

A modified Poincaré inequality and its application to First Passage Percolation

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Abstract: We extend a functional inequality for the Gaussian measure on \mathbb{R}^n to the one on $\mathbb{R}^{\mathbb{N}}$. This inequality improves on the classical Poincaré inequality for Gaussian measures. As an application, we prove that First Passage Percolation has sublinear variance when the edge times distribution belongs to a wide class of continuous distributions, including the exponential one. This extends a result by Benjamini, Kalai and Schramm [3], valid for positive Bernoulli edge times.

Résumé: On étend une inégalité fonctionnelle pour la mesure gaussienne sur \mathbb{R}^n à celle sur $\mathbb{R}^{\mathbb{N}}$. Cette inégalité peut être vue comme une amélioration de l'inégalité de Poincaré classique pour la mesure gaussienne. Comme application, nous montrons que la percolation de premier passage a une variance sous-linéaire pour une large classe de distributions des temps d'arêtes, incluant les lois exponentielles. Ceci étend un résultat de Benjamini, Kalai et Schramm [3], valable pour des temps d'arêtes strictement positifs suivant une loi de Bernoulli.

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1. Introduction

First Passage Percolation was introduced by Hammersley and Welsh [5] to model the flow of a fluid in a randomly porous material (see [6] for a recent account on the subject). We will consider the following model of First Passage Percolation in \mathbb{Z}^d , where $d \geq 2$ is an integer. Let $E = E(\mathbb{Z}^d)$ denote the set of edges in \mathbb{Z}^d . The passage time of the fluid through the edge e is denoted by x_e and is supposed to be nonnegative. Randomness of the porosity is given by a product probability measure on \mathbb{R}_+^E . Thus, \mathbb{R}_+^E is equipped with the probability measure $\mu = \nu^{\otimes E}$, where ν is a probability measure on \mathbb{R}_+ according to which each passage time is distributed, independently from the others. If u, v are two vertices of \mathbb{Z}^d , the notation $\alpha : \{u, v\}$ means that α is a path with end points u and v . When $x \in \mathbb{R}_+^E$, $d_x(u, v)$ denotes the first passage time, or equivalently the distance from u to v in the metric induced by x ,

$$d_x(u, v) = \inf_{\alpha: \{u, v\}} \sum_{e \in \alpha} x_e .$$

The study of $d_x(0, nu)$ when n is an integer which goes to infinity is of central importance. Kingman's subadditive ergodic theorem implies the existence, for each fixed u , of a "time constant" $t(u)$ such that:

$$\frac{d_x(0, nu)}{n} \xrightarrow[n \rightarrow +\infty]{\nu\text{-a.s.}} t(u) .$$

It is known (see Kesten [9], p.127 and 129) that if $\nu(\{0\})$ is strictly smaller than the critical probability for Bernoulli bond percolation on \mathbb{Z}^d , then $t(u)$ is positive for every u which is not the origin. Thus, in this case, one can say that the random variable $d_x(0, nu)$ is located around $nt(u)$, which is of order $O(|nu|)$, where we denote by $|\cdot|$ the L^1 -norm of vertices in \mathbb{Z}^d . In this paper, we are interested in the fluctuations of this quantity. Precisely, we define, for any vertex v ,

$$\forall x \in \mathbb{R}_+^E, f_v(x) = d_x(0, v) .$$

It is widely believed that the fluctuations of f_v are of order $|v|^{1/3}$. Apart from some predictions made by physicists, this faith relies on recent results for related growth models [2, 7, 8]. Until recently, the best results rigorously obtained for the fluctuations of f_v were some moderate deviation estimates of order $O(|v|^{1/2})$ (see [10, 15]). In 2003, Benjamini, Kalai and Schramm [3] proved that for Bernoulli edge times, the variance of f_v is of order $O(|v|/\log|v|)$, and therefore, the fluctuations are of order $O(|v|^{1/2}/(\log|v|)^{1/2})$.

The main result of the present paper is Theorem 4.4, where we extend the result of Benjamini et al. to a large class of probability measures, absolutely

continuous with respect to the Lebesgue measure. This class includes all the gamma and beta distributions (and therefore, the exponential distribution), but also any probability measure whose density is bounded away from 0 on its support, and notably the uniform distribution on $[a, b]$, with $0 \leq a < b$. The result of Benjamini et al. relies mainly on an inequality on the discrete cube due to Talagrand [14]. In their article, they suggested to extend their work to other edge-times distributions by using the tools developed by Ledoux [11], and they pointed out a Gaussian version of Talagrand's inequality found by Bobkov and Houdré [4]. Our strategy to extend Benjamini, Kalai and Schramm's result is thus to start from the Gaussian version of Talagrand's inequality, and then to adapt their argument to a continuous distribution ν via this inequality and a standard change of variable from the Gaussian distribution to ν . Following Ledoux [12], p. 41, we shall give another proof of the Gaussian inequality that we need in Proposition 2.2. Since it follows (almost) entirely Ledoux's argument, and since the result is implied, up to a multiplicative constant, by Bobkov and Houdré [4], we don't claim that this proposition is original at all. Nevertheless, we choosed to write it down, because in the precise form which we state it in, it really improves upon the classical Poincaré inequality for the Gaussian measure. By this, we mean that it implies the classical Poincaré inequality with the true, optimal constant.

This article is organized as follows. In Section 2, we derive the Gaussian analog to Talagrand's inequality [14], and extend it to a countable product of Gaussian measures. Section 3 is devoted to the obtention of similar inequalities for other continuous measures by a simple mean of change of variable. This allows us to adapt the argument of Benjamini et al. to some continuous settings in section 4.

Notation Given a probability space $(\Omega, \mathcal{A}, \mu)$ and a real valued measurable function f defined on Ω we let

$$\|f\|_{p,\mu} = \left(\int \|f\|^p d\mu \right)^{1/p} \in [0, \infty]$$

and we let $L^p(\mu)$ denote the set of f such that $\|f\|_{p,\mu} < \infty$. The *mean* of $f \in L^1(\mu)$ is denoted

$$\mathbb{E}_\mu(f) = \int f d\mu$$

and the *variance* of $f \in L^2(\mu)$ is

$$\mathbf{Var}_\mu(f) = \|f - \mathbb{E}_\mu(f)\|_{2,\mu}^2.$$

When the choice of μ is unambiguous we may write $\|f\|_p$ (respectively L^p , $\mathbb{E}(f)$, $\mathbf{Var}(f)$) for $\|f\|_{p,\mu}$ (respectively $L^p(\mu)$, $\mathbb{E}_\mu(f)$, $\mathbf{Var}_\mu(f)$).

2. A modified Poincaré inequality for Gaussian measures on $\mathbb{R}^{\mathbb{N}}$

The purpose of this section is to rewrite Ledoux's argument (see [12]) in deriving a Poincaré type inequality for a product measure having the form $\lambda \otimes \gamma^{\mathbb{N}}$ where λ is uniform on $\{0, 1\}^S$ with S finite and $\gamma^{\mathbb{N}}$ is the standard Gaussian measure on $\mathbb{R}^{\mathbb{N}}$. Such an inequality, which will prove to be very useful in the sequel, can be seen as a ‘‘Gaussian’’ version of an inequality proved by Talagrand in [14] for Bernoulli measures. Please note that it is implied, up to a multiplicative constant, by an inequality from Bobkov and Houdré [4]. We shall first prove it for the standard Gaussian on \mathbb{R}^n and then extend it to the case of $\mathbb{R}^{\mathbb{N}}$.

The case of \mathbb{R}^n

We let

$$\gamma(dy) = \frac{1}{\sqrt{2\pi}} e^{-y^2/2} dy$$

denote the standard Gaussian measure on \mathbb{R} and

$$\beta = \frac{1}{2}(\delta_0 + \delta_1)$$

the Bernoulli measure on $\{0, 1\}$. If $n \in \mathbb{N}$, and S is a finite set we let $\gamma^n = \gamma^{\otimes n}$ and $\lambda = \beta^{\otimes S}$ denote the associated product measures on \mathbb{R}^n and $\{0, 1\}^S$.

