Weak convergence of particle swarm optimization

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Abstract

Particle swarm optimization algorithm is a stochastic meta-heuristic solving global optimization problems appreciated for its efficiency and simplicity. It consists in a swarm of interacting particles exploring a domain and searching a global optimum. The trajectory of the particles has been well-studied in a deterministic case. More recently, authors shed some light on the stochastic approach to PSO, considering particles as random variables and studying them with probabilistic and statistical tools. These works paved the way to the present article. We focus here on weak convergence, the kind of stochastic convergence that appears in the Central Limit Theorem. We obtain three main results related to three different frameworks. These three settings depend either on the regime of the particles (oscillation or fast convergence) or on the sampling strategy (along the trajectory or in the swarm). The main application of these results is the construction of confidence intervals around the local optimum found by the swarm. The theorems are illustrated by a simulation study.

Keywords: Particle swarm optimization; Convergence; Central limit theorem

1 Introduction

The particle swarm optimization algorithm (PSO), based on social interactions (behaviors of birds) was introduced in Eberhart and Kennedy (1995). Since then, PSO has known a great popularity in many domains and gave birth to many variants of the original algorithm (see Zhang et al. (2015) for a survey of variants and applications). PSO is a stochastic meta-heuristic solving an optimization problem without any evaluation of the gradient. The algorithm explores the search space in an intelligent way thanks to a population of particles interacting with each other and updating at each step their position and their velocity. The dynamic of the particles relies on two attractors: their personal best position (historical best position of the particle denoted \(p_s^n\) below), and the neighborhood best position (corresponding to the social component of the particles, denoted \(g_s^n\)). In the dynamic equation, the attractors are linked with a stochastic process in order to explore the search space. Algorithm \(\text{Algorithm 1}\) refers to the classical version of PSO with \(S\) particles and \(N\) iterations.

\begin{algorithm}
\caption{Classical PSO}
\begin{algorithmic}
\State Initialize the swarm of \(S\) particles with random positions \(x_0^s\) and velocities \(v_0^s\) over the search space.
\For {\(n = 1\) to \(N\)}
\State Evaluate the optimization fitness function for each particle.
\State Update \(p_n^s\) (personal best position) and \(g_n^s\) (neighborhood best position).
\State Change velocity \((v_n^s)\) and position \((x_n^s)\) according to the dynamic equation.
\EndFor
\end{algorithmic}
\end{algorithm}

The convergence and stability analysis of PSO are important matters. In the literature, there are two kinds of convergence:

- the convergence of the particles towards a local or global optimum. This convergence is not obtained with the classical version of PSO. Van den Bergh and Engelbrecht (2010) and Schmitt and Wanka (2015) proposed a modified version of PSO to obtain the convergence.
- the convergence of each particle to a point (e.g. Poli (2009)).

If we focus on the convergence of each particle to a point, a prerequisite is the stability of the trajectory of the particles. In a deterministic case, Clerc and Kennedy (2002) dealt with the stability of the particles with some conditions on the parametrization of PSO. Later, Kadirkamanathan et al. (2006) used the Lyapunov stability theorem to study the stability. About the convergence of PSO, Van Den Bergh and Engelbrecht (2006) looked at the trajectories of the particles.
and proved that each particle converges to a stable point (deterministic analysis). Under stagnation hypotheses (no improvement of the personal and neighborhood best positions), Xu and Yu (2018) has the exact formula of the second moment. More recently, Bonardi and Michalewicz (2016) or Cleghorn and Engelbrecht (2018) provided results for the order-1 and order-2 stabilities with respectively stagnant and non-stagnant distribution assumptions (both weaker than the stagnation hypotheses). Since our main results rely on martingale convergence theorems the other recent work, Xu and Yu (2018) has to be mentioned as well.

Let us introduce some notations. We consider here a cost function \( f : \mathbb{R}^d \to \mathbb{R}^+ \) that should be minimized on a compact set \( \Omega \). Consequently the particles evolve in \( \Omega \subset \mathbb{R}^d \). Let \( x^s_n \in \mathbb{R}^d \; 1 \leq s \leq S \) be sequences of independent random vectors in \( \mathbb{R}^d \) whose margins are uniformly distributed over \([0,1]\) and denote by \( c_1 \) and \( c_2 \) three positive constants which will be discussed later. Then the PSO algorithm considered in the sequel is defined by the following equations (or dynamic equations, Poli (2009)):

\[
\begin{align*}
\dot{v}^s_n &= \omega \cdot v^s_n + c_1 r_1 n \odot (p^s_n - x^s_n) + c_2 r_2 n \odot (g^s_n - x^s_n), \\
\dot{x}^s_n &= x^s_n + v^s_{n+1}
\end{align*}
\]

where \( \odot \) stands for the Hadamard product:

\[ u \odot v = (u_1 v_1, \ldots, u_d v_d) \]

and \( p^s_n \) (resp. \( g^s_n \)) is the best personal position (resp. the best neighborhood position of the particle \( s \)):

\[
\begin{align*}
p^s_n &= \arg\min_{t \in \{x^s_0, \ldots, x^s_n\}} f(t), \\
g^s_n &= \arg\min_{t \in \{p^s_0, \ldots, p^s_n\} \cap \Omega} f(t)
\end{align*}
\]

with \( \Omega(s) \) the neighborhood of particle \( s \). This neighborhood depends on the swarm’s topology: if the topology is called global (all the particles communicate between each other) then \( g^s_n = g_n = \arg\min_{t \in \{p^s_0, \ldots, p^s_n\}} f(t) \) (see Lane et al. (2008)).

Our main objective is to provide (asymptotic) confidence sets for the global or local optimum of the cost function \( f \). If \( g = \arg\min_{t \in \Omega} f(t) \) for some domain \( \Omega \), a confidence region at level \( 1 - \alpha \) (with \( \alpha \in (0,1) \)) for \( g \) is a random set \( \Lambda \subset \mathbb{R}^d \) such that:

\[
P(g \in \Lambda) \geq 1 - \alpha
\]

The set \( \Lambda \) depends on the swarm and is consequently random due to the random evolution of the particles. The probability symbol above, \( \mathbb{P} \), depends on the distribution of the particles in the swarm. Let us illustrate the use of confidence interval for a real-valued PSO. Typically the kind of results we expect is: \( g \in [m, M] \) with a probability larger than, say 99 \%. This does not aim at yielding a precise estimate for \( g \) but defines a “control area” for \( g \), as well as a measure of the imprecision and variability of PSO.

Convergence of the swarm will not be an issue here. In fact we assume that the personal and global best converge: see assumptions \( A_2 \) and \( B_2 \) below. We are interested in the “step after”: localizing the limit of the particles with high probability, whatever their initialization and trajectories.

Formally, confidence set estimation forces us to inspect order two terms (i.e. the rate of convergence), typically convergence of the empirical variance. The word asymptotic just means that the sample size increases to infinity.
2 Main results

The usual euclidean norm and associated inner product for vectors in \( \mathbb{R}^d \) are denoted respectively \( \| \cdot \| \) and \( \langle \cdot, \cdot \rangle \). If \( X \) is a random vector with null expectation then \( \mathbb{E}(X \otimes X) = \mathbb{E}(XX') \) is the covariance matrix of \( X \). The covariance matrix is crucial since it determines the limiting Gaussian distribution in the Central Limit Theorem. We will need two kinds of stochastic convergence in the sequel:  convergence in probability of \( X_n \) to \( X \) is denoted \( X_n \to P X \). The arrow \( \hookrightarrow \) stands for convergence in distribution (weak convergence).

Except in section 2.3, we consider a single particle in order to alleviate notations. We drop the particle index so that \( x_n^\ell = x_n, p_n^\ell = p_n \) and \( g_n^\ell = g_n \).

At last we take for granted that particles are warm, reached an area of the domain were they fluctuate without exiting (condition \( A_1 \) below).

2.1 First case: oscillatory \((p \neq g)\)

Denote \( \bar{\xi}_n = \max \{|p_n - p|, |g_n - g|\} \) and:

\[
\theta = \frac{c_1 p + c_2 g}{c_1 + c_2}, \quad c = \frac{c_1 + c_2}{2}.
\]

The following assumptions are required and discussed after the statement of Theorem 1.

- \( A_1 \) : For all \( n, x_n \in \Omega \) where \( \Omega \) is a compact subset of \( \mathbb{R}^d \).
- \( A_2 \) : \( \sqrt{N} \bar{\xi}_n \to 0 \) when \( N \to +\infty \).
- \( A_3 \) : The inequality below connects \( c_1, c_1 \) and \( \omega \):

\[
2c \frac{1 - \omega}{1 + \omega} (1 + \omega - \frac{c_1}{2}) > \frac{c_1^2 + 2}{12}.
\]

Before stating the Theorem we need a last notation. Let \( \delta = (\delta_1, \ldots, \delta_d) \in \mathbb{R}^d \). The notation \( \text{diag}(\delta) \) stands for the diagonal \( d \times d \) matrix with entries \( \delta_1, \ldots, \delta_d \) and \( \delta^\otimes 2 \) is the vector in \( \mathbb{R}^d \) defined by \( \delta^\otimes 2 = (\delta_1^2, \ldots, \delta_d^2) \).

