Random walks in non homogeneous Poissonian environment

Youri Davydov 1 and Valentin Konakov2, *

1 Université Lille 1 and Saint Peterbourg State University
2 Higher Scool of Economics

1 Introduction

We consider the moving particle process in $\mathbb{R}^d$ which is defined in the following way. There are two independent sequences $(T_k)$ and $(\varepsilon_k)$ of random variables. The variables $T_k$ are non negative and $\forall k \ T_k \leq T_{k+1}$, while variables $\varepsilon_k$ form an i.i.d sequence with common distribution concentrated on the unit sphere $S^{d-1}$.

The values $\varepsilon_k$ are interpreted as the directions, and $T_k$ as the moments of change of directions. A particle starts from zero and moves in the direction $\varepsilon_1$ up to the moment $T_1$. It then changes direction to $\varepsilon_2$ and moves on within the time interval $T_2 - T_1$, etc. The speed is constant at all sites. The position of the particle at time $t$ is denoted by $X(t)$.

Study of the processes of this type has a long history. The first work dates back probably to Pearson and continued by Kluyer (1906) and Rayleigh (1919). Mandelbrot (1982) considered the case where the increments $T_n - T_{n-1}$ form i.i.d. sequence with the common law having a heavy tail. He also introduced the term ”Levy flights” later changed to ”Random flights”.

To date, a large number of works were accumulated, devoted to the study of such processes, we mention here only articles by Kolesnik (2009), Orsingher and De Gregorio (2012, 2015) and Orsingher and Garra (2014) which contain an extensive bibliography and where for different assumptions on $(T_k)$ and $(\varepsilon_k)$ the exact formulas for the distribution of $X(t)$ were derived.

Our goals are different.

Firstly, we are interested in the global behavior of the process $X = \{X(t), \ t \in \mathbb{R}^+\}$, namely, we are looking for conditions under which the processes $\{Y_T, \ T > 0\}$,

$$Y_T(t) = \frac{1}{B(T)}X(tT), \ t \in [0,1],$$

weakly converges in $C[0,1]$ : $Y_T \Rightarrow Y, \ B_T \to \infty, \ T \to \infty$.

From now on we suppose that the points $(T_k), \ T_k \leq T_{k+1}$, form a Poisson point process in $\mathbb{R}^+$ denoted by $\mathbf{T}$.

It is clear that in the homogeneous case the process $X(t)$ is a conventional random walk because the spacings $T_{k+1} - T_k$ are independent, and then the limit process is Brownian motion.

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In the non homogeneous case the situation is more complicated as these spacings are not independent. Nevertheless it was possible to distinguish three modes that determine different types of limiting processes.

For a more precise description of the results it is convenient to assume that $T_k = f(\Gamma_k)$, where $\Pi = (\Gamma_k)$ is a standard homogeneous Poisson point process on $R_+$ with intensity 1. In this case

$$(\Gamma_k) \overset{\mathcal{L}}{=} (\gamma_1 + \gamma_2 + \cdots + \gamma_k)$$

where $(\gamma_k)$ are i.i.d standard exponential random variables.

If the function $f$ has power growth,

$$f(t) = t^\alpha, \alpha \geq 1,$$

the behavior of the process is analogous to the uniform case and then in the limit we obtain a Gaussian process which is a linearly transformed Brownian motion

$$Y(t) = \int_0^t K_\alpha(s)dW(s),$$

where $W$ is a process of Brownian motion, for which the covariance matrix of $W(1)$ coincides with the covariance matrix of $\varepsilon_1$ and $K_\alpha(s)$ is a nonrandom kernel, its exact expression is given below.

In the case of exponential growth,

$$f(t) = e^{t\beta}, \beta > 0,$$

the limiting process is piecewise linear with an infinite number of units, but $\forall \varepsilon > 0$ the number of units in the interval $[\varepsilon, 1]$ will be a.s. finite.

Finally, with the super exponential growth of $f$, the process degenerates: its trajectories are linear functions:

$$Y(t) = \varepsilon t, \quad t \in [0, 1], \quad \varepsilon \overset{\text{Law}}{=} \varepsilon_1.$$

In the second part of the paper the process $X(t)$ is considered as a Markov chain. We construct diffusion approximations for this process and investigate their accuracy. To prove the weak convergence we use the approach of Stroock and Varadhan (1979). Under our assumptions the diffusion coefficients $a$ and $b$ have the property that for each $x \in R^d$ the martingale problem for $a$ and $b$ has exactly one solution $P_x$ starting from $x$ (that is well posed). It remains to check the conditions from Stroock and Varadhan (1979) which imply the weak convergence of our sequence of Markov chains to this unique solution $P_x$. We consider also more general model which may be called as ”random walk over ellipsoids in $R^d$ “. For this model we establish the convergence of transition densities and obtain Edgeworth type expansion up to the order $n^{-3/2}$, where $n$ is a number of switching. The main tool in this part is the paramertix method (Konakov (2012), Konakov and Mammen (2009)).

2 Random flights in Poissonian environment

The reader is reminded that we suppose $T_k = f(\Gamma_k)$, where $(\Gamma_k)$ is a standard homogeneous Poisson point process on $R_+$. Assume also that $E\varepsilon_1 = 0$. 

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It is more convenient to consider at first the behavior of processes

\[ Z_n(t) = Y_{T_n}(t), \]

as for \( T = T_n \) the paths of \( Z_n \) have an integer number of full segments on the interval \([0,1]\). The typical path of \( \{Z_n(t), t \in [0,1]\} \) is a continuous broken line with vertices \( \{t_{n,k}, \frac{k}{B_n}\}, \ k = 0, 1, \ldots, n \}, \) where \( t_{n,k} = \frac{k}{B_n}, \ T_0 = 0, \ B_n = B(T_n), \ S_k = \sum_i \epsilon_i(T_i - T_{i-1}). \)

**Theorem 1.** Under previous assumptions

1) If the function \( f \) has power growth: \( f(t) = t^\alpha, \ \alpha \geq 1/2, \) we take \( B(T) = T^{2\alpha-1}. \)

Then \( Z_n \Rightarrow Y, \) where \( Y \) is a Gaussian process

\[ Y(t) = \sqrt{2\alpha} \int_0^t s^{\frac{\alpha-1}{\alpha}} dW(s), \]

and \( W \) is a process of Brownian motion, for which the covariance matrix of \( W(1) \) coincides with the covariance matrix of \( \epsilon_1. \)

2) If the function \( f \) has exponential growth: \( f(t) = e^{\beta t}, \ \beta > 0, \) we take \( B(T) = T. \)

Then \( Z_n \Rightarrow Y, \) where \( Y \) is a continuous piecewise linear process with the vertices at the points \( (t_k, Y(t_k)), \)

\[ t_k = e^{-\beta \Gamma_{k-1}}, \ \Gamma_0 = 0, \]

\[ Y(t_k) = \sum_{i=k}^{\infty} \epsilon_i(e^{-\beta \Gamma_{i-1}} - e^{-\beta \Gamma_i}), \ \ Y(0) = 0. \]