Given a measurable mapping $f : \{0, 1\}^S \times \mathbb{R}^n \mapsto \mathbb{R}$ and $q \in S$ we set

$$\nabla_q f(x, y) = f(x, y) - \int_{\{0,1\}} f(x, y) d\lambda(x_q).$$

A map f is said to be *weakly differentiable* (in the y variables) provided there exist locally integrable (in the y variables) functions denoted $(\frac{\partial f}{\partial y_j})$, $j = 1, \dots, n$, such that

$$\int_{\mathbb{R}^n} \frac{\partial f}{\partial y_j}(x, y) g(y) dy = - \int_{\mathbb{R}^n} \frac{\partial g}{\partial y_j}(y) f(x, y) dy$$

for every smooth function $g : \mathbb{R}^n \mapsto \mathbb{R}$ with compact support.

The *weighted Sobolev* space $H_1^2(\lambda \otimes \gamma^n)$ is defined to be the space of weakly differentiable functions f on $\{0, 1\}^S \times \mathbb{R}^n$ such that

$$\|f\|_{H_1^2}^2 = \|f\|_2^2 + \sum_{j=1}^n \left\| \frac{\partial f}{\partial y_j} \right\|_2^2 < \infty.$$

The following result is essentially a Gaussian version of Talagrand's Theorem 1.5 in [14].

Proposition 2.1 For every function f in $H_1^2(\lambda \otimes \gamma^n)$,

$$\text{Var}(f) \leq \sum_{q \in S} \|\nabla_q f\|_2^2 + \sum_{j=1}^n \left\| \frac{\partial f}{\partial y_j} \right\|_2^2 \phi \left(\frac{\left\| \frac{\partial f}{\partial y_j} \right\|_1}{\left\| \frac{\partial f}{\partial y_j} \right\|_2} \right),$$

where $\phi : [0, 1] \mapsto [0, 1]$ is defined as

$$\phi(u) = 2 \int_0^1 \frac{u^{2t}}{(1+t)^2} dt.$$

Remark 1 Function ϕ is continuous nondecreasing with $\phi(0) = 0$ (in fact $\phi(u) \sim_{u=0} -1/(\log(u))$) and $\phi(1) = 1$. In particular, Proposition 2.1 implies the standard Poincaré inequality for γ^n . That is

$$\text{Var}(f) \leq \sum_{i=1}^n \left\| \frac{\partial f}{\partial y_i} \right\|_2^2,$$

for all $f \in H_2^1(\gamma^n)$.

Proof: For every $f \in L^2$

$$\text{Var}_{\lambda \otimes \gamma^n}(f) = \mathbb{E}_{\gamma^n}(\text{Var}_\lambda(f)) + \text{Var}_{\gamma^n}(\mathbb{E}_\lambda(f)). \quad (1)$$

Furthermore, variance satisfies the following well known tensorisation property (see for instance Ledoux [11], Proposition 5.6 p. 98):

$$\forall g \in L^2(\lambda), \text{Var}_\lambda(g) \leq \sum_{q \in S} \|\nabla_q g\|_{2,\lambda}^2.$$

Hence

$$\mathbb{E}_{\gamma^n}(\text{Var}_\lambda(f)) \leq \sum_{q \in S} \|\nabla_q f\|_{2,\lambda \otimes \gamma^n}^2. \quad (2)$$

Following Benjamini et al. [3], we shall use an hypercontractivity property of a semi-group $(P_t)_{t \geq 0}$ with invariant measure γ^n in order to bound the last term in the right hand side of inequality (1). A most natural choice is to take for $(P_t)_{t \geq 0}$ the semi-group induced by n independent copies of an Ornstein-Uehlenbeck process.

Let $C_0^\infty(\mathbb{R}^n)$ be the space of smooth real valued functions on \mathbb{R}^n that go to zero at infinity and let $\mathcal{A} \subset C_0^\infty(\mathbb{R}^n)$ be the set of functions whose partial derivatives of all order $\frac{\partial^k f}{\partial x_{i_1} \dots \partial x_{i_k}}$ lie in $C_0^\infty(\mathbb{R}^n)$. For any $g \in \mathcal{A}$ and $t \geq 0$,

$$\frac{dP_t(g)}{dt} = \mathbf{L}P_t(g),$$

where

$$\mathbf{L}g(x) = \Delta g(x) - x \cdot \nabla g(x) .$$

If f is defined on $\{0, 1\}^S \times \mathbb{R}^n$ and $x \in \{0, 1\}^S$, we let f_x denote the function on \mathbb{R}^n defined as

$$\forall y \in \mathbb{R}^n, f_x(y) = f(x, y) .$$

Let \mathcal{A}_S be the set of functions f on $\{0, 1\}^S \times \mathbb{R}^n$ such that, for every $x \in \{0, 1\}^S$, the function f_x belongs to \mathcal{A} .

Without loss of generality we may (hence do) assume that $f \in \mathcal{A}_S$. For, by standard approximation results \mathcal{A}_S is a dense subset of $H_1^2(\lambda \otimes \gamma^n)$ equipped with the norm $\|f\|_{H_1^2}$ and the inequality to be proved is a closed condition in $H_1^2(\lambda \otimes \gamma^n)$.

Given such an f , let $F = \mathbb{E}_\lambda(f)$, which belongs to \mathcal{A} . Then

$$\begin{aligned} P_0 F &= F , \\ P_t F &\xrightarrow[t \rightarrow \infty]{L^2} \int F(y) d\gamma^n(y) . \end{aligned}$$

Therefore,

$$\begin{aligned} \text{Var}_{\gamma^n}(F) &= \int_0^{+\infty} -\frac{d}{dt} \mathbb{E}_{\gamma^n}((P_t F)^2) dt , \\ &= 2 \int_0^{+\infty} \mathbb{E}_{\gamma^n}(-P_t F \mathbf{L} P_t F) dt . \end{aligned}$$

Notice that for every function g in \mathcal{A} , it follows from integration by parts that:

$$\mathbb{E}_{\gamma^n}(-g \mathbf{L}g) = \sum_{i=1}^n \left\| \frac{\partial g}{\partial y_i} \right\|_2^2 .$$

Consequently,

$$\text{Var}_{\gamma^n}(F) = 2 \sum_{i=1}^n \int_0^{+\infty} \mathbb{E}_{\gamma^n} \left[\left(\frac{\partial}{\partial y_i} P_t F \right)^2 \right] dt . \quad (3)$$

Recall that, by Mehler's formula

$$P_t g(y) = \int g \left(ye^{-t} + z\sqrt{1 - e^{-2t}} \right) d\gamma_n(z) .$$

This expression implies that, for any $i = 1, \dots, n$ and $g \in \mathcal{A}$,

$$\frac{\partial}{\partial y_i} P_t g = e^{-t} P_t \left(\frac{\partial g}{\partial y_i} \right) .$$

Thus,

$$\mathbb{E}_{\gamma^n} \left[\left(\frac{\partial}{\partial y_i} P_t F \right)^2 \right] = e^{-2t} \mathbb{E}_{\gamma^n} \left[\left(P_t \frac{\partial F}{\partial y_i} \right)^2 \right].$$

Nelson's Theorem asserts that $(P_t)_{t \geq 0}$ is a hypercontractive semi-group with hypercontractive function $q(t) = 1 + \exp(2t)$ (see Nelson [13] or Ané et al. [1] p. 22):

$$\forall t \geq 0, \forall g \in L^2(\gamma^n), \|P_t g\|_2 \leq \|g\|_{q^*(t)},$$

where $q^*(t) = 1 + \exp(-2t)$ is the conjugate exponent of $q(t)$. Thus,

$$\begin{aligned} \mathbb{E}_{\gamma^n} \left[\left(\frac{\partial}{\partial y_i} P_t F \right)^2 \right] &\leq e^{-2t} \left[\mathbb{E}_{\gamma^n} \left(\frac{\partial F}{\partial y_i} \right)^{q^*(t)} \right]^{\frac{2}{q^*(t)}}, \\ &= e^{-2t} \left[\mathbb{E}_{\gamma^n} \left(\int \frac{\partial f}{\partial y_i} d\lambda \right)^{q^*(t)} \right]^{\frac{2}{q^*(t)}}, \\ &\leq e^{-2t} \left[\mathbb{E}_{\lambda \otimes \gamma^n} \left(\int \frac{\partial f}{\partial y_i} \right)^{q^*(t)} \right]^{\frac{2}{q^*(t)}}, \end{aligned} \quad (4)$$

which follows from Jensen's inequality. Now, according to Hölder inequality,

$$\mathbb{E} \left[\left(\frac{\partial f}{\partial y_i} \right)^{q^*(t)} \right] \leq \mathbb{E} \left[\left(\frac{\partial f}{\partial y_i} \right)^2 \right]^{q^*(t)-1} \mathbb{E} \left[\left\| \frac{\partial f}{\partial y_i} \right\| \right]^{2-q^*(t)}.$$