Theorem 1. Set \( \mathcal{C} = c \left( \frac{1 - \omega}{1 + \omega} \right) (1 + \omega - \frac{\delta^2}{2}) \) and:

\[
\mathcal{C} = \frac{1}{24} \frac{c_1 c_2}{\delta^2} \frac{\mathcal{C}}{\delta^2 - \frac{c_1^2 + 2}{12}}.
\]

Denote finally \( \Gamma = \mathcal{C} \cdot \text{diag}(p - g)^\otimes 2 \) then:

\[
\sqrt{N} \left( \frac{1}{N} \sum_{n=1}^N x_n - \theta \right) \hookrightarrow \mathcal{N}(0, \Gamma)
\]

where \( \mathcal{N}(0, \Gamma) \) denotes the Gaussian centered random vector of \( \mathbb{R}^d \) with covariance matrix \( \Gamma \).

Discussion of the Assumptions:

We avoid here the assumption of stagnation: the personal and local best are not supposed to be constant but they oscillate around their expectation. The convergence occurs at a rate ensuring that neither \( g_n \) nor \( p_n \) are involved in the weak convergence of the particles \( x_n \). Condition \( A_2 \) is specific of what we intend by a convergent PSO. It ensures that \( p_n \) and \( g_n \) have no impact on the weak convergence behavior of the particles. With other words Assumption \( A_2 \) requires that the oscillations of \( p_n \) and \( g_n \) around their expectations are negligible. We tried here to model the stagnation phenomenon which consists in sequence of iterations during which \( g_n \) (resp. \( p_n \)) remain constant for \( n \) in \( [N, \bar{N}] \) where \( N < \bar{N} \) are two integers.

Note that assumption \( A_3 \) is exactly the condition found in Poli (2009) (see the last paragraph of section III) for defining order 2 stability. This condition may be extended to the case when \( c_1 \neq c_2 \), see Cleghorn and Engelbrecht (2018) and references therein. At last \( A_3 \) holds for the classical calibration appearing in Clerc and Kennedy (2002) (constriction constraints) with \( c = 1.496172 \) and \( \omega = 0.72984 \).

The next Corollary provides finally two kinds of by-products: asymptotic confidence intervals, in \( \mathbb{R} \), for the coordinates of \( \theta \) and confidence regions for \( \theta \) in \( \mathbb{R}^d \). Let \( \alpha \in [0, 1] \), \( q_\alpha \) be the \( 1 - \alpha \) quantile of the standard Gaussian distribution and \( X_{d-1}^{(d)} \) be the \( 1 - \alpha \) quantile of the Chi-square distribution with \( d \) degrees of freedom. Set also \( \theta = (\theta_1, \ldots, \theta_d) \) and \( \bar{x}_N = (\bar{x}_{N,1}, \ldots, \bar{x}_{N,d}) \).

Corollary 2. Pick \( \ell \in \{1, \ldots, d\} \). An asymptotic (in \( N \)) confidence interval at level \( 1 - \alpha \) for \( \theta_\ell \) is directly derived from Theorem 1:

\[
\mathcal{J}_{1-\alpha}(\theta_\ell) = [\bar{x}_{N,\ell} - s_{\ell}(N, \alpha), \bar{x}_{N,\ell} + s_{\ell}(N, \alpha)]]
\]
with $s_t(N, \alpha) = |p_t - g_t| \sqrt{\frac{-q}{N}}$.
An asymptotic confidence (in $N$) region at level $1 - \alpha$ for the vector $\theta = (p + g)/2$ is:

$$\Lambda_d(1 - \alpha) = \left\{ t \in \mathbb{R}^d : N \| \Gamma^{-1/2} (t - \bar{x}_N) \|^2 \leq \chi^2_{d}(\alpha) \right\}$$

$$= \left\{ t = (t_1, \ldots, t_d) : \sum_{\ell=1}^d \left( \frac{t_\ell - \bar{x}_N,\ell}{p_\ell - g_\ell} \right)^2 \leq \frac{\chi^2_{d}(\alpha)}{N^{1-\alpha}} \right\}.$$

We note however that the vector $\theta$ may not be of crucial interest for the initial optimization problem conversely to $g$.
This point will be discussed in the last section.

### 2.2 Second case: non-oscillatory and stagnant ($p = g$)

In this section we study again a single particle and suppose once and for all that $x_n \in \mathbb{R}$. We assume throughout this subsection that the particle is under stagnation that is $p_n = p$ for $n$ sufficiently large (see assumption $B_2$ below). This assumption is strong but a more general framework leads to theoretical developments out of our scope. Starting from Equation (1), the PSO equation becomes this time:

$$x_{n+1} = (1 + \omega) x_n - \omega x_{n-1} + c (r_1 + r_2) (p - x_n).$$

Change the centering and consider $x_n - p = y_n$. The previous equation becomes:

$$y_{n+1} = (1 + \omega - c + c\varepsilon_n) y_n - \omega y_{n-1}, \quad (2)$$

where $\varepsilon_n$ is the sum of two independent random variables with $\mathcal{U}[-1/2; 1/2]$ distribution.
Assuming that for all $n y_n \neq 0$, we have then:

$$\frac{y_{n+1}}{y_n} = (1 + \omega - c + c\varepsilon_n) - \omega \frac{y_{n-1}}{y_n} = \frac{y_{n+1}}{y_n} - \omega \frac{y_{n-1}}{y_n} \quad (3)$$

It is plain that $y_{n+1}/y_n$ defines a Markov chain (more precisely: a non-linear auto-regressive process) which will play a crucial role in the forthcoming results. It is shown in the proof section that $y_{n+1}/y_n$ is Harris recurrent and has...
consequently a stationary distribution denoted $\pi$. The definition of Harris recurrence needed is given for instance in [Meyn and Tweedie (2012)], beginning of Chapter 9. Take $Z_n$ a copy of $y_n/y_{n-1}$ with $Z_0$ a realization of $\pi$. Then define:

$$
\mu_x = \mathbb{E}_\pi \log |Z_0|,
$$

$$
\sigma^2_x = \text{Var}_\pi (\log |Z_0|) + 2 \sum_{k=1}^{+\infty} \text{Cov}_\pi (\log |Z_0|, \log |Z_k|).
$$

We are ready to introduce a new set of assumptions.

- **B$_1$**: $1 + \omega - c < \omega/c < (1 + c)/4$.
- **B$_2$**: For sufficiently large $n$ $g_n = p_n = p = g$ is constant.

Before stating next Theorem notice that

$$
\sum_{n=1}^{N} \log \left| \frac{y_n}{y_{n-1}} \right| = \log |y_N| - \log |y_0|.
$$

**Theorem 3.** Let $\omega \in (0,1)$, $c > 1$, when $\text{B$_1 - 2$}$ hold, then:

$$
\frac{1}{\sqrt{N}} (\log |x_N - g_N| - N\mu_x) \Rightarrow \mathcal{N} (0, \sigma^2_x)
$$

when $N$ tends to infinity.

**Remark 1.** The theorem above is not a Central Limit Theorem for $x_N$. It is derived thanks to a CLT but it shows that the asymptotic distribution of $|x_N - g_N|$ is asymptotically log-normal with approximate parameters $N\mu_x$ and $N\sigma^2_x$.

**Remark 2.** The mean and variance $\mu_x$ and $\sigma^2_x$ are usually unknown but may be approximated numerically. We refer to the simulation section for more details.

**Corollary 4.** If $p_n = p$ for all $n$ (pure stagnation) an asymptotic non convex confidence region for $g$ at level $1 - \alpha$ denoted $\Lambda_{1-\alpha}$ below may be derived from the preceding Theorem:

$$
\Lambda_{1-\alpha} (g) = \Lambda_1^{+} \cup \Lambda_1^{-}
$$

$$
\Lambda_1^{+} = \left[ x_n + \exp \left( \mu_x + \frac{\sigma_x}{\sqrt{n}} \varepsilon^{1/2} \right), x_n + \exp \left( \mu_x - \frac{\sigma_x}{\sqrt{n}} \varepsilon^{1-\alpha/2} \right) \right]
$$

$$
\Lambda_1^{-} = \left[ x_n - \exp \left( \mu_x + \frac{\sigma_x}{\sqrt{n}} \varepsilon^{1-\alpha/2} \right), x_n - \exp \left( \mu_x - \frac{\sigma_x}{\sqrt{n}} \varepsilon^{1/2} \right) \right]
$$

Figure 2: Display of constraint B$_1$ in the plane $(c, \omega)$. 
Remark 3. Under matrix form the equation (2) is purely linear but driven by a random matrix:

\[
\begin{pmatrix}
    y_{n+1} \\
y_n
\end{pmatrix} =
\begin{pmatrix}
    1 + \omega - c (1 + \varepsilon_n) & -\omega \\
    1 & 0
\end{pmatrix}
\begin{pmatrix}
    y_n \\
y_{n-1}
\end{pmatrix}
\]

with \( T_n = S_1 S_{n-1} \ldots S_2 \). It is plain here that a classical Central Limit Theorem cannot hold for the sequence \((y_n)_{n \in \mathbb{N}}\). In the proofs we turn to asymptotic theory for the product of random matrices. We refer to the historical references: Furstenberg and Kesten (1960) and Berge (1984) who proved Central Limit Theorems for the regularity index of the product of i.i.d random matrices. Later Hennion (1997) generalized their results. But the assumptions of (almost surely) positive entries is common to all these papers. Other authors obtain similar results under different sets of assumptions (see Le Page (1982), Benoist and Quint (2016), and references therein), typically revolving around characterization of the semi-group spanned by the distribution of \( S_n \). These assumptions are uneasy to check here and we carried out to a direct approach with Markov chain fundamental tools.

