Interdisciplinary application of nonlinear time series methods

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Abstract

This paper reports on the application to field measurements of time series methods developed on the basis of the theory of deterministic chaos. The major difficulties are pointed out that arise when the data cannot be assumed to be purely deterministic and the potential that remains in this situation is discussed. For signals with weakly nonlinear structure, the presence of nonlinearity in a general sense has to be inferred statistically. The paper reviews the relevant methods and discusses the implications for deterministic modeling. Most field measurements yield nonstationary time series, which poses a severe problem for their analysis. Recent progress in the detection and understanding of nonstationarity is reported. If a clear signature of approximate determinism is found, the notions of phase space, attractors, invariant manifolds, etc., provide a convenient framework for time series analysis. Although the results have to be interpreted with great care, superior performance can be achieved for typical signal processing tasks. In particular, prediction and filtering of signals are discussed, as well as the classification of system states by means of time series recordings.

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1. Introduction

The most direct link between chaos theory and the real world is the analysis of time series data in terms of nonlinear dynamics. Most of the fundamental properties of nonlinear dynamical systems have by now been observed in the laboratory. However, the usefulness of chaos theory in cases where the system is not manifestly deterministic is much more controversial. In particular, evidence for chaotic behaviour in field measurements has been claimed – and disputed – in many areas of science, including biology, physiology, and medicine; geo- and astrophysics, as well as the social sciences and finance.

This article will take a critical look at the published literature, evaluating the perspectives and the limitations of the approach. While common misconceptions will be elucidated, I will try to adopt a constructive point of view by highlighting those cases where information in fact has been gained by the application of methods from chaos theory.

Along with the treatment of conceptual issues, instructive practical examples from the literature and from work done in our group will be evaluated. I will try to depict the state of the art of the application of chaos theory to real time series. Neither naive enthusiasm to explain all kinds of unsolved time series problems by nonlinear determinism is justified, nor is the pessimistic view that no real system is ever sufficiently deterministic and thus out of reach for analysis. At least, chaos theory has inspired a new set of useful time series tools and provides a new language to formulate time series problems – and to find their solutions.

Previous works of review character are Grassberger et al. [1], Abarbanel et al. [2], as well as Kugiumtzis et al. [3,4]. Apart from a collection of research articles by Ott et al. [5], two books on nonlinear time series from the point of view of chaos theory are available so far, one by Abarbanel [6] and one by Kantz and Schreiber [7]. While in the former volume chaoticity is usually assumed – as already reflected in the title – the latter book puts some emphasis on practical applications to time series that are not manifestly found, nor simply assumed, to be deterministic and chaotic. Apart from these works, a number of conference proceedings volumes are devoted to chaotic time series, including Refs. [8–12]. Nonlinear time series methods that arise as extensions and generalisations of linear tools can be mostly found in the statistical literature. Major references are the books by Tong [13] and by Priestley [14].

The application of nonlinear time series methods to field measurements has been accompanied by considerable controversy in the literature. Early enthusiasm has led to straightforward attempts to find, and even quantify, deterministic chaos in many types of systems, ranging from atmospheric dynamics [15–17] and financial markets [18–20] to heart [21,22] and brain activity [23]. In Ref. [24] it is even claimed that cigarette smoking optimises the “dimensional complexity” of an indicator of brain function. This wave of publications has been followed by a number of critical papers pointing out the methodological deficiencies of the former. Some of these will be cited below in their proper contexts. This review is however not the place to repeat the known arguments or the discussion. In my opinion, part of the controversy and the resulting frustration is due to the misconception that low-dimensional chaos is such an appealing theory that it can be expected to be present generically in nature. Of course, most researchers would deny that they have made such an a priori assumption. Nevertheless, the amount of evidence we require for a “climate attractor”, etc., does depend on how likely, that is, how convincing, we find such a concept. For example, in
medical research it is extremely tempting to have a means of measuring the “complexity” of the cardiac rhythm or even the brain function.

Now, if we assume chaotiocity in the sense of low-dimensional determinism as a starting point of our analysis, we can directly justify the use of delay coordinates by Takens' theorem. The number of degrees of freedom in the system is readily estimated as the embedding dimension where the number of false neighbours drops below the noise floor, the rate of increase of uncertainty is identified with the Lyapunov exponent and so forth. This rationale is pursued for example in Ref. [6]. An experience I share with many other researchers is that this way to proceed is quite dangerous since the possibilities for spurious results and wrong conclusions are overwhelming. As an example, take the analysis of the Salt Lake area data in Ref. [6]. Taking for granted that nonlinear dynamics is at work, locally linear phase space predictions seem most appropriate to forecast future values. The predictions shown in Ref. [6] seem fair enough but a closer inspection of the data shows that already the simple linear rule to follow the trend of the last two observations, $x_{n+1} = 2x_n - x_{n-1}$, is more appropriate than the local linear approach in that it gives forecast errors of about half the rms magnitude.

A second possibility is to try to establish low-dimensional chaos by positive evidence. We could for example look for self-similar geometry over a reasonable range of length scales and demonstrate that uncertainties indeed grow exponentially over a certain period of time. Eventually, we should be able to extract empirical deterministic models that predict future values and that can be iterated to yield time series with statistical properties comparable to the data. This approach is preferred by most theoretical researchers and has been emphasised for example in Ref. [7]. The major drawback is that only exceptional time series show such a clear signature, all of which are from laboratory experiments set up specifically for the study of chaotic phenomena. If this were the only way to go, applications to real world problems would have to be largely abandoned.

Finally, we can move the focus of our study from the question of whether deterministic chaos is really present to the question of whether deterministic chaos provides a useful language for the evaluation of a given signal. The concept that superior performance alone is a valid argument for the use of a particular method is not as surprising for an engineer or clinician as it may be for a physicist. In particular in time series analysis, very few people actually believe that the stock market or the brain actually are linear autoregressive machines. Nevertheless, linear time series methods have been applied to time series from these systems with considerable success. Evidence for the practical superiority of chaotic time series methods has so far been rather scarce.

In the following I will at no point assume deterministic chaos. However, I will first review in Section 2 the fundamentals of dynamical systems theory as the theoretical basis for nonlinear time series methods. In many cases this will only amount to a theoretical motivation; very few facts are rigorously proven for finite, noisy time series. I will briefly review the concepts of dynamical systems, strange attractors, phase space embedding, and the invariant characteristics of a process.

Next I will try to give an understanding of the signatures of determinism in finite observations. Section 3 will discuss what happens to the theoretical concepts when they are applied to real data of finite resolution and length. Some limitations are known rigorously, others can be understood heuristically. Some problems seem to be of purely technical nature but nevertheless may prove to be serious in practice. In particular, extended embedding theorems and amendments of the embedding procedure will be discussed. Estimators for characteristic quantities like dimension, entropy, and Lyapunov exponents are studied with respect to their practical viability. This material
will allow us to gauge for a given time series problem how far we are from the linear case and how close we are to a nonlinear deterministic situation. Accordingly, we will choose either linear methods, coarse but robust nonlinear tools, or more refined phase space methods.

One formal requirement for almost all time series methods is stationarity. Specific tests for nonstationarity in a nonlinear context that have been proposed in the literature are discussed in Section 4, where also hints will be given what can be done in the presence of nonstationarity apart from choosing a different time series problem. In Section 5, formal statistical tests for nonlinearity in a time series will be set up, with particular emphasis on possible nonlinear determinism in the data. The section will state what has to be done in order to perform such a test correctly, but it will also discuss what can (and what cannot) be learned from such a test. While standard solutions for the forecasting and filtering of linearly correlated but otherwise random sequences exist, and methods for strongly deterministic but chaotic systems are also well established, signals of a mixed character are more difficult to deal with. Section 6 will discuss the problems that arise in practice and give some specific successful applications, including medical data analysis. While most estimates of invariants based on short and noisy data are dubious as absolute numbers, many authors have found the comparison of such numbers, or of nonlinear qualitative features, across an ensemble of systems (e.g. sick and healthy patients) quite promising. As discussed in Section 7, some of these works are of little use since the discrimination task could have been solved equally well by standard methods, other examples seem to give results that are by far superior to previous approaches.

Pure low-dimensional determinism is quite special and can be found in nature only to a crude approximation. The range of potential practical applications of nonlinear theory can only by increased significantly if the underlying paradigm is generalised in some respects. Current efforts concerning the analysis of data from extensively chaotic (e.g. spatio-temporal) systems, as well as from mixed, nonlinear and stochastic, sources are discussed in Section 8.

In this paper, I will not give many technical details of the practical implementation of the methods. These will be reviewed in a forthcoming article [25] which is accompanied by the publicly available software package TISEAN\(^1\) which contains many of the methods discussed here.

2. Theoretical foundation

This section will briefly recall the definitions and properties of some concepts of chaos theory, insofar as they are relevant to applied time series problems. Most of the basic concepts are usually formulated in a purely deterministic setting; that is, without any external noise. The time evolution is then given by a dynamical system in phase space. Since usually the state points cannot be observed directly but only through a measurement function, typically involving a projection onto fewer variables than phase space dimensions, we have to recover the missing information in some way. This can be done by time delay embeddings and related methods. We can then quantify

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\(^1\) The TISEAN software package is publicly available for download from either http://www.mpiipsdresden.mpg.de/\(^7\)tsa/TISEAN/docs/welcome.html or http://wptu38.physik.uni-wuppertal.de/Chaos/DOCS/welcome.html. The distribution includes an on-line documentation system.
properties of the system through measurements made on the embedded time series. Since it is
eventually the underlying system we want to characterise, these properties should ideally be
unaffected by the measurement and the embedding procedure. The presentation here serves the
main purpose of fixing the notation for the following section and is therefore extremely brief.
Theoretical issues that are directly related to time series methods are discussed in more detail in the
monographs \cite{6,7}. More general references on the theory of deterministic dynamical systems
include the volumes by Ott \cite{26}, as well as the older books by Bergé et al. \cite{27} and by Schuster
\cite{28}. More advanced material is contained in the work by Katok and Hasselblatt \cite{29}. A gentle
introduction to dynamics is given by Kaplan and Glass \cite{30}. The volume by Tsonis \cite{31} puts more
emphasis on the applied side.

2.1. Dynamical systems and predictability

When we are trying to understand an irregular (which essentially here means non-periodic)
sequence of measurements, an immediate question is what kind of process can generate such
a series. In the deterministic picture, irregularity can be autonomously generated by the nonlin-
earity of the intrinsic dynamics. Let the possible states of a system be represented by points in
a finite-dimensional phase space, say some $\mathbb{R}^d$. The transition from the system’s state $x(t_1)$ at time
$t_1$ to its state at time $t_2$ is then governed by a deterministic rule: $x(t_2) = T_{t_2-t_1}(x(t_1))$. This can be
realised either in continuous time by a set of ordinary differential equations:

$$\dot{x}(t) = F(x(t)),$$  \hspace{1cm} (1)

or in discrete time $t = n\Delta t$ by a map of $\mathbb{R}^d$ onto itself:

$$x_{n+1} = f(x_n).$$ \hspace{1cm} (2)

The family of transition rules $T_t$, or its realisation in the forms (1) or (2), are referred to as
a dynamical system. The particular choice of $F$ (resp. $f$) allows for many types of solutions, ranging
from fixed points and limit cycles to irregular behaviour. If the dynamics is dissipative (area
contracting, the case assumed throughout this work), the points visited by the system after transient
behaviour has died out will be concentrated on a subset of Lebesgue measure zero of phase space.
This set is referred to as an attractor, the set of points that are mapped onto it as its basin
of attraction. Since not all points on an attractor are visited with the same frequency, one defines
a measure $\mu(x)\,dx$, the average fraction of time a typical trajectory spends in the phase space
element $dx$. In an ergodic system, $\mu(x)$ is the same for almost all initial conditions. Phase space
averages taken with respect to $\mu(x)dx$ are then equal to time averages taken over a typical
trajectory.$^2$

In real-world systems, pure determinism is rather unlikely to be realised since all systems
somehow interact with their surroundings. Thus, the deterministic picture should be regarded only
as a limiting case of a more general framework involving fluctuations in the environment and in the

$^2$The notion of an attractor is mathematically difficult to define satisfactorily, see Milnor \cite{96}. The existence of
a natural measure has been proven only for hyperbolic systems, see Eckmann and Ruelle \cite{227}, and for a small number of
specific systems, see for example Benedicks \cite{228}.
system itself. However, it is the limiting case that is best studied theoretically and that is expected to show the clearest signatures in observations.

The deterministic approach is not the most common way to explain irregularity in a time series. The traditional answer given by the time series literature is that external random influences may be acting on the system. The external randomness explains the irregularity, while linear dynamical rules may be sufficient to explain structure found in the sequence. The most general linear (univariate) model is the autoregressive moving average (ARMA) process, given by

\[ x_n = \sum_{i=1}^{M} a_i x_{n-i} + \sum_{i=0}^{N} b_i \eta_{n-i}, \]

where \( \{\eta_n\} \) are Gaussian uncorrelated random increments. The linear stochastic description is attractive mainly because many rigorous results are available, including the properties of finite sample estimators.

Most sources of irregular signals, take for example the brain or the atmosphere, are known to be nonlinear. Nevertheless, if many weakly coupled degrees of freedom are active, their evolution may be averaged to quantities that are to a good approximation Gaussian random variables. If this approximation is valid, it is also reasonable to assume that an observed degree of freedom interacts with the averaged variable in a mean field way, justifying the linear dynamics in Eq. (3). However, there are many situations where this approximation fails, for example if the degrees of freedom of a system act in a coherent way, which can happen in nonlinear systems even when the coupling among the degrees of freedom is weak.