Thus, inequality (4) implies:

$$\begin{aligned} &\int_0^{+\infty} \mathbb{E}_{\gamma^n} \left[\left(\frac{\partial}{\partial y_i} P_t F \right)^2 \right] dt \\ &\leq \int_0^{+\infty} e^{-2t} \left[\mathbb{E} \left[\left(\frac{\partial f}{\partial y_i} \right)^2 \right]^{q^*(t)-1} \mathbb{E} \left[\left\| \frac{\partial f}{\partial y_i} \right\| \right]^{2-q^*(t)} \right]^{\frac{2}{q^*(t)}} dt, \\ &= \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 \int_0^{+\infty} e^{-2t} \left(\frac{\left\| \frac{\partial f}{\partial y_i} \right\|_1}{\left\| \frac{\partial f}{\partial y_i} \right\|_2} \right)^{\frac{2(2-q^*(t))}{q^*(t)}} dt, \\ &= \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 \int_0^1 \left(\frac{\left\| \frac{\partial f}{\partial y_i} \right\|_1}{\left\| \frac{\partial f}{\partial y_i} \right\|_2} \right)^{2s} \frac{1}{(1+s)^2} ds. \end{aligned}$$

where last equality follows from the change of variable $s = \frac{1-e^{-2t}}{1+e^{-2t}}$. A combination of this upper bound with equality (3) and inequalities (2) and (1) gives the desired result. \square

The case of $\mathbb{R}^{\mathbb{N}}$

Let now $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$ be the space of functions $f \in L^2(\lambda \otimes \gamma^{\mathbb{N}})$ verifying the two following conditions:

(a) For all $i \in \mathbb{N}$, there exists a function h_i in $L^2(\lambda \otimes \gamma^{\mathbb{N}})$ such that

$$-\int_{\mathbb{R}} g'(y_i) f(x, y) dy_i = \int_{\mathbb{R}} g(y_i) h_i(x, y) dy_i, \lambda \otimes \gamma^{\mathbb{N}} \text{ a.s.}$$

for every smooth function $g : \mathbb{R} \mapsto \mathbb{R}$ having compact support. The function h_i is called the partial derivative of f with respect to y_i , and is denoted by $\frac{\partial f}{\partial y_i}$.

(b) The sum of the L^2 norms of the partial derivatives of f is finite:

$$\sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 < \infty.$$

Remark 2 *It is not hard to verify that for f depending on finitely many variables, say y_1, \dots, y_n , then $f \in H_1^2(\lambda \otimes \gamma^{\mathbb{N}}) \Leftrightarrow f \in H_1^2(\lambda \otimes \gamma^n)$.*

Let this Sobolev space be equipped with the norm:

$$\|f\|_H = \sqrt{\|f\|_2^2 + \sum_{q \in S} \|\nabla_q f\|_2^2 + \sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2}.$$

Proposition 2.1 extends to functions of countably infinite Gaussian variables as follows:

Proposition 2.2 *For every function f in $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$,*

$$\text{Var}(f) \leq \sum_{q \in S} \|\nabla_q f\|_2^2 + \sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 \phi \left(\frac{\left\| \frac{\partial f}{\partial y_i} \right\|_1}{\left\| \frac{\partial f}{\partial y_i} \right\|_2} \right),$$

where ϕ is as in Proposition 2.1.

The proof of Proposition 2.2 relies on the following simple approximation lemma combined with Proposition 2.1.

Lemma 2.3 *Let \mathcal{F}_n be the σ -algebra generated by the first n coordinate functions in $\mathbb{R}^{\mathbb{N}}$. Let $f \in H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$ and $f_n = \mathbb{E}(f|\mathcal{F}_n)$ be the conditional expectation of f with respect to \mathcal{F}_n . Then*

- (i) $f_n \in H_1^2(\lambda \otimes \gamma^n)$,
- (ii) For every i in $\{1, \dots, n\}$,

$$\frac{\partial f_n}{\partial y_i} = \mathbb{E} \left(\frac{\partial f}{\partial y_i} | \mathcal{F}_n \right) ,$$

- (iii) The following convergence takes place in $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$:

$$f_n \xrightarrow[n \rightarrow \infty]{} f .$$

In particular,

$$\sum_{q \in S} \|\nabla_q f_n\|_2^2 \xrightarrow[n \rightarrow +\infty]{} \sum_{q \in S} \|\nabla_q f\|_2^2 , \quad (5)$$

$$\text{Var}(f_n) \xrightarrow[n \rightarrow +\infty]{} \text{Var}(f) , \quad (6)$$

and

$$\sum_{i=1}^n \left\| \frac{\partial f_n}{\partial y_i} \right\|_2^2 \phi \left(\frac{\left\| \frac{\partial f_n}{\partial y_i} \right\|_1}{\left\| \frac{\partial f_n}{\partial y_i} \right\|_2} \right) \xrightarrow[n \rightarrow +\infty]{} \sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 \phi \left(\frac{\left\| \frac{\partial f}{\partial y_i} \right\|_1}{\left\| \frac{\partial f}{\partial y_i} \right\|_2} \right) . \quad (7)$$

Proof : Of course, $\mathbb{E}(f|\mathcal{F}_n)$ belongs to $L^2(\lambda \otimes \gamma^n)$. Let $g : \mathbb{R}^n \mapsto \mathbb{R}$ be a smooth function with compact support and $i \in \{1, \dots, n\}$. According to Fubini's Theorem,

$$\begin{aligned} & - \int \frac{\partial g}{\partial y_i} \mathbb{E}(f|\mathcal{F}_n) \, dy_1 \dots dy_n \\ &= - \int \mathbb{E} \left(\int \frac{\partial g}{\partial y_i} f \, dy_i | \mathcal{F}_n \right) \, dy_1 \dots dy_{i-1} dy_{i+1} \dots dy_n . \end{aligned}$$

Then, it follows from the definition of $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$ that:

$$- \int \frac{\partial g}{\partial y_i} f \, dy_i = \int g \frac{\partial f}{\partial y_i} \, dy_i ,$$

with probability 1. Therefore,

$$- \int \frac{\partial g}{\partial y_i} \mathbb{E}(f|\mathcal{F}_n) \, dy_1 \dots dy_n = \int g \mathbb{E} \left(\frac{\partial f}{\partial y_i} | \mathcal{F}_n \right) \, dy_1 \dots dy_n .$$

proving assertions (i) and (ii) of the Lemma.

Since for any function g in $L^2(\lambda \otimes \gamma^{\mathbb{N}})$, the martingale $(\mathbb{E}(g|\mathcal{F}_n))_{n \geq 0}$ converges to g in $L^2(\lambda \otimes \gamma^{\mathbb{N}})$, one has

$$\mathbb{E} \left(\mathbb{E} \left(\frac{\partial f}{\partial y_i} \middle| \mathcal{F}_n \right)^2 \right) \xrightarrow{n \rightarrow \infty} \mathbb{E} \left(\frac{\partial f}{\partial y_i} \right)^2 .$$

In addition, by Jensen's inequality,

$$\mathbb{E} \left(\mathbb{E} \left(\frac{\partial f}{\partial y_i} \middle| \mathcal{F}_n \right)^2 \right) \leq \mathbb{E} \left(\left(\frac{\partial f}{\partial y_i} \right)^2 \right) ,$$

and since f belongs to $H_1^2(\lambda \times \gamma^{\mathbb{N}})$,

$$\sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 < \infty .$$

By Lebesgue convergence theorem, it follows that:

$$\sum_{i \in \mathbb{N}} \left\| \frac{\partial \mathbb{E}(f|\mathcal{F}_n)}{\partial y_i} \right\|_2^2 \xrightarrow{n \rightarrow \infty} \sum_{i \in \mathbb{N}} \left\| \frac{\partial f}{\partial y_i} \right\|_2^2 < \infty . \quad (8)$$

Thus,

$$\sum_{i \in \mathbb{N}} \left\| \frac{\partial}{\partial y_i} [f - \mathbb{E}(f|\mathcal{F}_n)] \right\|_2^2 \xrightarrow{n \rightarrow \infty} 0 .$$