3) In super exponential case, suppose that \( f \) is increasing absolutely continuous and such that

\[ \lim_{t \to \infty} \frac{f'(t)}{f(t)} = +\infty. \]

We take \( B(T) = T. \)

Then \( \frac{T_n}{T_{n+1}} \to 0 \) in probability, and \( Z_n \Rightarrow Y, \) where the limiting process \( Y \) degenerates:

\[ Y(t) = \epsilon_1 t, \ \ t \in [0,1]. \]

**Remark 1.** In the case of power growth the limiting process admits the following representation:

\[ Y(t) \overset{\mathcal{D}}{=} \alpha \sqrt{\frac{2}{2\alpha - 1}} W(t^{\frac{\alpha-1}{\alpha}}), \]

where, as before, \( W \) is a Brownian motion, for which the covariance matrix of \( W(1) \) coincides with the covariance matrix of \( \epsilon_1. \)

It is clear that we can also express \( Y \) in another way:

\[ Y(t) \overset{\mathcal{D}}{=} \alpha \sqrt{\frac{2}{2\alpha - 1}} K^{\frac{1}{2}} w(t^{\frac{\alpha-1}{\alpha}}), \]

where \( w \) is a standard Brownian motion and \( K \) is the covariance matrix of \( \epsilon_1. \)
Remark 2. In the case of exponential growth it is possible to describe the limiting process $Y$ in the following way:

We take a p.p.p. $T = (t_k)$, $t_k = e^{-\beta T_{k-1}}$, defined on $(0, 1]$, and define a step process

$$Z(t) = \varepsilon_k \quad \text{for} \quad t \in (t_{k+1}, t_k].$$

Then

$$Y(t) = \int_0^t Z(s) \, ds.$$

3 Diffusion approximation

In this section firstly we consider a model of random flight which is equivalent to the study of random broken lines $\{X_n(t), \ t \in [0, 1]\}$ with the vertices $(\xi_n, X_n(\xi_n))$, and such that $(h = \frac{1}{n})$

$$X_n((k + 1)h) = X_n(\xi_n) + hb(X_n(\xi_n)) + \sqrt{h}\xi_k(X(\xi_n)),$$

$$X_n(0) = x_0, \quad \xi_k(X_n(\xi_n)) = b_k\sigma(X_n(\xi_n))\varepsilon_k,$$

where $\{\varepsilon_k\}$ and $\{\rho_k\}$ are two independent sequences and

$\{\varepsilon_k\}$ are i.i.d. r. v. uniformly distributed on the unit sphere $S^{d-1}$,

$\{\rho_k\}$ are i.i.d. r. v. having a density, $\rho_k \geq 0$, $E\rho_k^2 = d$,

$b : R^d \to R^d$ is a bounded measurable function and $\sigma : R^d \to R^d 	imes R^d$ is a bounded measurable matrix function.

Theorem 2. Let $X = \{X(t), \ t \in [0, 1]\}$ be a solution of stochastic equation

$$X(t) = x_0 + \int_0^t b(X(s)) \, ds + \int_0^t \sigma(X(s)) \, dw(s).$$

Suppose that $b$ and $\sigma$ are continuous functions satisfying Lipschitz condition

$$|b(t) - b(s)| + |\sigma(t) - \sigma(s)| \leq K|t - s|.$$

Moreover it is supposed that $b(x)$ and $\frac{1}{\det(\sigma(x))}$ are bounded.

Then,

$$X_n \Rightarrow X \quad \text{in} \quad C[0, 1].$$

Our next result is about approximation of transition density. We consider now more general models given by a triplet $(b(x), \sigma(x), f(r; \theta))$, $x \in R^d$, $r \geq 0$, $\theta \in R^+$, where $b(x)$ is a vector field, $\sigma(x)$ is a $d \times d$ matrix, $a(x) := \sigma a^T(x) > \delta I$, $\delta > 0$, and $f(r; \theta)$ is a radial density depending on a parameter $\theta$ controlling the frequency of changes of directions, namely, the frequency increases when $\theta$ decreases. Suppose $X(0) = x_0$. The vector $b(x_0)$ acts shifting a particle from $x_0$ to $x_0 + \Delta(\theta)b(x_0)$, where $\Delta(\theta) = c_0\theta^2$, $c_0 > 0$. Several examples of such functions $\Delta(\theta)$ for different models will be given below. Define

$$\mathcal{E}_{x_0}(r) := \{x : |a^{-1/2}(x_0)(x - x_0 - \Delta(\theta)b(x_0))|^2 = r^2\},$$

$$\mathcal{S}_{x_0}(r) := \{y : |y - x_0 - \Delta(\theta)b(x_0)|^2 = r^2\}. $$
The initial direction is defined by a random variable $\xi_0$, the law of $\xi_0$ is a pushforward of the spherical measure on $S^d_{x_0}(1)$ under affine change of variables

$$x - x_0 - \Delta(\theta)b(x_0) = a^{1/2}(x_0)(y - x_0 - \Delta(\theta)b(x_0)).$$

Then particle moves along the ray $l_{x_0}$ corresponding to the directional unit vector

$$\varepsilon_0 := \frac{\xi_0 - x_0 - \Delta(\theta)b(x_0)}{|\xi_0 - x_0 - \Delta(\theta)b(x_0)|},$$
and changes the direction in $(r, r + dr)$ with probability

$$\det(a^{-1/2}(x_0)) \cdot f(r \,|\, a^{-1/2}(x_0)\varepsilon_0|)dr.$$  \hspace{1cm} (2)

Let $\rho_0$ be a random variable independent on $\xi_0$ and distributed on $l_{x_0}$ with the radial density (2). We consider the point $x_1 = x_0 + \Delta(\theta)b(x_0) + \rho_0 \varepsilon_0$. Let $(\varepsilon_k, \rho_k)$ be independent copies of $(\varepsilon_0, \rho_0)$. Starting from $x_1$ we repeat the previous construction to obtain $x_2 = x_1 + \Delta(\theta)b(x_1) + \rho_1 \varepsilon_1$. After $n$ switching we get a point $x_n$,

$$x_n = x_{n-1} + \Delta(\theta)b(x_{n-1}) + \rho_{n-1} \varepsilon_{n-1}.$$  

To obtain the one-step characteristic function $\Psi_1(t)$ we make use of formula (6) from Yadrenko (1980):

$$\Psi_1(t) = E e^{i\langle t, \rho \varepsilon \rangle} = \int_0^\infty \int_{E_{x_0}(r)} e^{i\langle t, a^{1/2}(x_0)a^{-1/2}(x_0)\xi \rangle} \mu_{E_{x_0}(r)}(d\xi) d\Phi_{\rho}(r) =$$

$$= \int_0^\infty \int_{S^d_{x_0}(r)} e^{i\langle a^{1/2}(x_0)\xi, y \rangle} \lambda_{r}^d(dy) f(r; \theta) dr =$$

$$= 2^{d/2} \frac{\Gamma(d/2)}{\Gamma((d-2)/2)} \int_0^\infty J_{d-2}(r \,|\, a^{1/2}(x_0)\xi \rangle) \|a^{1/2}(x_0)\xi\|^{d-2} f(r; \theta) dr, \hspace{1cm} (3)$$

where $J_{\nu}(z)$ is the Bessel function and $d\Phi_{\rho}(r)$ is the $F$ - measure of the layer between $E_{x_0}(r)$ and $E_{x_0}(r + dr), F$ is the law of $\rho_0 \varepsilon_0$. Now we make our main assumption about the radial density:

**A1** The function $f(r; \theta)$ is homogenous of degree $-1$, that is

$$f(\lambda r; \lambda \theta) = \lambda^{-1} f(r; \theta), \forall \lambda \neq 0.$$  

Denote by $p_{\varepsilon}(n, x, y)$ the transition density after $n$ switching in the RF-model described above. To obtain the one-step transition density $p_{\varepsilon}(1, x, y)$ (we write $(x, y)$ instead of $(x_0, x_1)$) we use the inverse Fourier transform, (3) and **A1**. After easy calculations we get

$$p_{\varepsilon}(1, x, y) = \Delta'^{-d/2}(\theta) q_x \left( \frac{y - x - \Delta(\theta)b(x)}{\sqrt{\Delta(\theta)}} \right), \hspace{1cm} (4)$$

where

$$q_x(z) = \frac{2^{d/2} \Gamma(d/2)}{(2\pi)^d} \int_{\mathbb{R}^d} \cos(r, z) \left[ \int_0^\infty \frac{J_{d-2}(\rho \,|\, a^{1/2}(x)\tau \rangle)}{(\rho \,|\, a^{1/2}(x)\tau \rangle)^{d-2}} f(\rho; c_\theta) d\rho \right] d\tau. \hspace{1cm} (5)$$

Consider two examples.
**Example 1.** We put $\Delta(\theta) = (d+1)^2\theta^2$ and

$$f(r; \theta) = \frac{1}{\Gamma(d)} r^{-1} \left( \frac{r}{\theta} \right)^d \exp \left( -\frac{r^2}{\theta} \right).$$

Using (3), formula 6.623 (2) on page 726 from Gradshtein and Ryzhik (1963), and the doubling formula for the Gamma function we obtain

$$p_{\xi}(1, x, y) = \Delta^{-d/2}(\theta) q_x \left( \frac{y - x - \Delta(\theta)b(x)}{\sqrt{\Delta(\theta)}} \right),$$

where

$$q_x(z) = \frac{(d+1)^{d/2}}{2^d \pi^{(d-1)/2} \Gamma \left( \frac{d+1}{2} \right)} e^{-\sqrt{\pi n} |a^{-1/2}(x)z|}.$$

It is easy to check that

$$\int z_i q_x(z) = 0, \quad \int z_i z_j q_x(z) dz = a_{ij}(x).$$

**Example 2.** We put $\Delta(\theta) = \theta^2/2$ and

$$f(r; \theta) = C_d r^{-1} \left( \frac{r}{\theta} \right)^d \exp \left( -\frac{r^2}{\theta} \right),$$

where $C_d = \frac{2^{(d+1)/2}}{(d-2)! \sqrt{\pi}}$ if $d$ is odd, and $C_d = \frac{2}{(|d-2)/2|!}$ if $d$ is even. From (3) and formula 6.631 (4) on page 731 of Gradshtein and Ryzhik (1963) we obtain

$$p_{\xi}(1, x, y) = \Delta^{-d/2}(\theta) \phi_x \left( \frac{y - x - \Delta(\theta)b(x)}{\sqrt{\Delta(\theta)}} \right),$$

where

$$\phi_x(z) = \frac{1}{(2\pi)^d/2 \sqrt{\det a(x)}} \exp \left( -\frac{1}{2} \langle a^{-1}(x)z, z \rangle \right).$$

It is easy to see that the transition density (4) corresponds to the one step transition density in the following Markov chain model

$$X_{(k+1)\Delta(\theta)} = X_{k\Delta(\theta)} + \Delta(\theta) b(X_{k\Delta(\theta)}) + \sqrt{\Delta(\theta)} \xi_{(k+1)\Delta(\theta)},$$

where the conditional density (under $X_{k\Delta(\theta)} = x$) of the innovations $\xi_{(k+1)\Delta(\theta)}$ is equal to $q_x(\cdot)$. If we put $\theta = \theta_n = \sqrt{\frac{2}{n}}$, then $\Delta(\theta_n) = \frac{1}{n}$ and we obtain a sequence of Markov chains defined on an equidistant grid

$$X_{\frac{k+1}{n}} = X_{\frac{k}{n}} + \frac{1}{n} b(X_{\frac{k}{n}}) + \frac{1}{\sqrt{n}} \xi_{\frac{k+1}{n}}, \quad X_0 = x_0. \quad (6)$$

Note that the triplet $(b(x), \sigma(x), f(r; \theta))$, $x \in \mathbb{R}^d$, $r \geq 0$, $\theta \in \mathbb{R}^+$, of the Example 2 corresponds to the classical Euler scheme for the $d$-dimensional SDE

$$dX(t) = b(X_t) dt + \sigma(X_t) dW(t), \quad X(0) = x_0. \quad (7)$$
Let \( p(1, x, y) \) be transition density from 0 to 1 in the model (7). We make the following assumptions

\[ \text{(A2)} \] The function \( a(x) = \sigma \sigma^T(x) \) is uniformly elliptic.

\[ \text{(A3)} \] The functions \( b(x) \) and \( \sigma(x) \) and their derivatives up to the order six are continuous and bounded uniformly in \( x \). The 6-th derivative is globally Lipschitz.

**Theorem 3.** Under assumptions (A2), (A3) we have the following expansion: for any positive integer \( S \) as \( n \to \infty \)

\[
\sup_{x,y \in \mathbb{R}^d} \left( 1 + |y - x|^S \right) \cdot \left| p_\varepsilon(n, x, y) - p(1, x, y) - \frac{1}{2n^2} p \otimes (L_*^2 - L^2) p(1, x, y) \right| = O(n^{-3/2}), \tag{8}
\]

where

\[
L = \frac{1}{2} \sum_{i,j=1}^d a_{ij}(x) \partial^2_{x_i x_j} + \sum_{i=1}^d b_i(x) \partial_{x_i}. \tag{9}
\]

The operator \( L_* \) in (3) is the same operator as in (9) but with coefficients "frozen" at \( x \). Clearly, \( L = L_* \) but, in general, \( L^2 \neq L_*^2 \). The convolution type binary operation \( \otimes \) is defined for functions \( f \) and \( g \) in the following way

\[
(f \otimes g)(t, x, y) = \int_0^t ds \int_{\mathbb{R}^d} f(s, x, z) g(t - s, z, y) dz.
\]

**Proof.** It follows immediately from Theorem 1 of Konakov and Mammen (2009).

4 Proof of Th. 1

4.1 Asymptotic behaviour in case 3)

We have, taking \( B_n = B(T_n) = T_n \):

\[
\sup_{t \in [0, T_n]} \|X_n(t)\|_{\infty} \leq \sum_{k=1}^{n-1} \frac{T_k - T_{k-1}}{T_n} = \frac{T_{n-1}}{T_n} \xrightarrow{n \to \infty} 0 \text{ a.s.}
\]

At the same time,

\[
X_n(1) = \frac{S_n - \varepsilon_n(T_n - T_{n-1})}{T_n} = \varepsilon_n + o(1) \Rightarrow P_{\varepsilon_1}
\]

Therefore the process \( X_n \) converges weakly to the process \( \{Y(t)\} \), \( Y(t) = \varepsilon_1 t, t \in [0, 1] \). This process is in some sense degenerate. Hence this case is not very interesting.

