ANOMALOUS SLOW DIFFUSION FROM PERPETUAL HOMOGENIZATION

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This paper is concerned with the asymptotic behavior of solutions of stochastic differential equations $d\gamma_t = d\omega_t - \nabla V(\gamma_t)dt$, $\gamma_0 = 0$. When $d = 1$ and $V$ is not periodic but obtained as a superposition of an infinite number of periodic potentials with geometrically increasing periods $V(x) = \sum_{k=0}^{\infty} U_k(x/R_k)$, where $U_k$ are smooth functions of period 1, $U_k(0) = 0$, and $R_k$ grows exponentially fast with $k$, we can show that $\gamma_t$ has an anomalous slow behavior and we obtain quantitative estimates on the anomaly using and developing the tools of homogenization. Pointwise estimates are based on a new analytical inequality for subharmonic functions. When $d \geq 1$ and $V$ is periodic, quantitative estimates are obtained on the heat kernel of $\gamma_t$, showing the rate at which homogenization takes place. The latter result proves Davies’ conjecture and is based on a quantitative estimate for the Laplace transform of martingales that can be used to obtain similar results for periodic elliptic generators.

1. Introduction. It is now well known that natural Brownian motions on various disordered or complex structures are anomalously slow.

These mechanisms of the slow diffusion for instance are well understood for very regular strictly self-similar fractals. The archetypical specific example of a deep problem being the one solved in Barlow and Bass (1999) on the Sierpinski carpet [which is infinitely ramified, a codeword for hard to understand rigorously: for a survey on diffusions on fractals we refer to Barlow (1998), for an alternative approach to Osada (1995) and for the random Sierpinski carpet to Hambly, Kumagai, Kusuoka and Zhou (1998)]. It appears that the main feature is the existence of an infinite number of scales of obstacle (with proper size) for the diffusion.

It is our object to show that one can implement the common idea that this last feature (infinitely many scales) is the key to the possibility of anomalous diffusion, in a general context using the tools of homogenization.

The strategy of the proof might appear paradoxical: it is not a priori very sensible to try to prove that the diffusion is anomalous by the use of homogenization theory which is a vast mathematical machine destined to prove an opposite result.
that is, a central limit theorem and thus normal diffusion. But it will be shown that when the homogenization process is not finished, an anomalous behavior whose characteristics are controlled by homogenization theory might appear.

This paper will focus on the subdiffusive behavior in dimension one (Section 2.1), which will allow the introduction of a concept of differentiation between spatial scales that can be applied to a more general framework.

The proof of the anomaly of the exit times is based on a new quantitative analytical inequality for subharmonic functions (Section 2.3) that is linked with stability properties of elliptic divergence form operators.

The extension of those results to higher dimensions has been done in Ben Arous and Owhadi (2001) and to the super-diffusive case in Ben Arous and Owhadi (2002) and Owhadi (2001b).

The control of the anomalous heat kernel tail is based on sharp quantitative estimates for the Laplace transform of a martingale. These estimates allow us to put into evidence the rate at which homogenization takes place on the behavior of the heat kernel of an elliptic generator in any dimension (Section 2.2). The quantitative control of the heat kernel in homogenization theory outside any asymptotic regime has been recognized as difficult and important Norris (1997). For instance, this problem is at the center of Davies’ conjecture emphasized as “well beyond existing results” [Davies (1993)]. With Theorem 2.8 we give a proof of that conjecture in any dimension for elliptic operators with only bounded coefficients.

1.1. History. The idea of associating homogenization (or renormalization) on large number of scales with the anomaly of a physical system has already been applied from an heuristic point of view to several physical models.

Maybe one of the oldest of such applications is to differential effective medium theories which was first proposed by Bruggeman to calculate the conductivity of a two-component composite structure formed by successive substitutions [Bruggerman (1935) and Garland and Tanner (1977)] and generalized in Norris (1985) to materials with more than two phases. For instance this theory has been applied to compute the anomalous electrical and acoustic properties of fluid-saturated sedimentary rocks [Sen, Scala and Cohen (1981)]. More recently this problem has been analyzed from a rigorous point of view in Avellaneda (1987), Kozlov (1995), Allaire and Briane (1996) and Jikov and Kozlov (1999).

The heuristic application of this idea to prove the anomalous behavior of diffusion seems to have been done only for the super-diffusive case that is to say for a diffusion evolving among a large number of divergence-free drifts. Maybe this is explained by the strong motivation to explore convective transports in turbulent flows which are known to be characterized by a large number of scales of eddies. The first observation was empirical: in Richardson (1926) Richardson empirically conjectured that the diffusion coefficient $D_\lambda$ in turbulent air depends on the scale length $\lambda$ of the measurement. More recently physicists and mathematicians have started to investigate on the super-diffusive phenomenon (from both heuristic and
rigorous points of view) using the tools of homogenization or renormalization (the first cousin of multi-scale homogenization): we refer to Avellaneda and Majda (1990), Glimm and Zhang (1992), Avellaneda (1996), Bhattacharya (1999), Fannjiang and Papanicolaou (1994) and Fannjiang and Komorowski (2000).

1.2. The model. Let us consider in dimension one a Brownian motion with a drift given by the gradient of a potential $V$, that is, the solution of the stochastic differential equation:

$$dy_t = d\omega_t - \nabla V(y_t) \, dt, \quad y_0 = 0.$$  

The multi-scale potential $V$ is given by a sum of infinitely many periodic functions with (geometrically) increasing periods:

$$V = \sum_{n=0}^{\infty} U_n \left( \frac{x}{R_n} \right).$$

In this formula we have two important ingredients: the potentials $U_k$ and the scale parameters $R_k$. We will now describe the hypothesis we make on these two items of our model.

1. Hypotheses on the potentials $U_k$. We will assume that

$$U_k \in C^\infty(\mathbb{T}),$$  

$$U_k(0) = 0.$$  

Here $C^\infty(\mathbb{T}^d)$ denotes the space of smooth functions on the torus $\mathbb{T} := \mathbb{R} \setminus \mathbb{Z}$. We will also assume that the first derivate of the $U_k$ are uniformly bounded, that is,

$$K_1 := \sup_{k \in \mathbb{N}} \sup_{x \neq y} |U_k(x) - U_k(y)| / |x - y| < \infty.$$  

We will also need the notation

$$K_0 := \sup_{k \in \mathbb{N}} \text{Osc}(U_k),$$

where the oscillation of $U_k$ is given by $\text{Osc}(U) := \sup U - \inf U$. We write $D(U_k)$ for the effective diffusivities associated with the potentials $U_k$: if $z_t$ is the solution of $dz_t = d\omega_t - \nabla U_k(z_t) \, dt$, it is well known [Olla (1994)] that as $\varepsilon \downarrow 0$, $\varepsilon z_{t/\varepsilon^2}$ converges in law toward a Brownian motion with covariance matrix $D(U_k)$ given by

$$D(U_n) = \left( \int_{\mathbb{T}} e^{2U_n(x)} \, dx \int_{\mathbb{T}} e^{-2U_n(x)} \, dx \right)^{-1}.$$  

We also assume that the effective diffusivity matrices of the $U_k$'s are uniformly bounded away from 0 and 1.

$$\lambda_{\min} = \inf_{n \in \mathbb{N}} D(U_n) > 0 \quad \text{and} \quad \lambda_{\max} = \sup_{n \in \mathbb{N}} D(U_n) < 1.$$
2. Hypotheses on the scale parameters $R_k$. $R_k$ is a spatial scale parameter growing exponentially fast with $k$, more precisely we will assume that $R_0 = r_0 = 1$ and that the ratios between scales defined by (we write $\mathbb{N}^*$ the set of integers different from 0)

$$r_k = R_k/R_{k-1} \in \mathbb{N}^*$$

for $k \geq 1$, are integers uniformly bounded away from 1 and $\infty$: we will denote by

$$\rho_{\min} := \inf_{k \in \mathbb{N}^*} r_k \quad \text{and} \quad \rho_{\max} := \sup_{k \in \mathbb{N}^*} r_k$$

and assume that

$$\rho_{\min} \geq 2 \quad \text{and} \quad \rho_{\max} < \infty. \quad (11)$$

Since $\|\nabla V\|_{\infty} < \infty$, it is well known that the solution of (1) exists; is unique up to sets of measure 0 with respect to the Wiener measure and is a strong Markov continuous Feller process.

REMARK 1.1. Note that if $\forall n, U_n \in \{W_1, \ldots, W_p\}$, the $(W_i)$ being nonconstant, then the conditions (8) and (5) are trivially satisfied.

2. Main results.

2.1. Subdiffusive behavior. Our first objective is to show that the solution of (1) is abnormally slow and the asymptotic subdiffusivity will be characterized in three ways:

- as an anomalous behavior of the expectation of $\tau(0, r)$ (the exit time from a ball of radius $r$, for $r \to \infty$, i.e., $\mathbb{E}_0[\tau(0, r)] \sim r^{2+\nu}$);
- as an anomalous behavior of the variance at time $t$, that is, $\mathbb{E}_0[\gamma_t^2] \sim t^{1-\nu}$ as $t \to \infty$;
- as an anomalous (non-Gaussian) behavior of the tail of the transition probability of the process.

More precisely there exists a constant $\rho_0(K_0, K_1, \lambda_{\max})$ such that:

THEOREM 2.1. If $\rho_{\min} > \rho_0$ and $\tau(0, r)$ is the exit time associated with the solution of (1) then there exists a constant $C_1$ depending on $K_0, K_1$ such that

$$\mathbb{E}_0[\tau(0, r)] = r^{2+\nu_1(r)+\varepsilon(r)}, \quad (12)$$

where $\varepsilon(r) \to 0$ as $r \to \infty$ and

$$0 < -\frac{\ln \lambda_{\max}}{\ln \rho_{\max}} - \frac{C_1}{\rho_{\min} \ln \rho_{\max}} \leq \nu_1(r) \leq -\frac{\ln \lambda_{\min}}{\ln \rho_{\min}} + \frac{C_1}{\rho_{\min} \ln \rho_{\min}}. \quad (13)$$
THEOREM 2.2. If $\rho_{\text{min}} > \rho_0$ and $y_t$ is a solution of (1) then there exists a constant $C_2$ depending on $K_0, K_1$ and a time $t_0$ depending on $K_1, \rho_{\text{min}}, \rho_{\text{max}}, \lambda_{\text{max}}$ such that, for $t > t_0$,

$$E[y_t^2] = t^{1-v_2(t)/2},$$

where

$$0 < -\frac{\ln \lambda_{\text{max}}}{\ln \rho_{\text{max}}} - \frac{C_2}{\ln \rho_{\text{min}} \ln \rho_{\text{max}}} \leq v_2(t) \leq -\frac{\ln \lambda_{\text{min}}}{\ln \rho_{\text{min}}} + \frac{C_2}{(\ln \rho_{\text{min}})^2}. \tag{15}$$

THEOREM 2.3. If $\rho_{\text{min}} > \rho_0$ and $y_t$ is a solution of (1) then there exist constants $C_5$ depending on $K_0, K_1, R_2; C_3$ on $K_0, K_1, \rho_{\text{min}}; C_4, C_6$ and $C_7$ on $K_0, K_1$ such that if $t, h > 0$ and

$$\frac{t}{h} \geq C_5 \quad \text{and} \quad \frac{h^2}{t} \geq C_3 \left( \frac{t}{h} \right)^{(\ln \lambda_{\text{max}})/(2 \ln \rho_{\text{max}})+C_4/(\ln \rho_{\text{min}})^2},$$

then

$$\ln \mathbb{P}[|y_t| \geq h] \leq -C_6 \frac{h^2}{t} \left( \frac{t}{h} \right)^{v_3} \tag{17}$$

with

$$v_3 = -\frac{\ln \lambda_{\text{max}}}{\ln \rho_{\text{max}}} - \frac{C_7}{\ln \rho_{\text{min}} \ln \rho_{\text{max}}} > 0. \tag{18}$$

REMARK 2.4. The second condition in (16) is really needed since the leading exponent associated with $(t/h)$ is $\left( \frac{\ln \lambda_{\text{max}}}{2 \ln \rho_{\text{max}}} \right)$, that is, half the one associated to $v_3$. This condition corresponds to a frontier with a heat kernel diagonal regime.

2.1.1. Description of the proofs. Before discussing the results further we want to describe the proof. A perpetual homogenization process takes place over the infinite number of scales $0, \ldots, n, \ldots$. The idea is to distinguish, when one tries to estimate (12), (14) or (17), the smaller scales which have already been homogenized ($0, \ldots, n_{\text{ef}}$ called effective scales), the bigger scales which have not had a visible influence on the diffusion ($n_{\text{dri}}, \ldots, \infty$ called drift scales because they will be replaced by a constant drift in the proof) and some intermediate scales that manifest their particular shapes in the behavior of the diffusion ($n_{\text{dri}} - 1 = n_{\text{ef}} + n_{\text{per}}$ called perturbation scales because they will enter in the proof as a perturbation of the homogenization process over the smaller scales). To estimate (12) for instance, if one considers the periodic approximation of the potential

$$V_0^n(x) = \sum_{k=0}^n U_k(x/R_k), \tag{19}$$
the corresponding process $y_t^{(n)}$ will have an asymptotic (homogenized) variance 
[Olla (1994)]:

$$D(V_0^n) = \left( \int_T e^{2V_0^n(R_n,x)} \, dx \int_T e^{-2V_0^n(R_n,x)} \, dx \right)^{-1}.$$  

(20)

The variance $D(U_0)$ is smaller than 1 and because of the geometric growth of the periods $R_n$ and a minimal separation between them (i.e., $\rho_{\text{min}} > \rho_0$), $D(V_0^n)$ decreases exponentially fast in $n$.