Finally, it is trivial to check that $\nabla_q \mathbb{E}(f|\mathcal{F}_n)$ converges to $\nabla_q f$ in $L^2(\lambda \times \gamma^{\mathbb{N}})$ as n tends to infinity. Therefore, the convergence takes place in $H_1^2(\lambda \times \gamma^{\mathbb{N}})$. Assertions (5) and (6) follow while the proof of assertion (7) follows from Lebesgue convergence theorem just as the proof of (8). \square

3. Extension to other measures

As usual, we can deduce from Proposition 2.2 other inequalities by mean of change of variables. To make this precise, let Ω be a measurable space and $\Psi : \mathbb{R}^{\mathbb{N}} \mapsto \Omega$ a measurable isomorphism (meaning that Ψ is one to one with Ψ and Ψ^{-1} measurables). Let $\Psi^* \gamma^{\mathbb{N}}$ denote the image of $\gamma^{\mathbb{N}}$ by Ψ . That is $\Psi^* \gamma^{\mathbb{N}}(A) = \gamma^{\mathbb{N}}(\Psi^{-1}(A))$. For $g : S \times \Omega \mapsto \mathbb{R}$ such that $g \circ (Id, \Psi) \in H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$, one obviously has

$$\text{Var}_{\lambda \otimes \Psi^* \gamma^{\mathbb{N}}}(g) = \text{Var}_{\lambda \otimes \gamma^{\mathbb{N}}}(g \circ \Psi)$$

and

$$\|\partial_{i,\Psi}g\|_{p,\Psi^*\gamma^{\mathbb{N}}} = \left\| \frac{\partial g \circ \Psi}{\partial y_i} \right\|_{p,\gamma^{\mathbb{N}}}$$

where $\partial_{i,\Psi}g$ is defined as

$$\partial_{i,\Psi}g(x, \omega) = \frac{\partial(x \circ \Psi)}{\partial y_i}(q, \Psi^{-1}(\omega)).$$

Hence inequality in Proposition 2.2 for $f = g \circ (Id, \Psi)$ transfers to the same inequality for g provided $\frac{\partial f}{\partial y_i}$ is replaced by $\partial_{i,\Psi}g$.

Example 1 Let $k \geq 2$ be an integer, $\mathbb{S}^{k-1} \subset \mathbb{R}^k$ the unit $k - 1$ dimensional sphere, $\Omega = (\mathbb{R}_*^+ \times \mathbb{S}^{k-1})^{\mathbb{N}}$, and let $E = (\mathbb{R}_*^k)^{\mathbb{N}}$. A typical point in Ω will be written as $(\rho, \theta) = (\rho^i, \theta^i)$ and a typical point in E as $y = (y^j)$. Now consider the change of variables $\Psi : E \mapsto \Omega$ given by $\Psi(y) = (\Psi(y)^j)$ with

$$\Psi^j(y) = \left(\|y^j\|^2, \frac{y^j}{\|y^j\|} \right).$$

The image of $(\gamma^k)^{\mathbb{N}} = \gamma^{\mathbb{N}}$ by Ψ is the product measure $\tilde{\gamma}^{\mathbb{N}}$ where $\tilde{\gamma}$ is the probability measure on $\mathbb{R}_*^+ \times S^{k-1}$ defined by

$$\tilde{\gamma}(dtdv) = \frac{1}{r_k} e^{-t/2} t^{k/2-1} \mathbf{1}_{t>0} dt dv$$

Here $r_k = \int_0^\infty e^{-t/2} t^{k/2-1} dt$, and dv stands for the uniform probability measure on \mathbb{S}^{k-1} . For $g : \{0, 1\}^S \times \Omega \mapsto \mathbb{R}$ with $g \circ (Id, \Psi) \in H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$, let

$$\left(\frac{\partial g}{\partial \rho^j}(x, \rho, \theta), \nabla_{\theta^i} g(x, \rho, \theta) \right) \in \mathbb{R} \times T_{\theta^i} \mathbb{S}^{k-1} \subset \mathbb{R} \times \mathbb{R}^k$$

denote the partial gradient of g with respect to the variable (ρ^i, θ^i) where $T_{\theta^i} \mathbb{S}^{k-1} \subset \mathbb{R}^k$ stands for the tangent space of \mathbb{S}^{k-1} at θ^i . It is not hard to verify that for all $i \in \mathbb{N}$ and $j \in \{1, \dots, k\}$,

$$\partial_{i,j,\Psi}g(x, \rho, \theta) = 2 \frac{\partial g}{\partial \rho^i}(x, \rho, \theta) \sqrt{\rho^i} \theta_j^i + \frac{1}{\sqrt{\rho^i}} [\nabla_{\theta^i} g(x, \rho, \theta)]_j. \quad (9)$$

Corollary 3.1 (χ^2 distribution) Let $\alpha > 0$, $k \geq 2$ an integer, and let \mathbb{R}_*^+ be equipped with the distribution

$$\nu(dt) = \frac{e^{-\alpha t} t^{k/2-1}}{\int_0^\infty e^{-\alpha s} s^{k/2-1} ds} \mathbf{1}_{t>0} dt,$$

and $\mathbb{R}_*^{+\mathbb{N}}$ with the product measure $\nu^{\mathbb{N}}$. For $g \in H_1^2(\nu^{\mathbb{N}})$ define

$$\nabla_i g(y) = \frac{\partial g}{\partial y_i}(y) \sqrt{y_i}.$$

Then

$$\text{Var}(g) \leq \frac{2}{\alpha} \sum_{i \in \mathbb{N}} \|\nabla_i g\|_2^2 \phi \left(c(k) \frac{\|\nabla_i g\|_1}{\|\nabla_i g\|_2} \right) \quad (10)$$

with

$$c(k) = \sqrt{k} \frac{\int_0^\pi |\cos(t)| \sin(t)^{k-2} dt}{\int_0^\pi \sin(t)^{k-2} dt} = \frac{2\sqrt{k}}{(k-1) \int_0^\pi \sin(t)^{k-2} dt}.$$

Proof: follows from (9) applied to the map $(x, \rho, \theta) \rightarrow g(\frac{\rho}{2\alpha})$. Details are left to the reader. \square

Corollary 3.2 (Uniform distribution on \mathbb{S}^n) Let dv_n denote the normalized Riemannian probability measure on $\mathbb{S}^n \subset \mathbb{R}^{n+1}$. For $g \in H_2^1(dv_n)$ and $i = 1, \dots, n+1$ let $\nabla_i g(\theta)$ denote the i^{th} component of $\nabla g(\theta)$ in \mathbb{R}^{n+1} (we see $T_\theta \mathbb{S}^n$ as the vector space of \mathbb{R}^{n+1} consisting of vector that are orthogonal to θ). Then, for $n \geq 2$,

$$\text{Var}(g) \leq \frac{1}{n-1} \sum_{i \in \mathbb{N}} \|\nabla_i g\|_2^2 \phi \left(\frac{\|\nabla_i g\|_1}{\|\nabla_i g\|_2} \right). \quad (11)$$

Proof: follows from (9) applied to the map $(x, \rho, \theta) \rightarrow g(\theta)$. Details are left to the reader. \square

Example 2 If one wants to get a result similar to Proposition 2.2 with γ replaced by another probability measure ν , one may of course perform the usual change of variables through inverse of repartition function. In the sequel, we denote by

$$g(x) = \frac{1}{\sqrt{2\pi}} e^{-\frac{x^2}{2}} \quad (12)$$

the density of the normalized Gaussian distribution, and by

$$G(x) = \int_{-\infty}^x g(u) du \quad (13)$$

its repartition function. For any function ϕ from \mathbb{R} to \mathbb{R} , we shall note $\tilde{\phi}$ the function from $\mathbb{R}^{\mathbb{N}}$ to $\mathbb{R}^{\mathbb{N}}$ such that $(\tilde{\phi}(x))_j = \phi(x_j)$.