4.2 Asymptotic behaviour in case 2)

Take \( B_n = T_n \) and show that the limit process \( Y \) is not trivial. For simplicity fix \( \beta = 1 \). We have now \( t_{n,k} = \frac{T_k - T_{k-1}}{T_n} = e^{-T_n r_k} = e^{-(\gamma_{k+1} + \cdots + \gamma_n)}, \) and

\[
X_n(t_{n,k}) = \sum_{i=1}^k \varepsilon_i (e^{-\gamma_{i+1} - \cdots - \gamma_n} - e^{-\gamma_i - \cdots - \gamma_n}), \quad k = 1, \ldots, n.
\]
The process $X_n$ is completely defined by 2 independent vectors $(\varepsilon_1, \ldots, \varepsilon_n)$ and $(\gamma_1, \ldots, \gamma_n)$. Hence its distribution will be the same if we replace these vectors by $(\varepsilon_n, \ldots, \varepsilon_1)$ and $(\gamma_n, \ldots, \gamma_1)$. In another words, the process $(X_n(\cdot)) \overset{d}{=} (Y_n(\cdot))$, where $Y_n(\cdot)$ is a broken line with vertices $(\tau_{n,k}, Y_n(\tau_{n,k}))$, $(\tau_{n,k})$, $\tau_{n,1} = 1$, $\tau_{n,k} = e^{-(\gamma_1+\cdots+\gamma_{k-1})}$, $k = 2, \ldots, n$, and

$$Y_n(\tau_{n,k}) = \sum_{i=k}^{n} \varepsilon_i \left( e^{-(\gamma_1+\cdots+\gamma_{i-1})} - e^{-(\gamma_1+\cdots+\gamma_i)} \right) + \varepsilon_n e^{-(\gamma_1+\cdots+\gamma_{n-1})};$$

$Y_n(0) = 0$, and $\gamma_0 := 0$.

Using the notation $\Gamma_k = \gamma_1 + \cdots + \gamma_k$ we get the more compact formula:

$$Y_n(\tau_{n,k}) = \sum_{i=k}^{n-1} \varepsilon_i \left( e^{-\Gamma_{i-1}} - e^{-\Gamma_i} \right) + \varepsilon_n e^{-\Gamma_{n-1}}.$$

Consider now the process $\{Y(t), t \in [0,1]\}$ defined as follows:

$$Y(0) = 0, \quad Y(t_k) = \sum_{i=k}^{\infty} \varepsilon_i \left( e^{-\Gamma_{i-1}} - e^{-\Gamma_i} \right); \quad \text{for } t \in [t_{k+1}, t_k] \ Y(t) \text{ is defined by linear interpolation}.$$

where $t_k = e^{-\Gamma_{k-1}}$, $k = 2, 3, \ldots$, $t_1 = 1$; for $t \in [t_{k+1}, t_k]$ $Y(t)$ is defined by linear interpolation. The paths of $Y$ are continuous broken lines, starting at 0 and having an infinite number of segments in the neighborhood of zero.

The evident estimation

$$\begin{align*}
\sup_{t \in [0,1]} |Y(t) - Y_n(t)| &\leq \sum_{i=n}^{\infty} \varepsilon_i \left( e^{-\Gamma_{i-1}} - e^{-\Gamma_i} \right) + e^{-\Gamma_{n-1}} \\
&\leq \sum_{i=n}^{\infty} \left( e^{-\Gamma_{i-1}} - e^{-\Gamma_i} \right) + e^{-\Gamma_{n-1}} = 2e^{-\Gamma_{n-1}} \to 0 \quad \text{a.s.}
\end{align*}$$

shows that a.s. $Y_n(\cdot) \overset{C[0,1]}{\to} Y(\cdot)$.

Conclusion: In case 2), the process $X_n$ converges weakly to $Y(\cdot)$.

**Remark 3.** In the case where $\beta \neq 1$ it is simply necessary replace $e^{-\Gamma_k}$ by $e^{-\frac{\Gamma_k}{\beta}}$.

**Remark 4.** It seems that the last result could be expanded by considering more general sequences $(\varepsilon_k)$.

Interpretation: $\frac{\varepsilon_k}{|\varepsilon_k|}$ defines direction, $|\varepsilon_k|$ defines the velocity of displacement in this direction on the step $S_k$.

### 4.3 Asymptotic behaviour in case of power growth

In this case $T_k = \Gamma_k^\alpha$, $\alpha > 1/2$, $t_{n,k} = \frac{T_k}{T_n} = \left( \frac{\Gamma_k}{\Gamma_n} \right)^\alpha$, and

$$X_n(t_{n,k}) = \frac{1}{B_n} \sum_{i=1}^{k} \varepsilon_i (\Gamma_i^\alpha - \Gamma_{i-1}^\alpha); \quad \Gamma_0 = 0, \ k = 0, 1, \ldots, n.$$
Let $x \in \mathbb{R}^d$ be such that $|x| = 1$. We will show below that

$$\text{Var} \left( \sum_{i=1}^{n} \langle \varepsilon_i, x \rangle (\Gamma_i^\alpha - \Gamma_{i-1}^\alpha) \right) = E\langle \varepsilon_1, x \rangle^2 \sum_{i=1}^{n} E(\Gamma_i^\alpha - \Gamma_{i-1}^\alpha)^2 \sim C(x)n^{2\alpha - 1}, \ n \to \infty,$$

where $C(x) = \frac{2\alpha}{2\alpha - 1} E\langle \varepsilon_1, x \rangle^2$. Therefore it is natural to take $B_n^2 = n^{2\alpha - 1}$.

We proceed in 5 steps:

**Step 1:** Lemmas

**Step 2:** We compare $X_n(\cdot)$ with $Z_n(\cdot)$ where $Z_n(t_{n,k}) = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i \Gamma_i^\alpha - \Gamma_{i-1}^\alpha$ and show that $\|X_n - Z_n\|_\infty \xrightarrow{p} 0$.

**Step 3:** We compare $Z_n(\cdot)$ with $W_n(\cdot)$ where $W_n(t_{n,k}) = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i(i - 1)^{\alpha - 1}$ and state that $\|Z_n - W_n\|_\infty \xrightarrow{p} 0$.

**Step 4:** We show that process $U_n(\cdot)$,

$$U_n \left( \left( \frac{k}{n} \right)^\alpha \right) = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i(i - 1)^{\alpha - 1},$$

converges weakly to the limiting process

$$Y(t) = \sqrt{2\alpha} \int_0^t s^{\frac{\alpha - 1}{2\alpha}} \, dW(s);$$

here $W(\cdot)$ is a process of Brownian motion, for which the covariance matrix of $W(1)$ coincides with the covariance matrix of $\varepsilon_1$.

