Generalized dimensions, large deviations and the distribution of rare events

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**Abstract**

Generalized dimensions of multifractal measures are usually seen as static objects, related to the scaling properties of suitable partition functions, or moments of measures of cells. When these measures are invariant for the flow of a chaotic dynamical system, generalized dimensions take on a dynamical meaning, as they provide the rate function for the large deviations of the first hitting time, which is the (average) time required to connect any two different regions in phase space. We prove this result rigorously under a set of stringent assumptions. As a consequence, the statistics of hitting times provides new algorithms for the computation of the spectrum of generalized dimensions. Numerical examples, presented along with the theory, suggest that the validity of this technique reaches far beyond the range covered by the theorem.

We state our result within the framework of extreme value theory. This approach reveals that hitting times are also linked to dynamical indicators such as stability of the motion and local dimensions of the invariant measure. This suggests that one can use local dynamical indicators from finite time series to gather information on the multifractal spectrum of generalized dimension. We show an application of this technique to experimental data from climate dynamics.

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1. Introduction and summary of the paper

Generalized dimensions are a primary tool for the analysis of multifractal measures, i.e., measures whose local densities feature a range of different scaling exponents. Interest in these quantities originated in the eighties of the last century [1–3], primarily for the study of chaotic attractors and fully developed turbulence [4,5] and rapidly became important also from the mathematical viewpoint. The combined effort of physicists and mathematicians leads to the development of the so-called thermodynamical formalism [6–8] in which generalized dimensions play a major role. For the sake of numerical experiments, but also of application to empirical data, many different techniques have been proposed along the years for the numerical calculation of generalized dimensions (we just quote [2,9–14] because even a partial list of references would require a full paper).

In this work we link generalized dimensions and the recurrence properties of the dynamics. In the same line of thought of Kac’s theorem [15], and more generally of ergodic theory, we study the connection between a dynamical quantity, the hitting time of a small set, and a static quantity, the statistical distribution of the measure of small balls.

This approach has been initiated in [1,16,17]: generalized dimensions can be derived from the moments of the so-called first return time, the length of time required for the dynamics to return close to a chosen initial point on the attractor. Numerical experiments in [18] showed applicability and failures of this technique: in [19] the relation between dimensions generated via return times and the original quantities has been examined rigorously and in full generality, also providing explicit examples and counter-examples.

A related result will be derived herein, by considering hitting times [20], rather than return times: hitting times are related to targeting small regions of phase space, starting from a different, arbitrary point. Proposition 1 in Section 3 shows that dimensions indexed by $q \leq 2$ (the meaning of this index, which is related to the order of moments, will be explained in the next section) can be computed in this way. In Section 4.1, the abstract theory...
is translated into a numerical algorithm and applied to the test cases of the Arnol'd cat map and the Hénon attractor.

A second fundamental point of this paper is to show that generalized dimensions yield the rate function for the large deviations of the first hitting time of a ball of given radius. They quantify the rate at which the probability of observing “non typical” values diminishes when the radius goes to zero. Our result, Proposition 2 in Section 3, parallels a previous investigation [20] where the target set in phase space is a dynamical cylinder, rather than a ball. Rigorous theory is presented for the case of conformal repellers but we believe that similar results hold in more general settings. A numerical illustration of the large deviation statistics is presented in Section 4.2, in the case of a dynamical system evolving on the Sierpinski gasket.

To obviate the limitation of the hitting technique to a part of the spectrum of generalized dimensions, in Section 5 we continue the investigation of a preceding paper by some of the present authors [23], in which the correlation dimension (the dimension of index \( q = 2 \)) was computed by means of extreme value theory (EVT). The dynamical extremal index (DEI) appearing in the Gumbel’s limiting law followed as a by-product and it was interpreted as the rate of backward contraction on the unstable subspaces, a quantity closely related to positive Lyapunov exponents. We now extend this technique to the case of arbitrary, positive, integer values \( q \neq 1 \) of the associated DEI is again related to Lyapunov exponents, but in Section 6, via Proposition 3 and the subsequent analysis, we show that it is also affected by the variation of the invariant density and by the lack of uniform hyperbolicity of the system. We consider the rate function of the hitting time as a way to detect and quantify the presence of rare events in the dynamics. These events are produced by the presence of points where the local dimensions and the first hitting time (our statistical indicators) do not assume their typical values. Although these non-typical points (like e.g. unstable fixed points, periodic points) have null probability to be attained by the dynamics, their influence in a finite region around their location affects the convergence of statistical indicators via an exponentially small probability of deviations from the typical values. The rate function measures the intensity of such deviations.

To show a realization of this scenario, in Section 7 we analyze experimental data coming from climate dynamics. In fact, our broader goal is to implement statistical tools to investigate and to interpret data coming from various physical situations, like climate dynamics, but also turbulence, neuroscience and biology. Further applications of the methods developed in this paper will appear in forthcoming publications.

2. Definitions and statement of the principal results

2.1. Review of related literature on multifractals

Let \((M, \mu, T)\) be a dynamical system given by a map \( T \) acting on a metric space \( M \) with distance \( d(\cdot, \cdot) \) and preserving a Borel probability measure \( \mu \). If we denote by \( B(z, r) \) the ball of radius \( r \) centered at \( z \in M \), we define the spectrum \( D_q \) \((q \neq 1)\) of the generalized dimensions of \( \mu \) by the scaling relation of the \( q \)-correlation integral with respect to the radius \( r \)

\[
\Gamma_q(r) = \int_M \mu(B(z, r))^q d\mu \sim r^{\tau(q)} \mu(B(z, r)) = r^{\tau(q)},
\]

For \( q = 1 \) the above is replaced by:

\[
\int_M \log(\mu(B(z, r))) d\mu \sim r \mu(B(z, r)) = r \mu(B(z, r)).
\]

See [1] and [16] for the introduction of the theory, [8] for a formal rigorous definition of the previous scaling behavior and the books [9,21,24] for several applications to non-linear systems.

The scaling equation (1) can be made mathematically precise by the following limit that defines the real function \( \tau \) of the index \( q \) (employing liminf and limsup to define upper and lower quantities, when they differ):

\[
\tau(q) = \lim_{r \to 0} \frac{\log \Gamma_q(r, q)}{\log r}.
\]

Generalized dimensions \( D_q \) are obtained from the function \( \tau(q) \) via the equation

\[
\tau(q) = D_q(q - 1)
\]

when \( q \neq 1 \), and by the l’Hôpital rule when \( q = 1 \). It is well known that, for a large class of dynamical systems, the Legendre transform of \( \tau(q) \), namely

\[
f(\alpha) = \min(\alpha q - \tau(q)),
\]

is the Hausdorff dimension of the set of points \( z \in M \) verifying:

\[
\lim_{r \to 0} \frac{\log \mu(B(z, r))}{\log r} = \alpha,
\]

provided the limit exists, see [8,21] and references therein. This limit is called the local dimension of the measure \( \mu \) at the point \( z \).

The so-called exact dimensional measures \( \mu \) have a local dimension that is constant \( \mu - \) a.e., This dimension coincides with the information dimension \( D_1 \) [25]. Several dynamical systems with hyperbolic properties possess an invariant measure that is exact dimensional, whose information dimension can be expressed in terms of the Lyapunov exponents and of the metric entropy. When the function \( f(\alpha) \) is not singular, the space \( M \) can be parted into uncountably many subsets characterized by the same local dimension, which are of zero measure (except of course that of dimension \( D_1 \)) but of positive Hausdorff dimension, yielding what is called a multifractal.

When the generalized dimensions \( D_q \) vary with \( q \), they imply deviations of the local dimensions defined in Eq. (6) from its expected value \( D_1 \). To gauge the deviation of the observable

\[
\frac{\log \mu(B(z, r))}{\log r}
\]

from its expected value, at the finite resolution \( r > 0 \), one considers the quantity \( \mu \left( \left\{ z \in M \text{ s.t. } \frac{\log \mu(B(z, r))}{\log r} \in I \right\} \right) \), where \( I \) is any interval in \( R \), including or not the expected value \( D_1 \). It has been recently proven [26] the interesting result that in the family of systems known as conformal repellers the previous deviations decrease exponentially when \( r \) tends to zero, with a rate that is given in terms of the generalized dimensions:

\[
\lim_{r \to 0} \frac{1}{\log r} \log \mu \left( \left\{ z \in M \text{ s.t. } \frac{\log \mu(B(z, r))}{\log r} \in I \right\} \right) = \inf_{x \in I} Q(s).
\]

The rate function \( Q(s) \) is determined again by \( \tau(q) \) of Eq. (3):

\[
Q(s) = \sup_{q \in R} \{-qs + \tau(q + 1)\}.
\]