Corollary 3.3 (Unidimensional change of variables) *Let ν be a probability on \mathbb{R}^+ absolutely continuous with respect to the Lebesgue measure, with density h and repartition function*

$$H(t) = \int_0^t h(u) du ,$$

then, for every function f on $\{0,1\}^S \times \mathbb{R}^n$ such that $f \circ (Id, \widetilde{H^{-1} \circ G}) \in H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$,

$$Var_{\lambda_S \otimes \nu^{\otimes \mathbb{N}}}(f) \leq \sum_{q \in S} \|\nabla_q f\|_2^2 + 2 \sum_{i \in \mathbb{N}} \|\nabla_i f\|_2^2 \phi \left(\frac{\|\nabla_i f\|_1}{\|\nabla_i f\|_2} \right) ,$$

where for every integer i ,

$$\nabla_i f(x, y) = \psi(y_i) \frac{\partial f}{\partial y_i}(x, y) ,$$

and ψ is defined on $I = \{t \geq 0 \text{ s.t. } h(t) > 0\}$:

$$\forall t \in I, \psi(t) = \frac{g \circ G^{-1}(H(t))}{h(t)} .$$

Proof : It is a straightforward consequence of Proposition 2.2, applied to $f \circ (Id, \widetilde{H^{-1} \circ G})$. □

4. Application to First Passage Percolation

It turns out that Corollary 3.1 is particularly well suited to adapt the argument of Benjamini, Kalai and Schramm [3] to show that First Passage Percolation has sublinear variance when the edges have a χ^2 distribution. This includes the important case of exponential distribution, for which First Passage Percolation becomes equivalent to a version of Eden growth model (see for instance Kesten [9], p.130). We do not want to restrict ourselves to those distributions. Nevertheless, due to the particular strategy that we adopt, we can only prove the result for some continuous edge distributions which behave roughly like a gamma distribution.

Definition 4.1 Let ν be a probability on \mathbb{R}^+ absolutely continuous with respect to the Lebesgue measure, with density h and repartition function

$$H(t) = \int_0^t h(u) du .$$

Such a measure will be said to be *nearly gamma* provided it satisfies the following set of conditions:

- (i) The set of $t \geq 0$ such that $h(t) > 0$ is an interval I ;
- (ii) h restricted to I is continuous;
- (iii) The map $\psi : I \mapsto \mathbb{R}$ defined by

$$\psi(y) = \frac{g \circ G^{-1}(H(y))}{h(y)},$$

is such that :

- (a) There exists a positive real number A such that

$$\forall y \in I, \psi(y) \leq A\sqrt{y}.$$

- (b) There exists $\varepsilon > 0$ such that, as a goes to zero,

$$\nu(y : \psi(y) \leq a) = O(a^\varepsilon).$$

In Definition 4.1, condition (iii) is of course the most tedious to check. A simple sufficient condition for a probability measure to be nearly gamma will be given in Lemma 4.3, the proof of which relies on the following asymptotics for the Gaussian repartition function G .

Lemma 4.2 *As x tends to $-\infty$,*

$$G(x) = g(x) \left(\frac{1}{|x|} + o\left(\frac{1}{x}\right) \right),$$

and as x tends to $+\infty$,

$$G(x) = 1 - g(x) \left(\frac{1}{x} + o\left(\frac{1}{x}\right) \right).$$

Consequently,

$$g \circ G^{-1}(y) \stackrel{y \rightarrow 0}{\sim} y\sqrt{-2 \log y},$$

and

$$g \circ G^{-1}(y) \stackrel{y \rightarrow 1}{\sim} (1-y)\sqrt{-2 \log(1-y)}.$$

Proof: A simple change of variable $u = x - t$ in G gives:

$$G(x) = g(x) \int_0^{+\infty} e^{-\frac{t^2}{2} + xt} dt,$$

Integrating by parts, we get:

$$\begin{aligned} G(x) &= g(x) \left(-\frac{1}{x} + \frac{1}{x} \int_0^{+\infty} t e^{-\frac{t^2}{2} + xt} dt \right), \\ &= g(x) \left(\frac{1}{|x|} + o\left(\frac{1}{x}\right) \right), \end{aligned}$$

as x goes to $-\infty$. Since $G(-x) = 1 - G(x)$, we get that, as x goes to $+\infty$:

$$G(x) = 1 - g(x) \left(\frac{1}{x} + o\left(\frac{1}{x}\right) \right) .$$

Let us turn to the asymptotic of $g \circ G^{-1}(y)$ as y tends to zero. Let $x = G^{-1}(y)$, so that “ y tends to zero” is equivalent to “ x tends to $-\infty$ ”. One has therefore,

$$\begin{aligned} G(x) &= \frac{g(x)}{|x|} (1 + o(1)) , \\ \log G(x) &= \log g(x) - \log |x| + o(1) , \\ &= -\frac{x^2}{2} - \log |x| + O(1) , \\ \log G(x) &= -\frac{x^2}{2} (1 + o(1)) , \\ |x| &= \sqrt{-2 \log G(x)} . \end{aligned}$$

Since $g(x) = |x|G(x)(1 + o(1))$,

$$g(x) = G(x) \sqrt{-2 \log G(x)} (1 + o(1)) ,$$

and therefore,

$$g \circ G^{-1}(y) = y \sqrt{-2 \log y} (1 + o(1)) ,$$

as y tends to zero. The asymptotic of $g \circ G^{-1}(y)$ as y tends to 1 is derived in the same way. \square

Given two functions r and l , we write $l(x) = \Theta(r(x))$ as x goes to x^* provided there exist positive constants $C_1 \leq C_2$ such that

$$C_1 \leq \liminf_{x \rightarrow x^*} \frac{r(x)}{l(x)} \leq \limsup_{x \rightarrow x^*} \frac{r(x)}{l(x)} \leq C_2 .$$

Lemma 4.3 *Assume that condition (i) and (ii) of Definition 4.1 hold. Let $0 \leq \underline{\nu} < \bar{\nu} \leq \infty$ denote the endpoints of I . Assume furthermore condition (iii) is replaced by conditions (iv) and (v) below.*

(iv) *There exists $\alpha > -1$ such that as x goes to $\underline{\nu}$,*

$$h(x) = \Theta((x - \underline{\nu})^\alpha) ,$$

(v) *$\bar{\nu} < \infty$ and there exists $\beta > -1$ such that as x goes to $\bar{\nu}$,*

$$h(x) = \Theta((\bar{\nu} - x)^\beta) ,$$

or $\bar{\nu} = \infty$ and

$$\exists A > \underline{\nu}, \forall t \geq A, C_1 h(t) \leq \int_t^\infty h(u) du \leq C_2 h(t),$$

where C_1 and C_2 are positive constants.

Then, ν is nearly gamma.

Proof: Since h is a continuous function on $]\underline{\nu}, \bar{\nu}[$, it attains its minimum on every compact set included in $]\underline{\nu}, \bar{\nu}[$. The minimum of h on $[a, b]$ is therefore strictly positive as soon as $\underline{\nu} < a \leq b < \bar{\nu}$. In order to show that condition (iii) holds, we thus have to concentrate on the behaviour of the function ψ near $\underline{\nu}$ and $\bar{\nu}$. Condition (iv) implies that, as x goes to $\underline{\nu}$,

$$H(x) = \Theta((x - \underline{\nu})^{\alpha+1}). \quad (14)$$

This, via Lemma 4.2, leads to

$$\psi(x) = \Theta\left((x - \underline{\nu})\sqrt{-\log(x - \underline{\nu})}\right), \quad (15)$$

as x goes to $\underline{\nu}$. Similarly, if $\bar{\nu} < \infty$, condition (v) implies that, as x goes to $\bar{\nu}$,

$$H(x) = \Theta((\bar{\nu} - x)^{\beta+1}), \quad (16)$$

which leads via Lemma 4.2 to

$$\psi(x) = \Theta\left((\bar{\nu} - x)\sqrt{-\log(\bar{\nu} - x)}\right), \quad (17)$$

as x goes to $\bar{\nu}$. Therefore, if $\bar{\nu} < \infty$, part (a) of condition (iii) holds. Besides, asymptotics (15) and (17) imply that for a small enough,

$$\begin{aligned} \nu(y \text{ s.t } \psi(y) \leq a) &\leq \mathbb{P}((x - \underline{\nu})\sqrt{-\log(x - \underline{\nu})} \leq Ca) \\ &\quad + \mathbb{P}((\bar{\nu} - x)\sqrt{-\log(\bar{\nu} - x)} \leq Ca), \end{aligned}$$

for a positive constant C . Thus,

$$\begin{aligned} \nu(y : \psi(y) \leq a) &\leq \mathbb{P}(x - \underline{\nu} \leq Ca) + \mathbb{P}(\bar{\nu} - x \leq Ca), \\ &= H(\underline{\nu} + Ca) + 1 - H(\bar{\nu} - Ca), \\ &= O(a^{\alpha+1}) + O(a^{\beta+1}), \end{aligned}$$

as a goes to zero, according to equations (14) and (16). Therefore, condition (iii) holds when $\bar{\nu}$ is finite.