By homogenization theory, $y_t^n$ is characterized by a mixing length $\xi_m(V_0^n)$ such that if one writes $\tau_n$ as its associated exit times, then for $r > \xi_m(V_0^n)$,

$$E_0[\tau^n(0, r)] \sim r^2 D(V_0^n).$$  

(21)

Writing $n_{\text{ef}}(r) = \sup\{n : R_n \leq r\}$ one proves that $E_0[\tau(0, r)] \sim E_0[\tau^{n_{\text{ef}}(r)}(0, r)]$ by showing the stability of $E_0[\tau(0, r)]$ under the influence of $V_{n_{\text{ef}}(r)+1} = \sum_{k=n_{\text{ef}}(r)+1}^\infty U_k(x/R_k)$. This control is based on a new analytical inequality which shall be described in the sequel and allows us to obtain that

$$E_0[\tau^{n_{\text{ef}}(r)}(0, r)] e^{-6 \text{Osc}_r(V_{n_{\text{ef}}(r)+1}^\infty)} \leq E_0[\tau(0, r)] \leq E_0[\tau^{n_{\text{ef}}(r)}(0, r)] e^{6 \text{Osc}_r(V_{n_{\text{ef}}(r)+1}^\infty)}.$$  

(22)

In these inequalities, $\text{Osc}_r(V_{n_{\text{ef}}(r)+1}^\infty)$ stands for $\sup_{B(0,r)} V_{n_{\text{ef}}(r)+1}^\infty - \inf_{B(0,r)} V_{n_{\text{ef}}(r)+1}^\infty$ and is controlled by

$$\text{Osc}_r(V_{n_{\text{ef}}(r)+1}^\infty) \leq \text{Osc}(U_{n_{\text{ef}}(r)+1}^\infty) + \|\nabla V_{n_{\text{ef}}(r)+1}^\infty\|_{\infty} r,$$

that is, $n_{\text{ef}}(r) + 1$ acts as a perturbation scale and $n_{\text{ef}}(r) + 2, \ldots, \infty$ as drift scales. From this,

$$E_0[\tau(0, r)] \sim \frac{r^2}{D(V_0^{n_{\text{ef}}(r)})}.$$  

(23)

Thus, if

$$- \lim_{r \to \infty} \inf \frac{1}{\ln r} \ln D(V_0^{n_{\text{ef}}(r)}) > 0$$

one has subdiffusivity, in the sense as defined above.

The proof of (14) follows similar lines by the introduction mixing times $\tau_m(V_0^n)$ and visibility times $\tau_v(V_{\infty}^p)$ [such that for $\tau_m(V_0^n) < t < \tau_v(V_{\infty}^p)$; $V_{\infty}^p$ does not have a real influence on the behavior of the diffusion $y_t$ and $V_0^n$ has been homogenized]. Then choosing $n_{\text{ef}}(t) = \sup\{n : \tau_m(V_0^n) \leq t\}$ one obtains the following.
**Proposition 2.5.** Letting $\nu_2(t)$ be the function associated with (14), one has, for $t > t_1, \rho_{\min}, \rho_{\max}, \lambda_{\max}$,

\begin{equation}
\nu_{\text{ef}}(t) \left(1 - \frac{C_{K_1}}{\ln \rho_{\min}}\right) \leq \nu_2(t) \leq \nu_{\text{ef}}(t) \left(1 + \frac{C_{K_1}}{\ln \rho_{\min}}\right),
\end{equation}

with

\begin{equation}
\nu_{\text{ef}}(t) = \frac{\ln(1/\lambda_{\text{ef}}(t))}{\ln \rho_{\text{ef}}(t)}
\end{equation}

This proposition shows that this separation between scales is more than a conceptual tool: it does reflect the underlying phenomenon. Indeed the anomalous function $\nu_2(t)$ is given in the first order in $1/(\ln \rho_{\min})$ by the number of effective scales by $E[y_2^2] \sim t D(V_0^{n_{\text{ef}}})$, and in this approximation $\nu_2(t) \sim \nu_{\text{ef}}(t)$ where $\nu_{\text{ef}}(t)$ corresponds to a medium in which the ratios $r_n$ and the effective diffusivities $D(U_n)$ have been replaced by their geometric mean over the $n_{\text{ef}} + 1$ effective scales. The origin of the constant $C_{K_1}/(\ln \rho_{\min})$ in (24) is the perturbation scales. More precisely, one has to fix the drift scales by $n_{\text{dri}}(t) = \inf \{n : \tau_\nu(V_{\infty}^n) \geq t\}$, and in general there is a gap between $n_{\text{ef}}(t)$ and $n_{\text{dri}}(t)$, the scales $U_n$ situated in this gap manifest their particular shape in the behavior of $\nu_2(t)$ and since no hypothesis have been made on those shapes one has to take into account their influence as a perturbation.

One may notice that in many papers on diffusions on fractals [see, e.g., Barlow (1998), Section 3] obtaining estimates on hitting times is essentially the key to the whole problem and the same is true here: this strategy has been adapted in Ben Arous and Owhadi (2001). In this paper we have chosen to not use this strategy in order to put an emphasis on the role played by the never-ending homogenization process taking place on these diffusions on fractals. Indeed one might wonder why the estimates of the behavior of Brownian motion on fractals are of the form

\begin{equation}
E[y_2^2] \sim t^{2/d_w},
\end{equation}

\begin{equation}
E[\tau(0, r)] \sim r^{d_w},
\end{equation}

\begin{equation}
\ln p(t, x, y) \sim -\left(\frac{|x - y|^{d_w}}{t}\right)^{1/(d_w - 1)}.
\end{equation}

One explanation is given here by the number of effective scales hidden in the estimates (26)–(28). Let us assume the model to be self similar [for all $k$, $r_k = \rho$ and $U_k = U$, $D(U_k) = \lambda$]. In Table 1 we have summarized formulae giving (in the first approximation in $1/\ln \rho$) the number of effective scales and the formulae linking them with those anomalous estimates (appearing in the proof, the influence of the perturbation scales will be neglected). This gives three values of $d_w$ corresponding to (26)–(28) and the interesting point is to compare them.

Let us observe that the multi-scale homogenization technique gives back the right forms for the mean squared displacement, the exit times and the
transition probability densities; they are explained by the number of scales which homogenization can be considered as complete associated with each observation. Moreover \( d_{w,1}, d_{w,2} \) and \( d_{w,3} \) are equal up to the first-order approximation in \( 1/\ln \rho \); nevertheless, they are not equal and this is not surprising. Indeed when \( \rho \) is small the second-order term in \( 1/(\ln \rho)^2 \) cannot be neglected since the perturbation scales becomes more and more dominant [and the influence of the perturbation scales is of the order of \( 1/(\ln \rho)^2 \)].

2.1.2. Strong overlap between the spatial scales. The anomaly is based on a minimal separation between spatial scales, that is, \( \rho_{\min} > \rho_0 \) and one might wonder what happens below this boundary. The answer will be given for a self similar case, that is, \( V \) is said to be self similar if for all \( n \), \( U_n = U \) and \( \rho_{\min} = \rho_{\max} = \rho \).

**Theorem 2.6.** If the potential \( V \) in (1) is self similar then, for all \( \rho \geq 2, \)
\[
E_0[\tau(0, r)] = r^{2 + \nu(r)}
\]
with
\[
\nu(r) = \frac{P_\rho(2U) + P_\rho(-2U)}{\ln \rho} + \epsilon(r)
\]
with \( \epsilon(r) \to 0 \) as \( r \to \infty. \)

Here \( P_\rho \) is the topological pressure associated with the shift operator \( s_\rho: x \in \mathbb{T} \to \rho x \in \mathbb{T} \) [see (129) for its definition].

Using the convexity properties of the topological pressure one has \( P_\rho(2U) + P_\rho(-2U) \geq 0 \) and the following proposition.

**Proposition 2.7.** \( P_\rho(2U) + P_\rho(-2U) = 0 \) if and only if
\[
\lim_{n \to \infty} \left\| \frac{1}{n} \sum_{k=0}^{n-1} \left( U(\rho^k x) - \int_{\mathbb{T}^d} U(x) \, dx \right) \right\|_\infty = 0.
\]
From this one deduces that for the simple example \( U(x) = \sin(x) - \sin(81x) \), \( \mathbb{E}[\tau(0, r)] \) is anomalous (subdiffusive \( \sim r^{2+\nu} \) with \( \nu > 0 \) for \( \rho \in \{2\} \cup \{4, \ldots, 26\} \cup \{28, \ldots, 80\} \cup \{82, \ldots, +\infty\} \) and normal \( \sim r^2 \) for \( \rho = 3, 27, 81 \). Thus if \( U \) is not a constant function, there exists \( \rho_0(K_0, K_1, D(U)) \) such that for \( \rho > \rho_0 \), \( y_t \) has a clear anomalous behavior (\( \mathbb{E}_0[\tau(0, r)] \sim r^{2+\nu} \) with \( \nu > 0 \)) but in the interval \( (1, \rho_0] \) both cases are possible: \( y_t \) may show a normal or an anomalous behavior according to the value of the ratio between scales \( \rho \) and the regions of normal behavior (characterized by Proposition 2.7) might be separated by regions of anomalous behavior.

What creates this phenomenon is a strong overlap or interaction between scales: that is, why the region \( (1, \rho_0] \) will be called “overlapping ratios,” that is, in this region the fluctuation of \( V \) at a size \( \xi > 0 \) is not represented by a single \( U_n(x/R_n) \) but by several ones and to characterize the behavior of \( y_t \) in that region one must introduce additional parameters describing the shapes of the fluctuations \( U_n \), elsewhere a normal or a subdiffusive behavior are both possible.

### 2.2. Davies’ conjecture and quantitative estimates on rate of convergence toward the limit process in homogenization.

The proof of Theorem 2.3 has not been described yet. The strategy is still to distinguish effective, perturbation and drift scales nevertheless it is not obvious to determine how many scales have been homogenized in the estimation of \( \mathbb{P}_0(y_t \geq h) \). The answer is directly linked with the rate at which the transition probability densities associated with a periodic elliptic operator do pass from a large deviation behavior to a homogenized behavior.

Consider for instance in any dimension \( d \geq 1 \), \( U \in \mathbb{L}^\infty(\mathbb{T}^d) \) and the Dirichlet form

\[
\mathcal{E}(f, f) = \frac{1}{2} \int_{\mathbb{R}^d} |\nabla f(x)|^2 \frac{e^{-2U(x)}}{\int_{\mathbb{T}^d} e^{-2U(z)} dz} dx, \quad f \in \mathcal{D}[\mathcal{E}] = H^1(\mathbb{R}^d).
\]

Write \( p(t, x, y) \) its associated heat kernel with respect to

\[
m_U(dx) = \frac{e^{-2U(x)}}{\int_{\mathbb{T}^d} e^{-2U(z)} dz} dx
\]

the invariant measure associated to (32). Note that when \( U \) is smooth the associated operator can be written \( L = 1/2\Delta - \nabla U \nabla \) and it is well known that:

- **Large deviation regime**: for \( |x - y| \gg t \) the paths of the diffusion concentrate on the geodesics and

\[
\ln p(t, x, y) \sim -\frac{|x - y|^2}{2t}.
\]
• **Heat kernel diagonal regime**: for $|x - y|^2 \ll t$, the behavior is fixed by the diagonal of the heat kernel and

\[
p(t, x, y) \sim \frac{C_0(x)}{t^{d/2}}.
\]

Davies conjectured that [we refer to Davies (1993); he considers periodic operators of divergence form nevertheless the idea remains unchanged] that $p(t, x, y)$ should have a homogenized behavior ($\ln p(t, x, y) \sim -(x - y)D(U)^{-1}(x - y)/(2t)$) for $t$ large enough.

Norris (1997) has shown that the homogenized behavior of the heat kernel $p(t, x, y)$ corresponding to a periodic operator on the torus $\mathbb{T}^d$ (dimension $d$ side 1) starts at least for $t \ln t \gg |x - y|^2$ (with $|x - y|^2 \ll t$); in this paper it will be shown that it starts for $t \gg |x - y|$ in any dimension.