In addition to the point \( z \), let us now consider a second point \( x \in M \), and let us denote by \( H_{B(x, r)}(X) \) the first hitting time of the
point \( x \) in the ball \( B(z, r) \):
\[
\mathcal{H}_{B(z, r)}(x) = \min \{n > 0 \text{ s.t. } T^n(x) \in B(z, r) \}. \tag{9}
\]
A particular situation happens when the point \( x \) belongs to the ball \( B(z, r) \). In this case one calls \( \mathcal{H}_{B(z, r)}(x) \) the time of first return of \( x \) into \( B(z, r) \). It is convenient to define \( \mu_{B(z, r)}(\cdot) \), the restriction of \( \mu \) to \( B(z, r) \):
\[
\mu_{B(z, r)}(A) = \frac{\mu(A)}{\mu(B(z, r))},
\]
where \( A \) is any measurable set. When the invariant measure \( \mu \) is ergodic, it is well known that the first return time satisfies Kac’s theorem [15]:
\[
\mathbb{E}_{\mu_{B(z, r)}}(\mathcal{H}_{B(z, r)}) = \int_{B(z, r)} \mathcal{H}_{B(z, r)}(x) d\mu_{B(z, r)}(x) = \frac{1}{\mu(B(z, r))}. \tag{10}
\]
By using the previous result in association with the Ornstein and Weiss theorem [27], it is possible to show that the first return time satisfies
\[
\lim_{r \to 0} \frac{\log \mathcal{H}_{B(z, r)}(x)}{-\log r} = D_1, \tag{11}
\]
for \( z \) chosen \( \mu \)-a.e., (see [28,29]). Observe that in the above equation we are considering the first return of the center of the ball into the ball itself, i.e. \( x = z \). The first return time enjoys exponential large deviations, namely it was proven in [26] that:
\[
\mu \left( \left\{ z \in M \text{ s.t. } \frac{\log \mathcal{H}_{B(z, r)}(z)}{-\log r} \in I \right\} \right) \sim e^{\inf_{x \in I} Q^a(x)}. \tag{12}
\]
The rate function \( Q^a \) is slightly different from the \( Q \) given above. For its precise definition we refer again to [26] Theorem 2.5, where the large deviation property is proven.\(^2\)

The question has been asked whether a multifractal description of the first return time could be meaningful, by considering the set of points where the limit in Eq. (11) is different from the typical value \( D_1 \). Yet, in the case of conformal repellers it has been proven [30,31], that all level sets with a value different from \( D_1 \) have the same Hausdorff dimension of the ambient space, see also [32]. This is a further point in favor of using large deviations instead of the multifractal description for studying recurrence quantities.

### 2.2. Mathematical achievements of this paper

We are now ready to summarize our mathematical relevant results and the numerical algorithms that can be deduced from them and that will be extensively used in the next parts of the paper. We defer to these parts for precise statements and technical details.

Let us first return in full generality to hitting times, when the initial condition \( x \) does not necessarily belong to a neighborhood of the final state \( z \). For systems with super-polynomial decay of correlations a result analogous to Eq. (11) holds:
\[
\lim_{r \to 0} \frac{\log \mathcal{H}_{B(z, r)}(x)}{-\log r} = D_1, \tag{13}
\]
for \( x \) and \( z \) chosen \( \mu \)-a.e. [33].

The first question that we will discuss in this paper is whether hitting times enjoy exponentially large deviations, that is, whether and under what conditions it holds true that
\[
\mu \times \mu \left( \left\{ (x, z) \in M \times M \text{ s.t. } \frac{\log \mathcal{H}_{B(z, r)}(x)}{-\log r} \in I \right\} \right) \sim e^{\inf_{x \in I} Q^a(x)},
\]
where the rate function \( \hat{Q} \) presumably involves again the generalized dimensions. Notice that we weighted the event with the product measure since there are two sources of alea, in the choice of the starting point \( x \) and of the target point \( z \). Such large deviation property has interesting physical consequences. As anticipated above, we expect that the presence of points \( x \) and \( z \) giving limits different from \( D_1 \) in Eq. (13) – which we interpret respectively as exceptional initial conditions (\( x \)) and rare target regions (\( z \)) – yields deviations in the limit of Eq. (13) for small \( r \). These deviations go to zero exponentially fast, in a well-defined limit procedure, with a rate which is measurable and that can be linked to the intensity of the extreme events.

The large deviations bound (14) is proven in Section 3. It can be stated as follows:

- For a large class of systems including conformal repellers we have
  \[
  \lim_{r \to 0} \frac{1}{\log r} \log \int_M \int_M \mathcal{H}_{B(z, r)}(x) e^{-q} d\mu(x) d\mu(z) = \left\{ \begin{array}{ll}
  \tau(2 - q), & q > 0 \\
  \tau(2), & q \leq 0
  \end{array} \right.
  \tag{14}
  \]
  where \( \hat{Q} = \sup_q \{ \log \mathcal{H}_{B(z, r)}(x) - \log r \} \) and \( \hat{Q}(z) = \hat{Q}(D_1 + s) \).

Our second mathematical result establishes an extreme value distribution for a specific observable, chosen appropriately so that the parameters entering such a distribution are related to generalized dimensions and provide a new dynamical indicator via the evaluation of the associated extremal index. To introduce this observable, we let \( T \) be a map acting on the space \( M \), preserving the probability measure \( \mu \). Also consider the direct product of \( q \) systems of the same type and call \( T_q \) the direct product map acting on \( M^q \) and preserving the direct product measure \( \mu_q \). Define the following observable on \( M^q \):
\[
\phi(x_1, x_2, \ldots, x_q) = -\log \left( \max_{i=1, \ldots, q} d(x_i, x_1) \right), \tag{15}
\]
where each \( x_i \in M \); we also write \( X_q = (x_1, x_2, \ldots, x_q) \) and \( T_q(X_q) = (T(x_1), \ldots, T(x_q)) \). We consider the observable \( M_q(\phi; X_q) = \max(\phi(x_i), \ldots, \phi(T_q(X_q))) \), and its distribution function \( F_{\mu_q}(u) = \mu_q(\{X_q \in M^q \text{ s.t. } M_q(\phi; X_q) \leq u \}) \). We can prove the following result:

- For a large class of systems including conformal repellers, suppose that \( t > 0 \) is a real positive number and \( u_q \) a sequence verifying \( \mu_q(\phi > u_q) \to \frac{1}{t^q} \), then there exists a number \( \theta_q \in [0, 1] \) such that:
  \[
  \lim_{\theta_q \to \infty} F_{\mu_q}(u_q) = e^{-\theta_q t}.
  \]
  The quantity \( \theta_q \) is called the Dynamical Extremal Index.

A sketch of the proof of this result is provided in the Appendix and it follows analogous results by two of the authors.

The relationship between the preceding proposition and generalized dimensions is obtained by observing that
\[
\mu_q(\phi > u_q) = \int_M d\mu(x_1) \mu(B(x_1, e^{-u_q} y^{-1})),
\]
which allows to write the asymptotic relation, for large \( n \),
\[
\frac{u_n}{D_q(q - 1)} \sim -\frac{\log t}{D_n(q - 1)} + \frac{\log n}{D_n(q - 1)} = -\frac{\log t}{a_n} + b_n. \tag{16}
\]
The real quantities \(a_n\) and \(b_n\) can be obtained by a maximum likelihood estimation as we will show in Section 5. Our third mathematical result concerns the computation of the dynamical extremal index. In the Appendix we give a general formula (60) of geometric type. It is applied to a special case in Section 6, to obtain the result:

- For a large class of maps including conformal repellers preserving a measure absolutely continuous with respect to Lebesgue and with density \(h\) we have

\[
\theta_q = 1 - \frac{\int h(x)^q dx}{\int h(x)^q dx}.
\]

2.3. Numerical algorithms

The numerical techniques for the evaluation of the various quantities considered in this paper are described where needed: in particular in Sections 4 and 5 we tackle the computation of generalized dimensions.

A different case is that of Sections 5.2 and 6, where the method employed to compute both the generalized dimensions and the spectrum of DEI is substantially the same used in [23] for the determination of the correlation dimension and of \(\theta_2\). It can be summarized here in the following steps:

1. Initialization of \(q\) subsystems with different starting points chosen at random with respect to \(\mu\). This can be done by choosing a point at random in the basin of attraction of the map and iterating a large number of times the dynamics.

2. Computation of a trajectory of length \(N\) by simultaneous iteration of the \(q\) subsystem, followed at each iteration by the computation of the observable \(\Phi\). This produces a time series of length \(N\).

3. Estimation from the obtained time series of \(\theta_q\), using (for example) the S"u"veges estimate [34]. This estimate requires the choice of a high threshold that can be taken as a high quantile of the time series.

4. Division of time into blocks of length \(n\), selection of the maximum in each of these blocks and computation of the empirical distribution of this maximum.

5. Maximum likelihood estimation of the scale parameter \(\sigma\) of the empirical maximum distribution (for \(n\) large enough, the real value \(a_n\) in Eq. (16) is close to the inverse of this parameter). This can be done using the Matlab function \texttt{gevfit}.

6. Computation of \(D_q = 1/(q-1)\sigma\).

7. Averaging the obtained values of \(D_q\) obtained over \(l\) trajectories.

3. Large deviations for the first hitting times

Extreme value theory (EVT) can be used to determine the probability that the system enters for the first time a small region of the phase space (rare event) after a certain amount of time [36]. Instead of looking directly to the probability of the first occurrence of such rare events, one could ask whether the presence of those events influence the convergence of the indicators towards their expected values. This can be achieved by looking at the deviations from typical values and the rate of such deviations can be obtained using EVT.