Now, suppose that $\bar{\nu} = \infty$. Condition (v) implies:

$$\forall t \geq A, 1/C_2 \leq \frac{h(t)}{\int_t^\infty h(u)} du \leq 1/C_1 .$$

Integrating this inequality between A and y leads to the existence of three positive constants B, C'_1 and C'_2 such that:

$$\forall y \geq B, C'_1 y \leq \log \frac{1}{1 - H(y)} \leq C'_2 y .$$

Thus,

$$\forall y \geq B, C_1 \sqrt{C'_1 y} \leq \psi(y) \leq C_2 \sqrt{C'_2 y} . \tag{18}$$

This, combined with equation (15) proves that part (a) of condition (iii) holds. Now, equations (15) and (18) imply that for a small enough, there is a constant C such that,

$$\nu(y : \psi(y) \leq a) \leq \mathbb{P}((x - \underline{\nu})\sqrt{-\log(x - \underline{\nu})} \leq Ca) ,$$

which was already proved to be of order $O(a^{\alpha+1})$. This concludes the proof of Lemma 4.3. \square

Remark 3 With the help of Lemma 4.3, it is easy to check that most usual distributions are nearly gamma. This includes all gamma and beta distributions, as well as any probability measure whose density is bounded away from 0 on its support, and notably the uniform distribution on $[a, b]$, with $0 \leq a < b$. Nevertheless, remark that some distributions which have a sub-exponential upper tail may not satisfy the assumptions of Lemma 4.3, and be nearly gamma, though. For example, this is the case of the distribution of $|N|$, where N is a standard Gaussian random variable.

Now, we can state the main result of this article.

Theorem 4.4 *Let ν be a nearly gamma probability measure, with finite moment of order 2. Let μ denote the measure $\nu^{\otimes E}$. Then,*

$$\text{Var}_\mu(f_\nu) = O\left(\frac{|v|}{\log|v|}\right) ,$$

as $|v|$ tends to infinity.

In order to prove Theorem 4.4, we use the same averaging argument as in Benjamini et al. [3]. It relies on the following lemma:

Lemma 4.5 *There exists a constant $c > 0$, such that, for every $m \in \mathbb{N}^*$, there exists a function g_m from $\{0, 1\}^{m^2}$ to $\{0, \dots, m\}$ such that:*

$$\forall q \in \{1, \dots, m^2\}, |\nabla_q g_m| \in \{0, 1/2\},$$

and

$$\max_{y \in \{0, \dots, m\}} \lambda(x \text{ s.t. } g_m(x) = y) \leq \frac{c}{m}.$$

Since Benjamini et al. do not give a proof for this lemma, we offer the following one.

Proof: From Stirling's Formula,

$$\binom{m^2}{\lfloor m^2/2 \rfloor} \cdot \frac{m}{2^{m^2}} \xrightarrow{n \rightarrow \infty} \frac{1}{\sqrt{2\pi}},$$

and this implies that the following supremum is finite:

$$c_1 = \sup \left\{ 2 \binom{m^2}{\lfloor m^2/2 \rfloor} \cdot \frac{m}{2^{m^2}} \text{ s.t. } m \in \mathbb{N}^* \right\}.$$

Notice also that $c_1 \geq 1$. Now, let \preceq denote the alphabetical order $\{0, 1\}^{m^2}$, and let us list the elements in $\{0, 1\}^{m^2}$ as follows:

$$(0, 0, \dots, 0) = x_1 \preceq x_2 \preceq \dots \preceq x_{2^{m^2}} = (1, 1, \dots, 1).$$

For any m in \mathbb{N}^* , we define the following integer:

$$k(m) = \left\lceil \frac{2^{m^2}}{m} \right\rceil,$$

and the following function on $\{0, 1\}^{m^2}$:

$$\forall i \in \{1, \dots, 2^{m^2}\}, g_m(x_i) = \left\lfloor \frac{i}{k(m)} \right\rfloor.$$

Remark that $g_m(x_{2^{m^2}}) \leq m/c_1 \leq 1$. Therefore, g is a function from $\{0, 1\}^{m^2}$ to $\{0, \dots, m\}$. Now, suppose that x_i and x_l differ from exactly one coordinate. Then,

$$\begin{aligned} |i - l| &\leq \binom{m^2}{i} + \binom{m^2}{l}, \\ &\leq 2 \binom{m^2}{\lfloor m^2/2 \rfloor}, \\ &\leq \frac{2^{m^2}}{m}, \\ &\leq k(m). \end{aligned}$$

Consequently,

$$\begin{aligned} g_m(x_i) - g_m(x_l) &\leq \left\lfloor \frac{l}{k(m)} + 1 \right\rfloor - \left\lfloor \frac{l}{k(m)} + 1 \right\rfloor, \\ &= 1, \end{aligned}$$

which implies that $|\nabla_q g| \in \{0, 1/2\}$. Finally, for any $y \in \{0, \dots, m\}$, g takes the value y at most $k(m)$ times, and

$$\begin{aligned} \lambda(x \text{ s.t. } g_m(x) = y) &\leq \frac{k(m)}{2m^2}, \\ &\leq \frac{c_1}{m} + \frac{1}{2m^2}, \\ &\leq \frac{2c_1}{m}. \end{aligned}$$

So the lemma holds with $c = 2c_1$. □

We are now ready to prove Theorem 4.4.

Proof of Theorem 4.4 : In the whole proof, Y shall denote a random variable with distribution ν . We denote its second moment by σ^2 :

$$\sigma^2 = \mathbb{E}(Y^2).$$

Let $m = \lceil |v|^{\frac{1}{4}} \rceil$, and $S = \{1, \dots, d\} \times \{1, \dots, m^2\}$. Let $c > 0$ and g_m be as in Lemma 4.5. As in [3], we define a random vertex in $\{0, \dots, m\}^d$ by the following mean. For any $a = (a_{i,j})_{(i,j) \in S} \in \{0, 1\}^S$, let

$$z = z(a) = \sum_{i=1}^d g_m(a_{i,1}, \dots, a_{i,m^2}) \mathbf{e}_i,$$

where $(\mathbf{e}_1, \dots, \mathbf{e}_d)$ denotes the standard basis of \mathbb{Z}^d . We now equip the space $\{0, 1\}^S \times \mathbb{R}_+^E$ with the probability measure $\lambda \otimes \mu$, where λ is the uniform measure on $\{0, 1\}^S$, and we define the following function \tilde{f} on $\{0, 1\}^S \times \mathbb{R}_+^E$:

$$\forall (a, x) \in \{0, 1\}^S \times \mathbb{R}_+^E, \tilde{f}(a, x) = d_x(z(a), v + z(a)).$$

The first important point to notice is that f and \tilde{f} are not too far apart. Indeed, let $\alpha(a)$ be a path from 0 to $z(a)$, such that $|\alpha(a)| = |z(a)|$ (here, $|\alpha|$ is the number of edges in α). Let $\beta(a)$ denote the path $v + \alpha(a)$, which goes from v to $v + z(a)$. Then:

$$\begin{aligned} |\tilde{f}(a, x) - f(x)| &\leq d_x(0, z(a)) + d_x(v, v + z(a)), \\ &\leq \sum_{e \in \alpha(a)} x_e + \sum_{e \in \beta(a)} x_e, \end{aligned}$$

Thus, using $|z| \leq m$,

$$\|f - \tilde{f}\|_2 \leq \left\| \sum_{e \in \alpha(a)} x_e \right\|_2 + \left\| \sum_{e \in \beta(a)} x_e \right\|_2 \leq 2\mathbb{E}(|z(a)|\sigma) \leq 2m\sigma.$$

Therefore,

$$\begin{aligned} \sqrt{\text{Var}_\mu(f)} &= \|f - \mathbb{E}(f)\|_2 \leq \|f - \tilde{f}\|_2 + \|\tilde{f} - \mathbb{E}(\tilde{f})\|_2 + \|\mathbb{E}(\tilde{f}) - \mathbb{E}(f)\|_2 \\ &\leq 2\|f - \tilde{f}\|_2 + \|\tilde{f} - \mathbb{E}(\tilde{f})\|_2 \leq 4m\sigma + \sqrt{\text{Var}(\tilde{f})} \end{aligned} \quad (19)$$

It remains to bound $\text{Var}(\tilde{f})$. To this end, we will use Corollary 3.3. Using the notations of Corollary 3.3 and Definition 4.1, we need to prove that $\tilde{f} \circ (Id, \widetilde{H^{-1} \circ G})$ belongs to $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$ when ν is nearly gamma. This, and the application of Corollary 3.3 to \tilde{f} is contained in the following Lemma.

