)} d \epsilon<\infty . \tag{6.5}
\end{equation*}
$$

Here the supremum is taken over all finitely discrete probability measures $Q$ on $(\mathcal{X}, \mathcal{A})$ with $\|F\|_{Q, 2}^{2}=\int F^{2} d Q>0$. These conditions are by no means necessary, but they suffice for many examples. The finiteness of the previous integral will be referred to as the uniform entropy condition.

Theorem 6.10: Let $\mathcal{F}$ be a class of measurable functions that satisfies the uniform entropy bound (6.5). Let the class $\mathcal{F}_{\delta}=\left\{f-g: f, g \in \mathcal{F},\|f-g\|_{P, 2}<\delta\right\}$ and $\mathcal{F}_{\infty}^{2}$ be $P$-measurable for every $\delta>0$. If $P^{*} F^{2}<\infty$, then $\mathcal{F}$ is $P$-Donsker.

Proof: Let $\delta_{n} \downarrow 0$ be arbitrary. By Markov's inequality and the symmetrization lemma,

$$
P^{*}\left(\left\|G_{n}\right\|_{\mathcal{F}_{\delta_{n}}}>x\right) \leq \frac{2}{x} E^{*}\left\|\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}_{\delta_{n}}} .
$$

Since the supremum in the right-hand side is measurable by assumption, Fubini's theorem applies and the outer expectation can be calculated as $E_{X} E_{\epsilon}$. Fix $X_{1}, \ldots, X_{n}$. By Hoeffding's inequality, the stochastic process $f \mapsto\left\{n^{-1 / 2} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\}$ is sub-Gaussian for the $L_{2}\left(P_{n}\right)$-seminorm

$$
\|f\|_{n}=\sqrt{\frac{1}{n} \sum_{i=1}^{n} f^{2}\left(X_{i}\right)}
$$

Use the second part of the maximal inequality Corollary 6.2 to find that

$$
\left.E_{\epsilon}\left\|\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \epsilon_{i} f\left(X_{i}\right)\right\|_{\mathcal{F}_{\delta_{n}}} \leq \int_{0}^{\infty} \sqrt{\log N\left(\epsilon, \mathcal{F}_{\delta_{n}}, L_{2}\left(P_{n}\right)\right.}\right) d \epsilon
$$

For large values of $\epsilon$, the set $\mathcal{F}_{\delta_{n}}$ fits in a single ball of radius $\epsilon$ around the origin, in which case the integrand is zero. This is certainly the case for values of $\epsilon$ larger than $\theta_{n}$, where

$$
\theta_{n}^{2}=\sup _{f \in \mathcal{F}_{\delta_{n}}}\|f\|_{n}^{2}=\left\|\frac{1}{n} \sum_{i=1}^{n} f^{2}\left(X_{i}\right)\right\|_{\mathcal{F}_{\delta_{n}}}
$$

Furthermore, covering numbers of the class $\mathcal{F}_{\delta}$ are bounded by covering numbers of $\mathcal{F}_{\infty}=\{f-g: f, g \in \mathcal{F}\}$. The latter satisfy $N\left(\epsilon, \mathcal{F}_{\infty}, L_{2}(Q)\right) \leq N^{2}\left(\epsilon / 2, \mathcal{F}, L_{2}(Q)\right)$ for every measure $Q$.

Limit the integral in (6.5) to the interval $\left(0, \theta_{n}\right)$, make a change of variables, and bound the integrand to obtain the bound

$$
\int_{0}^{\theta_{n} /\|F\|_{n}} \sup _{Q} \sqrt{\log N\left(\epsilon\|F\|_{Q, 2}, \mathcal{F}, L_{2}(Q)\right)} d \epsilon\|F\|_{n} .
$$

Here the supremum is taken over all discrete probability measures. The integrand is integrable by assumption. Furthermore, $\|F\|_{n}$ is bounded below by $\left\|F_{*}\right\|_{n}$, which converges almost surely to its expectation, which may be assumed positive. Use the Cauch-Schwarz inequality and the dominated convergence theorem to see that the expectation of this integral converges to zero provided $\theta_{n} \rightarrow_{p^{*}} 0$. This would conclude the proof of asymptotic equicontinuity.