**Step 5:** We show that the convergence $W_n \Rightarrow Y$ follows from the convergence $U_n \Rightarrow Y$.

**Finally:** We get the convergence $X_n \Rightarrow Y$.

### 4.3.1 Step 1

**Lemma 1.** Let $\alpha > 0$ and $m \geq 1$. Then $\forall x > 0, h > 0$

$$(x + h)^\alpha - x^\alpha = \sum_{k=1}^{m} a_k h^k x^{\alpha - k} + R(x, h), \quad (12)$$

where

$$a_k = \frac{\alpha(\alpha - 1)\ldots(\alpha - k + 1)}{k!},$$

and

$$|R(x, h)| \leq |a_{m+1}| h^{m+1} \max\{x^{\alpha - (m+1)}, (x + h)^{\alpha - (m+1)}\}. \quad (13)$$

**Proof.** By the formula of Taylor-Lagrange we have (12) with

$$|R(x, y)| \leq \frac{1}{(m + 1)!} h^{m+1} \sup_{x \leq t \leq x + h} |f^{(m+1)}(t)|,$$

where $f(t) = t^\alpha$. As $f^{(m+1)}(t) = \alpha(\alpha - 1)\ldots(\alpha - m)t^{\alpha - (m+1)}$, we get the result.
Lemma 2. For $\alpha \geq 0$ and $k \to \infty$
\[
\left(1 + \frac{\alpha}{k}\right)^k = e^\alpha + O\left(\frac{1}{k}\right).
\] (14)

Proof. It follows from the inequalities:
\[
0 \leq e^\alpha - \left(1 + \frac{\alpha}{k}\right)^k \leq \frac{e^\alpha \alpha^2}{k^2}.
\]
\[\Box\]

Lemma 3. Let $\Gamma$ be the Gamma function. Then as $k \to \infty$
\[
\frac{\Gamma(k + \alpha)}{\Gamma(k)} = k^\alpha + O(k^{\alpha-1}).
\]

Proof. It follows from Lemma 2 and well known asymptotic
\[
\Gamma(t) = t^{t-\frac{1}{2}}e^{-t}\sqrt{2\pi} \left(1 + \frac{1}{12t} + O\left(\frac{1}{t^2}\right)\right), \ t \to \infty.
\]
\[\Box\]

Lemma 4. For any real $\beta$ we have as $k \to \infty$
\[
E(\Gamma_k^\beta) = k^\beta + O(k^{\beta-1}).
\]

Proof. The result follows from the well known fact that
\[
E(\Gamma_k^\beta) = \frac{\Gamma(k + \beta)}{\Gamma(k)}
\]
and Lemma 3. \[\Box\]

Lemma 5. Let $\alpha \geq 0$. The following relations take place as $k \to \infty$:
\[
\Gamma_{k+1}^\alpha - \Gamma_k^\alpha = \alpha \gamma_{k+1} \Gamma_k^{\alpha-1} + \rho_k,
\] (15)
where $|\rho_k| = O(k^{\alpha-2})$ in probability;
\[
E|\Gamma_{k+1}^\alpha - \Gamma_k^\alpha|^2 = 2\alpha^2 k^{2\alpha-2} + O(k^{2\alpha-3});
\] (16)
\[
E|\Gamma_{k+1}^\alpha - \Gamma_k^\alpha - \alpha \gamma_{k+1} \Gamma_k^{\alpha-1}|^2 = O(k^{2\alpha-4}).
\] (17)
We deduce immediately from (16) the following relation.
Corollary 1. We have
\[ \sum_{k=1}^{n-1} E|\Gamma_{k+1}^\alpha - \Gamma_k^\alpha|^2 = \frac{2\alpha^2}{2\alpha - 1} n^{2\alpha - 1} + O(n^{2\alpha - 2}). \]

**Proof of Lemma 5.** We find, applying Lemma 1,
\[ \Gamma_{k+1}^\alpha - \Gamma_k^\alpha = \alpha \gamma_{k+1} \Gamma_k^{\alpha-1} + R(\Gamma_k, \gamma_{k+1}), \]
where
\[ R(\Gamma_k, \gamma_{k+1}) \leq \frac{1}{2} \gamma_{k+1}^2 \max_{\Gamma_k \leq s \leq \Gamma_{k+1}} |\alpha (\alpha - 1)| s^{\alpha - 2} \leq \frac{\alpha (\alpha - 1)}{2} \gamma_{k+1}^2 \max\{\Gamma_{k+1}^{\alpha-2}, \Gamma_k^{\alpha-2}\}. \]

As \( \Gamma_k \sim k \) a.s. when \( k \to \infty \), we get (15). The proofs of (16) and (17) follow directly from (18), (19) and Lemma 4. \( \square \)

### 4.3.2 Step 2

We show that \( \|X_n - Z_n\|_\infty \overset{p}{\to} 0 \), where
\[ Z_n(t_{n,k}) = \frac{\alpha}{B_n} \sum_{i=1}^k \varepsilon_i \gamma_i \Gamma_i^{\alpha-1}. \]

It is clear that
\[ \delta_n := \|X_n - Z_n\|_\infty = \sup_{t \in [0,1]} |X_n(t) - Z_n(t)| = \max_{k \leq n} |X(t_{n,k}) - Z_n(t_{n,k})| = \max_{k \leq n} |r_k|, \]
where
\[ r_k = \frac{1}{B_n} \sum_{i=1}^k \varepsilon_i \left[ \Gamma_i^\alpha - \Gamma_i^{\alpha-1} - \alpha \gamma_i \Gamma_i^{\alpha-1} \right] = \sum_{i=1}^k \varepsilon_i \xi_i, \]
and
\[ \xi_i = \left( \Gamma_i^\alpha - \Gamma_i^{\alpha-1} - \alpha \gamma_i \Gamma_i^{\alpha-1} \right) \frac{1}{B_n}. \]

Let \( \mathcal{M} = \sigma(\xi_1, \xi_2, \ldots, \xi_n) = \sigma(\gamma_1, \gamma_2, \ldots, \gamma_n) \). Under condition \( \mathcal{M} \) the sequence \( (r_k) \) is the sequence of sums of independent random variables with mean zero. By Kolmogorov’s inequality
\[ \mathbb{P}\{\max_{k \leq n} |r_k| \geq t\} \leq \mathbb{E}\{\max_{k \leq n} |r_k| \geq t \mid \mathcal{M}\} \leq \mathbb{E} \left( \frac{1}{t^2} \sum_{j=1}^n \xi_j^2 \right) = \frac{1}{t^2} \sum_{j=1}^n \mathbb{E}\xi_j^2. \]

By Lemma 5 \( \mathbb{E}\xi_j^2 = O(j^{-3}) \). Therefore,
\[ \sum_{j=1}^n \mathbb{E}\xi_j^2 = O(n^{-2}). \]

Finally we get from (20): \( \forall t > 0 \)
\[ \mathbb{P}\{\delta_n \geq t\} \overset{n \to \infty}{\to} 0, \]
which gives the convergence \( \|X_n - Z_n\|_\infty \overset{p}{\to} 0 \).
4.3.3 Step 3

We show now that \( \|Z_n - W_n\|_\infty \xrightarrow{p} 0 \); where \( W_n(t_{n,k}) = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i (i - 1)^{a-1} \).