Lemma 4.6 *If ν is nearly gamma, the function $f_v \circ \widetilde{H^{-1} \circ G}$ belongs to $H_1^2(\gamma^{\mathbb{N}})$, $\tilde{f} \circ (Id, \widetilde{H^{-1} \circ G})$ belongs to $H_1^2(\lambda \otimes \gamma^{\mathbb{N}})$ and,*

$$\text{Var}(\tilde{f}) \leq \sum_{q \in S} \left\| \nabla_q \tilde{f} \right\|_2^2 + 2 \sum_{e \in E} \left\| \nabla_e \tilde{f} \right\|_2^2 \phi \left(\frac{\left\| \nabla_e \tilde{f} \right\|_1}{\left\| \nabla_e \tilde{f} \right\|_2} \right),$$

where for every edge e ,

$$\nabla_e \tilde{f}(a, x) = \psi(x_e) \frac{\partial \tilde{f}}{\partial x_e}(a, x). \quad (20)$$

Furthermore, conditionally to z , there is almost surely only one x -geodesic from z to $z + v$, denoted by $\gamma_x(z)$, and:

$$\frac{\partial \tilde{f}}{\partial x_e}(a, x) = \mathbb{1}_{e \in \gamma_x(z(a))}.$$

Proof: The fact that $f_v \circ \widetilde{H^{-1} \circ G}$ and $\tilde{f} \circ (Id, \widetilde{H^{-1} \circ G})$ are in L^2 follows from the basic fact that ν has a finite moment of order 2, and that $d_x(u, v)$ is dominated by a sum of $|v - u|$ independent variables with distribution ν , which is the length of a deterministic path of length $|v - u|$. We shall prove that $f_v \circ \widetilde{H^{-1} \circ G}$ satisfies the integration by part formula **(a)** of the definition of H_1^2 . The similar result for \tilde{f} is obtained in the same way.

Let e be an edge of \mathbb{Z}^d , and x^{-e} be an element of $(\mathbb{R}^+)^{E(\mathbb{Z}^d)\setminus\{e\}}$. For every y in \mathbb{R}^+ , we denote by (x^{-e}, y) the element x of $(\mathbb{R}^+)^{E(\mathbb{Z}^d)}$ such that:

$$x_e = y \text{ and } \forall e' \neq e, x_{e'} = x_{e'}^{-e} .$$

Now, we fix x^{-e} in $(\mathbb{R}^+)^{E(\mathbb{Z}^d)}$. We denote by g_e the function defined on \mathbb{R}^+ by:

$$g(y) = f_v(x^{-e}, y) .$$

We will show that there is a nonnegative real number y_∞ such that:

$$\begin{cases} \forall y \leq y_\infty, g(y) = g(0) + y \\ \text{and} \\ \forall y > y_\infty, g(y) = g(y_\infty) \end{cases} \quad (21)$$

For any $n \geq |v|$, let us denote by Γ_n the set of paths from 0 to v whose number of edges is not greater than n . We have:

$$g(y) = \inf_{n \geq |v|} g_n(y) ,$$

where

$$g_n(y) = \inf_{\gamma \in \Gamma_n} \sum_{e' \in \gamma} (x^{-e}, y)_{e'} .$$

The functions g_n form a nonincreasing sequence of nondecreasing functions:

$$\forall n \geq |v|, \forall y \in \mathbb{R}^+, \forall y' \geq y, g_{n+1}(y) \leq g_n(y) \leq g_n(y') .$$

In particular, this implies that for every y in \mathbb{R}^+ ,

$$g(y) = \lim_{n \rightarrow \infty} g_n(y) .$$

Now, we claim that, for every $n \geq |v| + 3$, there exists $y_n \in \mathbb{R}^+$ such that:

$$\begin{cases} \forall y \leq y_n, g_n(y) = g_n(0) + y \\ \text{and} \\ \forall y > y_n, g_n(y) = g_n(y_n) \end{cases} \quad (22)$$

and furthermore,

$$\text{the sequence } (y_n)_{n \geq |v|+3} \text{ is nonincreasing.} \quad (23)$$

Indeed, since Γ_n is a finite set, the infimum in the definition of g_n is attained. Let us call a path which attains this infimum an (n, y) -geodesic and let $\tilde{\Gamma}(n, y, e)$ be the set of (n, y) -geodesics which contain the edge e . Remark

that as soon as $n \geq |v| + 3$, there exists a real number A such that e does not belong to any (n, A) -geodesic: it is enough to take A greater than the sum of the length of three edges forming a path between the end-points of the edge e . Therefore, the following supremum is finite:

$$y_n = \sup\{y \in \mathbb{R}^+ \text{ s.t. } \tilde{\Gamma}(n, y, e) \neq \emptyset\} .$$

Now, if e belongs to an (n, y) -geodesic γ , for any $y' \leq y$, γ is an (n, y') -geodesic to which e belongs, and $g_n(y) - g_n(y') = y - y'$. If $\tilde{\Gamma}(n, y, e)$ is empty, then for any $y' \geq y$, e does not belong to any (n, y') -geodesic, and $g_n(y) = g_n(y')$. This proves that:

$$\forall y < y_n, g_n(y) = g_n(0) + y ,$$

$$\forall y, y' > y_n, g_n(y) = g_n(y') .$$

Since g_n is continuous, we have proved claim (22). Now remark that if e does not belong to any (n, y) -geodesic, then e does not belong to any $(n + 1, y)$ -geodesic, since $\Gamma_n \subset \Gamma_{n+1}$. Therefore, $y_{n+1} \leq y_n$, and this proves claim (23). Since $(y_n)_{n \geq |v|+3}$ is nonnegative, it converges to a nonnegative number y_∞ as n tends to infinity. Now, let n be a integer greater than $|v| + 3$:

$$\forall n \geq N, \forall y, y' > y_n, g_n(y) = g_n(y') .$$

Since $y_n \leq y_N$,

$$\forall n \geq N, \forall y, y' > y_N, g_n(y) = g_n(y') .$$

Letting n tend to infinity in the last equation, we get:

$$\forall N \geq |v| + 3, \forall y, y' > y_N, g(y) = g(y') .$$

Therefore,

$$\forall y, y' > y_\infty, g(y) = g(y') .$$

On the other side,

$$\forall n \geq |v| + 3, \forall y \leq y_n, g_n(y) = g_n(0) + y .$$

Since $y_n \geq y_\infty$,

$$\forall n \geq |v| + 3, \forall y \leq y_\infty, g_n(y) = g_n(0) + y .$$

Letting n tend to infinity in the last expression, we get:

$$\forall y \leq y_\infty, g(y) = g(0) + y .$$

Finally, g is continuous. Indeed, the convergent sequence (g_n) is uniformly equicontinuous, since all these functions are 1-Lipschitz, and the continuity of g follows from Arzelà-Ascoli Theorem. We have proved claim (21). Remark that $y_\infty = y_\infty(x^{-e})$ depends on x^{-e} . We define, for any x^{-e} ,

$$h_e(x^{-e}, x_e) = \begin{cases} 1 & \text{if } x_e \leq y_\infty(x^{-e}) \\ 0 & \text{if } x_e > y_\infty(x^{-e}) \end{cases} .$$

It is easy to see that, for any smooth function $F : \mathbb{R} \rightarrow \mathbb{R}$ having compact support, for any x^{-e} ,

$$- \int_{\mathbb{R}} F'(x_e) f_v(x^{-e}, x_e) dx_e = \int_{\mathbb{R}} F(x_e) h_e(x^{-e}, x_e) dx_e . \quad (24)$$

It is known that there is almost surely a geodesic from 0 to v (see [6] for instance), i.e the infimum in the definition of f_v is attained with probability 1. Furthermore, in this setting, where the distribution of the lengths is continuous, there is almost surely only one x -geodesic from 0 to v . For any x , we shall denote by $\gamma_x(0)$ the unique x -geodesics from 0 to $0 + v$. Then, with ν -probability 1, one can see from the definitions of y_n and y_∞ that:

$$h_e(x^{-e}, x_e) = \mathbb{1}_{e \in \gamma_x(0)} . \quad (25)$$

Performing the change of variable $x \mapsto \widetilde{H^{-1} \circ G}$ in equation (24), one gets the integration by parts formula **(a)** for $f_v \circ \widetilde{H^{-1} \circ G}$, with the following partial derivative with respect to x_e :

$$x \mapsto \psi(x_e) h_e(\widetilde{H^{-1} \circ G}(x)) .$$

The fact that the sum of the L^2 norms of the partial derivatives of f_v and \tilde{f} is finite follows from estimate (28). Then, inequality (20) is a consequence of Corollary 3.3, and the expression of $\frac{\partial \tilde{f}}{\partial x_e}(a, x)$ is derived in the same way than (25). \square