The two paradigms, nonlinear deterministic and linear stochastic behaviour, are the extreme positions in the area spanned by the properties “nonlinearity” and “stochasticity”. They are singled out not because they are particular realistic for most situations, but rather because of their paradigmatic role and their solid mathematical background. Since the literature abounds with premature conclusions like that a system that is not found to be linear must be deterministic instead, let us emphasise that we are dealing with rather narrow limiting cases by drawing a sketch of the world of time series models that also contains all kinds of mixtures, see Fig. 1.

The description of a particular time series by an empirical model will of course be guided by the paradigm adopted for the study. The idea is that one cannot possibly do a better modelling job than recovering the equations that actually govern the system under observation. One should note, however, that a description of a system which includes external influences is only complete if the input sequence (for example in Eq. (3) the sequence of increments \( \{\eta_n\} \)) is known.

One can ask what constitutes a complete description in the case of a deterministic system. In the mathematical sense, the system equations (1) or (2) together with the initial conditions are sufficient. For this to be true, both parts must be known with infinite precision, which is unphysical. If the system is chaotic, then even in the noise free case, errors in the initial condition and in the specification of the model will grow in time and will have to be corrected. Of course, as soon as noise is present, even if only in the measurement procedure, the situation becomes worse. If we could recover \( F \) (resp. \( f \)) from the observations correctly, we still could not simply generate future values by applying \( f \) to the observed present state because we could not account for the noise. Thus, we may face the confusing situation that the original equations of motion do not necessarily give the best model in terms of predictions. (See Refs. [32,33] for illustrative examples.)
Fig. 1. Sketch of the variety of systems spanned by the properties “nonlinearity” and “stochasticity”. Areas where theoretical knowledge and technology for the analysis of time series are available are outlined. Besides the (possibly noisy) periodic oscillations (a), these are mainly the deterministic chaotic and the linear stochastic areas. Also the common routes to chaos (c) and (d), and extensions for small nonlinearity (b) or small noise (e) are marked. There are a few “islands” (f), like hidden Markov models and a few others, where a connection can be made between a nonlinear stochastic model approach and particular real world phenomenon. This sketch has been inspired by a similar representation by O. Kaplan.

Since there is not much that can be done by way of modelling the noise in the system, modelling deterministic systems will still attempt to fit a function \( f \) (or \( F \)) to the data such that Eq. (2) (or (1), resp.) will hold to the best available approximation. The most widespread approach is to solve Eq. (2) in the least-squares sense, minimising

\[
\sum_{n=1}^{N-1} [x_{n+1} - \hat{f}(x_n)]^2.
\]  

(4)

Historical references on the prediction problem include the papers by Farmer and Sidorowich [34], and by Casdagli [35]. Advanced procedures are available that take better care of the noise problem, see Refs. [36,37], as well as Section 3.3. As an alternative to explicitly fitting equations of motion to the data, it has been proposed [38] to synchronise a model system with the observed phenomenon (may be given a posteriori by a time series) as a means of identifying the correct model equations. This approach is conceptually less unequivocal since there exist examples where systems lock into a generalised synchronised mode although the systems are quite different in structure.

2.2. Phase space, embedding, Poincaré sections

An immediate consequence of the formulation of the dynamics in a vector space is that when analysing time series, we will almost always have only incomplete information. Although more and more multi-probe measurements are being carried out, still the vast majority of time series taken outside the laboratory are single-valued. But even if multiple simultaneous measurements are available, they will not typically cover all the degrees of freedom of the system. Fortunately, however, the missing information can be recovered from time delayed copies of the available signal, if certain requirements are fulfilled. The theoretical framework for this approach is set by a number
of theorems, all of which specify the precise conditions when an attractor in delay coordinate space is equivalent to the original attractor of a dynamical system in phase space.

Let \( \{x(t)\} \) be a trajectory of a dynamical system in \( \mathbb{R}^d \) and \( \{s(t) = s(x(t))\} \) the result of a scalar measurement on it. Then a delay reconstruction with delay time \( \tau \) and embedding dimension \( m \) is given by

\[
s(t) = (s(t - (m - 1)\tau), s(t - (m - 2)\tau), \ldots, s(t)) .
\]

The celebrated delay embedding theorem by Takens [39] states that among all delay maps of dimension \( m = 2d + 1 \), those that form an embedding of a compact manifold with dimension \( d \) are dense, provided that the measurement function \( s: \mathbb{R}^d \rightarrow \mathbb{R} \) is \( C^2 \) and that both the dynamics and the measurement function is generic in the sense that it couples all degrees of freedom. In the original version by Takens, \( d \) is the integer dimension of a smooth manifold, the phase space containing the attractor. Thus \( d \) can be much larger than the attractor dimension.

Sauer et al. [40] were able to generalise the theorem to what they call the fractal delay embedding prevalence theorem. It states that under certain genericity conditions, the embedding property is already given when \( m > 2d_{\text{box}} \) where \( d_{\text{box}} \) is the box counting dimension of the attractor of the dynamical system. Further generalisations in Ref. [40] assert that, provided sufficiently many coordinates are used, also more general schemes than simple delay embeddings are allowed. For practical purposes, filtered and SVD embeddings have interesting properties. Generalisations to periodically or stochastically driven systems (where the driving force is assumed to be known) are heuristically straightforward but the relevant genericity requirements are rather involved. The proofs have been worked out by Stark and coworkers [41]. Depending on the application, a reconstruction of the state space up to ambiguities on sets of measure zero may be tolerated. For example, for the determination of the correlation dimension, events of measure zero can be neglected and thus any embedding with a dimension larger than the (box counting) attractor dimension is sufficient [42,229].

Although the embedding theorems provide an important means of understanding the reconstruction procedure, none of them is formally applicable in practice. The reason is that they all deal with infinite, noise free trajectories of a dynamical system.\(^3\) It is not obvious that the theorems should be “approximately valid” if the requirements are “approximately fulfilled”, for example, if the data sequence is long but finite and reasonably clean but not noise free – the best we can hope for in time series analysis. Some of these issues will be discussed in Section 3.3.

When analysing time continuous systems, Poincaré sections are an attractive alternative to the reconstruction with fixed delay times. Instead of the time continuous trajectory, only its intersection points with a fixed surface of section are regarded. Generically, the resulting set has a dimension which is exactly one less than the attractor dimension. This concept is particularly useful if the system is driven periodically and the surface of section can be taken as the hyperplane defined by a fixed phase of the driving force. In this case, in the intersection points are equally spaced in time. Reducing the dimensionality of the problem can be an advantage but it comes at the price of reducing the number of available points for a statistical analysis.

\(^3\) A finite piece of trajectory is a one-dimensional curve which is generically and trivially embedable in three dimensions.
2.3. Quantitative description

A time series is usually not a very compact representation of a time evolving phenomenon. It is necessary to condense the information and find a parametrisation that contains the features that are most relevant for the underlying system. Most ways to quantitatively describe a time series are derived from methods to describe an assumed underlying process. Thus, for example, measures of chaoticity in a time series are usually derived from measures of chaoticity in a dynamical system. The rationale is that a certain class of processes is assumed to have generated the time series and then the measure quantifying that process is estimated from the data. Therefore, it is often necessary to distinguish between the abstract quantity, for example the power spectrum of a stochastic process, and its estimate from a time series, for example the periodogram.

Since the underlying process is only observed through some measurement procedure, it is most useful to attempt to estimate quantities that are invariant under reasonable changes in the measurement procedure. As will be seen below (in Section 3.4), the finite resolution and duration of time series recordings damage the invariance properties of quantities which are formally invariant for infinite data. If the value of an observable depends on the observation procedure it looses its value as an absolute characteristic. While in some cases we can still make approximate statements, the interpretation of results has to be undertaken with great care. If we want to compare the results between different experiments, at least we have to unify the measurement procedure and the details of the analysis. If in a realistic situation invariance has been given up anyway, the quantities discussed in the following are no longer singled out that strongly among all possible ways of turning a time series into a number. Consequently, there is no lack of ad hoc definitions and characteristics that have been used in the literature. Since they are invariably defined for time series, rather than the underlying processes, some of them will be discussed later in Section 3.

2.3.1. Linear observables

In the linear approach to time series analysis, a quantitative characterisation of a process is done on the basis of either the two-point autocovariance function or the power spectrum. If only a finite time series \( \{s_n, n = 1, \ldots, N\} \) is available, the autocovariance function can be estimated, e.g. by

\[
C(\tau) = \frac{1}{N - \tau} \sum_{n=\tau+1}^N s_n s_{n-\tau}.
\]

Depending on the circumstances, other estimators may be preferable. A whole branch of research is devoted to the proper estimation of the power spectrum from a time series. The simplest estimator, known as the periodogram \( P_k \), is based on the Fourier transform of \( \{s_n\} \),

\[
S_k = \sum_{n=0}^{N-1} s_n e^{i 2\pi kn/N}
\]

through \( P_k = |S_k|^2 \). Issues of spectral estimation will not be discussed here. An introduction and pointers to the literature can be found for example in *Numerical Recipes* [43].

According to the Wiener–Khinchin theorem, the power spectrum of a process equals the Fourier transform of its autocovariance function. For finite time series this is only true if either \( C(\tau) \) is computed on a periodically continued version of \( \{s_n\} \), or \( P_k \) is computed on a version of \( \{s_n\} \) that is
extended to \( n = -N, \ldots, N \) by padding with \( N \) zeroes. Nevertheless, both descriptions contain basically the same information, only that it is presented in different forms. Furthermore, there is a direct connection between the power spectrum and the coefficients of an ARMA model, Eq. (3), yielding a third possible representation.

The power spectrum of a process (and its autocovariance function) is unchanged by the time evolution of the system (if it is stationary, see Section 4 below). However, it is affected by smooth coordinate changes, e.g., by the characteristics of a measurement device. Usually, the noninvariance of the power spectrum is not a serious drawback. The power spectrum is most useful for the study of oscillatory signals with sharp frequency peaks. The location of these peaks is conserved, only their relative magnitude may be affected by the change of coordinates, see Refs. [230, 231].

Sharp peaks in the power spectrum indicate oscillatory behaviour and are useful indicators in linear as well as in nonlinear signals. Broad band contributions, however, have a less clear interpretation since they can be either due to deterministic or stochastic irregularity. Therefore, the power spectrum is only of limited use for the study of signals with possible nonlinear deterministic structure.

2.3.2. Lyapunov exponents

The hallmark of deterministic chaos is the sensitive dependence of future states on the initial conditions. An initial infinitesimal perturbation will typically grow exponentially, the growth rate is called the Lyapunov exponent. Let \( x_k \) and \( x_l \) be two points in state space with distance \( \| x_k - x_l \| = \delta_0 \leq 1 \). Denote by \( \delta_{\Delta n} \) the distance after a time \( \Delta n \) between the two trajectories emerging from these points, \( \delta_{\Delta n} = \| x_k + \Delta n - x_l + \Delta n \| \). Then the Lyapunov exponent \( \lambda \) is determined by

\[
\delta_{\Delta n} \approx \delta_0 e^{\lambda \Delta n}, \quad \delta_{\Delta n} \ll 1, \quad \Delta n \gg 1.
\]

A positive, finite, value of \( \lambda \) means an exponential divergence of nearby trajectories, which defines chaos. A mathematically more rigorous definition will have to involve a first limit \( \delta_0 \to 0 \) such that a second limit \( \Delta n \to \infty \) can be performed without saturation due to the finite size of the attractor.

Here, only the single (maximal) Lyapunov exponent will be discussed. Lyapunov spectra can be defined that take into account the different growth rates in different local directions of phase space. However, the nonleading exponents are notoriously difficult to estimate from time series data. Only in very few cases of clean laboratory time series trustworthy results have been obtained so far. (See [7] for a discussion of the arising problems.) For field data, Lyapunov spectra beyond the first exponent have not so far been demonstrated to be a useful concept.

There have been a number of attempts to generalise the Lyapunov exponent to systems which are not purely deterministic. For the usual definition, an arbitrarily small amount of noise leads to a diffusive separation of initially close trajectories and a divergent Lyapunov exponent (mind the order of the two limits involved). For very small noise levels, there may still be a range of length scales where the separation proceeds exponentially, until the finite size saturation is reached. This is the behaviour that is probed by the real space methods of estimating Lyapunov exponents from data, in particular the two very similar algorithms introduced independently by Rosenstein et al. [44] and by Kantz [45]. From the theoretical point of view, intermediate length scale definitions are less attractive since the resulting quantities are no longer invariant under smooth coordinate transformations.
While the connection between chaoticity and fractality arises naturally, there are counterexamples, both for strange but non-chaotic and for chaotic but non-fractal attractors.

An alternative way to introduce noise into the definition of Lyapunov exponents is to study the separation of initially close trajectories of two identical copies of a system which are evolving subject to the same noise realisation. Then the Lyapunov exponent quantifies the contribution to the divergence that originates in the intrinsic instability of the deterministic part of the system. This is essentially the kind of instability probed by the tangent space methods to obtain Lyapunov exponents from data, most prominently Refs. [46–49].