This allows to complete the picture describing the behavior of $p(t, x, y)$:

• **Homogenization regime**: for $1 \ll |x - y| \ll t$ and $|x - y|^2 \gg t$, homogenization takes place and

\[
\ln p(t, x, y) \sim -|x - y|_{D^{-1}(U)}^2/(2t)
\]

with

\[
|x - y|_{D^{-1}(U)}^2 := (x - y)D(U)^{-1}(x - y).
\]

More precisely we will prove the following:

**Theorem 2.8.** Consider $p(t, x, y)$ the heat kernel associated with the Dirichlet form (32) with respect to the measure $m_U$. Then there exist constants $C, C_2$ depending only on $d$ and $\text{Osc}(U)$ such that for

\[
C|x - y| < t, \quad C\sqrt{t} < |x - y|, \quad C < |x - y|,
\]

one has

\[
p(t, x, y) \leq \frac{1}{(2\pi t)^{d/2}(\det(D(U)))^{1/2}} \exp(-(1 - E)|y - x|_{D^{-1}(U)}^2/(2t)),
\]

\[
p(t, x, y) \geq \frac{1}{(2\pi t)^{d/2}(\det(D(U)))^{1/2}} \exp(-(1 + E)|y - x|_{D^{-1}(U)}^2/(2t)),
\]

with

\[
E(t, x, y) := C_2 \left( \frac{|x - y|}{t} + \frac{\sqrt{t}}{|x - y|} \right) \leq \frac{1}{10}.
\]

Theorem 2.8 proves Davies’ conjecture, moreover $E(t, x, y)$ acts as a quantitative error term putting into evidence the rate at which homogenization takes place for the heat kernel, and it also acts as the inverse of a distance from the domains
associated with the large deviation regime and the heat kernel diagonal regime. Observe that all the constants do depend only on $d$ and $\text{Osc}(U)$. It is straightforward to extend those estimates to any periodic elliptic operator. They can be likened to results obtained by Dembo (1996) for discrete martingales with bounded jumps based on moderate deviations techniques.

2.2.1. A note on the proof of Theorem 2.3. Those estimates basically say that the homogenized behavior of the heat kernel associated with a periodic medium of period $R$ starts for $t > R|x - y|$. Thus in the proof of Theorem 2.3 the number of the smaller scales that can be considered as homogenized is fixed by $n_{ef}(t/h) = \sup_n \{R_n \leq t/h\}$, which [assume $D(U_n) = \lambda$ and $R_n = \rho^n$ for simplification] leads to an anomaly of the form

$$\ln \mathbb{P}(y_t \geq h) \leq -C \frac{h^2}{t^{\lambda n_{ef}(t/h)}}$$

with $d_w \sim 2 - \frac{\ln \lambda}{\ln \rho}$. Equation (42) suggests that the origin of the anomalous shape of the heat kernel for the reflected Brownian motion on the Sierpinski carpet can be explained by the formula linking the number of effective scales and the ratio $t/h$.

The first condition in (16) can be translated into “homogenization has started on at least the first scale” and the second one into “the heat kernel associated with (1) is far from its diagonal regime” (one can have $h^2/t \ll 1$ before reaching that regime; this is explained by the slow down of the diffusion).

2.2.2. A quantitative inequality for exponential martingales. The core of the proofs of Theorems 2.3 and 2.8 is an inequality giving a quantitative estimate for the Laplace transform of a martingale:

Consider $M_t$ a continuous square integrable $\mathcal{F}_t$ adapted martingale such that $M_0 = 0$ and for $\lambda \in \mathbb{R}$, $\mathbb{E}[e^{\lambda M_t}] < \infty$. Assume that there exists a function $f : \mathbb{R}^+ \rightarrow \mathbb{R}^+$ such that for all $t_2 > t_1 \geq 0$ one has, a.s.,

$$\mathbb{E} \left[ \int_{t_1}^{t_2} d\langle M, M \rangle_s | \mathcal{F}_{t_1} \right] \leq \int_{0}^{t_2-t_1} f(s) \, ds$$

with $f(s) = f_1$ for $s < t_0$ and $f(s) = f_2$ for $s \geq t_0$ with $t_0 > 0$ and $0 < f_2 < f_1$.

**Theorem 2.9.** Let $M_t$ be the martingale described above.

1. For all $0 < |\lambda| < (2e(f_1 - f_2)t_0)^{-1/2}$ one has

$$\mathbb{E}[\exp(\lambda M_t)] \leq e^{3(1-1/g(\lambda))^2} \exp \left( \frac{g(\lambda)}{\lambda} f_2 t \right)$$

with $g(\lambda) = \frac{1}{1-\lambda^2(f_1-f_2)\ln e}$ that verifies $1 \leq g \leq 2$. 

2. For all $0 < v < (2e(f_1 - f_2)t_0)^{-1}$ one has

\[ \mathbb{E}[\exp(v\langle M, M\rangle_t)] \leq \exp(vf_2t)\frac{\exp(v_0(f_1 - f_2))}{((f_1 - f_2)v_0)^2} . \]  

This theorem uses the knowledge on the conditional behavior of the quadratic variation of a martingale to upper bound its Laplace transform, and it is well known that a quantitative control on the Laplace transform leads to a quantitative control on the heat kernel tail. The condition $\lambda$ small enough marks the boundary between the large deviation regime and the homogenization regime. A direct application of the key theorem is the following result.

**Corollary 2.10.** Let $M_t$ be the martingale given in Theorem 2.9. Write $C_1 = (2e(f_1 - f_2)t_0)^{1/2}/f_2$. For $r = \frac{C_1}{t} < 1$ one has

\[ \mathbb{P}(M_t \geq x) \leq e^{3/2r^2}\exp\left(-(1 - r^2)\frac{x^2}{2f_2t}\right). \]  

This corollary gives a quantitative control on the tail of the law of $M_t$ from the asymptotic behavior of its conditional brackets.

2.3. **An analytical inequality for subharmonic functions.** The stability property (22) is based on the following analytical inequality:

**Theorem 2.11.** Let $\Omega$ be an open bounded subset of $\mathbb{R}$, $d = 1$, for $\lambda \in C^\infty(\overline{\Omega})$ such that $\lambda > 0$ on $\overline{\Omega}$ and $\phi, \psi \in C^2(\overline{\Omega})$ null on $\partial\Omega$ and both subharmonic with respect to the operator $-\nabla(\lambda\nabla)$, one has

\[ \int_{\Omega} \lambda(x)|\nabla\phi(x) \cdot \nabla\psi(x)| \, dx \leq 3 \int_{\Omega} \lambda(x)|\nabla\phi(x) \cdot \nabla\psi(x)| \, dx. \]  

The constant 3 in this theorem is the optimal one. We believe that this inequality might also be true in higher dimensions, that is, we have the following conjecture.

**Conjecture 2.12.** For $\Omega \subset \mathbb{R}^d$ an open subset with smooth boundary, there exist a constant $C_{d,\Omega}$ depending only on the dimension of the space and the open set such that for $\lambda \in C^\infty(\overline{\Omega})$ such that $\lambda > 0$ on $\overline{\Omega}$ and $\phi, \psi \in C^2(\overline{\Omega})$ null on $\partial\Omega$ and both subharmonic with respect to the operator $-\nabla(\lambda\nabla)$, one has

\[ \int_{\Omega} \lambda(x)|\nabla\phi(x) \cdot \nabla\psi(x)| \, dx \leq C_{d,\Omega} \int_{\Omega} \lambda(x)|\nabla\phi(x) \cdot \nabla\psi(x)| \, dx. \]  

This conjecture is equivalent to the stability of the Green functions of divergence form elliptic operators under a deformation. More precisely write $G_\lambda$ the Green function associated with $-\nabla(\lambda\nabla)$ with Dirichlet conditions on $\partial\Omega$. 


Proposition 2.13. Conjecture 2.12 is true with the constant $C_{d, \Omega}$ if and only if for all $\lambda, \mu$ bounded and strictly positive on $\Omega$,

\[
(49) \quad \left( \sup_{\Omega} \max \left( \frac{\mu}{\lambda}, \frac{\lambda}{\mu} \right) \right)^{-C_{d, \Omega}} \leq \frac{G_\mu(x, y)}{G_\lambda(x, y)} \leq \left( \sup_{\Omega} \max \left( \frac{\mu}{\lambda}, \frac{\lambda}{\mu} \right) \right)^{C_{d, \Omega}}.
\]

Remark 2.14. It would be interesting to prove this proposition since it allows to obtain sharp quantitative estimates on the comparison of elliptic operators with non-Laplacian principal part. By Proposition 2.13 it is easy to check that Conjecture 2.12 implies the Harnack inequality. One might think that one would be able to obtain (49) using Aronson’s estimates and keeping track of the dependence of the constants in the Harnack inequality, but this is not the case since Harnack inequality is an isotropic inequality and (49) compares in an optimal way Green functions of operators which can be strongly anisotropic.

Let us recall that the Harnack inequality associated with the operator $L = -\nabla \lambda \nabla$ says that for all $L$-harmonic functions $u$ in $B(0, r)$ one has

\[
\sup_{x \in B(0, r/2)} u(x) \leq C_L \inf_{x \in B(0, r/2)} u(x),
\]

where the optimal constant $C_L$ grows toward infinity as $\sup \lambda / \inf \lambda \to \infty$ whereas the constant associated with Conjecture 2.12 is independent of $\lambda$. That is why the Harnack inequality strategy, which has already been used to obtain quantitative results for the comparison with the Laplace operator [we refer to Stampacchia (1965), Ancona (1997), Gruter and Widman (1982) and Pinchover (1989)] allows us to obtain

\[
(50) \quad \frac{G_\lambda(x, y)}{G_0(x, y)} \leq C_H
\]

but with a constant $C_H$ exploding like $C_d \exp(C_d (\sup \lambda / \inf \lambda)^C_d)$.

Remark 2.15. Since the conjecture is true in dimension one with $C_{d, \Omega} = 3$ (this constant is an homotopy invariant), it is through Proposition 2.13 that one obtains stability property (22).

Remark 2.16. It is easy to deduce from Theorem 2.11 that if $\Omega$ is a bounded open subset of $\mathbb{R}^d$ and $\phi, \psi$ are both convex or both concave functions on $\Omega$ and null on $\partial \Omega$, then

\[
(51) \quad \int_{\Omega} |\nabla x \phi(x) \cdot \nabla x \psi(x)| \ dx \leq 3 \int_{\Omega} \nabla x \phi(x) \cdot \nabla x \psi(x) \ dx.
\]

Remark 2.17. Conjecture 2.12 (Theorem 2.11 when $d = 1$) has an interesting signification (and consequences) in the framework of electrostatic theory, we refer to Chapter 13 of Owhadi (2001a).
2.4. Remark: fast separation between scales. The feature that distinguishes a strong slow behavior from a weak one is the rate at which spatial scales do separate. Indeed one can follow the proofs given above, changing the condition $\rho_{\text{max}} < \infty$ into $\mathbb{R}_n = [0, \rho^{\alpha}]$ ($\rho, \alpha > 1$) and $\lambda_{\text{max}} = \lambda_{\text{min}} = \lambda < 1$ to obtain:

- A weak slow behavior of the exit times
  \[
  C_1 r^2 e^{g(r)} \leq \mathbb{E}_0[\tau(0, r)] \leq C_2 r^2 e^{g(r)}
  \]
  with $g(r) = (\ln r)^{1/\alpha}(\ln 1/\lambda)(\ln \rho)^{-1/\alpha}$.

- A weak slow behavior of the mean squared displacement
  \[
  C_1 t e^{-f(t)} \leq \mathbb{E}_0[y_t^2] \leq C_2 t e^{-f(t)}
  \]
  with $f(t) = (\ln r)^{1/\alpha}(\ln 1/\lambda)(2 \ln \rho)^{-1/\alpha}(1 + \varepsilon(t))$.

- A weak slow behavior of the heat kernel tail: for $h > 0$, $C_1 < t/h < C_2 h$
  \[
  \mathbb{P}[y_t \geq h] \leq C_3 e^{-C_4(h^2/\mu)k(t/h)}
  \]
  with $k(x) = -((x/(\ln \rho))^{1/\alpha}(1+\varepsilon(x))$.

And as $\alpha \downarrow 1$, the behavior of the solution of (1) passes from weakly anomalous to strongly anomalous.

3. Proofs.

3.1. Davies’ conjecture and quantitative estimates on rate of convergence toward the limit process in homogenization.

3.1.1. Quantitative control of the Laplace transform of a martingale: Theorem 2.9. The core of the proof of the anomalous heat kernel tail (Theorem 2.3) and the quantitative estimates on the heat kernel associated with an elliptic operator (Theorem 2.8) is Theorem 2.9 that will be proven in this section.

Let $M_t$ be the martingale described in Theorem 2.9. Let $q > 1$. Using the Hölder inequality and Itô’s formula it is easy to obtain that, with $h_q = \frac{q^2}{2(q-1)}$,

\[
\mathbb{E}[\exp(h_q \lambda^2 s_t^2)] \leq \mathbb{E}[\exp(h_q \lambda^2 s_t^2)]^{1/q}.
\]

Thus the quantitative control of the Laplace transform of the martingale shall follow from this control on its bracket.