2.3 The swarm at a fixed step

In this section we change our viewpoint. Instead of considering a single particle and sampling along its trajectory we will take advantage of the whole swarm but at a fixed and common iteration step. Our aim here is to localize the minimization of the cost function based on \((x_1^n, \ldots, x_S^n)\). This time the particle index \( s \) varies up to \( S \) the swarm size, whereas the index \( n \) is fixed. In this subsection we assume that \( S \uparrow +\infty \) and asymptotic is with respect to \( S \). We do not drop \( n \) in order to see how the iteration steps influence the results. We still address only the case \( x_i^n \in \mathbb{R} \) even if our results may be straightforwardly generalized to \( x_i^n \in \mathbb{R}^d \). We provide below a Central Limit Theorem suited to the case when the number of particles in the swarm becomes large. In order to clarify the method, we assume that for all particles \( x_i^n \) in the swarm \( p^n = g_n = p \). In other words, no local minimum stands in the domain \( \Omega \), which implies additional smoothness or convexity assumptions on the cost function \( f \). This may be possible by a preliminary screening of the search space. Indeed a first (or several) run(s) of preliminary PSO(s) on the whole domain identifies an area where a single optimum lies. Then a new PSO is launched with initial values close to this optimum and with parameters ensuring that most of the particles will stay in the identified area.

So we are given \((x_1^n, \ldots, x_S^n)\) where \( S \) is the sample size. Basically, the framework is the same as in the non oscillatory case studied above for a single particle. From (4) we get with \( y_i^n = x_i^n - p \):

\[
\begin{align*}
y_i^n &= T_i^n y_i^n, \\
T_i^n &= \Pi_{j=2}^S S_j, \\
S_j &= \begin{pmatrix} 1 + \omega - c (1 + \varepsilon_j) & -\omega \\
1 & 0 \end{pmatrix}.
\end{align*}
\]

Assume that the domain \( \Omega \) contains 0 and that for all \( s \) \((x_0, x_1^s)_{s < S}\) are independent, identically distributed and centered then from the decomposition above, for all \( n \) and \( s \), \( E y_i^n = 0 \) and the \((y_i^n)_{1 \leq i \leq S}\) are i.i.d too.

The assumptions we need to derive Theorem 5 below are:

- \( C_1 \) : The operational domain \( \Omega \) contains 0 (and is ideally a symmetric set).
- \( C_2 \) : The couples \((x_0, x_1^s)_{s < S}\) are i.i.d and centered.
- \( C_3 \) : For all \( i \in \{1, \ldots, S\} \) \( p^n = g_n = p \).

When \( S \) is large the following Theorem may be of interest and is a simple consequence of the i.i.d. CLT.

**Theorem 5.** Under assumptions \( C_{1-3} \) a Central Limit Theorem holds when \( S \) the number of particles in the swarm tends to \( +\infty \):

\[
\frac{1}{\sqrt{S}} \sum_{i=1}^S (x_i^n - g_n) \xrightarrow{S \to +\infty} \mathcal{N}\left(0, \sigma_n^2\right),
\]

where \( \sigma_n^2 = E (x_i^n - g_n)^2 \) is estimated consistently by:

\[
\hat{\sigma}_n^2 = \frac{1}{S} \sum_{i=1}^S (x_i^n - g_n)^2.
\]

**Remark 4.** The convergence of \( \hat{\sigma}_n^2 \) to \( \sigma_n^2 \) is a straightforward consequence of the weak and strong laws of large numbers.

Denote \( \Xi = S_{1 \leq i \leq S} x_i^n \). The Theorem above paves the way towards an asymptotic confidence interval.

**Corollary 6.** An asymptotic confidence interval at level \( 1 - \alpha \) for \( g \) is:

\[
\Lambda_n (g) = \left[ \Xi \pm \frac{\hat{\sigma}_n}{\sqrt{S}} q_{1-\alpha/2}, \Xi \pm \frac{\hat{\sigma}_n}{\sqrt{S}} q_{1-\alpha/2} \right].
\]
3 Simulation and numerical results

All along this section we take $c_1 = c_2 = c$.

The Himmelblau’s function is chosen as example for our experiments. It is a 2-dimensional function with four local optima in $[-10, 10]^2$ defined by:

$$f(x, y) = (x^2 + y - 11)^2 + (x + y^2 - 7)^2.$$  

Figure 3 illustrates the contour of this function.

With the Himmelblau’s function, we can observe the two different behaviors of the particles: oscillatory and non-oscillatory. The Himmelblau’s function is positive and has four global minima in:

$$(3, 2), (-2.81, 3.13), (-3.77, -3.28), (3.58, -1.84)$$ where $f(x, y) = 0$. We use a ring topology (for a quick review of the different topologies of PSO see Lane et al. (2008)) for the algorithm in order to have both oscillating and non-oscillating particles. The latter converge quickly. The former go on running between two groups of particles converging to two distinct local optima.

3.1 Oscillatory case

We select particles oscillating between $(3.58, -1.84)$ and $(3, 2)$, these values could be both their personal best position or their neighborhood best position. In this case, the convergence of the $p_n$ and $g_n$ to $(3.58, -1.84)$ or $(3, 2)$ satisfies the conditions of Theorem 1. We have to verify that the Gaussian asymptotic behavior of $H_s^1(N) = \sqrt{N} \left( \frac{1}{N} \sum_{n=1}^{N} s_n - \frac{p_s}{2} \right)$ for each $s$ oscillating particle.

We launch PSO with a population of 200 particles and with 2000 iterations, $\omega = 0.72984$ and $c = 1.496172$. A ring topology was used to ensure the presence of oscillating particles. A particle is said oscillating if between the 500th and the 2000th iteration, Assumptions $A_{1-3}$ holds.

A visual tool to verify the normality of $H_s^1(N)$ for a particle is a normal probability plot. Figures 4 and 5 displays the normal probability plot of $H_s^1(N)$ respectively for the $x$ axis and $y$ axis. For each axis, the normality is confirmed: $H_s^1(N)$ fits well the theoretical quantiles.

To check the formula of the covariance matrix $\Gamma$, the confidence ellipsoid is also a good indicator (see Figure 6). For a single particle, $H_s^1(N)$ is not necessarily always inside the confidence ellipsoid and does not respect the percentage of
Figure 6: Trajectory of $H_1^t(N)$ for an oscillating particle in $[-10, 10]^2$. The confidence ellipsoid at a level of 85% is displayed in red. Around 99% of the trajectory of $H_1^t(N)$ is inside the ellipse.

the defined confidence level. Figure 7 shows the trajectory of $x^t_n$ and $H_1^t(N)$ on the y axis, $H_1^t(N)$ remains bounded.

With 200 Monte-Carlo simulations of PSO (200 particles, and 2000 iterations), we select all the particles oscillating between (3.58, −1.84) and (3, 2), and for each of them we compute $H_1^t(2000)$. Figure 8 displayed the density of $H_1^t(2000)$ using 1150 oscillating particles. Almost 95% of the particles are inside the confidence ellipsoid of level 95% (represented in red).
Figure 7: Top track: Trajectory of a oscillating particle on the y axis. The particle is oscillating between 2 and -1.84. Bottom track: corresponding trajectory of $H_1^s(N)$ on the y axis, the red dot are corresponding to the 95% confidence interval. The trajectory of $H_1^s(N)$ is well bounded.

Figure 8: Density of $H_1(2000)$ with 1150 particles issued from Monte-Carlo simulations. The red ellipse is the 95% confidence ellipsoid.
3.2 Non-oscillatory case

We study now the behaviors of non-oscillating particles on the Himmelblau’s function. We launch PSO with a population of 1000 particles and with 2000 iterations, \( \alpha = 0.72984 \) and \( c = 1.496172 \). A ring topology was used to ensure the presence of enough particle converging to each local optimum. We select particles converging to (3, 2), meaning that \( p_n = g_n = p \) for a sufficiently large \( n \). For the weak convergence of the particle, we consider:

\[
H_2^3(N) = \frac{1}{\sqrt{N}} (\log |x_n^g - g_n| - N\mu_x).
\]

First, it is easy to check the linear dependency of \( \log |x_n^g - g_n| \) with a single display of the trajectory. Figure 9 illustrates this phenomenon for a single particle. We observed numerical issues when we reach the machine precision, but a numerical approximation of \( \mu_x \) can be performed thanks to a linear regression.

