Stochastics II

Lecture notes

Prof. Dr. Evgeny Spodarev
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1 General theory of random functions

1.1 Random functions

Let $(\Omega, \mathcal{A}, P)$ be a probability space and $(\mathcal{S}, \mathcal{B})$ a measurable space, $\Omega, \mathcal{S} \neq \emptyset$.

**Definition 1.1.1**

A random element $X : \Omega \rightarrow \mathcal{S}$ is a $\mathcal{A}|\mathcal{B}$-measurable mapping (Notation: $X \in \mathcal{A}|\mathcal{B}$), i.e.,

$$X^{-1}(B) = \{\omega \in \Omega : X(\omega) \in B\} \in \mathcal{A}, \quad B \in \mathcal{B}.$$  

If $X$ is a random element, then $X(\omega)$ is a realization of $X$ for arbitrary $\omega \in \Omega$.

We say that the $\sigma$-algebra $\mathcal{B}$ of subsets of $\mathcal{S}$ is induced by the set system $\mathcal{M}$ (Elements of $\mathcal{M}$ are also subsets of $\mathcal{S}$), if

$$\mathcal{B} = \sigma(\mathcal{F}) \supseteq \mathcal{F} \quad (\text{Notation: } \mathcal{B} = \sigma(\mathcal{M})).$$

If $\mathcal{S}$ is a topological or metric space, then $\mathcal{M}$ is often chosen as a class of all open sets of $\mathcal{S}$ and $\sigma(\mathcal{M})$ is called the Borel $\sigma$-algebra (Notation: $\mathcal{B} = \mathcal{B}(\mathcal{S})$).

**Example 1.1.1**

1. If $\mathcal{S} = \mathbb{R}$, $\mathcal{B} = \mathcal{B}(\mathbb{R})$, then a random element $X$ is called a random variable.

2. If $\mathcal{S} = \mathbb{R}^m$, $\mathcal{B} = \mathcal{B}(\mathbb{R}^m)$, $m > 1$, then $X$ is called a random vector. Random variables and random vectors are considered in the lectures „Elementare Wahrscheinlichkeitsrechnung und Statistik“ and „Stochastik I“.

3. Let $\mathcal{S}$ be the class of all closed sets of $\mathbb{R}^m$. Let

$$\mathcal{M} = \{\{A \in \mathcal{S} : A \cap B \neq \emptyset\} : B - \text{arbitrary compactum of } \mathbb{R}^m\}.$$  

Then $X : \Omega \rightarrow \mathcal{S}$ is a random closed set.

As an example we consider $n$ independent uniformly distributed points $Y_1, \ldots, Y_n \in [0, 1]^m$ and $R_1, \ldots, R_n > 0$ (almost surely) independent random variables, which are defined on the same probability space $(\Omega, \mathcal{A}, P)$ as $Y_1, \ldots, Y_n$. Consider $X = \bigcup_{i=1}^n B_{R_i}(Y_i)$, where $B_r(x) = \{y \in \mathbb{R}^m : ||y - x|| \leq r\}$. Obviously, this is a random closed set. An example of a realization is provided in Figure 1.1.

**Exercise 1.1.1**

Let $(\Omega, \mathcal{A})$ and $(\mathcal{S}, \mathcal{B})$ be measurable spaces, $\mathcal{B} = \sigma(\mathcal{M})$, where $\mathcal{M}$ is a class of subsets of $\mathcal{S}$. Prove that $X : \Omega \rightarrow \mathcal{S}$ is $\mathcal{A}|\mathcal{B}$-measurable if and only if $X^{-1}(C) \in \mathcal{A}$, $C \in \mathcal{M}$.

**Definition 1.1.2**

Let $T$ be an arbitrary index set and $(\mathcal{S}_t, \mathcal{B}_t)_{t \in T}$ a family of measurable spaces. A family $X = \{X(t), t \in T\}$ of random elements $X(t) : \Omega \rightarrow \mathcal{S}_t$ defined on $(\Omega, \mathcal{A}, P)$ and $\mathcal{A}|\mathcal{B}_t$-measurable for all $t \in T$ is called a random function (associated with $(\mathcal{S}_t, \mathcal{B}_t)_{t \in T}$).
Therefore it holds \( X : \Omega \times T \to (S_t, t \in T) \), i.e. \( X(\omega, t) \in S_t \) for all \( \omega \in \Omega, t \in T \) and \( X(\cdot, t) \in \mathcal{A}|B_t, t \in T \). We often omit \( \omega \) in the notation and write \( X(t) \) instead of \( X(\omega, t) \).

Sometimes \( (S_t, B_t) \) does not depend on \( t \in T \) as well: \( (S_t, B_t) = (S, B) \) for all \( t \in T \).

Special cases of random functions:

1. \( T \subseteq \mathbb{Z} : X \) is called a random sequence or stochastic process in discrete time.
   
   Example: \( T = \mathbb{Z}, \mathbb{N} \).

2. \( T \subseteq \mathbb{R} : X \) is called a stochastic process in continuous time.
   
   Example: \( T = \mathbb{R}_+, [a,b], -\infty < a < b < \infty, \mathbb{R} \).

3. \( T \subseteq \mathbb{R}^d, d \geq 2 : X \) is called a random field.
   
   Example: \( T = \mathbb{Z}^d, \mathbb{R}^d_+, [a,b]^d \).

4. \( T \subseteq \mathcal{B}(\mathbb{R}^d) : X \) is called a set-indexed process.

   If \( X(\cdot) \) is almost surely non-negative and \( \sigma \)-additive on the \( \sigma \)-algebra \( T \), then \( X \) is called a random measure.

The tradition of denoting the index set with \( T \) comes from the interpretation of \( t \in T \) for the cases 1 and 2 as time parameter.

For every \( \omega \in \Omega \), \( \{X(\omega, t), t \in T\} \) is called a trajectory or path of the random function \( X \).

We would like to prove that the random function \( X = \{X(t), t \in T\} \) is a random element within the corresponding function space, which is equipped with a \( \sigma \)-algebra that now is specified.

Let \( \mathcal{S}_T = \prod_{t \in T} \mathcal{S}_t \) be the cartesian product of \( \mathcal{S}_t, t \in T \), i.e., \( x \in \mathcal{S}_T \) if \( x(t) \in \mathcal{S}_t, t \in T \). The elementary cylindric set in \( \mathcal{S}_T \) is defined as

\[
C_T(B, t) = \{ x \in \mathcal{S}_T : x(t) \in B \},
\]

where \( t \in T \) is a selected point from \( T \) and \( B \in \mathcal{B}_t \) a subset of \( \mathcal{S}_t \). \( C_T(B, t) \) therefore contains all trajectories \( x \), which go through the „gate” \( B \), see Figure 1.2.

Definition 1.1.3

The \textit{cylindric \( \sigma \)-algebra} \( \mathcal{B}_T \) is introduced as a \( \sigma \)-algebra induced in \( \mathcal{S}_T \) by the family of all
elementary cylinders. It is denoted by $B_T = \bigotimes_{t \in T} B_t$. If $B_t = B$ for all $t \in T$, then $B^T$ is written instead of $B_T$.

**Lemma 1.1.1**

The family $X = \{X(t), t \in T\}$ is a random function on $(\Omega, A, P)$ with phase spaces $(S_t, B_t)_{t \in T}$ if and only if for every $\omega \in \Omega$ the mapping $\omega \mapsto X(\omega, \cdot)$ is $A|B_T$-measurable.

**Exercise 1.1.2**

Proof Lemma 1.1.1.

**Definition 1.1.4**

Let $X$ be a random element $X : \Omega \to S$, i.e. $X$ be $A|B$-measurable. The distribution of $X$ is the probability measure $P_X$ on $(S, B)$ such that $P_X(B) = P(X^{-1}(B))$, $B \in B$.

**Lemma 1.1.2**

An arbitrary probability measure $\mu$ on $(S, B)$ can be considered as the distribution of a random element $X$.

**Proof** Take $\Omega = S$, $A = B$, $P = \mu$ and $X(\omega) = \omega$, $\omega \in \Omega$.

When does a random function with given properties exist? A random function, which consists of independent random elements always exists. This assertion is known as

**Theorem 1.1.1 (Lomnicki, Ulam):**

Let $(S_t, B_t, \mu_t)_{t \in T}$ be a sequence of probability spaces. It exists a random sequence $X = \{X(t), t \in T\}$ on a probability space $(\Omega, A, P)$ (associated with $(S_t, B_t)_{t \in T}$) such that

1. $X(t), t \in T$ are independent random elements.
2. $P_{X(t)} = \mu_t$ on $(S_t, B_t)$, $t \in T$.

A lot of important classes of random processes is built on the basis of independent random elements; cf. examples in Section 1.2.

**Definition 1.1.5**

Let $X = \{X(t), t \in T\}$ be a random function on $(\Omega, A, P)$ with phase space $(S_t, B_t)_{t \in T}$. The finite-dimensional distributions of $X$ are defined as the distribution law $P_{t_1, \ldots, t_n}$ of $(X(t_1), \ldots, X(t_n))^T$ on $(S_{t_1, \ldots, t_n}, B_{t_1, \ldots, t_n})$, for arbitrary $n \in \mathbb{N}$, $t_1, \ldots, t_n \in T$, where $S_{t_1, \ldots, t_n} = S_{t_1} \times \ldots \times S_{t_n}$ and

![Abb. 1.2: Trajectories which pass a „gate“ $B_t$.](image)
Let $X$ be a random function with given finite-dimensional distributions. Then there exists a random function $X(\omega, t)$ such that

\[ P_{t_1, \ldots, t_n}(B_1 \times \ldots \times B_n) = P(X(t_1) \in B_1, \ldots, X(t_n) \in B_n). \]

**Exercise 1.1.3**
Prove that $X_{t_1, \ldots, t_n} = (X(t_1), \ldots, X(t_n))^T$ is a $\mathcal{A} | \mathcal{B}_{t_1, \ldots, t_n}$-measurable random element.

**Definition 1.1.6**
Let $\mathcal{S}_T = \mathbb{R}$ for all $t \in T$. The random function $X = \{X(t), t \in T\}$ is called symmetric, if all of its finite-dimensional distributions are symmetric probability measures, i.e., $P_{t_1, \ldots, t_n}(A) = P_{t_1, \ldots, t_n}(-A)$ for $A \in \mathcal{B}_{t_1, \ldots, t_n}$ and all $n \in \mathbb{N}$, $t_1, \ldots, t_n \in T$, whereby

\[ P_{t_1, \ldots, t_n}(-A) = P((-X(t_1), \ldots, -X(t_n))^T \in A). \]

**Exercise 1.1.4**
Prove that the finite-dimensional distributions of a random function $X$ have the following properties: for arbitrary $n \in \mathbb{N}$, $n \geq 2$, $\{t_1, \ldots, t_n\} \subset T$, $B_k \in \mathcal{S}_{T_k}$, $k = 1, \ldots, n$ and an arbitrary permutation $(i_1, \ldots, i_n)$ of $(1, \ldots, n)$ it holds:

1. **Symmetry**: $P_{t_{i_1}, \ldots, t_{i_n}}(B_1 \times \ldots \times B_n) = P_{t_1, \ldots, t_n}(B_{i_1} \times \ldots \times B_{i_n})$

2. **Consistency**: $P_{t_{1, \ldots, t_{n-1}}}(B_1 \times \ldots \times B_{n-1} \times \mathcal{S}_{t_n}) = P_{t_1, \ldots, t_{n-1}}(B_1 \times \ldots \times B_{n-1})$

The following theorem evidences that these properties are sufficient to prove the existence of a random function $X$ with given finite-dimensional distributions.

**Theorem 1.1.2 (Kolmogorov):**
Let $\{P_{t_1, \ldots, t_n}, n \in \mathbb{N}, \{t_1, \ldots, t_n\} \subset T\}$ be a family of probability measures on

\[ (\mathbb{R}^m \times \ldots \times \mathbb{R}^m, \mathcal{B}(\mathbb{R}^m) \otimes \ldots \otimes \mathcal{B}(\mathbb{R}^m)), \]

which fulfill conditions 1 and 2 of Exercise 1.1.4. Then there exists a random function $X = \{X(t), t \in T\}$ defined on a probability space $(\Omega, \mathcal{A}, P)$ with finite-dimensional distributions $P_{t_1, \ldots, t_n}$.

**Proof** See [13], Section II.9. \(\square\)

This theorem also holds on more general (however not arbitrary!) spaces than $\mathbb{R}^m$, on so-called Borel spaces, which are in a sense isomorphic to $([0, 1], \mathcal{B}[0, 1])$ or a subspace of that.

**Definition 1.1.7**
Let $X = \{X(t), t \in T\}$ be a random function with values in $(\mathcal{S}, \mathcal{B})$, i.e., $X(t) \in \mathcal{S}$ almost surely for arbitrary $t \in T$. Let $(T, C)$ be itself a measurable space. $X$ is called measurable if the mapping $X : (\omega, t) \mapsto X(\omega, t) \in \mathcal{S}$, $(\omega, t) \in \Omega \times T$, is $\mathcal{A} \otimes C | \mathcal{B}$-measurable.

Thus, Definition 1.1.7 not only provides the measurability of $X$ with respect to $\omega \in \Omega$: $X(\cdot, t) \in \mathcal{A} | \mathcal{B}$ for all $t \in T$, but $X(\cdot, \cdot) \in \mathcal{A} \otimes C | \mathcal{B}$ as a function of $(\omega, t)$. The measurability of $X$ is of significance if $X(\omega, t)$ is considered at random moments $\tau : \Omega \to T$, i.e., $X(\omega, \tau(\omega))$. This is in particular the case in the theory of martingales if $\tau$ is a so-called stopping time for $X$. The distribution of $X(\omega, \tau(\omega))$ might differ considerably from the distribution of $X(\omega, t)$, $t \in T$. 


1.2 Elementary examples

The theorem of Kolmogorov can be used directly for the explicit construction of random processes only in few cases, since for a lot of random functions their finite-dimensional distributions are not given explicitly. In these cases a new random function \( X = \{X(t), t \in T\} \) is built as \( X(t) = g(t, Y_1, Y_2, \ldots) \), \( t \in T \), where \( g \) is a measurable function and \( \{Y_n\} \) a sequence of random elements (also random functions), whose existence has already been ensured. For that we give several examples.

Let \( X = \{X(t), t \in T\} \) be a real-valued random function on a probability space \((\Omega, \mathcal{A}, P)\).

1. White noise:

**Definition 1.2.1**

The random function \( X = \{X(t), t \in T\} \) is called white noise, if all \( X(t), t \in T \), are independent and identically distributed (i.i.d.) random variables.

White noise exists according to the Theorem 1.1.1. It is used to model the noise in (electromagnetic or acoustical) signals. If \( X(t) \sim Ber(p), p \in (0, 1), t \in T \), one means Salt-and-pepper noise, the binary noise, which occurs at the transfer of binary data in computer-networks. If \( X(t) \sim N(0, \sigma^2), \sigma^2 > 0, t \in T \), then \( X \) is called Gaussian white noise. It occurs e.g. in acoustical signals.

2. Gaussian random function:

**Definition 1.2.2**

The random function \( X = \{X(t), t \in T\} \) is called Gaussian, if all of its finite-dimensional distributions are Gaussian, i.e. for all \( n \in \mathbb{N}, t_1, \ldots, t_n \subset T \) it holds

\[
X_{t_1, \ldots, t_n} = ((X(t_1), \ldots, X(t_n))^{\top} \sim N(\mu_{t_1, \ldots, t_n}, \sum_{t_1, \ldots, t_n}^n),
\]

where the mean is given by \( \mu_{t_1, \ldots, t_n} = (\mathbb{E}X(t_1), \ldots, \mathbb{E}X(t_n))^{\top} \) and the covariance matrix is given by \( \sum_{t_1, \ldots, t_n} = ((\text{cov}(X(t_i), X(t_j)))^{n}_{i,j=1} \).

**Exercise 1.2.1**

Proof that the distribution of an Gaussian random function \( X \) is uniquely determined by its mean value function \( \mu(t) = \mathbb{E}X(t), t \in T, \) and covariance function \( C(s, t) = \mathbb{E}[X(s)X(t)] \), \( s, t \in T \), respectively.

An example for a Gaussian process is the so-called Wiener process (or Brownian motion) \( X = \{X(t), t \geq 0\} \), which has the expected value zero \( (\mu(t) \equiv 0, t \geq 0) \) and the covariance function \( C(s, t) = \min \{s, t\}, s, t \geq 0 \). Usually it is additionally required that the paths of \( X \) are continuous functions.

We shall investigate the regularity properties of the paths of random functions in more detail in Section 1.3. Now we can say that such a process exists with probability one (with almost surely continuous trajectories).

**Exercise 1.2.2**

Prove that the Gaussian white noise is a Gaussian random function.

3. Lognormal- and \( \chi^2 \)-functions:

The random function \( X = \{X(t), t \in T\} \) is called lognormal, if \( X(t) = e^{Y(t)}, \) where \( Y = \)
\{ Y(t), t \in T \} \text{ is a Gaussian random function. } X \text{ is called } \chi^2\text{-function, if } X(t) = \| Y(t) \|^2,

where \( Y = \{ Y(t), t \in T \} \) is a Gaussian random function with values in \( \mathbb{R}^n \), for which \( Y(t) \sim \mathcal{N}(0, I) \), \( t \in T \); here \( I \) is the \((n \times n)\)-unit matrix. Then it holds that \( X(t) \sim \chi^2_n \), \( t \in T \).

4. Cosine wave:
\( X = \{ X(t), t \in \mathbb{R} \} \) is defined by \( X(t) = \sqrt{2} \cos(2\pi Y + tZ) \), where \( Y \sim \mathcal{U}([0, 1]) \) and \( Z \) is a random variable, which is independent of \( Y \).

Exercise 1.2.3
Let \( X_1, X_2, \ldots \) be i.i.d. cosine waves. Determine the weak limit of the finite-dimensional distributions of the random function \( \left\{ \frac{1}{\sqrt{n}} \sum_{k=1}^{n} X_k(t), \ t \in \mathbb{R} \right\} \) for \( n \to \infty \).

5. Poisson process:
Let \( \{ Y_n \}_{n \in \mathbb{N}} \) be a sequence of i.i.d. random variables \( Y_n \sim \text{Exp}(\lambda), \lambda > 0 \). The stochastic process \( X = \{ X(t), t \geq 0 \} \) defined as \( X(t) = \max \{ n \in \mathbb{N} : \sum_{k=1}^{n} Y_k \leq t \} \) is called Poisson process with intensity \( \lambda > 0 \). \( X(t) \) counts the number of certain events until the time \( t > 0 \), where the typical interval between two of these events is \( \text{Exp}(\lambda) \)-distributed. These events can be claim arrivals of an insurance portfolio the records of elementary particles in the Geiger counter, etc. Then \( X(t) \) represents the number of claims or particles within the time interval \([0, t]\).

1.3 Regularity properties of trajectories

The theorem of Kolmogorov provides the existence of the distribution of a random function with given finite-dimensional distributions. However, it does not provide a statement about the properties of the paths of \( X \). This is understandable since all random objects are defined in the almost surely sense (a.s.) in probability theory, with the exception of a set \( A \subset \Omega \) with \( P(A) = 0 \).

Example 1.3.1
Let \( \Omega, \mathcal{A}, P = ([0, 1], \mathcal{B}([0, 1]), \nu_1) \), where \( \nu_1 \) is the Lebesgue measure on \([0, 1]\). We define \( X = \{ X(t), \ t \in [0, 1] \} \) by \( X(t) \equiv 0, \ t \in [0, 1] \) and \( Y = \{ Y(t), \ t \in [0, 1] \} \) by

\[
Y(t) = \begin{cases} 1, & t = U, \\ 0, & \text{sonst}, \end{cases}
\]

where \( U(\omega) = \omega, \omega \in [0, 1], \) is a \( \mathcal{U}([0, 1]) \)-distributed random variable defined on \( (\Omega, \mathcal{A}, P) \). Since \( P(Y(t) = 0) = 1, \ t \in T, \) because of \( P(U = t) = 0, \ t \in T, \) it is clear that \( X \overset{d}{=} Y \). Nevertheless, \( X \) and \( Y \) have different path properties since \( X \) has continuous and \( Y \) has discontinuous trajectories, and \( P(X(t) = 0, \ \forall t \in T) = 1, \) where \( P(Y(t) = 0, \ \forall t \in T) = 0. \)

It may well be that the „set of exceptions“ \( A \) (see above) is very different for \( X(t) \) for every \( t \in T \). Therefore, we require that all \( X(t), t \in T, \) are defined simultaneously on a subset \( \Omega_0 \subset \Omega \) with \( P(\Omega_0) = 1 \). The so defined random function \( \tilde{X} : \Omega_0 \times T \to \mathbb{R} \) is called modification of \( X : \Omega \times T \to \mathbb{R} \). \( X \) and \( \tilde{X} \) differ on a set \( \Omega \setminus \Omega_0 \) with probability zero. Therefore we indicate later when stating that „random function \( X \) possesses a property \( C \)“ that it exists a modification of \( X \) with this property \( C \). Let us hold it in the following definition:
1 General theory of random functions

Definition 1.3.1
The random functions \( X = \{ X(t), \ t \in T \} \) and \( Y = \{ Y(t), \ t \in T \} \) defined on the same probability space \((\Omega, \mathcal{A}, \mathbb{P})\) associated with \((S_t, \mathcal{B}_t)_{t \in T}\) have equivalent trajectories (or are called stochastically indistinguishable) if

\[
A = \{ \omega \in \Omega : X(\omega, t) \neq Y(\omega, t) \text{ for a } t \in T \} \in \mathcal{A}
\]

and \( \mathbb{P}(A) = 0 \).

This term implies that \( X \) and \( Y \) have paths, which coincide with probability one.

Definition 1.3.2
The random functions \( X = \{ X(t), \ t \in T \} \) and \( Y = \{ Y(t), \ t \in T \} \) defined on the same probability space \((\Omega, \mathcal{A}, \mathbb{P})\) are called (stochastically) equivalent, if

\[
B_t = \{ \omega \in \Omega : X(\omega, t) \neq Y(\omega, t) \} \in \mathcal{A}, \ t \in T,
\]

and \( \mathbb{P}(B_t) = 0, \ t \in T \).

We also can say that \( X \) and \( Y \) are versions or modifications of one and the same random function. If the space \((\Omega, \mathcal{A}, \mathbb{P})\) is complete (i.e. the implication of \( A \in \mathcal{A} : \mathbb{P}(A) = 0 \) is for all \( B \subset A: B \in \mathcal{A} \) (and then \( \mathbb{P}(B) = 0 \)), then the indistinguishable processes are stochastically equivalent, but vice versa is not always true (it is true for so-called separable processes. This is the case if \( T \) is countable).

Exercise 1.3.1
Prove that the random functions \( X \) and \( Y \) in Example 1.3.1 are stochastically equivalent.

Definition 1.3.3
The random functions \( X = \{ X(t), \ t \in T \} \) and \( Y = \{ Y(t), \ t \in T \} \) (not necessarily defined on the same probability space) are called equivalent in distribution, if \( \mathbb{P}_X = \mathbb{P}_Y \) on \((S_t, \mathcal{B}_t)_{t \in T}\). Notation: \( X \overset{d}{=} Y \).

According to Theorem 1.1.2 it is sufficient for the equivalence in distribution of \( X \) and \( Y \) that they possess the same finite-dimensional distributions. It is clear that stochastic equivalence implies equivalence in distribution, but not the other way around.

Now, let \( T \) and \( S \) be Banach spaces with norms \( | \cdot |_T \) and \( | \cdot |_S \), respectively. The random function \( X = \{ X(t), \ t \in T \} \) is now defined on \((\Omega, \mathcal{A}, \mathbb{P})\) with values in \((S, \mathcal{B})\).

Definition 1.3.4
The random function \( X = \{ X(t), \ t \in T \} \) is called

a) stochastically continuous on \( T \), if \( X(s) \overset{\mathbb{P}}{\underset{s \rightarrow t}{\rightarrow}} X(t) \), for arbitrary \( t \in T \), i.e.

\[
\mathbb{P}(|X(s) - X(t)|_S > \varepsilon) \overset{s \rightarrow t}{\rightarrow} 0, \text{ for all } \varepsilon > 0.
\]

b) \( L^p \)-continuous on \( T \), \( p \geq 1 \), if \( X(s) \overset{L^p}{\underset{s \rightarrow t}{\rightarrow}} X(t) \), \( t \in T \), i.e. \( \mathbb{E}|X(s) - X(t)|^p \overset{s \rightarrow t}{\rightarrow} 0 \). For \( p = 2 \) the specific notation „continuity in the square mean “is used.

c) a.s. continuous on \( T \), if \( X(s) \overset{f.s.}{\underset{s \rightarrow t}{\rightarrow}} X(t) \), \( t \in T \), i.e., \( \mathbb{P}(X(s) \underset{s \rightarrow t}{\rightarrow} X(t)) = 1 \), \( t \in T \).

d) continuous, if all trajectories of \( X \) are continuous functions.
In applications one is interested in the cases c) and d), although the weakest form of continuity is the stochastic continuity.

\[ L^p\text{-continuity} \implies \text{stochastic continuity} \iff \text{a.s. continuity} \iff \text{continuity of all paths} \]

Why are cases c) and d) important? Let us consider an example.

**Example 1.3.2**
Let \( T = [0, 1] \) and \((\Omega, \mathcal{A}, \mathbb{P})\) be the canonical probability space with \( \Omega = \mathbb{R}^{[0, 1]} \), i.e. \( \Omega = \prod_{t \in [0, 1]} \mathbb{R} \). Let \( X = \{X(t), \ t \in [0, 1]\} \) be a stochastic process on \((\Omega, \mathcal{A}, \mathbb{P})\). Not all events are elements of \( \mathcal{A} \), like e.g. \( A = \{\omega \in \Omega : X(\omega, t) = 0 \text{ for all } t \in [0, 1]\} = \cap_{t \in [0, 1]} \{X(\omega, t) = 0\} \), since this is an intersection of measurable events from \( \mathcal{A} \) in an uncountable number. If however \( X \) is continuous, then all of its paths are continuous functions and one can write \( A = \cap_{t \in D} \{X(\omega, t) = 0\} \), where \( D \) is a dense countable subset of \([0, 1]\), e.g., \( D = \mathbb{Q} \cap [0, 1] \). Then it holds that \( A \in \mathcal{A} \).

However, in many applications (like e.g. in financial mathematics) it is not realistic to consider stochastic processes with continuous paths as models for real phenomena. Therefore, a bigger class of possible trajectories of \( X \) is allowed: the so-called \( \text{càdlàg-class} \) (càdlàg = continue à droite, limitée à gauche (fr.)).

**Definition 1.3.5**
A stochastic process \( X = \{X(t), \ t \in \mathbb{R}\} \) is called \( \text{càdlàg} \), if all of its trajectories are right-side continuous functions, which have left-side limits.

Now, we would like to consider the properties of the notion of continuity (introduced above) in more detail. One can note e.g. that the stochastic continuity is a property of the two-dimensional distribution \( \mathbb{P}_{s,t} \) of \( X \). This is shown by the following lemma.

**Lemma 1.3.1**
Let \( X = \{X(t), \ t \in T\} \) be a random function associated with \((\mathcal{S}, \mathcal{B})\), where \( \mathcal{S} \) and \( T \) are Banach spaces. The following statements are equivalent:

a) \( X(s) \xrightarrow{\mathbb{P}} Y \),

b) \( \mathbb{P}_{s,t} \xrightarrow{d} \mathbb{P}_{(Y,Y)} \),

where \( t_0 \in T \) and \( Y \) is a \( \mathcal{S} \)-valued \( \mathcal{B} \)-random element. For the stochastic continuity of \( X \), one should choose \( t_0 \in T \) arbitrarily and \( Y = X(t_0) \).