We now state and prove a general result on large deviations in the statistics of the first hitting time. The result relies on a set of assumptions that hold true for several dynamical systems possessing some sort of hyperbolicity and exponential decay of correlations. We therefore consider dynamical systems \((M, \mu, T)\) that verify the following assumptions:

- **A-1:** Exponential distribution of hitting times with error.
  There is a constant \(C > 0\) such that for \(\mu\)-a.e. \(z \in M\) and \(t > 0\) we have
  \[
  \mu\left(\left\{ x \in M \text{ s.t. } H_{B(z,r)}(x) \geq \frac{t}{\mu(B(z,r))}\right\}\right) - e^{-t} \leq C \delta(t, 1) e^{-t}.
  \]
  (17)
  where \(\delta_t = O(\mu(B(z, r))|\log \mu(B(z, r))|)\).

In particular, for \(t > \mu(B(z, r))|\log \mu(B(z, r))|^C\) we have:

\[
\mu\left(\left\{ x \in M \text{ s.t. } H_{B(z,r)}(x) \geq \frac{t}{\mu(B(z,r))}\right\}\right) = \exp[-t(1+\delta_t)](1+\eta_t).
\]

(19)

with \(\eta_t = O(\mu(B(z, r))\mu(B(z, r)))\), while for \(t \leq \mu(B(z, r))\log \mu(B(z, r))\) we have:

\[
\mu\left(\left\{ x \in M \text{ s.t. } H_{B(z,r)}(x) \geq \frac{t}{\mu(B(z,r))}\right\}\right) \geq 1 - \frac{t}{C}.
\]

(20)

Notice that the above implies that both \(\eta_t\) and \(\delta_t\) depend on the point \(z \in M\).

- **A-2:** Exact dimensionality. The measure \(\mu\) verifies Eq. (6) and the limit value is \(\alpha = D_\gamma\) for \(\mu\)-a.e. \(z\).

- **A-3:** Uniform bound for the local measure. There exists \(d^* > 0\) such that for all \(z \in M\) we have

\[
\mu(B(z, r)) \leq r^{d^*}.
\]

(21)

- **A-4:** Existence and analyticity of the correlation integrals.
  For all \(q \in \mathbb{R}\) the limit defining \(r(q)\), Eq. (3), exists. Moreover the function \(r(q)\) is real analytic for all \(q \in \mathbb{R}\), \(r(0) = -D\), \(r(1) = 0\), \(r(q) \geq 0\) and \(r(q) \leq 0\). In particular \(r'(q) < 0\) if and only if \(\mu\) is not a measure of maximal entropy.

We derive the first assumption from Keller’s paper [37], where condition (17) is proven for the so-called REPPO maps, which include a large class of mixing systems with exponential decay of correlations. Keller’s derivation of (17) contains fine estimations of the quantity \(\mu\left(H_{B(z,r)}(x) \geq \frac{t}{\mu(B(z,r))}\right)\) for \(t \sim \mu(B(z, r))\log \mu(B(z,r))\), which we adopted in formulas (19) and (20).

Moreover, Keller’s conditions are even more general, since they hold for any point \(z\), provided the rescaled time \(\frac{t}{\mu(B(z,r))}\) is modified as \(\frac{t}{\nu(B(z,r))}\), where the factor \(\nu\) depends on the target point \(z\) and converges to the extremal index \(\alpha\) when \(r\) goes to zero.

For the sake of “nice” expanding systems we are considering, including the REPPO ones, this extremal index is equal to one almost everywhere, see also [36] for an extensive discussion, and this explains our choice in Eq. (20). Assumption A-4 is a strong one and it has been proven to hold for conformal mixing repellors endowed with Gibbs measures in [8]. This condition has also been assumed in [26], to prove the large deviation result (12). For the same class of conformal repellors Assumption A-3 holds too, see Lemma 3.15 in [26]. As remarked in the Introduction, we consider these ideal systems interesting models to establish rigorous results that might also hold in more general settings.

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3 Details on this function can be found at [35] and in [23].

4 In the proof of Proposition 1 we set the constant exponent \(C\) to unity because its value is irrelevant for the proof.

5 REPPO stands for Rare events Perron–Frobenius operators, since the conditions are given in terms of the spectral properties of the transfer operator.
Proposition 1. Let us suppose that the dynamical system $(M, \mu, T)$ verifies Assumptions (A-1)–(A-4). Then:

- For $q > 0$,
  \begin{equation}
  \lim_{r \to 0} \frac{1}{\log r} \log \int_{M} \int_{M} \mathcal{H}(x,Y(x))^{q-1} d\mu(x)d\mu(y) \\
  = \lim_{r \to 0} \frac{1}{\log r} \log \int_{M} \mu(B(z,r))^{q-1} d\mu(z).
  \end{equation}

- For $q \leq 0$,
  \begin{equation}
  \lim_{r \to 0} \frac{1}{\log r} \log \int_{M} \int_{M} \mathcal{H}(x,Y(x))^{q-1} d\mu(x)d\mu(y) \\
  = \lim_{r \to 0} \frac{1}{\log r} \log \int_{M} \mu(B(z,r))^{q-1} d\mu(z).
  \end{equation}

Proof. We follow the scheme of the proof of Theorem 3.1 in [20], by translating its argument from cylinders to balls. We use a simple lemma whose proof is a standard exercise:

Lemma 1. Consider a function $f$ from $M$ to the integer numbers larger than or equal to one: $f : M \to \mathbb{N}_+$. Let $0 < A \leq 1$ and define

\[ I(q) = \int_{M} f^q(x)d\mu(x). \]

Then, when $q > 0$,

\[ I(q) = 1 - \lim_{r \to \infty} \frac{q}{A^q} \int_{A}^{\infty} t^q f(x) \left( \frac{x}{A} \right) dt \]  

\[ < 1 \text{ condition A-1 simply becomes } |\mu\left( \left\{ x \in M \text{ s.t. } H(x) < \frac{t}{A} \right\} \right)| - e^{-t} \leq O(\delta_t). \]  

On the other hand, when $q < 0$,

\[ I(q) = -\frac{q}{A^q} \int_{A}^{\infty} t^q f(x) \left( \frac{x}{A} \right) dt. \]  

To prove Proposition 1 we need to apply the previous lemma to $f(x) = H(x)$ and consider the integral

\[ I(q - 1, z, r) = \int_{M} H(x,Y(x))^{q-1} d\mu(x), \]

which is of the form studied in the Lemma. There are three cases to consider.

1. Case 1: $q > 1$, that is, $q - 1 > 0$ that allows to apply formula [25]. Because of assumption A-1, Eq. (17), the limit in Eq. (25) is null. Moreover, using again Eq. (17) for $t > 1$,

\[ e^{-1}(1 - C_\delta t) \leq \mu\left( \left\{ x \in M \text{ s.t. } H(x) \geq \frac{t}{\mu(B(z,r))} \right\} \right) \leq e^{-1}(1 + C_\delta t). \]  

we can bound the integrand at r.h.s. in Eq. (25):

\[ I(q - 1, z, r) \geq \frac{q - 1}{\mu(B(z,r))^{q-1}} \int_{1}^{\infty} t^{q-2} e^{-t} dt - C_\delta \int_{0}^{\infty} t^{q-1} e^{-t} dt. \]

The last two integrals are convergent. Moreover, observe that Eq. (18) and A-3 imply that $\delta_t$ is uniformly bounded from above. Taking $r$ small enough the term into brackets becomes positive and larger than a quantity independent of $r$. A similar reasoning yields an upper bound for $I(q - 1, z, r)$ of the same form. Since the double integral

\[ \int_{M} \int_{M} H(x,Y(x))^{q-1} d\mu(x)d\mu(y) \]

is the single integral of $I(q - 1, z, r)$ with respect to $d\mu(z)$, the equality (22) follows.

2. Case 2: $0 < q < 1$. Since $q - 1 < 0$ we employ Eq. (26):

\[ I(q - 1, z, r) = \frac{1 - q}{\mu(B(z,r))^{q-1}} \int_{0}^{\infty} t^{q-2} \mu\left( \left\{ x \in M \text{ s.t. } H(x) < \frac{t}{\mu(B(z,r))} \right\} \right) dt. \]

We use again Eq. (28) to get

\[ I(q - 1, z, r) \geq \frac{1 - q}{\mu(B(z,r))^{q-1}} \int_{1}^{\infty} t^{q-2} \left( 1 - e^{-t} \right) dt = \frac{1 - q}{\mu(B(z,r))^{q-1}} S(z, r). \]

In the above we have put

\[ S(z, r) = \int_{1}^{\infty} t^{q-2} \left( 1 - e^{-t} \right) dt + C_\delta \int_{0}^{\infty} t^{q-1} e^{-t} dt. \]

The term $S(z, r)$ is again composed of a constant (the first convergent integral) and of a vanishing quantity (when $r \to 0$), which puts us in the position of using the previous technique to obtain a first inequality between the two terms of Eq. (22).