Using many converging particles, a Monte Carlo approximation of \( \mu_x \) is done. For the approximation of \( \sigma_x \), a possibility is:

\[
\sigma_x^2 = \text{Var}_\pi (\log |X_0|) + 2 \sum_{k=1}^T \text{Cov}_\pi (\log |X_0|, \log |X_k|),
\]

where \( T = 20 \). With near 240 converging particles to (3, 2), we found that for the first coordinate:

\[
\mu_x = -0.032, \quad \sigma_x = 0.156.
\]

We verify the asymptotic normality of \( H_2(N) \) with a normal probability plot using the approximation of \( \mu_x \). Figure 10 displays the normal probability plot of \( H_2(N) \) on the first coordinate, the theoretical quantiles are well fitted by \( H_2(N) \). Figure 11 illustrates different trajectories of \( H_2(N) \) on the first coordinate which are bounded by the 95% confidence interval deduced from \( \sigma_x \).

3.3 Swarm at a fixed step

To check the Theorem 5, we study:

\[
H_2^3(S) = \frac{1}{\sqrt{S\sigma_n}} \sum_{i=1}^S (s_i^g - g_n).
\]

In practice, we encountered some difficulties to verify Theorem 5 because of the convergence rate of the particles. Indeed, when \( p_n = g_n = p \), the particle \( s \) converges exponentially to \( g_n \) but the spread of the rate of convergence is large. As a consequence, at a fixed step of PSO, some particles could be considered as outliers because of a lower rate of convergence. Because of these particles qualified as belated, the asymptotic Gaussian behavior of \( H_2^3(S) \) is not verified. A solution is to filter the particles and remove the belated particles. Figures 12 and 13 illustrate this phenomenon for the Himmelblau’s function in 2D and with near 1500 converging particles to (3, 2) over 500 iterations. In Figure 12, we compute without any filtering \( H_2^3(S) \) and we notice that the Gaussian behavior is not verified and some jumps appeared. The presence of these “jumps” is due to belated particles which have a lower rate of convergence in comparison to the swarm. When we remove these particles with a classical outliers detection algorithm in Figure 13, Theorem 5 seems to be verified.

4 Conclusion

Our main theoretical contribution revolves around the three CLTs and the confidence regions derived from sampling either a particle’s path or the whole swarm. Practically the confidence set \( \Lambda_{1-\alpha} \) localizes, say \( g \), with probability \( 1 - \alpha \). The simulations carried out in Python language tend to foster the theory.

This work was initiated in order to solve a practical issue in connection with oil industry and petrophysics. Yet in the previous section we confined ourselves to simulated data for several reasons. It appears that our method should be applied on real-data for a clear validation.

A second limitation of our work is the asymptotic set-up. The results are stated for samples large enough. This may not have an obvious meaning for people not familiar with statistics or probability. Practitioners usually claim that the Central Limit Theorem machinery works well for a sample size larger than thirty/forty whenever data are stable, say stationary (path with no jumps...). As a consequence the first iterations of the algorithm should be avoided to ensure a warming. When studying PSO the behavior of \( p_n \) and \( g_n \) turns out to be crucial too in order to ensure this stability, hence the validity of the theoretical results. However the control of \( p_n \) and \( g_n \) is a difficult matter, which explains our current assumptions on stagnation or “almost stagnation”. However, in order to fix the question of the asymptotic versus non-asymptotic approach a work is on progress. We expect finite-sample concentration inequalities: the non-asymptotic counterparts of CLT.
Figure 9: $|x_n - g_n|$ over 500 iterations in a logarithmic scale. After 300 iterations, we reach the computer precision. We take advantage of the linear behavior of $\log(|x_n - g_n|)$ on the 200 first iterations to perform a linear regression estimating $\mu_x$.

Figure 10: Normal probability plot of $H_2(N)$ on the first coordinate.

Figure 11: Trajectories of $H_2(N)$ on the first coordinate for five particles. The red dot represents the 95% confidence interval deduced from $\sigma_x$. Trajectories stop around the 400 iterations (after an heating phase) due to numerical precision.
In the oscillating framework our results involve the parameter denoted \( \theta \), the center of the segment \([p, g]\). Here we miss the target \( g \) for instance. We claim that our method may be adapted to tackle this issue but we will not elaborate on it. Briefly speaking: running two independent PSOs on the same domain with two distinct pairs \((c_1, c_2)\) and \((c_1', c_2')\) will provide, under additional assumptions on the local and global minima of the cost function, two CLTs involving:

\[
\theta = \frac{c_1 g + c_2 p}{c_1 + c_2} \quad \text{and} \quad \theta' = \frac{c_1' g + c_2' p}{c_1' + c_2'}.
\]

A simple matrix inversion will then recover a joint CLT for the couple \((g, p)\) since, when \(c_2c_1' \neq c_1c_2'\) then:

\[
p = \frac{c_2' \theta - c_1 \theta'}{c_2c_1' - c_1c_2'}
\]

for instance.

## 5 Derivations of the results

We start with some notations. First we recall that the sup-norm for square matrices of size \(d\) defined by:

\[
\|M\|_{\sup} = \sup_{x \neq 0} \frac{\|Mx\|}{\|x\|}.
\]

The tensor product notation is appropriate when dealing with special kind of matrices for instance covariance matrices. Let \(u\) and \(v\) be two vectors of \(\mathbb{R}^d\) then \(u \otimes v = uv'\) (where \(v'\) is the transpose of vector \(v\)) stands for the rank one matrix defined for all vector \(x\) by \((u \otimes v)(x) = (v,x)u\). Besides \(\|u \otimes v\|_{\sup} = \|u\|\|v\|\). The Hadamard product between vectors was mentioned earlier. Its matrix version may be defined a similar way. Let \(M\) and \(S\) be two matrices with same size then \(M \circ S\) is the matrix whose \((i, j)\) cell is \((M \circ S)_{i,j} = m_{i,j} s_{i,j}\). We recall without proof the following computation rule mixing Hadamard and tensor product. Let \(\eta, \epsilon, u\) and \(v\) be four vectors in \(\mathbb{R}^d\). Then:

\[
(\eta \otimes u) \otimes (\epsilon \otimes v) = (\eta \otimes \epsilon) \otimes (u \otimes v),
\]

and the reader must be aware that the Hadamard product on the left-hand side operates between vectors whereas on the right-hand side it operates on matrices.

We will need \(\mathcal{F}_n^q\) the filtration generated by the path of particle number \(s\) up to step \(n\) : \(\{x_{0}^{s}, ..., x_{n}^{s}\}\) and \(\mathcal{F}_n^p\) the filtration generated by the swarm up to step \(n\) : \(\{x_{0}^{n}, ..., x_{n}^{n} : s = 1, ..., S\}\).

We will also denote later \(g_n = \mathbb{E}(g_n) + \xi_n\) and \(p_n = \mathbb{E}(p_n) + \nu_n\) the expectation-variance decomposition of \(g_n\) and \(p_n\) where \(\xi_n\) and \(\nu_n\) are centered random vectors and support all the variability of \(g_n\) and \(p_n\) respectively.

### 5.1 First case: oscillatory

In this subsection we prove Theorem[1]. We start with two Lemmas and a Proposition who will be invoked later.
**Lemma 7.** Let \( \varepsilon_n^{(i)} \) and \( \eta_n^{(j)} \) be any coordinate of the random vectors \( \varepsilon_n \) and \( \eta_n \) appearing in ???. Clearly, \( \varepsilon_n^{(i)} \) and \( \eta_n^{(j)} \) are not independent but not correlated and follow the same type of distribution. Besides:

\[
\mathbb{E}\varepsilon_n^{(i)} = \mathbb{E}\eta_n^{(j)} = 1/6, \quad \mathbb{E}\varepsilon_n^{(i)} \eta_n^{(j)} = 0,
\]

\[
\mathbb{E}\varepsilon_n^{(i)} \eta_n^{(j)} = \mathbb{E}\varepsilon_n^{(i)} \eta_n^{(j)} = 0.
\]

The proof is very simple hence omitted.

The next Lemma is also straightforward but will be frequently invoked later.

**Lemma 8.** Let \( E_n \) be a sequence of i.i.d centered random matrices with finite moment of order 4, let \( u_n \) and \( v_n \) two sequence of random vectors almost surely bounded and such that \( (u_n, v_n) \) is for all \( n \) independent from \( E_n \) then for the \( \| \cdot \|_\infty \) norm:

\[
\frac{1}{N} \sum_{n=1}^{N} E_n \to u_n \to v_n \rightarrow 0.
\]

The proof is a simple application of Cauchy-Schwartz inequality.

**Proof of Theorem 1**

**First step:** We show that the problem of convergence in distribution for \( \frac{1}{N} \sum_{n=1}^{N} (x_n - \theta) \) may be shifted to proving a CLT for a martingale difference array.