We have
\[
\Delta_n = \sup_{t \in [0,1]} |Z_n(t) - W_n(t)| = \max_{k \leq n} |Z_n(t_{n,k}) - W_n(t_{n,k})| = \max_{k \leq n} \{ |\beta_k| \},
\]
where \( \beta_k = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i (\Gamma_i^{a-1} - (i - 1)^{a-1}) \).

Similar to the previous case \( (\beta_k) \) under condition \( \mathcal{M} \) is the sequence of sums of independent random variables with mean zero. Therefore
\[
\mathbb{P}\{ \max_{k \leq n} \{ |\beta_k| \} \geq t \} = \mathbb{E}\left( \mathbb{P}\{ \max_{k \leq n} \{ |\beta_k| \} \geq t \} \right) \leq \frac{1}{t^2} \sum_{j=1}^{n} \mathbb{E}\eta_j^2,
\]
where \( \eta_j = \frac{\alpha}{B_n} \gamma_j (\Gamma_j^{a-1} - (j - 1)^{a-1}) \).

**Estimation of \( \mathbb{E}\eta_j^2 \).**

By independence of \( \gamma_j \) and \( \Gamma_j \-
\[
\mathbb{E}\eta_j^2 = \frac{2\alpha^2}{B_n^2} \mathbb{E}(\Gamma_j^{a-1} - (j - 1)^{a-1})^2
\]
Let us change \( j - 1 \) to \( k \)
\[
\mathbb{E}(\Gamma_k^{a-1} - k^{a-1})^2 = \mathbb{E}(\Gamma_k^{2a-2}) + k^{2a-2} - 2k^{a-1} \mathbb{E}(\Gamma_k^{a-1}) = \\
= \frac{\Gamma(k + 2\alpha - 2)}{\Gamma(k)} + k^{2a-2} - 2k^{a-1} \frac{\Gamma(k + \alpha - 1)}{\Gamma(k)} \quad \text{(by Lemma 3)} = \\
= [k^{2a-2} + O(k^{2a-3}) + k^{2a-2} - 2k^{2a-2}] = O(k^{2a-3}).
\]

Hence
\[
\mathbb{E}\eta_j^2 \leq C\frac{j^{2a-3}}{n^{2a-1}}
\]
and
\[
\sum_{j=1}^{n} \mathbb{E}\eta_j^2 \leq C\frac{1}{n^2}
\]

We have finally \( \mathbb{P}\{ \max_{k \leq n} \{ |\beta_k| \} \geq t \} \to 0, \ n \to \infty \), which gives the convergence \( \|W_n - Z_n\| \xrightarrow{p} 0 \).

4.3.4 Step 4

Let \( U_n \) be the process defined at the points \( \frac{k}{n} \) by
\[
U_n\left(\left(\frac{k}{n}\right)^{\alpha}\right) = \frac{\alpha}{B_n} \sum_{i=1}^{k} \varepsilon_i \gamma_i (i - 1)^{a-1}, \quad k = 1, 2, \ldots, n,
\]
and by linear interpolation on the intervals \([\frac{k}{n}, \frac{k+1}{n}]\), \( k = 0, \ldots, n - 1 \). We now state weak convergence of the processes \( U_n \) to the process \( Y \),
\[
Y(t) = \sqrt{2\alpha} \int_{0}^{t} s^{\frac{a-1}{2a}} \, dW(s),
\]

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W is a Brownian motion, for which the covariance matrix of \( W(1) \) coincides with the covariance matrix of \( \varepsilon_1 \).

The proof is standard because \( U_n(\cdot) \) represents a (more or less) usual broken line constructed by the consecutive sums of independent (non-identically distributed) random variables. One could apply Prokhorov's theorem (see Gikhman and Skorohod (1996), ch.IX, sec. 3, Th.1).

Only one thing must be checked: that for any \( 0 < s < t \leq 1 \), and for any \( x \in \mathbb{R}^d, |x| = 1 \), we have the convergence \( (U_n(t) - U_n(s), x) \implies (Y(t) - Y(s), x) \).

It is clear that

\[
[U_n(t) - U_n(s)] = \left[ U_n \left( \frac{k}{n} \right)^\alpha - U_n \left( \frac{l}{n} \right)^\alpha \right] \to 0,
\]

if \( \left( \frac{k}{n} \right)^\alpha \to t, \left( \frac{l}{n} \right)^\alpha \to s \).

Let \( l < k \). As

\[
\left\langle U_n \left( \frac{k}{n} \right)^\alpha - U_n \left( \frac{l}{n} \right)^\alpha, x \right\rangle = \frac{\alpha}{B_n} \sum_{i=l+1}^k \langle \varepsilon_i, x \rangle \gamma_i(i-1)^{\alpha-1},
\]

by the theorem of Lindeberg-Feller it is sufficient to state the convergence of variances.

We have

\[
\text{Var}\left\langle U_n \left( \frac{k}{n} \right)^\alpha - U_n \left( \frac{l}{n} \right)^\alpha, x \right\rangle = \frac{2\alpha^2}{n^{2\alpha-1}}E\langle \varepsilon_1, x \rangle^2 \sum_{i=l+1}^k (i-1)^{2\alpha-2} \to \frac{2\alpha^2}{2\alpha-1}E\langle \varepsilon_1, x \rangle^2\left[t^{2\alpha-1} - s^{2\alpha-1}\right],
\]

and

\[
\text{Var}(Y(t) - Y(s), x) = 2\alpha E\langle \varepsilon_1, x \rangle^2 \int_s^t u^{\frac{a-1}{\alpha}} du = \frac{2\alpha^2}{2\alpha-1}E\langle \varepsilon_1, x \rangle^2\left[t^{\frac{2\alpha-1}{\alpha}} - s^{\frac{2\alpha-1}{\alpha}}\right],
\]

which are the same.

### 4.3.5 Step 5: Convergence \( X_n \Rightarrow Y \).

Due to the steps 2 and 3 it is sufficient to show that \( W_n \Rightarrow Y \).

Let \( f_n : [0, 1] \to [0, 1] \), be a piecewise linear continuous function such that \( f_n(t_{n,k}) = \left( \frac{k}{n} \right)^\alpha \); \( t_{n,k} = \left( \frac{k}{n} \right)^\alpha \); \( k = 0, 1, \ldots, n \).

By definition of \( W_n \) and \( U_n \) we have

\[ W_n(t) = U_n(f_n(t)), \ t \in [0, 1]. \]

By the corollary to Lemma 6 (see below) the function \( f_n \) converges in probability uniformly to \( f, f(t) = t \), and by previous step \( U_n \Rightarrow Y \).

It means that we can apply Lemma 7 which gives the necessary convergence.