Since $\sup \left\{P f^{2}: f \in \mathcal{F}_{\delta_{n}}\right\} \rightarrow 0$ and $\mathcal{F}_{\delta_{n}} \subset \mathcal{F}_{\infty}$, it is certainly enough to prove that

$$
\left\|P_{n} f^{2}-P f^{2}\right\|_{\mathcal{F}_{\infty}} \rightarrow_{P^{*}} 0
$$

This is a uniform law of large numbers for the class $\mathcal{F}_{\infty}^{2}$. This class has integrable envelope $(2 F)^{2}$ and is measurable by assumption. For any pair $f, g$ of functions in $\mathcal{F}_{\infty}$,

$$
P_{n}\left|f^{2}-g^{2}\right| \leq P_{n}|f-g| 4 F \leq\|f-g\|_{n}\|4 F\|_{n} .
$$

It follows that the covering number $N\left(\epsilon\|2 F\|_{n}^{2}, \mathcal{F}_{\infty}^{2}, L_{1}\left(P_{n}\right)\right)$ is bounded by the covering number $N\left(\epsilon\|F\|_{n}, \mathcal{F}_{\infty}, L_{2}\left(P_{n}\right)\right)$. By assumption, the latter number is bounded by a fixed number, so its logarithm is certainly $o_{P}^{*}(n)$, as required for the uniform law of large numbers, Theorem 6.8. This concludes the proof of asymptotic equicontinuity.

Finally, we show that $\mathcal{F}$ is totally bounded in $L_{2}(P)$. By the result of the last paragraph, there exists a sequence of discrete measures $P_{n}$ with $\left\|\left(P_{n}-P\right) f^{2}\right\|_{\mathcal{F}_{\infty}}$ converging to zero. Take $n$ sufficiently large so that the supremum is bounded by $\epsilon^{2}$. by assumption, $N\left(\epsilon, \mathcal{F}, L_{2}\left(P_{n}\right)\right)$ is finite. Any $\epsilon$-net for $\mathcal{F}$ in $L_{2}\left(P_{n}\right)$ is a $\sqrt{2} \epsilon$-net in $L_{2}(P)$. This completes the proof.

### 6.5 Lindberg-type theorem and its applications

In this section, we want to show how the methods developed earlier of symmetrization and entropy bounds can be used to prove a limit theorem for sums of independent stochastic processes $Z_{n i}, f \in \mathcal{F}$ with bounded sample paths indexed by an arbitrary set $\mathcal{F}$. We need some preliminarily results for this purpose.

Lemma 6.5.1: Let $Z_{1}, Z_{2}, \ldots Z_{n}$ be independent stochastic processes with mean 0 . With $\left\|\|_{\mathcal{F}}\right.$ denoting the sup norm on $\mathcal{F}$, we get

$$
E^{*} \Phi\left(\frac{1}{2}\left\|\sum_{i=1}^{n} \epsilon_{i} Z_{i}\right\|_{\mathcal{F}}\right) \leq E^{*} \Phi\left(\left\|\sum_{i=1}^{n} Z_{i}\right\|_{\mathcal{F}}\right) \leq E^{*} \Phi\left(2\left\|\sum_{i=1}^{n} \epsilon_{i}\left(Z_{i}-\mu_{i}\right)\right\|_{\mathcal{F}}\right)
$$

for every non-decreasing, convex $\Phi: \mathbb{R} \rightarrow \mathbb{R}$ and arbitrary functions $\mu_{i}: \mathcal{F} \rightarrow \mathbb{R}$.
Proof: The inequality on the right can be proved using techniques as in the proof of symmetrization Lemma and is left to the reader. For the inequality on the left, let $Y_{1}, Y_{2}, \ldots Y_{n}$ be independent copy of $Z_{1}, \ldots Z_{n}$ defined on $\left(\prod_{i=1}^{n}\left(\mathcal{X}_{i}, a_{i}, P_{i}\right) x(\mathcal{Z}, \mathcal{G}, \mathcal{Q}) x \prod_{1}^{n}\left(\mathcal{X}_{i}, a_{i}, P_{i}\right)\right.$ (the $\epsilon_{i}$ 's are defined on $(\mathcal{Z}, \mathcal{G}, \mathcal{Q})$ ) and depend on the last n co-ordinates exactly as $Z_{1}, \ldots Z_{n}$ depend on the first $n$ coordinates. Since $E Y_{i}(f)=0$, the LHS of the above expression is the average of

$$
E^{*} \Phi\left(\| \frac{1}{2} \sum_{i=1}^{n} \epsilon_{i}\left[Z_{i}(f)-E Y_{i}(f) \|_{\mathcal{F}}\right)\right.
$$