Starting from the initial PSO equation:

\[
x_{n+1} - x_n = \omega (x_n - x_{n-1}) + c_1 r_{1,n} \circ (p_n - x_n) + c_2 r_{2,n} \circ (g_n - x_n)
\]

it is simple to derive:

\[
x_{n+1} - x_n = \omega (x_n - x_{n-1}) + c_1 r_{1,n} \circ (p_n - x_n) + c_2 r_{2,n} \circ (g_n - x_n) + R_n
\]

with

\[
R_n = c_1 r_{1,n} \circ (p_n - p) + c_2 r_{2,n} \circ (g_n - g).
\]

Denote \( z_n = (x_n - \theta) \) then after some additional calculations we get:

\[
z_{n+1} - (1 + \omega - c) z_n + \omega z_{n-1} = c_1 e_{1,n} \circ \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right) - c_2 e_{2,n} \circ \left( \frac{c_1}{c_1 + c_2} (p - g) + z_n \right) + R_n
\]

with, for \( i = 1, 2 \):

\[
r_i,n = \mathbb{E} r_{i,n} + \varepsilon_{i,n} = 1/2 + \varepsilon_{i,n} \text{ and } \varepsilon_{i,n} \text{ is the centered version of } r_{i,n}.
\]

Summing the above equation from \( n = 1 \) to \( n = N - 1 \):

\[
\sum_{n=1}^{N} z_n = \frac{1}{c} \sum_{n=1}^{N} u_n + \mathcal{R}_N
\]

with:

\[
u_n = c_1 e_{1,n} \circ \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right) - c_2 e_{2,n} \circ \left( \frac{c_1}{c_1 + c_2} (p - g) + z_n \right),\]

\[
\mathcal{R}_N = \frac{1}{c} \left[ \sum_{n=1}^{N} R_n + z_1 + (c - 1) z_N + \omega (z_{N-1} - z_0) \right]
\]

Notice that \( u_n \) just above is a martingale difference with respect to the nested filtration \( \mathcal{F}_n : \mathbb{E} (u_n | \mathcal{F}_{n-1}) = 0 \) because \( e_{1,n} \) and \( e_{2,n} \) are centered and independent from \( z_n \).

In view of conditions \( A_1 \) and \( A_2 \):

\[
\mathcal{R}_N \frac{1}{\sqrt{N}} \rightarrow p 0,
\]

then from (7) \( \frac{1}{\sqrt{N}} \sum_{n=1}^{N} z_n \) converges in distribution if and only if \( \frac{1}{\sqrt{N}} \sum_{n=1}^{N} u_n \) does.

**Second step:** it is shown below that weak convergence for \( \frac{1}{\sqrt{N}} \sum_{n=1}^{N} u_n \) holds under the assumptions of the Theorem.

We aim at proving a Levy-Lindeberg version of the CLT for the series of \( u_n \) in two steps (Theorem 2.1.9 p. 46 and its corollary 2.1.10 in Duflo (1997)): first checking the Lyapunov condition holds (hence the Lindeberg’s uniform integrability that ensures uniform tightness of the sequence) then ensuring convergence in probability of the conditional covariance structure of \( \frac{1}{\sqrt{N}} \sum_{n=1}^{N} u_n \).

Here the Lyapunov condition holds trivially because we are faced with bounded martingale difference sequences. Indeed by assumption \( A_1 \) and since the random variables \( e_{1,n} \) and \( e_{2,n} \) are almost surely bounded \( u_n \) is itself almost surely bounded.
We turn to the conditional covariance sequence of $u_n$:

$$\Gamma_N = \frac{1}{N} \sum_{n=1}^{N} \mathbb{E}(u_n \otimes u_n | \mathcal{F}_{n-1}).$$

Starting from the definition of $u_n$ just below (7) leads after simple calculations to:

$$\mathbb{E}(u_n \otimes u_n | \mathcal{F}_{n-1}) = \frac{c_1^2 + c_2^2}{6} \mathbb{I} \odot [(p - g) \otimes (p - g)] + \frac{c_1 + c_2}{12} \mathbb{I} \odot (z_n \otimes z_n)$$

$$- 2 \frac{c_1 - c_2}{c_1 + c_2} \mathbb{I} \odot [(p - g) \otimes z_n]$$

Notice that, due to the operation $\mathbb{I} \odot (.) \mathbb{E}(u_n \otimes u_n | \mathcal{F}_{n-1})$ is a diagonal matrix with:

$$\tau_{i,n} = \frac{c_1^2 + c_2^2}{6(c_1 + c_2)^2}(p_i - g_i)^2 + \frac{c_1 + c_2}{12} z_{i,n}^2 - \frac{c_1 - c_2}{c_1 + c_2} \frac{c_1 c_2}{6} [(p_i - g_i) z_n]$$

Averaging the previous results in $n$ we get:

$$\Gamma_N = \frac{c_1^2 + c_2^2}{6} \mathbb{I} \odot [(p - g) \otimes (p - g)] + \frac{c_1 + c_2}{12} \mathbb{I} \odot \left( \frac{1}{N} \sum_{n=1}^{N} z_n \otimes z_n \right)$$

$$- 2 \frac{c_1 - c_2}{c_1 + c_2} \mathbb{I} \odot \left( \frac{1}{N} \sum_{n=1}^{N} z_n \right) \otimes \left( \frac{1}{N} \sum_{n=1}^{N} z_n \right)$$

The last term is negligible: indeed Lemma[8] shows that $\frac{1}{N} \sum_{n=1}^{N} z_n \to p \neq 0$. Our only task consists in studying convergence for $\mathbb{I} \odot \left( \frac{1}{N} \sum_{n=1}^{N} z_n \otimes z_n \right)$. More precisely we just need to prove convergence for $(1/N) \sum_{n=1}^{N} z_{i,n}^2$ for all $i \in \{1, \ldots, d\}$.

**Third (and last) step**: Proof of convergence of $\frac{1}{N} \sum_{n=1}^{N} z_{i,n}^2$ for all $i$.

Let us restart from (6):

$$z_{i,n+1} = (1 + \omega - c) z_{i,n} - \omega z_{i,n-1} + u_{i,n}$$

Clearly $R_n$ is a negligible term due to assumption $A_2$. All the authors’ computations were effectively carried out with $R_n$ but in order to alleviate this proof we will proceed here as if:

$$z_{i,n+1} = (1 + \omega - c) z_{i,n} - \omega z_{i,n-1} + u_{i,n}$$

From now on we drop the dimension index $i$ in the previous equation:

$$z_{i,n+1}^2 = (1 + \omega - c)^2 z_{i,n}^2 + \omega^2 z_{i,n-1}^2 + u_{i,n}^2$$

$$+ 2 [(1 + \omega - c) z_{i,n} - \omega z_{i,n-1}] u_{i,n} - 2 \omega (1 + \omega - c) [z_{i,n} z_{i,n-1}]$$

The following equation is derived from the preceding one after summing and considering the $j^{th}$ component of each vector:

$$\left(1 - (1 + \omega - c)^2 - \omega^2\right) \sum_{n=1}^{N} z_{i,n}^2 = \sum_{n=1}^{N} u_{i,n}^2 - 2 \omega (1 + \omega - c) \sum_{n=1}^{N} z_{i,n} z_{i,n-1}$$

$$+ 2 \sum_{n=1}^{N} [(1 + \omega - c) z_{i,n} - \omega z_{i,n-1}] u_{i,n} + z_{i,n}^2 - z_{i,n+1}^2 + \omega^2 \left(z_{i,0}^2 - z_{i,N}^2\right)$$

(8)

It is easily seen that Lemma[8] may be applied to the series with term $[(1 + \omega - c) z_{i,n} - \omega z_{i,n-1}] u_{i,n}$ and its average tends to zero.

From now on we remove all indices $i$ for components of the vectors in order to alleviate notations.

Our next task consists in studying $\sum_{n=1}^{N} z_{i,n} z_{i,n-1}$. To that purpose we turn back again to equation (6) and multiply by $z_{i,n}$ the whole equation and sum over $n$:

$$\sum_{n=1}^{N} z_{i,n+1} z_{i,n} - (1 + \omega - c) \sum_{n=1}^{N} z_{i,n}^2 + \omega \sum_{n=1}^{N} z_{i,n} z_{i,n-1} = \sum_{n=1}^{N} z_{i,n} u_{i,n}$$

hence

$$(1 + \omega) \sum_{n=1}^{N} z_{i,n} z_{i,n-1} = (1 + \omega - c) \sum_{n=1}^{N} z_{i,n}^2 + \sum_{n=1}^{N} z_{i,n} u_{i,n} + z_{i,0} - z_{i,N+1}$$

$$= (1 + \omega) \sum_{n=1}^{N} z_{i,n} z_{i,n-1}$$

14
The same arguments as above show that \((1/N) \sum_{n=1}^N z_n \eta_n\) tends to 0 in probability by the weak law of large numbers. This proves that when \(N \to +\infty:\)

\[
\frac{1}{N} \left\{ \sum_{n=1}^N z_n^\omega - \frac{1 + \omega - c}{1 + \omega} \sum_{n=1}^N z_n^2 \right\} \to_p 0
\]

Plugging this result in (8) and taking into account the remarks above:

\[
\lim_{N \to +\infty} \frac{1}{N} \left\{ \left( 1 - (1 + \omega - c)^2 - \omega^2 \right) + 2\omega \frac{(1 + \omega - c)^2}{1 + \omega} \right\} \sum_{n=1}^N \frac{z_n^2}{z_n^2} = 0 \tag{9}
\]

Our last task consists in finding \(\lim_{N \to +\infty} (1/N) \sum_{n=1}^N u_n^2\). Calculations are closed to those carried out at the beginning of the second step but this time without conditioning with respect to \(\mathcal{F}_n\). Remind that the index \(i\) was removed for the sake of clarity:

\[
u_n^2 = \left[ c_1 \varepsilon_{1,n} \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right) - c_2 \varepsilon_{2,n} \left( \frac{c_1}{c_1 + c_2} (p - g) + z_n \right) \right]^2
\]