**Proof** a) \(\Rightarrow\) b) \( X(s) \xrightarrow{\mathbb{P}} Y \) means \( (X(s), X(t))^\top \xrightarrow{\mathbb{P}} (Y, Y)^\top \).

\[
\mathbb{P}(\{ |X(s) - X(t)|_2 > \varepsilon \}) \leq \mathbb{P}(\{ |X(s) - Y|_\mathcal{S} > \varepsilon/2 \}) + \mathbb{P}(\{ |X(t) - Y|_\mathcal{S} > \varepsilon/2 \}) \xrightarrow{s,t \rightarrow t_0} 0
\]

This results in \( \mathbb{P}_{s,t} \xrightarrow{d} \mathbb{P}_{(Y,Y)} \), since \( \mathbb{P} \)-convergence is stronger than \( d \)-convergence.

b) \(\Rightarrow\) a) For arbitrary \( \varepsilon > 0 \) we consider a continuous function \( g_\varepsilon : \mathbb{R} \rightarrow [0, 1] \) with \( g_\varepsilon(0) = 0 \), \( g_\varepsilon(x) = 1 \),
The assertion results from the Cauchy-Schwarz inequality:

\[ \text{Proof} \]

\[ \text{covariance function of } \]

\[ \begin{align*}
\text{hence } P(|X(s) - X(t)|_S > \varepsilon) & \leq \text{E}_c(|X(s) - X(t)|_S) = \int_S \int_S g_c(|x - y|_S) P_{s,t}(d(x, y)) \xrightarrow{s \to t_0} \\
\text{f}_g \text{P}_{s,t}(d(x, y)) = 0, \text{ since } P(Y, Y) \text{ is concentrated on } \{(x, y) \in S^2 : x = y\} \text{ and } g_c(0) = 0. \text{ Thus } \{X(s)\}_{s \to t_0} \text{ is a fundamental sequence (in probability), therefore } X(s) \xrightarrow{P \text{ s \to t_0}} Y. \]
\]

\[ \square \]

It may be that \( X \) is stochastically continuous, although all of the paths of \( X \) have jumps, i.e. \( X \) cannot possess any a.s. continuous modification. The descriptive explanation for that is that such \( X \) may have a jump at concrete \( t \in T \) with probability zero. Therefore jumps of the paths of \( X \) always occur at different locations.

**Exercise 1.3.2**

Prove that the Poisson process is stochastically continuous, although it does not possess any a.s. continuous modification.

**Exercise 1.3.3**

Let \( T \) be compact. Prove that if \( X \) is stochastically continuous on \( T \), then it also is uniformly stochastically continuous, i.e., for all \( \varepsilon, \eta > 0 \) \( \exists \delta > 0 \), such that for all \( s, t \in T \) with \( |s - t|_T < \delta \) it holds that \( P(|X(s) - X(t)|_S > \varepsilon) < \eta. \)

Now let \( S = \mathbb{R}, \text{EX}^2(t) < \infty, t \in T, \text{EX}(t) = 0, t \in T \). Let \( C(s, t) = \text{E}[X(s)X(t)] \) be the covariance function of \( X \).

**Lemma 1.3.2**

For all \( t_0 \in T \) and a random variable \( Y \) with \( \text{EY}^2 < \infty \) the following assertions are equivalent:

a) \( X(s) \xrightarrow{L^2_{s \to t_0}} Y \)

b) \( C(s, t) \xrightarrow{s,t \to t_0} \text{EY}^2 \)

**Proof**

a) \( \Rightarrow \) b)

The assertion results from the Cauchy-Schwarz inequality:

\[ |C(s, t) - \text{EY}^2| = |\text{E}(X(s)X(t)) - \text{EY}^2| = |\text{E}[(X(s) - Y)(X(t) - Y)]| - \text{EY}^2| \leq \text{E}|X(s) - Y| |X(t) - Y| + \text{E}|(X(s) - Y)Y| + \text{E}|(X(t) - Y)Y| \]

\[ \leq \text{E} \left( \frac{|X(s) - Y|^2 |X(t) - Y|^2}{||X(s) - Y||_2^2 ||X(t) - Y||_2^2} + \frac{\text{EY}^2 |X(s) - Y|^2}{||X(s) - Y||_2^2} + \frac{\text{EY}^2 |X(t) - Y|^2}{||X(t) - Y||_2^2} \right) \xrightarrow{s,t \to t_0} 0 \]

with assumption a).

b) \( \Rightarrow \) a)

\[ \text{E}(X(s) - X(t))^2 = \text{E}(X(s))^2 - 2\text{E}[X(s)X(t)] + \text{E}(X(t))^2 \]

\[ = C(s, s) + C(t, t) - 2C(s, t) \xrightarrow{s,t \to t_0} 2\text{EY}^2 - 2\text{EY}^2 = 0. \]
Thus, \( \{X(s), \ s \to t_0\} \) is a fundamental sequence in the \( L^2 \)-sense, and we get \( X(s) \xrightarrow{L^2} Y \).

A random function \( X \), which is continuous in the mean-square sense, may still have uncontinuous trajectories. In most of the cases which are practically relevant, \( X \) however has an a.s. continuous modification. Later on this will become more precise by stating a corresponding theorem.

**Corollary 1.3.1**
The random function \( X \), which satisfies the conditions of Lemma 1.3.2, is continuous on \( T \) in the mean-square sense if and only if its covariance function \( C : T^2 \to \mathbb{R} \) is continuous on the diagonal \( \text{diag} T^2 = \{(s, t) \in T^2 : s = t\} \), i.e., \( \lim_{s,t \to t_0} C(s, t) = C(t_0, t_0) = \text{Var} X(t_0) \) for all \( t_0 \in T \).

**Proof** Choose \( Y = X(t_0) \) in Lemma 1.3.2.

**Remark 1.3.1**
If \( X \) is not centered, then the continuity of \( \mu(\cdot) \) together with the continuity of \( C \) on \( \text{diag} T^2 \) is required to ensure the \( L^2 \)-continuity of \( X \) on \( T \).

**Exercise 1.3.4**
Give an example of a stochastic process with a.s. discontinuous trajectories, which is \( L^2 \)-continuous.

Now we consider the property of (a.s.) continuity in more detail. As mentioned before, we can merely talk about continuous modification or version of a process. The possibility to possess such a version also depends on the properties of the two-dimensional distributions of the process. This is proven by the following theorem (originally proven by A. Kolmogorov).

**Theorem 1.3.1**
Let \( X = \{X(t), \ t \in [a, b]\}, -\infty < a < b \leq +\infty \) be a real-valued stochastic process. \( X \) has a continuous version, if there exist constants \( \alpha, c, \delta > 0 \) such that

\[
\mathbb{E}|X(t + h) - X(t)|^\alpha < c|h|^{1+\delta}, \ t \in (a, b),
\]

for sufficiently small \( |h| \).

**Proof** See, e.g. [7], Theorem 2.23.

Now we turn to processes with càdlàg-trajectories. Let \((\Omega, \mathcal{A}, \mathbb{P})\) be a complete probability space.

**Theorem 1.3.2**
Let \( X = \{X(t), \ t \geq 0\} \) be a real-valued stochastic process and \( D \) a countable dense subset of \([0, \infty)\). If

a) \( X \) is stochastically right-handside continuous, i.e., \( X(t + h) \xrightarrow{P_{h \to +0}} X(t), \ t \in [0, +\infty) \),

b) the trajectories of \( X \) at every \( t \in D \) have finite right- and left-handside limits, i.e.,

\( \lim_{h \to \pm 0} X(t + h) < \infty, \ t \in D \) a.s.,

then \( X \) has a version with a.s. càdlàg-paths.

Without proof.
Lemma 1.3.3
Let \( X = \{X(t), \ t \geq 0\} \) and \( Y = \{Y(t), \ t \geq 0\} \) be two versions of a random function, both defined on the probability space \((\Omega, \mathcal{A}, \mathbb{P})\), with property that \( X \) and \( Y \) have a.s. right-hand-side continuous trajectories. Then \( X \) and \( Y \) are indistinguishable.

Proof Let \( \Omega_X, \Omega_Y \) be „sets of exception“ for which the trajectories of \( X \) and \( Y \), respectively are not right-sided continuous. It holds that \( \mathbb{P}(\Omega_X) = \mathbb{P}(\Omega_Y) = 0 \). Consider \( A_t = \{\omega \in \Omega : X_0(\omega, t) \neq Y_0(\omega, t)\}, \ t \in [0, +\infty) \) and \( A = \cup_t A_t \), where \( Q_+ = Q \cap [0, +\infty) \). Since \( X \) and \( Y \) are stochastically equivalent, it holds that \( \mathbb{P}(A) = 0 \) and therefore

\[
\mathbb{P}(\tilde{A}) \leq \mathbb{P}(A) + \mathbb{P}(\Omega_X) + \mathbb{P}(\Omega_Y) = 0,
\]

where \( \tilde{A} = A \cup \Omega_X \cup \Omega_Y \). Therefore \( X(\omega, t) = Y(\omega, t) \) holds for \( t \in Q_+ \) and \( \omega \in \Omega \setminus \tilde{A} \). Now, we prove this for all \( t \geq 0 \). For arbitrary \( t \geq 0 \) a sequence \( \{t_n\} \subset Q_+ \) exists, such that \( t_n \downarrow t \). Since \( X_0(\omega, t_n) = Y_0(\omega, t_n) \) for all \( n \in \mathbb{N} \) and \( \omega \in \Omega \setminus \tilde{A} \), it holds that \( X(\omega, t) = \lim_{n \to \infty} X_0(\omega, t_n) = \lim_{n \to \infty} Y_0(\omega, t_n) = Y(\omega, t) \) for \( t \geq 0 \) and \( \omega \in \Omega \setminus \tilde{A} \). Therefore \( X \) and \( Y \) are indistinguishable.

Corollary 1.3.2
If càdlàg-processes \( X = \{X(t), \ t \geq 0\} \) and \( Y = \{Y(t), \ t \geq 0\} \) are versions of the same random function then they are indistinguishable.

1.4 Differentiability of trajectories

Let \( T \) be a linear normed space.

Definition 1.4.1
A real-valued random function \( X = \{X(t), \ t \in T\} \) is differentiable on \( T \) in direction \( h \in T \) stochastically, in the \( L^p \)-sense, \( p \geq 1 \), or a.s., if

\[
\lim_{l \to 0} \frac{X(t + hl) - X(t)}{l} = X'_h(t), \ t \in T
\]

exists in the corresponding sense, namely stochastically, in the \( L^p \)-space or a.s..

The Lemmas 1.3.1 - 1.3.2 show that the stochastic differentiability is a property that is determined by three-dimensional distributions of \( X \) (because the common distribution of \( \frac{X(t+hl)-X(t)}{l} \) and \( \frac{X(t+hl')-X(t)}{l'} \) should converge weakly), whereas the differentiability in the mean-square sense is determined by the smoothness of the covariance function \( C(s, t) \).

Exercise 1.4.1
Show that

1. the Wiener process is not stochastically differentiable on \( [0, \infty) \).
2. the Poisson process is stochastically differentiable on \( [0, \infty) \), however not in the \( L^p \)-mean, \( p > 1 \).

Lemma 1.4.1
A centered random function \( X = \{X(t), \ t \in T\} \) (i.e., \( \mathbb{E}X(t) = 0, \ t \in T \)) with \( \mathbb{E}[X^2(t)] < \infty, t \in T \) is \( L^2 \)-differentiable in \( t \in T \) in direction \( h \in T \) if its covariance function \( C \) is differentiable.
twice in \((t, t)\) in direction \(h\), i.e., if \(\exists C''_{hh}(t, t) = \frac{\partial^2 C(s, t)}{\partial s \partial h_h} \bigg|_{s=t}\). \(X'_h(t)\) is \(L^2\)-continuous in \(t \in T\) if \(C''_{hh}(s, t) = \frac{\partial^2 C(s, t)}{\partial s \partial h_h}\) is continuous in \(s = t\). Moreover, \(C''_{hh}(s, t)\) is the covariance function of \(X'_h = \{X'_h(t), t \in T\}\).

**Proof** According to Lemma 1.3.2 it is enough to show that

\[
I = \lim_{l, l' \to 0} E \left( \frac{X(t + lh) - X(t)}{l} \cdot \frac{X(s + l'h) - X(s)}{l'} \right)
\]

exists for \(s = t\). Indeed we get

\[
I = \left. \frac{1}{ll'} \left( C(t + lh, s + l'h) - C(t + lh, s) - C(t, s + l'h) + C(t, s) \right) \right|_{l, l' \to 0} = \frac{C''_{hh}(s, t)}{l l'} \to C''_{hh}(s, t).
\]

All other statements of the lemma result from this relation.

**Remark 1.4.1**

The properties of the \(L^2\)-differentiability and a.s. differentiability of random functions are disjoint in the following sense: there are stochastic processes that have \(L^2\)-differentiable paths, although they are a.s. discontinuous, and vice versa, processes with a.s. differentiable paths are not always \(L^2\)-differentiable, since e.g. the first derivative of their covariance function is not continuous.

**Exercise 1.4.2**

Give appropriate examples!

### 1.5 Moments und covariance

Let \(X = \{X(t), t \in T\}\) be a random function that is real-valued, and let \(T\) be an arbitrary index space.

**Definition 1.5.1**

The mixed moment \(\mu^{(j_1, \ldots, j_n)}(t_1, \ldots, t_n)\) of \(X\) of order \((j_1, \ldots, j_n) \in \mathbb{N}^n\), \(t_1, \ldots, t_n \in T\) is given by \(\mu^{(j_1, \ldots, j_n)}(t_1, \ldots, t_n) = E[X^{j_1}(t_1) \cdot \ldots \cdot X^{j_n}(t_n)]\), where it is required that the expected value exists and is finite. Then it is sufficient to assume that \(E|X(t)|^j < \infty\) for all \(t \in T\) and \(j = j_1 + \ldots + j_n\).

Important special cases:

1. \(\mu(t) = \mu^{(1)}(t) = E X(t), t \in T\) – mean value function of \(X\).
2. \(\mu^{(1,1)}(s, t) = E [X(s) X(t)] = C(s, t) - (\text{non-centered})\ covariance\ function\ of\ \(X\).\) Whereas the centered covariance function is: \(K(s, t) = \text{cov}(X(s), X(t)) = \mu^{(1,1)}(s, t) - \mu(s) \mu(t), s, t \in T\).

**Exercise 1.5.1**

Show that the centered covariance function of a real-valued random function \(X\)

1. is symmetric, i.e., \(K(s, t) = K(t, s), s, t \in T\).
2. is positive semidefinite, i.e., for \( n \in \mathbb{N}, t_1, \ldots, t_n \in T, z_1, \ldots, z_n \in \mathbb{R} \) it holds that
\[
\sum_{i,j=1}^{n} K(t_i, t_j)z_iz_j \geq 0.
\]

3. satisfies \( K(t, t) = \text{Var} X(t), t \in T \).

Property 2 also holds for the non-centered covariance function \( C(s, t) \).

The mean value function \( \mu(t) \) shows a (non random) trend. If \( \mu(t) \) is known, the random function \( X \) can be centered by considering a random function \( Y = \{Y(t), t \in T\} \) with \( Y(t) = X(t) - \mu(t), t \in T \).

The covariance function \( K(s, t) \) (\( C(s, t) \), respectively) contains information about the dependence structure of \( X \). Sometimes the correlation function \( R(s, t) = \frac{K(s, t)}{\sqrt{K(s,s)K(t,t)}} \) for all \( s, t \in T \):

\[
K(s, s) = \text{Var} X(s) > 0, K(t, t) = \text{Var} X(t) > 0 \text{ is used instead of } K \text{ and } C, \text{ respectively. Because of the Cauchy-Schwarz inequality it holds that } |R(s, t)| \leq 1, s, t \in T. \text{ The set of all mixed moments in general does not (uniquely) determine the distribution of a random function.}

**Exercise 1.5.2**

Give an example of different random functions \( X = \{X(t), t \in T\} \) und \( Y = \{Y(t), t \in T\} \), for which it holds that \( \mathbb{E}X(t) = \mathbb{E}Y(t), t \in T \) and \( \mathbb{E}(X(s)X(t)) = \mathbb{E}(Y(s)Y(t)), s, t \in T \).

**Exercise 1.5.3**

Let \( \mu : T \rightarrow \mathbb{R} \) be a measurable function and \( K : T \times T \rightarrow \mathbb{R} \) be a positive semidefinite symmetric function. Prove that a random function \( X = \{X(t), t \in T\} \) exists with \( \mathbb{E}X(t) = \mu(t), \text{cov}(X(s), X(t)) = K(s, t), s, t \in T \).

Let now \( X = \{X(t), t \in T\} \) be a real-valued random function with \( \mathbb{E}|X(t)|^k < \infty, t \in T, \) for a \( k \in \mathbb{N} \).

**Definition 1.5.2**

The mean increment of order \( k \) of \( X \) is given by \( \gamma_k(s, t) = \mathbb{E}(X(s) - X(t))^k, s, t \in T \).

Special attention is paid to the function \( \gamma(s, t) = \frac{1}{2}\gamma_2(s, t) = \frac{1}{2}\mathbb{E}(X(s) - X(t))^2, s, t \in T \), which is called variogram of \( X \). In geostatistics the variogram is often used instead of the covariance function. A lot of times we discard the condition \( \mathbb{E}X^2(t) < \infty, t \in T, \) instead we assume that \( \gamma(s, t) < \infty \) for all \( s, t \in T \).

**Exercise 1.5.4**

Prove that there exist random functions without finite second moments with \( \gamma(s, t) < \infty, s, t \in T \).

**Exercise 1.5.5**

Show that for a random function \( X = \{X(t), t \in T\} \) with mean value function \( \mu \) and covariance function \( K \) it holds that:

\[
\gamma(s, t) = \frac{K(s, s) + K(t, t)}{2} - K(s, t) + \frac{1}{2}(\mu(s) - \mu(t))^2, s, t \in T.
\]

If the random function \( X \) is complex-valued, i.e., \( X : \Omega \times T \rightarrow \mathbb{C} \), with \( \mathbb{E}|X(t)|^2 < \infty, t \in T \), then the covariance function of \( X \) is introduced as \( K(s, t) = \mathbb{E}(X(s) - \mathbb{E}(X))(X(t) - \mathbb{E}(X)), s, t \in T \), where \( \overline{z} \) is the complex conjugate of \( z \in \mathbb{C} \). Then it holds that \( K(s, t) = \overline{K(t, s)}, s, t \in T, \) and \( K \) is positive semidefinite, i.e, for all \( n \in \mathbb{N}, t_1, \ldots, t_n \in T, z_1, \ldots, z_n \in \mathbb{C} \) it holds that \( \sum_{i,j=1}^{n} K(t_i, t_j)z_i\overline{z_j} \geq 0. \)


1.6 Stationarity and Independence

Let $T$ be a subset of the linear vector space with operations $+$, $-$ over space $\mathbb{R}$.

**Definition 1.6.1**

The random function $X = \{X(t), \ t \in T\}$ is called *stationary (strict sense stationary)* if for all $n \in \mathbb{N}, h, t_1, \ldots, t_n \in T$ with $t_1 + h, \ldots, t_n + h \in T$ it holds that:

$$P(X(t_1), \ldots, X(t_n)) = P(X(t_1 + h), \ldots, X(t_n + h)),$$

i.e., all finite-dimensional distributions of $X$ are invariant with respect to translations in $T$.

**Definition 1.6.2**

A (complex-valued) random function $X = \{X(t), \ t \in T\}$ is called *second-order stationary* (or *wide sense stationary*) if $E|X(t)|^2 < \infty$, $t \in T$, and $\mu(t) \equiv EX(t) \equiv \mu$, $t \in T$, $K(s, t) = \text{cov}(X(s), X(t)) = K(s + h, t + h)$ for all $h, s, t \in T : s + h, t + h \in T$.

If $X$ is second-order stationary, it is convenient to introduce a function $K(t) := K(0, t), \ t \in T$ whereby $0 \in T$.

Strict sense stationarity and wide sense stationarity do not result from each other. However it is clear that if a complex-valued random function is strict sense stationary and possesses finite second-order moments, then the function is also second-order stationary.

**Definition 1.6.3**

A real-valued random function $X = \{X(t), \ t \in T\}$ is *intrinsic second-order stationary* if $\gamma_k(s, t), s, t \in T$ exist for $k \leq 2$, and for all $s, t, h \in T$, $s + h, t + h \in T$ it holds that $\gamma_1(s, t) = 0$, $\gamma_2(s, t) = \gamma_2(s + h, t + h)$.

For real-valued random functions, intrinsic second-order stationarity is more general than second-order stationarity since the existence of $E|X(t)|^2$, $t \in T$ is not required.

The analogue of the stationarity of increments of $X$ also exists in strict sense.

**Definition 1.6.4**

Let $X = \{X(t), \ t \in T\}$ be a real-valued stochastic process, $T \subset \mathbb{R}$. It is said that $X$

1. possesses *stationary increments* if for all $n \in \mathbb{N}, h, t_0, t_1, t_2, \ldots, t_n \in T$, with $t_0 < t_1 < t_2 < \ldots < t_n, t_i + h \in T, i = 0, \ldots, n$ the distribution of

$$(X(t_1 + h) - X(t_0 + h), \ldots, X(t_n + h) - X(t_{n-1} + h))^\top$$

does not depend on $h$.

2. possesses *independent increments* if for all $n \in \mathbb{N}, t_0, t_1, \ldots, t_n \in T$ with $t_0 < t_1 < \ldots < t_n$ the random variables $X(t_0), X(t_1) - X(t_0), \ldots, X(t_n) - X(t_{n-1})$ are pairwise independent.

Let $(S_1, B_1)$ and $(S_2, B_2)$ be measurable spaces. In general it is said that two random elements $X : \Omega \to S_1$ and $X : \Omega \to S_2$ are *independent* on the same probability space $(\Omega, \mathcal{A}, P)$ if $P(X \in A_1, Y \in A_2) = P(X \in A_1)P(Y \in A_2)$ for all $A_1 \in B_1, A_2 \in B_2$.

This definition can be applied to the independence of random functions $X$ and $Y$ with phase space $(S_T, B_T)$, since they can be considered as random elements with $S_1 = S_2 = S_T, B_1 = B_2 = B_T$ (cf. Lemma 1.1.1). The same holds for the independence of a random element (or a random function) $X$ and of a sub-$\sigma$-algebra $\mathcal{G} \subset \mathcal{A}$: this is the case if $P(\{X \in A\} \cap G) = P(X \in A)P(G)$, for all $A \in B_1, G \in \mathcal{G}$ (or $A \in B_T, G \in \mathcal{G}$).
1 General theory of random functions

1.7 Processes with independent increments

In this section we concentrate on the properties and existence of processes with independent increments.

Let $\{\varphi_{s,t}, \ s, t \geq 0\}$ be a family of characteristic functions of probability measures $Q_{s,t}$, $s, t \geq 0$ on $\mathcal{B}(\mathbb{R})$, i.e., for $z \in \mathbb{R}$, $s, t \geq 0$ it holds that $\varphi_{s,t}(z) = \int_{\mathbb{R}} e^{izx} Q_{s,t}(dx)$.

**Theorem 1.7.1**

There exists a stochastic process $X = \{X(t), \ t \geq 0\}$ with independent increments with the property that for all $s, t \geq 0$ the characteristic function of $X(t) - X(s)$ is equal to $\varphi_{s,t}$ if and only if

$$\varphi_{s,t} = \varphi_{s,u} \varphi_{u,t}$$

(1.7.1)

for all $0 \leq s < u < t < \infty$. Thereby the distribution of $X(0)$ can be chosen arbitrarily.

**Proof**

The necessity of the condition (1.7.1) is clear since for all $s \in (0, \infty): s < u < t$ it holds

$$X(t) - X(s) = X(t) - X(u) + X(u) - X(s),$$

and $X(t) - X(u)$ and $X(u) - X(s)$ are independent. Then it holds

$$\varphi_{s,t} = \varphi_{Y_1} \varphi_{Y_2} = \varphi_{s,u} \varphi_{u,t}.$$

Now we prove the sufficiency.

If the existence of a process $X$ with independent increments and property $\varphi_{X(t) - X(s)} = \varphi_{s,t}$ on a probability space $(\Omega, \mathcal{A}, P)$ had already been proven, one could define the characteristic functions of all of its finite-dimensional distributions with the help of $\{\varphi_{s,t}\}$ as follows.

Let $n \in \mathbb{N}, 0 = t_0 < t_1 < \ldots < t_n < \infty$ and $Y = (X(t_0), X(t_1) - X(t_0), \ldots, X(t_n) - X(t_{n-1}))^\top$. The independence of increments results in

$$\varphi_Y(z_0, z_1, \ldots, z_n) = \mathbb{E} e^{iz^\top Y} = \varphi_{X(t_0)}(z_0) \varphi_{t_0,t_1}(z_1) \ldots \varphi_{t_{n-1},t_n}(z_n), \ z \in \mathbb{R}^{n+1},$$

where the distribution of $X(t_0)$ is an arbitrary probability measure $Q_0$ on $\mathcal{B}(\mathbb{R})$. For $X_{t_0,\ldots,t_n} = (X(t_0), X(t_1), \ldots, X(t_n))^\top$ however it holds that $X_{t_0,\ldots,t_n} = AY$, where

$$A = \begin{pmatrix} 1 & 0 & 0 & \ldots & 0 \\ 1 & 1 & 0 & \ldots & 0 \\ 1 & 1 & 1 & \ldots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 1 & 1 & 1 & \ldots & 1 \end{pmatrix}.$$ 

Then $\varphi_{X_{t_0,\ldots,t_n}}(z) = \varphi_{AY}(z) = \mathbb{E} e^{iz^\top AY} = \mathbb{E} e^{iz^\top (A^\top z, Y)} = \varphi_Y(A^\top z)$ holds. Therefore the distribution of $X_{t_0,\ldots,t_n}$ possesses the characteristic function $\varphi_{X_{t_0,\ldots,t_n}}(z) = \varphi_{Q_0(l_0)} \varphi_{t_0,t_1}(l_1) \ldots \varphi_{t_{n-1},t_n}(l_n)$, where $l = (l_1, l_1, \ldots, l_n)^\top = A^\top z$, thus

$$\begin{cases} l_0 = z_0 + \ldots + z_n \\ l_1 = z_1 + \ldots + z_n \\ \vdots \\ l_n = z_n \end{cases}$$
Thereby $\varphi_{X(t_0)} = \varphi_{Q_0}$ and $\varphi_{X_{t_1, \ldots, t_n}}(z_1, \ldots, z_n) = \varphi_{X_{t_0, \ldots, t_n}}(0, z_1, \ldots, z_n)$ holds for all $z_i \in \mathbb{R}$.

Now we prove the existence of such a process $X$.

For that we construct the family of characteristic functions

$$\{\varphi_{t_0}, \varphi_{t_0, t_1, \ldots, t_n}, \varphi_{t_1, \ldots, t_n}, 0 = t_0 < t_1 < \ldots < t_n < \infty, \ n \in \mathbb{N}\}$$

from $\varphi_{Q_0}$ and $\{\varphi_{s,t}, 0 \leq s < t\}$ as above, thus

$$\varphi_{t_0} = \varphi_{Q_0}, \ \varphi_{t_1, \ldots, t_n}(0, z_1, \ldots, z_n) = \varphi_{t_0, t_1, \ldots, t_n}(0, z_1, \ldots, z_n), \ z_i \in \mathbb{R},$$

$$\varphi_{t_0, \ldots, t_n}(z) = \varphi_{t_0}(z_1 + \ldots + z_n)\varphi_{t_0, t_1}(z_1 + \ldots + z_n)\ldots\varphi_{t_{n-1}, t_n}(z_n).$$

Now we have to check whether the corresponding probability measures of these characteristic
functions fulfill the conditions of Theorem 1.1.2. We will do that in equivalent form since by Exercise 1.8.1 the conditions of symmetry and consistency in Theorem 1.1.2 are equivalent to:

a) $\varphi_{t_{m-1}, t_m}(z_0, \ldots, z_n, \ldots, z_{m-1}) = \varphi_{t_0, \ldots, t_n}(z_0, \ldots, z_{m-1}, \ldots, z_n)$ for an arbitrary permutation $(0, 1, \ldots, n) \mapsto \{0, 1, \ldots, t_n\}$,

b) $\varphi_{t_0, \ldots, t_{m-1}, t_m}(z_0, \ldots, z_{m-1}, z_{m+1}, \ldots, z_n) = \varphi_{t_0, \ldots, t_n}(z_0, \ldots, z_{m-1}, \ldots, z_n)$, for all $z_0, \ldots, z_n \in \mathbb{R}$, $m \in \{1, \ldots, n\}$.

Condition a) is obvious. Condition b) holds since

$$\varphi_{t_{m-1}, t_m}(0 + z_{m+1} + \ldots + z_n)\varphi_{t_m, t_{m+1}}(z_{m+1} + \ldots + z_n) = \varphi_{t_{m-1}, t_{m+1}}(z_{m+1}, \ldots, z_n)$$

for all $m \in \{1, \ldots, n\}$. Thus, the existence of $X$ is proven.

**Example 1.7.1**  
1. If $T = \mathbb{N}_0 = \mathbb{N} \cup \{0\}$, then $X = \{X(t), \ t \in \mathbb{N}_0\}$ has independent increments if and only if $X(n) = \sum_{i=0}^{n} Y_i$, where $\{Y_i\}$ are independent random variables and $Y_n \overset{d}{=} X(n) - X(n-1), \ n \in \mathbb{N}$. Such a process $X$ is called random walk. It also may be defined for $Y_i$ with values in $\mathbb{R}_m$.