Lyapunov exponents quantify the average exponential growth rate of infinitesimal initial errors. Their natural units are therefore inverse times. However, this does not justify to quote inverse Lyapunov exponents as average predictability horizons or predictability times. (The two processes of averaging and taking the inverse of a quantity do not commute.) In fact, the degree of instability and predictability can vary considerably throughout phase space, as it has been pointed out for example by Abarbanel et al. [50] and by Smith [51]. The Lyapunov exponents constitute a particular way of averaging over these variations. They are constructed in a way such that the average becomes independent of the initial condition and invariant under smooth coordinate changes. For typical prediction times, one has to form different averages which cannot be expected to be invariant. It has been argued that the loss of information about the system by averaging in a specific way over local variations of the instability or predictability is too severe. Several people [52,53] have therefore proposed concepts of local Lyapunov exponents and predictabilities. Local Lyapunov exponents are defined in a quite similar way as the usual exponents, except that the limit $\Delta n \to \infty$ is omitted, whence they become position dependent, or local. In particular, Bailey et al. [54] have studied the statistical properties of these exponents. They consider the case that dynamical noise is perturbing the system. In that case they can prove a central limit theorem about the existence and convergence of finite time Lyapunov exponents. Unfortunately, local quantities are almost never invariant in any useful sense. Quite trivially, they will change whenever the positions are transformed. These changes may easily been kept track of. But as soon as the coordinate changes are not isometries, the statistical weights of different areas in phase space are changed. Thus, the values, and distributions of values, of local Lyapunov exponents and local predictability times are manifestly non-invariant.

2.3.3. Dimension and entropy

Besides the exponential divergence of trajectories, the most striking feature of chaotic dynamical systems is the irregular geometry of the sets in phase space visited by the system state point in the course of time. This fractal geometry is a natural consequence of the divergence of trajectories which can be realised in a finite phase space only through some folding mechanism. Stretching, folding, and volume contraction lead to statistically self-similar structure on small length scales.\(^4\)

While the average stretching rate is quantified by the Lyapunov exponent, the loss of information due to the folding is reflected by the entropy of the process. The self-similar character of the resulting point sets and measures defined on them can be characterised by fractal dimensions. Several definitions of non-integer dimensions have been proposed in the literature. Most well

\(^4\) While the connection between chaoticity and fractality arises naturally, there are counterexamples, both for strange but non-chaotic and for chaotic but non-fractal attractors.
known is the Hausdorff dimension of a set and the more easily computable box counting (or capacity) dimension. We can also weight the points in the set by the frequency with which they are visited on average. Then we need a definition of the dimension in terms of the natural measure \( \mu(x)dx \) defined on the set.

One way to proceed is to take weighted averages of the number of points contained in the elements of a partition of phase space and study their dependence on the refinement of the partition. The translation of this scheme into a time series context leads to the box-counting methods of dimension estimation. The practical problems that arise when a space of moderate dimensionality must be covered by boxes of small length can be overcome by sophisticated bookkeeping algorithms. However, these methods make rather inefficient use of the statistics available and suffer from severe finite size effects on the larger length scales. They are therefore not recommended for the study of invariant properties of real world time series.

An alternative way to define the dimension of a measure \( \mu(x)dx \) is by means of correlation integrals \( C_q(x) \). Let us define the locally averaged density \( \rho_e \) to be the convolution of \( \mu \) with a kernel function \( K_e(r) = K(r/e) \) of bandwidth \( e \) that falls off sufficiently fast for the convolution to exist:

\[
\rho_e(x) = \int_y d y \mu(y) K_e(||x - y||) .
\] (9)

Most commonly, the kernel is chosen to be \( K_e(r) = \Theta(1 - r/e) \) where \( \Theta(\cdot) \) is the Heaviside step function, \( \Theta(x) = 0 \) if \( x \leq 0 \) and \( \Theta(x) = 1 \) for \( x > 0 \). Other kernels are popular in statistical density estimation. The correlation integral of order \( q \) is given by the order \( q \) average of \( \rho_e \):

\[
C_q(e) = \int_x d x \mu(x)[\rho_e(x)]^{q-1} .
\] (10)

For a self-similar measure we have

\[
C_q(e) \propto e^{(q-1)D_q}, \quad e \to 0 .
\] (11)

In the literature, \( D_q \) is called the order-\( q \) dimension. This definition includes the dimension \( D_0 \) that has been shown to coincide with the Hausdorff dimension in many cases, and the information dimension \( D_1 \) through l’Hospital’s rule. Although \( D_1 \) is the most relevant because of its information theoretic meaning – it quantifies the scaling of the amount of information needed to specify the state of the system with the required accuracy – we will usually at most be able to estimate a lower bound on it, the correlation dimension \( D_2 \). The correlation dimension as a means of quantifying the “strangeness” of an attractor has been introduced by Grassberger and Procaccia [55]. For finite samples, the double integral in \( C_2 \) can be evaluated down to much smaller scales than the other \( C_q \)'s. Although generic attractors are expected to be multifractal, that is, \( D_q \) depends on \( q \), this property is difficult to study in real time series. Only for exceptionally long, clean signals, \( D_q \) can be obtained for \( q \neq 1,2 \). For real world recordings it is already an ambitious goal to establish a possible fractal nature by finding a scaling region of \( C_2 \).

When analysing time series we are usually dealing with distributions of delay vectors with delay \( \tau \) in an \( m \)-dimensional reconstructed phase space. The \( m \) dependence of \( C_q \) in the limit of large \( m \) can then be expressed as

\[
C_q(m,e) = c(m) e^{-(q-1)h_{str} e^{(q-1)D_q}}, \quad e \to 0, \ m \to \infty
\] (12)
which defines the order $q$ entropy $h_q$. The pre-factor $x(m)$ depends on the norm $\| \cdot \|$ and the kernel function. Although $x(m)$ does not affect the asymptotic value of the entropy $h_q$, the convergence for finite $m$ can be dramatically different, as demonstrated in Ref. [56]. (See also Ref. [57].) Again, the case $q = 1$ is singled out since the (Shannon or Kolmogorov) entropy $h_1$ is additive when independent processes are joined. Also, $h_1$ is related to the Lyapunov exponents via Pesin’s identity, a fact that can be used for consistency checks. However, as with the dimensions, the case $q = 2$ is much more accessible with time series data. See for example Ref. [58]. An algorithm for the determination of the Kolmogorov entropy is given in [59].

Equivalent scaling behaviour to that of Eq. (12) is valid for a large class of kernel functions in the average Eq. (9), see Refs. [60,61]. Besides the hard kernel given by the Heaviside function, the most natural choice is a Gaussian, so that for example $C_2$ reads:

$$C_2^G(\varepsilon) = \int_{x,y} d\varepsilon \mu(y) \mu(x) e^{-|x-y|^1/2\varepsilon^2} .$$

Apart from yielding smoother curves for finite sample estimates, Gaussian kernel correlation integrals have some other attractive properties. For example, $\log C_2^G$ is additive under the pointwise summation of independent variables, in particular, a deterministic signal and measurement noise. It is not so much in use mainly since its numerical implementation seems quite awkward when the definition, Eq. (13) is used straightforwardly. However, it is quite easily obtained from the usual (step kernel) correlation integral by

$$C_2^G(\varepsilon) = \frac{1}{2\varepsilon^2} \int_0^\infty d\tilde{\varepsilon} e^{-\tilde{\varepsilon}^1/4\varepsilon^2} \tilde{C}_2(\tilde{\varepsilon}) .$$

### 2.4. Comparing dynamics and attractors

The quantities considered so far were all meant to characterise a single process. Different processes can of course be compared by comparing these numbers. It may however be interesting to have some means to answer the question of how different two processes are directly, without going through the reduction to a small number of characteristics. Recently, several authors independently have begun to use relative measures for the classification of systems through time series [62,63] and for the study of nonstationary signals [64–66]. Therefore, some theoretical background will be given for the less ad hoc measures of dissimilarity. Practical aspects as well as some more informal but useful quantities will be taken up in Section 3.6.

Let us first make a clear distinction between the problem of defining a measure of dissimilarity between attractors (dynamics, measures, probability distributions) and between trajectories. The latter is related to the question if two systems are dynamically synchronised in a general sense. Generalised synchronisation means that there is a smooth mapping that relates the states of the two systems at any time. The present work does not address the latter question. Relevant references include Refs. [67–71].

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5 It can be easily seen that $C_2^G(\varepsilon) = \int_0^\infty d\tilde{\varepsilon} e^{-\tilde{\varepsilon}^1/4\varepsilon^2} d/d\varepsilon C_2(\tilde{\varepsilon})$ from which the result follows by partial integration.
Information theory provides a measure of distance between two probability densities \( \mu(x) \) and \( \nu(x) \) which is based on the Kullback entropy [72]. Let

\[
H_K(\mu, \nu) = \int x \nu(x) \log \frac{\mu(x)}{\nu(x)} .
\] (15)

Then

\[
\gamma_K(\mu, \nu) = H_K(\mu, \nu) + H_K(\nu, \mu) = \int x (\mu(x) - \nu(x)) (\log \mu(x) - \log \nu(x))
\] (16)

is positive definite, symmetric, and fulfills the triangle inequality. Thus \( \gamma_K(\cdot, \cdot) \) is a metric. For the same reasons discussed above, expression (16) will be replaced by its analog based on second order correlation integrals. A distance \( \gamma_2(\mu, \nu) \) can then be defined by\(^6\)

\[
\gamma_2(\mu, \nu)^2 = \int x \left[ \mu(x) - \nu(x) \right]^2 = \lim_{\varepsilon \to 0} \left[ C_2(\varepsilon; \nu) + C_2(\varepsilon; \mu) - 2C_2(\varepsilon; \mu, \nu) \right] ,
\] (17)

where \( C_2(\varepsilon; \mu, \nu) \) is the cross-correlation integral

\[
C_2(\varepsilon; \mu, \nu) = \int \int_{x,y} \mu(x) \nu(y) K_{\varepsilon}(||x - y||) .
\] (18)

The case that \( K_{\varepsilon}(r) = \Theta(\varepsilon - r) \) as the generalisation of the Grassberger–Procaccia correlation integral [55] has been introduced by Kantz [73]. The Gaussian kernel case \( K_{\varepsilon}(r) = e^{-r^2/4\varepsilon^2} \) together with its finite sample properties has been studied by Diks and coworkers [74]. If the proper limit \( \varepsilon \to 0 \) is taken, \( \gamma_2(\mu, \nu) \) is nothing but the \( L^2 \) distance of the two probability densities. For finite \( \varepsilon, C_2(\varepsilon; \nu) + C_2(\varepsilon; \mu) - 2C_2(\varepsilon; \mu, \nu) \) is no longer formally a distance, except for particular kernels including the Gaussian case.

One drawback of the second order distance \( \gamma_2(\mu, \nu) \) as compared to the Kullback distance \( \gamma_K(\mu, \nu) \) is that it is no longer invariant under smooth coordinate transformations acting on both distributions. Other measures of distance, in particular for discrete point sets, have been discussed by Moeckel and Murray [75]. One could further ask what happens if two distributions are observed but they may have been obtained with different measurement functions. This amounts to the question if the two distributions are absolutely continuous with respect to an unknown reference distribution. The hope for a time-series-based answer seems unrealistic at this stage.

The quantities mentioned so far are all based on geometrical ideas. Apart from asking for similar phase space geometry one can also ask for similar dynamical evolution laws. If approximate predictive models can be established for the time series, one can derive measures of dissimilarity that often allow stable estimates for rather short sequences. Kadtke [76] uses global models of the form

\[
s_{n+1} = f(s_n) = \sum_{i=1}^{M} a_i f_i(s_n)
\] (19)

\(^6\) Strictly speaking, this formula is only valid for smooth distributions \( \mu, \nu \). For measured data we can safely assume that smoothness is imposed on fractal distributions by measurement errors.
to fit several time series or segments individually. Then changes and differences in the dynamics are monitored by changes in the model parameters \( a_i \). Technically, it is important to choose a model class with as few basis functions \( f_i \) as possible. Otherwise the values of individual coefficients \( a_i \) may strongly depend on unimportant details of the signal. Two numerically distinct sets of parameters may equally well model the same data. This approach requires the dynamical models to involve globally adjustable parameters, excluding locally constant or locally linear methods.

There are at least two other ways to compare the dynamics of two predictive models \( f \) and \( g \). One approach that has been taken by Hernández and coworkers [62] is to use both models to make predictions on a time series \( \{ s_n \} \) and compare them for each time step \( n \):

\[
\gamma_p(f, g; \{ s_n \})^2 = \frac{1}{N} \sum_{n=1}^{N} [f(s_n) - g(s_n)]^2 .
\]

(20)

Since \( \gamma_p(f, g) \) is nothing but the \( L^2 \) distance of the vectors formed by the individual predictions, it is a distance measure in the mathematical sense. In general its value depends on the choice of time series \( \{ s_n \} \). If \( f \) has been obtained by a fit to the signal \( \{ s_n^{(1)} \} \) and \( g \) by a fit to \( \{ s_n^{(2)} \} \), a symmetric distance measure between \( \{ s_n^{(1)} \} \) and \( \{ s_n^{(2)} \} \) is given by

\[
\gamma_p(f, g)^2 = \gamma_p(f, g; \{ s_n^{(1)} \})^2 + \gamma_p(f, g; \{ s_n^{(2)} \})^2 .
\]

(21)

The cross-prediction error used in Refs. [63,66] is quite similar, but it compares the individual predictions to the observed values rather than to each other:

\[
\gamma_c(f; \{ s_n \})^2 = \frac{1}{N-1} \sum_{n=1}^{N-1} [s_{n+1} - f(s_n)]^2 .
\]

(22)

A small value of \( \gamma_c(f; \{ s_n^{(2)} \}) \) indicates that the dynamics on \( \{ s_n^{(2)} \} \) is a subset of the dynamics found on \( \{ s_n^{(1)} \} \) and modelled by \( f \). A symmetric measure of dissimilarity (not a formal distance measure in general) is given by

\[
\gamma_c(f, g)^2 = \gamma_c(f; \{ s_n^{(2)} \})^2 + \gamma_c(g; \{ s_n^{(1)} \})^2 .
\]

(23)

In the last two schemes, in principle any method of prediction can be used. In Refs. [63,66] stable results have been obtained with simple locally constant phase space predictors. Prediction errors will be discussed in Section 3.5. Section 3.6 will discuss a few practical aspects of the above approaches.