Developing the square into bracket and taking into account that the cross product will tend to zero after averaging in probability, we can focus only on:

\[
v_{1,N} = \frac{c_1}{N} \sum_{n=1}^N \varepsilon_{1,n}^2 \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right)^2 \quad \text{and} \quad v_{2,N} = \frac{c_2}{N} \sum_{n=1}^N \varepsilon_{2,n}^2 \left( \frac{c_1}{c_1 + c_2} (p - g) + z_n \right)^2
\]

Both \(v_{1,N}\) and \(v_{2,N}\) may be managed the same way. We develop only for \(v_{1,N}\):

\[
v_{1,N} = \frac{c_1}{N} \sum_{n=1}^N \left( \varepsilon_{1,n}^2 - \mathbb{E} \varepsilon_{1,n}^2 \right) \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right)^2 + \frac{c_1}{N} \sum_{n=1}^N \mathbb{E} \varepsilon_{1,n}^2 \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right)^2
\]

The first series above tends in probability to 0 again by the weak law of large numbers for martingale difference arrays already mentioned several times earlier in the proof. The second series depends only on the \(z_i\)’s. We finally get:

\[
\lim_{N \to +\infty} v_{1,N} = \lim_{N \to +\infty} \frac{c_1^2}{12N} \sum_{n=1}^N \left( \frac{c_2}{c_1 + c_2} (p - g) - z_n \right)^2
\]

Similarly:

\[
\lim_{N \to +\infty} v_{2,N} = \lim_{N \to +\infty} \frac{c_2^2}{12N} \sum_{n=1}^N \left( \frac{c_1}{c_1 + c_2} (p - g) + z_n \right)^2
\]

Putting the two last equations with (9) and (10) we conclude that:

\[
\lim_{N \to +\infty} \left\{ \left( 1 - (1 + \omega - c)^2 - \omega^2 \right) + 2\omega \frac{(1 + \omega - c)^2}{1 + \omega} \right\} \sum_{n=1}^N \frac{z_n^2}{z_n^2} = \frac{1}{6} \left( \frac{c_1c_2}{c_1 + c_2} \right)^2 (p - g)^2
\]

However the previous equation is valid only if the big term between brackets in the right hand side is strictly positive. This term depends only on \(c_1, c_2\) and \(\omega\). This compatibility condition ensures just that a non-null limit for \(\sum_{n=1}^N z_{ni}^2\) exists. Besides:

\[
1 - (1 + \omega - c)^2 - \omega^2 + 2\omega \frac{(1 + \omega - c)^2}{1 + \omega} - \frac{c_1^2 + c_2^2}{12} = 2c \frac{1 - \omega - c}{1 + \omega} \left( 1 + \omega - \frac{c}{2} \right) - \frac{c_1^2 + c_2^2}{12}
\]

and the positivity condition matches assumption \(A_3\) within the Theorem.

This yields:

\[
\lim_{N \to +\infty} \frac{1}{N} \sum_{n=1}^N z_{ni}^2 = \frac{1}{6} \left( \frac{c_1c_2}{c_1 + c_2} \right)^2 \left\{ 2c \frac{1 - \omega - c}{1 + \omega} \left( 1 + \omega - \frac{c}{2} \right) - \frac{c_1^2 + c_2^2}{12} \right\}^{-1} (p - g)^2
\]

hence from (??) and with simple algebra:

\[
\Gamma_N \to_p \mathbb{R} \text{diag} (p - g)^2
\]
with 
\[ R = \frac{1}{24} \left( \frac{c_1 c_2}{c} \right)^2 \left[ 2c \frac{1 - \omega}{1 + \omega} \left( 1 + \omega - \frac{c}{2} \right) \right]^{-1} \]

All that was done until this point ensures that \( \frac{1}{N} \sum_{n=1}^{N} u_n \rightarrow Z \) where \( Z \) is a centered Gaussian random vector in \( \mathbb{R}^d \) with covariance operator \( (p - g)^{\otimes 2} \).

At last (7) leads to the desired result : \( \frac{1}{N} \sum_{n=1}^{N} z_n \rightarrow Z/c. \)

### 5.2 Second case: non-oscillatory and stagnant

We start from (3) and set \( X_n = y_n/y_{n-1} \). Then we get:

\[ X_{n+1} = 1 + \omega - c - \frac{\omega}{X_n} + c \epsilon_n, \]

where \( \epsilon_n = r_{1,n} + r_{2,n} - 1 \) has a “witch hat” distribution (convolution of two uniform distributions) with support \([-1, +1]\).

We focus now on the above homogeneous Markov chain \( X_n \) and we aim at proving that a CLT holds for \( h(X_n) = \log |X_n| \) namely that for some \( \mu \) and \( \sigma^2 \):

\[ \sqrt{N} \left[ \frac{1}{N} \sum_{n=1}^{N} \log |X_n| - \mu \right] \rightarrow \mathcal{N}(0, \sigma^2), \]

which will yield:

\[ \sqrt{N} \left[ \frac{1}{N} \sum_{n=1}^{N} \log |y_n| - \mu \right] \rightarrow \mathcal{N}(0, \sigma^2). \]

We aim at applying Theorem 1 p. 302 in Jones (2004) (see also Meyn and Tweedie (2012), section 17.5 for similar results). We need to check three points: (i) \( X_n \) is Harris ergodic, (ii) the existence of a small set \( \mathcal{C} \) and (iii) of a function \( g \) with a drift condition (see Meyn and Tweedie (2012)) such that (12) below holds. Following Hairer and Mattingly (2011), section 3.2 it turns out that proving (ii) and (iii) above with a drift function \( g \) but without the minoration bound (12) is sufficient to ensure (i).

Denote \( P(t, x) \) the transition kernel of \( X_n \). It is plain that \( P(t, x) \) coincides with the density of the uniform distribution on the set:

\[ \mathcal{E}_\tau = \left[ 1 + \omega - 2c - \frac{\omega}{x}, 1 + \omega - \frac{\omega}{x} \right]. \]

Theorem 3 is a consequence of the two Lemmas below coupled with the above-mentioned Theorem 1 p.302 in Jones (2004).

**Lemma 9.** Take \( M_x = \omega / (c - \tau) \) with any \( 0 < \tau < c \) then the set \( \mathcal{G} = (-\infty, -M_x] \cup [M_x, +\infty) \) is a small set for the transition kernel of \( X_n \).

**Proof:**

We have to show that for all \( x \in \mathcal{G} \) and Borel set \( A \) in \( \mathbb{R} \):

\[ P(A, x) \geq \epsilon Q(A), \]

where \( \epsilon > 0 \) and \( Q \) is a probability distribution. The main problem here comes from the compact support of \( P(t, x) \). Take \( x \) such that \( |x| \geq M \) then:

\[ 1 + \omega - c - \frac{\omega}{M} + c \epsilon_n \leq 1 + \omega - c - \frac{\omega}{x} + c \epsilon_n \leq 1 + \omega - c + \frac{\omega}{M} + c \epsilon_n, \]

where \( \epsilon_n \) has compact support \([-1, +1]\). It is simple to see that with \( M = M_x = \omega / (c - \tau) \) the above bound becomes:

\[ 1 + \omega - 2c + \tau + c \epsilon_n \leq 1 + \omega - c - \frac{\omega}{x} + c \epsilon_n \leq 1 + \omega + \tau + c \epsilon_n. \]

The intersection of the supports of \( 1 + \omega - 2c + \tau + c \epsilon_n \) and \( 1 + \omega + \tau + c \epsilon_n \) is the set \([1 + \omega - c - \tau, 1 + \omega - c + \tau]\) whatever the value of \( x \) in \( \mathcal{G} \). The probability measure \( Q \) mentioned above may be chosen as the uniform distribution with support \([1 + \omega - c - \tau, 1 + \omega - c + \tau]\).

Now we turn to the drift condition. Our task consists in constructing a function \( g : \mathbb{R} \rightarrow [1, +\infty) \) such that for all \( x \):

\[ \int_{\mathbb{R}} g(t) P(t, x) dt \leq \rho_1 g(x) + \rho_2 1_{x \in \mathcal{G}}, \]   

(11)
where $0 < \rho_1 < 1$ and $\rho_2 \geq 0$. Besides, in order to get a CLT on $\log |X_n|$ we must further ensure that for all $x$:

$$\log^2 |x| \leq g(x).$$  \hspace{1cm} (12)

Note however that, if (11) holds for $g$ but (12) fails, then both conditions will hold for updated function $g^* = \eta g$ with constant $\eta > 1$ and $\rho_2' = \eta \rho_2$ such that (12) holds.

The next Lemma constructs the function $g$ mentioned above.

**Lemma 10.** Take for $g$ the even function defined by $g(x) = C_1 / \sqrt{|x|}$ for $|x| \leq M \tau$ and $g(x) = C_2 (\log |x|)^2$ for $|x| > M \tau$.

Assume that $B_1$ holds. Then it is always possible to choose three constants $\tau$, $C_1$ and $C_2$ such that (11) holds for a specific choice of $\rho_1$ and $\rho_2$.

The proof of Lemma 10 is postponed to the Appendix.

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Appendix

This appendix provides a mathematical derivation of Lemma 10.