2. The Poisson process with intensity $\lambda$ has independent increments (we will show that later).

3. The Wiener process possesses independent increments.

**Exercise 1.7.1**  
Prove that!

**Exercise 1.7.2**  
Let $X = \{X(t), \ t \geq 0\}$ be a process with independent increments and $g : [0, \infty) \to \mathbb{R}$ an arbitrary (deterministic) function. Show that the process $Y = \{Y(t), \ t \geq 0\}$ with $Y(t) = X(t) + g(t), \ t \geq 0$, also possesses independent increments.

### 1.8 Additional exercises

**Exercise 1.8.1**  
Prove the following assertion: The family of probability measures $P_{t_1, \ldots, t_n}$ on $(\mathbb{R}^n, \mathcal{B}(\mathbb{R}^n)), \ n \geq 1, \ t = (t_1, \ldots, t_n) \in T^n$ fulfills the conditions of the theorem of Kolmogorov if and only if for $n \geq 2$ and for all $s = (s_1, \ldots, s_n) \in \mathbb{R}^n$ the following conditions are fulfilled:
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a) \( \varphi_{P_{t_1,\ldots,t_n}}((s_1,\ldots,s_n)^\top) = \varphi_{P_{t_{\pi(1)},\ldots,t_{\pi(n)}}}((s_{\pi(1)},\ldots,s_{\pi(n)})^\top) \) for all \( \pi \in S_n \).

b) \( \varphi_{P_{t_1,\ldots,t_{n-1}}}((s_1,\ldots,s_{n-1})^\top) = \varphi_{P_{t_1,\ldots,t_n}}((s_1,\ldots,s_{n-1},0)^\top) \).

Remark: \( \varphi(\cdot) \) denotes the characteristic function of the corresponding measure. \( S_n \) denotes the group of all permutations \( \pi : \{1,\ldots,n\} \rightarrow \{1,\ldots,n\} \).

Exercise 1.8.2
Show the existence of a random function whose finite-dimensional distributions are multivariate-normally distributed and explicitly give the measurable spaces \( (E_{t_1,\ldots,t_n},E_{t_1,\ldots,t_n}) \).

Exercise 1.8.3
Give an example of a family of probability measures \( P_{t_1,\ldots,t_n} \), which do not fulfill the conditions of the theorem of Kolmogorov.

Exercise 1.8.4
Let \( X = \{X(t), t \in T\} \) and \( Y = \{Y(t), t \in T\} \) be two stochastic processes which are defined on the same complete probability space \( (\Omega,F,P) \) and which take values in the measurable space \( (S,B) \).

a) Proof that: \( X \) and \( Y \) are stochastically equivalent \( \Rightarrow P_X = P_Y \).

b) Give an example of two processes \( X \) and \( Y \) for which holds: \( P_X = P_Y \), but \( X \) and \( Y \) are not stochastically equivalent.

c) Proof that: \( X \) and \( Y \) are stochastically indistinguishable \( \Rightarrow X \) and \( Y \) are stochastically equivalent.

d) Proof in the case of countability of \( T \): \( X \) and \( Y \) are stochastically equivalent \( \Rightarrow \) \( X \) and \( Y \) are stochastically indistinguishable.

e) Give in the case of uncountability of \( T \) an example of two processes \( X \) and \( Y \) for which holds: \( X \) and \( Y \) are stochastically equivalent but not stochastically indistinguishable.

Exercise 1.8.5
Let \( W = \{W(t), t \in \mathbb{R}\} \) be a Wiener Process. Which of the following processes are Wiener processes as well?

a) \( W_1 = \{W_1(t) := -W(t), t \in \mathbb{R}\} \),

b) \( W_2 = \{W_2(t) := \sqrt{t}W(1), t \in \mathbb{R}\} \),

c) \( W_3 = \{W_3(t) := W(2t) - W(t), t \in \mathbb{R}\} \).

Exercise 1.8.6
Given a stochastic process \( X = \{X(t), t \in [0,1]\} \) which consists of independent and identically distributed random variables with density \( f(x), x \in \mathbb{R} \). Show that such a process cannot be continuous in \( t \in [0,1] \).

Exercise 1.8.7
Give an example of a stochastic process \( X = \{X(t), t \in T\} \) which is stochastically continuous on \( T \), and prove why this is the case.
Exercise 1.8.8
In connection with the continuity of stochastic processes the so-called criterion of Kolmogorov plays a central role. (see also theorem 1.3.1 in the lecture notes): Let \( X = \{ X(t), t \in [a, b] \} \) be a real-valued stochastic process. If constants \( \alpha, \varepsilon > 0 \) and \( C := C(\alpha, \varepsilon) > 0 \) exist such that
\[
E|X(t + h) - X(t)|^\alpha \leq C|h|^{1+\varepsilon}
\] (1.8.1)
for sufficient small \( h \), then the process \( X \) possesses a continuous modification. Show that:

a) If you fix the variable \( \varepsilon = 0 \) in condition (1.8.1), then in general the condition is not sufficient for the existence of a continuous modification. *Hint: Consider the Poisson process.*

b) The Wiener process \( W = \{ W(t), t \in [0, \infty) \} \) possesses a continuous modification. *Hint: Consider the case \( \alpha = 4 \).*

Exercise 1.8.9
Show that the Wiener process \( W \) is not stochastically differentiable at any point \( t \in [0, \infty) \).

Exercise 1.8.10
Show that the covariance function \( C(s, t) \) of a complex-valued stochastic process \( X = \{ X(t), t \in T \} \)

a) is symmetric, i.e. \( C(s, t) = \overline{C(t, s)} \), \( s, t \in T \),

b) fulfills the identity \( C(t, t) = \text{Var} X(t) \), \( t \in T \),

c) is positive semidefinite, i.e. for all \( n \in \mathbb{N}, t_1, \ldots, t_n \in T, z_1, \ldots, z_n \in \mathbb{C} \) it holds that:
\[
\sum_{i=1}^{n} \sum_{j=1}^{n} C(t_i, t_j) z_i \overline{z_j} \geq 0.
\]

Exercise 1.8.11
Show that it exists a random function \( X = \{ X(t), t \in T \} \) which simultaneously fulfills the conditions:

- The second moment \( E X^2 \) does not exist.
- The variogram \( \gamma(s, t) \) is finite for all \( s, t \in T \).

Exercise 1.8.12
Give an example of a stochastic process \( X = \{ X(t), t \in T \} \) whose paths are simultaneously \( L^2 \)-differentiable but not almost surely differentiable, and prove why this is the case.

Exercise 1.8.13
Give an example of a stochastic process \( X = \{ X(t), t \in T \} \) whose paths are simultaneously almost surely differentiable but not \( L^1 \)-differentiable, and prove why this is the case.

Exercise 1.8.14
Proof that the Wiener process possesses independent increments.

Exercise 1.8.15
Proof: A (real-valued) stochastic process \( X = \{ X(t), t \in [0, \infty) \} \) with independent increments already has stationary increments if the distribution of the random variable \( X(t + h) - X(h) \) is independent of \( h \).
2 Counting processes

In this chapter we consider several examples of stochastic processes which model the counting of events and thus possess piecewise constant paths. Let $(\Omega, \mathcal{A}, P)$ be a probability space and $\{S_n\}_{n \in \mathbb{N}}$ a non-decreasing sequence of a.s. non-negative random variables, i.e. $0 \leq S_1 \leq S_2 \leq \ldots \leq S_n \leq \ldots$.

**Definition 2.0.1**
The stochastic process $N = \{N(t), \ t \geq 0\}$ is called a counting process if

$$N(t) = \sum_{n=1}^{\infty} 1(S_n \leq t),$$

where $1(A)$ is the indicator function of the event $A \in \mathcal{A}$.

$N(t)$ counts the events which occur at $S_n$ until time $t$. $S_n$ e.g. may be the time of occurrence of

1. the n-th elementary particle in the Geiger counter, or
2. a damage in the insurance of material damage, or
3. a data packet at a server within a computer network, etc.

A special case of the counting processes are the so-called renewal processes.

2.1 Renewal processes

**Definition 2.1.1**
Let $\{T_n\}_{n \in \mathbb{N}}$ be a sequence of i.i.d. non-negative random variables with $P(T_1 > 0) > 0$. A counting process $N = \{N(t), \ t \geq 0\}$ with $N(0) = 0$ a.s., $S_n = \sum_{k=1}^{n} T_k$, $n \in \mathbb{N}$, is called a renewal process. Thereby $S_n$ is called the time of the n-th renewal, $n \in \mathbb{N}$.

The name „renewal process“ is given by the following interpretation. The „interarrival times“ $T_n$ are interpreted as the lifetime of a technical spare part or mechanism within a system, thus $S_n$ is the time of the n-th break down of the system. The defective part is immediately replaced by a new part (comparable with the exchange of a lightbulb). Thus, $N(t)$ is the number of repairs (the so-called „renewals“) of the system until time $t$.

**Remark 2.1.1**
1. It is $N(t) = \infty$ if $S_n \leq t$ for all $n \in \mathbb{N}$.

2. Often it is assumed that only $T_2, T_3, \ldots$ are identically distributed with $\mathbb{E}T_n < \infty$. The distribution of $T_1$ is freely selectable. Such a process $N = \{N(t), \ t \geq 0\}$ is called delayed renewal process (with delay $T_1$).

3. Sometimes the requirement $T_n \geq 0$ is omitted.
4. It is clear that \( \{S_n\}_{n \in \mathbb{N}_0} \) with \( S_0 = 0 \) a.s., \( S_n = \sum_{k=1}^{n} T_k, n \in \mathbb{N} \) is a random walk.

5. If one requires that the \( n \)-th exchange of a defective part in the system takes a time \( T_n' \), then by \( \tilde{T}_n = T_n + T_n', n \in \mathbb{N} \) a different renewal process is given. Its stochastic property does not differ from the process which is given in definition 2.1.1.

In the following we assume that \( \mu = E T_n \in (0, \infty), n \in \mathbb{N} \).

**Theorem 2.1.1 (Individual ergodic theorem):**

Let \( N = \{N(t), t \geq 0\} \) be a renewal process. Then it holds that:

\[
\lim_{t \to \infty} \frac{N(t)}{t} = \frac{1}{\mu} \quad \text{a.s.}
\]

**Proof** For all \( t \geq 0 \) and \( n \in \mathbb{N} \) it holds that \( \{N(t) = n\} = \{S_n \leq t < S_{n+1}\} \), therefore \( S_{N(t)} \leq t < S_{N(t)+1} \) and

\[
\frac{S_{N(t)}}{N(t)} \leq \frac{t}{N(t)} \leq \frac{S_{N(t)+1}}{N(t)+1} = \frac{N(t)+1}{N(t)}.
\]

If we can show that \( \frac{S_{N(t)}}{N(t)} \xrightarrow{a.s.} \mu \) and \( N(t) \xrightarrow{a.s.} \infty \), then \( t \xrightarrow{\frac{N(t)}{N(t)}} \mu \) holds and therefore the assertion of the theorem.

According to the strong law of large numbers of Kolmogorov (cf. lecture notes „Wahrscheinlichkeitsrechnung“ (WR), theorem 7.4) it holds that \( \frac{S_n}{n} \xrightarrow{n \to \infty} \mu \), thus \( S_n \xrightarrow{n \to \infty} \infty \) and therefore

\[
P(N(t) < \infty) = 1 \quad \text{since} \quad P(N(t) = \infty) = P( S_n \leq t, \forall n) = 1 - P(\exists n : \forall m \in \mathbb{N}_0 \, S_{n+m} > t) = 1 - 1 = 0.
\]

Then \( N(t), t \geq 0 \), is a real random variable.

We show that \( N(t) \xrightarrow{t \to \infty} \infty \). All trajectories of \( N(t) \) are monotonously non-decreasing in
\[ t \geq 0, \text{ thus } \exists \lim_{t \to \infty} N(\omega, t) \text{ for all } \omega \in \Omega. \text{ Moreover it holds that } \]
\[
P(\lim_{t \to \infty} N(t) < n) = \lim_{n \to \infty} P(\lim_{t \to \infty} N(t) < n) = \lim_{n \to \infty} \lim_{t \to \infty} P(N(t) < n) \]
\[
= \lim_{n \to \infty} \lim_{t \to \infty} P(S_n > t) = \lim_{n \to \infty} \lim_{t \to \infty} P(\sum_{k=1}^{n} T_k > t) \]
\[
\leq \lim_{n \to \infty} \lim_{t \to \infty} \frac{1}{n} \sum_{k=1}^{n} P(T_k > \frac{n}{t}) = 0. \]

The equality (⋆) holds since \( \{\lim_{t \to \infty} N(t) < n\} = \\{\exists t_0 \in \mathbb{Q}_+: \forall t \geq t_0 N(t) < n\} = \bigcup_{t_0 \in \mathbb{Q}_+} \cap_{t \geq t_0} \{N(t) < n\} = \liminf_{t \to \infty} \{N(t) < n\}, \) then the continuity of the probability measure is used, where \( \mathbb{Q}_+ = \mathbb{Q} \cap \mathbb{R}_+ = \{q \in \mathbb{Q} : q \geq 0\}. \) Since for every \( \omega \in \Omega \) it holds that \( \lim_{n \to \infty} \frac{S_n}{n} = \lim_{n \to \infty} \frac{\sum_{k=1}^{n} T_k}{n} \) (the codomain of a realization of \( N(\cdot) \) is a subsequence of \( \mathbb{N} \)), it holds that \( \lim_{t \to \infty} \frac{S_n}{n} = \mu. \)

**Remark 2.1.2**

One can generalize the ergodic theorem to the case of non-identically distributed \( T_n. \) Thereby we require that \( \mu_n = \mathbb{E}T_n, \{T_n - \mu_n\}_{n \in \mathbb{N}} \) are uniformly integrable and \( \frac{1}{n} \sum_{k=1}^{n} \mu_k \xrightarrow{n \to \infty} \mu > 0. \)

Then we can prove that \( \frac{N(t)}{t} \xrightarrow{p} \frac{1}{\mu} \) (cf. [2], page 276).

**Theorem 2.1.2 (Central limit theorem):**

If \( \mu \in (0, \infty), \sigma^2 = \text{Var}T_1 \in (0, \infty) \), it holds that
\[
\mu^2 \cdot \frac{N(t) - \frac{t}{\mu}}{\sigma \sqrt{t}} \xrightarrow{d} Y, \]
where \( Y \sim \mathcal{N}(0, 1). \)

**Proof** According to the central limit theorem for sums of i.i.d. random variables (cf. theorem 7.5, WR) it holds that
\[
\frac{S_n - n\mu}{\sqrt{n\sigma^2}} \xrightarrow{n \to \infty} Y. \tag{2.1.1} \]

Let \( [x] \) be the whole part of \( x \in \mathbb{R}. \) It holds for \( a = \frac{\sigma^2}{\mu} \) that
\[
P\left( \frac{N(t) - \frac{t}{\mu}}{\sqrt{at}} \leq x \right) = P\left( N(t) \leq x\sqrt{at} + \frac{t}{\mu} \right) = P\left( S_{m(t)} > t \right), \]
where \( m(t) = \left[x\sqrt{at} + \frac{t}{\mu}\right] + 1, t \geq 0, \text{ and } \lim_{t \to \infty} m(t) = \infty. \) Therefore we get that
\[
\left| P\left( \frac{N(t) - \frac{t}{\mu}}{\sqrt{at}} \leq x \right) - \varphi(x) \right| = \left| P\left( S_{m(t)} > t \right) - \varphi(x) \right| = \left| P\left( \frac{S_{m(t)} - \mu m(t)}{\sigma \sqrt{m(t)}} > \frac{t - \mu m(t)}{\sigma \sqrt{m(t)}} \right) - \varphi(x) \right| := I_t(x)\]
for arbitrary \( t \geq 0 \) and \( x \in \mathbb{R} \), where \( \varphi \) is the distribution function of the \( \mathcal{N}(0,1) \)-distribution. For fixed \( x \in \mathbb{R} \) we introduce \( Z_t = \frac{t - \mu m(t)}{\sigma \sqrt{m(t)}} - x, \ t \geq 0 \). Then it holds that

\[
I_t(x) = \left| P\left( \frac{S_{m(t)} - \mu m(t)}{\sigma \sqrt{m(t)}} + Z_t > -x \right) - \varphi(x) \right|.
\]

If we can prove that \( Z_t \xrightarrow{t \to \infty} 0 \), then applying (2.1.1) and the theorem of Slutsky (theorem 6.4.1, WR) would result in \( \frac{S_{m(t)} - \mu m(t)}{\sigma \sqrt{m(t)}} + Z_t \xrightarrow{d \ t \to \infty} Y \sim \mathcal{N}(0,1) \) since \( Z_t \xrightarrow{t \to \infty} 0 \) a.s. results in \( Z_t \xrightarrow{t \to \infty} 0 \). Therefore we could write \( I_t(x) \xrightarrow{t \to \infty} |\tilde{\varphi}(-x) - \varphi(x)| = |\varphi(x) - \varphi(x)| = 0 \), where \( \tilde{\varphi}(x) = 1 - \varphi(x) \) is the tail function of the \( \mathcal{N}(0,1) \)-distribution, and the property of symmetry of \( \mathcal{N}(0,1) \): \( \tilde{\varphi}(-x) = \varphi(x) \), \( x \in \mathbb{R} \) was used.

Now we show that \( Z_t \xrightarrow{t \to \infty} 0 \), thus \( \frac{t - \mu m(t)}{\sigma \sqrt{m(t)}} \xrightarrow{t \to \infty} -x \). It holds that \( m(t) = x \sqrt{a t} + \frac{1}{\mu} + \varepsilon(t) \), where \( \varepsilon(t) \in [0,1) \). Then it holds that

\[
\frac{t - \mu m(t)}{\sigma \sqrt{m(t)}} = \frac{t - \mu x \sqrt{a t} - t - \mu \varepsilon(t)}{\sigma \sqrt{m(t)}} = \frac{x}{\sigma \sqrt{a t} + \frac{1}{\mu} + \varepsilon(t)} \mu - \mu \varepsilon(t) - \mu \varepsilon(t).
\]

\[\xrightarrow{t \to \infty} -x.\]

\[\Box\]

**Remark 2.1.3**

In Lineberg form, the central limit theorem can also be proven for non-identically distributed \( T_n \), cf. [2], pages 276 - 277.

**Definition 2.1.2**

The function \( H(t) = \mathbb{E} N(t), t \geq 0 \) is called renewal function of the process \( N \) (or of the sequence \( \{S_n\}_{n \in \mathbb{N}} \)).

Let \( F_T(x) = \mathbb{P}(T_1 \leq x), x \in \mathbb{R} \) be the distribution function of \( T_1 \). For arbitrary distribution functions \( F, G : \mathbb{R} \to [0,1] \) the convolution \( F \ast G \) is defined as \( F \ast G(x) = \int_{-\infty}^{\infty} F(x-y) dG(y) \). The \( k \)-fold convolution \( F^k \) of the distribution \( F \) with itself, \( k \in \mathbb{N}_0 \), is defined inductive:

\[
F^0(x) = 1(x \in [0,\infty)), \ x \in \mathbb{R},
\]

\[
F^1(x) = F(x), \ x \in \mathbb{R},
\]

\[
F^{k+1}(x) = F^k \ast F(x), \ x \in \mathbb{R}.
\]

**Lemma 2.1.1**

The renewal function \( H \) of a renewal process \( N \) is monotonously non-decreasing and right-sided continuous on \( \mathbb{R}_+ \). Moreover it holds that

\[
H(t) = \sum_{n=1}^{\infty} \mathbb{P}(S_n \leq t) = \sum_{n=1}^{\infty} F_T^n(t), \ t \geq 0.
\]

(2.1.2)
Proof The monotony and right-sided continuity of $H$ are consequences from the almost surely monotony and right-sided continuity of the trajectories of $N$. Now we prove (2.1.2):

$$H(t) = \mathbb{E}N(t) = \mathbb{E}\sum_{n=1}^{\infty} \mathbf{1}(S_n \leq t) = \sum_{n=1}^{\infty} \mathbb{E}\{S_n \leq t\} = \sum_{n=1}^{\infty} P(S_n \leq t) = \sum_{n=1}^{\infty} F_T^n(t),$$

since $P(S_n \leq t) = P(T_1 + \ldots + T_n \leq t) = F_T^n(t)$, $t \geq 0$. The equality ($\ast$) holds for all partial sums on both sides, therefore in the limit as well. \qed

Except for exceptional cases it is impossible to calculate the renewal function $H$ by the formula (2.1.2) analytically. Therefore the Laplace transform of $H$ is often used in calculations. For a monotonously (e.g. monotonously non-decreasing) right-sided continuous function $G : [0, \infty) \to \mathbb{R}$ the Laplace transform is defined as $\hat{g}(s) = \int_0^{\infty} e^{-sx}dG(x)$, $s \geq 0$. Here the integral is to be understood as the Lebesgue-Stieltjes integral, thus as a Lebesgue integral with respect to the measure $\mu_G$ on $\mathcal{B}_{\mathbb{R}}$, defined by $\mu_G([x, y]) = G(y) - G(x)$, $0 \leq x < y < \infty$, if $G$ is monotonously non-decreasing.

Just to remind you: the Laplace transform $\hat{g}$ of a random variable $X \geq 0$ is defined by

$$\hat{g}(s) = \int_0^{\infty} e^{-sx}dF_X(x), s \geq 0.$$

Lemma 2.1.2

For $s > 0$ it holds that:

$$\hat{h}(s) = \frac{\hat{T}_1(s)}{1 - \hat{T}_1(s)}.$$

Proof It holds that:

$$\hat{h}(s) = \int_0^{\infty} e^{-sx}dH(x)^{\ast} = \int_0^{\infty} e^{-sx}d\left(\sum_{n=1}^{\infty} F_T^n(x)\right) = \sum_{n=1}^{\infty} \int_0^{\infty} e^{-sx}dF_T^n(x)$$

$$= \sum_{n=1}^{\infty} \hat{T}_{n+1}\ldots\hat{T}_n(s) = \sum_{n=1}^{\infty} \left[\frac{\hat{T}_1(s)}{1 - \hat{T}_1(s)}\right]^n = \frac{\hat{T}_1(s)}{1 - \hat{T}_1(s)},$$

where for $s > 0$ it holds that $\hat{T}_1(s) < 1$ and thus the geometric series $\sum_{n=1}^{\infty} \left[\frac{\hat{T}_1(s)}{1 - \hat{T}_1(s)}\right]^n$ converges. \qed

Remark 2.1.4

If $N = \{N(t), t \geq 0\}$ is a delayed renewal process (with delay $T_1$), the statements of lemmas 2.1.1 - 2.1.2 hold in the following form:

1. $$H(t) = \sum_{n=0}^{\infty} (F_{T_1} \ast F_{T_2}^n)(t), t \geq 0,$$

where $F_{T_1}$ and $F_{T_2}$, respectively are the distribution functions of $T_1$ and $T_n$, $n \geq 2$, respectively.

2. $$\hat{h}(s) = \frac{\hat{T}_1(s)}{1 - \hat{T}_2(s)}, s \geq 0,$$ (2.1.3)

where $\hat{T}_1$ and $\hat{T}_2$ are the Laplace transforms of the distribution of $T_1$ and $T_n$, $n \geq 2$. 
For further observations we need a theorem (of Wald) about the expected value of a sum (with random number) of independent random variables.

**Definition 2.1.3**

Let \( \nu \) be a \( \mathbb{N} \)-valued random variable and be \( \{X_n\}_{n \in \mathbb{N}} \) a sequence of random variables defined on the same probability space. \( \nu \) is called independent of the future, if for all \( n \in \mathbb{N} \) the event \( \{\nu \leq n\} \) does not depend on the \( \sigma \)-algebra \( \sigma(\{X_k, \ k > n\}) \).

**Theorem 2.1.3 (Wald’s identity):**

Let \( \{X_n\}_{n \in \mathbb{N}} \) be a sequence of random variables with \( \sup \mathbb{E}|X_n| < \infty, \mathbb{E}X_n = a, \ n \in \mathbb{N} \) and be \( \nu \) a \( \mathbb{N} \)-valued random variable which is independent of the future, with \( \mathbb{E}\nu < \infty \). Then it holds that

\[
\mathbb{E}\left( \sum_{n=1}^{\nu} X_n \right) = a \cdot \mathbb{E}\nu.
\]

**Proof** Calculate \( S_n = \sum_{k=1}^{n} X_k, \ n \in \mathbb{N} \). Since \( \mathbb{E}\nu = \sum_{n=1}^{\infty} \mathbb{P}(\nu \geq n) \), the theorem follows from Lemma 2.1.3.

**Lemma 2.1.3 (Kolmogorov-Prokhorov):**

Let \( \nu \) be a \( \mathbb{N} \)-valued random variable which is independent of the future and it holds that

\[
\sum_{n=1}^{\infty} \mathbb{P}(\nu \geq n) \mathbb{E}|X_n| < \infty.
\]

(2.1.4)

Then \( \mathbb{E}S_\nu = \sum_{n=1}^{\infty} \mathbb{P}(\nu \geq n) \mathbb{E}X_n \) holds. If \( X_n \geq 0 \) a.s., then condition (2.1.4) is not required.

**Proof** It holds that \( S_\nu = \sum_{n=1}^{\nu} X_n = \sum_{n=1}^{\infty} X_n 1(\nu \geq n) \). We introduce the notation \( S_{\nu,n} = \sum_{k=1}^{n} X_k 1(\nu \geq k), \ n \in \mathbb{N} \). First, we prove the lemma for \( X_n \geq 0 \) f.s., \( n \in \mathbb{N} \). It holds \( S_{\nu,n} \uparrow S_\nu, \ n \to \infty \) for every \( \omega \in \Omega \), and thus according to the monotone convergence theorem it holds that: \( \mathbb{E}S_\nu = \lim_{n \to \infty} \mathbb{E}S_{\nu,n} = \lim_{n \to \infty} \sum_{k=1}^{n} \mathbb{E}(X_k 1(\nu \geq k)) \). Since \( \{\nu \geq k\} = \{\nu \leq k-1\}^c \) does not depend on \( \sigma(X_k) \subset \sigma(\{X_n, \ n \geq k\}) \) it holds that \( \mathbb{E}(X_k 1(\nu \geq k)) = \mathbb{E}X_k 1(\nu \geq k), \ k \in \mathbb{N}, \) and thus \( \mathbb{E}S_\nu = \sum_{n=1}^{\infty} \mathbb{P}(\nu \geq n) \mathbb{E}X_n \).

Now, let \( X_n \) be arbitrary. Take \( Y_n = |X_n|, \ Z_n = \sum_{n=1}^{n} Y_n, \ Z_{\nu,n} = \sum_{k=1}^{n} Y_k 1(\nu \geq k), \ n \in \mathbb{N} \). Since \( Y_n \geq 0, \ n \in \mathbb{N} \), it holds that \( \mathbb{E}Z_\nu = \sum_{n=1}^{\infty} \mathbb{E}(X_n 1(\nu \geq k)) < \infty \) from (2.1.4). Since \( |S_{\nu,n}| \leq Z_{\nu,n} \leq Z_\nu, \ n \in \mathbb{N}, \) according to the dominated convergence theorem of Lebesgue it holds that \( \mathbb{E}S_\nu = \lim_{n \to \infty} \mathbb{E}S_{\nu,n} = \sum_{n=1}^{\infty} \mathbb{E}X_n \mathbb{P}(\nu \geq n) \), where this series converges absolutely.

**Conclusion 2.1.1**

1. \( H(t) < \infty, \ t \geq 0 \).

2. For an arbitrary Borel measurable function \( g : \mathbb{R}_+ \to \mathbb{R}_+ \) and the renewal process \( N = \{N(t), \ t \geq 0\} \) with interarrival times \( \{T_n\}, T_n \text{ i.i.d.}, \mu = \mathbb{E}T_n \in (0, \infty) \) it holds that

\[
\mathbb{E}\left( \sum_{k=1}^{N(t)+1} g(T_k) \right) = (1 + H(t))\mathbb{E}g(T_1), \ t \geq 0.
\]

**Proof** 1. For every \( t \geq 0 \) it is obvious that \( \nu = 1 + H(t) \) does not depend on the future of \( \{T_n\}_{n \in \mathbb{N}} \), the rest follows from theorem 2.1.3 with \( X_n = g(T_n), \ n \in \mathbb{N} \).