Write $\mu = \frac{t}{t_0}$ ($[\mu]$ shall stand for the integer part of $\mu$). Using the Hölder inequality and the control (43), one obtains, for $1 < z < \infty$,

\[
\mathbb{E}[\exp(h_q \lambda^2 |s_t^2 - [\mu]s_{t_0}|)]^{1/q} \leq \mathbb{E}[\exp\left(z h_q \lambda^2 |s_t^2 - [\mu]s_{t_0}|\right)]^{1/(zq)} \exp\left((h_q/q)\lambda^2 (t - [\mu]t_0) f_1\right).
\]
Then by taking the limit $z \downarrow 1$, one easily obtains that
\[
\mathbb{E}[\exp(h_q \lambda^2 \langle M, M \rangle_t)]^{1/q} \leq \mathbb{E}[\exp(h_q \lambda^2 \langle M, M \rangle_{[\mu]t_0})]^{1/q} \exp((h_q/q)\lambda^2(t - [\mu]t_0)f_1).
\]
Write $a = f_2/f_1$. We will need the following lemma:

**Lemma 3.1.** Let $M_t$ be the martingale described in Theorem 2.9 and $\eta > 0$; for $a = f_2/f_1$ and $\mu = t/t_0$ one has
\[
\mathbb{E}[\exp(\eta \langle M, M \rangle_t)] \leq 1 + \sum_{n=1}^{+\infty} \frac{(\eta f_1 t_0)^n}{n!} \sum_{0 \leq m \leq n \wedge \mu} (\mu - m)^m C_n^m (a - 1)^m.
\]

**Proof.** By the Taylor expansion of the exponent one obtains
\[
\exp(\eta \langle M, M \rangle_t) = 1 + \sum_{n=1}^{+\infty} \eta^n W_n
\]
with $W_n = \int_{u_i > 0} 1(0 < u_1 + \cdots + u_n < t) f(u_1) \cdots f(u_n) du_1 \cdots du_n$. Using the control (43) on the conditional brackets of the martingale it is easy to obtain by induction on the integrand and the Markov property that
\[
\mathbb{E}[W_n] \leq \int_{u_i > 0} 1(0 < u_1 + \cdots + u_n < t) f(u_1) \cdots f(u_n) du_1 \cdots du_n.
\]
Combining this with (59), and using the fact that $f(s) \leq f_1 g(s/t_0)$ with $g(z) = 1(z < 1) + a 1(z \geq 1)$, one obtains that
\[
\mathbb{E}[\exp(\eta \langle M, M \rangle_t)] \leq 1 + \sum_{n=1}^{+\infty} (\eta f_1 t_0)^n G_n
\]
with $G_n = \int_{z_i > 0} 1(0 < z_1 + \cdots + z_n < \mu) \prod_{k=1}^{n} (1(z_k < 1) + a 1(z_k \geq 1)) dz_1 \cdots dz_n$. Developing the product in $G_n$ one obtains by integration, induction and straightforward combinatorial computation that
\[
G_n = \frac{1}{n!} \sum_{0 \leq m \leq \mu \wedge n} C_n^m (\mu - m)^m (a - 1)^m.
\]
Which leads to (58) by inequality (60). □

Using Lemma 3.1 one obtains
\[
\mathbb{E}[\exp(h_q \lambda^2 \langle M, M \rangle_{[\mu]t_0})] \leq \sum_{n=0}^{+\infty} \frac{(h_q \lambda^2 f_1 t_0)^n}{n!} \sum_{0 \leq m \leq n \wedge [\mu]} ([\mu] - m)^m C_n^m (a - 1)^m.
\]
Changing the order of summation, one obtains
\[
\mathbb{E} \left[ \exp \left( hq \lambda^2 \langle M, M \rangle_{[\mu]} \right) \right] \\
\leq \exp(hq \lambda^2 f_1 t_0[\mu]) \sum_{0 \leq m \leq [\mu]} ([\mu] - m m! (h q (a - 1) \lambda^2 f_1 t_0)^m).
\]

Now we will need the following lemma:

**Lemma 3.2.** For \(-1/e < y < 0\),
\[
\sum_{0 \leq m \leq [\mu]} \frac{([\mu] - m m! y^m)}{m!} \leq \frac{\exp(y[\mu])}{y^2}.
\]

**Proof.** Put \(-1/e < x < 0\) and write, for \(n \in \mathbb{N}\),
\[
I_n = \sum_{0 \leq m \leq n} \frac{x^m}{m!} (n - m)^m.
\]
It will be shown here that \(\forall p \in \mathbb{N}^* , \forall n \in \mathbb{N}\),
\[
I_n \leq (u_p(x))^{-n} (1 - u_p \exp(x u_p))^{-1}
\]
where \(u_p\) the increasing sequence defined by \(u_0 = 0\) and \(u_{p+1} = \exp(-x u_p)\) and converging to \(y_0\) the smallest positive solution of \(y \exp(xy) = 1\).

Inequality (62) is then obtained for \(u_p(y) = u_2(y) = \exp(-y)\) and using \(\exp(-y) - 1 \geq -y\) and \(-1/e < y < 0\).

Write \(y_1 = \inf\{y > 0 : y \exp(|x|y) = 1\}\) (note that \(0 < y_1 < 1\)) and consider for \(-y_1 < y < y_1\) the function
\[
f : y \rightarrow (1 - y \exp(xy))^{-1}.
\]
By Taylor’s expansion, for \(y \in (-y_1, y_1)\),
\[
f(y) = \sum_{n=0}^{+\infty} \frac{(n x y)^n}{n!} \sum_{m=0}^{+\infty} \frac{(n x y)^m}{m!} = \frac{1}{1 - y \exp(|x|y)} < \infty
\]
with a normal convergence of the series, the order of the limits can be changed, which leads to
\[
f(y) = \sum_{m=0}^{+\infty} \frac{(n x y)^m}{m!} \sum_{n=0}^{+\infty} n^m y^n = \sum_{m=0}^{+\infty} \frac{x^m}{m!} \sum_{n=m}^{+\infty} (n - m)^m y^n
\]
\[
= \sum_{n=0}^{+\infty} y^n \sum_{m=0}^{n} (n - m)^m x^m \frac{1}{m!} = \sum_{n=0}^{+\infty} y^n I_n.
\]
It follows that \(\forall n \in \mathbb{N}, I_n = \frac{f^{(n)}(0)}{n!}\). Now, for \(-1/e < x < 0\), the constant \(y_0 = \inf\{y > 0 : y \exp(xy) = 1\}\) does exist and \(\forall y \in ]-y_1, y_0[ , \forall n, f^{(n)}(y) \geq 0\) (thus \(I_n \geq 0\)).
Thus from the classical theorem of Taylor’s expansion, the series \( \sum_{n=0}^{+\infty} y^n \frac{f^{(n)}(0)}{n!} \) converges toward \( f \) for \( y \in ]-y_1, y_0[ \) and in that interval

\[
\sum_{n=0}^{\infty} y^n I_n = (1 - y \exp(xy))^{-1}
\]

from which one deduces that \( \forall y \in ]0, y_0[ , \forall n \in \mathbb{N} , I_n \leq y^{-n}(1 - y \exp(xy))^{-1} \).

On the other hand, if one considers the sequence \( u_0 = 0, u_{p+1} = \exp(-x u_p) \), then it is an exercise to show that \( u_p \) is increasing and will converge toward \( y_0 \), which leads to (63). \( \square \)

Applying (62) to (61) with \( y = h_q(a - 1)\lambda^2 f_1 t_0 \) one obtains that, for \( 0 < |\lambda| < (eh_q(f_1 - f_2)t_0)^{-1/2} \), one has

\[
\mathbb{E}[\exp(h_q \lambda^2 (M, M)_{[\mu]_0})]^{1/q} \leq \exp\left(\frac{h_q \lambda^2 f_2 t_0[\mu]}{q}\right)(h_q(1 - a)\lambda^2 f_1 t_0)^{-2/q}.
\]

Writing \( \nu = \lambda^2 h_q \) and combining (64) with (57), one obtains inequality (45) of Theorem 2.9.

Combining (64) with (57) and (55) one obtains inequality (44) of Theorem 2.9 by choosing \( q = (\lambda^2(f_1 - f_2)t_0 e)^{-1} (q > 2 \) under the condition imposed on \( \lambda \).

3.1.2. Upper bound estimate (39) of Theorem 2.8. Theorem 2.9 can be used to give quantitative estimates on any operator as soon as a cell problem is well defined. Consider \( y_t \) as a diffusion on \( \mathbb{R}^d \) that may be decomposed for \( t > 0 \) as

\[
y_t = x + \chi(t) + M_t
\]

where \( \chi(t) \) is a uniformly (in \( t \)) bounded random vector process (\( \|\chi\|_{\infty} \leq C_\chi \)) and \( M_t \) is a continuous square integrable \( \mathcal{F}_t \) adapted martingale such that \( M_0 = 0 \).

Assume that for all \( l \in \mathbb{R}^d \) with \( ||l|| = 1 \) there exists a function \( f : \mathbb{R}^+ \to \mathbb{R}^+ \) such that for all \( t_2 > t_1 \geq 0 \) one has, a.s.,

\[
\mathbb{E}\left[ \int_{t_1}^{t_2} d(M \cdot l, M \cdot l)_s |\mathcal{F}_{t_1} \right] \leq \int_0^{t_2-t_1} f(s) \, ds
\]

with \( f(s) = f_1 \) for \( s < t_0 \) and \( f(s) = t^d Dl < f_1 \) for \( s \geq t_0 \) with \( t_0 > 0 \) where \( D \) is a positive definite symmetric matrix.

Assume that the diffusion \( y_t \) has symmetric Markovian probability densities \( p(t, x, y) \) with respect to the measure \( m(dy) \) such that, for all \( x, y \in \mathbb{R}^d \) and \( t > 0 \),

\[
p(t, x, y) \leq \frac{C_2}{t^{d/2}}
\]
and for $\delta > 0$,
\begin{equation}
\mathbb{P}_x(|y_t - x| \geq \delta) \leq C_3 e^{-C_4 \delta^2 / t}
\end{equation}
where $C_2, C_3, C_4$ are constants.

**Theorem 3.3.** Let $y_t$ be the diffusion described above. Then with $k_1 = 30(e(f_1 - \lambda_{\min}(D))t_0)^{1/2}/\lambda_{\min}(D)$ and $k_2 = 30 + 10d\lambda_{\max}(D)(1 + C_4)$,
\begin{equation}
k_1 |x - y| < t, \quad k_2 < \frac{|x - y|}{\sqrt{t}}, \quad |x - y| > 4C\chi,
\end{equation}
one has
\begin{equation}
p(t, x, y) \leq \frac{E_1}{t^{d/2}} \exp \left(-\left(1 - E\right) \frac{|y - x|^2}{2t^{d-1}}\right)
\end{equation}
with $E_1 = C_2(5(\lambda_{\min}(D)C_4)^{-1} + 2^d C_3)$ and $E = 3((\frac{E|x - y|}{t})^2 + \frac{\sqrt{t}}{|x - y|}) \leq \frac{1}{10}$.

**Proof.** The estimate on the heat kernel $p(t, x, y)$ will follow from the chain rule and decomposing it the probability of moving away from $x$ to “a well chosen set containing $y$ in the time $tq$” and its complement. More precisely, writing $e_{y - x} := (y - x)/|y - x|$ and $A_\delta = \{ z \in \mathbb{R}^d : (z - x) \cdot e_{y - x} \geq (1 - \delta)|x - y|\}$, using (67) one obtains that, for $t > 0, x, y \in \mathbb{R}^d$ and $0 < q < 1$,
\begin{align}
p(t, x, y) &= \int_{A_\delta} p(tq, x, z) p(t(1 - q), z, y) m(dz) \\
&\quad + \int_{A_\delta^c} p(tq, x, z) p(t(1 - q), z, y) m(dz) \\
&\leq \frac{C_2}{t^{d/2}} \left[ \frac{1}{(1 - q)^{d/2}} \mathbb{P}_x(y_{tq} \cdot e_{y - x} \geq |x - y|(1 - \delta)) \right. \\
&\quad \left. + \frac{1}{q^{d/2}} \mathbb{P}_y(|y_t(1 - q)| \geq \delta |x - y|) \right].
\end{align}
Let us choose $\delta = \exp(-|x - y|(dD(e_{x - y})/\sqrt{t})^{-1})$ and $q = 1 - 2D(e_{x - y})C_4\delta$ [we will use the notation $D(l) := \|lD\|$].

For $|x - y|/\sqrt{t} > \max(dD(e_{x - y}) \ln(4D(e_{x - y})C_4), 3dD(e_{x - y}))$ (which basically says that the heat kernel is far from its diagonal behavior) one has $\delta < \frac{1}{10}$ and $\frac{1}{2} < q < 1$. Using the Aronson type estimate (68) one controls the second term in (71):
\begin{equation}
\mathbb{P}_y(|y_{t(1 - q)}| \geq \delta |x - y|) \leq C_3 \exp\left(-\frac{|x - y|^2}{2D(e_{x - y})t}\right).
\end{equation}
By properties (65), (66) and Corollary 2.10, one controls the first term in (71): for $r < 1$ with $r = \frac{C_1 \rho}{\alpha^2}$, $\rho = |x - y|(1 - \delta) - C_\chi$ and $C_1 = (2e(f_1 - D(e_{x-y}))t_j)^{1/2}/D(e_{x-y})$, one has

$$
(73) \quad \mathbb{P}_x(y_t \cdot e_{x-y} \geq |x - y|(1 - \delta)) \leq e^{(3/2)r^2} \exp\left(-\frac{\rho^2}{2D(e_{x-y})t_j}\right).
$$

Combining (73), (72), (71) and using the value of $q$ and $\delta$ given above, one easily obtains estimate (70) of Theorem 3.3 under Conditions (69).