### 3. Nonlinear analysis of limited data

In the previous section, some definitions and theoretical motivation were given for a number of concepts which now have to be adapted to the case that instead of a measure, an attractor, a dynamical system, all we have is a finite, noisy time series. The way to proceed crucially depends on the point of view we want to assume about the nature of the system. As said earlier, we cannot assume deterministic chaos for any measured time series. If we want to use the theoretical results available, we need to establish it from the data, maybe backed up by additional considerations. Often, we will not be able to find low-dimensional structure, but we may still borrow some concepts just because they give a convenient framework for certain problems.
3.1. Embedding finite, noisy time series

How much information can be recovered from time delayed copies of finite sets of noisy measurements is quite a complicated question and a general answer is not available. The embedding theorems mentioned previously all assume that the observations are available with arbitrary precision. For some results, in particular those concerning the attractor dimension, it is also assumed that arbitrarily small length scales can be accessed which implies that an infinite amount of data is available. A mathematical theorem cannot simply be expected to be almost valid if the conditions are almost fulfilled. Consequently, several authors have investigated what happens to the embedding procedure when noise is present and the sequence is of finite length. For the embedding procedure, noise seems to be the dominant limiting factor.

Only a few theoretical results relevant for practical work are available on the embedding of noisy signals. First of all, we have to make a fundamental distinction between noise due to measurement error and noise that is intrinsic to the dynamics. In the first case, we suppose that there is a deterministic dynamical system underlying the signal. Thus, it is clear what we want to reconstruct by the embedding procedure. If the noise is coupled to the system we have to specify in what sense we want to use an embedding in the first place. Unfortunately, the nature of the noise is usually not known independently. There is no general straightforward way to infer its properties from a time series without making strong assumptions about the dynamical system or the spectral properties of the noise.

One remarkable paper about the effect of measurement noise on the embedding procedure is that by Casdagli and coworkers [77]. Their main result is that a reconstruction technique that leads to a formally valid embedding with noise-free data can nevertheless amplify noise even in a singular way. That means that in such a case not all degrees of freedom of the system can be recovered from a scalar time series even for arbitrarily small amounts of noise. The examples studied in Ref. [77] suggest that this situation is quite typical and not just found in constructed pathological examples. Thus, bold interpretations of Takens' theorem, for example, that we can recover the full dynamics of the human body from a recording of a single variable, is not only in contradiction with common sense but also disproven by mathematical arguments.

Some results are available on the embedding of noise-driven signals. One line of thought supposes that the driving noise sequence is known. The dynamical system then becomes a nonlinear input–output device. Casdagli [78] and Stark et al. [41] formalise the idea that the observations of the output of such a system can be embedded in the sense that time delayed copies of the observation sequence together with the state of the input variable specify the state of the system equally well as the full output together with the input state. These results are maybe more useful for time series analysis than they may seem, given the fact that we almost never know the noise sequence. There are certain signals where the external influence can be inferred to some extent from the observed output. Consider, for example, a recording of the cardiac cycle, for example an electrocardiogram (ECG). The cycle itself is fairly regular but the initiation of a new cycle seems not to be fully determined by degrees of freedom of the heart itself. But even if the beat times were random, we could always infer a posteriori that triggering must have occurred once we observe a new cycle. An illustration of this point will be given in an example below.

There are other theoretical works that also follow the idea that dynamical noise can be isolated in certain cases where the observations contain sufficient redundancy. Muldoon and coworkers
study the case that more probes are available than necessary to cover the degrees of freedom of the system. They demonstrate in a number of examples that sufficient redundancy in the measurements allows for a distinction between the deterministic part of the signal and the dynamical noise. This allows also to recover missing variables by an embedding procedure.

3.2. Practical aspects of embedding

In most interdisciplinary applications we do not know much about the nature of the noise. For example, biological systems are almost never isolated, and measurements are always of finite accuracy. Observational noise is not always white and Gaussian, although this is often the case. If we make any assumption about the statistical properties of the noise, we have to carefully check the consistency of the results.

One of the most immediate restrictions of the embedding theorems for finite data is that the information contained in a time delay representation of real data is influenced by the choice of embedding parameters. While the theorems do not restrict the delay time \( \tau \) (only a few exceptional cases must be excluded), the proper choice of \( \tau \) does matter for practical applications. Also, there are many cases where the theoretically sufficient embedding dimension \( m \) is not optimal for a certain purpose. Larger (but sometimes also smaller) values may give superior results. The literature on this issue is quite confusing and at times contradictory. Part of the confusion is due to the fact that optimality can only be assessed with respect to a specific application. When fitting the dynamics by a global polynomial model, the embedding dimension should be as small as possible in order to limit the number of coefficients in the model. On the other hand, for local projective noise reduction, the redundancy of an embedding with small \( \tau \) and large \( m \) allows for better noise averaging. For signal classification we do not even need a formal embedding since the difference between states may be statistically better defined in a low-dimensional projection where small neighbourhoods tend to contain more points.

The discussion will therefore be cut short by giving some pointers to the literature and by proposing simply to carry out each study with several embedding strategies and to compare the results. Theoretical work on the embedability of noisy sequences is found in Refs. [77,80,81]. More heuristic studies are Refs. [82—88]. One general remark is that if one attempts to formally optimise the performance of an embedding, one should use a large enough class of possible embeddings. There is no theoretical reason to restrict the study to time delay embeddings with equally weighted lags that are integer multiples of a common lag time. In fact, it was already reported in Ref. [1] that minimising the redundancy in a reconstruction \( s(t - \tau_{m-1}), s(t - \tau_{m-2}), \ldots, s(t) \) does not necessarily yield \( \tau_n = n\tau_1 \) but for example \( \tau_1 < \tau_2 < 2\tau_1 \). More general functions of \( s(t') \), \( t - w \leq t' \leq t \) in a time window of length \( w \) may be considered. The singular value decomposition constitutes the special case of maximising the variance among all linear combinations within a time window. But neither is it necessary to consider linear functions only, nor is the variance always the most interesting characteristic.

Let us finish this section with a particular type of signal where a time delay embedding proves useful even though the signal has a strong stochastic (or high-dimensional, in any case unpredictable) component, the electrocardiogram (ECG). The ECG records the electro-chemical activity of the heart which is essential for its pumping mechanism. The cardiac muscle can be regarded as a spatially extended excitable medium with an excitable, an excited, and a refractory phase. At the
onset of a cardiac cycle, a stimulus is initiated at the *sino-atrial (SA) node*, a specialised collection of muscle cells. The resulting depolarisation wave proceeds along a well-defined pathway, first through the atria and then to the ventricles. The excited tissue contracts and thereby ejects blood to the body and the lungs. Eventually, all cardiac tissue has been excited and is refractory whence the stimulus dies out. (If this condition fails, then the *re-entry* phenomenon can occur which is the cause of serious arrhythmiae.) The pathway of the depolarisation wave is quite similar from cycle to cycle, the variation over a few cycles can be parametrised approximately by a one- or two-parameter family of ECG curves. However, the onset of a new cycle fluctuates from beat to beat in a way that is not well understood. Certainly, the interbeat fluctuations cannot be modelled successfully by a low-dimensional deterministic approach. Coupling to the breath activity, blood pressure, as well as more complex control signals from the central nervous system have to be taken into account. With this picture in mind, it could not be expected from the embedding theorems that a delay coordinate representation of the ECG is useful at all.

Let us consider a stochastically driven, damped harmonic oscillator as a toy model for such an input–output system with an unknown, fluctuating input sequence:

\[
\ddot{x} + \dot{x} + x = a(t) .
\]  

(24)

The driving term is taken to be zero except for kicks of random strength at times \( t_i \) such that the inter-beat intervals, \( t_i - t_{i-1} \) are random in the interval \([p, q]\). Fig. 2 shows two trajectories of such a system with different choices of the inter-beat time interval \([p, q]\). The kicks are realised by finite jumps by a random amount in the interval \([0, 1]\). To the left, solutions of Eq. (24) are plotted versus

![Fig. 2. Trajectories of a kicked, damped harmonic oscillator. Left: signal plotted versus time. Middle: true two-dimensional phase space. Right: delay embedding. Upper: kicks occur close in time. Lower: kicks are well separated in time and the system can relax between kicks. See text for discussion.](image)
time. In the middle, the true phase space spanned by $x(t)$ and $\dot{x}(t)$ is shown while to the right a delay representation is used with a delay of one time unit. In the upper row, beats are initiated with a time separation of $p = 0$ to $q = T/2$ where $T = 4\pi/\sqrt{3} \approx 7.26$ time units is the period of oscillation. This does not allow the system to relax sufficiently between kicks in order to form characteristic structure in phase space. Consequently, an embedding provides no clear picture and additional information would be needed for an analysis of such a time series. In the lower row, no kicks were closer in time than $p = T$, the maximal separation being $q = 3T$. The inter-beat parts of the trajectory are distinct because of the different kick strength, but since this is the only fluctuating parameter except for the inter-beat interval, they are essentially restricted to a two-dimensional manifold. This manifold is preserved under time delay embedding although neither the sequence of beat times nor the beat amplitudes are used in any way. The randomness acts locally around the origin in the indeterminacy as to when the next beat will occur and how far the system will be taken by the kick.

Similarly to this toy example, most of the lacking information in the ECG, that is, the times at which a new beat is triggered and possible other parameters of the new cycle, can be deduced from the recording itself. Once the cycle is on its way, we can find its origin quite easily. Thus, the redundancy in the ECG trace explains why delay representations of ECGs are found to be approximately confined to low-dimensional manifolds, see for example Fig. 3. In the left panel delay has been set to 12 ms in order to resolve best the large spike (the QRS-complex) that corresponds to the depolarisation of the ventricle (the large loop in Fig. 3). In the right panel, a longer delay of 24 ms has been used in order to unfold the smaller structures around the origin which represent the atrial depolarisation (P-wave) and ventricular re-polarisation (T-wave) phases. After each beat, the signal retires to the baseline, or the origin in delay space, where it may spend some time until a new beat occurs. Since at that point the future is indeterminate, the set visited by the trajectory is essentially finite dimensional, without being formally deterministic.

3.3. Estimating dynamics and predicting

Perhaps the most fundamental idea behind the approach to time series analysis taken in this paper is that an irregular signal may have been generated by a nonlinear dynamical system with...
only a few active degrees of freedom. Therefore, one of the most important goals should be to establish effective equations of motion that follow this principle and are consistent with the data. The ability to generate a time series that is equivalent to the measured one can be taken as evidence for the validity of the approach, and is therefore interesting in its own right. But there are many other situations where effective model equations can be of great value.

Most properties of chaotic systems are much more easily determined from equations than from a time series. Thus, if a time series can be well represented by model equations, one might even abandon the analysis of the series in favour of an analysis of the model. But this situation is rather rare. Except for well controlled laboratory experiments, dynamical modelling is seldom faithful enough to justify such an approach. Nevertheless, analysing an empirical model, and maybe synthetic time series data generated from it, can provide a valuable consistency test for the results of time series analysis. The best we can hope for when fitting a model to data is that the result comes close to the real underlying dynamics. However, chaotic dynamical systems generically show the phenomenon of *structural instability*. This means that models with very similar parameters may exhibit qualitatively different global dynamics, for example close to an attractor *crisis* (see for example Ref. [26]). Therefore, if we simply iterate fitted model equations, we may see substantially different behaviour from the actual system even if the model in itself is faithful. One way to moderate this danger is to introduce a small amount of dynamical noise comparable to the modelling error when iterating the equations. Dynamical noise softens the sensitive dependence on parameters to some extent. Alternatively, or additionally, one could study *ensembles* of models which are compatible with the data. The ensemble variation of statistical properties can then be taken as an indicator for the expected effect of the remaining modelling error.

The most obvious reason to reconstruct model equations from a time series is that one may be interested in predictions of future values. This task is of course quite common in meteorology and finance and several other fields. Moreover, in many situations the average error when predicting a time series can be taken as an indicator of the structure present in the signal. Thus, we will often use a nonlinear prediction error as a quantifier for the comparison of signals. The use of prediction errors will be discussed in Section 3.5. A choice of models that can be employed will be described below.

In the context of nonlinear dynamics, the modelling task consists in estimating the function \( F \) or \( f \) that is supposed to generate the data via Eqs. (1) and (2) respectively. This may look like the common problem of estimating a nonlinear function, and there is indeed a close relation. However, the information available is not quite what one would like to have for that purpose. All we usually have is a noisy scalar time series:

\[
\begin{align*}
    s_n &= s(x_n) + \xi_n, \\
    x_n &= f(x_{n-1}) + \eta_n.
\end{align*}
\]  

(25)

Here, also an intrinsic noise term \( \eta \) has been included, since no real system is ever really isolated. Since we cannot completely recover \( \{x_n\} \) from \( \{s_n\} \), the best we can do is to use some kind of embedding of \( \{s_n\} \) and look for a mapping \( f_s \) that acts on the embedding vectors. If the original phase space is \( d \)-dimensional this mapping may have to be defined in up to \( 2d \) dimensions according to the theory of embeddings. However, information about \( f_s \) is only given through the data, that is, on the attractor. This may render the estimation problem singular, depending on the model class from which \( f \) shall be estimated.