2. For $s > 0$ consider $T_n^{(s)} = \min\{T_n, s\}, n \in \mathbb{N}$. Choose $s > 0$ such that for freely selected (but fixed) $\varepsilon > 0: \mu^{(s)} = E T_1^{(s)} \geq \mu - \varepsilon > 0$. Let $N^{(s)}$ be the renewal process which is based on the sequence $\{T_n^{(s)}\}_{n \in \mathbb{N}}$ of interarrival times: $N^{(s)}(t) = \sum_{n=1}^{\infty} I(T_n^{(s)} \leq t), t \geq 0$. It holds $N(t) \leq N^{(s)}(t), t \geq 0, a.s.$, according to conclusion 2.1.1:

$$(\mu - \varepsilon)(EN^{(s)}(t) + 1) \leq \mu^{(s)}(EN^{(s)}(t) + 1) = ES_{N^{(s)}(t)+1}^{(s)} + T_{N^{(s)}(t)+1}^{(s)} \leq t + s,$$

$t \geq 0$, where $S_n^{(s)} = T_1^{(s)} + \ldots + T_n^{(s)}, n \in \mathbb{N}$. Thus $H(t) = E N^{(s)}(t) \leq EN^{(s)}(t) \leq \frac{t + s}{\mu - \varepsilon}$, $t \geq 0$. Since $\varepsilon > 0$ is arbitrary, it holds that $\lim_{t \to \infty} \frac{H(t)}{t} \leq \frac{1}{\mu}$, and also our assertion $H(t) < \infty, t \geq 0$.

\[ \square \]

**Conclusion 2.1.2 (Elementary renewal theorem):**

For a renewal process $N$ as defined in conclusion 2.1.1, 1) it holds:

$$\lim_{t \to \infty} \frac{H(t)}{t} = \frac{1}{\mu}.$$ 

**Proof** In conclusion 2.1.1, part 2) we already proved that $\lim_{t \to \infty} \frac{H(t)}{t} \leq \frac{1}{\mu}$. If we show $\lim_{t \to \infty} \frac{H(t)}{t} \geq \frac{1}{\mu}$, our assertion would be proven. According to theorem 2.1.1 it holds that

$$\frac{1}{\mu} = E \lim_{t \to \infty} \frac{N(t)}{t} \leq \lim_{t \to \infty} E \frac{N(t)}{t} = \lim_{t \to \infty} \frac{H(t)}{t}.$$

\[ \square \]

**Remark 2.1.5** 1. We can prove that in the case of the finite second moment of $T_n$ ($\mu_2 = E T_1^2 < \infty$) we can derive a more exact asymptotics for $H(t), t \to \infty$:

$$H(t) = \frac{t}{\mu} + \frac{\mu_2}{2\mu^2} + o(1), t \to \infty.$$ 

2. The elementary renewal theorem also holds for delayed renewal processes, where $\mu = E T_2$.

We define the renewal measure $H$ on $\mathcal{B}(\mathbb{R}_+)$ by $H(B) = \sum_{n=1}^{\infty} \int_B dF_T^n(x), B \in \mathcal{B}(\mathbb{R}_+)$. It holds $H((-\infty, t]) = H((0, t]), H((s, t]) = H(t) - H(s), s, t \geq 0$, if $H$ is the renewal function as well as the renewal measure.

**Theorem 2.1.4 (Fundamental theorem of the renewal theory):**

Let $N = \{N(t), t \geq 0\}$ be a (delayed) renewal process associated with the sequence $\{T_n\}_{n \in \mathbb{N}}$, where $T_n, n \in \mathbb{N}$ are independent, $\{T_n, n \geq 2\}$ identically distributed, and the distribution of $T_2$ is not arithmetic, thus not concentrated on a regular lattice with probability 1. The distribution of $T_1$ is arbitrary. Let $E T_2 = \mu \in (0, \infty)$. Then it holds that

$$\int_0^t g(t-x) dH(x) \overset{t \to \infty}{\longrightarrow} \frac{1}{\mu} \int_0^\infty g(x) dx,$$

where $g : \mathbb{R}_+ \to \mathbb{R}$ is Riemann integrable $[0, n]$, for all $n \in \mathbb{N}$, and $\sum_{n=0}^{\infty} \max_{n \leq x \leq n+1} |g(x)| < \infty$. 
Without proof.

In particular $H((t - u, t]) \xrightarrow{t \to \infty} \frac{u}{\mu}$ holds for an arbitrary $u \in \mathbb{R}_+$, thus $H$ asymptotically (for $t \to \infty$) behaves as the Lebesgue measure.

**Abb. 2.2:**

**Definition 2.1.4**
The random variable $\chi(t) = S_{N(t)+1} - t$ is called *excess* of $N$ at time $t \geq 0$.

Obviously $\chi(0) = T_1$ holds. We now give an example of a renewal process with stationary increments.

Let $N = \{N(t), t \geq 0\}$ be a delayed renewal process associated with the sequence of interarrival times $\{T_n\}_{n \in \mathbb{N}}$. Let $F_{T_1}$ and $F_{T_2}$ be the distribution functions of the delays $T_1$ and $T_n$, $n \geq 2$. We assume that $\mu = E T_2 \in (0, \infty)$, $F_{T_2}(0) = 0$, thus $T_2 > 0$ a.s. and

$$F_{T_1}(x) = \frac{1}{\mu} \int_0^x \bar{F}_{T_2}(y) dy, \quad x \geq 0. \quad (2.1.5)$$

In this case $F_{T_1}$ is called the *integrated tail distribution function* of $T_2$.

**Theorem 2.1.5**
Under the conditions we mentioned above, $N$ is a process with stationary increments.

**Abb. 2.3:**

**Proof** Let $n \in \mathbb{N}$, $0 \leq t_0 < t_1 < \ldots < t_n < \infty$. Because $N$ does not depend on $T_n$, $n \in \mathbb{N}$ the common distribution of $(N(t_1 + t) - N(t_0 + t), \ldots, N(t_n + t) - N(t_{n-1} + t))^T$ does not depend on $t$, if the distribution of $\chi(t)$ does not depend on $t$, thus $\chi(t) \overset{d}{=} \chi(0) = T_1$, $t \geq 0$, see Figure ....
We show that \( F_{T_1} = F_{X(t)}, t \geq 0 \).

\[
F_{X(t)}(x) = P(\chi(t) \leq x) = \sum_{n=0}^{\infty} P(S_n \leq t, t < S_{n+1} \leq t + x)
\]
\[
= P(S_0 = 0 \leq t, t < S_1 = T_1 \leq t + x)
\]
\[
+ \sum_{n=1}^{\infty} \mathbb{E}(\mathbb{E}(1(S_n \leq t, t < S_n + T_{n+1} \leq t + x) \mid S_n))
\]
\[
= F_{T_1}(t + x) - F_{T_1}(t) + \sum_{n=1}^{\infty} \int_0^t P(t - y < T_{n+1} \leq t + x - y) dF_{S_n}(y)
\]
\[
= F_{T_1}(t + x) - F_{T_1}(t) + \int_0^t P(t - y < T_2 \leq t + x - y) d\left(\sum_{n=1}^{\infty} F_{S_n}(y)\right).
\]

If we can prove that \( H(y) = \frac{y}{\mu}, y \geq 0 \), then we would get

\[
F_{X(t)}(x) = F_{T_1}(t + x) - F_{T_1}(t) + \frac{1}{\mu} \int_0^t (F_{T_2}(z + x) - 1 + 1 - F_{T_2}(z)) d(-z)
\]
\[
= F_{T_1}(t + x) - F_{T_1}(t) + \frac{1}{\mu} \int_0^t (\bar{F}_{T_2}(z) - \bar{F}_{T_2}(z + x)) dz
\]
\[
= F_{T_1}(t + x) - F_{T_1}(t) + F_{T_1}(t) - \frac{1}{\mu} \int_x^{t+x} \bar{F}_{T_2}(y) dy
\]
\[
= F_{T_1}(t + x) - F_{T_1}(t + x) + F_{T_1}(x) = F_{T_1}(x), x \geq 0,
\]
according to the form (2.1.5) of the distribution of \( T_1 \).

Now we like to show that \( H(t) = \frac{t}{\mu}, t \geq 0 \). For that we use the formula (2.1.4): it holds that

\[
\hat{l}_{T_1}(s) = \frac{1}{\mu} \int_0^\infty e^{-st}(1 - F_{T_2}(t)) dt = \frac{1}{\mu} \int_0^\infty e^{-st} dt - \frac{1}{\mu} \int_0^\infty e^{-st} F_{T_2}(t) dt
\]
\[
= \frac{1}{\mu s} \left(1 + \int_0^\infty F_{T_2}(t) de^{-st}\right) = \frac{1}{\mu s} (1 + e^{-st} F_{T_2}(t) \bigg|_{-F_{T_2}(0)=0}^{\infty} - \int_0^\infty e^{-st} dF_{T_2}(t))
\]
\[
= \frac{1}{\mu s} (1 - \hat{l}_{T_2}(s)) = \hat{l}_{T_1}(s), s \geq 0.
\]

Using the formula (2.1.4) we get

\[
\hat{l}_H(s) = \frac{\hat{l}_{T_1}(s)}{1 - \hat{l}_{T_2}(s)} = \frac{1}{\mu s} = \frac{1}{\mu} \int_0^\infty e^{-st} dt = \hat{l}_{\frac{1}{\mu}}(s), s \geq 0.
\]

Since the Laplace transform of a function uniquely determines this function, it holds that \( H(t) = \frac{t}{\mu}, t \geq 0 \).

**Remark 2.1.6**

In the proof of Theorem 2.1.5 we showed that for the renewal process with delay which possesses
the distribution (2.1.5), \( H(t) \sim \frac{t}{\mu} \) not only asymptotically for \( t \to \infty \) (as in the elementary renewal theorem) but it holds \( H(t) = \frac{t}{\mu} \), for all \( t \geq 0 \). This means, per unit of the time interval we get an average of \( \frac{1}{\mu} \) renewals. For that reason such a process \( N \) is called \textit{homogeneous renewal process}.

We can prove the following theorem:

**Theorem 2.1.6**

If \( N = \{N(t), \ t \geq 0\} \) is a delayed renewal process with arbitrary delay \( T_1 \) and non-arithmetic distribution of \( T_n, \ n \geq 2, \mu = ET_2 \in (0, \infty) \), then it holds that

\[
\lim_{t \to \infty} F_{\chi(t)}(x) = \frac{1}{\mu} \int_0^x \tilde{F}_{T_2}(y) dy, \ x \geq 0.
\]

This means, the limit distribution of excess \( \chi(t), t \to \infty \) is taken as the distribution of \( T_1 \) when defining a homogeneous renewal process.

### 2.2 Poisson processes

**2.2.1 Poisson processes**

In this section we generalize the definition of a homogeneous Poisson process (see section 1.2, example 5)

**Definition 2.2.1**

The counting process \( N = \{N(t), \ t \geq 0\} \) is called Poisson process with intensity measure \( \Lambda \) if

1. \( N(0) = 0 \) a.s.
2. \( \Lambda \) is a locally finite measure \( \mathbb{R}_+ \), i.e., \( \Lambda : \mathcal{B}(\mathbb{R}_+) \to \mathbb{R}_+ \) possesses the property \( \Lambda(B) < \infty \) for every bounded set \( B \in \mathcal{B}(\mathbb{R}_+) \).
3. \( N \) possesses independent increments.
4. \( N(t) - N(s) \sim \text{Pois}(\Lambda((s,t])) \) for all \( 0 \leq s < t < \infty \).

Sometimes the Poisson process \( N = \{N(t), \ t \geq 0\} \) is defined by the corresponding random Poisson counting measure \( N = \{N(B), \ B \in \mathcal{B}(\mathbb{R}_+)\} \), i.e., \( N = (\{0,t\}), t \geq 0 \), where a counting measure is a locally finite measure with values in \( \mathbb{N}_0 \).

**Definition 2.2.2**

A random counting measure \( N = \{N(B), \ B \in \mathcal{B}(\mathbb{R}_+)\} \) is called Poissonsh with locally finite intensity measure \( \Lambda \) if

1. For arbitrary \( n \in \mathbb{N} \) and for arbitrary pairwise disjoint bounded sets \( B_1, B_2, \ldots, B_n \in \mathcal{B}(\mathbb{R}_+) \) the random variables \( N(B_1), N(B_2), \ldots, N(B_n) \) are independent.
2. \( N(B) \sim \text{Pois}(\Lambda(B)), \ B \in \mathcal{B}(\mathbb{R}_+), \ B\text{-bounded.} \)

It is obvious that properties 3 and 4 of definition 2.2.1 follow from properties 1 and 2 of definition 2.2.2. Property 1 of definition 2.2.1 however is an autonomous assumption. \( N(B), B \in \mathcal{B}(\mathbb{R}_+) \) is interpreted as the number of points of \( N \) within the set \( B \).
Remark 2.2.1
As stated in definition 2.2.2, a Poisson counting measure can also be defined on an arbitrary topological space $E$ equipped with the Borel-$\sigma$-algebra $\mathcal{B}(E)$. Very often $E = \mathbb{R}^d$, $d \geq 1$ is chosen in applications.

Lemma 2.2.1
For every locally finite measure $\Lambda$ on $\mathbb{R}_+$ there exists a Poisson process with intensity measure $\Lambda$.

Proof If such a Poisson process had existed, the characteristic function $\varphi_{N(t)-N(s)}(\cdot)$ of the increment $N(t) - N(s)$, $0 \leq s < t < \infty$ would have been equal to $\varphi_{s,t}(z) = \varphi\text{Pois}(\Lambda((s,t]))(z) = e^{\Lambda((s,t])(e^{iz} - 1)}$, $z \in \mathbb{R}$ according to property 4 of definition 2.2.1. We show that the family of characteristic functions $\{\varphi_{s,t}, 0 \leq s < t < \infty\}$ possesses property 1.7.1: for all $n : 0 \leq s < u < t$, $\varphi_{s,u}(z)\varphi_{u,t}(z) = e^{\Lambda((s,u])(e^{iz} - 1)}e^{\Lambda((u,t])(e^{iz} - 1)} = e^{\Lambda((s,u])+(\Lambda((u,t]))(e^{iz} - 1)} = e^{\Lambda((s,t])(e^{iz} - 1)} = \varphi_{s,t}(z)$, $z \in \mathbb{R}$ since the measure $\Lambda$ is additive. Thus, the existence of the Poisson process $N$ follows from theorem 1.7.1.

Remark 2.2.2
The existence of a Poisson counting measure can be proven with the help of the theorem of Kolmogorov, yet in a more general form than in theorem 1.1.2.

From the properties of the Poisson distribution it follows that $\mathbb{E}N(B) = \text{Var} N(B) = \Lambda(B)$, $B \in \mathcal{B}(\mathbb{R}_+)$. Thus $\Lambda(B)$ is interpreted as the mean number of points of $N$ within the set $B$, $B \in \mathcal{B}(\mathbb{R}_+)$. We get an important special case if $\Lambda(dx) = \lambda dx$ for $\lambda \in (0, \infty)$, i.e., $\Lambda$ is proportional to the Lebesgue measure $\nu_1$ on $\mathbb{R}_+$. Then we call $\lambda = \mathbb{E}N(1)$ the intensity of $N$.

Soon we will prove that in this case $N$ is a homogeneous Poisson process with intensity $\lambda$. To remind you: In section 1.2 the homogeneous Poisson process was defined as a renewal process with interarrival times $T_N \sim \text{Exp}(\lambda)$: $N(t) = \sup\{n \in \mathbb{N} \mid S_n \leq t\}$, $S_n = T_1 + \ldots + T_n$, $n \in \mathbb{N}$, $t \geq 0$.

Exercise 2.2.1
Show that the homogeneous Poisson process is a homogeneous renewal process with $T_1 \overset{d}{=} T_2 \sim \text{Exp}(\lambda)$. Hint: you have to show that for an arbitrary exponential distributed random variable $X$ the integrated tail distribution function of $X$ is equal to $F_X$.

Theorem 2.2.1
Let $N = \{N(t), t \geq 0\}$ be a counting process. The following statements are equivalent.

1. $N$ is a homogeneous Poisson process with intensity $\lambda > 0$.
2. a) $N(t) \sim \text{Pois}(\lambda t)$, $t \geq 0$
   b) for an arbitrary $n \in \mathbb{N}$, $t \geq 0$, it holds that the random vector $(S_1, \ldots, S_n)$ under condition $\{N(t) = n\}$ possesses the same distribution as the order statistics of i.i.d. random variables $U_i \in \mathcal{U}([0,t])$, $i = 1, \ldots, n$.
3. a) $N$ has independent increments,
   b) $\mathbb{E}N(1) = \lambda$, and
   c) property 2b) holds.
4. a) $N$ has stationary and independent increments, and  
   b) $P(N(t) = 0) = 1 - \lambda t + o(t)$, $P(N(t) = 1) = \lambda t + o(t)$, $t \downarrow 0$ holds.

5. a) $N$ has stationary and independent increments,  
   b) property 2a) holds.

**Remark 2.2.3**  
1. It is obvious that Definition 2.2.1 with $\Lambda(dx) = \lambda dx$, $\lambda \in (0, \infty)$ is an equivalent definition of the homogeneous Poisson process according to Theorem 2.2.1.

2. The homogeneous Poisson process $N$ was introduced in the beginning of the 20th century from the physicists A. Einstein and M. Smoluchovsky to be able to model the counting process of elementary particle in the Geiger counter.

3. From 4b) it follows $P(N(t) > 1) = o(t)$, $t \downarrow 0$.

4. The intensity of $N$ has the following interpretation: $\lambda = EN(1) = \frac{1}{\mathbb{E}T_1}$, thus the mean number of renewals of $N$ within a time interval with length 1.

5. The renewal function of the homogeneous Poisson process is $H(t) = \lambda t$, $t \geq 0$. Thereby $H(t) = \Lambda([0,t])$, $t > 0$ holds.

**Proof** Structure of the proof: 1) $\Rightarrow$ 2) $\Rightarrow$ 3) $\Rightarrow$ 4) $\Rightarrow$ 5) $\Rightarrow$ 1)

1) $\Rightarrow$ 2):

From 1) follows $S_n = \sum_{k=1}^{n} T_k \sim \text{Erl}(n, \lambda)$ since $T_k \sim \text{Pois}(\lambda)$, $n \in \mathbb{N}$, thus $P(N(t) = 0) = P(T_1 > t) = e^{-\lambda t}$, $t \geq 0$, and for $n \in \mathbb{N}$

$$P(N(t) = n) = P\{\{N(t) \geq n\} \setminus \{N(t) \geq n + 1\}\} = P(N(t) \geq n) - P(N(t) \geq n + 1)$$

$$= P(S_n \leq t) - P(S_{n+1} \leq t) = \int_0^t \frac{\lambda^n x^{n-1}}{(n-1)!} e^{-\lambda x} dx - \int_0^t \frac{\lambda^{n+1} x^n}{n!} e^{-\lambda x} dx$$

$$= \int_0^t \frac{d}{dx} \left( \frac{(\lambda x)^n}{n!} e^{-\lambda x} \right) dx = \frac{(\lambda t)^n}{n!} e^{-\lambda t}, \quad t \geq 0.$$

Thus 2a) is proven.

Now let’s prove 2b). According to the transformation theorem of random variables (cf. theorem 3.6.1, WR), it follows from

$$\begin{align*}
S_1 &= T_1 \\
S_2 &= T_1 + T_2 \\
&\vdots \\
S_{n+1} &= T_1 + \ldots + T_{n+1}
\end{align*}$$

that the density $f(S_1,\ldots,S_{n+1})$ of $(S_1,\ldots,S_{n+1})^\top$ can be expressed by the density of $(T_1,\ldots,T_{n+1})^\top$, $T_i \sim \text{Exp}(\lambda)$, i.i.d.:

$$f(S_1,\ldots,S_{n+1})(t_1,\ldots,t_{n+1}) = \prod_{k=1}^{n+1} f_{T_k}(t_k - t_{k-1}) = \prod_{k=1}^{n+1} \lambda e^{-\lambda(t_k-t_{k-1})} = \lambda^{n+1} e^{-\lambda t_{n+1}}$$

for arbitrary $0 \leq t_1 \leq \ldots \leq t_{n+1}$, $t_0 = 0$.

For all other $t_1,\ldots,t_{n+1}$ it holds $f(S_1,\ldots,S_{n+1})(t_1,\ldots,t_{n+1}) = 0$. 
Therefore

\[
\begin{align*}
  f(s_1, \ldots, s_n|N(t) = n) &= f(s_1, \ldots, s_n)(t_1, \ldots, t_n|S_k \leq t, \ k \leq n, \ S_{n+1} > t) \\
  &= \frac{\int_t^\infty f(s_1, \ldots, s_{n+1})(t_1, \ldots, t_{n+1})dt_{n+1}}{f_0(t_1) \cdots f_{t-1}(t_1, \ldots, t_n) \int_t^\infty f(s_1, \ldots, s_{n+1})(t_1, \ldots, t_{n+1})dt_{n+1}dt_1} \\
  &= \frac{\int_t^\infty \lambda_{n+1}e^{-\lambda_{n+1}t}dt_{n+1}}{\int_0^t f(s_1) \cdots f_{t-1}(s_1, \ldots, s_n)\lambda_n e^{-\lambda_n t}dt_1} \\
  &= \frac{\int_t^\infty \lambda_{n+1}e^{-\lambda_{n+1}t}dt_{n+1}}{\lambda_n e^{-\lambda_n t}} \times I(0 \leq t_1 \leq t_2 \leq \ldots \leq t_n \leq t) \\
  &= \frac{n!}{t_n}I(0 \leq t_1 \leq t_2 \leq \ldots \leq t_n \leq t).
\end{align*}
\]

This is exactly the density of \( n \) i.i.d. \( \mathcal{U}([0, t]) \)-random variables.

**Exercise 2.2.2**

Proof this.

2) \( \Rightarrow \) 3)

From 2a) obviously follows 3b). Now we just have to prove the independence of the increments of \( N \). For an arbitrary \( n \in \mathbb{N}, x_1, \ldots, x_n \in \mathbb{N}, t_0 = 0 < t_1 < \ldots < t_n \) for \( x = x_1 + \ldots + x_n \) it holds that

\[
P(\bigcap_{k=1}^n \{N(t_k) - N(t_{k-1}) = x_k\}) = \prod_{k=1}^n \left( \frac{\lambda(t_k - t_{k-1})^{x_k}}{x_k!} e^{-\lambda(t_k - t_{k-1})} \right) \]

since the probability of \((*)\) belongs to the polynomial distribution with parameters \( n, \left\{ \frac{t_k - t_{k-1}}{t_n} \right\}_{k=1}^n \).

Because the event \((*)\) is that at the independent uniformly distributed toss of \( x \) points on \([0, t]\), exactly \( x_k \) points occur within the basket of length \( t_k - t_{k-1}, k = 1, \ldots, n; \)

\[
\text{Abb. 2.4:}
\]

Thus 3a) is proven since \( P(\bigcap_{k=1}^n \{N(t_k) - N(t_{k-1}) = x_k\}) = \prod_{k=1}^n P(\{N(t_k) - N(t_{k-1}) = x_k\}). \)
We prove that $N$ possesses stationary increments. For an arbitrary $n \in \mathbb{N}_0$, $x_1, \ldots, x_n \in \mathbb{N}$, $t_0 = 0 < t_1 < \ldots < t_n$ and $h > 0$ we consider $I(h) = P(\cap_{k=1}^n \{ N(t_k + h) - N(t_{k-1} + h) = x_k \})$ and show that $I(h)$ does not depend on $h \in \mathbb{R}$. According to the formula of the total probability it holds that

$$I(h) = \sum_{m=0}^{\infty} P(\cap_{k=1}^n \{ N(t_k + h) - N(t_{k-1} + h) = x_k \} | N(t_n + h) = m) \cdot P(N(t_n + h) = m)$$

$$= \sum_{m=0}^{\infty} \frac{m!}{x_1! \cdots x_n!} \prod_{k=1}^{n} \left( \frac{t_k + h - t_{k-1} - h}{t_n + h - h} \right)^{x_k} e^{-\lambda(t_n+h)} (\lambda(t_n+h))^m \frac{1}{m!}$$

$$= \sum_{m=0}^{\infty} P(\cap_{k=1}^n \{ N(t_k) - N(t_{k-1}) = x_k \} | N(t_n + h) = m) \times P(N(t_n + h) = m) = I(0).$$

We now show property 4b) for $h \in (0, 1)$:

$$P(N(h) = 0) = \sum_{k=0}^{\infty} P(N(h) = 0, N(1) = k) = \sum_{k=0}^{\infty} P(N(h) = 0, N(1) - N(h) = k)$$

$$= \sum_{k=0}^{\infty} P(N(1) - N(h) = k, N(1) = k)$$

$$= \sum_{k=0}^{\infty} P(N(1) = k) P(N(1) - N(h) = k | N(1) = k)$$

$$= \sum_{k=0}^{\infty} P(N(1) = k)(1-h)^k.$$ 

We have to show that $P(N(h) = 0) = 1 - \lambda h + o(h)$, i.e., $\lim_{h \to 0} \frac{1}{h} (1 - P(N(h) = 0)) = \lambda$. Indeed it holds that

$$\frac{1}{h} (1 - P(N(h) = 0)) = \frac{1}{h} \left( 1 - \sum_{k=0}^{\infty} P(N(1) = k)(1-h)^k \right) = \sum_{k=1}^{\infty} P(N(1) = k) \cdot \frac{1 - (1-h)^k}{h}$$

$$\xrightarrow{h \to 0} \sum_{k=1}^{\infty} P(N(1) = k) \lim_{h \to 0} \frac{1 - (1-h)^k}{h} = \sum_{k=0}^{\infty} P(N(1) = k) k = E N(1) = \lambda,$$

since the series uniformly converges in $h$ because it is dominated by $\sum_{k=0}^{\infty} P(N(1) = k) k = \lambda < \infty$ because of the inequality $(1-h)^k \geq 1 - kh$, $h \in (0, 1)$, $k \in \mathbb{N}$.

Similarly one can show that $\lim_{h \to 0} \frac{P(N(h) = 1)}{h} = \lim_{h \to 0} \sum_{k=1}^{\infty} P(N(1) = k) k(1-h)^{k-1} = \lambda$. 4) ⇒ 5)

We have to show that for an arbitrary $n \in \mathbb{N}$ and $t \geq 0$

$$p_n(t) = P(N(t) = n) = e^{-\lambda t} \frac{(\lambda t)^n}{n!} \quad (2.2.1)$$

holds. We will prove that by induction with respect to $n$. First we show that $p_0(t) = e^{-\lambda t}$, $h = 0$. For that we consider $p_0(t + h) = P(N(t + h) = 0) = P(N(t) = 0, N(t + h) - N(t) =
0) = p_0(t)p_0(h) = p_0(t)(1 - \lambda h + o(h))$, $h \to 0$. Similarly one can show that $p_0(t) = p_0(t-h)(1 - \lambda h + o(h))$, $h \to 0$. Thus $p_0'(t) = \lim_{h \to 0} \frac{p_0(t+h) - p_0(t)}{h} = -\lambda p_0(t)$, $t > 0$ holds. Since $p_0(0) = \mathbb{P}(N(0) = 0) = 1$, it follows from
\[
\begin{cases}
  p_0'(t) = -\lambda p_0(t) \\
p_0(0) = 1,
\end{cases}
\]
that it exists an unique solution $p_0(t) = e^{-\lambda t}$, $t \geq 0$. Now for $n$ the formulair (2.2.1) be approved. Let’s prove it for $n + 1$.
\[
\begin{align*}
p_{n+1}(t+h) &= \mathbb{P}(N(t + h) = n + 1) \\
 &= \mathbb{P}(N(t) = n, N(t + h) - N(t) = 1) + \mathbb{P}(N(t) = n + 1, N(t + h) - N(t) = 0) \\
 &= p_n(t) \cdot p_1(h) + p_{n+1}(t) \cdot p_0(h) \\
 &= p_n(t)(\lambda h + o(h)) + p_{n+1}(t)(1 - \lambda h + o(h)), 
\end{align*}
\]
Thus
\[
\begin{align*}
p_{n+1}'(t) &= -\lambda p_{n+1}(t) + \lambda p_n(t), \\
p_{n+1}(0) &= 0 @2.2.2.
\end{align*}
\]
Since $p_n(t) = e^{-\lambda t}(\frac{\lambda t}{n!})^n$, we obtain $p_{n+1}(t) = e^{-\lambda t}(\frac{\lambda t}{(n+1)!})$ as solution of (2.2.2). (Indeed
\[
p_{n+1}(t) = C(t)e^{-\lambda t} \Rightarrow C'(t)e^{-\lambda t} = \lambda C(t)e^{-\lambda t} \ldots \ldots + \lambda p_n(t)
\]
\[
C'(t) = \frac{\lambda^{n+1}t^n}{n!} \Rightarrow C(t) = \frac{\lambda^{n+1}t^{n+1}}{(n+1)!}, C(0) = 0
\]
5) $\Rightarrow 1)$. Let $N$ be a counting process $N(t) = \max\{n : S_n \leq t\}$, $t \geq 0$, which fulfills conditions 5a) and 5b). We show that $S_n = \sum_{k=1}^{n} T_k$, where $T_k$ i.i.d. with $T_k \sim \text{Exp} (\lambda)$, $k \in \mathbb{N}$. Since $T_k = S_k - S_{k-1}, k \in \mathbb{N}, S_0 = 0$, we consider for $b_0 = 0 \leq a_1 < b_1 \leq \ldots \leq a_n < b_n$
\[
\mathbb{P}\left(\cap_{k=1}^{n} \{a_k < S_k \leq b_k\}\right) = \mathbb{P}(\cap_{k=1}^{n-1} \{N(a_k) - N(b_{k-1}) = 0, N(b_k) - N(a_k) = 1\} \\
\cap \{N(a_n) - N(b_{n-1}) = 0, N(b_n) - N(a_n) \geq 1\}) \\
= \prod_{k=1}^{n-1} \left(\mathbb{P}(N(a_k) - b_{k-1}) = 0) \mathbb{P}(N(b_k) - a_k) = 1) \times \mathbb{P}(N(a_n) - b_{n-1}) = 0) \mathbb{P}(N(b_n) - a_n) \geq 1\right) \\
= e^{-\lambda(a_n-b_{n-1})} \left(1-e^{-\lambda(b_n-a_n)}\right) \prod_{k=1}^{n-1} \lambda(b_k-a_k)e^{-\lambda(b_k-b_{k-1})} \\
= \lambda^{n-1}(e^{-\lambda a_n} - e^{-\lambda b_n}) \prod_{k=1}^{n-1} (b_k - a_k) = \int_{a_1}^{b_1} \ldots \int_{a_n}^{b_n} \lambda^n e^{-\lambda y_n} dy_n \ldots y_1.
\]
The common density of $(S_1, \ldots, S_n)^\top$ therefore is given by $\lambda^n e^{-\lambda y_n} 1(y_1 \leq y_2 \leq \ldots \leq y_n)$. \hfill \square

### 2.2.2 Compound Poisson process

**Definition 2.2.3**

Let $N = \{N(t), t \geq 0\}$ be a homogeneous Poisson process with intensity $\lambda > 0$, build by
means of the sequence \( \{T_n\}_{n \in \mathbb{N}} \) of interarrival times. Let \( \{U_n\}_{n \in \mathbb{N}} \) be a sequence of i.i.d. random variables, independent of \( \{T_n\}_{n \in \mathbb{N}} \). Let \( F_U \) be the distribution function of \( U_1 \). For an arbitrary \( t \geq 0 \) let \( X(t) = \sum_{k=1}^{N(t)} U_k \). The stochastic process \( X = \{X(t), \ t \geq 0\} \) is called **compound Poisson process with parameters** \( \lambda, F_U \). The distribution of \( X(t) \) thereby is called **compound Poisson distribution with parameters** \( \lambda t, F_U \).