Now Theorem 2.8 is a straightforward application of Theorem 3.3 and a trivial adaptation of the constants appearing in Theorem 3.3. Consider $p(t, x, y)$ the heat kernel associated with the Dirichlet form (32). Since $p(t, x, y)$ is continuous in $L^\infty(\mathbb{R}^d)$ norm with respect to $U$ [we refer to Chen, Qian, Hu and Zheng (1998) whose result can easily be adapted to our case] and $C^\infty(\mathbb{R}^d)$ is dense in $L^\infty(\mathbb{R}^d)$ with respect to that norm, one can assume $U$ to be smooth and the general result follows by observing that the estimates in Theorem 3.3 depend only on $\text{Osc}(U)$.

By definition $y_t$ has symmetric probability densities with respect to the measure $m_U$ and the following Aronson-type upper bound is available [Seignourel (1998)]:

$$
(74) \quad p(t, x, y)e^{-2U(y)} \leq Ce^{(4+d)\text{Osc}(U)} \frac{1}{t^{d/2}} \exp\left(-\frac{|x - y|^2}{4t}\right).
$$

It follows that conditions (67) and (68) are satisfied with constants $C_2, C_3$ and $C_4$ depending only on $d$ and $\text{Osc}(U)$. Now write $\chi_l$ the solution of the associated cell problem: for $l \in S^d$, $L_U \chi_l = -l \cdot \nabla U$ with $\chi(0) = 0$.

Using Stampacchia [(1965), Theorem 5.4, Chapter 5] on elliptic equations with discontinuous coefficients [see also Stampacchia (1966)] and using the periodicity of $\chi$ and observing that $\chi_l(x) = l \cdot x - F_l(x)$ where $F_l$ is harmonic with respect to $L_U$, one easily obtains that

$$
(75) \quad C_\chi = \|\chi_l\|_\infty \leq C_d \exp((3d + 2)\text{Osc}(U)).
$$

From Itô’s formula one has $l \cdot y_t = x + \chi_l(y_t) - \chi_l(x) + \int_0^t (l - \nabla \chi_l) d\omega_s$, which corresponds to the decomposition given in (65). The martingale can be written $l \cdot M_t = \int_0^t (l - \nabla \chi_l) d\omega_s$ and its bracket is equal to $\langle l \cdot M, l \cdot M \rangle_t = \int_0^t |l - \nabla \chi_l(y_s)|^2 ds$. It is easy to obtain from Theorem 3.9 of Gilbarg and Trudinger (1983) that

$$
(76) \quad f_1 = \|\nabla \chi_l\|_\infty \leq C_d(1 + \|\nabla U\|_\infty) \exp((3d + 2)\text{Osc}(U)) < \infty.
$$

Writing $\phi_l$ the periodic solution of the ergodicity problem $L_U = |l \nabla \chi_l|^2 - l^1 D(U)l$, $\phi_l(0) = 0$, and observing that $\phi_l = F_l^2 - l^1 D(U)l \phi_l$ where $L_U \phi_l = 1$
it is easy to obtain from (75), Theorem 5.4, Chapter 5 of Stampacchia (1965) and the periodicity of $\phi_l$ that

\begin{equation}
C_\phi = \|\phi_l\|_\infty \leq C_d \exp\left((9d + 4)\text{Osc}(U)\right).
\end{equation}

Since, from Itô’s formula

\begin{equation}
E_x[\langle l \cdot M_l, l \cdot M_l \rangle_t] = E[\phi(y_t) - \phi(x)] + t'tD(U)l,
\end{equation}

the martingale satisfies the conditions of Theorem 3.3 with $f_2 = t'D(U)l$ and $t_0 = C_\phi/(f_1 - \lambda_{\min}(D))$. Now one can use Theorem 3.3 to obtain a quantitative control on the heat kernel. It is important to note that all the constants appearing in that theorem only depend on $d$ and $\text{Osc}(U)$ except maybe $k_{10} = 30(e(f_1 - \lambda_{\min}(D))t_0)^{1/2}/\lambda_{\min}(D)$ in which $f_1$ appears. This is where the trick operates, indeed $(f_1 - \lambda_{\min}(D))t_0 = C_\phi$ which is a constant depending only on $\text{Osc}(U)$ and $d$. Thus in reality all the constants only depend on the dimension and on $\text{Osc}(U)$, which proves the upper bound in Theorem 2.8.

3.1.3. Lower bound estimate (40) of Theorem 2.8. Let $y_t$ the diffusion associated with the Dirichlet form (32). As has been done in Section 3.1.2 one can prove estimate (40) assuming that $U$ is smooth and the general case will follow by the continuity of the heat kernel with respect to $U$ in $L^\infty(\mathbb{R}^d)$ norm.

First, we will need the following estimate.

PROPOSITION 3.4. For $l \in \mathbb{S}^d$, $\lambda > k_{5,d,\text{Osc}(U)}$ and $k_{6,d,\text{Osc}(U)} \lambda < t$ one has

\begin{equation}
\mathbb{P}[y^U_t \cdot l \geq \lambda] \geq \frac{1}{4\sqrt{2\pi}} \int_X e^{-z^2/2} \, dz
\end{equation}

with $X = \frac{\lambda}{\sqrt{t'D(U)lt}}(1 + F)$ and $F = \frac{k_{7,d,\text{Osc}(U)} \lambda}{\lambda} + k_{8,d,\text{Osc}(U)}\sqrt{2} t \leq \frac{1}{10}$.

PROOF. For $l \in \mathbb{S}^d$, let $F_l$, $\chi_l$, $\phi_l$ be the functions introduced in Section 3.1.2. Write $\mathcal{F}_t$ the filtration associated with Brownian motion appearing in the SDE solved by $y_t$. $F_l(y_t)$ is a $(\mathbb{P}, \mathcal{F}_t)$-continuous local martingale vanishing at 0 such that (by Itô calculus)

\begin{equation}
\langle F_l, F_l \rangle_t = tD(l) + \phi_l(y_t) + M_t
\end{equation}

with $M_t = -\int_0^t \nabla \phi_l(y_s) \, d\omega_s$. Since $\langle F_l, F_l \rangle_\infty = \infty$ a.s. by the Dambis, Dubins–Schwarz representation theorem, $B_t = F_l(y_{T_t})$ is a $(\mathcal{F}_t)$-Brownian motion with $F_l(y_t) = B_{\langle F_l, F_l \rangle_t}$ and

\begin{equation}
T_t = \inf\{s : \langle F_l, F_l \rangle_s > t\}.
\end{equation}

The idea of the proof is then to show that the probability of $y_t$ to move away from 0 behaves like the probability of a BM of variance $D(l)$ to move away. To achieve this it will be sufficient to show that $M_t$ becomes negligible in front of $tD(l)$ using Corollary 2.10 to control $\mathbb{P}(M_t \geq x)$. More precisely we will use the following lemma.
**Lemma 3.5.** For
\[ \lambda > 0, \nu > \| \phi_l \|_\infty, \mu > 0, \quad \lambda + \| \chi_l \|_\infty + \mu \leq C_2 \mu \sqrt{D(l) t \nu^{-1}}, \]
one has
\[ \mathbb{P}[y_t \cdot l \geq \lambda] \geq \frac{1}{2} \mathbb{P}[B_D(l)t \geq \lambda + \| \chi_l \|_\infty + \mu] - \mathbb{P}[|M_t| \geq \nu - \| \phi_l \|_\infty]. \]

**Proof.** Let \( \lambda > 0 \) from the representation theorem \( \mathbb{P}[F_l(y_t) \geq \lambda] = \mathbb{P}[B_D(l)t + E_t \geq \lambda] \) with \( E_t = B_{\langle F_l, F_l \rangle} - B_D(l)t \). It follows that, for \( \mu > 0, \)
\[ \mathbb{P}[F_l(y_t) \geq \lambda] \geq \mathbb{P}[B_D(l)t \geq \lambda + \mu] - \mathbb{P}[|E_t| > \mu]. \]

It follows from (80) that, for \( \nu > 0, \)
\[ \mathbb{P}[|E_t| \geq \mu] \leq \mathbb{P}[|\phi_l(y_t) + M_t| \geq \nu] + \mathbb{P}[\sup_{|z| < \nu} |B_D(l)t + z - B_D(l)|t \geq \mu], \]
from which one deduces
\[ \mathbb{P}[|E_t| \geq \mu] \leq \mathbb{P}[|M_t| \geq \nu - \| \phi_l \|_\infty] + 2 \mathbb{P}[|B_v| > \mu]. \]

Combining (84) and (85) one obtains that \( \nu > \| \phi_l \|_\infty: \)
\[ \mathbb{P}[y_t \cdot l \geq \lambda] \geq \mathbb{P}[B_D(l)t \geq \lambda + \| \chi_l \|_\infty + \mu] - 4 \mathbb{P}[B_D(l)t \geq \mu \sqrt{\frac{D(l)t}{\nu}}] - \mathbb{P}[|M_t| \geq \nu - \| \phi_l \|_\infty], \]
which leads to (83) under the last condition in (82). \( \square \)

Now let us show the following.

**Lemma 3.6.** For \( C_M x < t \) one has
\[ \mathbb{P}(M_t \geq x) \leq 3 \exp\left( -\frac{x^2}{f_2 l} \right) \]
where \( f_2 \) and \( C_M \) depend only on \( d \) and \( \text{Osc}(U) \).

**Proof.** Write \( G(x) = \frac{1}{2} \phi_l^2 - \| \phi_l \|_\infty \phi_l \). Since
\[ L_U G(x) = |\nabla \phi_l|^2 - (\| \phi_l \|_\infty - \phi_l)(|\nabla F_l|^2 - D(l)) \]
one obtains from Itô’s formula that
\[ \mathbb{E}[\langle M, M \rangle_t] \leq 2 \| \phi_l \|_\infty \mathbb{E}\left[ \int_0^t |\nabla F_l|^2(y_s) ds + D(l)t \right] + \| G \|_\infty \]
\[ \leq 2 \| \phi_l \|_\infty (\| \phi_l \|_\infty + 2D(l)t) + 2 \| \phi_l \|_\infty^2. \]
Thus \( M_t \) satisfies the conditions of Corollary 2.10 with \( f_2 = 4 \| \phi_l \|_\infty D(l), \)
\( f_1 = |\nabla \phi_l|^2_\infty \) and \( t_0 = 4 \| \phi_l \|_\infty^2/(f_1 - f_2), \)
which leads to (86) by observing that \((f_1 - f_2)t_0)^{1/2}/f_2\) is upper bounded by a constant depending only on \( \text{Osc}(U) \) and \( d \). \( \square \)
It follows from equation (83) that under the additional conditions
(88) \[ C_M (v - \| \phi \|_\infty) < t \quad \text{and} \quad \lambda + \| \chi \|_\infty + \mu < C_3 (v - \| \phi \|_\infty) \]
where \( C_3 \) depends only on \( d \) and \( \text{Osc} (U) \), one has
(89) \[ P[y_l \cdot l \geq \lambda] \geq \frac{1}{4} P[B_{D(l)t} \geq \lambda + \| \chi \|_\infty + \mu]. \]
Choosing \( v = \| \phi \|_\infty + 2/C_3 (\lambda + \| \chi \|_\infty + \mu) \) and
\[ \mu = 4(\lambda + \| \chi \|_\infty) \frac{3}{2} \left( \frac{C_2 \sqrt{D(l)/C_3}}{t} \right)^{-1} \]
for \( \lambda > \| \chi \|_\infty \) and \( t > C_4 (d, \text{Osc} (U), \lambda) \), conditions (82) and (88) are satisfied and
(90) \[ \mu < C_5 (d, \text{Osc} (U)) \lambda \sqrt{\frac{\lambda}{t}} \leq \frac{\lambda}{10}, \]
and it follows from (89) that
(91) \[ P[y_l \cdot l \geq \lambda] \geq \frac{1}{4} P[B_{D(l)t} \geq \lambda \left( 1 + C_5 \sqrt{\frac{\lambda}{t}} \right) + \| \chi \|_\infty], \]
which proves Proposition 3.4. \( \square \)

Now let \( t > 0, x, y \in \mathbb{R}^d \) and \( p(t, x, y) \) be the heat kernel associated with the Dirichlet form (32). Using the chain rule one obtains that, for \( 0 < q < 1 \) and \( \delta > 0 \),
(92) \[ p(t, x, y) \geq C_{d, \text{Osc} (U)} P_x (y_{1q} \in B(y, \delta \sqrt{t})) \inf_{z \in B(y, \delta \sqrt{t})} p((1 - q)t, z, y). \]
It follows by Aronson’s estimates that
(93) \[ p(t, x, y) \geq C_{d, \text{Osc} (U)} P_x (y_{1q} \in B(y, \delta \sqrt{t})) \times (t(1 - q))^{-d/2} \exp(-C_{d, \text{Osc} (U)}, 2\delta^2/(1 - q)). \]