where ( $e_{1} \ldots e_{n}$ ) range over $\{-1,1\}^{n}$
By Jenssen's inequality

$$
\begin{array}{r}
E_{Z, Y}^{*} \Phi\left(\| \frac{1}{2} \sum_{i=1}^{n} e_{i}\left[Z_{i}(f)-Y_{i}(f) \|_{\mathcal{F}}\right)\right. \\
\leq E_{Z, Y}^{*} \Phi\left(\left\|\frac{1}{2} \sum_{i=1}^{n}\left[Z_{i}(f)-Y_{i}(f)\right]\right\|_{\mathcal{F}}\right) .
\end{array}
$$

Apply triangle inequality and convexity of $\Phi$ to complete the proof.

Lemma 6.5.2: For arbitrary stochastic processes $Z_{1}, \ldots Z_{n}$, arbitrary functions $\mu_{1}, \ldots \mu_{n}$ : $\mathcal{F} \rightarrow \mathbb{R}$, and $x>0$

$$
\begin{aligned}
& \beta_{n}(x) P^{*}\left(\left\|\sum_{i=1}^{n} Z_{i}\right\|_{\mathcal{F}}>x\right) \\
& \leq 2 P^{*}\left(4\left\|\sum_{i=1}^{n} \epsilon_{i}\left(Z_{i}-\mu_{i}\right)\right\|_{\mathcal{F}}>x\right)
\end{aligned}
$$

and $\beta_{n}(x) \leq \inf _{f} P\left(1 \sum_{i=1}^{n} Z_{i}(f) \mid<x / 2\right)$. In particular, this is true for i.i.d. mean 0 processes and $\beta_{n}(x)=1-4 \frac{n}{x^{2}} \sup _{f} \operatorname{var}\left(Z_{1}(f)\right)$.

Proof: Let $Y_{1}, \ldots Y_{n}$ be independent copy of $Z_{1}, Z_{2}, \ldots Z_{n}$ as defined above. If $\left\|\sum_{1}^{n} Z_{i}\right\|_{\mathcal{F}}>$ $x$, then there is some $f$ for which

$$
\left|\sum_{1}^{n} Z_{i}(f)\right|>x .
$$

Fix a realization of $Z_{1}, \ldots Z_{n}$ and corresponding $f$. For this realization,

$$
\begin{aligned}
\beta & \leq P_{Y}^{*}\left(1 \sum_{i=1}^{n} Y_{i}(f) \mid<x / 2\right) \\
& \leq P_{Y}^{*}\left(1 \sum_{i=1}^{n} Y_{i}(f)-\sum_{i=1}^{1} Z_{i}(f) \mid>x / 2\right) \\
& \leq P_{Y}^{*}\left(\left\|\sum_{i=1}^{n}\left(Y_{i}-Z_{i}\right)\right\|_{\mathcal{F}}>x / 2\right)
\end{aligned}
$$

The far left and far right sides do not depend on a particular $f$, and the inequality between them is valid on the set

$$
\left\{\left\|\sum_{1}^{n} Z_{i}\right\|_{\mathcal{F}}>x\right\} .
$$

Integrate both sides with respect to $Z_{1}, \ldots Z_{n}$ over the set to obtain

$$
\beta P^{*}\left(\left\|\sum_{i=1}^{n} Z_{i}\right\|>x\right) \leq P_{Z}^{*} P_{Y}^{*}\left(\| \sum_{i=1}^{n}\left(Y_{i}-Z_{i}\right) \mid>x / 2\right) .
$$

By symmetry, the RHS equals

$$
E_{\epsilon} P_{Z}^{*} P_{Y}^{*}\left(\left\|\sum_{i=1}^{n} \epsilon_{i}\left(Y_{i}-Z_{i}\right)\right\|_{\mathcal{F}}>x / 2\right) .
$$

In view of triangle inequality, this is less than or equal to

$$
2 P^{*}\left(\left\|\sum_{1}^{n} \epsilon_{i}\left(Y_{i}-\mu_{i}\right)\right\|_{\mathcal{F}}>x / 4\right) .
$$

Assuming i.i.d. $Z_{1}, \ldots Z_{n}$, one can derive the inequality from Chebyshev inequality.