### Lemma 6. Let

\[ M_n = \max_{k \leq n} \left\{ \frac{\Gamma_k}{\Gamma_n} - \frac{k}{n} \right\}. \]

Then \( M_n \xrightarrow{p} 0, \ n \to \infty. \)
Proof of Lemma 6. We have
\[
\mathbb{P}\{M_n > \varepsilon\} = \mathbb{E}\left\{\mathbb{P}\left\{\max_{k \leq n} \left| \frac{\Gamma_k}{\Gamma_n} - \frac{k}{n} \right| > \varepsilon \mid \Gamma_n\right\}\right\} = \\
\int_0^\infty \mathbb{P}\left\{\max_{k \leq n} \left| \frac{\Gamma_k}{\Gamma_n} - \frac{k}{n} \right| > \varepsilon \mid \Gamma_n = t\right\} \mathcal{P}(dt) = \\
\int_0^\infty \mathbb{P}\left\{\max_{k \leq n} \left| \xi_{n,k} - \frac{k}{n} \right| > \varepsilon \right\} \mathcal{P}(dt) = \mathbb{P}\left\{\max_{k \leq n} \left| \xi_{n,k} - \frac{k}{n} \right| > \varepsilon \right\},
\]
where \((\xi_{n,k})_{k=1,\ldots,n}\) are the order statistics from \([0,1]\)-uniform distribution.

Let \(\delta_n := \max_{k \leq n} |\xi_{n,k} - \frac{k}{n}|\). Evidently, \(\delta_n \leq \sup_{[0,1]} |F_n^*(x) - x|\), where \(F_n^*\) is a uniform empirical distribution function. By Glivenko-Cantelli theorem, \(\sup_{[0,1]} |F_n^*(x) - x| \to 0\) a.s., which gives the convergence \(M_n \to 0\) in probability. \(\square\)

Corollary 2. \(M^{(1)}_n = \max_{k \leq n} \left| \left(\frac{\Gamma_k}{\Gamma_n}\right)^\alpha - \left(\frac{k}{n}\right)^\alpha \right| \overset{p}{\to} 0\), \(n \to \infty\).

The proof follows directly from Lemma 6 due to the uniform continuity of the function \(h(x) = x^\alpha\), \(x \in [0,1]\).

Lemma 7. Let \(\{U_n\}\) be a sequence of continuous processes on \([0,1]\] weakly convergent to some limit process \(U\). Let \(\{f_n\}\) be a sequence of random continuous bijections \([0,1]\) on \([0,1]\] which in probability uniformly converges to the identity function \(f(t) \equiv t\). Then the process \(W_n\), \(W_n(t) = U_n(f_n(t))\), \(t \in [0,1]\], will converge weakly to \(U\).

Proof of Lemma 7. By theorem 4.4 from Billingsley (1968) we have the weak convergence in \(\mathbb{M} := C[0,1] \times C[0,1]\)
\[(U_n, f_n) \Longrightarrow (U, f).\]

By Skorohod representation theorem we can find a random elements \((\tilde{U}_n, \tilde{f}_n)\) and \((\tilde{U}, \tilde{f})\) of \(\mathbb{M}\) (defined probably on a new probability space) such that
\[(U_n, f_n) \overset{D}{=} (\tilde{U}_n, \tilde{f}_n), \quad (U, f) \overset{D}{=} (\tilde{U}, \tilde{f}),\]
and \((\tilde{U}_n, \tilde{f}_n) \to (\tilde{U}, \tilde{f})\) a.s. in \(\mathbb{M}\).

As the last convergence implies evidently the a.s. uniform convergence of \(\tilde{U}_n(\tilde{f}_n(t))\) to \(\tilde{U}(\tilde{f}(t))\), we get the convergence in distribution of \(U(f_n(\cdot))\) to \(U(f(\cdot)) = U(\cdot)\). \(\square\)

5 Proof of Th. 2

Proof of Th. 2. We need some facts from Stroock and Varadhan (1979). Consider \((\Omega, \mathcal{M})\), where \(\Omega = C([0,\infty); R^d)\) be the space of continuous trajectories from \([0, \infty)\) into \(R^d\). Given \(t \geq 0\) and \(\omega \in \Omega\) let \(x(t, \omega)\) denote the position of \(\omega\) in \(R^d\) at time \(t\). If we put
\[D(\omega, \omega') = \sum_{n=1}^\infty \frac{1}{2^n} \frac{\sup_{0 \leq t \leq n} |x(t, \omega) - x(t, \omega')|}{1 + \sup_{0 \leq t \leq n} |x(t, \omega) - x(t, \omega')|},\]

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then it is well known that $D$ is a metric on $\Omega$ and $(\Omega, D)$ is a Polish space. The convergence induced by $D$ is uniform convergence on bounded $t$ - intervals. For simplicity, we will omit $\omega$ in the future and we will be assuming that all our processes are homogeneous in time. Analogous results for time-inhomogeneous processes may be obtained by simply considering the time-space processes.

We will use $\mathcal{M}$ to denote the Borel $\sigma$ - field of subsets of $(\Omega, D)$, \( \mathcal{M} = \sigma[x(t) : t \geq 0] \). We also will consider an increasing family of $\sigma$-algebras $\mathcal{M}_t = \sigma[x(s) : 0 \leq s \leq t]$. Classical approach to the construction of diffusion processes corresponding to given coefficients $a$ and $b$ involves a transition probability function $P(s, x, t, \cdot)$ which allows to construct for each $x \in R^d$, a probability measure $P_x$ on $\Omega = C([0, \infty); R^d)$ with the properties that

$$P_x(x(0) = x) = 1$$

and

$$P_x(x(t_2) \in \Gamma | \mathcal{M}_{t_1}) = P(t_1, x(t_1); t_2, \Gamma) \text{ a.s.} P_x$$

for all $0 \leq t_1 < t_2$ and $\Gamma \in B_{R^d}$ (the Borel $\sigma$ - algebra in $R^d$). It appears that this measure is a martingale measure for a special martingale related with the second order differential operator

$$L = \frac{1}{2} \sum_{i,j=1}^{d} a^{ij}(\cdot) \frac{\partial^2}{\partial x_i \partial x_j} + \sum_{i=1}^{d} b^i(\cdot) \frac{\partial}{\partial x_i},$$

namely, for all $f \in C_0^\infty(R^d)$

$$P_x(x(0) = x) = 1,$$

$$(f(x(t)) - \int_{0}^{t} Lf(x(u))du, \mathcal{M}_t, P_x)$$

is a martingale. We will say that the martingale problem for $a$ and $b$ is well-posed if, for each $x$ there is exactly one solution to that martingale problem starting from $x$. We will be working with the following set up. For each $h > 0$ let $\Pi_h(x, \cdot)$ be a transition function on $R^d$. Given $x \in R^d$, let $P^h_x$ be the probability measure on $\Omega$ characterized by the properties that

$$(i) \quad P^h_x(x(0) = x) = 1,$$

$$(ii) \quad P^h_x\left\{ x(t) = \frac{(k+1)h - t}{h} x(kh) + \frac{t - kh}{h} x((k + 1)h), \; kh \leq t < (k + 1)h \right\} = 1 \quad (24)$$

for all $k \geq 0$,

$$(iii) \quad P^h_x(x((k + 1)h) \in \Gamma | \mathcal{M}(kh)) = \Pi_h(x(kh), \Gamma), \text{ } P^h_x - \text{ a.s.} \quad (25)$$