Even if the embedding problem can be solved (or avoided, if multiple simultaneous measurements are available) the estimation problem remains difficult. The standard approach would be to
choose some parameter-dependent model for $\hat{f}$, and optimise the parameters using a maximum likelihood or least-squares procedure. This however implies that the value $y = f_s(x)$ is known at a number of locations, perhaps with some uncertainty. But for the usual procedure to work, the locations $x$ have to be given without error. This cannot be assumed in time series analysis because $f_s$ is sampled only at the noisy data points. This and other practical problem in estimating dynamics from a time series are discussed in Kostelich [36]. In Refs. [32,89] illustrative material and a partial solution can be found. The most thorough discussion, the one that also comes closest to a satisfactory solution of the problem, is offered by Jaeger and Kantz [37]. The solution involves two ingredients. The first is to replace the ordinary least-squares procedure by a procedure that also optimises the positions $M_x$ (sometimes called total least-squares). Unfortunately, this renders the fitting problem nonlinear, even if the model class is a linear combination of basis functions. The second part requires to minimise the one-step error $(y - f_s(x))^2$ but simultaneously optimising the precision of $f_s(x), f_s(f_s(x)), \ldots$, a difficult nonlinear minimisation problem which is quite computer time intensive.

As for the model class from which $\hat{f}$ is to be determined, a number of different propositions have been made. One possibility favoured by many authors is to expand the dynamics in Taylor series locally in phase space. This has first been proposed in the context of Lyapunov exponent estimation by Eckmann and coworkers [46]. They perform local linearisation on time series data to obtain the dynamics in tangent space. In an classical paper [34], Farmer and Sidorowich port the idea to the prediction problem.

In practice, the expansion is carried out up to at most linear order. Since one has to work in several dimensions, the number of coefficients in higher order approximations becomes too large for a local treatment. In $m$-dimensional delay coordinates, the local model is then quite simply

$$s_{n+\Delta n} = a_0^{(m)} + \sum_{j=1}^{m} a_j^{(m)} s_{n-(j-1)\tau},$$

where $\Delta n$ is the time over which predictions are being made and $\tau$ is the time delay as usual. The coefficients $a_j^{(m)}, j = 0, \ldots, m$, may be determined by a least-squares procedure, involving only points $s_n$ within a small neighbourhood around the reference point $s_m$. Thus, the coefficients will vary throughout phase space. The fit procedure amounts to solving $m + 1$ linear equations for the $m + 1$ unknowns.

When fitting the parameters $a$, several problems are encountered that seem purely technical in the first place but are related to the nonlinear properties of the system. If the system is low-dimensional, the data that can be used for fitting will locally not span all the available dimensions but only a subspace, typically. Therefore, the linear system of equations to be solved for the fit will be ill conditioned. However, we are only interested in that part of the linear map (26) which relates points on the attractor to their future. There are several ways to regularise the least-squares problem. In the presence of noise, the equations are not formally ill-conditioned but still the part of the solution that relates the noise directions to the future point is meaningless (and uninteresting). Equivalently to adding a small amount of noise one can add a small factor times the unit matrix before the singular matrix is inverted. The most appealing approach however is to restrict the fitting procedure to the directions spanned by the data which can locally be identified with the principal components or singular vectors of the data distribution. These and a few
other regularisation schemes for locally linear predictions are discussed in great detail by Kugiumtzis and coworkers [90]. In his contribution to the Santa Fe Institute time series contest in 1991, Sauer [91] has emphasised the close interplay between phase space embedding and fitting of the dynamics.

The optimal degree of locality of a locally linear modelling approach has been used by Casdagli [92] as a measure for nonlinearity in a time series. He compares the predictive quality of models fitted with using different numbers of neighbours. In the absence of nonlinearity, the globally linear fit using all available points as neighbours should give best results since it uses the largest number of points and is structurally more robust. For increasing degrees of nonlinearity, the tradeoff between lack of statistics with few neighbours and curvature error with large neighbourhoods should move the optimum towards smaller and smaller length scales. The rationale of this paper is quite attractive – to define the degree of nonlinearity by what is the most useful assumption for modelling.

If one is interested in a robust, low variance measure of nonlinear predictability without necessarily aiming at optimal forecasting power, one should consider using locally constant approximations to the dynamics. The idea is simply that determinism will cause similar present states to evolve into similar future states. (This idea has been first used for predictions by Lorenz [93] who called it the method of analogues.) Since we expect the signal to be noisy, it is advantageous to consider a collection of similar states rather than the single most similar state observed so far (Lorenz’ analogue). Thus, in order to make a prediction on the point \( s_n \), we form a neighbourhood \( \mathcal{U}_n \), either with a fixed radius or a fixed number of elements. The prediction model is that \( s_{n+\Delta n} = a^{(n)} \) where \( a^{(n)} \) may vary throughout phase space. The fitting problem then degenerates to finding the local (in phase space) average over the future points of \( s_k \), \( k \in \mathcal{U}_n \):

\[
s_{n+\Delta n} = \frac{1}{|\mathcal{U}_n|} \sum_{k \in \mathcal{U}_n} s_{k + \Delta n}.
\]

Here, \( |\mathcal{U}_n| \) denotes the number of points in the neighbourhood. In its simplest implementation, all that has to be set are the embedding parameters and a length scale, the radius of phase space neighbourhoods. Usually, none of these parameters need to be determined by a fit. Therefore the locally constant model can be regarded as almost parameter free. Very similar algorithms has first been used by Pikovsky [94] and later by Sugihara and May [95]. It is quite popular in the context of nonlinearity testing [85], see also Section 5. Extensions to the locally constant or linear multivariate function interpolation procedure by spline smoothing and adaptive parameter optimisation are implemented in the MARS (multivariate adaptive regression splines) package which is described in [97].

An approach that is quite different from those mentioned so far attempts to fit the dynamics by a nonlinear function that is globally defined in phase space. The price for having a single expression for the model function is that it has to accommodate the nonlinear structure appropriately. The most straightforward generalisation of the linear autoregressive models is to include higher-order polynomial terms in the ansatz for \( \mathbf{f}(s) \). This is sometimes called a Volterra series expansion. Another popular model class are radial basis functions. Their use is documented by a large body of papers, for example Powell [98], Broomhead and Lowe [99], Casdagli [35], and Smith [100]. As usual, the model is given by a linear combination of orthogonal functions \( \Phi_i(s) \). Here, the functions are essentially of the same form, radially symmetric about some centre point \( s_i^{(c)} \):

\[
\Phi_i(s) = \Phi(||s - s_i^{(c)}||).
\]
This yields
\[ \hat{f}(s) = a_0 + \sum_{i=1}^{M} a_i \Phi(||s - s_i^{(c)}||) . \] (28)

Quite a variety of functional forms of \( \Phi(r) \) can be used, including Gaussians or powers. For the art of choosing the centre points \( s_i^{(c)} \) (which do not have to be points from the data set), reference is made to the literature – individual experimentation is also recommended.

After giving examples of model classes which are most popular in the nonlinear time series community, it should be stressed again that multivariate function estimation is a common problem in statistics and has a rich literature. Neural networks have been very fashionable over the past few years. They have proven to have remarkable capabilities, but, as the huge literature indicates, they need a lot of expertise and experience to be used reliably. Another branch of research pursues the idea of a regression tree in order to organise multidimensional structure for prediction, see [101] for a classical reference.

If any parameter dependent model is used, for example one like Eq. (28), or a neural network, care has to be taken to avoid overfitting. Overfitting means that a larger model class can always increase the accuracy of a fit, even though at some point only the details of the particular realisation of the process that is available for fitting are accommodated. This problem can be avoided by limiting the number of adjustable coefficients in the model. There are at least two ways to look at the problem. Both result in a penalty for the number of coefficients in the cost function that is to be minimised, but the exact form of this penalty differs for the two approaches. Akaike [102] observes that if we are planning to use the model for making predictions, then what we want to minimise is not the least-squares error of the model on the data but the expectation of that error for the case that the model is applied to new data from the same source. This can mean that a model obtained from such a fit may outperform the original equations when prediction errors are compared. The second approach is due to Rissanen [103]. The idea here is that modelling provides a way to represent a data set in a more compact way than storing the measured data. If the data can be reproduced by a simple model, one can store the model and the errors instead, which is enough to recover all the information. In such a context, fitting is a tradeoff between reducing errors and increasing the size of the model. Practically usable formulas have only been derived in this context for linear (AR) models. One of the reasons is that it is difficult to compare the importance of parameters across different model classes.

In order to avoid overfitting in practical problems, one has to validate the predictive power of a model on yet unused data. Simply splitting the available data into two parts, one for fitting and one for testing, is the cleanest possibility, but unfortunately is quite wasteful in terms of data usage. Alternatively, one can use a cross-validation technique, as they are common in the statistical literature. For \( k \)-fold cross-validation, the data set is split into two segments in \( k \) different ways and the fitting and testing is repeated \( k \) times. At each time, a different part of the data is used as a test set and the remaining points are used for fitting. The errors of the \( k \) tests may then be averaged together. (For the correct averaging, in particular if the different test sets overlap, consult the statistical literature.) The advantage is that testing is done on all available points, making best use of the data base. At the same time, each fit is based on a large part of the time series. If the expense in computer time that is necessary to repeat the fit \( k \) times is feasible, \( N \)-fold cross-validation can be used on \( N \) data points, thereby optimising the available statistics for the fits. This case is sometimes
called *take-one-out* statistics for obvious reasons. Strictly speaking, \( k \)-fold cross-validation assumes that there are no serial dependencies between data in the different segments, which may or may not be true. If in doubt, one can ensure that the training section and the test section are sufficiently far apart in time. With this restriction, locally constant and locally linear predictors provide take-one-out out-of-sample errors automatically since the current point has to be excluded from the neighbourhoods. In the presence of serial correlations, one should also exclude temporally close points that may still be dynamically related.

It should be stressed that cross-validation is a technique for model verification and not for model optimisation. If several models are proposed and their out-of-sample errors are compared, the error quoted for the best of these models can no longer be regarded as an out-of-sample statistic since it is obtained by optimisation over a known training set.

Section 3.5 will discuss a few issues that arise when the error with respect to a nonlinear prediction scheme is to be used as a quantifier for the predictability in a system, or a measure of “complexity”. If such a quantifier is only used in a relative way to compare different signals, it is not formally necessary to use out-of-sample errors. In fact, unless cross-validation at high \( k \) is carried out, in-sample errors usually have lower sample variance and may therefore give better discriminative power. For the case of nonlinearity testing, this has been discussed by Theiler and Prichard [104].

### 3.4. Estimating invariants

All quantitative indicators of chaos involve in their definitions some kind of limit. If these indicators are to be estimated from a finite time series measurement, none of these limits can actually be carried out. Formally, the desirable theoretical properties of these indicators, in particular their invariance under smooth coordinate transformations, will be lost. If the indicator is to be computed for the purpose of a comparative study, lack of invariance may be compensated for by standardising the characterisation procedure. Nonlinear indicators which are not necessarily invariant but are optimised for their power to discriminate between different dynamical states are discussed in Section 3.5 below.

With the typical data quality in nonlaboratory experiments, and given that pure low-dimensional determinism is quite a particular phenomenon in nature, we will very seldom be able to reliably estimate the proper dimension or Lyapunov exponent of a real-world phenomenon. The issues we have to consider in such an attempt will be discussed below — most of them are well covered by the literature, for example the book by Kantz and Schreiber [7]. One of the sharpest critics of naive use of the Grassberger–Procaccia correlation dimension describes the situation roughly like that: There is little use in computing the correlation dimension; If it is less than three one does not need it because the structure of the attractor is obvious, if it is larger than three it cannot be estimated reliably.\(^7\) Very similar statements can be made for the other invariants from chaos theory as well. Later, the discussion will proceed to quantities that have less theoretical value but are easier to compute or give statistically more powerful results.

\(^7\) P. Grassberger, private communication.
3.4.1. Lyapunov exponents

Lyapunov exponents measure the rate of divergence of initially close trajectories. A positive but finite Lyapunov exponent is therefore a sharp criterion for the existence of deterministic chaos. The older literature on the determination of exponents from time series can be seen as extensions of techniques that have been developed for the analysis of systems with known evolution equations. From these they inherit the assumption that there actually exist such dynamical equations and trajectory separation evolves indeed exponentially. Sano and Sawada [47] as well as Eckmann and coworkers [46] introduce locally linear fits to the dynamics in order to follow the evolution in tangent space. The algorithm by Wolf et al. [105] follows several nearby trajectories to measure the average increase of local volume. Many refinements of these methods have been proposed, see Ref. [106] for a comparative discussion of all but the most recent Lyapunov algorithms. It is not wise to use these algorithms when it cannot be taken for granted that the dynamics is deterministic since none of them actually verifies the exponential behaviour of trajectories.

More recently, the emphasis has shifted from the estimation of exponents under the assumption of determinism to the verification of exponential growth of errors. Very similar algorithms for this purpose have been proposed independently by Rosenstein et al. [44] and by Kantz [45]. We follow the latter reference here. The key idea is that initially close trajectories do not necessarily diverge exactly exponentially, but only on average. In order to cancel fluctuations around the general exponential growth, one has to average appropriately over many trajectory segments. Let $W$ denote a set of delay reconstructed points $s_k$ selected at random from a long trajectory such that they approximate the true probability distribution. Let $|W|$ denote the number of members in $W$. The set of points in an $\varepsilon$-neighbourhood of $s_k$ is denoted by $U_k$. Now define

$$S(\Delta n) = \frac{1}{|W|} \sum_{k \in W} \ln \left( \frac{1}{|U_k|} \sum_{i \in U_k} |s_{k+\Delta n} - s_{i+\Delta n}| \right).$$

(29)

If the distances $|s_{k+\Delta n} - s_{i+\Delta n}|$ grow like $e^{\lambda \Delta n}$, then so does $\exp S(\Delta n)$, but with less fluctuations. Kantz [45] discusses why this is the correct (that is, unbiased) way to average. A plot of $S(\Delta n)$ versus $\Delta n$ must show a reasonably straight line over a range of length scales before we accept its slope as an estimate of the Lyapunov exponent $\lambda$. In order to find such a scaling range one has to choose the radius $\varepsilon$ of the neighbourhoods $W$ as small as possible, but not so small that too few neighbours are found or that distances are dominated by noise.