Proposition 6.5.3: Let $0<p<\infty$ and $X_{1}, \ldots X_{n}$ be independent stochastic processes indexed by $T$. Then there exist constants $C_{p}$ and $0<\mu_{p}<1$ such that

$$
E^{*} \max _{k \leq n}\left\|S_{k}\right\|^{p} \leq C_{p}\left(E^{*} \max _{k \leq n}\left\|X_{k}\right\|^{p}+F^{-1}\left(\mu_{p}\right)^{P}\right),
$$

where $F^{-1}$ is quantile function of $\max _{k \leq n}\left\|S_{k}\right\|^{*}$. If $X_{1}, \ldots X_{n}$ are symmetric, then there exist $K_{p}$ and $\left(0<v_{p}<1\right)$ such that

$$
E^{*}\left\|S_{n}\right\|^{p} \leq K_{p}\left(E^{*} \max _{k \leq n}\left\|x_{k}\right\|^{p}+G^{-1}\left(v_{p}\right)^{p}\right)
$$

where $G^{-1}$ is quantile function of $\left\|S_{n}\right\|^{*}$. For each $p \leq 1$, the last inequality has a version for mean zero processes.

Proof: This is based on the following fact due to Höffman-Jørgensen.
Fact: Let $X_{1}, \ldots X_{n}$ be independent stochastic processes indexed by an arbitrary set. Then for $\lambda, \eta>0$
(1) $P^{*}\left(\max _{k \leq n}\left\|S_{k}\right\|>3 \lambda+\eta\right) \leq P^{*}\left(\max _{i S K \leq n}\left\|S_{k}\right\|>\lambda\right)^{2}+P^{*}\left(\max _{k \leq n}\left\|x_{k}\right\|>\eta\right)$.

If $X_{1}, \ldots X_{n}$ are independent symmetric, then

$$
\text { (2) } P^{*}\left(\left\|S_{n}\right\|>2 \lambda+\eta\right) \leq 4 P^{*}\left(\left\|S_{n}\right\|>\lambda\right)^{2}+P^{*}\left(\max _{k \leq n}\left\|x_{k}\right\|>\eta\right) \text {. }
$$

## Proof of Proposition:

Take $\lambda=\eta=t$ in the above inequality (1) to find that for $t>0$

$$
\begin{align*}
E^{*} \max _{k \leq n}\left\|S_{k}\right\|^{p} & \leq 4^{p} \int P^{*}\left(\max _{k \leq n}\left\|S_{k}\right\|>4 t\right) d\left(t^{p}\right) \\
& \leq(4 t)^{p}+4^{p} \int_{t}^{\infty} P^{*}\left(\left\|S_{k}\right\|>t\right)^{2} d\left(t^{p}\right)+4^{p} \int_{t}^{\infty} P^{*}\left(\max _{k \leq n}\left\|X_{k}\right\|^{*}>t d\left(t^{p}\right)\right. \\
& \leq(4 t)^{p}+4^{p} P^{*}\left(\max _{k \leq n}\left\|S_{k}\right\|>t\right) E^{*} \max _{k \leq n}\left\|S_{n}\right\|^{p}+4^{p} E^{*} \max _{k \leq n}\left\|X_{k}\right\|^{p} . \tag{6.6}
\end{align*}
$$

Choose $t$ such that $4^{p} P^{*}\left(\max _{k \leq n}\left\|S_{k}\right\|>t\right)<1 / 2$ we get the inequality. Similar arguments using inequality (2) gives the second inequality.

The inequality for mean-zero processes follows using symmetrization using the inequality above. Then one can using desymmetrization as by Jensen inequality $E^{*}\left\|S_{n}\right\|^{p}$ is bounded by $E^{*}\left(\left\|S_{n}-T_{n}\right\|^{p}\right)$ if $T_{n}$ is sum of $n$ independent copies of $X_{1}, \ldots X_{n}$.

For each $n$, let $Z_{1}, \ldots Z_{n}, m_{n}$ be independent stochastic processes indexed by a common semi-metric space $(\mathcal{F}, \mathcal{P})$. These processes are defined on a product probability space of dimension $m_{n}$. Define random semi-metric

$$
d_{n}^{2}(f, g)=\sum_{i=1}^{m_{n}}\left(Z_{n i}(f)-Z_{n i}(g)\right)^{2} .
$$

The condition below in Theorem 6.5.1 is random entropy condition.