Our next goal is to find a good upper-bound for $\left\| \nabla_e \tilde{f} \right\|_2^2$, for any edge e . Let A be as in Definition 4.1. It follows from Lemma 4.6 and the fact that ν is nearly gamma that:

$$\begin{aligned} \left\| \nabla_e \tilde{f} \right\|_2^2 &= \mathbb{E} \left(\psi(x_e)^2 \mathbb{1}_{e \in \gamma_x(z)} \right) , \\ &\leq A^2 \mathbb{E}_\lambda \left(\mathbb{E}_\mu \left(x_e \mathbb{1}_{e \in \gamma_x(z)} \right) \right) . \end{aligned}$$

Now, we use the fact that for any fixed z , μ is invariant under translation by z .

$$\begin{aligned}
 \left\| \nabla_e \tilde{f} \right\|_2^2 &\leq A^2 \mathbb{E}_\lambda \left(\mathbb{E}_\mu \left(x_{e-z} \mathbb{1}_{e-z \in \gamma_x(0)} \right) \right), \\
 &= A^2 \mathbb{E}_\mu \left(\sum_{e' \in \gamma_x(0)} \mathbb{E}_\lambda \left(x_{e'} \mathbb{1}_{e-z=e'} \right) \right), \\
 &= A^2 \mathbb{E}_\mu \left(\sum_{e' \in \gamma_x(0)} x_{e'} \mathbb{P}_\lambda(z = e - e') \right), \\
 &\leq A^2 \sup_{z_0} \mathbb{P}(z = z_0) \mathbb{E}_\mu \left(\sum_{e' \in \gamma_x(0) \cap \mathcal{Q}_e} x_{e'} \right),
 \end{aligned}$$

where $\mathcal{Q}_e = \{e' \in E(\mathbb{Z}^d) \text{ s.t. } \mathbb{P}(z = e - e') > 0\}$. Using Lemma 4.5,

$$\sup_{z_0} \mathbb{P}(z = z_0) \leq \left(\frac{c}{m} \right)^d.$$

Let e^- and e^+ denote the end-points of e , and $\mathcal{B}(0, dm)$ be the L^1 -ball with center 0 and radius dm . Remark now that \mathcal{Q}_e is included in the set of edges $e + \mathcal{B}(0, dm)$, which is itself included in $\mathcal{B}(e^-, dm + 1)$. This simply follows from the fact that g_m takes its values in $\{0, \dots, m\}$. Let r be a deterministic path going through every vertex of the surface of the ball $\mathcal{B}(e^-, dm + 1)$, and such that there is a constant C (depending only on d) such that $|r| \leq Cm^{d-1}$. From the definition of a geodesic, we get:

$$\begin{aligned}
 \mathbb{E}_\mu \left(\sum_{e' \in \gamma_x(0) \cap \mathcal{Q}_e} x_{e'} \right) &\leq \mathbb{E} \left(\sum_{e' \in \mathcal{B}(e^-, dm+1)} x_{e'} \right), \\
 &\leq \mathbb{E} \left(\sum_{e' \in r} x_{e'} \right), \\
 &\leq Cm^{d-1} \mathbb{E}(Y).
 \end{aligned}$$

Therefore,

$$\left\| \nabla_e \tilde{f} \right\|_{2, \lambda \otimes \mu}^2 \leq CA^2 dm^{d-1} \mathbb{E}(Y) \left(\frac{c}{m} \right)^d \leq A^2 d \mathbb{E}(Y) \frac{c^d}{|v|^{\frac{1}{4}}}. \quad (26)$$

In order to use Corollary 3.1, one needs to bound from below the quotient $\left\| \nabla_e \tilde{f} \right\|_2 / \left\| \nabla_e \tilde{f} \right\|_1$. Let a be a positive real number.

$$\begin{aligned} \left\| \nabla_e \tilde{f} \right\|_1 &= \mathbb{E} \left(\psi(x_e) \mathbb{1}_{e \in \gamma_x(z)} \right), \\ &= \mathbb{E} \left(\psi(x_e) \mathbb{1}_{e \in \gamma_x(z)} \mathbb{1}_{\psi(x_e) \leq a} \right) + \mathbb{E} \left(\psi(x_e) \mathbb{1}_{e \in \gamma_x(z)} \mathbb{1}_{\psi(x_e) > a} \right), \\ &\leq \mathbb{E} \left(\psi(x_e) \mathbb{1}_{\psi(x_e) \leq a} \right) + \frac{1}{a} \mathbb{E} \left(\psi(x_e)^2 \mathbb{1}_{e \in \gamma_x(z)} \mathbb{1}_{\psi(x_e) > a} \right), \\ &\leq a\nu(\psi(x_e) \leq a) + \frac{1}{a} \left\| \nabla_e \tilde{f} \right\|_2^2. \end{aligned}$$

Since ν is nearly gamma, there exists $\varepsilon > 0$ such that:

$$\nu(\psi(x_e) \leq a) = O(a^\varepsilon),$$

as a goes to zero. Now, let us choose $a = \left\| \nabla_e \tilde{f} \right\|_2^{2/(2+\varepsilon)}$. Thus,

$$\left\| \nabla_e \tilde{f} \right\|_1 = O \left(\left\| \nabla_e \tilde{f} \right\|_2^{1+\frac{\varepsilon}{2+\varepsilon}} \right).$$

This, via inequality (26), leads to

$$\log \frac{\left\| \nabla_e \tilde{f} \right\|_1}{\left\| \nabla_e \tilde{f} \right\|_2} = O(\log |v|) \tag{27}$$

In addition, if one chooses a particular x -geodesic γ from z to $z+v$, and a L^1 -geodesic $\alpha(z)$ from z to $z+v$,

$$\begin{aligned} \sum_{e \in E} \left\| \nabla_e \tilde{f} \right\|_2^2 &= \mathbb{E} \left(\sum_{e \in \gamma} \psi(x_e)^2 \right), \\ &\leq A \mathbb{E} \left(\sum_{e \in \gamma} x_e \right), \\ &\leq A \mathbb{E}(Y) \times \mathbb{E}_\lambda(|\alpha(z)|), \\ \sum_{e \in E} \left\| \nabla_e \tilde{f} \right\|_2^2 &\leq A \mathbb{E}(Y) |v|. \end{aligned} \tag{28}$$

Collecting estimates (27) and (28),

$$\begin{aligned} \sum_{e \in E} \left\| \nabla_e \tilde{f} \right\|_2^2 \phi \left(\frac{2\sqrt{2} \left\| \nabla_e \tilde{f} \right\|_1}{\pi \left\| \nabla_e \tilde{f} \right\|_2} \right) &\leq \sum_{e \in E} \left\| \nabla_e \tilde{f} \right\|_2^2 \sup_{e \in E} \phi \left(\frac{2\sqrt{2} \left\| \nabla_e \tilde{f} \right\|_1}{\pi \left\| \nabla_e \tilde{f} \right\|_2} \right), \\ &= O \left(\frac{|v|}{\log |v|} \right), \end{aligned} \tag{29}$$

where we used the fact that $\phi(u) \sim -1/(2 \log(u))$ when u goes to zero. According to Lemma 4.5, for any $q \in S$, $|\nabla_q g_m| \in \{0, 1/2\}$. Therefore,

$$\forall q \in S, \|\nabla_q f\|_2^2 \leq \sigma^2,$$

and thus,

$$\sum_{q \in S} \|\nabla_q f\|_2^2 \leq |S| \sigma^2 = dm^2 \sigma^2.$$

Inequality (29), the assumption $|m| = \lceil |v|^{1/4} \rceil$ and Lemma 4.6 lead therefore to:

$$\text{Var}(\tilde{f}) = O\left(\frac{|v|}{\log |v|}\right).$$

This, together with inequality (19) implies that:

$$\text{Var}(f_v) = O\left(\frac{|v|}{\log |v|}\right).$$

The proof of Theorem 4.4 is complete. \square

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