The compound Poisson process \( X(t), t \geq 0 \) can be interpreted as the sum of „marks“ \( U_n \) of a homogeneous marked Poisson process \( (N, U) \) until time \( t \).

In queueing theory \( X(t) \) is interpreted as the overall workload of a server until time \( t \) if the requests to the service occur at times \( S_n = \sum_{k=1}^{n} T_k, n \in \mathbb{N} \) and represent the amount of work \( U_n, n \in \mathbb{N} \).

In actuarial mathematics \( X(t), t \geq 0 \) is the total damage in a portfolio until time \( t \geq 0 \) with number of damages \( N(t) \) and amount of loss \( U_n, n \in \mathbb{N} \).

**Theorem 2.2.2**

Let \( X = \{X(t), t \geq 0\} \) be a compound Poisson process with parameters \( \lambda, F_U \). The following properties hold:

1. \( X \) has independent and stationary increments.

2. If \( \hat{m}_U(s) = E e^{s U_1}, s \in \mathbb{R} \), is the moment generating function of \( U_1 \), such that \( \hat{m}_U(s) < \infty, \ s \in \mathbb{R} \), then it holds that

\[
\hat{m}_{X(t)}(s) = e^{\lambda t (\hat{m}_U(s) - 1)}, \ s \in \mathbb{R}, \ t \geq 0, \quad E X(t) = \lambda t E U_1, \quad \text{Var} X(t) = \lambda t E U_1^2, \ t \geq 0.
\]

**Proof**

1. We have to show that for arbitrary \( n \in \mathbb{N} \), \( 0 \leq t_0 < t_1 < \ldots < t_n \) and \( h \)

\[
P \left( \sum_{i_1=N(t_0+h)+1}^{N(t_1+h)} U_{i_1} \leq x_1, \ldots, \sum_{i_n=N(t_{n-1}+h)+1}^{N(t_n+h)} U_{i_n} \leq x_n \right) = \prod_{k=1}^{n} P \left( \sum_{i_k=N(t_{k-1}+h)+1}^{N(t_k)} U_{i_k} \leq x_k \right)
\]

for arbitrary \( x_1, \ldots, x_n \in \mathbb{R} \). Indeed it holds that

\[
P \left( \sum_{i_1=N(t_0+h)+1}^{N(t_1+h)} U_{i_1} \leq x_1, \ldots, \sum_{i_n=N(t_{n-1}+h)+1}^{N(t_n+h)} U_{i_n} \leq x_n \right)
\]

\[
= \sum_{k_1, \ldots, k_n=0}^{\infty} \left( \prod_{j=1}^{n} F^{sk_j}(x_j) \right) P \left( \cap_{m=1}^{n} \{N(t_m + h) - N(t_{m-1} + h) = k_m\} \right)
\]

\[
= \sum_{k_1, \ldots, k_n=0}^{\infty} \left( \prod_{j=1}^{n} F^{sk_j}(x_j) \right) \left( \prod_{m=1}^{n} P(N(t_m) - N(t_{m-1}) = k_m) \right)
\]

\[
= \prod_{m=1}^{n} \sum_{k_m=0}^{\infty} F^{sk_m}(x_m) P(N(t_m) - N(t_{m-1}) = k_m)
\]

\[
= \prod_{m=1}^{n} P \left( \sum_{k_m=N(t_{m-1})+1}^{N(t_m)} \leq x_m \right)
\]

2. **Exercise 2.2.3**

\[\square\]
2.2.3 Cox process

A Cox process is a (in general inhomogeneous) Poisson process with intensity measure $\Lambda$ which as such is a random measure. The intuitive idea is stated in the following definition.

**Definition 2.2.4**

Let $\Lambda = \{\Lambda(B), \ B \in B(\mathbb{R}_+)\}$ be a random a.s. locally finite measure. The random counting measure $N = \{N(B),\ B \in B(\mathbb{R}_+)\}$ is called Cox counting measure (or doubly stochastic Poisson measure) with random intensity measure $\Lambda$ if for arbitrary $n \in \mathbb{N}$, $k_1,\ldots,k_n \in \mathbb{N}_0$ and $0 \leq a_1 < b_1 \leq a_2 < b_2 \leq \ldots \leq a_n < b_n$ it holds that $P(\bigcap_{i=1}^n \{N((a_i,b_i]) = k_i\}) = \mathbb{E}\left(\prod_{i=1}^n e^{-\Lambda((a_i,b_i])}\frac{\Lambda^k((a_i,b_i])}{k!}\right)$. The process $\{N(t), t \geq 0\}$ with $N(t) = N((0,t])$ is called Cox process (or doubly stochastic Poisson process) with random intensity measure $\Lambda$.

**Example 2.2.1**  
1. If the random measure $\Lambda$ is a.s. absolutely continuous with respect to the Lebesgue measure, i.e., $\Lambda(B) = \int_B \lambda(t)dt$, $B$ - bounded, $B \in B(\mathbb{R}_+)$, where $\{\lambda(t), t \geq 0\}$ is a stochastic process with a.s. Borel-measurable Borel-integrable trajectories, then $\lambda(t) \geq 0$ a.s. for all $t \geq 0$ is called the *intensity process* of $N$.

2. In particular, it can be that $\lambda(t) \equiv Y$ where $Y$ is a non-negative random variable. Then it holds that $\Lambda(B) = Y\nu_1(B)$, thus $N$ has a random intensity $Y$. Such Cox processes are called *mixed Poisson processes*.

A Cox process $N = \{N(t), t \geq 0\}$ with intensity process $\{\lambda(t), t \geq 0\}$ can be build explicitly as the following. Let $\tilde{N} = \{\tilde{N}(t), t \geq 0\}$ be a homogeneous Poisson process with intensity 1, which is independent of $\{\lambda(t), t \geq 0\}$. Then $N \overset{d}{=} N_1$, where the process $N_1 = \{N_1(t), t \geq 0\}$ is given by $N_1(t) = \tilde{N}(\int_0^t \lambda(y)dy), t \geq 0$. The assertion $N \overset{d}{=} N_1$ of course has to be proven. However, we shall assume it without proof. It is also the basis for the simulation of the Cox process $N$.

2.3 Additional exercises

**Exercise 2.3.1**

Let $\{N(t)\}_{t \geq 0}$ be a renewal process with interarrival times $T_i$, which are exponentially distributed, i.e. $T_i \sim \text{Exp}(\lambda)$.

- a) Prove that $N(t)$ is Poisson distributed for every $t > 0$.
- b) Determine the parameter of this Poisson distribution.
- c) Determine the renewal function $H(t) = E\ N(t)$.

**Exercise 2.3.2**

Prove that a (real-valued) stochastic process $X = \{X(t), t \in [0,\infty)\}$ with independent increments already has stationary increments if the distribution of the random variable $X(t + h) - X(h)$ does not depend on $h$.

**Exercise 2.3.3**

Let $N = \{N(t), t \in [0,\infty)\}$ be a Poisson process with intensity $\lambda$. Calculate the probabilities that within the interval $[0, s]$ exactly $i$ events occur under the condition that within the interval $[0, t]$ exactly $n$ events occur, i.e. $P(N(s) = i \mid N(t) = n)$ for $s < t, i = 0, 1,\ldots, n$. 
Exercise 2.3.4
Let \( N^{(1)} = \{N^{(1)}(t), t \in [0, \infty)\} \) and \( N^{(2)} = \{N^{(2)}(t), t \in [0, \infty)\} \) be independent Poisson processes with intensities \( \lambda_1 \) and \( \lambda_2 \). In this case the independence indicates that the sequences \( T^{(1)}_1, T^{(1)}_2, \ldots \) and \( T^{(2)}_1, T^{(2)}_2, \ldots \) are independent. Show that \( N = \{N(t) := N^{(1)}(t) + N^{(2)}(t), t \in [0, \infty)\} \) is a Poisson process with intensity \( \lambda_1 + \lambda_2 \).

Exercise 2.3.5 (Queuing paradox):
Let \( N = \{N(t), t \in [0, \infty)\} \) be a renewal process. Then \( T(t) = S_{N(t)+1} - t \) is called the time of excess, \( C(t) = t - S_{N(t)} \) the current lifetime and \( D(t) = T(t) + C(t) \) the lifetime at time \( t > 0 \).

Now let \( N = \{N(t), t \in [0, \infty)\} \) be a Poisson process with intensity \( \lambda \).

a) Calculate the distribution of the time of excess \( T(t) \).

b) Show that the distribution of the current lifetime is given by \( P(C(t) = t) = e^{-\lambda t} \) and the density is given by \( f_{C(t)\mid N(t)>0}(s) = \lambda e^{-\lambda s}1\{s \leq t\} \).

c) Show that \( P(D(t) \leq x) = (1 - (1 + \lambda \min\{t, x\})e^{-\lambda x})1\{x \geq 0\} \).

d) To determine \( ET(t) \), one could argue like this: On average \( t \) lies in the middle of the surrounding interval of interarriving time \( (S_{N(t)}, S_{N(t)+1}) \), i.e. \( ET(t) = \frac{1}{2} E(S_{N(t)+1} - S_{N(t)}) = \frac{1}{2} E(T_{N(t)+1}) = \frac{1}{2\lambda} \). Considering the result from part (a) this reasoning is false. Where is the mistake in the reasoning?

Exercise 2.3.6
Let \( X = \{X(t) := \sum_{i=1}^{N(t)} U_i, t \geq 0\} \) be a compound Poisson process. Let \( M_{N(t)}(s) = ES^{N(t)} \), \( s \in (0, 1) \), be the generating function of the Poisson processes \( N(t) \), \( \mathcal{L}\{U\}(s) = \exp\{-sU\} \) the Laplace Transform of \( U_i, i \in \mathbb{N} \), and \( \mathcal{L}\{X(t)\}(s) \) the Laplace Transform of \( X(t) \). Prove that
\[
\mathcal{L}\{X(t)\}(s) = M_{N(t)}(\mathcal{L}\{U\}(s)), \quad s \geq 0.
\]

Exercise 2.3.7
Let \( X = \{X(t), t \in [0, \infty)\} \) be a compound Poisson process with \( U_i \) i.i.d., \( U_i \sim \text{Exp}(\gamma) \), where the intensity of \( N(t) \) is given by \( \lambda \). Show that for the Laplace transform \( \mathcal{L}\{X(t)\}(s) \) of \( X(t) \) it holds:
\[
\mathcal{L}\{X(t)\}(s) = \exp\left\{ -\frac{\lambda s}{\gamma + s} \right\}.
\]

Exercise 2.3.8
Write a function in \( \text{R} \) (alternatively: Java) to which we pass time \( t \), intensity \( \lambda \) and a value \( \gamma \) as parameters. The return of the function is a random value of the compound Poisson process with characteristics \((\lambda, \text{Exp}(\gamma))\) at time \( t \).

Exercise 2.3.9
Let the stochastic process \( N = \{N(t), t \in [0, \infty)\} \) be a Cox process with intensity function \( \lambda(t) = Z \), where \( Z \) is a discrete random variable which takes values \( \lambda_1 \) and \( \lambda_2 \) with probabilities 1/2. Determine the moment generating function as well as the expected value and the variance of \( N(t) \).

Exercise 2.3.10
Let \( N^{(1)} = \{N^{(1)}(t), t \in [0, \infty)\} \) and \( N^{(2)} = \{N^{(2)}(t), t \geq 0\} \) be two independent homogeneous Poisson processes with intensities \( \lambda_1 \) and \( \lambda_2 \). Moreover, let \( X \geq 0 \) be an arbitrary non-negative
random variable which is independent of $N^{(1)}$ and $N^{(2)}$. Show that the process $N = \{N(t), t \geq 0\}$ with

$$N(t) = \begin{cases} N^{(1)}(t), & t \leq X, \\ N^{(1)}(X) + N^{(2)}(t - X), & t > X \end{cases}$$

is a Cox process whose intensity process $\lambda = \{\lambda(t), t \geq 0\}$ is given by

$$\lambda(t) = \begin{cases} \lambda_1, & t \leq X, \\ \lambda_2, & t > X. \end{cases}$$
3 Wiener process

3.1 Elementary properties

In Example 2) of Section 1.2 we defined the Brownian motion (or Wiener process) $W = \{W(t), \ t \geq 0\}$ as an Gaussian process with $\mathbb{E}W(t) = 0$ and $\text{cov}(W(s), W(t)) = \min\{s, t\}$, $s, t \geq 0$. The Wiener process is called after the mathematician Norbert Wiener (1894 - 1964).

Why does the Brownian motion exist? According to theorem of Kolmogorov (Theorem 1.1.2) it exists a real-valued Gaussian process $X = \{X(t), \ t \geq 0\}$ with mean value $\mathbb{E}X(t) = \mu(t)$, $t \geq 0$, and covariance function $\text{cov}(X(s), X(t)) = C(s, t)$, $s, t \geq 0$ for every function $\mu: \mathbb{R}_+ \to \mathbb{R}$ and every positive semidefinite function $C: \mathbb{R}_+ \times \mathbb{R}_+ \to \mathbb{R}$. We just have to show that $C(s, t) = \min\{s, t\}$, $s, t \geq 0$ is positive semidefinite.

Exercise 3.1.1
Prove this!

We now give a new (equivalent) definition.

Definition 3.1.1
A stochastic process $W = \{W(t), \ t \geq 0\}$ is called Wiener process (or Brownian motion) if

1. $W(0) = 0$ a.s.
2. $W$ possesses independent increments
3. $W(t) - W(s) \sim \mathcal{N}(0, t - s)$, $0 \leq s < t$

The existence of $W$ according to Definition 3.1.1 follows from Theorem 1.7.1 since $\varphi_{s,t}(z) = \mathbb{E}e^{iz(W(t) - W(s))} = e^{-(t-s)^2/2}$, $z \in \mathbb{R}$, and $e^{-(t-s)^2/2}e^{-(u-s)^2/2} = e^{-(t-u)^2/2}$ for $0 \leq s < u < t$, thus $\varphi_{s,u}(z)\varphi_{u,t}(z) = \varphi_{s,t}(z)$, $z \in \mathbb{R}$. From Theorem 1.3.1 the existence of a version with continuous trajectories follows.

Exercise 3.1.2
Show that Theorem 1.3.1 holds for $\alpha = 3$, $\sigma = \frac{1}{2}$.

Therefore, it is often assumed that the Wiener process possesses continuous paths (just take its corresponding version).

Theorem 3.1.1
Both definitions of the Wiener process are equivalent.

Proof 1. From definition in Section 1.2 follows Definition 3.1.1. $W(0) = 0$ a.s. follows from $\text{Var}(W(0)) = \min\{0, 0\} = 0$. Now we prove that the increments of $W$ are independent. If $Y \sim \mathcal{N}(\mu, K)$ is a n-dimensional Gaussian random vector and $A$ a $(n \times n)$-matrix, then $AY \sim \mathcal{N}(A\mu, AKAT)$ holds, this follows from the explicit form of the characteristic function of $Y$. Now let $n \in \mathbb{N}$, $0 = t_0 < t_1 < \ldots < t_n$, $Y =$
The process \( W(t_0), W(t_1), \ldots, W(t_n) \)\( ^\top \). For \( Z = (W(t_0), W(t_1) - W(t_0), \ldots, W(t_n) - W(t_{n-1})) \)\( ^\top \) it holds that \( Z = AY \), where

\[
A = \begin{pmatrix}
1 & 0 & 0 & \ldots & \ldots & 0 \\
-1 & 1 & 0 & \ldots & \ldots & 0 \\
0 & -1 & 1 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & 0 & \ldots & -1 & 1
\end{pmatrix}.
\]

Thus \( Z \) is also Gaussian with a covariance matrix which is diagonal. Indeed, it holds\( \text{cov}(W(t_{i+1}) - W(t_i), W(t_{j+1}) - W(t_j)) = \min\{t_{i+1}, t_{j+1}\} - \min\{t_i, t_j\} + \min\{t_i, t_j\} = 0 \) for \( i \neq j \). Thus the coordinates of \( Z \) are uncorrelated, which means independence in case of a multivariate Gaussian distribution. Thus the increments of \( W \) are independent. Moreover, for arbitrary \( 0 \leq s < t \) it holds that \( W(t) - W(s) \sim \mathcal{N}(0, t-s) \). The normal distribution follows since \( \mathbb{E}W(t) = \mathbb{E}W(s) = 0 \) and \( \text{Var}(W(t) - W(s)) = \text{Var}(W(t)) - 2 \text{cov}(W(s), W(t)) + \text{Var}(W(s)) = t - 2 \min\{s, t\} + s = t - s. \)

2. From Definition 3.1.1 the definition in Section 1.2 follows. Since \( W(t) - W(s) \sim \mathcal{N}(0, t-s) \) it holds\( \text{cov}(W(s), W(t)) = \mathbb{E}[W(s)(W(t) - W(s) + W(s))] = \mathbb{E}(W(t) - W(s)) + \text{Var}(W(s)) = s, \)

thus it holds \( \text{cov}(W(s), W(t)) = \min\{s, t\} \). From \( W(t) - W(s) \sim \mathcal{N}(0, t-s) \) and \( W(0) = 0 \) it also follows that \( \mathbb{E}W(t) = 0, t \geq 0. \) The fact that \( W \) is a Gaussian process, follows from point 1) of the proof, relation \( Y = A^{-1}Z \).

\[\square\]

**Definition 3.1.2**

The process \( \{W(t), t \geq 0\}, W(t) = (W_1(t), \ldots, W_d(t)) \)\( ^\top \), \( t \geq 0 \), is called \textit{d-dimensional Brownian motion} if \( W_i = \{W_i(t), t \geq 0\} \) are independent Wiener processes, \( i = 1, \ldots, d. \)

The definitions above and Exercise 3.1.2 ensure the existence of a Wiener process with continuous paths. How do we find an explicit way of building these paths? We will show that in the next section.

### 3.2 Explicit construction of the Wiener process

First we construct the Wiener process on the interval \([0, 1]\). The main idea of the construction is to introduce a stochastic process \( X = \{X(t), t \in [0, 1]\} \) which is defined on a probability subspace of \((\Omega, \mathcal{A}, \mathbb{P})\) with \( X \overset{d}{=} W \), where \( X(t) = \sum_{n=1}^{\infty} c_n(t)Y_n, t \in [0, 1], \{Y_n\}_{n \in \mathbb{N}} \) is a sequence of i.i.d. \( \mathcal{N}(0, 1) \)-random variables and \( c_n(t) = \int_0^t H_n(s)ds, t \in [0, 1], n \in \mathbb{N} \). Here, \( \{H_n\}_{n \in \mathbb{N}} \) is the orthonormed Haar basis in \( L_2([0, 1]) \) which is introduced shortly now.

#### 3.2.1 Haar- and Schauder-functions

**Definition 3.2.1**

The functions \( H_n : [0, 1] \rightarrow \mathbb{R}, n \in \mathbb{N} \), are called \textit{Haar functions} if \( H_1(t) = 1, t \in [0, 1], H_2(t) = 1_{[0, \frac{1}{2}]}(t) - 1_{[\frac{1}{2}, 1]}(t), H_k(t) = 2^k(1_{[1, 2]}(t) - 1_{[2n, 2n+1]}(t)), t \in [0, 1], 2^n < k \leq 2^{n+1}, \) where \( I_{n,k} = [a_{n,k}, a_{n,k} + 2^{-n-1}], J_{n,k} = (a_{n,k} + 2^{-n-1}, a_{n,k} + 2^{-n}], a_{n,k} = 2^{-n}(k - 2^n - 1), n \in \mathbb{N}. \)
Lemma 3.2.1
The function system \( \{H_n\}_{n \in \mathbb{N}} \) is an orthonormal basis in \( L^2([0,1]) \) with scalar product 
\[
\langle f, g \rangle = \int_0^1 f(t)g(t)dt, \quad f, g \in L^2([0,1]).
\]

Proof The orthonormality of the system \( \langle H_k, H_n \rangle = \delta_{kn}, \quad k, n \in \mathbb{N} \) directly follows from definition 3.2.1. Now we prove the completeness of \( \{H_n\}_{n \in \mathbb{N}} \). It is sufficient to show that for arbitrary function \( g \in L^2([0,1]) \) with \( \langle g, H_n \rangle = 0, \quad n \in \mathbb{N} \), it holds \( g = 0 \) almost everywhere on \([0,1]\). In fact, we always can write the indicator function of an interval \( 1_{[a_{n,k}, a_{n,k} + 2^{-n}]} \) as a linear combination of \( H_n, \quad n \in \mathbb{N} \).

\[
1_{[0, \frac{1}{2}]} = \frac{H_1 + H_2}{2},
\]
\[
1_{(\frac{1}{2}, 1]} = \frac{H_1 - H_2}{2},
\]
\[
1_{[0, \frac{3}{4}]} = \frac{(1_{[0, \frac{1}{2}]} + \frac{1}{\sqrt{2}}H_2)}{2},
\]
\[
1_{(\frac{1}{4}, \frac{1}{2}]} = \frac{(1_{[0, \frac{1}{2}]} - \frac{1}{\sqrt{2}}H_2)}{2},
\]
\[
\vdots
\]
\[
1_{[a_{n,k}, a_{n,k} + 2^{-n}]} = \frac{(1_{[a_{n,k}, a_{n,k} + 2^{-n}]} + 2^{-n}H_k)}{2}, \quad 2^n < k \leq 2^{n+1}.
\]

Therefore it holds \( \int_0^1 g(t)dt = 0, \quad n \in \mathbb{N}_0, \quad k = 1, \ldots, 2^n - 1 \), and thus \( G(t) = \int_0^t g(s)ds = 0 \) for \( t = \frac{k}{2^n}, \quad n \in \mathbb{N}_0, \quad k = 1, \ldots, 2^n - 1 \). Since \( G \) is continuous on \([0,1]\), it follows \( G(t) = 0, \quad t \in [0,1] \), and thus \( g(s) = 0 \) for almost every \( s \in [0,1] \).

From lemma 3.2.1 it follows that two arbitrary functions \( f, g \in L^2([0,1]) \) have expansions 
\[
f = \sum_{n=1}^{\infty} \langle f, H_n \rangle H_n \quad \text{and} \quad g = \sum_{n=1}^{\infty} \langle g, H_n \rangle H_n \quad \text{(these series converge in } L^2([0,1]) \text{)}.
\]

\[
\langle f, g \rangle = \sum_{n=1}^{\infty} \langle f, H_n \rangle \langle g, H_n \rangle \quad \text{(Parseval identity)}.
\]

Definition 3.2.2
The functions \( S_n(t) = \int_0^t H_n(s)ds = \langle 1_{[0,t]}, H_n \rangle, \quad t \in [0,1], \quad n \in \mathbb{N} \) are called Schauder functions.
Lemma 3.2.2
It holds:

1. $S_n(t) \geq 0$, $t \in [0, 1]$, $n \in \mathbb{N} \setminus \{1\}$,
2. $\sum_{k=1}^{2^n} S_{2^n+k}(t) \leq \frac{1}{2} 2^{-\frac{n}{2}}$, $t \in [0, 1]$, $n \in \mathbb{N}$,
3. Let $\{a_n\}_{n \in \mathbb{N}}$ be a sequence of real numbers with $a_n = O(n^\varepsilon)$, $\varepsilon < \frac{1}{2}$, $n \to \infty$. Then the series $\sum_{n=1}^{\infty} a_n S_n(t)$ converges absolutely and uniformly in $t \in [0, 1]$ and therefore is a continuous function on $[0, 1]$.

Proof

1. follows directly from definition 3.2.2.
2. follows since functions $S_{2^n+k}$ for $k = 1, \ldots, 2^n$ have disjoint supports and $S_{2^n+k}(t) \leq S_{2^n+k}(\frac{2k-1}{2^n}) = 2^{-\frac{n}{2}-1}$, $t \in [0, 1]$.
3. It suffices to show that $R_n = \sup_{t \in [0, 1]} \sum_{k>2^n} |a_k| S_k(t) \xrightarrow{n \to \infty} 0$. For every $k \in \mathbb{N}$ and $c > 0$ it holds $|a_k| \leq ck^\varepsilon$. Therefore it holds for all $t \in [0, 1]$, $n \in \mathbb{N}$

$$\sum_{2^n < k \leq 2^{n+1}} |a_k| S_k(t) \leq c \cdot 2^{(n+1)\varepsilon} \sum_{2^n < k \leq 2^{n+1}} S_k(t) \leq c \cdot 2^{(n+1)\varepsilon} \cdot 2^{-\frac{n}{2}-1} \leq c \cdot 2^{\varepsilon - n(\frac{1}{2} - \varepsilon)}.$$

Since $\varepsilon < \frac{1}{2}$, it holds $R_m \leq c \cdot 2^{\varepsilon} \sum_{n \geq m} 2^{-n(\frac{1}{2} - \varepsilon)} \xrightarrow{m \to \infty} 0.$

Lemma 3.2.3
Let $\{Y_n\}_{n \in \mathbb{N}}$ be a sequence of (not necessarily independent) random variables defined on $(\Omega, \mathcal{A}, \mathbb{P})$, $Y_n \sim \mathcal{N}(0, 1)$, $n \in \mathbb{N}$. Then it holds $|Y_n| = O((\log n)^{\frac{1}{2}})$, $n \to \infty$, a.s.