Now for \( l \in \mathbb{R}^d \) let us define the probability measure \( \bar{P}_x \) as
(94) \[ \frac{d\bar{P}_x}{dP_x} = \frac{e^{l \cdot y_t}}{E_x [e^{l \cdot y_t}]} \]
Henceforth we can assume \( x := 0 \) and we will fix
(95) \[ l := D(U)^{-1} y/(qt) \]
and assume
(96) \[ |l| \leq 1. \]
Writing \( \bar{E}_x \) as the expectation associated with \( \bar{P}_x \), one has
(97) \[ P_0 (y_{1q} \in B(y, \delta \sqrt{t})) = \bar{E}_0 [e^{-l \cdot y_{1q}} 1_{y_{1q} \in B(y, \delta \sqrt{t})}] \bar{E}_0 [e^{-l \cdot y_{1q}}] \]
\[ \geq e^{-y D(U)^{-1} y/(qt) - C_{d, \text{Osc} (U)}, 3|y|/(qt)^{1/2}} \bar{E}_0 [y_{1q} \in B(y, \delta \sqrt{t})] \bar{E}_0 [e^{-l \cdot y_{1q}}]. \]
Now it is trivial to check that the generator of $\gamma_t$ with respect to $\bar{P}_x$ is
\[ \bar{L} = \Delta/2 - \nabla U \nabla + l \cdot \nabla. \]

Let us write $\bar{p}$ the heat kernel associated with that generator. It is trivial to obtain from (95), (96) and Theorem 1.4 of Norris (1997) that for $z \in B(y, \delta \sqrt{t})$ one has
\[ \bar{p}(tq,0,z) \geq C_{d,\text{Osc}(U),4}(qt)^{-d/2} \exp(-C_{d,\text{Osc}(U),5}\delta^2/q). \]

It follows that
\[ \bar{P}_0[\gamma_{tq} \in B(y, \delta \sqrt{t})] \geq C_{d,\text{Osc}(U),6}\delta dq^{-d/2} \exp(-C_{d,\text{Osc}(U),5}\delta^2/q). \]

Moreover, for $\lambda > 0$,
\[ \mathbb{E}_0[e^{l \cdot \gamma_{tq}}] \geq \mathbb{P}_0[l \cdot \gamma_{tq} \geq \lambda] e^{\lambda}, \]

and choosing $\lambda = lD^{-1}(U)ltq$ one easily obtains from Proposition 3.4 that there exist constants $C_1, C_2, C_3$ depending on $d$ and $\text{Osc}(U)$ such that for $|y| > C_1$ and $C_2|y| < tq$ one has
\[ \mathbb{E}_0[e^{l \cdot \gamma_{tq}}] \geq \exp\left(\frac{yD^{-1}(U)y}{2qt}(1 - F)\right) \]

with
\[ F := C_3(qt/y^2 + |y|/(qt)). \]

Now let us choose
\[ q := 1 - \exp(-|x - y|t^{-1/2}) \]

and
\[ \delta := (1 - q)^{1/2}. \]

With these values for $q$ and $\delta$ and combining (102) with (97) and (93) one obtains that, for $|y - x| > C_{7,d,\text{Osc}(U)}$ and $C_{8,d,\text{Osc}(U)}|y - x| < t$,
\[ p_t(x,y) \geq C_{9,d,\text{Osc}(U)}t^{-d/2} \exp(-(1 - F_2)|x - y|^2_{D^{-1}(U)}/(2t)) \]

with
\[ F_2 := C_{10,d,\text{Osc}(U)}(t/|x - y|^2 + |y - x|/(t)). \]

It is then easy to deduce the lower bound of Theorem 2.8 by an appropriate shift of the constants.
3.2. An analytical inequality for subharmonic functions.

3.2.1. The inequality: Theorem 2.11. There is no loss of generality by assuming \( \Omega \) to be the segment \((0, 1)\). We will give a geometrical proof Theorem 2.11 explaining why we expect the existence of an homotopy invariant constant \( C_{d, \Omega} \) in Conjecture 2.12. The Theorem 2.11 is proven if the inequality (47) is true when \( \phi \) and \( \psi \) are Green functions \( G_\lambda(x, z) \) of \(-\nabla(\lambda \nabla)\) with Dirichlet condition on \( \partial(0, 1) \).

Let \((x, y) \in (0, 1)^2, x < y\). Write \( \Omega_1 = \{z \in \Omega : \nabla_z G(x, z) \nabla_z G(y, z) < 0\} \). Inequality (47) is true if

\[
- \int_{\Omega_1} \nabla_z G(x, z) \lambda(z) \nabla_z G(y, z) \, dz \leq \int_{\Omega} \nabla_z G(x, z) \lambda(z) \nabla_z G(y, z) \, dz.
\]

Write \( A_x = \{z \in \Omega : G(x, z) > G(x, y)\} \) and \( A_y = \{z \in \Omega : G(y, z) > G(x, y)\} \). Integrating by parts one obtains

\[
\int_{A_x} \nabla_z G(x, z) \lambda(z) \nabla_z G(y, z) \, dz = \int_{A_y} \nabla_z G(x, z) \lambda(z) \nabla_z G(y, z) \, dz = 0.
\]

Now the one-dimensional specificity shall be used. Since \( G(x, z) \) is increasing from 0 to \( x \) and decreasing from \( x \) to 1, it follows that \( \Omega_1 = (x, y) \) and \((A_x / \Omega_1) \cap (A_y / \Omega_1) = \emptyset\). Combining this with (109) one obtains (108), which proves the theorem. Let us note that a simple computation shows that the constant 3 is sharp.

3.2.2. Equivalence with the stability of the Green functions: Proposition 2.13. Write, for \( \varepsilon \in [0, 1] \), \( \lambda_\varepsilon(x) = e^{U(x) + \varepsilon T(x)} \). Write \( \psi_\varepsilon \) for the solution of \(-\nabla(\lambda_\varepsilon \nabla \psi_\varepsilon) = g\) with Dirichlet condition on \( \overline{\Omega} \) and \( g \in C^\infty(\overline{\Omega}), g > 0\).

Assume Conjecture 2.12 to be true, then Proposition 2.13 is proven if

\[
e^{-C_{d, \Omega}} \|\psi\|_\infty \leq \left\| \frac{\psi_1}{\psi_0} \right\|_\infty \leq e^{C_{d, \Omega}} \|T\|_\infty.
\]

One obtains by differentiation (writing \( L_{\lambda_\varepsilon} = -\nabla \lambda_\varepsilon \nabla \)) \( L_{\lambda_\varepsilon} \partial_\varepsilon \psi_\varepsilon = -L_{\partial_\varepsilon \lambda_\varepsilon} \psi_\varepsilon \), which leads, by integration by parts, to

\[
\partial_\varepsilon \psi_\varepsilon = -\int_{\Omega} \nabla_y G_{\lambda_\varepsilon}(x, y) \lambda_\varepsilon(y) \nabla_y G_{\lambda_\varepsilon}(y, z) T(y) g_\varepsilon(z) \, dy \, dz.
\]

Using Conjecture 2.12,

\[
|\partial_\varepsilon \psi_\varepsilon| \leq \|T\|_\infty \int_{\Omega} \left| \nabla_y G_{\lambda_\varepsilon}(x, y) \lambda_\varepsilon(y) \nabla_y G_{\lambda_\varepsilon}(y, z) \right| g_\varepsilon(z) \, dy \, dz \leq \|T\|_\infty C_{d, \Omega} \psi_\varepsilon.
\]

And integrating \( \partial_\varepsilon \ln \psi_\varepsilon \leq \|T\|_\infty C_{d, \Omega} \) one obtains the upper bound in (110) (the lower bound being proven in a similar way).
Conversely, if conjecture 2.12 is false one can find $\delta > 0, x, z \in \Omega^2$ and $g$ being a smooth approximation of a Dirac around $z$ such that, if
\[ T(y) = -\text{Sign}(\nabla_y G_{\lambda_c}(x, y)\lambda_c(y)\nabla_y G_{\lambda_c}(y, z)), \]
one has
\[ \partial_\varepsilon \ln \psi_\varepsilon(x) > \|T\|_\infty (1 + \delta) C_{d,\Omega}, \]
which leads to a contradiction of (49).

3.3. Subdiffusive behavior.

3.3.1. Exit times: Theorems 2.1, 2.6 and Proposition 2.7. For $r > 1$, write the number of effective scales
\[ n_{\text{ef}}(r) = \sup\{n \geq 0 : R_n \leq r\}. \]
First, let us prove that the exit time from $B(0, r)$ is controlled by the homogenization on those first $n_{\text{ef}}(r)$ scales:

**Lemma 3.7.**
\[ \frac{r^2}{D(V^{n_{\text{ef}}}_{0}(r))} \leq \mathbb{E}_0[\tau(0, r)] \leq \frac{r^2}{D(V^{n_{\text{ef}}}_{0}(r))} C_\tau \]
with $C_\tau = 4 e^{6(K_0 + K_1/(\rho_{\text{min}} - 1))}$.

**Proof.** Write $\mathbb{E}^U$, the expectation with respect to the law of probability associated with the generator $\frac{1}{2} \Delta - \nabla U \nabla$. By Theorem 2.11 and Proposition 2.13, one obtains that
\[ e^{-6 \text{Osc}_{r}(V^{\infty}_{n_{\text{ef}}(r)+1})} \leq \mathbb{E}_0[\tau(0, r)]/\mathbb{E}_0^{V^{n_{\text{ef}}}_{0}(r)}[\tau(0, r)] \leq e^{6 \text{Osc}_{r}(V^{\infty}_{n_{\text{ef}}(r)+1})}. \]
Bounding $U_{n_{\text{ef}}+1}(x)$ by $\text{Osc}(U_n) \leq K_0$, and for $k \geq n_{\text{ef}} + 2$, $U_k(x)$ by $\|\nabla U_k\|_{\infty} \|x\| \leq K_1 \|x\|/R_k$ one obtains that, for $x \in B(0, r),
\[ \left|V^{\infty}_{n_{\text{ef}}(r)+1}(x)\right| \leq K_0 + K_1/(\rho_{\text{min}} - 1). \]
Writing $p_{\text{ef}}$ corresponds to the maximum number of periods of the scale $n_{\text{ef}}$ included in the segment $[0, r]$: $p_{\text{ef}}(r) = \sup\{p \geq 1 : p R_{n_{\text{ef}}(r)} \leq r\}$, one obtains
\[ \mathbb{E}_0^{V^{n_{\text{ef}}}_{0}(r)}[\tau(0, p_{\text{ef}}(r) R_{n_{\text{ef}}(r)})] \leq \mathbb{E}_0^{V^{n_{\text{ef}}}_{0}(r)}[\tau(0, r)] \leq \mathbb{E}_0^{V^{n_{\text{ef}}}_{0}(r)}[\tau(0, (p_{\text{ef}}(r) + 1) R_{n_{\text{ef}}(r)})]. \]
Using $\mathbb{E}_0^{V^{n_{\text{ef}}}_{0}(r)}[\tau(0, k R_{n_{\text{ef}}(r)})] = (k R_{n_{\text{ef}}(r)})^2/D(V^{n_{\text{ef}}}_{0}(r))$ and (116)–(118) one obtains (115).

We will need the following mixing lemma.
Lemma 3.8. Let \((g, f) \in C^1(\mathbb{T}^d)^2\) and \(R \in \mathbb{N}^*\):
\[
\left| \int_{\mathbb{T}^d} g(x) f(Rx) \, dx - \int_{\mathbb{T}^d} g(x) \int_{\mathbb{T}^d} f(x) \, dx \right| \leq \| \nabla g \|_\infty / R \int_{\mathbb{T}^d} |f| \, dx.
\]

Proof. The proof follows trivially from the following equation:
\[
\int_{\mathbb{T}^d} g(x) f(Rx) \, dx - \int_{\mathbb{T}^d} g(x) \int_{\mathbb{T}^d} f(x) \, dx = \int_{y \in [0, 1]^d, x \in \mathbb{T}^d} f(Rx + y) \left( g(x + y/R) - g(x) \right).
\]
\[(119)\]

From Lemma (3.8) we will deduce a quantitative estimate on the multiscale effective diffusivities:

Lemma 3.9.
\[
(\lambda_{\min} e^{-4K_1/\rho_{\min}})^n \leq D(V^{n-1}) \leq (\lambda_{\max} e^{4K_1/\rho_{\min}})^n.
\]

Proof. The proof of (120) is based on the following functional mixing estimate (obtained from Lemma 3.8): for \(U, W \in C^1(\mathbb{T})\) and \(R \in \mathbb{N}^*\) one has
\[
e^{-\|\nabla W\|_\infty / R} \lesssim \int_{\mathbb{T}} e^{U(Rx) + W(x)} \frac{d}{dx} \left( \int_{\mathbb{T}} e^{U(x)} \, dx \int_{\mathbb{T}} e^{W(x)} \, dx \right)
\leq e^{\|\nabla W\|_\infty / R}.
\]
(121)

Then by the explicit formula (20) and a straightforward induction on \(n\), one obtains that [using (5)]
\[
\prod_{k=0}^{n-1} \left( e^{4K_1/r_k} \int_{\mathbb{T}} e^{2U_k(x)} \, dx \int_{\mathbb{T}} e^{-2U_k(x)} \, dx \right)^{-1} \leq D(V^{n-1}_0),
\]
(122)

\[
D(V^{n-1}_0) \leq \prod_{k=0}^{n-1} \left( e^{-4K_1/r_k} \int_{\mathbb{T}} e^{2U_k(x)} \, dx \int_{\mathbb{T}} e^{-2U_k(x)} \, dx \right)^{-1},
\]
(123)

which leads to (120) by (8) and (10). \(\square\)

Combining (120) with (115), (114) and (10), one obtains Theorem 2.1.
When the medium is self-similar, we will need the following lemma:

Lemma 3.10.
\[
\lim_{n \to \infty} -\frac{1}{n} \ln (D(V^{n-1})) = \mathcal{P}_\rho(2U) + \mathcal{P}_\rho(-2U).
\]
(124)
PROOF. The limit (124) is a direct consequence of the following theorem that is an application of the theory of level-3 large deviations [we refer to Ellis (1985) for a sufficient reminder].