Theorem 6.5.1 (Lindberg): For each $n$, let $Z_{1}, \ldots Z_{n}, m_{n}$ be independent stochastic processes indexed by $(\mathcal{F}, \mathcal{P})$ (a totally bounded semi-metric space).

Assume
(a) $\sum_{1}^{m_{n}} E^{*}\left(\left\|Z_{n i}\right\|_{\mathcal{F}}^{2}\left\{\left\|Z_{n i}\right\|_{\mathcal{F}}>\eta\right) \rightarrow 0\right.$ for each $\eta>0$
(b) For each $(f, g) \epsilon, \mathcal{F} \otimes \mathcal{F}$
$\sum_{1}^{m_{n}} E\left(Z_{n i}(f), Z_{n i}(g)\right) \rightarrow C(f, g)$ finite
(c) $\sup _{p(f, g)<\delta_{n}} \sum_{1}^{n} E\left(Z_{n i}(f)-Z_{n i}(g)\right)^{2} \rightarrow$ for $\delta_{n} \downarrow 0$
(d) $\int_{0}^{\delta_{n}} \sqrt{\log N(\epsilon, f, d n)} d \epsilon \xrightarrow{P^{*}} 0$ for $\delta_{n} \downarrow 0$

Then the sequence $\sum_{i=1}^{n}\left(Z_{n i}-E Z_{n i}\right)$ is asymptotically $p$-equicontinuous. It converges in distrubtion in $l^{\infty}(\mathcal{F})$.

Proof: The condition (a) implies Lindberg condition for $\left\{Z_{n i}(f), i=1,2 \ldots m_{n}\right\}$ and using (b) we get marginal distributions of sum converge to a Gaussian process.

Set $Z_{n i}^{0}=Z_{n i}-E Z_{n i}$ and let $\delta_{n} \downarrow 0$ arbitrary. For fixed $t$ and $n$ large, Chebychev inequality and (c) give the bound

$$
P\left(\left|\sum_{i=1}^{n}\left(Z_{n i}^{0}(f)-Z_{n i}^{0}(g)\right)\right|>t / 2\right)<1 / 2
$$

for $f, g$ with $P(f, g)<\delta_{n}$. By Lemma 6.8, for sufficiently large $n$,

$$
P^{*}\left(\sup _{p(f, g)<\delta_{n}}\left|\sum_{1}^{m_{n}}\left(Z_{n i}^{0}(f)-Z_{n i}^{0}(g)\right)\right|>t\right)
$$

$$
\leq 4 P\left(\sup _{p(f, g)<\delta_{n}}\left|\sum_{1}^{m_{n}} \epsilon_{i}\left(Z_{n i}(f)-Z_{n i}(g)\right)\right|>t / 4\right)
$$

For fixed values of processes $Z_{1}, \ldots Z_{n}, m_{n}$ define $A \leq \mathbb{R}^{m_{n}}$ as a set of all vectors.

$$
\left(Z_{n i}(f)-Z_{n i}(g), \ldots Z_{n m_{n}}(f)-Z_{n m_{n}}(g)\right)
$$

when $f, g$ are in the set $\left\{(f, g) \epsilon \mathcal{F} x \mathcal{F}, P(f, g)<\delta_{n}\right\}$. By Höffding inequality (Theorem 6.7), the stochastic process $\left\{\sum \epsilon_{i} a_{i}, a \epsilon A\right\}$ is Gaussian for Euclidean metric on $A_{n}$. By Corollary 6.9

$$
\begin{aligned}
& E_{\epsilon} \sup _{p(f, g)<\delta_{n}}\left|\sum_{i=1}^{m_{n}} \epsilon_{i}\left(Z_{n i}(f)-Z_{n i}(g)\right)\right| \\
& \quad \leq \int_{0}^{\infty} \sqrt{\log N\left(\epsilon, A_{n},\| \|\right)} d \epsilon
\end{aligned}
$$

where \|\| is the Euclidean norm. The integrand can be bounded using $N\left(\epsilon, A_{n},\| \|\right) \leq$ $N^{2}(\epsilon / 2, \mathcal{F}, d n)$. If, in addition,

$$
\Theta_{n}^{2}=\sup _{a \in A_{n}} \sum_{i=1}^{n} a_{i}^{2}=\left\|\sum_{i=1}^{m_{n}} a_{i}^{2}\right\|_{A_{n}},
$$

then for $\epsilon>\Theta_{n}$, the set $A_{n}$ fits in the ball of radius $\epsilon$ around the origin and the integrand vanishes. From entropy condition (d) and this, we conclude that the integral converges to zero in outer probability if $\Theta_{n} \rightarrow 0$ in probability. Under the measurability assumption, this gives equicontinuity of $\sum_{1}^{m_{n}} Z_{n i}^{0}$.