Proof
We have to show that for $c > \sqrt{2}$ and almost all $\omega \in \Omega$ it exists a $n_0 = n_0(\omega, c) \in \mathbb{N}$ such that $|Y_n| \leq c(\log n)^{\frac{1}{2}}$ for $n \geq n_0$. If $Y \sim \mathcal{N}(0, 1)$, $x > 0$, it holds

$$P(Y > x) = \frac{1}{\sqrt{2\pi}} \int_x^{\infty} e^{-\frac{y^2}{2}} dy = \frac{1}{\sqrt{2\pi}} \int_x^{\infty} \left(-\frac{1}{y}\right) d\left(e^{-\frac{y^2}{2}}\right) = \frac{1}{\sqrt{2\pi}} \left(\frac{1}{x} e^{-\frac{x^2}{2}} - \int_x^{\infty} e^{-\frac{y^2}{2}} \frac{1}{y^2} dy\right) \leq \frac{1}{\sqrt{2\pi}} \frac{1}{x} e^{-\frac{x^2}{2}}.$$
(We also can show that $\Phi(x) \sim \frac{1}{\sqrt{2\pi}} \frac{1}{x} e^{-\frac{x^2}{2}}$, $x \to \infty$.) Thus for $c > \sqrt{2}$ it holds

$$\sum_{n \geq 2} P(|Y_n| > c (\log n)^{\frac{1}{2}}) \leq c^{-1} \frac{2}{\sqrt{2\pi}} \sum_{n \geq 2} (\log n)^{-\frac{1}{2}} \log n = \frac{c^{-1}\sqrt{2}}{\sqrt{\pi}} \sum_{n \geq 2} (\log n)^{-\frac{1}{2}} n^{-\frac{1}{2}} < \infty.$$ 

According to the Lemma of Borel-Cantelli (cf. WR, Lemma 2.2.1) it holds $P(\cap_{k \geq 2} A_k) = 0$ if $\sum_k P(A_k) < \infty$ with $A_k = \{|Y_k| > e \cdot (\log k)^{\frac{1}{2}}\}, k \in \mathbb{N}$. Thus $A_k$ occurs in infinite number only with probability 0, with $|Y_n| \leq c (\log n)^{\frac{1}{2}}$ for $n \geq n_0$.

3.2.2 Wiener process with a.s. continuous paths

Lemma 3.2.4

Let $\{Y_n\}_{n \in \mathbb{N}}$ be a sequence of independent $\mathcal{N}(0,1)$-distributed random variables. Let $\{a_n\}_{n \in \mathbb{N}}$ and $\{b_n\}_{n \in \mathbb{N}}$ be sequences of numbers with $\sum_{k=1}^{m} |a_{2m+k}| \leq 2^{-\frac{m}{2}}$, $\sum_{k=1}^{m} |b_{2m+k}| \leq 2^{-\frac{m}{2}}$, $m \in \mathbb{N}$. Then the limits $U = \sum_{n=1}^{\infty} a_n Y_n$ and $V = \sum_{n=1}^{\infty} b_n Y_n$, $U \sim \mathcal{N}(0, \sum_{n=1}^{\infty} a_n^2)$, $V \sim \mathcal{N}(0, \sum_{n=1}^{\infty} b_n^2)$ exist a.s., where $\text{cov}(U, V) = \sum_{n=1}^{\infty} a_n b_n$. $U$ and $V$ are independent if and only if $\text{cov}(U, V) = 0.$

Proof Lemma 3.2.2 and 3.2.3 reveal the a.s. existence of the limits $U$ and $V$ (replace $a_n$ by $Y_n$ and $S_n$ by e.g. $b_n$ in Lemma 3.2.2). From the stability under convolution of the normal distribution it follows for $U^{(m)} = \sum_{n=1}^{m} a_n Y_n$, $V^{(m)} = \sum_{n=1}^{m} b_n Y_n$, that $U^{(m)} \sim \mathcal{N}(0, \sum_{n=1}^{m} a_n^2)$, $V^{(m)} \sim \mathcal{N}(0, \sum_{n=1}^{m} b_n^2)$. Since $U^{(m)} \xrightarrow{d} U$, $V^{(m)} \xrightarrow{d} V$ it follows $U \sim \mathcal{N}(0, \sum_{n=1}^{\infty} a_n^2)$, $V \sim \mathcal{N}(0, \sum_{n=1}^{\infty} b_n^2)$. Moreover, it holds

$$\text{cov}(U, V) = \lim_{m \to \infty} \text{cov}(U^{(m)}, V^{(m)})$$

$$= \lim_{m \to \infty} \sum_{i,j=1}^{m} a_i b_j \text{cov}(Y_i, Y_j)$$

$$= \lim_{m \to \infty} \sum_{i=1}^{m} a_i b_i = \sum_{i=1}^{\infty} a_i b_i,$$

according to the dominated convergence theorem of Lebesgue, since according to Lemma 3.2.3 it holds $|Y_n| \leq c (\log n)^{\frac{1}{2}}$, for $n \geq \mathbb{N}_0$, and the dominated series converges according to Lemma 3.2.2:

$$\sum_{n,k=2^m}^{2^{m+1}} a_n b_k Y_n Y_k \xrightarrow{a.s.} \sum_{n,k=2^m}^{2^{m+1}} a_n b_k c^2 n^\frac{1}{2} k^\frac{1}{2} \leq 2^2 \varepsilon (m+1) \cdot 2^{-\frac{m}{2}} \cdot 2^{-\frac{m}{2}} = 2^{-(1-2\varepsilon)m}, \quad 1 - 2\varepsilon > 0.$$ 

For sufficient large $m$ it holds $\sum_{n,k=m}^{\infty} a_n b_k Y_n Y_k \leq \sum_{j=m}^{\infty} 2^{-(1-2\varepsilon)j} < \infty$, and this series converges a.s.

Now we show

$$\text{cov}(U, V) = 0 \iff \text{U and V are independent}$$

Independence always results in the uncorrelation of random variables. We prove the other
direction. From \((U^{(m)}, V^{(m)}) \xrightarrow{\frac{d}{m \to \infty}} (U, V)\) it follows \(\varphi_{(U^{(m)}, V^{(m)})} \xrightarrow{m \to \infty} \varphi(U, V)\), thus

\[
\begin{align*}
\varphi_{(U^{(m)}, V^{(m))}}(s, t) & = \lim_{m \to \infty} \mathbb{E} \exp\{i \left( t \sum_{k=1}^{m} a_k Y_k + s \sum_{n=1}^{m} b_n Y_n \right) \} \\
& = \lim_{m \to \infty} \mathbb{E} \exp\{i \sum_{k=1}^{m} (t a_k + s b_k) Y_k \} = \lim_{m \to \infty} \prod_{k=1}^{m} \mathbb{E} \exp\{i (t a_k + s b_k) Y_k \} \\
& = \lim_{m \to \infty} \prod_{k=1}^{m} \exp\{-\frac{(t a_k + s b_k)^2}{2}\} = \exp\{-\sum_{k=1}^{\infty} \frac{(t a_k + s b_k)^2}{2}\} \\
& = \exp\left\{ -\frac{t^2}{2} \sum_{k=1}^{\infty} a_k^2 \right\} \exp\left\{ t s \sum_{k=1}^{\infty} a_k b_k \right\} \exp\left\{ -\frac{s^2}{2} \sum_{k=1}^{\infty} b_k^2 \right\} = \varphi_U(t) \varphi_V(s),
\end{align*}
\]

\(s, t \in \mathbb{R}\). Thus, \(U\) and \(V\) are independent if \(\text{cov}(U, V) = 0\).

\[\square\]

**Theorem 3.2.1**

Let \(\{Y_n, n \in \mathbb{N}\}\) be a sequence of i.i.d. random variables that are \(\mathcal{N}(0, 1)\)-distributed, defined on a probability space \((\Omega, \mathcal{A}, \mathbb{P})\). Then there exists a probability space \((\Omega_0, \mathcal{A}_0, \mathbb{P})\) of \((\Omega, \mathcal{A}, \mathbb{P})\) and a stochastic process \(X = \{X(t), t \in [0, 1]\}\) on it such that \(X(t, \omega) = \sum_{n=1}^{\infty} Y_n(\omega) S_n(t), t \in [0, 1], \omega \in \Omega_0\) and \(X \overset{d}{=} W\). Here, \(\{S_n\}_{n \in \mathbb{N}}\) is the family of Schauder functions.

**Proof** According to Lemma 3.2.2, 2) the coefficients \(S_n(t)\) fulfill the conditions of Lemma 3.2.4 for every \(t \in [0, 1]\). In addition to that it exists according to Lemma 3.2.3 a subset \(\Omega_0 \subset \Omega, \Omega_0 \in \mathcal{A}\) with \(\mathbb{P}(\Omega_0) = 1\), such that for every \(\omega \in \Omega_0\) the relation \(|Y_n(\omega)| = O(\sqrt{\log n}), n \to \infty\), holds. Let \(\mathcal{A}_0 = \mathcal{A} \cap \Omega_0\). We restrict the probability space to \((\Omega_0, \mathcal{A}_0, \mathbb{P})\). Then condition \(a_n = Y_n(\omega) = O(n^\varepsilon), \varepsilon < \frac{1}{2}\), is fulfilled since \(\sqrt{\log n} < n^\varepsilon\) for sufficient large \(n\), and according to Lemma 3.2.2, 3) the series \(\sum_{n=1}^{\infty} Y_n(\omega) S_n(t)\) converges absolutely and uniformly in \(t \in [0, 1]\) to the function \(X(\cdot, \omega), \omega \in \Omega_0\), which is a continuous function in \(t\) for every \(\omega \in \Omega_0\). \(X(\cdot, t)\) is a random variable since in Lemma 3.2.4 the convergence of this series holds almost surely. Moreover it holds \(X(t) \sim \mathcal{N}(0, \sum_{n=1}^{\infty} S_n^2(t)), t \in [0, 1]\).

We show that this stochastic process, defined on \((\Omega_0, \mathcal{A}_0, \mathbb{P})\), is a Wiener process. For that we check the conditions of Definition 3.1.1. We consider arbitrary times \(0 \leq t_1 < t_2, t_3 < t_4 \leq 1\)
and evaluate
\[
\text{cov}(X(t_2) - X(t_1), X(t_4) - X(t_3)) = \text{cov}\left(\sum_{n=1}^{\infty} Y_n(S_n(t_2) - S_n(t_1)), \sum_{n=1}^{\infty} Y_n(S_n(t_4) - S_n(t_3))\right)
\]
\[
= \sum_{n=1}^{\infty} (S_n(t_2) - S_n(t_1))(S_n(t_4) - S_n(t_3))
\]
\[
= \sum_{n=1}^{\infty} \left(\langle H_n, 1_{[0,t_2]} \rangle - \langle H_n, 1_{[0,t_1]} \rangle\right) \times \left(\langle H_n, 1_{[0,t_4]} \rangle - \langle H_n, 1_{[0,t_2]} \rangle\right)
\]
\[
= \sum_{n=1}^{\infty} \langle H_n, 1_{[0,t_2]} - 1_{[0,t_1]} \rangle \langle H_n, 1_{[0,t_4]} - 1_{[0,t_3]} \rangle
\]
\[
= \min\{t_2, t_4\} - \min\{t_1, t_4\} - \min\{t_2, t_3\} + \min\{t_1, t_3\},
\]
by Parseval inequality and since \(1_{[0,s]} > \int_0^s \min\{s,t\} \, du = \min\{s,t\}, s,t \in [0,1]\). If \(0 \leq t_1 < t_2 \leq t_3 < t_4 < 1\), it holds \(\text{cov}(X(t_2) - X(t_1), X(t_4) - X(t_3)) = t_2 - t_1 - t_2 + t_1 = 0\), thus the increments of \(X\) (according to Lemma 3.2.4) are uncorrelated. Moreover it holds \(X(0) \sim \mathcal{N}(0, \sum_{n=1}^{\infty} S_n^2(0)) = \mathcal{N}(0,0)\), therefore \(X(0) \overset{a.s.}{=} 0\). For \(t_1 = 0\), \(t_2 = t\), \(t_3 = 0\), \(t_4 = t\) it follows that \(\text{Var}(X(t)) = t\), \(t \in [0,1]\), and for \(t_1 = t_3 = s\), \(t_2 = t_4 = t\), that \(\text{Var}(X(t) - X(s)) = t - s - s + s = t - s\), \(0 \leq s < t \leq 1\). Thus it holds \(X(t) - X(s) \sim \mathcal{N}(0,t-s)\), and according to Definition 3.1.1 it holds \(X \overset{d}{=} W\).

**Remark 3.2.1**

1. Theorem 3.2.1 is the basis for an approximative simulation of the paths of a Brownian motion through the partial sums \(X^{(n)}(t) = \sum_{k=1}^{n} Y_k S_k(t), t \in [0,1]\), for sufficient large \(n \in \mathbb{N}\).

2. The construction in Theorem 3.2.1 can be used to construct the Wiener process with continuous paths on the interval \([0,t_0]\) for arbitrary \(t_0 > 0\). If \(W = \{W(t), t \in [0,1]\}\) is a Wiener process on \([0,1]\) then \(Y = \{Y(t), t \in [0,t_0]\}\) with \(Y(t) = \sqrt{t_0} W(\frac{t}{t_0}), t \in [0,t_0]\), is a Wiener process on \([0,t_0]\).

**Exercise 3.2.1**

Prove that.

3. The Wiener process \(W\) with continuous paths on \(\mathbb{R}_+\) can be constructed as follows. Let \(W^{(n)} = \{W^{(n)}(t), t \in [0,1]\}\) be independent copies of the Wiener process as in Theorem 3.2.1. Define \(W(t) = \sum_{n=1}^{\infty} 1(t \in [n-1,n])1(\sum_{k=1}^{n-1} W^{(k)}(1) + W^{(n)}(t-(n-1))), t \geq 0\), thus,
\[
W(t) = \begin{cases} 
W^{(1)}(t), & t \in [0,1], \\
W^{(1)}(1) + W^{(2)}(t-1), & t \in [1,2], \\
W^{(1)}(1) + W^{(2)}(1) + W^{(3)}(t-2), & t \in [2,3],
\end{cases}
\]

**Exercise 3.2.2**

Show that the introduced stochastic process \(W = \{W(t), t \geq 0\}\) is a Wiener process on \(\mathbb{R}_+\).
3 Wiener process

3.3 Distribution and path properties of Wiener processes

3.3.1 Distribution of the maximum

Theorem 3.3.1
Let \( W = \{ W(t), \ t \in [0, 1] \} \) be the Wiener process defined on a probability space \( (\Omega, \mathcal{F}, P) \). Then it holds:

\[
P \left( \max_{t \in [0, 1]} W(t) > x \right) = \frac{2}{\pi} \int_x^\infty e^{-y^2/2} dy \tag{3.3.1}
\]

for all \( x \geq 0 \).

The mapping \( \max_{t \in [0, 1]} W(t) : \Omega \to [0, \infty) \) given in relation (3.3.1) is a well-defined random variable since it holds: \( \max_{t \in [0, 1]} W(t, \omega) = \lim_{n \to \infty} \max_{i=1, \ldots, k} W(t, \omega) \) for all \( \omega \in \Omega \) since the trajectories of \( \{ W(t), \ t \in [0, 1] \} \) are continuous. From 3.3.1 it follows that \( \max_{t \in [0, 1]} W(t) \) has an exponential bounded tail: thus \( \max_{t \in [0, 1]} W(t) \) has finite \( k \)-th moments.

Useful ideas for the proof of Theorem 3.3.1
Let \( \{ W(t), \ t \in [0, 1] \} \) be a Wiener process and \( Z_1, Z_2, \ldots \) a sequence of independent random variables with \( P(Z_i = 1) = P(Z_i = -1) = \frac{1}{2} \) for all \( i \geq 1 \). For every \( n \in \mathbb{N} \) we define \( \{ \tilde{W}^n(t), \ t \in [0, 1] \} \) by:

\[
\tilde{W}^n(t) = \frac{S_{\lfloor nt \rfloor}}{\sqrt{n}} + (nt - \lfloor nt \rfloor) \frac{Z_{\lfloor nt \rfloor + 1}}{\sqrt{n}}, \text{ where } S_i = Z_1 + \ldots + Z_i, \ i \geq 1, \ S_0 = 0.
\]

Lemma 3.3.1
For every \( k \geq 1 \) and arbitrary \( t_1, \ldots, t_k \in [0, 1] \) it holds:

\[
\left( \tilde{W}^{(n)}(t_1), \ldots, \tilde{W}^{(n)}(t_k) \right) \overset{d}{\to} (W(t_1), \ldots, W(t_k))^\top.
\]

Proof Consider the special case \( k = 2 \) (for \( k > 2 \) the proof is analogous). Let \( t_1 < t_2 \). For all
\[ s_1 \hat{W}^{(n)}(t_1) + s_2 \hat{W}^{(n)}(t_2) = \left( s_1 + s_2 \right) \frac{S_{nt_1}}{\sqrt{n}} + s_2 \left( \frac{S_{nt_2} - S_{nt_1} + 1}{\sqrt{n}} \right) \]

since \( S_{nt_2} = S_{nt_1} + S_{nt_2} - S_{nt_1} + 1 + S_{nt_1} + 1 \).

Now observe that the 4 summands on the right-hand-side of the previous equation are independent and moreover that the latter two summands converge (a.s. and therefor particularly in distribution) to zero.

Consequently, it holds

\[
\lim_{n \to \infty} \mathbb{E} e^{(s_1 \hat{W}^{(n)}(t_1) + s_2 \hat{W}^{(n)}(t_2))} = \lim_{n \to \infty} \mathbb{E} e^{s_1 \hat{W}^{(n)}(t_1)} \mathbb{E} e^{s_2 \hat{W}^{(n)}(t_2) - S_{nt_1} + 1}
\]

\[
= \lim_{n \to \infty} \mathbb{E} e^{s_1 \hat{W}^{(n)}(t_1)} \mathbb{E} e^{s_2 \hat{W}^{(n)}(t_2) - S_{nt_1} + 1}
\]

\[
\text{CLT, CMT} = e^{\frac{s_1}{\sqrt{n}} (s_1 + s_2)^2} e^{-\frac{s_1^2}{2} s_2^2} = e^{\frac{s_1}{2} (s_1^2 + 2s_1s_2 + s_2^2)} = e^{\frac{s_1}{2} (s_1^2 + 2s_1s_2 \min(t_1, t_2) + s_2^2)} = \varphi(W(t_1), W(t_2))(s_1, s_2),
\]

where \( \varphi(W(t_1), W(t_2)) \) is the characteristic function of \( (W(t_1), W(t_2)) \).

**Lemma 3.3.2**

Let \( \hat{W}^{(n)} = \max_{t \in [0,1]} W^{(n)}(t) \). Then it holds:

\[
\hat{W}^{(n)} \overset{d}{=} \frac{1}{\sqrt{n}} \max_{k=1,\ldots,n} S_k, \quad \text{for all } n \in \mathbb{N}
\]

and

\[
\lim_{n \to \infty} \mathbb{P}(\hat{W}^{(n)} \leq x) = \sqrt{\frac{2}{\pi}} \int_{0}^{x} e^{-\frac{y^2}{2}} dy, \quad \text{for all } x \geq 0.
\]

Without proof

**Proof** of Theorem 3.3.1. We shall prove only the upper bound in Theorem 3.3.1.

From Lemma 3.3.1 and the continuous mapping theorem it follows for \( x \geq 0, k \geq 1 \) and \( t_1, \ldots, t_k \in [0,1] \) that

\[
\lim_{n \to \infty} \mathbb{P} \left( \max_{t \in \{t_1, \ldots, t_k\}} W^{(n)}(t) > x \right) = \mathbb{P} \left( \max_{t \in \{t_1, \ldots, t_k\}} W(t) > x \right),
\]

since \( (x_1, \ldots, x_k) \mapsto \max(x_1, \ldots, x_k) \) is continuous.

Consequently, it holds

\[
\liminf_{n \to \infty} \mathbb{P} \left( \max_{t \in [0,1]} \hat{W}^{(n)}(t) > x \right) \geq \mathbb{P} \left( \max_{t \in \{t_1, \ldots, t_k\}} W(t) > x \right),
\]
since \( \{ \max_{t \in \{t_1, \ldots, t_k\}} \hat{W}^{(n)}(t) > x \} \subseteq \{ \max_{t \in [0,1]} \hat{W}^{(n)}(t) > x \} \).

With \((t_1, \ldots, t_k)^\top = \left( \frac{1}{k}, \ldots, \frac{k}{k} \right)^\top\) and \(\max_{t \in [0,1]} W(t) = \lim_{k \to \infty} \max_{i=1, \ldots, k} W\left( \frac{i}{k} \right)\) a.s. (and therefore particularly in distribution) it holds
\[
\lim_{n \to \infty} \inf \mathbb{P}\left( \max_{t \in [0,1]} \hat{W}^{(n)}(t) > x \right) \geq \lim_{k \to \infty} \mathbb{P}\left( \max_{i=1, \ldots, k} W\left( \frac{i}{k} \right) > x \right) = \mathbb{P}\left( \max_{t \in [0,1]} W(t) > x \right).
\]

Conclusively, the assertion follows from lemma 3.3.2. □

### 3.3.2 Invariance properties

Specific transformations of the Wiener process again reveal the Wiener process.

**Theorem 3.3.2**

Let \( \{W(t), \ t \geq 0\} \) be a Wiener process. Then the stochastic processes \( \{Y^{(i)}(t), \ t \geq 0\}, i = 1, \ldots, 4, \) with
\[
\begin{align*}
Y^{(1)}(t) &= -W(t), \quad \text{(Symmetry)} \\
Y^{(2)}(t) &= W(t + t_0) - W(t_0) \quad \text{for a } t_0 > 0, \quad \text{Translation of the origin} \\
Y^{(3)}(t) &= \sqrt{c}W\left( \frac{t}{c} \right) \quad \text{for a } c > 0, \quad \text{(Scaling)} \\
Y^{(4)}(t) &= \left\{ \begin{array}{ll}
tW\left( \frac{1}{t} \right), & t > 0, \\
0, & t = 0. \end{array} \right. \quad \text{(Reflection at } t=0) 
\end{align*}
\]
are Wiener processes as well.

**Proof**

1. \( Y^{(i)}, i = 1, \ldots, 4, \) have independent increments with \( Y^{(i)}(t_2) - Y^{(i)}(t_1) \sim \mathcal{N}(0, t_2 - t_1) \).

2. \( Y^{(i)}(0) = 0, i = 1, \ldots, 4. \)

3. \( Y^{(i)}, i = 1, \ldots, 3, \) have continuous trajectories. \( \{Y^{(i)}(t), \ t \geq 0\} \) has continuous trajectories for \( t > 0 \).

4. We have to prove that \( Y^{(4)}(t) \) is continuous at \( t = 0 \), i.e. that \( \lim_{t \to 0} tW\left( \frac{1}{t} \right) = 0. \)
\[
\lim_{t \to 0} tW\left( \frac{1}{t} \right) = \lim_{t \to \infty} \frac{W(t)}{t} \text{ a.s.} = 0 \text{ because of Corollary ??}. \]

\[\Box\]

**Corollary 3.3.1**

Let \( \{W(t), \ t \geq 0\} \) be the Wiener process. Then it holds:
\[
\mathbb{P}\left( \sup_{t \geq 0} W(t) = \infty \right) = \mathbb{P}\left( \inf_{t \geq 0} W(t) = -\infty \right) = 1.
\]

and consequently
\[
\mathbb{P}\left( \sup_{t \geq 0} W(t) = \infty, \inf_{t \geq 0} W(t) = -\infty \right).
\]
\textbf{Proof} For \(x, c > 0\) it holds:
\[
P \left( \sup_{t \geq 0} W(t) > x \right) = P \left( \sup_{t \geq 0} W \left( \frac{t}{c} \right) > \frac{x}{\sqrt{c}} \right) = P \left( \sup_{t \geq 0} W(t) > \frac{x}{\sqrt{c}} \right)
\]
\[
\Rightarrow P \left( \sup_{t \geq 0} W(t) = 0 \right) \cup \left( \sup_{t \geq 0} W(t) = \infty \right) = P \left( \sup_{t \geq 0} W(t) = 0 \right) + P \left( \sup_{t \geq 0} W(t) = \infty \right) = 1.
\]
Moreover it holds
\[
P \left( \sup_{t \geq 0} W(t) = 0 \right) = P \left( \sup_{t \geq 0} W(t) \leq 0 \right) \leq P \left( W(t) \leq 0, \sup_{t \geq 1} W(t) \leq 0 \right)
\]
\[
= P \left( W(1) \leq 0, \sup_{t \geq 1} (W(t) - W(1)) \leq -W(1) \right)
\]
\[
= \int_{-\infty}^{0} P \left( \sup_{t \geq 1} W(t) - W(1) \leq -W(1) \mid W(1) = x \right) P(W(1) \in dx)
\]
\[
= \int_{-\infty}^{0} P \left( \sup_{t \geq 0} (W(t) - W(1)) \leq -x \mid W(1) = x \right) P(W(1) \in dx)
\]
\[
= \int_{-\infty}^{0} P \left( \sup_{t \geq 0} W(t) = 0 \right) P(W(1) \in dx)
\]
\[
= P \left( \sup_{t \geq 0} W(t) = 0 \right) \frac{1}{2},
\]
thus \(P \left( \sup_{t \geq 0} W(t) = 0 \right) = 0\) and thus \(P \left( \sup_{t \geq 0} W(t) = \infty \right) = 1\).

Analogously one can show that \(P \left( \inf_{t \geq 0} W(t) = -\infty \right) = 1\).

The remaining part of the claim follows from \(P(A \cap B) = 1\) for any \(A, B \in \mathcal{F}\) with \(P(A) = P(B) = 1\).

\textbf{Remark 3.3.1}
\(P \left( \sup_{t \geq 0} X(t) = \infty, \inf_{t \geq 0} X(t) = -\infty \right) = 1\) implies that the trajectories of \(W\) oscillate between positive and negative values on \([0, \infty)\) an infinite number of times.

\textbf{Corollary 3.3.2}
Let \(\{W(t), t \geq 0\}\) be a Wiener process. Then it holds
\[
P \left( \omega \in \Omega : W(\omega) \text{ is nowhere differentiable in } [0, \infty) \right) = 1.
\]

\textbf{Proof}
\[
\{\omega \in \Omega : W(\omega) \text{ is nowhere differentiable in } [0, \infty)\}
\]
\[
= \bigcap_{n=0}^{\infty} \{\omega \in \Omega : W(\omega) \text{ is nowhere differentiable in } [n, n+1]\}.
\]
It is sufficient to show that \(P(\omega \in \Omega : W(\omega) \text{ is differentiable for } t_0 = t_0(\omega) \in [0, 1]) = 0\).
Define the set
\[
A_{nm} = \left\{ \omega \in \Omega : \text{it exists a } t_0 = t_0(\omega) \in [0, 1] \text{ with } |W(t_0(\omega) + h, \omega) - W(t_0(\omega), \omega)| \leq mh, \forall h \in \left[0, \frac{4}{n}\right] \right\}.
\]
Then it holds

\[ \{ \omega \in \Omega : W(\omega) \text{ differentiable for } t_0 = t_0(\omega) \} = \bigcup_{m \geq 1} \bigcup_{n \geq 1} A_{nm}. \]

We still have to show \( \mathbb{P}(\bigcup_{m \geq 1} \cup_{n \geq 1} A_{nm}) = 0. \)

Let \( k_0(\omega) = \min_{k=1,2,...} \{ \frac{k}{n} \geq t_0(\omega) \}. \) Then it holds for \( \omega \in A_{nm} \) and \( j = 0, 1, 2 \)

\[
\left| W\left( \frac{k_0(\omega) + j + 1}{n}, \omega \right) - W\left( \frac{k_0(\omega) + j}{n}, \omega \right) \right| \leq \left| W\left( \frac{k_0(\omega) + j + 1}{n}, \omega \right) - W(t_0(\omega), \omega) \right| \\
+ \left| W\left( \frac{k_0(\omega) + j}{n}, \omega \right) - W(t_1(\omega), \omega) \right| \\
\leq \frac{8m}{n}.
\]

Let \( \Delta_n(k) = W\left( \frac{k+1}{n} \right) - W\left( \frac{k}{n} \right). \) Then it holds

\[
\mathbb{P}(A_{nm}) \leq \mathbb{P}\left( \bigcap_{k=0}^{n} \bigcap_{j=0}^{2} |\Delta_n(k+j)| \leq \frac{8m}{n} \right) \\
\leq \sum_{k=0}^{n} \mathbb{P}\left( \bigcap_{j=0}^{2} |\Delta_n(k+j)| \leq \frac{8m}{n} \right) \\
= \left( \sum_{k=0}^{n} \mathbb{P}( |\Delta_n(0)| \leq \frac{8m}{n} ) \right)^3 \\
\leq (n+1) \left( \frac{16m}{\sqrt{2\pi n}} \right)^3 \to 0, \quad n \to \infty,
\]

by the independence of the increments of the Wiener Process.

Since \( A_{nm} \subset A_{n+1,m}, \) it follows \( \mathbb{P}(A_{nm}) = 0. \)

**Corollary 3.3.3**

With probability 1 it holds:

\[
\sup_{n \geq 1} \sup_{0 \leq t_0 < t_1 < t_2 < \cdots < t_n \leq 1} \sum_{i=1}^{n} |W(t_i) - W(t_{i-1})| = \infty,
\]

i.e. \( \{ W(t), \; t \in [0,1] \} \) possesses a.s. trajectories with unbounded variation.