**Theorem 3.11.** Let \( U \in C^\alpha(\mathbb{T}^d) \) (the Hölder continuous with exponent \( \alpha > 0 \)). Let \( R \in \mathbb{N}, R \geq 2 \). Then

\[
\lim_{n \to \infty} \frac{1}{n} \ln \int_{\mathbb{T}^d} \exp \left( \sum_{k=0}^{n-1} U(R^k x) \right) dx = \mathcal{P}_R(U).
\]

We have written \( \mathcal{P}_R \) as the pressure associated with the scaling shift induced by \( R \) on the torus. For \( R \in \mathbb{N} \setminus \{0, 1\} \) one can see the torus as a shift space equipped with the transformation \( s_R: \mathbb{T}^d \to \mathbb{T}^d \):

\[
s_R: \mathbb{T}^d \to \mathbb{T}^d
\]

where for each \( k \), \( x^k \) is a vector in \( B = \{0, 1, \ldots, R - 1\}^d \) and for each \( i \in \{1, \ldots, d\} \), \( \sum_{k=1}^\infty \frac{x^k_i}{R^k} \) is the expression of \( x_i \) in base \( R \) (\( x^k_i \in \{0, \ldots, R - 1\} \)).

Give \( B \) with the discrete topology and \( B^{\mathbb{N}^*} \) with the product topology. Write \( \mu \) the probability measure on \( B \) affecting identical weight \( 1/R^d \) to each element of \( B \) and write \( \mathbb{P}_\mu \) the associated product measure on \( B^{\mathbb{N}^*} \).

With respect to the probability space \( (B^{\mathbb{N}^*}, \mathscr{B}(B^{\mathbb{N}^*}), \mathbb{P}_\mu) \) the coordinate representation process \( x = (x^1, \ldots, x^n, \ldots) \) is a sequence of i.i.d. random variables distributed by \( \mu \). When \( x \) is seen as an element of the torus \( \mathbb{T}^d \) then the probability measure induced by \( \mu \) on the torus is the Lebesgue measure.

Define the empirical measure \( E_n \) associated with the process \( x \) by

\[
E_n(x, \cdot) = \frac{1}{n} \sum_{k=0}^{n-1} \delta_{s_R^{k} \text{cycle}(x,n)},
\]

where \( \text{cycle}(x, n) \) is the periodic point in \( B^{\mathbb{N}^*} \) obtained by repeating \( (x^1, \ldots, x^n) \) periodically. For each \( x \), \( E_n(x, \cdot) \) is an element of the space \( \mathcal{M}(B^{\mathbb{N}^*}) \) of measures on \( B^{\mathbb{N}^*} \) and invariant by the shift \( s_R \).

Then by Theorem 9.1.1 of Ellis (1985), \( \{Q_n^{(3)}\} \), the \( \mathbb{P}_\mu \) distribution on \( \mathcal{M}(B^{\mathbb{N}^*}) \) of the empirical process \( \{E_n\} \) has a large deviation property with speed \( n \) and entropy function \( I^{(3)}_{\mu} \).

Recall that for \( P \in \mathcal{M}(B^{\mathbb{N}^*}) \), \( I^{(3)}_{\mu}(P) = \int_{B^{\mathbb{N}^*}} I^{(2)}_{\mu}(\tilde{P}) dP \) where \( \tilde{P} \) denotes the marginal distribution of \( x^1 \) associated with \( P \) and \( I^{(2)}_{\mu} \) is the relative entropy of \( \tilde{P} \) with respect to \( \mu \): \( I^{(2)}_{\mu}(\eta) = \int_B \ln \frac{d\eta}{d\mu} d\mu \).
Choosing \( U \in C(\mathbb{T}^d) \), the Hölder continuous with exponent \( \alpha \), one deduces from the large deviation property of \( \{Q_n^{(3)}\} \) and Varadhan’s theorem that

\[
\lim_{n \to \infty} \frac{1}{n} \ln \int_{\mathbb{T}^d} \exp(nE_n(x, U)) \, dx = \mathcal{P}_R(U),
\]

where \( \mathcal{P}_R(U) \) is the pressure of \( U \). Recall that

\[
\mathcal{P}_R(U) = \sup_{P \in \mathcal{M}_{sR}(B^{N^*})} \left\{ \int U \, dP - I_{(3)}(P) \right\},
\]

where \( \mathcal{M}_{sR}(B^{N^*}) \) is the space of measures on \( B^{N^*} \) invariant by the shift \( s_R \).

Since \( U \) is Hölder continuous,

\[
\left| nE_n(x, U) - \sum_{k=0}^{n-1} U(R^k x) \right| \leq \sum_{k=0}^{n-1} \left( \frac{C_d}{R^{\alpha-k}} \right)^{\alpha} C(d, \alpha) \sum_{k=0}^{\infty} \frac{1}{R^{k\alpha}} \leq C(d, \alpha, R) < \infty,
\]

and one obtains Theorem 3.11 from (130) and (128) \( \square \)

Combining (124) with (115) and (114), one obtains Theorem 2.6.

Now let us prove Proposition 2.7. The basic properties of the pressure can be found in Keller (1998) Theorem 4.1.10. [note that the definition of the pressure given here differs from the standard one of the topological pressure by a constant that is \( d \ln R \), here \( \mathcal{P}_R(0) = 0 \)]. Let us remind that \( \mathcal{P}_R \) is a convex function on the space of upper semicontinuous functions on the torus to \((-\infty, \infty)\) thus \( \mathcal{P}_R(U) + \mathcal{P}_R(-U) \geq 0 \).

We will recall the strict convexity of the topological pressure on a well-defined equivalence space: with \( s_R \) is associated a scaling operator \( S_R \) acting on the periodic continuous functions on \( \mathbb{T}^d \):

\[
S_R : C(\mathbb{T}^d) \to C(\mathbb{T}^d),
\]

\[
(x \to f(x)) \to (x \to f(s_R x) = f(Rx)).
\]

Write \( \mathcal{I}_{s_R}(\mathbb{T}^d) \) the closed subspace of \( C(\mathbb{T}^d) \) generated by the elements \( V - S_R^k V \) with \( V \in C(\mathbb{T}^d) \) and \( k \in \mathbb{N} \). Write \([U]\) the equivalence class of \( U \), then by Proposition 4.7 of Ruelle (1978), the function

\[
\mathcal{P}_R : C(\mathbb{T}^d) / \mathcal{I}_{s_R}(\mathbb{T}^d) \to (-\infty, +\infty),
\]

\[
[U] \to \mathcal{P}_R(U)
\]

is well defined on the set of equivalence classes induced by \( \mathcal{I}_{s_R}(\mathbb{T}^d) \) on \( C(\mathbb{T}^d) \). Moreover, it is strictly convex on the subset

\[
\left\{ [U] \in C(\mathbb{T}^d) / \mathcal{I}_{s_R}(\mathbb{T}^d) : \int_{\mathbb{T}^d} U(x) \, dx = 0 \right\}.
\]
We will now prove Proposition 2.7, since for \( c \in \mathbb{R} \), \( \mathcal{P}(U + c) = \mathcal{P}(U) + c \), it is sufficient to assume \( \int_{\mathbb{T}^d} U(x) \, dx = 0 \) and show that

\[
\mathcal{P}_R(2U) + \mathcal{P}_R(-2U) = 0 \iff \lim_{n \to \infty} \frac{1}{n} \left\| \sum_{k=0}^{n-1} S_{R^k} U \right\|_\infty = 0.
\]

(\( \Leftarrow \)): This implication is easy since

\[
0 \leq \mathcal{P}_R(2U) + \mathcal{P}_R(-2U) \leq \lim_{n \to \infty} \frac{4}{n} \left\| \sum_{k=0}^{n-1} S_{R^k} U \right\|_\infty.
\]

(\( \Rightarrow \)): Assume \( \mathcal{P}_R(2U) + \mathcal{P}_R(-2U) = 0 \), then let \( \varepsilon > 0 \). Then by the strict convexity of the pressure as described above there exist \( W_1, \ldots, W_k \in C(T^d) \) and \( m_1, \ldots, m_k \in \mathbb{N} \setminus \{0, 1\} \), \( \lambda_1, \ldots, \lambda_k \in \mathbb{R} \) such that \( W = \sum_{p=1}^k \lambda_p (W_p - S_{R^p} W_p) \) and \( \|U - W\|_\infty \leq \varepsilon \). Since \( \sum_{p=0}^{n-1} S_{R^p} W \) remains bounded it follows that

\[
\lim_{n \to \infty} \frac{1}{n} \left\| \sum_{k=0}^{n-1} S_{R^k} U \right\|_\infty \leq \varepsilon,
\]

which leads to the proof.

3.3.2. Mean squared displacement: Proposition 2.5 Theorem 2.2. Let \( y_t \) be the solution of (1). Write

\[
n_{\text{flu}}(t) = \sup\{n \in \mathbb{N} : R^2_n \leq t\};
\]

\( n_{\text{flu}} \) shall be the number of fluctuating scales that have an influence on the mean squared displacement at the time \( t \) (the effective scales plus the perturbation scales). Choose the number of perturbation scales to be

\[
n_{\text{per}} = \inf\{n \in \mathbb{N} : R^2_{n_{\text{flu}}} - n_{\text{flu}} 14 n K_0 10^4 \leq t D(V^{n_{\text{flu}}}_{0})\}.
\]

We will now prove the following proposition:

**Proposition 3.12.** For \( \rho_{\min} > 10 e^{30 K_1} \) and \( t > R_g \), \( n_{\text{per}} \) is well defined and

\[
C_1 e^{-8 n_{\text{per}} K_0} D(V^{n_{\text{flu}}}_{0}) t \leq \mathbb{E}[y_t^2] \leq C_2 e^{8 n_{\text{per}} K_0} D(V^{n_{\text{flu}}}_{0}) t.
\]

**Proof.** The proof of (139) is based on analytical inequalities that allow us to control the stability of the homogenization process on the smaller scales under the perturbation of larger ones. More precisely we will first work on an abstract decomposition of \( V \) given by (2) into effective scales \( U \), perturbation scales \( P \) and drift scales \( T : V = U + P + T \) with \( (U, P, T) \in C^\infty(T^1_{R_U}) \times C^\infty(T^1_{R_W}) \times C^\infty(\mathbb{R}) \), \( R_U, R_W \in \mathbb{N} \setminus \{0, 1\} \), \( R_W / R_U = R_P \in \mathbb{N}^* \) and \( W = U + P \) shall correspond to fluctuating scales.
Write $\chi^W$ as the solution of the cell problem associated with $L_W$ [$L_W \chi_W = -\nabla W$, $\chi_W(0) = 0$] and $F_W(x) = x - \chi^W(x)$. Since $F_W$ is harmonic with respect to $L_W = L_V + \nabla T \nabla$, one obtains by Itô’s formula that $F_W(y_t) = \int_0^t \nabla F_W(y_s) d\omega_s - \int_0^t \nabla T \nabla F_W(y_s) ds$ from which one obtains that

$$(1) = (\frac{1}{2} - \|\nabla T\|_\infty^2) \mathbb{E} \left[ \int_0^t |\nabla F_W(y_s)|^2 ds \right]$$

$$\leq \mathbb{E} [F_W^2(y_t)] \leq 2(1+t\|\nabla T\|_\infty^2) \mathbb{E} \left[ \int_0^t |\nabla F_W(y_s)|^2 ds \right].$$

Write $\chi^P$ as the solution of the cell problem associated with $L_P$ and $F^P = x - \chi^P$. We will show the following:

**Lemma 3.13.** $F_W = F_P - H^U$ with

$$(11) e^{-4\text{Osc}(P)} \chi^2 \leq (F^P(x))^2 \leq e^{4\text{Osc}(P)} \chi^2$$

and

$$(12) \|H^U\|_\infty \leq 2(1 + 4\|\nabla P\|_\infty) e^{2\text{Osc}(P)} R_W/R_P.$$  

**Proof.** Inequality (11) is a direct consequence of the explicit formula

$$F^P(x) = R_W \int_0^x e^{2P(y)} dy / \int_0^{R_W} e^{2P(y)} dy.$$  

Inequality (12) follows from the explicit formula

$$H^U(x) = R_W \left( \frac{\int_0^x e^{2P(y)} dy}{\int_0^{R_W} e^{2P(y)} dy} - \frac{\int_0^x e^{2(P(y)+U(y))} dy}{\int_0^{R_W} e^{2(P(y)+U(y))} dy} \right)$$

noting that the period of $P$ and $U$ are $R_W$ and $R_W/R_P$ and Lemma 3.8.