To prove the reverse inequality we begin by observing that for $t < 1$ condition A-1 simply becomes $\mu\left( \left\{ x \in M \text{ s.t. } H(x) > \frac{t}{\mu(B(z,r))} \right\} \right) - e^{-t} \leq O(\delta_t)$; then we start from Eq. (29) and we part the integral in two, the first from $\mu(B(z,r))$ to one and the second from one to infinity:

\[ I(q - 1, z, r) = \frac{1 - q}{\mu(B(z,r))^{q-1}} (J_1 + J_2). \]

We use again Eq. (28) to get:

\[ J_1 = \int_{\mu(B(z,r))}^{1} t^{q-2} \mu\left( \left\{ x \in M \text{ s.t. } H(x) \right\} \right) dt \leq \frac{t}{\mu(B(z,r))} \int_{\mu(B(z,r))}^{1} t^{q-2} \left( 1 - e^{-t} \right) dt + O(\delta_t) \int_{\mu(B(z,r))}^{1} t^{q-2} dt. \]

The first integral in the above is a positive constant. Let us consider the second term:

\[ O(\delta_t) \int_{\mu(B(z,r))}^{1} t^{q-2} dt \leq \frac{O(\delta_t)}{1 - q} \mu(B(z,r))^{q-1} \]

\[ = O\left( \mu(B(z,r))^{q} \log \mu(B(z,r)) \right), \]

where the last equality follows from $\delta_t = O\left( \mu(B(z,r)) \log \mu(B(z,r)) \right)$, Eq. (18). Therefore, when $r$ tends to zero, this term vanishes.

The case of $J_2$ is easier:

\[ J_2 = \int_{1}^{\infty} t^{q-2} \mu\left( \left\{ x \in M \text{ s.t. } H(x) < \frac{t}{\mu(B(z,r))} \right\} \right) dt \]

\[ \leq \int_{1}^{\infty} t^{q-2} dt = C_2, \]

which is again bounded by a constant.

3. Case 3: $q < 0$. We use again Eq. (30), with $J_1$ and $J_2$ defined in Eqs. (31) and (33), respectively. For the latter integral the inequality (33) still holds, with a different constant $C_2$. To deal with the integral $J_1$ we further split its domain into the intervals $[\mu(B(z,r)), \mu(B(z,r)) \log \mu(B(z,r))]$ and
\[ \mu(B(z, r)) \log \mu(B(z, r)), 1 \], thereby defining the integrals \( J_{1,1} \) and \( J_{1,2} \), respectively.

At this point we use Eq. (20) to estimate from above the integrand of \( J_{1,1} \) and Eq. (19) to do the same for \( J_{1,2} \). Putting the two estimates together, we obtain

\[
J_1 \leq C_3 \int_{\mu(B(z, r))} t^{q-1} dt \leq C_3 \left\{ \frac{q}{q} \right\} (\mu(B(z, r))^q - 1) \leq C_3 \left\{ \frac{q}{q} \right\} (\mu(B(z, r))^q,
\]

(34)

where \( C_3 \) is another constant independent of \( r \) and \( z \).

To get a lower bound we write

\[
J_1 = \int_{\mu(B(z, r))} t^{q-2} \mu \left( H_{B(z, r)} < \frac{t}{\mu(B(z, r))} \right) dt \\
\geq \int_{\mu(B(z, r))} t^{q-2} \mu(H_{B(z, r)} \leq 1) dt \\
= \mu(B(z, r)) \int_{\mu(B(z, r))} t^{q-2} dt = |q - 1|^{-1} \mu(B(z, r))^q \\
\times \left[ 1 - \mu(B(z, r))/q - 1 \right].
\]

In the last step we have used Eq. (21), so that the term in the square brackets is positive and uniformly bounded for \( r \) small enough.

By collecting all the preceding estimates, we get the desired result for all \( q \).

We are now ready to state our result on large deviations of the first hitting time. We first recall that the free energy function \( R(q) \), \( q \in \mathbb{R} \) associated with the process \( \log H_{B(z, r)}(x) \), is given by

\[
R(q) = \lim_{r \to 0} \frac{1}{-\log r} \log \int_M \mu H_{B(z, r)}(x) d\mu(x),
\]

(35)

provided the limit exists. If \( R(q) \) is \( C^2 \) and strictly convex on \( \mathbb{R} \), its Legendre transform is called the rate function \( \hat{Q}(s) \). Our previous Proposition shows that the free energy for the first hitting time verifies \( R(q) = -\tau(1 - q) \) when \( q > 1 \). In this range of values of \( q \), \( R'(q) = \tau'(1 - q) > 0 \) (by Assumption 2) and therefore the supremum for the rate function \( Q \) is attained for positive \( s \) by a value of \( q \) satisfying \( R'(q) = s \). On the other hand, Assumptions A-2 and A-4 immediately imply that for positive \( s \), \( Q(s) \) is a smooth convex function with the minimum at \( D_1 \). Since the free energy is not smooth everywhere, being not differentiable in \( q = 1 \), we cannot use the standard Gärtnert-Ellis theorem, but a local version of it, as it is reported in Lemma XIII.2 in [39]. Let us put \( \Delta = D_1 + 1 \), the above proves the crucial

Proposition 2. Let us suppose that \( \mu \) is not a measure of maximal dimension, which ensures that \( R(q) \) is strictly convex. Then for all \( s \in (0, \Delta) / \Delta \) we have

\[
\lim_{r \to 0} \frac{1}{-\log r} \log \mu(x) \left\{ \frac{\log H_{B(z, r)}(x)}{-\log r} > D_1 + s \right\} = \hat{Q}(D_1 + s),
\]

(36)

where \( \hat{Q}(s) = \sup_q \{q s + \tau(1 - q)\} \).

Remark 1. It is interesting to observe that the Legendre–Fenchel transform of the free energy function \( R(q) \) introduced above allows us to get the rate functions of the large deviations of different processes, namely:

- \( R(q) \) gives the rate function \( Q(s) \) of the information dimension, see (7).
- \( R(q - 1) \) gives the function \( f(\alpha) \) expressing the Hausdorff dimension of the level sets with local dimension \( \alpha \), see (5).
- \( R(q) \) gives the rate function \( \hat{Q}(s) \) of the first hitting time, see (36).
- \( R(q - 1) \) gives the function \( f(\alpha) \) expressing the Hausdorff dimension of the level sets with local dimension \( \alpha \), see (5).

We therefore consider \( R(q) \) an important global tool to analyze and describe the geometric and recursive properties of dynamically invariant measures.

4. Numerical determination of generalized dimensions

The numerical determination of generalized dimensions is a principal concern when experimental data are to be examined, or theoretical hypotheses need to be tested on model cases. Many techniques have been proposed [10,12–14] especially to deal with the case of negative dimensions (i.e., those corresponding to a negative value of \( q \)) that call in case rarified regions of the invariant measure. Via Kač’s theorem, these latter are related to large return times, hence rare events. For this reason a return time approach seems particularly suited to treat this case.

4.1. Hitting time integral

In this section we follow this approach, based upon Proposition 1. We assume that the data at our disposal are finite trajectories of the dynamical system \( (M, \mu, T) \), which we label as \( x_1 = T^s(x_0) \), where the point \( x_0 \) is to be chosen on the attractor of the dynamical system. We set as reference technique the evaluation of the correlation integral in Eq. (1) via a Birkhoff summation. This is effected by first finding the Birkhoff estimate of the measure of a ball \( \mu(B(z_0, r)) \) via

\[
J_N(z_0, r) = \frac{1}{N} \sum_{j=0}^{N-1} \chi_{B(z_0, r)}(x_j) \approx \mu(B(z_0, r)).
\]

(37)

Here and in the following \( \chi_A \) is the indicator function of the set \( A \). The above quantity is then raised to the power \( q - 1 \), and a second average with respect to the point \( z_0 \) is performed:

\[
\frac{1}{N} \sum_{j=0}^{N-1} \chi_{B(z_0, r)}(x_j) \approx \mu(B(z_0, r)).
\]

(38)

While in principle \((x_0, N)\) and \((z_0, N')\) can be different, one can take advantage of the choice \( x_0 = x_0 \) so as to use a single trajectory for the computation. We elect not to do so, for reasons that we will explain later. A complete discussion of the method of correlation integrals is presented in [9].

Once the correlation integral has been estimated (remark that Eqs. (37) and (38) are not scaling relations, but estimates that can be made arbitrarily precise), the problem remains of finding the scaling exponents implied in Eq. (1) or equally the function \( \tau(q) \). Two main ways exist to do this from data at finite resolution \( r \). The first is to employ extrapolation techniques of the ratio log \( J_N(r, q) / \log r \) computed at a number of values \( r \), such as Levin’s algorithm or similar [40]. This is particularly useful when the measure has a hierarchical structure. The second, more conventional, is to try to find a linear least square fit of the log of the correlation integral with respect to the log of \( r \). It is immediate to see that this is a sort of Hôpital’s rule to find the limit for \( r \) tending to zero in Eq. (3).