Define

$$a^{ij}_h(x) = \frac{1}{h} \int_{|y-x| \leq 1} (y_i - x_i)(y_j - x_j) \Pi_h(x, dy),$$

$$b^i_h(x) = \frac{1}{h} \int_{|y-x| \leq 1} (y_i - x_i) \Pi_h(x, dy),$$

and

$$\Delta^i_h(x) = \frac{1}{h} \Pi_h(x, R^d \setminus B(x, \varepsilon)), \quad (28)$$
where $B(x, \varepsilon)$ is the open ball with center $x$ and radius $\varepsilon$. What we are going to assume is that for all $R > 0$

$$\lim_{k \to 0} \sup_{|x| \leq R} \|a_k(x) - a(x)\| = 0, \quad (29)$$

$$\lim_{k \to 0} \sup_{|x| \leq R} |b_k(x) - b(x)| = 0, \quad (30)$$

$$\sup_{h > 0} \sup_{x \in \mathbb{R}^d} (|a_k(x)| + |b_k(x)|) < \infty, \quad (31)$$

$$\lim_{k \to 0} \sup_{x \in \mathbb{R}^d} \Delta^k_h(x) = 0. \quad (32)$$

**Theorem A.** (Strook and Varadhan (1979), page 272, Theorem 11.2.3). Assume that in addition to (29)-(32) the coefficients $a$ and $b$ are continuous and have the property that for each $x \in \mathbb{R}^d$ the martingale problem for $a$ and $b$ has exactly one solution $P_x$ starting from $x$ (that is well posed). Then $P_x^h$ converges weakly to $P_x$ uniformly in $x$ on compact subsets of $\mathbb{R}^d$.

Sufficient conditions for the well-posedness is given by the following theorem.

Let $S_d$ be the set of symmetric non-negative definite $d \times d$ real matrices.

**Theorem B.** (Strook and Varadhan (1979), page 152, Theorem 6.3.4). Let $a : \mathbb{R}^d \to S_d$ and $b : \mathbb{R}^d \to \mathbb{R}^d$ be bounded measurable functions and suppose that $\sigma : \mathbb{R}^d \to \mathbb{R}^d \times \mathbb{R}^d$ is a bounded measurable function such that $a = \sigma\sigma^*$. Assume that there is an $A$ such that

$$\|\sigma(x) - \sigma(y)\| + |b(x) - b(y)| \leq A|x - y| \quad (33)$$

for all $x, y \in \mathbb{R}^d$. Then the martingale problem for $a$ and $b$ is well-posed and the corresponding family of solutions $\{P_x : x \in \mathbb{R}^d\}$ is Feller continuous (that is $P_{x_n} \to P_x$ weakly if $x_n \to x$).

Note that (33) and uniform ellipticity of $\sigma$ imply the existence of the transition density $p(s, x; t, y)$ (Strook and Varadhan (1979), Theorem 3.2.1, page 71).

Consider the model

$$X((k + 1)h) = X(kh) + h\bar{b}(X(kh)) + \sqrt{h}\xi(X(kh)), \quad (34)$$

where $\{\xi_k\}$ are i.i.d. random vectors uniformly distributed on the unit sphere $S^{d-1}$, and $\{\rho_k\}$ are i.i.d. random variables having a density, $\rho_k \geq 0$, $E\rho_k^2 = d$. Let us check the conditions (29)-(32). It is easy to see that

$$\Pi_h(x, dy) = p^\xi_h(y)dy, \quad \text{where} \quad p^\xi_h(y) = h^{-d/2}f_\xi\left(\frac{y - x - h\bar{b}(x)}{\sqrt{h}}\right). \quad (35)$$

Here $f_\xi$ denotes the density of the random vector $\xi$. Let us check (32). Note that $E\xi = 0$ and the covariance matrix of the vector $\xi$ is equal to

$$\text{Cov}(\xi, \xi^T) = E(\rho_k^2\sigma(x)\varepsilon_\varepsilon\varepsilon_k^T \sigma^T(x)) = a(x). \quad (36)$$

We have

$$h\Delta^k_h(x) = \Pi_h(x, \mathbb{R}^d \setminus B(x, \varepsilon)) = \int_{\mathbb{R}^d \setminus B(x, \varepsilon)} p^\xi_h(y)dy = \int_{\mathbb{R}^d \setminus B(x, \varepsilon)} f_\xi(v)dv = P\left\{\xi \in B\left(0, \frac{\varepsilon}{\sqrt{h}}\right) - \sqrt{h}\bar{b}(x)\right\} \leq 16$$
\[ P \left\{ |\xi|^2 \geq \frac{\varepsilon^2}{4h} \right\} = o(h). \tag{37} \]

The last equality is a consequence of the Markov inequality. The equality (36), uniform ellipticity of \( a(x) \) and (37) implies (32). To prove (29) note that by (33)

\[
\begin{align*}
a^i_j(x) &= \frac{1}{h} \int_{|y-x| \leq 1} (y_i - x_i)(y_j - x_j)p_h^e(y)dy = \\
&= \int_{|v+\sqrt{h}b^i(x)| \leq \frac{1}{\sqrt{h}}} (v_i + \sqrt{h}b^i(x))(v_j + \sqrt{h}b^j(x))f_\xi(v)dv = \\
&= \int_{|v+\sqrt{h}b(x)| \leq \frac{1}{\sqrt{h}}} v_i v_j f_\xi(v)dv + o(\sqrt{h}) = a(x) + o(1). \tag{38}
\end{align*}
\]

To check (30) note that

\[
\begin{align*}
b^i_h(x) &= \frac{1}{h} \int_{|y-x| \leq 1} (y_i - x_i)p_h^e(y)dy = \\
&= \frac{1}{\sqrt{h}} \int_{|v+\sqrt{h}b(x)| \leq \frac{1}{\sqrt{h}}} (v_i + \sqrt{h}b^i(x))f_\xi(v)dv = \\
&= b^i(x) \int_{|v+\sqrt{h}b(x)| \leq \frac{1}{\sqrt{h}}} f_\xi(v)dv - \frac{1}{\sqrt{h}} \int_{|v+\sqrt{h}b(x)| > \frac{1}{\sqrt{h}}} v_i f_\xi(v)dv. \tag{39}
\end{align*}
\]

To estimate the second integral in (39) we apply the Cauchy - Schwarz inequality

\[
\frac{1}{\sqrt{h}} \int_{|v+\sqrt{h}b(x)| > \frac{1}{\sqrt{h}}} |v| f_\xi(v)dv \leq \frac{1}{\sqrt{h}} \left( \int |v|^2 f_\xi(v)dv \right)^{1/2} \left( P(\{ |\xi|^2 \geq \frac{1}{4h} \}) \right)^{1/2} = o(1), \tag{40}
\]

and (39), (40) imply (30). Finally, (31) follows from our calculations and assumptions of Theorem B. Weak convergence \( P_x^h \) to \( P_x \) follows now from Theorems A and B cited above.

\[ \square \]

6 References

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