**Proof** Since every continuous function \( g : [0,1] \to \mathbb{R} \) with bounded variation is differentiable almost everywhere, the assertion follows from Corollary 3.3.2.

**Alternative proof**

It is sufficient to show that \( \lim_{n \to \infty} \sum_{i=1}^{2^n} |W\left( \frac{i}{2^n} \right) - W\left( \frac{(i-1)t}{2^n} \right) | = \infty. \)

Let \( Z_n = \sum_{i=1}^{2^n} \left( W\left( \frac{i}{2^n} \right) - W\left( \frac{(i-1)t}{2^n} \right) \right)^2 - t. \) Hence \( \mathbb{E}Z_n = 0 \) and \( \mathbb{E}Z_n^2 = t^22^{-n+1} \) and with Tchebycheff’s inequality

\[
\mathbb{P}(|Z_n| < \varepsilon) \leq \frac{\mathbb{E}Z_n^2}{\varepsilon^2} = \left( \frac{t}{\varepsilon} \right)^2 2^{-n+1}, \quad \text{i.e.} \quad \sum_{i=1}^{\infty} \mathbb{P}(|Z_n| > \varepsilon) \xrightarrow{a.s.} 0.
\]
From lemma of Borel-Cantelli it follows that \( \lim_{n \to \infty} Z_n = 0 \) almost surely and thus
\[
0 \leq t \leq \sum_{i=1}^{2^n} \left( W \left( \frac{it}{2^n} \right) - W \left( \frac{(i-1)t}{2^n} \right) \right)^2 \leq \liminf_{n \to \infty} \max_{1 \leq k \leq 2^n} \left| W \left( \frac{kt}{2^n} \right) - W \left( \frac{(k-1)t}{2^n} \right) \right| \sum_{i=1}^{2^n} \left| W \left( \frac{it}{2^n} \right) - W \left( \frac{(i-1)t}{2^n} \right) \right|.
\]
Hence the assertion follows since \( W \) has continuous trajectories and therefore
\[
\lim_{n \to \infty} \max_{1 \leq k \leq 2^n} \left| W \left( \frac{kt}{2^n} \right) - W \left( \frac{(k-1)t}{2^n} \right) \right| = 0.
\]

\[\square\]

### 3.4 Additional exercises

**Exercise 3.4.1**

Give an intuitive (exact!) method to realize trajectories of a Wiener process \( W = \{W(t), t \in [0,1]\} \). Thereby use the independence and the distribution of the increments of \( W \). Additionally, write a program in \( R \) for the simulation of paths of \( W \). Draw three paths \( t \mapsto W(t,\omega) \) for \( t \in [0,1] \) in a common diagram.

**Exercise 3.4.2**

Given are the Wiener process \( W = \{W(t), t \in [0,1]\} \) and \( L := \arg\max_{t \in [0,1]} W(t) \). Show that it holds:
\[
P(L \leq x) = \frac{2}{\pi} \arcsin \sqrt{x}, \quad x \in [0,1].
\]

Hint: Use relation \( \max_{r \in [0,t]} W(r) \overset{d}{=} |W(t)| \).

**Exercise 3.4.3**

For the simulation of a Wiener process \( W = \{W(t), t \in [0,1]\} \) we also can use the approximation
\[
W_n(t) = \sum_{k=1}^{n} S_k(t) z_k
\]
where \( S_k(t), t \in [0,1], k \geq 1 \) are the Schauder functions, and \( z_k \sim \mathcal{N}(0,1) \) i.i.d. random variables and the series converges almost surely for all \( t \in [0,1] \) \((n \to \infty)\).

a) Show that for all \( t \in [0,1] \) the approximation \( W_n(t) \) also converges in the \( L^2 \)-sense to \( W(t) \).

b) Write a program in \( R \) (alternative: C) for the simulation of a Wiener process \( W = \{W(t), t \in [0,1]\} \).

c) Simulate three paths \( t \mapsto W(t,\omega) \) for \( t \in [0,1] \) and draw these paths into a common diagram. Hereby consider the sampling points \( t_k = \frac{k}{n}, k = 0,\ldots,n \) with \( n = 2^8 - 1 \).

**Exercise 3.4.4**

For the Wiener process \( W = \{W(t), t \geq 0\} \) we define the process of the maximum that is given by \( M = \{M(t) := \max_{s \in [0,t]} W(s), t \geq 0\} \). Show that it holds:
a) The density $f_M(t)$ of the maximum $M(t)$ is given by

$$f_M(t)(x) = \sqrt{\frac{2}{\pi t}} \exp\left\{ -\frac{x^2}{2t} \right\} 1\{x \geq 0\}.$$  

Hint: Use property $P(M(t) > x) = 2P(W(t) > x)$.

b) Expected value and variance of $M(t)$ are given by

$$\mathbb{E}M(t) = \sqrt{\frac{2t}{\pi}}, \quad \text{Var} M(t) = t(1 - 2/\pi).$$

Now we define $\tau(x) := \arg\min_{s \in \mathbb{R}} \{W(s) = x\}$ as the first point in time for which the Wiener process takes value $x$.

c) Determine the density of $\tau(x)$ and show that: $\mathbb{E}\tau(x) = \infty$.

**Exercise 3.4.5**

Let $W = \{W(t), t \geq 0\}$ be a Wiener process. Show that the following processes are Wiener processes as well:

$$W_1(t) = \begin{cases} 0, & t = 0, \\ tW(1/t), & t > 0, \end{cases} \quad W_2(t) = \sqrt{c}W(t/c), \quad c > 0.$$

**Exercise 3.4.6**

The Wiener process $W = \{W(t), t \geq 0\}$ is given. Size $Q(a, b)$ denotes the probability that the process exceeds the half line $y = at + b$, $t \geq 0$, $a, b > 0$. Proof that:

a) $Q(a, b) = Q(b, a)$ and $Q(a, b_1 + b_2) = Q(a, b_1)Q(a, b_2)$,

b) $Q(a, b)$ is given by $Q(a, b) = \exp\{-2ab\}$. 
4 Lévy Processes

4.1 Lévy Processes

Definition 4.1.1
A stochastic process \( \{X(t), \ t \geq 0\} \) is called Lévy process, if

1. \( X(0) = 0 \),
2. \( \{X(t)\} \) has stationary and independent increments,
3. \( \{X(t)\} \) is stochastically continuous, i.e for an arbitrary \( \varepsilon > 0, \ t_0 \geq 0 \):

\[
\lim_{t \to t_0} P(|X(t) - X(t_0)| > \varepsilon) = 0.
\]

Remark 4.1.1
One can easily see, that compound Poisson processes fulfill the 3 conditions, since for arb. \( \varepsilon > 0 \) it holds

\[
P(|X(t) - X(t_0)| < \varepsilon) \geq P(|X(t) - X(t_0)| > 0) \leq 1 - e^{-\lambda|t-t_0|} \xrightarrow{t \to t_0} 0.
\]

Further holds for the Wiener process for arb. \( \varepsilon > 0 \)

\[
P(|X(t) - X(t_0)| > \varepsilon) = \sqrt{\frac{2}{\pi(t-t_0)}} \int_{t_0}^\infty \exp\left(-\frac{y^2}{2(t-t_0)}\right) dy
\]

\[
\xrightarrow{x = \frac{y}{\sqrt{t-t_0}}} \frac{2}{\pi} \int_{t_0}^\infty e^{-\frac{x^2}{2}} dx \xrightarrow{t \to t_0} 0.
\]

4.1.1 Infinitely Divisibility

Definition 4.1.2
Let \( X : \Omega \to \mathbb{R} \) be an arbitrary random variable. Then \( X \) is called infinitely divisible, if for arbitrary \( n \in \mathbb{N} \) there exist i.i.d. random variables \( Y_1^{(n)}, \ldots, Y_n^{(n)} \) with \( X \overset{d}{=} Y_1^{(n)} + \ldots + Y_n^{(n)} \).

Lemma 4.1.1
The random variable \( X : \Omega \to \mathbb{R} \) is infinitely divisible if and only if the characteristic function \( \varphi_X \) of \( X \) can be expressed for every \( n \geq 1 \) in the form

\[
\varphi_X(s) = (\varphi_n(s))^n \quad \text{for all } s \in \mathbb{R},
\]

where \( \varphi_n \) are characteristic functions of random variables.
Proof \(\Rightarrow\) 
Y_1^{(n)}, \ldots, Y_n^{(n)} \text{ i.i.d., } X \overset{d}{=} Y_1^{(n)} + \ldots + Y_n^{(n)}. \) Hence, it follows that \(\varphi_X(s) = \prod_{i=1}^{n} \varphi_{Y_i^{(n)}}(s) = (\varphi_n(s))^n.\)

\(\Leftarrow\) 
\(\varphi_X(s) = (\varphi_n(s))^n \Rightarrow \) there exist \(Y_1^{(n)}, \ldots, Y_n^{(n)} \) i.i.d. with characteristic function \(\varphi_n\) and \(\varphi_{Y_1^{(n)}, \ldots, Y_n^{(n)}}(s) = (\varphi_n(s))^n = \varphi_X(s).\) With the uniqueness theorem for characteristic functions it follows that \(X \overset{d}{=} Y_1^{(n)} + \ldots + Y_n^{(n)}.\)

\[\text{Theorem 4.1.1}\]

Let \(\{X(t), \; t \geq 0\}\) be a Lévy process. Then the random variable \(X(t)\) is infinitely divisible for every \(t \geq 0.\)

**Proof** For arbitrary \(t \geq 0\) and \(n \in \mathbb{N}\) it obviously holds that

\[
X(t) = X\left(\frac{t}{n}\right) + \left(X\left(\frac{2t}{n}\right) - X\left(\frac{t}{n}\right)\right) + \ldots + \left(X\left(\frac{nt}{n}\right) - X\left(\frac{(n-1)t}{n}\right)\right).
\]

Since \(\{X(t)\}\) has independent and stationary increments, the summands are obviously independent and identically distributed random variables.

**Lemma 4.1.2**

Let \(X_1, X_2, \ldots : \Omega \to \mathbb{R}\) be a sequence of random variables. If there exists a function \(\varphi : \mathbb{R} \to \mathbb{C},\) such that \(\varphi(s)\) is continuous in \(s = 0\) and \(\lim_{n \to \infty} \varphi_{X_n}(s) = \varphi(s)\) for all \(s \in \mathbb{R},\) then \(\varphi\) is the characteristic function of a random variable \(X\) and it holds that \(X_n \overset{d}{\to} X.\)

**Definition 4.1.3**

Let \(\nu\) be a measure on the measure space \((\mathbb{R}, \mathcal{B}(\mathbb{R})).\) Then \(\nu\) is called a Lévy measure, if \(\nu(\{0\}) = 0\) and

\[
\int_{\mathbb{R}} \min\{y^2, 1\} \nu(dy) < \infty. \tag{4.1.1}
\]

Note

- Apparently every Lévy measure is \(\sigma\)-finite and

\[
\nu((-\varepsilon, \varepsilon))^c < \varepsilon, \quad \text{for all } \varepsilon > 0, \tag{4.1.2}
\]

where \((-\varepsilon, \varepsilon)^c = \mathbb{R} \setminus (-\varepsilon, \varepsilon).\)
In particular every finite measure $\nu$ is a Lévy measure, if $\nu(\{0\}) = 0$.

If $\nu(dy) = g(y)dy$ then $g(y) \sim \frac{1}{|y|^{\delta}}$ for $y \to 0$, where $\delta \in [0, 3)$.

An equivalent condition to (4.1.2) is

$$\int_{\mathbb{R}} \frac{y^2}{1 + y^2} \nu(dy) < \infty, \quad \text{since} \quad \frac{y^2}{1 + y^2} \leq \min \{y^2, 1\} \leq 2 \frac{y^2}{1 + y^2}.$$  \hspace{1cm} (4.1.3)

**Theorem 4.1.2**

Let $a \in \mathbb{R}$, $b \geq 0$ be arbitrary and let $\nu$ be an arbitrary Lévy measure. Let the characteristic function of a random variable $X : \Omega \to \mathbb{R}$ be given through the function $\varphi : \mathbb{R} \to \mathbb{C}$ with

$$\varphi(s) = \exp \left\{ is - \frac{bs^2}{2} + \int_{\mathbb{R}} \left( e^{isy} - 1 - isy \mathbb{1}_{y \in (-1, 1)} \right) \nu(dy) \right\} \quad \text{for all} \quad s \in \mathbb{R}. \quad (4.1.4)$$

Then $X$ is infinitely divisible.

**Remark 4.1.2**

- The formula (4.1.4) is also called Lévy-Khintchine formula.

- The inversion of Theorem 4.1.2 also holds, hence every infinitely divisible random variable has such a representation. Therefore the characteristic triplet $(a, b, \nu)$ is also called Lévy characteristic of an infinitely divisible random variable.

**Proof of Theorem 4.1.2**

1st step

Show that $\varphi$ is a characteristic function.

For $y \in (-1, 1)$ it holds

$$|e^{isy} - 1 - isy| = \sum_{k=0}^{\infty} \frac{(isy)^k}{k!} - 1 - isy \leq y^2 \left| \sum_{k=2}^{\infty} \frac{s^k}{k!} \right| \leq y^2 c$$

Hence it follows from (4.1.1) that the integral in (4.1.4) exists and therefore it is well-defined.

- Let now $\{c_n\}$ be an arbitrary sequence of numbers with $c_n > c_{n+1} \ldots > 0$ and $\lim_{n \to \infty} c_n = 0$. Then the function $\varphi_n : \mathbb{R} \to \mathbb{C}$ with

$$\varphi_n(s) := \exp \left\{ is \left( a - \int_{[-c_n,c_n] \cap (-1,1)} y \nu(dy) \right) - \frac{bs^2}{2} \right\} \exp \left\{ \int_{[-c_n,c_n]} \left( e^{isy} - 1 \right) \nu(dy) \right\}$$

is the characteristic function of the sum of independent random variables $Z_1^{(n)}$ and $Z_2^{(n)}$, since

- the first factor is the characteristic function of the normal distribution with expectation $a - \int_{[-c_n,c_n] \cap (-1,1)} y \nu(dy)$ and variance $b$.
- the second factor is the characteristic function of a compound Poisson process with parameters

$$\lambda = \nu([-c_n, c_n]^c) \quad \text{and} \quad P_U(\cdot) = \nu(\cdot \cap [-c_n, c_n]^c) / \nu([-c_n, c_n]^c).$$
Furthermore \( \lim_{n \to \infty} \varphi_n(s) = \varphi(s) \) for all \( s \in \mathbb{R} \), where \( \varphi \) is obviously continuous in 0, since it holds for the function \( \psi : \mathbb{R} \to \mathbb{C} \) in the exponent of (4.1.4)

\[
\psi(s) = \int_{\mathbb{R}} \left( e^{isy} - 1 - isy \mathbb{1}_{y \in (-1, 1)} \right) \nu(dy) \quad \text{for all } s \in \mathbb{R}
\]

that \( |\psi(s)| = cs^2 \int_{(-1,1)} y^2 \nu(dy) + \int_{(-1,1)^c} |e^{isy} - 1| \nu(dy) \). Out of this and from (4.1.3) it follows by Lebesgue’s theorem that \( \lim_{s \to 0} \psi(s) = 0 \).

Lemma 4.1.2 yields that the function \( \varphi \) given in (4.1.4) is the characteristic function of a random variable.

2nd step
The infinite divisibility of this random variable follows from Lemma 4.1.1 and out of the fact, that for arbitrary \( n \in \mathbb{N} \) \( \frac{\nu}{n} \) is also a Lévy measure and that

\[
\varphi(s) = \exp \left\{ \frac{i}{n} s - \frac{b s^2}{2} + \int_{\mathbb{R}} \left( e^{isy} - 1 - isy \mathbb{1}_{y \in (-1, 1)} \right) \left( \frac{\nu}{n} \right)(dy) \right\} \quad \text{for all } s \in \mathbb{R}.
\]

Remark 4.1.3
The map \( \eta : \mathbb{R} \to \mathbb{C} \) with

\[
\eta(s) = ias - \frac{bs^2}{2} + \int_{\mathbb{R}} \left( e^{isy} - 1 - isy \mathbb{1}_{y \in (-1, 1)} \right) \nu(dy)
\]

from (4.1.4) is called Lévy exponent of this infinitely divisible distribution.

4.1.2 Lévy-Khintchine Representation

\( \{X(t), \ t \geq 0\} \) – Lévy process. We want to represent the characteristic function of \( X(t), \ t \geq 0 \), through the Lévy-Khintchine formula.

Lemma 4.1.3
Let \( \{X(t), \ t \geq 0\} \) be a stochastically continuous process, i.e. for all \( \varepsilon > 0 \) and \( t_0 \geq 0 \) it holds that \( \lim_{t \to t_0} \mathbb{P}(|X(t) - X(t_0)| > \varepsilon) = 0 \). Then for every \( s \in \mathbb{R} \), \( t \mapsto \varphi_{X(t)}(s) \) is a continuous map from \([0, \infty)\) to \( \mathbb{C} \).

Proof

- \( y \mapsto e^{isy} \) continuous in 0, i.e. for all \( \varepsilon > 0 \) there exists a \( \delta_1 > 0 \), such that

\[
\sup_{y \in (-\delta_1, \delta_1)} \left| e^{isy} - 1 \right| < \frac{\varepsilon}{2}.
\]

- \( \{X(t), \ t \geq 0\} \) is stochastically continuous, i.e. for all \( t_0 \geq 0 \) there exists a \( \delta_2 > 0 \), such that

\[
\sup_{t \geq 0, \ |t-t_0| < \delta_2} \mathbb{P}(|X(t) - X(t_0)| > \delta_1) < \frac{\varepsilon}{4}.
\]
Hence, it follows that for \( s \in \mathbb{R}, t \geq 0 \) and \( |t - t_0| < \delta_2 \) it holds
\[
|\varphi_X(t)(s) - \varphi_X(t_0)(s)| = |E \left( e^{isX(t)} - e^{isX(t_0)} \right)| \leq E \left| e^{isX(t_0)} \left( e^{is(X(t) - X(t_0))} - 1 \right) \right| \\
\leq E \left| e^{is(X(t) - X(t_0))} - 1 \right| = \int_{\mathbb{R}} |e^{isy} - 1| P_{X(t)-X(t_0)}(dy) \\
\leq \int_{(-\delta_1,\delta_1)} |e^{isy} - 1| P_{X(t)-X(t_0)}(dy) \\
+ \int_{(-\delta_1,\delta_1)^C} \left| e^{isy} - 1 \right| P_{X(t)-X(t_0)}(dy) \\
\leq \sup_{y \in (-\delta_1,\delta_1)} \left| e^{isy} - 1 \right| + 2P( |X(t) - X(t_0)| > \delta_1) \leq \varepsilon.
\]

\[ \square \]

**Theorem 4.1.3**

Let \( \{X(t), t \geq 0\} \) be a Lévy process. Then for all \( t \geq 0 \) it holds
\[ \varphi_{X(t)}(s) = e^{t\eta(s)}, \quad s \in \mathbb{R}, \]
where \( \eta : \mathbb{R} \to \mathbb{C} \) is a continuous function. In particular it holds that
\[ \varphi_{X(t)}(s) = e^{t\eta(s)} = \left( e^{\eta(s)} \right)^t = \left( \varphi_{X(1)}(s) \right)^t, \quad \text{for all } s \in \mathbb{R}, t \geq 0. \]

**Proof** Due to stationarity and independence of increments we have
\[ \varphi_{X(t+t')}(s) = E e^{isX(t+t')} = E \left( e^{isX(t)} e^{is(X(t+t') - X(t))} \right) = \varphi_{X(t)}(s) \varphi_{X(t')}(s), \quad s \in \mathbb{R}. \]

Let \( g_s : [0, \infty) \to \mathbb{C} \) be defined by \( g_s(t) = \varphi_{X(t)}(s), s \in \mathbb{R}, g_s(t + t') = g_s(t)g_s(t'), t, t' \geq 0. \)
\( X(0) = 0. \)
\[
\begin{cases} 
    g_s(t + t') = g_s(t)g_s(t'), & t, t' \geq 0, \\
    g_s(0) = 1, \\
    g_s : [0, \infty) \to \mathbb{C} \text{ continuous.}
\end{cases}
\]

Hence there exists \( \eta : \mathbb{R} \to \mathbb{C}, \) such that \( g_s(t) = e^{t\eta(s)} \) for all \( s \in \mathbb{R}, t \geq 0. \) \( \varphi_{X(1)}(s) = e^{\eta(s)} \) and it follows that \( \eta \) is continuous. \[ \square \]

**Lemma 4.1.4 (Prokhorov):**

Let \( \mu_1, \mu_2, \ldots \) be a sequence of finite measures (on \( B(\mathbb{R}) \)) with
1. \( \sup_{n \geq 1} \mu_n(\mathbb{R}) < c, c = \text{const} < \infty \) (uniformly bounded)
2. for all \( \varepsilon > 0 \) there exists \( B_\varepsilon \in B(\mathbb{R}) \) compact, such that fulfills the tightness condition \( \sup_{n \geq 1} \mu_n(B_\varepsilon) \leq \varepsilon. \) Hence follows that there exists a subsequence \( \mu_{n_1}, \mu_{n_2}, \ldots \) and a finite measure over \( B(\mathbb{R}) \), such that for all \( f : \mathbb{R} \to \mathbb{C} \), bounded, continuous, it holds that
\[
\lim_{k \to \infty} \int_{\mathbb{R}} f(y) \mu_{n_k}(dy) = \int_{\mathbb{R}} f(y) \mu(dy)
\]

**Proof** See [14], page 122 - 123. \[ \square \]
Theorem 4.1.4
Let \( \{X(t), t \geq 0\} \) be a Lévy process. Then there exist \( a \in \mathbb{R}, b \geq 0 \) and a Lévy measure \( \nu \), such that
\[
\varphi_{X(1)}(s) = e^{ias - \frac{b s^2}{2}} + \int_{\mathbb{R}} \left( e^{i s y} - 1 - i y 1(y \in (-1, 1)) \right) \nu(dy), \quad \text{for all } s \in \mathbb{R}.
\]

Proof For all sequences \( (t_n)_{n \in \mathbb{N}} \subseteq (0, \infty) \), with \( \lim_{n \to \infty} t_n = 0 \), it holds
\[
\eta(s) = \left( e^{i \eta(s)} \right)' \bigg|_{t=0} = \lim_{n \to \infty} \frac{e^{i \eta(s)} - 1}{i t_n} = \lim_{n \to \infty} \frac{\varphi_{X(t_n)}(s) - 1}{t_n}, \quad (4.1.5)
\]
since \( \eta : \mathbb{R} \to \mathbb{C} \) is continuous. The latter convergence is even uniform in \( s \in [-s_0, s_0] \) for any \( s_0 > 0 \), since Taylor’s theorem yields
\[
\begin{align*}
\lim_{n \to \infty} \left| \eta(s) - \frac{e^{i \eta(s)} - 1}{t_n} \right| &= \lim_{n \to \infty} \left| \eta(s) - \frac{1}{t_n} \sum_{k=1}^{\infty} \frac{(t_n \eta(s))^k}{k!} \right| \\
&= \lim_{n \to \infty} \left| 1 - \sum_{k=2}^{\infty} \frac{(t_n \eta(s))^k}{k!} \right| \\
&= \lim_{n \to \infty} \left| \eta''(s)t_n \sum_{k=1}^{\infty} \frac{(t_n \eta(s))^{k-1}}{(k + 1)!} \right| \\
&\leq \lim_{n \to \infty} M^2 t_n \sum_{k=1}^{\infty} \frac{|t_n M|^{k-1}}{(k + 1)!}, \quad \text{where } M := \sup_{s \in [-s_0, s_0]} |\eta(s)| < \infty \\
&= \lim_{n \to \infty} M^2 t_n \sum_{k=1}^{\infty} \frac{1}{k(k + 1)} \\
&= \lim_{n \to \infty} M^2 t_n e^{\int_{-s_0}^{s_0} M} \\
&= 0.
\end{align*}
\]

Now let \( t_n = \frac{1}{n} \) and \( P_n \) be the distribution of \( X(\frac{1}{n}) \). Hence it follows that
\[
\lim_{n \to \infty} n \int_{\mathbb{R}} (e^{i s y} - 1) P_n(dy) = \lim_{n \to \infty} n \frac{\varphi_{X(\frac{1}{n})}(s) - 1}{\frac{1}{n}} = \eta(s)
\]
\[
\lim_{n \to \infty} \int_{\mathbb{R}} \int_{-s_0}^{s_0} \left( e^{i s y} - 1 \right) P_n(dy) ds = \int_{-s_0}^{s_0} \eta(s) ds
\]
and consequently
\[
\lim_{n \to \infty} n \int_{\mathbb{R}} \left( 1 - \frac{\sin(s_0 y)}{s_0 y} \right) P_n(dy) = \lim_{n \to \infty} n \int_{\mathbb{R}} -\frac{1}{2s_0} \int_{-s_0}^{s_0} \left( e^{i s y} - 1 \right) ds P_n(dy) = -\frac{1}{2s_0} \int_{-s_0}^{s_0} \eta(s) ds.
\]
Since $\eta : \mathbb{R} \to \mathbb{C}$ is continuous with $\eta(0) = 0$ and it follows from the mean value theorem, that for all $\varepsilon > 0$ it exists $\delta_0 > 0$, such that $\left| -\frac{1}{2s_0} \int_{s_0}^{s_0} \eta(s)ds \right| < \varepsilon$. Since $1 - \frac{\sin(s_0y)}{sy} \geq \frac{1}{2}$, for $|s_0y| \geq 2$, it holds: for all $\varepsilon > 0$ there exist $s_0 > 0$, $n_0 > 0$, such that

$$\limsup_{n \to \infty} \frac{n}{2} \int_{\{y:|y| \geq \frac{2}{s_0}\}} P_n(dy) \leq \limsup_{n \to \infty} n \int_{\mathbb{R}} \left(1 - \frac{\sin(s_0)S}{sy}\right) P_n(dy) < \varepsilon.$$ 

Hence for all $\varepsilon > 0$ there exist $s_0 > 0$, $n_0 > 0$, such that

$$n \int_{\{y:|y| \geq \frac{2}{s_0}\}} P_n(dy) \leq 4\varepsilon, \quad \text{for all } n \geq n_0.$$ 

Decreasing $s_0$ gives

$$n \int_{\{y:|y| \geq \frac{2}{s_0}\}} P_n(dy) \leq 4\varepsilon, \quad \text{for all } n \geq 1.$$ 

$$\frac{y^2}{1+y^2} \leq c \left(1 - \frac{\sin y}{y}\right), \quad \text{for all } y \neq 0 \quad \text{and a } c > 0.$$ 

Hence, it follows that

$$\sup_{n \geq 1} \int_{\mathbb{R}} \frac{y^2}{1+y^2} P_n(dy) \leq c' \quad \text{for a } c' < \infty.$$ 

Let now $\mu_n : \mathcal{B}(\mathbb{R}) \to [0, \infty)$ be defined as

$$\mu_n(B) = n \int_{B} \frac{y^2}{1+y^2} P_n(dy) \quad \text{for all } B \in \mathcal{B}(\mathbb{R}).$$ 

It follows that $\{\mu_n\}_{n \in \mathbb{N}}$ is uniformly bounded, $\sup_{n \geq 1} \mu_n(\mathbb{R}) < c'$. Furthermore holds $\frac{y^2}{1+y^2} \leq 1$, $\sup_{n \geq 1} \mu_n \left(\{y: |y| > \frac{2}{s_0}\}\right) \leq 4\varepsilon$ and $\{\mu_n\}_{n \in \mathbb{N}}$ relatively compact. After lemma 4.1.3 it holds: there exists $\{\mu_{n_k}\}_{k \in \mathbb{N}}$, such that

$$\lim_{k \to \infty} \int_{\mathbb{R}} f(y) \mu_{n_k}(dy) = \int_{\mathbb{R}} f(y) \mu(dy)$$

for a measure $\mu$ and $f$ continuous and bounded. Let for $s \in \mathbb{R}$ the function $f_s : \mathbb{R} \to \mathbb{C}$ be defined as

$$f_s(y) = \begin{cases} (e^{isy} - 1 - is \sin y) \frac{1+y^2}{y^2}, & y \neq 0, \\ \frac{-s^2}{2}, & \text{otherwise}. \end{cases}$$