In the present context, Proposition 1 implies that, for \( q \leq 2 \), the function \( \tau(q) \) can be equally seen as the limit of the ratio on the left hand side of Eq. (22), after the substitution \( q \rightarrow 2 - q \). It requires the computation of the hitting double-integral

\[
\frac{1}{N} \int \int_M H_{B(z, r)}(x)^{1-q} d\mu(x) d\mu(z).
\]

(39)
A Birkhoff estimate of this quantity can be obtained as follows:
first consider the inner integral in the above equation (it was
defined as \( I(1-q, z, r) \) in Eq. (27) of the previous section). This
can be estimated as
\[
I_{y}(1 - q, z_0, r) = \frac{1}{N} \sum_{j=0}^{N-1} \mathcal{H}_{B(z_0, r)}(x_j) \simeq I(1 - q, z_0, r).
\] (40)

In practice, it is convenient to fix \( N \) (and to stop the evaluation
of the motion) as soon as the trajectory of \( x_0 \) has entered the
ball \( B(z_0, r) \) \( H \) times. In so doing, \( N \) becomes a function of \( H, z_0 \) and \( x_0 \). This can be done also when evaluating the
conventional correlation integral, Eqs. (37) and (38). Next, we estimate the
outer integral in (39), again by a Birkhoff summation
\[
\frac{1}{N^r} \sum_{i=0}^{N-1} I_{y}(1 - q, \Gamma(z_0), r) \simeq \mathcal{Y}_{\mu}(q, r).
\] (41)

This procedure has the advantage that the same set of data can be
used to determine an approximation to both the correlation
integral \( \mathcal{Y}_{\mu}(q, q) \) and the hitting integral \( \mathcal{Y}_{\mu}(q, r) \), so that the
two methods can be compared fairly.

As first example of this comparison we choose the Arnold’s cat
map on the two-torus [41], a primary example of chaotic dynamical
system, with the absolutely continuous invariant measure \( \mu \)
given by the Lebesgue uniform measure on this manifold, so
that all generalized dimensions of the measure are equal to two.
In Fig. 1, left panel, we plot the numerically estimated integrals
\( \mathcal{Y}_{\mu}(r, q) \) and \( \mathcal{Y}_{\mu}(q, r) \) versus \( r \) in double logarithmic scale, for
a selected set of values of \( q \) ranging from \( q = -1 \) to \( q = 2 \) and
of \( r \) ranging from \( r = 10^{-3} \) to \( r = 10^{-1} \). The (trivially constant)
data for \( q = 1 \) separate the integrals \( \mathcal{Y}_{\mu}(r, q) \), \( \mathcal{Y}_{\mu}(r, q) \) that grow
from those that diminish when \( r \) tends to zero. The almost linear
shape of the curves confirms the scaling in Eq. (1) and a linear fit
as described above provides an estimate of \( \tau(q) \) and hence of \( D_q \).
Yet, a finer analysis reveals that the asymptotic behavior is not yet
achieved at finite \( r \). In fact, in the right panel the results obtained
using the slope of each linear interpolation between successive
values of \( r \) in the figure are displayed: since values of \( r \) are equally
spaced in logarithmic scale, we define
\[
\sigma_q(r) = \frac{1}{q - 1} \log \mathcal{Y}_{\mu}(q, r) - \log \mathcal{Y}_{\mu}(q, r) - \log r.
\] (42)

with \( \rho < 1 \). The values obtained are not constant: the lowest
set of data, in particular, is related to \( q = 2 \), the highest value
for which generalized dimensions can be obtained in this way.
Its value for \( r = 1.77 \times 10^{-3} \) is still far from the theoretical value
\( D_2 = 2 \), even if \( \sigma_q(r) \) is an acceleration procedure of the limit
in Eq. (3). Nonetheless, a further extrapolation can be performed.
Typically, convergence in these estimates is rather slow, in the
sense that a behavior of the kind
\[
\sigma_q(r) \leftarrow \mu \pm B/\log(r)
\] (43)

holds. Therefore, using \( D_q \) and \( B \) as fitting parameters of the
experimental data, a better estimate of \( D_q \) can be obtained. The
continuous line in the right panel of Fig. 1 plots such approximation.
The obtained value of \( D_q \) is correct to three digits. Finally, for
comparison, the data obtained by using \( \Gamma \) in lieu of \( \tau \) – in other words, the conventional correlation integral – are also reported,
shifted upwards by a small quantity for clarity. Recall that they
have been computed on the same raw data (trajectories) than the
former. They show a reduced precision for small values of \( r \), in
correspondence with the smallest value in the range plotted in
figure, \( q = -1 \).

This instability is shown in Fig. 2, which reports the same
data of the right panel of Fig. 1, for \( q = -1 \). In the same
figure the analogue data obtained by using the first-return time
integral [18,19]
\[
\Gamma_1(q, r) = \int_{\mathcal{M}} \mathcal{H}_{B(r,x)}^{-1}(x) \, d\mu(x)
\] (44)
are also reported. They are evaluated using a small portion of the
data used in the other two cases: in fact, in this case only the
first return is concerned, rather than the \( H \) hits required by the
previous techniques. Despite the fact that a rigorous proof of this
procedure is lacking (see nonetheless [19]) data are consistent
with the expected result.

A second example, when the hypotheses of Proposition 1 are
certainly not verified, is given by the Hénon map, at standard
parameter values [42]. Fig. 3 is the analogue of Fig. 1 in this
second case, with the only difference that in the right panel a
dimensional figure displays the quantity \( \sigma_q(r) \) versus \( r \)
and \( q \), computed from the correlation integral \( \mathcal{Y}_{\mu}(q, r) \) and
the hitting time integral \( \mathcal{Y}_{\mu}(q, r) \). In both cases we observe the
large local fluctuations typical of the Hénon physical measure.
More importantly, a considerable agreement between the two sets of
data is observed, when \( q \) is smaller than two.

In the successive Fig. 4 the generalized dimension obtained
by linear least square fit over the full range of the data in Fig. 3,
left panel, is displayed. It is well known that, in the case of the physical measure on the Hénon attractor, generalized dimensions strongly depend on the range of the fit — as well as on the sampling point chosen in this range. We do not aim to resolve this issue, but we remark that the coincidence between the results obtained by the correlation integral and the hitting time integral suggests that the results of Proposition 1 hold also in this case, which is clearly outside the scope of the hypotheses put forward in the previous section.

Finally, still in Fig. 4, we also plot the curve \( D_2/(q-1) \), which follows from Proposition 1 and describes the scaling of the hitting time integral for \( q \) larger than or equal to two. As in the case of ArnoI’d cat, data for \( q \) approaching two from below are not at convergence, while those for \( q \) significantly larger than two fit the theoretical curve remarkably well.

These experimental data lead us to conclude that the theoretical method to determine generalized dimensions implied by Proposition 1 has a practical value, but, for values of \( q \) between one and two, convergence must be accelerated by suitable techniques. We finally remark that, as conjectured in [18,19], the same can be expected when using first return times.

### 4.2. Using local dimensions computed via EVT

As described above, the key to the computation of generalized dimensions is the estimate of the measure of balls of the same radius \( r \), raised to a power and averaged with respect to the invariant measure. While dimensions are obtained via a scaling relation of these quantities when the radius vanishes, the distribution of such measures at fixed, finite \( r \), is also important. This observation leads to the definition of a finite resolution local dimension, \( D_{1,r}(z) \), which is precisely defined by the equation:

\[
\mu(B(z, r)) = r^{D_{1,r}(z)}. \tag{45}
\]

It is interesting to note the relations of this quantity to extreme value theory. In fact, defining as observable the function

\[
\phi_z(x) = -\log d(x, z), \tag{46}
\]

where \( z \) is the center of the ball in Eq. (45), and computing this latter on a trajectory of the system, \( x_j = T^j(x_0) \), large values of \( \phi_z(x_j) \) correspond to passes of the motion close to the point \( z \). By looking at the statistics of these extreme events-near approaches, one defines the (complementary) distribution function

\[
\tilde{F}_z(u) = \mu(\{x \in M \text{ s.t. } \phi_z(x) > u\}), \tag{47}
\]
which coincides with the measure of the ball of radius \( e^{-u} \) around the point \( z \). From the numerical point of view, as described in [36, Chapters 4 and 6 and references therein], one studies the tail of this distribution, either defined by considering arguments larger than \( u_{\text{cut}} = -\log r_{\text{cut}} \), or by setting a cutoff value in the distribution itself: \( F_2 < 1 - p \), with \( p \) close to one. This second case yields a cutoff value \( r_{\text{cut}} \), which now depends on position. Extreme value theory predicts that the tail distribution, suitably renormalized and shifted, converges for small cutoff to an exponential distribution, whose mean and standard deviation are the inverse of \( D_{1,\text{cut}}(z) \). The latter is then numerically computed as the inverse of the mean of such distribution. In other words, this is an alternative procedure to Eq. (37) that can be used in two ways.