Examples for the successful use of this approach for computer generated sequences and for time series from low-dimensional systems in laboratory experiments have been given in the original articles by Rosenstein et al. [44] and by Kantz [45], as well as in Ref. [7]. In many time series from field measurements, initially close trajectories are found to diverge rapidly. Algorithms that assume that this divergence is due to an intrinsic instability of the dynamics will then issue a positive Lyapunov exponent. However, an intrinsic instability of a chaotic dynamical system should result in exponential growth of the discrepancies, which is difficult to establish. Consider, for example, the divergence rate plot shown in Fig. 4. For a sequence of 2000 time intervals between heartbeats in a normal human, the function $S(\Delta n)$ defined in (Eq. (29)) has been computed using unit delay, initial neighbourhoods $\mathcal{W}$ of diameter 0.015 s around each point were formed in $2-10$ dimensions. Indeed, the trajectories diverge quite fast. However, they reach a saturation soon. The curves are not fitted well by straight lines which would be the case if the divergence was exponential. If one were to
assign a slope to the lines, the result would strongly depend on the length scale, embedding dimension, etc., and would be quite useless as an estimator of the Lyapunov exponent. It should, however, be remarked that the growth is not simply diffusive either. A non-invariant parameter, like the time of growth from $S(0)$ to twice its value for given fixed embedding and neighbourhood parameters may well be useful for the comparison of different subjects but may not have much relation to a possible intrinsic instability of the system.

### 3.4.2. Correlation dimension and entropy

The problems that arise when correlation integrals and the correlation dimension are estimated from finite time series have been discussed extensively in the literature. Statistical estimators for fractal dimensions and their theoretical properties are studied in Refs. [107–111]. Original contributions pointing out potential sources for spurious results are found for example in Refs. [112–116]. Some of the material has been reviewed for example in Refs. [1,7,117,118]. I will therefore only briefly state the main points.

If the probability distribution implied by the natural measure is approximated by a sum of delta functions at $N$ points $\{x\}$ independently drawn from it, we can estimate the correlation integral $C_2$ by the correlation sum

$$\hat{C}_2(\varepsilon) = \frac{2}{N(N-1)} \sum_{i=1}^{N} \sum_{j=i+1}^{N} \Theta(\varepsilon - \|x_i - x_j\|).$$

(30)

Thus, $\hat{C}_2(\varepsilon)$ is just the fraction of all pairs of points that are closer than $\varepsilon$. It has been shown, for example, by Grassberger [119] that $\hat{C}_2$ is an unbiased estimator of $C_2$. The hat will henceforth be suppressed. In time series applications, the assumption that the $\{x\}$ are independently drawn from the underlying distribution is usually violated due to serial, also called temporal correlations.

It has been pointed out by several authors that serial correlations can lead to spurious results for the correlation dimension, see for example Refs. [1,112,113,117]. The necessary correction is also well known (it has been proposed by Grassberger [115] and by Theiler [117]): pairs of points $i,j$ that are closer than some correlation time $t_{\text{min}}$ have to be excluded from the double sum in
(Eq. (30)). The loss of statistics is not dramatic since the total of pairs grows like $N^2$ while only a number of terms $\propto N$ is suppressed. Therefore it is advisable to be generous when choosing $t_{\text{min}}$ the time scale given by the decay of the linear autocorrelation function is often not sufficient. A useful tool to determine the decay of nonlinear correlations is the space-time separation plot introduced by Provenzale et al. [120], see Section 4, (Eq. (37)) below. The effect of failure to exclude serially correlated pairs from the correlation sum can be seen by comparing Fig. 5 of the present section and Fig. 9 of Section 4.

It should be remarked that the literature might give the wrong impression that the sensitivity to serial correlations is a flaw specific to the Grassberger–Procaccia dimension algorithm. In fact, linear correlations and nonlinear determinism are sources of predictability which are detected by any algorithm that does not explicitly exclude the structure imposed by one of these sources. This is the reason why also, for example, prediction errors or false nearest-neighbour techniques have to be augmented by a comparison to linearly correlated random surrogates in tests for nonlinearity, unless similar corrections are carried out as for the correlation dimension. For the false nearest-neighbour approach [121] this has been pointed out for example in Ref. [122].

If one wants to estimate a correlation dimension, plotting $C_2(\varepsilon)$ in a log–log plot is not always the best thing to do since deviations from the desired power law scaling do not appear very pronounced in this representation. A better method might be to plot the local slopes of the log–log plot,

$$D_2(\varepsilon) = \frac{d \log C_2(\varepsilon)}{d \log \varepsilon} ,$$

versus $\log \varepsilon$. These slopes can, for example, be obtained by a straight line fit over a small range of values of $\varepsilon$.

Theiler [123] gives a maximum likelihood estimator of the Grassberger–Procaccia correlation dimension. Since maximum likelihood estimation of the correlation dimension goes back to Takens [109], such quantities are often referred to as Takens’ estimator. The estimator is given by

$$D_{\text{ML}}(\varepsilon) = \frac{C_2(\varepsilon)}{\int_0(\varepsilon') C_2(\varepsilon')/\varepsilon' \, d\varepsilon'} .$$

Fig. 5. Maximum likelihood estimator of the correlation dimension as a function of the cutoff length scale and the embedding dimension for an intracranial recording [131] of the neural electric potential in a human. No scaling region of approximately constant $D_{\text{ML}}$ can be found. The time series was provided by Lehnertz and coworkers.
This quantity can also be plotted against log $\varepsilon$ for different values of the embedding dimensions. Its advantage over $D_2(\varepsilon)$ is that it incorporates all the statistical information that is available below the length scale $\varepsilon$. It is however implied that the dominating contamination at the small length scale is given by the effect of the finite size of the data set. In many practical situations this is not quite the case since measurement errors destroy the self-similarity as well. The effect of noise on the correlation integral has been studied in a number of papers [108,124–127] which are reviewed and compared in Ref. [128]. Olofsen and coworkers [129] as well as Schouten and coworkers [130] derive maximum likelihood estimators of the correlation dimension for data which are contaminated with noise. However, the noise amplitude enters the analysis as an unknown parameter which complicates the application of their results in practical situations.

If all precautions are taken and a dimension estimate is attempted on a complex data set, one should not be surprised if the result is negative in the sense that no scaling region and no proper saturation can be found. In fact, few real systems are sufficiently low-dimensional for this kind of analysis. As an example that illustrates this statement but which also shows that estimates of the correlation sum can be useful nevertheless, let us study an intracranial recording of the electric field in the brain of a human. The data and its analysis is described thoroughly in Elger and Lehnertz [131]. The data has been taken in an epilepsy patient and one of the questions is whether individual seizures can be anticipated from these recordings a few minutes before they occur. I will report on this application in more detail in Section 4.4. Huge amounts of data are available since multiple channel recordings are taken at 173 Hz continuously over several days for pre-surgical screening purposes. Not surprisingly, however, the data are quite nonstationary which is in fact essential for seizure anticipation to be possible. If one wants to assign an attractor dimension to such a time series, one first has to select windows in time which are short enough so that the dynamics can be considered to be effectively stationary within each window. A reasonable tradeoff between approximate stationarity and time series length is obtained with windows of 30 s duration. From this data, the maximum likelihood estimator of the correlation dimension, (Eq. (32)), is obtained as a function of the upper cutoff length scale $\varepsilon$ and plotted versus log $\varepsilon$. Clearly, there are no scaling regions where $D_{\text{ML}}$ becomes independent of $\varepsilon$ and $m$. Thus, it may be concluded that a low-dimensional attractor is not a good model for this data set. However, it will be discussed in Section 4.4 below how the correlation sum, although not suitable for a dimension estimate, could be used for monitoring changes in the brain state between different segments of a long recording.

3.5. Non-invariant characterisation

So far, invariance of an observable has been emphasized as a requirement for the objective characterisation of time series data. However, we have also seen that estimation of truly invariant quantities is an ambitious goal that is worth pursuing only with sufficiently high data quality and for systems from the appropriate class. This does not imply that we have to give up a quantitative description in all other cases. An absolute, portable characterisation is not always indispensable. Non-invariant characterisation of time series data can, for example, be useful for the comparative
study of multiple data sets. Of course, if we want to use an observable for comparisons that is not invariant under changes in the measurement procedure, we have to standardize the measurement procedure. Practical issues concerning the comparison of time series are discussed in Section 7. Let us here only discuss some nonlinear time series measures which are not invariant under coordinate changes but which have been used in the literature for various reasons, for example because of their robustness to noise and to small sample sizes.

Both noise and limited numbers of data points constitute severe problems whenever properties at small length scales in phase space have to be probed. In the presence of noise, length scales below the noise amplitude cannot be accessed without an explicit noise reduction step. Also, on small length scales the discrepancy between the finite collection of points and an underlying probability distribution becomes most pronounced. The obvious way to get away from these problems is to use coarse-grained quantities which are defined on intermediate length scales. While giving up invariance, the statistical properties of these quantities are often favourable. The most extreme step in this direction is to encode the time series as a symbol string and then analyse this sequence of discrete values. This approach is often called symbolic dynamics, although in the dynamical systems literature this name is reserved to a symbolic description that results from an encoding according to a generating partition. (See, for example, [132–134].) In the latter case, the symbol string has the same entropy as the full time series and, in fact, the full trajectory of a dynamical system can be in principle recovered from the (bi-infinite) symbol sequence. The partitions which are commonly used for symbolic encodings of time series data are almost never generating in that sense. Since any further refinement of a partition that is generating preserves this property, one often replaces a generating partition by a very fine ad hoc partition. Upon further refinement, the partition becomes “approximately generating” to a higher and higher degree. However, if fine partitions are used, the main advantage of symbolic dynamics, to reduce the information in a signal to the essential part, is lost.

To obtain a symbolic encoding, one can directly partition the measurements into a small number of classes by defining suitable thresholds. Alternatively, a more general partitioning can be defined after a phase space reconstruction step in multidimensional space. The ergodic properties of symbol sequences are traditionally studied much more than continuous state dynamical systems. References from the dynamical systems context are, for example, the works by Herzel [135], Ebeling and coworkers [136], Schürmann and Grassberger [137], and Hao [138]. Applications are discussed in Refs. [139,140] and many others. Symbolic encoding constitutes a severe selection among the available information. This may be a desirable property in cases with high noise levels.

At a moderate level of coarse graining, any finite length scale version of the invariant quantities discussed above can in principle be used for comparative purposes. Most popular seem to be intermediate length scale estimates of the correlation dimension and of the maximal Lyapunov exponent. One example of such an estimate is the maximum likelihood estimator given by Theiler [123], see (Eq. (32)). In general, its value depends on the upper cutoff length scale \( \varepsilon \), and the embedding parameters. Another popular statistic based on the correlation integral goes back to a paper by Brock et al. [141] and is usually referred to as the BDS statistic. The authors of Refs. [141,142] make use of the fact that for a sequence of independent random numbers, \( C_m(\varepsilon) = C_1(\varepsilon)^m \) holds, where \( m \) is the embedding dimension. In these papers, also a formal test for this property is introduced. The original BDS statistic is specifically designed so that one can derive
the asymptotic distribution analytically. A simpler expression that contains the same information is \( C_m(e)/C_1(e)^m \).

A whole class of measures for nonlinearity is given by constructive measures of predictability that use a specific modelling approach to make forecasts of time series. If an estimate of the dynamics \( \hat{f} \) has been produced (see Section 3.3), one can define an average prediction error for example by

\[
e^2 = \frac{1}{N - 1 - (m - 1)r} \sum_{n=(m-1)r+1}^{N-1} (s_n + \hat{f}(s_n))^2.
\]

This, or some differently averaged error of prediction, is then used as an indicator for the unpredictability of the signal. If such an interpretation is put forth, it is essential to use some cross-validation technique to ensure that \( e \) is an out-of-sample error. An out-of-sample error is obtained if the data set that is used for the estimation of \( e \) is a different one from that used to fit \( \hat{f} \). See also Section 3.3.

One of the simplest nonlinear predictive models is the locally constant approximation given by (Eq. (27)). Since it does not involve the numerical optimisation of parameters, the danger of overfitting is rather small. Of course, the predictions obtained are often not optimal but on the other hand statistically quite stable. A similar nonparametric prediction scheme has been used by Sugihara et al. [143] and by Kennel and Isabelle [85]. Barahona and Poon [144] (among others) have used a global polynomial model (a Volterra series) for nonlinearity testing. In Ref. [145], a number of popular measures for nonlinearity are compared quantitatively for the task of discriminating noisy chaotic data from randomised surrogates with the same linear properties. The finding there was that in this particular setting at the edge of detectability, the most stable statistics outperform more subtle measures.

A number of test statistics for the detection and quantification of nonlinearity have been used in the literature which, while not explicitly called prediction errors, can be seen as specific ways to quantify nonlinear predictability in the sense used here. Among these, the test statistic proposed by Kaplan and Glass [146] is particularly suited to quantify deterministic structure in densely sampled data which permit the estimation of local flow vectors. Also the technique of false nearest neighbours advocated by Kennel et al. [121] can be regarded in this way.

Pompe [147], and Paluš [148,149] advocate the use of coarse-grained redundancies, generalisations of the time-delayed mutual information. Prichard and Theiler [58] (among others) have pointed out that it can be computationally advantageous to estimate information theoretic quantities like redundancy and mutual information by their second order generalisations. The latter can be obtained from correlation integrals, thus avoiding the common problems with box-counting approaches. Correlation integrals are also much easier to compute than the adaptive partitionings used, for example, by Fraser and Swinney [82]. The drawback of using second-order quantities is that generalised entropies lack the additivity property. The generalised mutual information is therefore no longer positive definite. This is however unproblematic as long as it is used only as a relative measure.