Hence follows that $f_s$ is bounded and continuous and

$$\eta(s) = \lim_{n \to \infty} n \int_{\mathbb{R}} \left(e^{isy} - 1\right) P_n(dy) = \lim_{n \to \infty} \left(\int_{\mathbb{R}} f_s(y) \mu_n(dy) + is \int_{\mathbb{R}} \sin y P_n(dy)\right) = \lim_{k \to \infty} \left(\int_{\mathbb{R}} f_s(y) \mu_{n_k}(dy) + isn_k \int_{\mathbb{R}} \sin y P_{n_k}(dy)\right) = \int_{\mathbb{R}} f_s(y) \mu(dy) + \lim_{k \to \infty} isn_k \int_{\mathbb{R}} \sin y P_{n_k}(dy)$$

$$\eta(s) = i a^* s - \frac{bs^2}{2} + \int_{\mathbb{R}} \left(e^{isy} - 1 - is \sin y\right) \nu(dy),$$
for all \( s \in \mathbb{R} \) with \( a' = \lim_{k \to \infty} isn_k \int_{\mathbb{R}} \sin y P_{n_k}(dy) < \infty \), \( b = \mu(\{0\}) \), \( \nu : \mathcal{B}(\mathbb{R}) \to [0, \infty) \),

\[
\nu(dy) = \begin{cases} 
\frac{1+y^2}{y^2} \mu(dy), & y \neq 0, \\
0, & y = 0.
\end{cases}
\]

\( \int_{\mathbb{R}} |y1(y \in (-1,1)) - \sin y| \nu(dy) < \infty. \)

\( |y1(y \in (-1,1)) - \sin y| \frac{1+y^2}{y^2} < c'', \) for all \( y \neq 0 \) and a \( c'' > 0. \)

Hence follows that

\[
\eta(s) = i as - \frac{bs^2}{2} + \int_{\mathbb{R}} \left(e^{isy} - 1 - isy1(y \in (-1,1))\right) \nu(dy), \quad \text{for all } s \in \mathbb{R}.
\]

\[
a = a' + \int_{\mathbb{R}} (y1(y \in (-1,1)) - \sin y) \nu(dy).
\]

\[\square\]

### 4.1.3 Examples

1. Wiener process (it is enough to look at \( X(1) \))

\( X(1) \sim \mathcal{N}(0,1), \varphi_{X(1)}(s) = e^{-\frac{s^2}{2}} \) and hence follows

\[
(a,b,\nu) = (0,1,0).
\]

Let \( X = \{X(t), t \geq 0\} \) be a Wiener process with drift \( \mu \), i.e. \( X(t) = \mu t + \sigma W(t), W = \{W(t), t \geq 0\} \) - Brownian motion. It follows

\[
(a,b,\nu) = (\mu, \sigma^2, 0).
\]

\[
\varphi_{X(1)}(s) = \mathbb{E} e^{isX(1)} = \mathbb{E} e^{(\mu + \sigma W(1))is} = e^{is\mu - \sigma^2 s^2/2}, \quad s \in \mathbb{R}.
\]

2. Compound Poisson process with parameters \((\lambda, \mathbb{P}_n)\)

\( X(t) = \sum_{i=1}^{N(t)} U_i, \quad N(t) \sim \text{Pois}(\lambda t), \quad U_i \text{ i.i.d.} \sim \mathbb{P}_U. \)

\[
\varphi_{X(1)}(s) = \exp \left\{ \lambda \int_{\mathbb{R}} \left(e^{isx} - 1 \right) \mathbb{P}_U(dx) \right\}
\]

\[
= \exp \left\{ \lambda is \int_{\mathbb{R}} x1(x \in [-1,1]) \mathbb{P}_U(dx) + \lambda \int_{\mathbb{R}} \left(e^{isx} - 1 - isx1(x \in [-1,1])\right) \mathbb{P}_U(dx) \right\}
\]

\[
= \exp \left\{ \lambda is \int_{-1}^{1} x \mathbb{P}_U(dx) + \lambda \int_{\mathbb{R}} \left(e^{isx} - 1 - isx1(x \in [-1,1])\right) \mathbb{P}_U(dx) \right\}, \quad s \in \mathbb{R}.
\]

Hence follows

\[
(a,b,\nu) = \left( \lambda \int_{-1}^{1} x \mathbb{P}_U(dx), 0, \mathbb{P}_U \right), \quad \mathbb{P}_U \text{ finite on } \mathbb{R}.
\]
3. Process of Gauss-Poisson type

\( X = \{X(t), t \geq 0\}, X(t) = X_1(t) + X_2(t), t \geq 0. \)

\( X_1 = \{X_1(t), t \geq 0\} \) and \( X_2 = \{X_2(t), t \geq 0\} \) independent.

\( X_1 \) – Wiener process with drift \( \mu \) and variance \( \sigma^2, \)
\( X_2 \) – Compound Poisson process with parameters \( \lambda, P_U. \)

\[ \phi_{X(t)}(s) = \phi_{X_1(t)}(s)\phi_{X_2(t)}(s) \]
\[ = \exp \left( is\mu - \frac{\sigma^2 s^2}{2} + \lambda \int_{\mathbb{R}} e^{isx} - 1 P_U(dx) \right) \]
\[ = \exp \left\{ is \left( \mu + \lambda \int_{-1}^{1} x P_U(dx) \right) - \frac{\sigma^2 s^2}{2} \right\} \]
\[ + \int_{\mathbb{R}} \lambda \left( e^{isx} - 1 - isx 1(x \in [-1,1]) \right) P_U(dx) \], \( s \in \mathbb{R}. \)

Hence follows
\[ (a, b, \nu) = \left( \mu + \lambda \int_{-1}^{1} x P_U(dx), \sigma^2, \lambda P_U \right). \]

4. Stable Lévy process

\( X = \{X(t), t \geq 0\} \) – Lévy process with \( X(t) \sim \alpha \) stable distribution, \( \alpha \in (0, 2]. \) To introduce \( \alpha \)-stable laws \( \nu, \) let us begin with an example.

If \( X = W \) (Wiener process), then \( X(1) \sim \mathcal{N}(0, 1). \) Let \( Y, Y_1, \ldots, Y_n \) be i.i.d. \( \mathcal{N}(\mu, \sigma^2)\)-variables. Since the normal distribution is stable w.r.t. convolution it holds

\[ Y_1 + \ldots + Y_n \sim \mathcal{N}(n\mu, n\sigma^2) \]
\[ d \]
\[ = \sqrt{n} Y + n\mu - \sqrt{n} \mu \]
\[ = \sqrt{n} Y + \mu \left( n - \sqrt{n} \right) \]
\[ = n^{\frac{1}{2}} Y + \mu \left( n^{\frac{3}{2}} - n^{\frac{1}{2}} \right) \]
\[ = n^{\frac{1}{2}} Y + \mu \left( n - n^{\frac{1}{2}} \right), \quad \alpha = 2. \]

**Definition 4.1.4**

The distribution of a random variable \( Y \) is called \( \alpha \)-stable, if for all \( n \in \mathbb{N} \) independent copies \( Y_1, \ldots, Y_n \) of \( Y \) exist, such that

\[ Y_1 + \ldots + Y_n \overset{d}{=} n^{\frac{1}{2}} Y + d_n, \]

where \( d_n \) is deterministic. The constant \( \alpha \in (0, 2] \) is called index of stability.

Moreover, one can show that

\[ d_n = \begin{cases} \mu \left( n - n^{\frac{1}{2}} \right), & \alpha \neq 1, \\ \mu n \log n, & \alpha = 1. \end{cases} \]

**Example 4.1.1**

- \( \alpha = 2: \) **Normal distribution**, with any mean and any variance.
- \( \alpha = 1: \) **Cauchy distribution** with parameters \( (\mu, \sigma^2) \). The density:

\[ f_Y(x) = \frac{\sigma}{\pi \left( (x - \mu)^2 + \sigma^2 \right)}, \quad x \in \mathbb{R}. \]

It holds \( EY^2 = \infty, EY \) does not exist.
\* \( \alpha = \frac{1}{2} \): Lévy distribution with parameters \((\mu, \sigma^2)\). The density:

\[
f_Y(x) = \begin{cases} 
\left( \frac{2}{\pi} \right)^{\frac{1}{2}} \frac{1}{(x-\mu)^{\frac{3}{2}}} \exp \left\{ -\frac{\sigma}{\pi(x-\mu)} \right\}, & x > \mu, \\
0, & \text{otherwise}.
\end{cases}
\]

These examples are the only examples of \(\alpha\)-stable distribution, where an explicit form of the density is available. For other \(\alpha \in (0, 2)\), \(\alpha \neq \frac{1}{2}, 1\), the \(\alpha\)-stable distribution is introduced through its characteristic function. In general holds: If \(Y\) \(\alpha\)-stable, \(\alpha \in (0, 2]\), then \(\mathbb{E}|Y|^p < \infty, 0 < p < \alpha\).

**Definition 4.1.5**
The distribution of a random variable is called symmetric, if \(Y \overset{d}{=} -Y\).
If \(Y\) has a symmetric \(\alpha\)-stable distribution, \(\alpha \in (0, 2]\), then

\[
\varphi_Y(s) = \exp \{-c|s|^\alpha\}, \ s \in \mathbb{R}.
\]

Indeed, it follows from the stability of \(Y\) that

\[
(\varphi_Y(s))^n = e^{id_n s} \varphi_Y \left( n^{\frac{\alpha}{2}} s \right), \ s \in \mathbb{R}.
\]

It follows that \(d_n = 0\), since \(\varphi_{-Y}(s) = \varphi_Y(s) = \text{varphi}_Y(-s)\). It holds: \(e^{id_n s} = e^{-id_n s}, s \in \mathbb{R}\) and \(d_n = 0\). The rest is left as an exercise.

**Lemma 4.1.5**
Lévy-Khintchine representation of the characteristic function of a stable distribution. Any stable law is infinitely divisible with the Lévy triplet \((a, b, \nu)\), where \(a \in \mathbb{R}\) arbitrary,

\[
b = \begin{cases} 
\sigma^2, & \alpha = 2, \\
0, & \alpha < 2.
\end{cases}
\]

and

\[
\nu(dx) = \begin{cases} 
0, & \alpha = 2, \\
\frac{c_1}{x^{a-1}} 1(x \geq 0)dx + \frac{c_2}{|x|^{a+1}} 1(x < 0)dx, & \alpha < 2, c_1, c_2 \geq 0 : c_1 + c_2 > 0
\end{cases}
\]

Without proof

**Exercise:** Prove that

\[
P (|Y| \geq x) \underset{x \to \infty}{\sim} \begin{cases} 
\frac{e^{-x^2}}{\pi^\frac{\alpha}{2}}, & \alpha = 2, \\
\frac{c}{\pi^\frac{\alpha}{2}}, & \alpha < 2.
\end{cases}
\]

**Definition 4.1.6**
The Lévy process \(X = \{X(t), t \geq 0\}\) is called stable, if \(X(1)\) has an \(\alpha\)-stable distribution, \(\alpha \in (0, 2]\) (\(\alpha = 2\): Brownian motion (with drift)).

### 4.1.4 Subordinators

**Definition 4.1.7**
A Lévy process \(X = \{X(t), t \geq 0\}\) is called subordinator, if for all \(0 < t_1 < t_2\), \(X(t_1) \leq X(t_2)\) a.s.

Since

\[
X(0) = 0 \text{ a.s. } \Rightarrow X(t) \geq 0, \ t \geq 0, \text{ a.s.}
\]
This class of Subordinators is important since you can easily introduce \( \int_a^b g(t) dX(t) \) as a Lebesgue-Stieltjes-integral.

**Theorem 4.1.5**

The Lévy process \( X = X(t), \ t \geq 0 \) is a subordinator if and only if the Lévy-Khintchine representation can be expressed in the form

\[
\varphi_{X(1)}(s) = \exp \left\{ ias + \int_{\mathbb{R}} \left( e^{i\sigma x} - 1 \right) \nu(dx) \right\}, \quad s \in \mathbb{R},
\]

where \( a \in [0, \infty) \) and \( \nu \) is the Lévy measure, with

\[
\nu((-\infty, 0)) = 0, \quad \int_0^\infty \min\{1, y^2\} \nu(dy) < \infty.
\]

**Proof** Sufficiency

It has to be shown that \( X(t_2) \geq X(t_1) \) a.s., if \( t_2 \geq t_1 \geq 0 \).

First of all we show that \( X(1) \geq 0 \) a.s.. If \( \nu = 0 \), then \( X(1) = a \geq 0 \) a.s., hence

\[
\varphi_{X(t)}(s) = \left( \varphi_{X(1)}(s) \right)^t = e^{iat}, \quad s \in \mathbb{R}.
\]

\( X(t) = at \) a.s. and therefore it follows that \( X(t) \uparrow \) and \( X \) is a subordinator.

If \( \nu([0, \infty)) > 0 \), then there exists \( N > 0 \) such that for all \( n \geq N \) it holds \( 0 < \nu \left( \left[ \frac{1}{n}, \infty \right) \right) < \infty \).

It follows

\[
\varphi_{X(1)}(s) = \exp \left\{ ias + \lim_{n \to \infty} \frac{1}{n} \int_{\left[ \frac{1}{n}, \infty \right)} \left( e^{i\sigma x} - 1 \right) \nu(dx) \right\} = e^{ias} \lim_{n \to \infty} \varphi_n(s), \quad s \in \mathbb{R},
\]

where \( \varphi_n(s) = \int_{\left[ \frac{1}{n}, \infty \right)} \left( e^{i\sigma x} - 1 \right) \nu(dx) \) is the characteristic function of a compound Poisson process distribution with parameters \( \nu \left( \left[ \frac{1}{n}, \infty \right) \right), \frac{\nu\left( \left[ \frac{1}{n}, \infty \right) \right)}{\nu\left( \left[ \frac{1}{n}, \infty \right) \right)} \) for all \( n \in \mathbb{N} \). Let \( Z_n \) be the random variable with characteristic function \( \varphi_n \). It holds: \( Z_n = \sum_{i=1}^{N_{t_i}} U_i, \ N_n \sim \text{Pois} \left( \nu \left( \left[ \frac{1}{n}, \infty \right) \right) \right), \)

\( U_i \sim \frac{\nu\left( \left[ \frac{1}{n}, \infty \right) \right)}{\nu\left( \left[ \frac{1}{n}, \infty \right) \right)} \), hence follows \( Z_n \geq 0 \) a.s. and \( X(1) = \frac{a}{\lim_{\mathbb{R} \geq 0} Z_n} \geq 0 \) a.s.. Since \( X \) is a Lévy process, it holds

\[
X(1) = X\left( \frac{1}{n} \right) + X\left( \frac{2}{n} \right) - X\left( \frac{1}{n} \right) + \ldots + X\left( \frac{n}{n} \right) - X\left( \frac{n-1}{n} \right),
\]

where, because of stationarity and independence of the increments, \( X\left( \frac{k}{n} \right) - X\left( \frac{k-1}{n} \right) \) a.s. \( \geq 0 \) for \( 1 \leq k \leq n \) for all \( n \). Hence \( X(q_2) - X(q_1) \geq 0 \) a.s. for all \( q_1, q_2 \in \mathbb{Q}, q_2 \geq q_1 \geq 0 \). Now let \( t_1, t_2 \in \mathbb{R} \) such that \( 0 \leq t_1 \leq t_2 \). Let \( \{q_1^{(n)}, q_2^{(n)}\} \) be sequences of numbers from \( \mathbb{Q} \) with \( q_1^{(n)} \leq q_2^{(n)} \) such that \( q_1^{(n)} \downarrow t_1, q_2^{(n)} \uparrow t_2, n \to \infty \). For \( \varepsilon > 0 \)

\[
\Pr\left( X(t_2) - X(t_1) < -\varepsilon \right) = \Pr\left( X(t_2) - X\left( q_2^{(n)} \right) + X\left( q_2^{(n)} \right) - X\left( q_1^{(n)} \right) + X\left( q_1^{(n)} \right) - X(t_1) < -\varepsilon \right)
\]

\[
\leq \Pr\left( X(t_2) - X\left( q_2^{(n)} \right) + X\left( q_1^{(n)} \right) - X(t_1) < -\varepsilon \right)
\]

\[
\leq \Pr\left( X(t_2) - X\left( q_2^{(n)} \right) < -\varepsilon \right) + \Pr\left( X\left( q_1^{(n)} \right) - X(t_1) < -\frac{\varepsilon}{2} \right) \to 0 \quad \text{as} \quad n \to \infty,
\]
since \( X \) is stochastically continuous. Then
\[
\begin{align*}
\mathbb{P} ( X(t_2) - X(t_1) < \varepsilon) &= 0 & \text{for all } \varepsilon > 0 \\
\mathbb{P} ( X(t_2) - X(t_1) < 0) &= \lim_{\varepsilon \to +0} \mathbb{P} ( X(t_2) - X(t_1) < \varepsilon) = 0 \\
\Rightarrow X(t_2) &\geq X(t_1) \quad \text{a.s.}
\end{align*}
\]

**Necessity**

Let \( X \) be a Lévy process, which is a subordinator. It has to be shown that \( \varphi_{X(1)}(\cdot) \) has the form (4.1.6).

After the Lévy-Khintchine representation for \( X(1) \) it holds that
\[
\varphi_{X(1)}(s) = \exp \left\{ is - \frac{b^2 s^2}{2} + \int_0^\infty \left( e^{isx} - 1 - isx1(x \in [-1, 1]) \right) \nu(dx) \right\}, \quad s \in \mathbb{R}.
\]

The measure \( \nu \) is concentrated on \([0, \infty)\), since \( X(t) \) a.s. for all \( t \geq 0 \) and from the proof of Theorem 4.1.4 \( \nu((-\infty, 0)) = 0 \) can be chosen.

\[
\varphi_{X(1)}(s) = \exp \left\{ is - \frac{b^2 s^2}{2} \right\} \exp \left\{ \int_0^\infty \left( e^{isx} - 1 - isx1(x \in [-1, 1]) \right) \nu(dx) \right\}
\]

Hence it follows that \( X(1) = Y_1 + Y_2 \), where \( Y_1 \) and \( Y_2 \) are independent, \( Y_1 \sim N(a, b^2) \) and therefore \( b = 0 \). (Otherwise \( Y_1 \) could attain negative values and consequently \( X(1) \) could attain negative values as well.)

For all \( \varepsilon \in (0, 1) \)
\[
\varphi_{X(1)}(s) = \exp \left\{ is \left( a - \int_0^1 x \nu(dx) \right) + \int_0^\infty \left( e^{isx} - 1 - isx \right) \nu(dx) \right\}.
\]

It has to be shown that for \( \varepsilon \to 0 \) it holds \( \int_\varepsilon^\infty \left( e^{isx} - 1 \right) \nu(dx) \to \int_0^\infty \left( e^{isx} - 1 \right) \nu(dx) < \infty \) with \( \int_0^1 \min \{ x, 1 \} \nu(dx) < \infty \). \( \varphi_{X(1)}(s) = \exp \left\{ is \left( a - \int_0^1 x \nu(dx) \right) \right\} \varphi_{Z_1}(s) \varphi_{Z_2}(s) \), where \( Z_1 \) and \( Z_2 \) are independent, \( \varphi_{Z_1}(s) = \exp \left\{ \int_\varepsilon^\infty \left( e^{isx} - 1 - isx \right) \nu(dx) \right\} \), \( \varphi_{Z_2}(s) = \exp \left\{ \int_0^{\infty} \left( e^{isx} - 1 \right) \nu(dx) \right\} \), \( s \in \mathbb{R} \). Then \( X(1) \) \( \overset{d}{=} \) \( a - \int_0^1 x \nu(dx) + Z_1 + Z_2 \). There exist \( \varphi_{Z_1}^{(2)}(0) = -\frac{EZ_1^2}{2} < \infty \), \( \varphi_{Z_1}^{(1)}(0) = 0 = i\varepsilon Z_1 \) and it therefore follows that \( EZ_1 = 0 \) and \( P(Z_1 \leq 0) > 0 \). On the other hand, \( Z_2 \) has a compound Poisson distribution with parameters \( \nu(\varepsilon, \infty)) \), \( \frac{\nu(\varepsilon, +\infty)}{\nu(\varepsilon, +\infty)} \), \( \varepsilon \in (0, 1) \).

\[
\Rightarrow P(Z_2 \leq 0) > 0, \quad \text{since } P(Z_2 = 0) > 0.
\]

\[
\Rightarrow P(Z_1 + Z_2 \leq 0) \geq P(Z_1 \leq 0, Z_2 \leq 0) = P(Z_1 \leq 0) P(Z_2 \leq 0) > 0
\]

For \( X(1) \) to be positive it follows that \( a - \int_0^1 x \nu(dx) \geq 0 \) for all \( \varepsilon \in (0, 1) \). Hence \( a \geq 0 \) and
\[
\int_0^\infty \min \{ x, 1 \} dx < \infty.
\]

Moreover, for \( \varepsilon \downarrow 0 \) it holds \( Z_1 \overset{d}{\to} 0 \) and consequently
\[
\varphi_{X(1)}(s) = \exp \left\{ is \left( a - \int_0^1 x \nu(dx) \right) + \int_0^\infty \left( e^{isx} - 1 \right) \nu(dx) \right\}, \quad s \in \mathbb{R}.
\]
Example 4.1.2 ($\alpha$-stable subordinator):

Let $X = \{X(t), t \geq 0\}$ be a subordinator, with $a = 0$ and the Lévy measure

$$
\nu(dx) = \begin{cases} 
\frac{\alpha}{\Gamma(1-\alpha)} \frac{1}{x^{1+\alpha}} dx, & x > 0, \\
0 \cdot \frac{1}{x^{1+\alpha}} dx = 0, & x \leq 0.
\end{cases}
$$

By Lemma 4.1.5 it follows that $X$ is an $\alpha$-stable Lévy process.

We show that $\hat{l}_{X(t)}(s) = \mathbb{E}e^{-sX(t)} = e^{-ts^\alpha}$ for all $s, t \geq 0$.

$$
\varphi_{X(t)}(s) = \left(\varphi_{X(1)}(s)\right)^t = \exp \left\{ t \int_0^\infty \left( e^{i\alpha x} - 1 \right) \frac{\alpha}{\Gamma(1-\alpha)} \frac{1}{x^{1+\alpha}} dx \right\}, \quad s \in \mathbb{R}.
$$

It has to be shown that

$$
ua = \frac{\alpha}{\Gamma(1-\alpha)} \int_0^\infty (1 - e^{-ux}) \frac{dx}{x^{1+\alpha}}, \quad u \geq 0.
$$

This is enough since $\varphi_{X(t)}(\cdot)$ can be continued analytically to $\{ z \in \mathbb{C} : \text{Im} z \geq 0 \}$, i.e. $\varphi_{X(t)}(iu) = \hat{l}_{X(t)}(u), u \geq 0$. In fact, it holds that

$$
\int_0^\infty \left(1 - e^{-ux}\right) \frac{dx}{x^{1+\alpha}} = \int_0^\infty u \int_0^\infty e^{-uy} \frac{dx}{x^{1+\alpha}} dy
$$

by Fubini,

$$
= \int_0^\infty \int_y^\infty x^{-\alpha} dx\frac{1}{u} dy
$$

Subst.

$$
= \frac{u}{\alpha} \int_0^\infty e^{-uy} y^{-\alpha} dy
$$

and hence follows $\hat{l}_{X(t)}(s) = e^{-ts^\alpha}, t, s \geq 0$.

4.2 Additional Exercises

Exercise 4.2.1

Let $X$ be a random variable with distribution function $F$ and characteristic function $\varphi$. Show that the following statements hold:

a) If $X$ is infinitely divisible, then it holds $\varphi(t) \neq 0$ for all $t \in \mathbb{R}$. Hint: Show that $\lim_{n \to \infty} |\varphi_n(s)|^2 = 1$ for all $s \in \mathbb{R}$, if $\varphi(s) = (\varphi_n(s))^n$. Note further that $|\varphi_n(s)|^2$ is again a characteristic function and $\lim_{n \to \infty} n \frac{x}{u} = 1$ holds for $x > 0$.

b) Give an example (with explanation) for a distribution, which is not infinitely divisible.
Exercise 4.2.2
Let \( X = \{X(t), t \geq 0\} \) be a Lévy process. Show that the random variable \( X(t) \) is then infinitely divisible for every \( t \geq 0 \).

Exercise 4.2.3
Show that the sum of two independent Lévy processes is again a Lévy process, and state the corresponding Lévy characteristic.

Exercise 4.2.4
Look at the following function \( \varphi : \mathbb{R} \to \mathbb{C} \) with
\[
\varphi(t) = e^{\psi(t)}, \quad \text{where} \quad \psi(t) = 2 \sum_{k=-\infty}^{\infty} 2^{-k}(\cos(2^k t) - 1).
\]
Show that \( \varphi(t) \) is the characteristic function of an infinitely divisible distribution. Hint: Look at the Lévy-Khintchine representation with measure \( \nu(\{\pm 2^k\}) = 2^{-k}, k \in \mathbb{Z} \).

Exercise 4.2.5
Let the Lévy process \( \{X(t), t \geq 0\} \) be a Gamma process with parameters \( b, p > 0 \), that is, for every \( t \geq 0 \) it holds \( X(t) \sim \Gamma(b, pt) \). Show that \( \{X(t), t \geq 0\} \) is a subordinator with the Laplace exponent \( \xi(u) = \int_{0}^{\infty} (1 - e^{-uy}) \nu(dy) \) with \( \nu(dy) = py^{-1}e^{-uy}dy, y > 0 \). (The Laplace exponent of \( \{X(t), t \geq 0\} \) is the function \( \xi : [0, \infty) \to [0, \infty) \), for which holds that \( \mathbb{E} e^{-uX(t)} = e^{-t\xi(u)} \) for arbitrary \( t, u \geq 0 \))

Exercise 4.2.6
Let \( \{X(t), t \geq 0\} \) be a Lévy process with characteristic Lévy exponent \( \eta \) and \( \{\tau(s), s \geq 0\} \) an independent subordinator with characteristic Lévy exponent \( \gamma \). The stochastic process \( Y \) be defined as \( Y = \{X(\tau(s)), s \geq 0\} \).

(a) Show that
\[
\mathbb{E} \left( e^{i\theta Y(s)} \right) = e^{\gamma(-i\eta(\theta))s}, \quad \theta \in \mathbb{R},
\]
where \( \text{Im} z \) describes the imaginary part of \( z \).

Hint: Since \( \tau \) is a process with non-negative values, it holds \( \mathbb{E} e^{i\theta \tau(s)} = e^{\gamma(\theta)s} \) for all \( \theta \in \{z \in \mathbb{C} : \text{Im} z \geq 0\} \) through the analytical continuation of Theorem 4.1.3.

(b) Show that \( Y \) is a Lévy process with characteristic Lévy exponent \( \gamma(-i\eta(\cdot)) \).

Exercise 4.2.7
Let \( \{X(t), t \geq 0\} \) be a compound Poisson process with Lévy measure
\[
\nu(dx) = \frac{\lambda}{\sigma \sqrt{\pi}} e^{-\frac{x^2}{2\sigma^2}} dx, \quad x \in \mathbb{R},
\]
where \( \lambda, \sigma > 0 \). Show that \( \{\sigma W(N(t)), t \geq 0\} \) has the same finite-dimensional distributions as \( X \), where \( \{N(s), s \geq 0\} \) is a Poisson process with intensity \( 2\lambda \) and \( W \) is a standard Wiener process independent from \( N \).

Hint to exercise 4.2.6 a) and exercise 4.2.7
- In order to calculate the expectation for the characteristic function, the identity \( \mathbb{E}(X) = \mathbb{E}(\mathbb{E}(X|Y)) = \int_{\mathbb{R}} \mathbb{E}(X|Y = y) F_Y(dy) \) for two random variables \( X \) and \( Y \) can be used. In doing so, it should be conditioned on \( \tau(s) \).
• \( \int_{-\infty}^{\infty} \cos(sy)e^{-\frac{y^2}{2\pi}} dy = \sqrt{2\pi a} \cdot e^{-\frac{a^2}{4}} \) for \( a > 0 \) and \( s \in \mathbb{R} \).

**Exercise 4.2.8**
Let \( W \) be a standard Wiener process and \( \tau \) an independent \( \frac{\alpha}{2} \)-stable subordinator, where \( \alpha \in (0, 2) \). Show that \( \{W(\tau(s)), s \geq 0\} \) is a \( \alpha \)-stable Lévy process.

**Exercise 4.2.9**
Show that the subordinator \( T \) with marginal density
\[
f_T(t)(s) = \frac{t}{2\sqrt{\pi}} s^{-\frac{3}{2}} e^{-\frac{t^2}{4s}} 1\{s > 0\}
\]
is a \( \frac{1}{2} \)-stable subordinator. (Hint: Differentiate the Laplace transform of \( T(t) \) and solve the differential equation)
Bibliography