Firstly, it can be turned into a determination of generalized dimensions. Eq. (38) is here replaced by

\[
\gamma_j(q; r, \mu) \approx \frac{1}{N} \sum_{j=0}^{N-1} f^{r(q-1)D_1(x_j)},
\]

(48)

where \( x \) is chosen \( \mu - \text{a.e.} \), \( N \) is supposed to be large and dimensions are obtained by Eqs. (3) and (4). It is sometimes both necessary and convenient not to take the limit of vanishing \( r \) in Eq. (3). In practical applications this is sometimes dictated by the finite resolution of the data and the limited span of time evolution at our disposal. In Section 7 this situation is illustrated by applying the above procedure to the spectrum of dimensions for the North Atlantic atmospheric circulation.

Secondly, the same computations permit to evaluate in a direct way the large deviation function of local dimension: by taking \( I = (s, \infty) \) with \( s > D_1 \) in Eq. (7), we find that \( \mu(D_{1,s} > s) \sim r^{Q(0)} \).

Similarly, we have that \( \mu(D_{1,s} < s) \sim r^{-Q(0)} \) for \( s < D_1 \). Fig. 5 shows the numerically computed rate function \( Q(s) \) for the motion on a Sierpinski gasket defined in Section 5.2, Eq. (58). By lowering the cutoff value of \( r \), approach to the theoretical curve is observed. This theoretical value is given by the rate function \( Q(s) \), which is computed as the Legendre–Fenchel transform of the free energy \( \mathcal{R}(q) = -r(1 - q) \). For the case of motion on the Sierpinski gasket, the function \( \tau(q) \) is explicitly given by formula (59). Numerically, it can be obtained by the techniques described in this article, yielding results for \( Q(s) \) more reliable than those obtained by the direct computation of the distribution of \( \mu(D_{1,s} < s) \). This is undoubtedly an interesting result with potential applications to a wide class of dynamical systems.

5. Generalized dimensions via extreme value theory

In this section we describe a further method to compute the spectrum of the generalized dimensions for positive, integer values of \( q \) larger than one. It has the advantage of using EVT intrinsically and, in addition, it reveals a second spectrum of extremal indices that also possesses a dynamical meaning. This approach is a direct generalization of the method introduced in [23] for the correlation dimension. It is based on the investigation of close encounters, when two or more trajectories of the system approach each other within a small distance. This defines the extreme event that we investigate.

Let us consider the \( q \)-fold \((q > 1)\) direct product \((M, \mu, T)^{\otimes q}\) with the direct product map \( T^q = T \otimes \cdots \otimes T \) acting on the product space \( M^q \) and the product measure \( \mu^q = \mu \otimes \cdots \otimes \mu \). Define the following observable on \( M^q \):

\[
\phi(x_1, x_2, \ldots, x_q) = -\log(\max_{i=2, \ldots, q} d(x_1, x_i)),
\]

(49)

where each \( x_i \in M \). We also write \( \bar{x}_q = (x_1, x_2, \ldots, x_q) \) and \( T_q(\bar{x}_q) = (T(x_1), \ldots, T(x_q)) \).

Fig. 5. Rate function \( Q(s) \) for the motion on a Sierpinski gasket, computed from 1000 sampling points, each of which required a trajectory consisting of \( 10^6 \) iterates.

5.1. Statistics of exceedances

Let us first investigate the statistical distribution of the function \( \phi \), via the (complementary) distribution function \( \bar{F}(u) \):

\[
\bar{F}(u) = \mu_q(\{\bar{x}_q \in M^q \text{ s.t. } \phi(\bar{x}_q) > u\}).
\]

(50)

It is easily seen that

\[
\bar{F}(u) = \int_{M^q} d\mu_q(\bar{x}_q)\chi_{B(\bar{x}_q, e^{-u})}(\bar{x}_q) \cdot \chi_{B(\bar{x}_q, e^{-u})}(\bar{x}_q)
= \int_{M^q} d\mu_q(\bar{x}_q)\mu(B(\bar{x}_q, e^{-u}))\bar{s}^{-1}.
\]

(51)

Comparing Eq. (51) with Eq. (1) yields \( \bar{F}(u) \sim e^{-uD_q/(q-1)} \), so that one can obtain \( r(q) = (q - 1)D_q \) from the asymptotic behavior of \( \bar{F}(u) \) for large \( u \). This quantity can be estimated by a Birkhoff sum, involving the trajectories of \( q \) different initial conditions of the original system. The results of this procedure in the case of the Arnol'd cat dynamical system are reported in Fig. 6, in simple logarithmic scale, for values of \( q \) ranging from \( q = 2 \) to \( q = 8 \). The linear parts of these graphs follow closely the theoretical result \( \bar{F}(u) = \mu_q(\phi > u) = \pi q^{a_u^{-1}}e^{-u\bar{s}(a_u^{-1})} \). The limitations of the procedure are evident from the picture: for large values of \( q \) multiple “encounters” become scarcer and scarcer, so that the linear part, from which generalized dimensions can be extracted by linear fitting, becomes increasingly narrow when the length of the sampling trajectory is finite. Observe that the number of iterations considered in our numerical simulation largely exceed those typically available in real-world applications. On the other hand, this technique is not directly affected by the curse of dimensionality which plagues box counting procedures (but not the correlation integral, or hitting/return times methods, for that matter).

To complete the analysis of the previous section we also examine the case of the Hénon physical measure. Results are reported in Fig. 7, in full analogy with Fig. 6. The exponential decay is evident also here, and the slopes of the curves, together with Eq. (56), permit to extract the data \( r(2) = 1.2, r(3) = 2.32, r(4) = 3.3 \), which imply the generalized dimensions \( D_2 = 1.2, D_3 = 1.16 \) and \( D_4 = 1.1 \). These values compare favorably with the extensive calculations in [43]. Although the exponential decay of the data for larger values of \( q \) is also evident, the data do not allow to estimate the associated dimensions with the same precision.
in [23], that the distribution \( F_n(u_n) \) will converge, for \( n \to \infty \), to the Gumbel’s law
\[
G(t) = e^{-e^{-t}},
\]
which is one of three limiting laws prescribed by the extreme value theory (see Chapter 3 in [36] for a general introduction to extreme value theory). The quantity \( \theta_q \) is called the dynamical extremal index DEI and it will be studied in the next section.

We sketch the argument in the Appendix; here we stress that this convergence can also be investigated numerically and it provides an estimate of the generalized dimensions. In fact, because of Eq. (54)
\[
\bar{F}(u_n) \sim e^{-u_n D_q(q-1)} \frac{t}{n},
\]
and
\[
u_n \sim -\log t + \log n \frac{D_q(q-1)}{D_q(q-1)} = \frac{-\log t}{a_n} + b_n.
\]
The real quantities \( a_n \) and \( b_n \) can be obtained as described in Section 2.3, by a maximum likelihood estimation of the GEV parameters in \( F_n \) [44]. This was also described in Section II-A of [23], which yields the generalized dimensions \( D_q \).

We apply this procedure to the case of an I.F.S. measure [22] on the Sierpinski gasket, defined by the stochastic process on the unit square \( M = [0, 1]^2 \) realized by iteration of the maps \( f_i, i = 1, 2, 3 \) chosen at random with probability \( p_i \):
\[
\begin{align*}
f_1(x, y) &= (x/2, (y + 1)/2), \quad p_1 = 1/4, \\
f_2(x, y) &= ((x + 1)/2, (y + 1)/2), \quad p_2 = 1/4, \\
f_3(x, y) &= (x/2, y/2), \quad p_3 = 1/2.
\end{align*}
\]
The distribution \( F_n \) with block size \( n = 5000 \) is estimated for \( l = 20 \) trajectories of \( 2 \cdot 10^8 \) points. In Fig. 8, left panel, the numerically obtained generalized dimensions are compared with the analytical values [13]:
\[
D_q = \frac{\log_2(p_1 + p_2 + p_3)}{1 - q}.
\]

Good agreement is found for small values of \( q \), which later worsens as expected and discussed earlier. In the same Fig. 8, right panel, we also plot the results for the case of the Lorenz 1963 model [45], a continuous-time dynamical system. Here, the distribution \( F_n \) (with \( n = 10^4 \)) is obtained from trajectories of \( 10^8 \) points, simulated by the Euler method (which is clearly not the best technique, but the focus of our investigation is different) with step size 0.013. Dimensions estimates are obtained by averages over 20 trajectories, uncertainties being the standard deviations of these results.