By the Lindberg condition, there exist a sequence $\eta_{n} \downarrow 0$ such that

$$
E^{*}\left\|\sum a_{i}^{2}\left\{\left\|Z_{n i}\right\|>\eta_{n}\right\}\right\|_{A_{n}} \rightarrow 0
$$

Thus, without loss of generality, we can assume $\left\|Z_{n i}\right\|_{\mathcal{F}} \leq \eta_{n}$ to show $\Theta_{n} \xrightarrow{P^{*}} 0$.
Fix $Z_{n 1}, \ldots Z_{n m_{n}}$ and take $\epsilon$-net $b_{n}$ for $A_{n}$ with Euclidean norm. For every $a \epsilon A_{n}$, there is a $b \epsilon B_{n}$ with

$$
\left|\sum_{i}^{n} \epsilon_{i} a_{i}^{2}\right|=\left|\sum \epsilon_{i}\left(a_{i}-b_{i}\right)^{2}+2 \sum\left(a_{i}-b_{i}\right) \epsilon b_{i}+\sum \epsilon_{i} b_{i}^{2}\right| \leq \epsilon^{2}+2 \epsilon| | b \|+\left|\sum \epsilon_{i} b_{i}^{2}\right|
$$

By Höffding inequality, $\sum \epsilon_{i} b_{i}^{2}$ has orlicz norm for $\Psi_{2}$ bounded by a multiple of $\left(\sum, b_{i}^{4}\right)^{1 / 2} \leq \eta_{n}\left(\sum_{i}^{n} b_{i}^{2}\right)^{1 / 2}$. Apply Lemma 6.7 to the third term on the right and substitute $\sup A_{n}$ for sup over $B_{n}$.

$$
E_{\epsilon}\left\|\sum \epsilon_{i} a_{i}^{2}\right\| \leq \epsilon^{2}+2 \epsilon\left\|\sum a_{i}^{2}\right\|_{A_{n}}^{y_{2}}+\sqrt{1+\log \left|B_{n}\right|} \eta_{n}\left\|\sum a_{i}^{2}\right\|_{A_{n}}^{y_{2}} .
$$

The size of the $\epsilon$-net can be chosen

$$
\left|B_{n}\right| \leq N^{2}(\epsilon / 2, \mathcal{F}, d n)
$$

by entropy condition $d$, this variable is bounded in probability.
Conclude for some constant, $K$

$$
P\left(\left\|\sum \epsilon_{i} a_{i}^{2}\right\|_{A_{n}}>t\right) \leq P^{*}\left(B_{n} \mid>M\right)+K / t\left[\epsilon^{2}+\left(\epsilon+\eta_{n} \sqrt{\log M}\right) x E\left\|\sum a_{i}^{2}\right\|_{A_{n}}^{1 / 2}\right] .
$$

For $M$, $t$ sufficiently large, the RHS is smaller than $1-v$, for the constant $v$, in Proposition 6.5.3. More precisely, this can be achieved for $M$ such that

$$
P^{*}\left(B_{n} \mid>M\right) \leq\left(1-v_{1}\right) / 2
$$

and $t=\left(1-v_{1}\right) / 2$ and $t$ is equal to $\left(1-v_{1}\right)$ times the numerator of second term. Then $t$ is bigger than $v$-quantile of $\left\|\sum \epsilon_{i} a_{i}^{2}\right\|_{A_{n}}$. Thus, Proposition 6.5.3 gives

$$
E\left\|\sum \epsilon_{i} a_{i}^{2}\right\|_{A_{n}} \leqslant E\left\|\max a_{i}^{2}\right\|_{A_{n}}+t \leqslant \eta_{n}^{2}+\epsilon^{2}+\left(\epsilon+\eta_{n} \sqrt{\log M}\right)\left(E\left\|\sum a_{i}^{2}\right\|_{A_{n}}\right)^{(1 / 2)} .
$$