With measured data we will never be able to carry out the proper limit of small length scales. As a rule, the necessary coarse graining leads to a loss of invariance properties. There is one notable exception to this rule. Unstable periodic orbits embedded in a strange attractor define
a family of invariant quantities which are accessible at finite length scales. In particular, the existence, length, and stability of each orbit are such invariants. Consequently, many people have pursued the analysis of unstable periodic orbits from time series. Periodic orbit expansions [150–153] are of great theoretical appeal but they require knowledge of the dynamical system or data of exceptional quality for useful results. See Ref. [154] for a review. Similar data requirements are valid for the topological analysis of time series and the extraction of templates [155]. Recently, the emergence of methods for the stabilisation of chaotic systems [5,156] has attracted renewed interest in the detection of unstable fixed points or low order periodic orbits [157,158]. Unstable fixed points have been found and stabilised in a number of real world systems. Controllability with methods from chaos theory has been often taken as an indication of the presence of chaos in these systems. However, Christini and Collins [159] have shown that also stochastic, nonchaotic systems can be successfully controlled by such methods. If the detection and analysis of unstable fixed points is used as a means to detect nonlinearity and chaos, the same issues of significance and possible spurious results have been considered as for other quantifiers of nonlinearity.

3.6. Measures of dissimilarity

The idea to use relative measures between time series or segments of a long sequence for signal classification and nonstationarity testing has been brought up independently in a number of recent publications [62–66]. In principle, it is desirable to use relative measures that can be interpreted as a distance or a dissimilarity. As we have seen in Section 2.4, one such measure can be derived from the cross-correlation integral, (Eq. (18)). For the practical estimation of the cross-correlation integral, refer to what has been said about the correlation sum (Section 3.4.2). There are at least two other ways to construct an informal measure of dissimilarity from an estimator \( C_{XY}(\varepsilon) \) of the cross-correlation integral \( C_2(\varepsilon; \mu, \nu) \). (This notation implies that \( \mu(x) \) (resp. \( \mu(y) \)) are the probability distributions of the random variables \( X \) and \( Y \), respectively.) Kantz [73] defines an informal distance between attractors by the minimal length scale \( \varepsilon_0 \) above which the attractors are indistinguishable up to an accuracy \( \delta \):

\[
\max(|\log C_{XX}(\varepsilon) - \log C_{XY}(\varepsilon)|, |\log C_{YY}(\varepsilon) - \log C_{XY}(\varepsilon)|) < \rho \quad \forall \varepsilon > \varepsilon_0.
\]

This approach is particularly useful when comparing clean model attractors and noisy measurements. In that case the length scale at which the two attractors start to differ indicates the noise level. Another possible definition of a dissimilarity based on the cross-correlation integral is

\[
1 - C_{XY}(\varepsilon)/\sqrt{C_{XX}(\varepsilon)C_{YY}(\varepsilon)}.
\]

Further, Albano et al. [160] use a Kolmogorov–Smirnov test to detect dissimilarity of two correlation integrals.

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8 Confusingly, the purely topological properties which one would expect to be invariants in the first place, are not in general. The topological length of a cycle in a flow system depends on the choice of Poincaré section and winding numbers, etc., are only invariant under families of transformations where the family depends smoothly on its parameters. For example, the knot structure changes under reflection. Reflection is smooth but cannot be connected with the identity by a smooth family of transformations.
The different measures of dissimilarity based on predictive models (see Section 2.4) have been introduced in a more ad hoc way. They have less of a theoretical foundation and weaker invariance properties than cross-correlation integrals. In practical work however, this is often the price to be paid for statistical robustness and modest data requirements. Locally constant phase space predictors can yield stable results with a few hundred points and global polynomial or radial basis function models with even less. It should be stressed here that it is not essential for the present purpose that the predictions are optimal in the usual sense, as long as the predictions are sensitive to differences in the dynamics. In other words, for comparative purposes it may be advantageous to trade a possible bias for a lower variance.

4. Nonstationarity

Almost all methods of time series analysis, traditional linear, or nonlinear, require some kind of stationarity. Therefore, changes in the dynamics during the measurement usually constitute an undesired complication of the analysis. There are however situations where such changes represent the most interesting structure in the recording. For example, electro-encephalographic (EEG) recordings are often taken with the main purpose of identifying changes in the dynamical state of the brain. Such changes occur e.g. between different sleep stages, or between epileptic seizures and normal brain activity.

In the past, emphasis has been put on the question how stationarity can be established. If nonstationarity was detected, often the time series was discarded as unsuitable for a detailed analysis, or it was split into segments that were short enough to be regarded as stationary. More recently, authors have begun to exploit the information contained in time-variable dynamics as an essential part of the underlying process. Thus, this section will discuss tests for stationarity but also report on the steps that have been taken towards a time resolved study of nonstationary signals.

The most common definition of a stationary process found in textbooks (often called strong stationarity) is that all conditional probabilities are constant in time. Note that this definition is only applicable to the abstract generating process, and not to a realisation that produces a time series. If we regard a deterministic system as the limiting case of a stochastic process where the conditional probability density for a transition from state $x$ to state $x'$ is given by $\delta(x' - f(x))$, the definition requires $f(\cdot)$ to be unchanged with time. In the study of time series, the transition probabilities are unknown and have to be estimated from the data, subject to statistical fluctuations. In some cases, for example in intermittent systems, these fluctuations are large and the properties of measured time series can change dramatically, even though the underlying process is formally stationary after the above definition. There is no agreement on a definition of stationarity for time series. It seems reasonable to require that the duration of the measurement is long compared to the time scales of the systems. If this is the case, all temporal changes can be modelled as part of the dynamics. For this reason, processes with power law correlations are often considered nonstationary since no length of measurement could ever cover all time scales. On the other hand, processes with very well separated time scales can lead to time series which are stationary for practical purposes. The heart beat of a resting person is often homogeneous over several minutes. Longer recordings, however, cover new elements due to slower biological cycles. Since the common 24 h ECG recordings cover just a single cycle of the
circadian rhythm, they are more problematic with respect to stationarity than shorter or longer sequences.

4.1. Moving windows

A number of statistical tests for stationarity in a time series have been proposed in the literature. Most of the tests I am aware of are based on ideas similar to the following: Estimate a certain parameter using different parts of the sequence. If the observed variations are found to be significant, that is, outside the expected statistical fluctuations, the time series is regarded as nonstationary. In many applications of linear (frequency-based) time series analysis, stationarity has to be valid only up to the second moments ("weak stationarity"). Then, the obvious approach is to test for changes in quantities up to second order, like the mean, the variance, or the power spectrum. See e.g. [14] and references therein.

In a nonlinear dynamical framework, weak stationarity is not an interesting property. Quite often, the linear properties of the processes do not carry much information anyway. It is therefore desirable to use some nonlinear quantifier in order to trace nonstationarity. In particular, Isliker and Kurths [161] use a binned probability distribution. The method proposed there, however, suffers from a problem that arises with most nonlinear quantifiers. Unless quite narrow assumptions are made, the probability distribution of these quantities is not known exactly. Therefore, we cannot usually assess the significance of changes in these quantities in a rigorous way. Also, a signal might be considered stationary for some purpose, but not for another. A typical case is dimension estimation which requires stationarity in the probability of close recurrences. The authors of Ref. [161] use a $\chi^2$-test for the difference of histograms on sections of the data. This test, however, assumes that the histogram is formed by independent draws from some probability distribution. In the presence of serial correlations or deterministic structure, this is usually not justified. A possible remedy is to exclude points close in time from the histograms, thereby however losing statistical stability.

Computing nonlinear indicators for moving windows of data is attractive because it allows for a time resolved study of possible changes. As we have seen, however, there is a tradeoff between time resolution and statistical accuracy. A different way to proceed therefore is to completely give up the detailed time information and concentrate on testing the null hypothesis that the sequence is stationary. Let us remark that stationarity is an awkward concept to test for. What we would like to have is the assertion that a given time series is stationary. The failure of some test to reject the hypothesis of stationarity is not sufficient – the test might just have no power against the particular kind of nonstationarity present. Quite generally, a statistical test can never prove the null hypothesis. Thus, we would rather like to test against the null hypothesis that the data is non-stationary. Unfortunately, this is such a hopelessly composite hypothesis that we do not know how to devise a statistical test for it.

Therefore, formal tests against stationarity, like the one set up by Kennel [162] have to be understood with a particular alternative hypothesis in mind. The alternative in Ref. [162] is that the phase space geometry of the time series, reflected by the nearest-neighbour structure, is changing in time. The basic idea is that the expectation value of the number of reconstructed phase space points that have their nearest neighbour in the same half of the sequence is minimal for a stationary sequence. When thinking of geometry in phase space, nonstationarity introduces a tendency that points close in space are also close in time.
4.2. Recurrence plots

The relation between closeness in time and in phase space is the most relevant manifestation of nonstationarity in a nonlinear dynamical setting. The basic graphical tool that evaluates temporal and phase space distance of states is the classical recurrence plot of Eckmann et al. [163]. In its original version, a pair of times $n_1, n_2$ is called a recurrence if $s_{n_1}$ is one of the $k$th nearest neighbours of $s_{n_2}$, for some predefined value of $k$. An alternative given by Koebbe [164] is to define a recurrence to occur at times $n_1, n_2$ at resolution $\varepsilon$ if $n_1 \leq n_2$ and $||s_{n_1} - s_{n_2}|| \leq \varepsilon$. Usually, the $s_n$ are delay embedding vectors and the results depend on the embedding parameters. A recurrence plot is generated by marking all recurrences at a given neighbour order $k$ or resolution $\varepsilon$ in a graph with coordinates $n_1$ and $n_2$. In the second form, a recurrence plot can be simply identified with the expression

$$r_e(n_1, n_2) = (1 - \delta_{n_1, n_2}) K_e(||s_{n_1} - s_{n_2}||), \quad (35)$$

where the kernel function is usually taken to be the Heaviside step function $K_e(r) = \Theta(\varepsilon - r)$. The full recurrence structure is contained in the recurrence matrix, which is simply defined by $R_{n_1, n_2} = ||s_{n_1} - s_{n_2}||$. Obviously, $R$ is invariant under isometries (translations, rotations, and reflections) in phase space. McGuire and coworkers [165] give an algorithm to explicitly reconstruct an attractor up to isometries from a recurrence matrix. Of course, since an $(N \times N)$ symmetric matrix has $N(N - 1)$ independent entries, a recurrence matrix is not a very economical representation of $N$ vectors, and the requirement that the entries are distances in a space of dimension $m$ poses a strong constraint.

Since recurrence plots are rather difficult to read they have not gained much popularity beyond the admiration of the intriguing patterns they exhibit [166]. Zbilut and coworkers [167] propose different parameters for the statistical quantification of recurrence plots but they give little clues on how to interpret these numbers. Nevertheless, the recurrence plot can be a useful starting point for the analysis of nonstationary sequences if the relevant information is extracted in a suitable way. The most detailed account of these techniques has been given by Casdagli [168], where also the interrelations to other methods are discussed thoroughly. Let us use a very simple nonstationary dynamical system as an illustration in the following. Consider a one-parameter family of sawtooth maps $[0,1] \mapsto [0,1]$:

$$x_{n+1} = f_{\phi_n}(x_n) = f(x_n + \phi_n \mod 1), \quad f(x) = \begin{cases} 2x, & x < 1/2, \\ 2 - 2x, & 1/2 \leq x < 1. \end{cases} \quad (36)$$

Take a time series of length $N = 20000$ and let $\phi_n$ vary with time such that it covers two oscillations of a damped sine function within the measurement period: $\phi_n = (1 + e^{-n/N} \sin 4\pi n/N)/2$. Recurrence plots of (a) the time series for two-dimensional embeddings, $\varepsilon = 0.1$ and 0.01, for a three-dimensional embedding, $\varepsilon = 0.005$, and (b) for the parameter $\phi_n$ at $\varepsilon = 0.001$ are shown in Fig. 6. Casdagli [168] has pointed out that for a faithful embedding and in the limit $\varepsilon \to 0$, $N \to \infty$, the recurrence plot of the time series from a system approaches that of the varying parameter. For the above example, this can be verified from Fig. 6.
Fig. 6. Recurrence plots of a tent map time series subject to a parameter drift (see text for details). The left panel was obtained with two-dimensional embeddings. Above the diagonal: \( \varepsilon = 0.01 \) (50% of all points chosen at random for better contrast). Below the diagonal: \( \varepsilon = 0.1 \) (0.2% shown). In the right panel, embedding was done in three dimensions. Above the diagonal: \( \varepsilon = 0.005 \). Below the diagonal, the recurrence plot of the parameter sequence \( \{ \phi_n \} \) is shown with \( \varepsilon = 0.001 \) (2% shown). For decreasing \( \varepsilon \), the recurrence plots of the signal indeed converge to that of the parameter sequence (right panel, below the diagonal).

Let us note that the total average \( \sum_{n_1,n_2} r_\varepsilon(n_1,n_2) \) equals the sample correlation integral \( C(\varepsilon) \) in Eq. (30). However, in the practical estimation of \( d_2 \) one has to exclude terms with \( |n_2 - n_1| < t_{\min} \). One way to estimate the correlation time \( t_{\min} \) is the following. The partial average

\[
C(\varepsilon, \Delta n) = \frac{1}{N - \Delta n} \sum_{n=\Delta n + 1}^N r_\varepsilon(n,n - \Delta n)
\]

yields the space–time separation plot introduced by Provenzale et al. [120]. Contour lines of \( C(\varepsilon, \Delta n) \) should not increase with \( \Delta n \) except for possible oscillatory variation. The minimal \( \Delta n \) for which this is the case yields a guideline for the minimal time separation \( t_{\min} \) to be used in the correlation sum. Due to the temporal averaging, the space–time separation plot does not allow for a time-resolved analysis. Also, the effect of nonstationarity may average out in certain cases, as for example for the oscillating parameter in the tent map example. Let us remark that for the same reason also the time averaged statistic used by Kennel [162] fails to reject the null hypothesis of stationarity since the nearest neighbour of each point can have any temporal distance with about equal probability.