6. The dynamical extremal index

It is now important to consider the parameter \( \theta_q \) appearing in the exponent of the Gumbel’s law (55); in [23] it was called the Dynamical Extremal Index (DEI). We show in the Appendix that
\[
\theta_q \approx 1 - \lim_{n \to \infty} \frac{\mu_q(\Delta_q^n \cap T_q^{-1} \Delta_q^n)}{\mu_q(\Delta_q^n)}.
\]

For \( C^2 \) expanding maps of the interval, which preserve an absolutely continuous invariant measure \( \mu = hdx \) with strictly positive density \( h \) of bounded variation, it is possible to compute the right hand side of (60) and get:
\[
\mu_q(\Delta_q^n \cap T_q^{-1} \Delta_q^n) \approx \int dx_q h(x_q) \int dx_q h(x_q) x_h(x_{T_q x_q} \cdot e^{-u_n})(T_q x_q) \cdots \int dx_q h(x_q) x_h(x_{T_q x_q} \cdot e^{-u_n})(T_q x_q). \tag{61}
\]
Each of the $q - 1$ integrals above factorizes, and they depend on the parameter $x_k$. Therefore they can be treated as in the proof of Proposition 5.5 in [46], yielding the rigorous result:

**Proposition 3.** Suppose that: the map $T$ belongs to $C^2$; it preserves an absolutely continuous invariant measure $\mu = h dx$, with strictly positive density $h$ of bounded variation; it verifies conditions $P1 = -P5$ and $P8$ in [46]. Then

$$\theta_q = 1 - \frac{\int h(x)^q dx}{\int h(x)^q dx}.$$ (62)

This formula uses the translational invariance of the Lebesgue measure: we refer to sections II-B and II-C in [23] for analogous extensions to more general invariant measures and to SRB measures for attractors. We here extend the DEI introduced in [23] to a spectrum of indicators for well behaved maps of the interval. When the derivative of $T$ is constant, we have a scaling of the kind:

$$\theta_q \sim 1 - e^{-\log(q-1)/h},$$ (63)

where $h = \log |DT|$ is the Lyapunov exponent of the system.

Whenever the derivative exhibits appreciable variations, we expect a deviation from the scaling in Eq. (63). Hence, the variation with $q$ of the quantity

$$H_q = \frac{\log(1 - \theta_q)}{q - 1},$$ (64)

reveals how far we are from the positive Lyapunov exponent.

In particular we expect that such deviation is magnified when:

- the density has a minimum at 0, which usually happens when the derivative blows up to infinity, or when the derivative vanishes somewhere, like in multimodal maps,
- the density is unbounded, which happens when the derivative has an eigenvalue of 1 on periodic points. This is a typical occurrence for intermittent maps.

In both cases it is not possible to bound the second term at the right hand side of $\theta_q$ in Proposition 62 between finite quantities. More importantly, it is not clear that Eq. (62) still holds, since it was proved under the assumptions of a bounded and strictly positive density. As an illustration of this theory, we now give a few examples, for which the DEI can be computed explicitly:

- \(3x - \mod 1\). This map has constant derivative, has Lebesgue as an invariant measure and we have:

$$\theta_q = 1 - 3^{1-q}.$$ (65)

- **Markov maps.** We consider the following piecewise linear Markov map $T$ [47]:

$$T(x) = \begin{cases} T_1(x) = 3x \text{ if } x \in I_1 = (0, 1/3), \\ T_2(x) = 5/3 - 2x \text{ if } x \in I_2 = [1/3, 2/3), \\ T_3(x) = -2 + 3x \text{ if } x \in I_3 = [2/3, 1). \end{cases}$$

The density $h$ of $T$ is given by:

$$h(x) = \begin{cases} h_1 = 3/5 \text{ if } x \in I_1, \\ h_2 = 6/5 \text{ if } x \in I_2, \\ h_3 = 6/5 \text{ if } x \in I_3. \end{cases}$$

The DEI $\theta_q$ can be easily computed by Eq. (62) and it reads:

$$\theta_q = 1 - \frac{h_1^3 + h_2^3 + h_3^3}{h_1^3 + h_2^3 + h_3^3}.$$ (66)

- **Gauss map.** The Gauss map $T(x) = \frac{1}{x} - \lfloor \frac{1}{x} \rfloor$, $x \in (0, 1]$, has a.c. invariant density $h(x) = \frac{1}{\log(2/1/x)}$, which yields

$$\theta_q = 1 - \frac{\sum_{k=0}^{q-1} (-1)^k q^{k-1} q^{k+1} - 1 (2q-1) (q-1) \log 2}{2^{q+1} - 1 - q}.$$ (67)

- **The Hemmer map.** This map [48] is defined on the interval $[-1, 1]$ by $T(x) = 1 - 2\sqrt{|x|}$ and it has the explicit density $h(x) = \frac{1}{2}(1 - x)$ and Lyapunov exponent 1/2. The particularity of this density is that it vanishes in $x = 1$. The DEI for this case reads

$$\theta_q = 1 - \frac{q + 1}{2^q} \sum_{k=0}^{q} \frac{q!}{k!} \left(\frac{1}{k} + (-1)^{k+q-1} \frac{1}{2k + q + 1} \right).$$ (68)

For maps of the interval where Eq. (62) holds, the $H_q$ spectrum can be extended in a first place to real values of $q$, for $q \neq 1$, and then $H_q$ can be defined by continuity applying l'Hôpital rule and we get that

$$H_q = \lim_{q \to 1} H_q = \int \log |DT(x)| h(x) dx,$$ (69)

which is the Lyapunov exponent of the system. We therefore have a full spectrum of such indicators. We show in Fig. 9 the graphical representation of the extended $H_q$ spectrum of the Markovian map defined above, that can be computed explicitly and is given
Fig. 9. Graphical representation of the extended $H_q$ spectrum for the Markovian map described in the text. The limit as $q$ goes to $\infty$ is $\log 2 \approx 0.693$ and $\log 3 \approx 1.099$ as $q$ goes to $-\infty$, that correspond to respectively the logarithms of the minimum and the maximum of the absolute value of the derivative.

by Eq. (66). It is easily seen that for this map, $H_{\infty} = \lim_{q \to \infty} H_q$ is the minimum of the logarithm of the absolute value of the derivative, whereas $H_{-\infty} = \lim_{q \to -\infty} H_q$ is the maximum of this quantity, confirming the fact that the $H_q$ spectrum carries information on the expanding/hyperbolic behavior of the system in different regions of phase space and therefore is a tool for the study of chaotic dynamical systems. In particular, it has the following properties:

- $H_1$ is the Lyapunov exponent of the system, and if the derivative does not vary too much, $H_2$ is a good approximation of the Lyapunov exponent, as stated in [23].
- For large positive values of $q$, $H_q$ characterizes those regions of the phase space where the absolute value of the derivative is low and/or the density is high. The opposite holds for large negative values of $q$.
- If the derivative is constant, so is the $H_q$ spectrum.

6.1. Numerical computations

In support of our results, we perform numerical computations of the $\theta_q$ and $H_q$ spectra, first for the map $x \mapsto 3x \mod 1$. The results are reported in Fig. 10. Extremal indices are computed following the procedure in Section 2.3, using the Süveges estimator [34]. For each $q$, we used as a threshold the 0.997 quantile of the distribution $F(u)$ computed on a pre-run trajectory of $10^6$ points. Results compare favorably with theory for $\theta_q$, but less satisfactorily for $H_q$ when $q$ becomes large. In both cases, it turns out that the fixed threshold corresponds to lower values of the function $\phi$. In the last case, less precise results are also probably due to the fact that $\log(1 - \theta_q)$ diverges as $q$ tends to infinity. Numerical computations for the Gauss map are performed with the same set of parameters as before and are reported in Fig. 11. The same discrepancy of results is found for large values of $q$, due to similar reasons.

6.2. The $H_q$ spectrum for higher dimensional systems

Although higher dimensional systems fall outside the scope of Eq. (62), it was discussed in [23] that $H_2$ is strongly related to the entropy of the system (the sum of the positive Lyapunov exponents), and appreciable variation of $H_q$ with $q$ is the sign of a diversity of local hyperbolic behaviors of the system. On the other hand, a constant Jacobian gives rise to a constant spectrum $H_q$. We check this statement by computing the $\theta_q$ and $H_q$ spectra of a 2-dimensional map of constant Jacobian: the Lozi map with standard parameters [23]. Results are shown in Fig. 12 and are compared with the constant function equal to the positive Lyapunov exponent of the map of around 0.47 [23]. Good agreement is found, although some small systematic fluctuations are present, probably due to finite effects and to the role of the contracting subspace. We have again a discrepancy of numerical results for $q$ larger than 8, limiting the range of computation of the method to small, positive and integer values of $q$.

7. North Atlantic atmospheric variability

In order to show the usefulness of our results in the study of many dimensional, complex systems, we compute the generalized dimensions associated with the atmospheric circulation over the North Atlantic. As observable, we consider the daily sea-level pressure fields observed in the region: [lat 22.5°N–70°N, lon 70°E–50°W] for the period 1948–2015, issued from the NCEP reanalysis dataset [49]. Indeed, the sea-level pressure field is a proxy of the mid-latitude circulation as it traces the position of cyclones–anticyclones thanks to the rotation/stratification properties of atmospheric flows [50].