Now (c) gives \| $\sum E a_{i}^{2} \| A_{n} \rightarrow 0$. Combining this with Lemma 6.5.1, we get

$$
E\left\|\sum a_{i}^{2}\right\|_{A_{n}} \leq E\left\|\sum \epsilon_{i} a_{i}^{2}\right\|_{A_{n}}+0(1) \leq \delta+\delta\left(E\left\|\sum a_{i}^{2}\right\|_{A_{n}}\right)
$$

for $\delta>\max \left(\epsilon^{2}, \epsilon\right)$ and sufficiently large $n$.
Now for $c \geq 0, c \leq \delta+\delta \sqrt{c}$ implies

$$
c \leq\left(\delta+\sqrt{\left.\delta^{2}+4 \delta\right)^{2}} .\right.
$$

Applying this to $c=E\left\|\sum a_{i}^{2}\right\|_{A_{n}}$ to conclude $E\left\|\sum a_{i}^{2}\right\|_{A_{n}} \rightarrow 0$ as $n \rightarrow \infty$.
Example 1: One can look for $X_{1} \ldots X_{n}$ i.i.d. random variables

$$
Z_{n i}(f)=n^{(-1 / 2)} f\left(X_{i}\right)
$$

with $f$ measurable functions on the sample space. From the above theorem, we get that $\mathcal{F}$ is Donsker if it is measurable and totally bounded in $L_{2}(P)$, possesses a square integrable envelop and satisfies with $P_{n}$ equal to distribution of $X_{1} \ldots X_{n}$

$$
\int_{0}^{\delta_{n}} \sqrt{N\left(\epsilon, \mathcal{F}, L_{2}\left(P_{n}\right)\right)} d \epsilon \xrightarrow{P} 0 .
$$

If $\mathcal{F}$ satisfies the uniform entropy condition in section 6.5, then it satisfies the above condition. Thus, Theorem 6.10 is the consequence of the above theorem.

Example 2: The Lindberg condition on norms is not necessary for the CLT. In combination with truncation the preceding theorem applies to more general processes. Choose stochastic processes $Z_{1} \ldots Z_{n}, m_{n}$ with

$$
\sum_{1}^{m_{n}} P\left(\left\|Z_{n i}\right\|_{\mathcal{F}}>\eta\right) \rightarrow 0 \text { for } \eta>0
$$

Then truncated processes $Z_{n i}, \eta(f)=Z_{n i}(f) 1\left\{\left\|Z_{n i}\right\| \leq \eta\right\}$ satisfy

$$
\sum_{1}^{m_{n}} Z_{n i}-\sum_{1}^{m_{n}} Z_{n i}, \eta \xrightarrow{P} 0 \text { in } l^{\infty}(\mathcal{F})
$$

Since this is true for every $\eta>0$, it is also true for $\eta_{n} \downarrow 0$ sufficiently slowly. The processes $Z_{n, i, \eta_{n}}$ satisfy Lindberg condition. If they or their centered version satisfy other conditions of Theorem 6.5.1, then the sequence

$$
\sum_{1}^{m_{n}}\left(Z_{n i}-E Z_{n i, \eta_{n}}\right)
$$

converges weakly in $l^{\infty}(\mathcal{F})$. The random semi-metrics $d_{n}$ decrease by truncation. Thus, one can get the result for truncated processes under weaker conditions.

We now define measure-like processes. In the previous theorem, consider index set $\mathcal{F}$ as a set of measurable functions $f: x \rightarrow \mathbb{R}$ on a measurable space $(x, a)$ and the distance be in $L_{2}(Q) . Q$ finite measure: assume

$$
\text { (1) } \int_{0}^{\infty} \sup _{Q} \sqrt{\log N\left(\epsilon\|F\|_{Q, 2} \mathcal{F} L_{2}(Q)\right)} d \epsilon<\infty \text {. }
$$

Then the preceding theorem yields a central limit theorem for processes with increments that are bounded by random $L_{2}$-metric on $\mathcal{F}$. We call $Z_{n i}$ measure-like with respect to random measures $\mu_{n i}: f$

$$
\left(Z_{n i}(f)-Z_{n i}(g)\right) \leq^{2} \int(f-g)^{2} d \mu_{n i} f, g, \epsilon \mathcal{F}
$$

For measure-like processes $d_{n}$ is bounded by $L_{2}\left(\sum \mu_{n i}\right)$-semi-metric and entropy condition there can be related to uniform entropy condition.