A Proof of Lemma 10

The proof of the Lemma just consists in an explicit construction of the above-mentioned \( \tau \), \( C_1 \), and \( C_2 \). This construction is detailed for the sake of completeness.

At this point and in order to simplify the computations below we will assume that the distribution of \( \varepsilon_n \) is uniform on \([-1, +1]\) instead of the convolution of two \( \mathcal{N}[-1/2, 1/2] \) distributions.

Set \( \lambda = 1 + \omega - c \), assume that \( \lambda > 0 \) (the case \( \lambda < 0 \) follows the same lines) and notice that:

\[
\int_{\mathbb{R}} g(t) P(t,x) dt = \frac{1}{2c} \int_{\lambda-(\omega/x)-c}^{\lambda-(\omega/x)+c} g(s) ds = \frac{1}{2c} \int_{\lambda-(\omega/x)-c}^{\lambda-(\omega/x)-\lambda+c} g(s) ds,
\]

the last inequality stemming from parity of \( g \). We should consider two cases \( x > 0 \) and \( x < 0 \).

The proof takes 2 parts (\( x > 0 \) and \( x < 0 \) respectively). Both are given again for completeness and because the problem is not symmetric. Each part is split in three steps: the two first steps deal with \( x \notin \mathcal{E} \), the third with \( x \in \mathcal{E} = (-\infty, -M_{\tau}] \cup [M_{\tau}, +\infty) \).

Part 1: \( x > 0 \)

First step: We split \([0,M_{\tau}]\) in two subsets, \([0,M_{\tau}] = [0,A_{\tau}] \cup [A_{\tau}, M_{\tau}]\) with:

\[ A_{\tau} = \omega / (M_{\tau} + 1 + \omega) \]

is chosen such that \( 0 \leq x \leq A_{\tau} \) implies the following inequality on the lower bound of the integral: \((\omega/x) - \lambda - c > M_{\tau}\). Clearly \( A_{\tau} \leq M_{\tau} \) because \( \lambda > 0 \) so \( \tau = M_{\tau} \).

Then:

\[
\frac{1}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} g(s) ds = \frac{C_2}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} \log^2|s| ds \leq C_2 \log^2 |(\omega/x) - \lambda + c|.
\]

Let:

\[
\sup_{0 \leq x \leq \omega/(c+M_{\tau})} \sqrt{|x|} (\log |(\omega/x) - \lambda + c|)^2 = K_1(\omega,c,\tau) < +\infty.
\]

The strictly positive \( K_1(\omega,c,\tau) \) exists because \( \sqrt{\log^2 |(\omega/x) - \lambda + c|} \) is bounded on \([0,A_{\tau}]\). The first condition reads:

\[
\frac{1}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} g(s) ds \leq \rho_1 C_1 / \sqrt{|x|}, \quad 0 \leq x \leq A_{\tau}
\]

whenever

\[
C_2 K_1(\omega,c,\tau) \leq \rho_1 C_1
\]

and \( \rho_1 \) will be fixed after the second step.

Second step: Now we turn to \( A_{\tau} \leq x \leq M_{\tau} \). We still have \( g(x) = C_1 / \sqrt{|x|} \) but we need to focus on the bounds of the integral.

This time the lower bound of the integral \((\omega/x) - \lambda - c \in [-\lambda - \tau, M_{\tau}]\) and the upper bound \((\omega/x) - \lambda + c \in [2c - \lambda - \tau, 2c + M_{\tau}]\). We are going to require that \((\omega/x) - \lambda - c \geq -M_{\tau}\) it suffices to take \( \lambda + \tau \leq M_{\tau} \) and this comes down to the following set of constraint on \( \tau \): \( \{ \tau \geq c - \omega \} \cup \{ \tau \leq c - 1 \} \). We keep the second and assume once and for all that:

\[
\tau \leq c - 1.
\]

Then for \( x \in [A_{\tau}, M_{\tau}]\),

\[
\frac{1}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} g(s) ds = \frac{1}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} g(s) ds + \frac{1}{2c} \int_{(\omega/x)-\lambda-c}^{(\omega/x)-\lambda+c} g(s) ds
\]

\[
\equiv \mathcal{A}_1 + \mathcal{A}_2.
\]

We want to make sure that the upper bound \((\omega/x) - \lambda + c \) is larger than \( M_{\tau} \). This will hold if \((\omega/M_{\tau}) - \lambda + c \geq M_{\tau}\) hence if \( 2c - \tau - \lambda \geq M_{\tau}\). We imposed previously that \( \lambda + \tau \leq M_{\tau}\). So:

\[
\tau < c - \frac{\omega}{c} \Rightarrow M_{\tau} < c \Rightarrow 2c - \tau - \lambda \geq M_{\tau}
\]
but the constraint $\tau < c - \omega/c$ is weaker than (14) consequently (15) holds.

Focus on the first term $\mathcal{I}_1$ in (15) and consider:

$$\mathcal{I}_1 = \frac{1}{2c} \int_{(\omega/x) - \lambda - c}^{M} g(s) \, ds = \frac{1}{2c} \int_{(\omega/x) - \lambda - c}^{M} C \sqrt{x} \, ds.$$  

Consider the (only) two situations on the sign of $(\omega/x) - \lambda - c = (\omega/x) - (1 + \omega)$. If $x < \omega/(1 + \omega)$ then $(\omega/x) - (1 + \omega) > 0$ and $\mathcal{I}_1 \leq \frac{C_1}{c} \sqrt{M}$. Notice by the way and for further purpose that:

$$\sup_{x \in [\tau, \infty)} \sqrt{|x|} \mathcal{I}_1 \leq \frac{C_1}{c} \sqrt{\omega/(1 + \omega)} \cdot M.$$  

If $x \geq \omega/(1 + \omega)$ then $(\omega/x) - (1 + \omega) \leq 0$ and:

$$\mathcal{I}_1 = \frac{1}{2c} \int_{(\omega/x) - \lambda - c}^{M} g(s) \, ds = \frac{1}{2c} \int_{(\omega/x) - (1 + \omega)}^{0} \frac{C_1}{\sqrt{|s|}} \, ds + \frac{1}{2c} \int_{0}^{M} \frac{C_1}{\sqrt{|s|}} \, ds$$

Again:

$$\sqrt{|x|} \mathcal{I}_1 \leq \frac{C_1}{c} \left[ \sqrt{x(1 + \omega) - \omega} + \sqrt{xM} \right].$$

From the bounds above we see that:

$$\sup_{x \in [\tau, M]} \sqrt{|x|} \mathcal{I}_1 \leq \frac{C_1}{c} \left[ \sqrt{M(1 + \omega)} - \omega + M \right].$$

The reader will soon understand why we need to make sure that the right hand side in equation above is strictly under $C_1$. It is not hard to see that the function $\tau \mapsto \sqrt{|M(1 + \omega) - \omega| + M}$ is increasing and continuous on $[0, c - 1]$. If we prove that for some $\delta \in ]0, 1[:$

$$\frac{1}{c} \left[ \sqrt{M(1 + \omega) - \omega} + M \right] = 1 - 3\delta < 1,$$

then the existence of some $\tau^+ > 0$ such that:

$$\frac{1}{c} \left[ \sqrt{M(1 + \omega) - \omega} + M \right] = 1 - 2\delta < 1$$  

will be granted. But $\frac{1}{c} \left[ \sqrt{M(1 + \omega) - \omega} + M \right] = \frac{1}{c} \left[ \sqrt{\omega \lambda + \omega^2} \right]$.

If we assume that $\lambda < \omega/c < (1 + c)/4$ (assumption $B_1$) then since $c > 1$:

$$\frac{1}{c} \left[ \sqrt{\omega/c + \omega/c} \right] < \frac{1}{2} \left( 1 + \frac{1}{c} \right) < 1.$$

We turn to $\mathcal{I}_2$ in (15):

$$\mathcal{I}_2 = \frac{1}{2c} \int_{M}^{(\omega/x) - \lambda + c} g(s) \, ds = \frac{C_2}{2c} \int_{M}^{(\omega/x) - \lambda + c} (\log s)^2 \, ds \leq \frac{C_2}{\sqrt{|x|}} K_2 (\omega, \tau^+),$$

where:

$$K_2 (\omega, \tau^+) = \sup_{x \in [\tau, M]} \frac{\sqrt{|x|}}{2c} \int_{M}^{(\omega/x) - \lambda + c} (\log s)^2 \, ds.$$  

Set finally $\rho_1^+ = 1 - \delta < 1$.

From (15) we get:

$$\frac{1}{2c} \int_{(\omega/x) - \lambda - c}^{M} g(s) \, ds \leq \frac{C_1}{\sqrt{|x|}} (1 - 2\delta) + \frac{C_2}{\sqrt{|x|}} K_2 (\omega, \tau^+)$$

$$\leq \rho_1^+ \frac{C_1}{\sqrt{|x|}},$$

whenever holds the new condition:

$$C_2 K_2 (\omega, \tau) \leq C_1 \delta.$$  

(17)
Finally comparing (13) and (17), we see that both conditions cannot be incompatible. Accurate choices of the couple \((C_1^+, C_2^+)\) are given by the summary bound:

\[
C_2^+ \leq C_1^+ \min \left( \frac{\delta}{K_2}, \frac{1 - \delta}{K_1} \right).
\]  

(18)

It is now basic to see that the quadruple \((C_1^+, C_2^+, \tau^+, \rho_1^+)\) yields the drift condition (11) for \(x \notin \mathscr{C}'\).