Fig. 10. Indicators $\theta_q$ and $H_q$ of the map $x \mapsto 3x \mod 1$, obtained from averaging the results over 20 trajectories of $2 \cdot 10^7$ points. The uncertainty is the standard deviation of the results. See text for further discussion.
Considering for this reason, we follow the approach described in [51] and [52] and use as threshold values $T_{p,z}$ for the observable $\phi_z$, the $p$-quantile of the distribution of $\phi_z$, where $p$ is fixed. This ensures that the extreme value statistics is computed with the same sample statistics at all points. The effective radius considered for the computation of the generalized dimensions is then taken to be the average of $e^{-T_{p,z}}$ over $z$. Applying formula (48) to the computation of $D_q$, one obtains the non-linear behavior pictured in Fig. 13. We give a summary of the results found when adopting the above procedure with different quantiles in Table 1. When the quantile is relatively low, a large sample of recurrences is used (corresponding to a larger average cutoff radius). This implies a lower spread of the distributions of $D_1(r,z)$. To the contrary, when the quantile is larger, the sample statistics contains fewer recurrences and the spread in $D_1(r,z)$ increases. Note that, although $\min(D_1(r,z))$ and $\max(D_1(r,z))$ seem to experience large variations with different quantiles, these have to be compared with the dimension of the phase space, which corresponds here to the number of grid points of the sea-level pressure fields used, 1060. The relative variation is therefore very small, less than 1%. Depending on the size of the datasets, one can then look for the best estimates of the $D_1(r,z)$ distribution for several values of $q$ and look for a range of stable estimates in $q$ space.

Generalized dimensions are a piece of information that is typical of the invariant, ultimate measure, that is, the mathematical object that ergodic theory defines in the infinite time limit; nonetheless, as shown in this paper and in the papers quoted in the references, different techniques exist to extrapolate their value from data generated from observations which are not yet

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7 Note that, in this context we also speak of daily dimension and persistence via EVT, meaning that each $z$ is the sea-level pressure field averaged during a day.

References [51] and [52] computed the finite time local7 dimensions $D_1,r(z)$ and persistence $\theta(z)$ of those fields with the method detailed in [51], using as threshold the 98th quantile of the observable distribution. It has been shown that those computations introduce a valuable piece of information in the description of the atmospheric flow at mid-latitudes. In particular, i) the minima of the local dimension $D_1,r(z)$ correspond to zonal flow circulation regimes, where the low pressure systems are confined to the polar regions, opposite to high pressure areas, and ii) the maxima of the local dimension $D_1,r(z)$ correspond to blocked flows, where high and low pressure structures are distributed in East–West direction. In this region, $D_1,r(z)$ takes values between 4 and 25, depending on the spatial resolution of the sea-level pressure field, with an average value of about 12.

In the previous sections we have described tools to compute the generalized dimensions and to link them to local dimensions. In principle, the method described in Section 4.2 requires the computation of a distribution of local dimensions at a uniform resolution $r_{\text{cut}}$. The nature of data at our disposal is not suited to an analysis at fixed resolution, since it may oversample, or undersample the distributions of extreme events of the observable $\phi_z(x) = -\log\{d(z,x)\}$, depending on the point $z$ that we are considering. For this reason, we follow the approach described in [51] and [52] and use as threshold values $T_{p,z}$ for the observable $\phi_z$ the $p$-quantile of the distribution of $\phi_z$, where $p$ is fixed. This ensures that the extreme value statistics is computed with the same sample statistics at all points. The effective radius considered for the computation of the generalized dimensions is then taken to be the average of $e^{-T_{p,z}}$ over $z$. Applying formula (48) to the computation of $D_q$, one obtains the non-linear behavior pictured in Fig. 13. We give a summary of the results found when adopting the above procedure with different quantiles in Table 1. When the quantile is relatively low, a large sample of recurrences is used (corresponding to a larger average cutoff radius). This implies a lower spread of the distributions of $D_1,r(z)$. To the contrary, when the quantile is larger, the sample statistics contains fewer recurrences and the spread in $D_1,r(z)$ increases. Note that, although $\min(D_1,r(z))$ and $\max(D_1,r(z))$ seem to experience large variations with different quantiles, these have to be compared with the dimension of the phase space, which corresponds here to the number of grid points of the sea-level pressure fields used, 1060. The relative variation is therefore very small, less than 1%. Depending on the size of the datasets, one can then look for the best estimates of the $D_1,r(z)$ distribution for several values of $q$ and look for a range of stable estimates in $q$ space.
extratropical storms or blocking) produce large excursions of the local dimensions \(D_{l, r}(z)\), which in turn are associated with large deviations of hitting and return times, in the proximity of special points in phase space. Since the computation of local dimensions is relatively feasible also for systems with a high number of degrees of freedom, this can be used to trace the location of singularities originating the multifractal \(D_q\) spectra.

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**Appendix**

We now sketch how the spectral theory developed in [37] allows to get Gumbel’s law for the distribution of the maxima considered in Section 5.2. We now define:

\[
\Delta_n^q = \{ (x_q^1, \ldots, x_q^n) \in M^n, d(x_q^1, x_q^2) < e^{-u_n}, \ldots, d(x_q^1, x_q^n) < e^{-u_n} \}.
\]

In order to establish rigorous results, we restrict ourselves to \(C^2\) expanding maps of the interval preserving an absolutely continuous invariant measure \(\mu = hdx\), with density \(h\) of bounded variation and strictly positive. The key observation is that the distribution \(\mu_q(M_n \leq u_n)\) can be written as

\[
\mu_q(M_n \leq u_n) = \int d\mu_q \prod_{j=1}^{\lfloor n^{1-q} \rfloor} h(x_j) \prod_{k=1}^{\lfloor n^{1-q} \rfloor} h(x_k),
\]

where \(S'\) denotes the complement of the set \(S\). We now rewrite the preceding formula in a more convenient way. At this regard we first introduce the transfer operator \(P_{n, q}\) for the map \(T_q\), which is defined, for an observable \(f\) in some suitable Banach space \(\mathcal{B}\), (for instance the space of bounded variation functions) by the duality relation:

\[
P_{n, q} f(x_q^1, \ldots, x_q^n) = \int f(x_q^{n+1}) d\mathcal{X}(x_q^{n+1}),
\]

where \(d\mathcal{X}\) denotes the \(q\)-Lebesgue measure. By repeatedly using this duality relation we can write (71) as

\[
\mu_q(M_n \leq u_n) = \int d\mathcal{X} \prod_{j=1}^{\lfloor n^{1-q} \rfloor} h(x_j) \prod_{k=1}^{\lfloor n^{1-q} \rfloor} h(x_k),
\]

where the perturbed operator \(\tilde{P}_{I_{\Delta_n}^q}\) is defined as \(\tilde{P}_{I_{\Delta_n}^q}(f) = P_{n, q}(f | I_{\Delta_n}^q)\). When \(n\) tends to infinity the characteristic function of \(\Delta_n^q\) tends to the identity and the operators \(\tilde{P}_{I_{\Delta_n}^q}\) converge to \(P_{h, q}\) in \(\mathcal{B}\). If the unperturbed operator \(P_{n, q}\) has a spectral gap, which means it allows exponential mixing for the observables in \(\mathcal{B}\), the same is true for the operator \(P_{h, q}\) and the maximal, isolated, eigenvalue of \(P_{h, q}\), which in our case is 1 is close to that of \(\tilde{P}_{I_{\Delta_n}^q}\), \(\lambda_{q,n}\), more precisely [37]:

\[
\tilde{\lambda}_{q,n} \sim 1 - (1 - \chi_0)\mu_q(\Delta_n^q),
\]

where the quantity \(\chi_0\) will be defined in a moment. The operator \(P_{h, q}\) now decomposes as the sum of a projection along the one dimensional eigenspace associated to the eigenvalue \(\tilde{\lambda}_{q,n}\) and spanned by the density \(\tilde{h}^{(n)}\) and an operator with a spectral radius exponentially decreasing to zero and which can be neglected in the limit of large \(n\). This allows us to write

\[
\mu_q(M_n \leq u_n) \sim \tilde{\lambda}_{q,n} \int h(x_1) \cdots h(x_n) \hat{\tilde{h}}^{(n)} d\mathcal{X}.
\]
where the integral on the right hand side converges to 1 for \( n \to \infty \). In conclusion we get by approximating \( \lambda_{q,n} \) as above:

\[
\mu_q(\mathcal{M}_n \leq u_n) \sim \left[ 1 - (1 - x_0)\mu_q(\Delta^q_1) \right]^n \sim \exp\left[ - (1 - x_0)\mu_q(\Delta^q_1) \right].
\]

We now recall that we are under the assumption that \( n\mu_q(\psi > u_n) = n\mu_q(\Delta^q_1) \to t, n \to \infty \); we therefore obtain the Gumbel law \( e^{-qt} \), provided the dynamical extremal index \( \theta_q \) is defined as

\[
\theta_q = 1 - x_0. 
\]  

The term \( x_0 \) is obtained by the previous perturbation theory under the assumption, which is verified in our setting, that the set \( \{ x_i \in M^q, x_1 = x_i, i = 2, \ldots, q \} \) in the product space is left invariant by the direct product of the q maps; in particular we have:

\[
x_0 = \lim_{n \to \infty} \frac{\mu_q(\Delta^q_1 \cap T^{-1}_q \Delta^q_1)}{\mu_q(\Delta^q_1)}. 
\]  

The deduction of the Gumbel’s law with this spectral technique is explained in more detail in the papers \([23, 46]\). In particular, in the paper \([46]\), Lemma 5.2, it was shown how the term \( x_0 \), which in general is given by an infinite series, reduces to the value given above. The adaptation to our case is straightforward and it can be done along the diagonal in each plane \( \{ x_i = x_j \}, j \neq 1 \).

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