Lemma 6.5.3: Let $\mathcal{F}$ be a class of measurable function $F$. Let $Z_{1}, \ldots Z_{n}, m_{n}$ be measurable processes indexed by $\mathcal{F}$. If $\mathcal{F}$ satisfies uniform entropy condition (1) for a set $Q$ that
contain measures $\mu_{n i}$ and $\sum_{1}^{m_{n}} \int F^{2} d \mu_{n i}=0_{p}^{*}(1)$ then entropy condition d) of Theorem 6.5 .1 is satisfied.

Proof: Let $\mu_{n}=\sum_{i=1}^{m_{n}} \mu_{n i}$. Since $d_{n}$ is bounded by $L_{2}\left(\mu_{n}\right)$-semi-metric we get

$$
(*) \int_{0}^{\delta_{n}} \sqrt{\log N\left(\epsilon, \mathcal{F}, d_{n}\right)} d \epsilon \leq \int_{0}^{\delta /\|F\|_{\mu n}} \sqrt{\log N\left(\epsilon\|F\|_{\mu n}, \mathcal{F}, L_{2}\left(\mu_{n}\right)\right)} d \epsilon .\|F\|_{\mu n}
$$

on the set where $\|F\|_{\mu_{n}}^{2}<\infty$. Denote by

$$
J(\delta)=\int_{0}^{\delta} \sup _{Q} \sqrt{\log N\left(\epsilon\|F\|_{Q, 2,} \mathcal{F}, L_{2}(Q)\right.} d \epsilon
$$

on the set where $\|F\|_{\mu_{n}}>\eta$ the RHS of the equation (*) is bounded by $J\left(\delta_{n} / \eta\right) 0_{P}(1)$ which converges to zero for $\eta>0$. On the set where $\|F\|_{\mu_{n}} \leq \eta$, we have the bound $J(\infty) \eta$ as $\eta$ is arbitrary we get the result.

Example 1: Suppose we have independent processes $Z_{1}, \ldots Z_{n}, m_{n}$ are measure line. Let $\mathcal{F}$ satisfy uniform entropy condition. Assume for some probability measure $P$ with $\int F^{2} d P^{*}<\infty$

$$
E^{*} \sum_{1}^{m_{n}} \int F_{n i}^{2} 1\left\{\int F^{2} d \mu_{n i}>\eta\right\} \rightarrow 0 \text { for } n \rightarrow 0
$$

and

$$
\sup _{\|f-g\|_{p, 2}<\delta_{n}} E^{*} \sum_{1}^{m_{n}} \int(f-g)^{2} d \mu_{n i} \rightarrow 0 \text { as } \delta_{n} \downarrow 0
$$

and

$$
\int F^{2} d \mu_{n i}=0_{p}^{*}(1)
$$

then we get that condition of Theorem 6.5.1 (d) is satisfied and other conditions imply the other conditions as

$$
\left\|Z_{n i}\right\|_{\mathcal{F}} \leq Z_{n i}(f) 1+4 \int F^{n} d \mu_{n i} \text { for any } f
$$

This gives $\sum_{1}^{m_{n}}\left(Z_{n i}-E Z_{n i}\right)$ converges in $l^{\infty}(F)$ if covariance functions converge point-wise and measureability conditions as satisfied.

Example 2: Consider in the above example $Z_{n i}=c_{n i} \delta_{x_{n i}}$ where for each $n$, ( $c_{n i}$ constant) $X_{1}, \ldots X_{n}, m_{n}$ are independent random variables is a measurable space ( $x, a$ ) with induced measures $P_{n i}$ with $\int f d P_{n i}$ exist for each element $f$ of a class of measurable functions $f: x \rightarrow \mathbb{R}$. We can consider weighted empirical process

$$
G_{n}(f)=\sum_{i=1}^{m_{n}} c_{n i}\left(f\left(X_{n i}\right)-\int f d P_{n i}\right) .
$$

Suppose $\max _{1 \leq i \leq m_{n}}\left|c_{n i}\right| \rightarrow 0$

$$
\sum_{1}^{m_{n}} c_{n i} P_{n i} \leq P
$$

for $P$ a probability measure $P$ with $E_{p}^{*} F^{2}<\infty$. Note that $Z_{i i}$ are measure-like with $\mu_{n i}=c_{n i}^{2} \delta_{x_{n i}}$. Then we get under measurability condition $G_{n}$ converges weakly to a Gaussian process on $l^{\infty}(\mathcal{F})$. In addition, we can prove that the limiting process has uniformly continuous sample paths with respect to $L_{2}(P)$-semi-metric.

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