**Third step:** The remaining step is to check the inequality for some \(\rho_2\):

\[
\int_\mathbb{R} g(t) P(t,x) dt \leq \rho_1^+ g(x) + \rho_2,
\]

for any \(x \in \mathscr{C}'\) - that is any \(|x| > M_\tau\) (rather \(x > M_\tau\) here as explained above since \(x > 0\)). We see that:

\[
0 \leq \frac{\omega}{x} \leq \frac{\omega}{M_\tau},
\]

and:

\[
\frac{1}{2c} \int_{(\omega/x)-\delta}^{(\omega/x)-\lambda+c} g(s) ds \leq \frac{1}{2c} \int_{-\delta}^{2c-\delta-c} g(s) ds \leq \frac{1}{2c} \int_{-c}^{c} g(s) ds.
\]

The values of the constants \(C_1\) and \(C_2\) were fixed above. Then denote:

\[
\rho_2^+ = \frac{1}{2c} \int_{-c}^{c} g(s) ds > 0,
\]

then clearly for any \(x \in \mathscr{C}'\):

\[
\int_\mathbb{R} g(t) P(t,x) dt \leq \rho_2^+,
\]

so that (11) holds.

**Part 2 \((x \leq 0)\)**

We go on with \(x < 0\) and \(\lambda > 0\), set \(y = -x \geq 0\),

\[
\int_\mathbb{R} g(t) P(t,x) dt = \frac{1}{2c} \int_{-\lambda-(\omega/x)}^{\lambda-(\omega/x)+c} g(s) ds = \frac{1}{2c} \int_{-\lambda}^{\lambda} g(s) ds.
\]

Since \(g\) is even and in view of the proposed \(\mathscr{C}'\) we just have to prove exactly the following drift condition with \(x > 0\):

\[
\frac{1}{2c} \int_{(\omega/x)+\lambda}^{(\omega/x)+\lambda+c} g(s) ds \leq \rho_1 g(x) + \rho_2 \mathbf{1}_{x \in \mathscr{C}'.}
\]

**First step:** Take \(x \notin \mathscr{C}'\). We split \([0,M_\tau]\) in two subsets, \([0,M_\tau] = [0,B_\tau] \cup [B_\tau,M_\tau]\) with \(B_\tau = \omega/(M_\tau - \lambda + c)\) is chosen such that \(0 \leq x \leq B_\tau\) implies the following inequality on the lower bound of the integral: \((\omega/x) + \lambda - c > M_\tau\). Clearly \(B_\tau \leq M_\tau\) for all \(\tau\). Then:

\[
\frac{1}{2c} \int_{(\omega/x)+\lambda}^{(\omega/x)+\lambda+c} g(s) ds = \frac{C_2}{2c} \int_{(\omega/x)+\lambda}^{(\omega/x)+\lambda+c} \log^2 |s| ds \leq C_2 \log^2 |(\omega/x) + \lambda + c|.
\]

Let:

\[
\sup_{0 \leq x \leq B_\tau} \sqrt{|x|} \log^2 |(\omega/x) + \lambda + c| = K_1(\omega, c, \tau) < +\infty.
\]

The strictly positive \(K_1(\omega, c, \tau)\) exists because \(\sqrt{|x|} \log^2 |(\omega/x) + \lambda + c|\) is bounded on \([0,B_\tau]\). The initial condition reads:

\[
\frac{1}{2c} \int_{(\omega/x)+\lambda}^{(\omega/x)+\lambda+c} g(s) ds \leq \rho_1 C_1 / \sqrt{|x|}, \quad 0 \leq x \leq B_\tau
\]

whenever \(C_2 K_1(\omega, c, \tau) \leq \rho_1 C_1\) and \(\rho_1\) will be fixed later.

**Second step:** Now we turn to \(B_\tau \leq x \leq M_\tau\). We still have \(g(x) = C_1 / \sqrt{|x|}\) but we need to focus on the bounds of the integral.
This time the lower bound of the integral \((\omega/x) + \lambda - c \in [\lambda - \tau, M_\tau]\) and the upper bound \((\omega/x) + \lambda + c \in [\lambda + 2c - \tau, M_\tau + 2c]\). If we assume that \(\tau \leq \lambda\) then \((\omega/x) + \lambda - c \geq 0 \geq -M_\tau\). Besides in order that \(\lambda + 2c - \tau \geq M_\tau\) we just need that \(2c \geq M_\tau\) or \(\tau \leq c - \omega/(2c)\). As a consequence the assumption:

\[
\tau \leq \min \left(\lambda, c - \frac{\omega}{2c}\right)
\]  

allows to write for \(x \in [B_\tau, M_\tau]\),

\[
\frac{1}{2c} \int_{(\omega/x) + \lambda - c}^{(\omega/x) + \lambda + c} g(s) \, ds = \frac{1}{2c} \int_{(\omega/x) + \lambda - c}^{M_\tau} g(s) \, ds + \frac{1}{2c} \int_{M_\tau}^{(\omega/x) + \lambda + c} g(s) \, ds
\]

\[
\equiv \mathcal{J}_1 + \mathcal{J}_2,
\]

with non-null \(\mathcal{J}_1\) and \(\mathcal{J}_2\). Focus on:

\[
\mathcal{J}_1 = \frac{C_1}{2c} \int_{(\omega/x) + \lambda - c}^{M_\tau} \frac{1}{\sqrt{s}} \, ds = \frac{C_1}{c} \left[ \sqrt{M_\tau} - \sqrt{(\omega/x) + \lambda + c} \right]
\]

\[
\sup_{x \in [B_\tau, M_\tau]} \sqrt{x} \mathcal{J}_1 \leq \frac{C_1}{c} M_\tau.
\]

At last we see that for \(M_\tau \leq 1\) i.e. \(\tau \leq c - \omega\), \(\sup_{x \in [B_\tau, M_\tau]} \sqrt{x} \mathcal{J}_1 \leq C_1/c\). This condition combined with (19) let us set in the sequel:

\[
\tau \leq \tau^- = \min(\lambda, c - \omega).
\]

We turn to \(\mathcal{J}_2\):

\[
\frac{1}{2c} \int_{M_\tau}^{(\omega/x) + \lambda + c} g(s) \, ds = \frac{C_2}{2c} \int_{M_\tau}^{(\omega/x) + \lambda + c} \log^2 |s| \, ds \leq \frac{C_2}{\sqrt{x}} K^2_2(\omega, c, \tau^-),
\]

with:

\[
K^2_2 = K_2(\omega, c, \tau^-) = \sup_{x \in [B_\tau, M_\tau]} \sqrt{x} \int_{M_\tau}^{(\omega/x) + \lambda + c} (\log s)^2 \, ds.
\]

Set finally \(\rho^-_1 = (1 + 1/c)/2 < 1\). From all that was done above we get:

\[
\int_{\mathbb{R}} g(t) P(t, x) \, dt \leq \frac{C_1}{1} \frac{1}{|x|} + \frac{C_2}{\sqrt{|x|}} K_2 \leq \rho^-_1 \frac{C_1}{\sqrt{|x|}}
\]

whenever:

\[
C_2 K_2(\omega, c, \tau^-) \leq \frac{1 - 1/c}{2} C_1.
\]

This will be combined with the constraint of the first step \(C_2 K_1^c \leq \rho_1 C_1\) (we denoted \(K_1(\omega, c, \tau^-) = K_1^c\)). The new condition:

\[
C^2_2 \leq C^-_1 \min \left( \frac{\rho^-_1}{K_1^c}, \frac{1 - 1/c}{2K_2} \right)
\]

ensures that

\[
\frac{1}{2c} \int_{(\omega/x) + \lambda - c}^{(\omega/x) + \lambda + c} g(s) \, ds \leq \rho^-_1 g(x) \text{ for } x \in [0, M_\tau].
\]

**Third step:** The remaining step is to check the inequality:

\[
\int_{\mathbb{R}} g(t) P(t, x) \, dt \leq \rho^-_1 g(x) + \rho_2,
\]

for any \(x \in \mathcal{G}\) -that is here any \(x > M_\tau\). Adapting the method given above is straightforward and leads to the desired result with a given \(\rho^-_1\).

We are ready to conclude. Take

\[
C^*_2 = C^-_1 \min \left( \frac{\delta}{K^*_2}, \frac{1 - \delta}{K^*_1}, \frac{1 - 1/c}{2K^*_2} \right).
\]

Conditions (18) and (20) hold for the couple \((C^*_1, C^*_2)\). For such a couple we have:

\[
\int_{\mathbb{R}} g(t) P(t, x) \, dt \leq \rho^+_1 g(x) + \rho^+_2, \quad x > 0,
\]

\[
\int_{\mathbb{R}} g(t) P(t, x) \, dt \leq \rho^-_1 g(x) + \rho^-_2, \quad x \leq 0.
\]
and for all $x$:

$$\int_R g(t) P(t,x) \, dt \leq \max (\rho_1^+, \rho_1^-) g(x) + \max (\rho_2^+, \rho_2^-).$$

This finishes the proof of the Lemma.