The long time behavior of $\mathbb{E} [\int_0^t |\nabla F_W(y_s)|^2 ds]$ is a perturbation of $D(W)t$ as shown in the following lemma:

**Lemma 3.14.** If

$$R_P > 16 e^{4\text{Osc}(P)} (\|\nabla P\|_\infty + \|\nabla T\|_\infty) e^{2\|\nabla T\|_\infty/R_P},$$

then for $t > 0$,

$$(13) -(R_W^2/R_P^2) (e^{10\text{Osc}(P)/R_P^2} 100 e^{4\|\nabla T\|_\infty/R_P} + \frac{1}{6} e^{-4\text{Osc}(P)} D(W)t)$$

$$\leq \mathbb{E} \left[ \int_0^t |\nabla F_W|^2(y_s) ds \right]$$

and

$$(14) \mathbb{E} \left[ \int_0^t |\nabla F_W|^2(y_s) ds \right] \leq 6 e^{4\text{Osc}(P)} D(W)t + (R_W^2/R_P^2) (e^{10\text{Osc}(P)/R_P^2} 900 e^{4\|\nabla T\|_\infty/R_P}).$$
Proof. For the proofs of (143) and (144) by scaling one can assume that $R_W = 1$ and $R_U = 1/R_P$. Write for $\zeta > 0$

\begin{equation}
\phi_\zeta = 2 \int_0^x \frac{e^{2V(y)}}{\int_0^1 e^{2W(y)}dy} \left[ \int_0^y \frac{e^{2(P-T)(z)}}{\int_0^1 e^{2P(z)}dz} dz - \zeta \int_0^y \frac{e^{-2(P+T)(z)}}{\int_0^1 e^{-2P(z)}dz} dz \right] dy.
\end{equation}

Using Lemma 3.8 to separate the scales in (145), it is an easy exercise to obtain that if $R_P > 16 e^{4 \text{Osc}(P)} (\|\nabla P\|_\infty + \|\nabla T\|_\infty) e^{2\|\nabla T\|_\infty}/R_P$ then:

- for $\zeta = 6 e^{4 \text{Osc}(P)}$ one has $\sup R_\phi \leq 900 e^{10 \text{Osc}(P)} R_2 e^{4\|\nabla T\|_\infty}/R$;
- for $\zeta = \frac{1}{6} e^{-4 \text{Osc}(P)}$ one has $\inf R_\phi \geq -100 e^{10 \text{Osc}(P)} R_2 e^{-4\|\nabla T\|_\infty}/R$.

Observing that $L_V \phi_\zeta = |l - \chi W|^2 - \zeta D(W)$ one deduces (143) and (144) by applying Itô's formula. □

Combining (140), (141), (143) and (144) and choosing $U = V_0^{n_{flu}} - n_{per}$, $P = V_{n_{flu}}^{n_{flu}} - n_{per} + 1$, $T = V^{n_{flu}}_{n_{flu}}$ ($R_W = R_{n_{flu}}$, $R_P = R_{n_{flu}}/R_{n_{flu} - n_{per}}$) and $n_{flu}$ as defined in (137), one obtains that, for $\rho_{\min} > C_{K_1, K_0}$,

\begin{equation}
D(V_0^{n_{flu}}) t e^{-8n_{per} K_0}/24 - R_{n_{flu} - n_{per}}^2 500 e^{6n_{per} K_0} \leq E[y_t^2],
\end{equation}

\begin{equation}
E[y_t^2] \leq (D(V_0^{n_{flu}}) t + R_{n_{flu} - n_{per}}^2) e^{8n_{per} K_0} 500,
\end{equation}

which leads to (139) by the choice (138) for $n_{per}$. □

By the uniform control of the ratios (10) one obtains quantitative estimates on the number of fluctuating and perturbation scales (137) and (138); combining them with the control (139) and the exponential speed of convergence of the multiscale effective diffusivities toward zero (120), one obtains Proposition 2.5 and Theorem 2.2.

3.3.3. Heat kernel tail: Theorem 2.3. As was done for the mean squared displacement, the proof of Theorem 2.3 shall follow from an abstract decomposition of the potential $V$. More precisely, let $R_W \in \mathbb{N} \setminus \{0, 1\}$, $(W, T) \in C^\infty(T_{R_W}^1) \times C^\infty(\mathbb{R})$, $(\|\nabla T\|_\infty < \infty)$ and write $V = W + T$ and $y_t$ the diffusion associated with $L_V$. It has been shown in the proof of Proposition 3.12 that by decomposing $W$ into $U + P$ where $U$ is of period $R_W/\rho_{\min} \in \mathbb{N}$, one has, for all $t > 0$ and all $x \in \mathbb{R}^d$,

\begin{equation}
E_x \left[ \int_0^t |\nabla F^W|^2(y_s) ds \right] \leq \zeta_2 D(W)t + \frac{R_{W}^2}{R_P^2} C_2^\phi,
\end{equation}

where the constants $C_2^\phi$, $\zeta_2$ are those given by equation (144). We will now show that from the control (148) (and $\|X^W\|_\infty \leq R_W$ that is given by the explicit formula of the solution of the cell problem) one can deduce the following lemma:
LEMMA 3.15. For

\[ R_W \leq h/2. \]  \hspace{1cm} (149)

\[ \| \nabla T \|_\infty 2^3 (\xi_2 D(W))^{1/2} \leq (h/t) \leq \left( \frac{R_P}{(R_W \sqrt{C_2^\phi})} \right) \xi_2 D(W) \]

and

\[ \left( \frac{R_P}{(R_W \sqrt{C_2^\phi})} \right) \xi_2 D(W) \leq (h/t), \]

\[ \frac{(R_P/(R_W \sqrt{C_2^\phi})) \xi_2 D(W) e^{-h^2/(2\xi_2 D(W)t)}}{\zeta_2 D(W) t} \leq (h/t), \]

one has

\[ \mathbb{P}[y_t \geq h] \leq C e^{-h^2/(2^9 \xi_2 D(W)t)}. \]  \hspace{1cm} (152)

**Proof.** The proof of (152) is based on a control of the Laplace transform of \( y_t \); more precisely it is well known that for \( \lambda > 0 \) and \( h > 0 \) one has \( \mathbb{P}[y_t \geq h] \leq \mathbb{E}[e^{\lambda (y_t - h)}] \). Observing that \( y_t = \chi_W(y_t) + \int_0^t \nabla F_W(y_s) d\omega_s - \int_0^t \nabla T \cdot \nabla F_W(y_s) ds \) and using \( \| \chi_W \|_\infty \leq R_W \) one deduces by the Cauchy–Schwarz inequality that

\[ \mathbb{P}[y_t \geq h] \leq e^{\lambda (R_W - h)} \mathbb{E}[e^{2\lambda \int_0^t \nabla F_W(y_s) d\omega_s}]^{1/2} \times \mathbb{E}[e^{2\sqrt{\lambda t} \| \nabla T \|_\infty \lambda (\int_0^t \| \nabla F_W(y_s) \|^2 ds)^{1/2}}]^{1/2}. \]  \hspace{1cm} (153)

If \( X \) is a positive bounded random variable, \( \mu' > 0 \) and \( \lambda' > 0 \), it is easy to show by integrating by parts over \( d\mathbb{P}(X \geq x) \) and using \( \mathbb{P}(X \geq x) \leq \mathbb{E}[\exp(\lambda'(X - x))] \) that

\[ \mathbb{E}[\exp(\mu' \sqrt{X})] \leq 1 + \mu' \exp \left( \frac{(\mu')^2}{4\lambda'} \right) \sqrt{\frac{4}{\lambda'}} \mathbb{E}[\exp(\lambda' X)]. \]

Applying this inequality to (153) with \( X = \int_0^t \| \nabla F_W(y_s) \|^2 ds \), \( \lambda' = 8\lambda^2 \) and \( \mu' = 2\lambda\sqrt{\lambda} \| \nabla T \|_\infty \) and observing by Itô’s formula that \( \mathbb{E}[e^{2\lambda \int_0^t \nabla F_W(y_s) d\omega_s}] \leq \mathbb{E}[e^{8\lambda^2 \int_0^t \| \nabla F_W(y_s) \|^2 ds}]^{1/2} \) one obtains

\[ \mathbb{P}[y_t \geq h] \leq C e^{\lambda (R_W - h)} e^{\| \nabla T \|_\infty h/4} \mathbb{E}[e^{8\lambda^2 \int_0^t \| \nabla F_W(y_s) \|^2 ds}]. \]

Now observe that \( \int_0^t \nabla F_W(y_s) d\omega_s \) satisfies the conditions of Theorem 2.9 with \( f_2 = \xi_2 D(W) \), and \( f_0(f_1 - f_2) = \frac{R_P^2}{R_P^2} C_2^\phi. \) It follows that for

\[ 8\lambda^2 \leq \left( \frac{R_P^2}{R_P^2} \right)/(2e R_W^2 C_2^\phi) \]

one has

\[ \mathbb{E}[e^{8\lambda^2 \int_0^t \| \nabla F_W(y_s) \|^2 ds}] \leq C R_P^4 (e^{8\lambda^2 \xi_2 D(W)t})/(\lambda^4 (C_2^\phi)^2 R_W^4). \]
Assuming $R_W < h/2$ and choosing $\lambda = \frac{h}{2\zeta D(W)t}$, the condition on $\lambda$ in (154) is satisfied under the right inequality in (150) and one obtains

$$\mathbb{P}[Y_t \geq h] \leq C e^{-h^2/(2^7\zeta^2 D(W)t)} e^{\|\nabla T\|_\infty 4t/(h^2(2^5 C^2_2) R^4_W)}.$$  

From this the result (152) follows easily by assuming the left inequality in (150) (which basically says that the influence of the drift scales $\|\nabla T\|_\infty$ is small in front of the influence of the fluctuating scales) and condition (151). \qed

Now let us choose $W = V_0^{n_{\text{flu}}}$, $P = V_{n_{\text{flu}} n_{\text{per}}+1}^{n_{\text{flu}}}$, $T = V_{n_{\text{flu}}+1}^{\infty}$ ($R_W = R_{n_{\text{flu}}}$, $R_p = R_{n_{\text{flu}}}/R_{n_{\text{flu}}-n_{\text{per}}}$) in Lemma 3.15. For $p \in \mathbb{N}^*$ define the function

$$n_{\text{per}}(p) = \inf\{ n \in \mathbb{N} : (R_p/R_{p-n}) e^{-3n K_0 (D(V_0^{p-1}))^{1/2}} \geq 2^9 e^{5K_1} \}. $$

$n_{\text{per}}(p)$ corresponds to the number of perturbation scales among $p$ fluctuating scales. We will from now assume that $\rho_{\text{min}} \geq 2^9 e^{11K_1}$, which implies that $n_{\text{per}}$ is well defined and $1 \leq n_{\text{per}}(p) \leq p$. Define

$$n_{\text{flu}}(t/h) = \inf\{ n \in \mathbb{N} : 2^6 (K_1/R_{n+1}) e^{2n_{\text{per}}(n) K_0 (D(V_0^n))^{1/2}} \leq h/t \}.$$ 

$n_{\text{flu}} - n_{\text{per}}$ corresponds to the number of fully homogenized scales given $t/h$. $n_{\text{flu}}$ is well defined and greater than 1 under the following assumption that basically says that homogenization has started on at least the first scale:

$$2 K_0 e^{2K_0} 2^{-6} \leq t/h. $$

By the definition of $n_{\text{flu}}$ the left inequality in (150) is satisfied. Using (156), the right inequality in (150) is implied by the definition of $n_{\text{per}}$. Inequality (149) is satisfied if $2 R_{n_{\text{flu}}} \leq h$; by the definition of $n_{\text{flu}}$ this is implied by the following inequality which basically says that the heat kernel behavior is far from its diagonal regime:

$$h^2/(D(V_0^{n_{\text{flu}}})^{1/2} t) \geq 2 K_1 e^{2K_0} 2^6 e^{2n_{\text{per}}(n_{\text{flu}}) K_0}.$$ 

By the definition of $n_{\text{flu}}$ and $n_{\text{per}}$, inequality (151) is satisfied by the following inequality which also says that the heat kernel is far from its diagonal regime:

$$2^{14} e^{4(n_{\text{per}}+1) K_0} \ln[R_{n_{\text{flu}}+1}] \leq h^2/(D(V_0^{n_{\text{flu}}}) t).$$ 

With this assumption, it follows by inequality (152) that

$$\mathbb{P}[Y_t \geq h] \leq C e^{-h^2/(2^{11} e^{4n_{\text{per}} K_0} D(V_0^{n_{\text{flu}}}) t)}. $$

Using the control (120) on $D(V_0^{n_{\text{flu}}})$, and (10) on the ratios, one obtains Theorem 2.3. Condition (157) is translated into the first inequality in (16) and conditions (158) and (159) into the second one.
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