4.3. Tracing parameter variation

A useful way to formulate nonstationary dynamics is by introducing a temporal variation of dynamical parameters into the system, as it was already done in the tent map example. If this variation is sufficiently slow, recurrence plots and similar techniques can assess these changes to some extent. In Ref. [168], it is shown by a scaling argument that for a dynamical system with time varying parameters, the recurrence plot in the limit of small \( \varepsilon \), large \( N \), and sufficient \( m \) approaches the recurrence plot of the fluctuating parameter. This can be seen in Fig. 6 for the time varying tent map. However, it is in general difficult to extract the time variation of the parameter from its
where the prediction \( s_{n_1+i} \) is given by the average over an \( \varepsilon \)-neighbourhood,

\[
\hat{s}_{n_1+i} = \begin{cases} 
\frac{1}{w} \sum_{j=0}^{w} r_s(n_1 + i, n_2 + j) s_{n_2+j+1} & \varepsilon > 0, \\
\frac{1}{w} \sum_{j=1}^{w} s_{n_2+j+1} & \varepsilon = 0,
\end{cases}
\]

where \( \varepsilon = \sum_{j=1}^{w} r_s(n_1 + i, n_2 + j) \) is the number of neighbours of \( s_{n_1+i} \) closer than \( \varepsilon \). Fig. 7 shows \( \gamma(\varepsilon, w, n_1, n_2) \) for the modulated tent map series with \( w = 1000 \) and \( n_1, n_2 \) in steps of 500. The variation of the prediction error with the segment location \( n_1, n_2 \) is clearly visible. The information contained in Fig. 7 can be processed in a way that makes the time dependence of the parameter more clear. The quantities \( \gamma(\varepsilon, w, n_1, n_2) \) can be regarded as a dissimilarity matrix and treated by a cluster algorithm. (This technique will be discussed in more detail below in Section 7.) If the analysis is successful, the clusters are localised in parameter space and can be used to define coordinates in that space. The time-varying “distances” of each segment \( i \) to the clusters \( \nu \) \((D_i)^\nu \), to be defined in Eq. (50) of Section 7.2 below) then reflect the time variation of the parameter(s). A successful example with two clusters and one parameter is shown in Fig. 8.
Fig. 7. Cross-prediction errors for the model (40) obtained with locally constant prediction in one dimension versus location in time of segments. Local neighbourhoods were formed with a radius of one-quarter of the variance of the sequence, $\varepsilon = 0.072$. Time windows of length $w = 1000$ overlapping by 500 time steps were used.

Fig. 8. Left panel: Time dependence of the parameter $\phi_n$ in the tent map example. Right panel: The information contained in Fig. 7 was used to form two clusters of similar time series. For each segment, the distances $D^{(1)}_i, D^{(2)}_i$ to clusters (1) and (2) are computed. The difference $D^{(2)}_i - D^{(1)}_i$ is plotted for each segment versus time. Such a plot cannot only reveal that there is a single changing parameter but also the form of its change. That the units in both panels are of comparable magnitude is purely coincidental. Clustering of time series is discussed in Section 7.2 below, where also $D^{(l)}_i$ is defined.

4.4. An application

Let us finish the discussion of nonstationary time series with an application that is currently studied with considerable effort in a number of research groups: the anticipation of epileptic seizure onset from intracranial recordings of neural potentials. Epileptic seizures manifest themselves in specific patterns in the neural electric field. While traditional electro-encephalograms (EEG) with electrodes placed on the surface of the scalp show such patterns when the epilepsy activity has reached a cortical region that is sufficiently close to the surface, for a detailed study of focal epilepsy in deeper regions of the brain electrodes have to be implanted in the epileptogenic region. This is a common clinical technique in pre-surgical screening. The specific activity during seizures is usually so pronounced that it can be detected visually and also automatically.\footnote{Many authors have discussed the question if there is evidence for low-dimensional chaos and strange attractors in normal, or, more likely, in epilepsy EEG data. References include [169–172] and many more, conclusions are controversial. More recently, the focus has shifted from the question of chaos versus noise to the quantification of changes in the EEG. Most authors agree that ictal (during seizures) and inter-ictal EEG can be distinguished with nonlinear, but also with spectral methods. References for the former are, for example, Refs. [65,173–176], and Ref. [177] for the latter.} A far more
challenging problem is to detect specific changes in the dynamics of the recordings just prior to the actual seizure. First of all, a reliable anticipation of a seizure several minutes ahead potentially allows for pharmacological or electrophysiological intervention. The insights into the mechanism that leads to the large scale pathological activity are of equally high interest. The problem however is very intricate. At any given time, the recorded neural activity is very rich and far from being understood. Although usually simultaneous recordings at several positions are available, it is not clear to what extent multivariate studies provide more insight at this stage [178]. Electrode spacings down to fractions of a millimeter are still much larger than typical coherence lengths. The dynamics is time variable and only part of this variability is specific for the generation of epileptic activity. It is not expected that the brain falls into a single typical state prior to a seizure but rather that the pre-seizure activity shows some characteristic yet variable behaviour. The task of time series analysis is to find and specify such features that allow for the detection of the critical state.

Elger and Lehnertz [131] claim statistically significant positive evidence for seizure predictability several minutes ahead of seizure onset. The authors use a sliding window version of the correlation integral that has been customised for this particular purpose. Since the intracranial EEG signal is nonstationary even in episodes without epileptic activity, the window length should be short enough for the segment to be effectively stationary but long enough to yield stable results. Half-overlapping windows of 30 s duration were chosen at a sampling rate of 173 Hz. The variance of the data segments is time dependent which makes it difficult to choose a length scale for the determination of an effective scaling index. The dominating contribution to the variance found in pathological regions often arises from spikes occurring at irregular intervals. Although these spikes are characteristic for epileptogenic tissue, they seem not to be specific precursors of seizures and should therefore not be overemphasised in the analysis. In Ref. [131] this is achieved by selecting approximate scaling regions for each individual window. These are usually found at much smaller scales than the spikes. Pre-seizure behaviour is found to be accompanied by epochs of smaller effective scaling indices as compared to the standard behaviour found in segments that are spatially and/or temporally well separated from the seizure.

Correlation dimension data for a time series recorded by Elger and Lehnertz has already been shown in Section 3.4.2 and it may surprise the reader that in Ref. [131] the authors do find small approximate scaling regions. This finding can be reproduced by limiting the correction for the dynamical correlations to the exclusion of pairs which are not more than three sampling intervals apart in time. For the determination of a proper dimension, or for the interpretation of the effect as a signature of a finite attractor dimension, this would be disastrous, but it is fully justified by the resulting discriminative power for the particular purpose at hand. Fig. 9 repeats the same as was done in Fig. 5 but with the limited correction indicated above. The left panel shows a data segment that was taken far away in time from any seizure. The right panel was calculated from the same data as Fig. 5, a segment taken about 10 min prior to the onset of an epileptic seizure. Indeed we see tiny “plateaus” not present in Fig. 5.\footnote{The particular shape of the deflection that could have been interpreted as a plateau resembles what is typically seen for intermittent systems, see [179].} Further, we see that there is a difference between the two segments in where the pseudo-scaling is found. It is such differences that have been studied systematically in Ref. [131]. The statistical material presented there is based on 16 patients and shows
Fig. 9. Maximum likelihood estimator of the correlation dimension for an intracranial EEG recording [131] of an epilepsy patient. The effect of dynamically correlated pairs was only incompletely corrected. Left: a data set measured long before the next seizure. Right: a data set measured about 10 min prior to seizure onset. The right panel was produced with the same time series as Fig. 5.

that the pre-seizure states and normal epochs follow significantly different distributions of approximate scaling indices. Certainly, a number of ad hoc decisions have been made in devising the algorithm and it is not fully clear from Ref. [131] to what extent the same sample of patients has been used to optimise parameters. Further research will have to evaluate whether the differences are strong enough to make reliable out-of-sample predictions for individual patients. For clinical applicability, it will eventually be necessary to compute and interpret the relevant quantities in real time.

The finding that the discriminative power declines upon full correction of the dynamical correlations indicates that it is not exactly phase space geometry that distinguishes the different states. It should however be stressed that the shown segments are not distinguishable by their autocorrelation functions or spectra. Thus, what is represented in Fig. 9 is a difference in the nonlinear dynamical correlation. This suggests that there may be nonlinear indicators which are more sensitive to changes in the dynamics and that might correlate even more strongly with the seizure onset.

5. Testing for nonlinearity

There are two distinct motivations to use a nonlinear approach when analysing time series data. It might be that the arsenal of linear methods has been exploited thoroughly but all the efforts left certain structures in the time series unaccounted for. It is also common that a system is known to include nonlinear components and therefore a linear description seems unsatisfactory in the first place. Such an argument is often heard for example in brain research – nobody expects the brain to be a linear device. In fact, there is ample evidence for nonlinearity in particular in small assemblies of neurons. Nevertheless, the latter reasoning is rather dangerous. That a system is known to contain nonlinear components does not prove that this nonlinearity is also reflected in a specific signal we measure from that system. In particular, we do not know if it is of any practical use to go beyond the linear approximation. After all, we do not want our data analysis to reflect our prejudice about the underlying system but to represent a fair account of the structures that are present in the data. Consequently, the application of nonlinear time series methods has to be justified by establishing nonlinearity in the time series data.
This section will discuss formal statistical tests for nonlinearity. First, a suitable null hypothesis for the underlying process will be formulated covering all Gaussian linear processes or a class that is somewhat wider. We will then attempt to reject this null hypothesis by comparing the value of a nonlinear parameter estimated on the data with its probability distribution for the hypothesis. Since only exceptional cases allow for the exact or asymptotic derivation of this distribution unless strong additional assumptions are made, we have to estimate it by a Monte Carlo resampling technique. This procedure is known in the nonlinear time series literature as the method of surrogate data, see Refs. [104,180,181]. Thus we have to face a two-fold task. We have to find a nonlinear parameter that is able to actually detect an existing deviation of the data from a given null hypothesis and we have to provide an ensemble of randomised time series that accurately represents the null hypothesis.

5.1. Detecting weak nonlinearity

In the preceding sections, several quantities have been discussed that can be used to characterise nonlinear time series. For the purpose of nonlinearity testing we need such quantities that are particular powerful in discriminating linear dynamics and weakly nonlinear signatures – strong nonlinearity is usually more easily detectable. Quite a number of such measures has been proposed and used in the literature. An important objective criterion that can be used to guide the preferred choice is the discrimination power of the resulting test. The power $\beta$ is defined as the probability that the null hypothesis is rejected when it is indeed false. It will obviously depend on how and how strongly the data actually deviates from the null hypothesis.

Traditional measures of nonlinearity are derived from generalisations of the two-point autocovariance function or the power spectrum. The use of higher-order cumulants and bi- and multi-spectra is discussed for example in Ref. [182]. One particularly useful third-order quantity is

$$
\phi^{rev}(\tau) = \frac{\sum_{n=1}^{N}(s_n - \bar{s}) s_n \sim s_{n-\tau})^3}{[\sum_{n=1}^{N}(s_n - \bar{s})^2]^{3/2}},
$$

(41)

since it measures the asymmetry of a series under time reversal. (Remember that the statistics of linear stochastic processes is always symmetric under time reversal. This can be most easily seen when the statistical properties are given by the power spectrum which contains no information about the direction of time.) Time reversibility as a criterion for discriminating time series is discussed in detail in Ref. [183].

When a nonlinearity test is performed with the question in mind if nonlinear deterministic modelling of the signal may be useful, it seems most appropriate to use a test statistic that is related to a nonlinear deterministic approach. We have to keep in mind however that a positive test result only indicates nonlinearity, not necessarily determinism. Since nonlinearity tests are usually performed on data sets which do not show unambiguous signatures of low-dimensional determinism (like clear scaling over several orders of magnitude), one cannot simply estimate one of the quantitative indicators of chaos, like dimension or Lyapunov exponent. The formal answer would almost always be that both are probably infinite. Still, some useful test statistics are at least inspired by these quantities. Usually, some effective value at a finite length scale has to be computed rather than attempting to take the proper limits. We can largely follow the discussion in
Section 3.5, considering the limiting case that the deterministic signature to be detected is probably weak. In that case the major limiting factor for the performance of a statistical indicator is its variance since possible differences between two samples may be hidden among the statistical fluctuations. In Ref. [145], a number of popular measures of nonlinearity are compared quantitatively. The results can be summarised by stating that in the presence of time-reversal asymmetry, the three-point autocorrelation (Eq. (41)) gives very reliable results. However, many non-linear evolution equations produce little or no time-reversal asymmetry in the statistical properties of the signal. In these cases, simple measures like a prediction error of a locally constant phase space predictor performed best. It was found to be advantageous to choose embedding and other parameters in order to obtain a quantity that has a small spread of values for different realisations of the same process, even if at these parameters no valid embedding could be expected.

5.2. Surrogate data tests

All of the measures of nonlinearity discussed above have in common that their probability distribution on finite data sets is not known analytically. Some authors have tried to give error bars for measures like predictabilities (e.g. see [144]) or averages of pointwise dimensions (e.g. see [184]) based on the observation that these quantities are averages (or medians) of many individual terms, in which case the variance (or quartile points) of the individual values yield an error estimate. This reasoning is however only valid if the individual terms are independent, which is usually not the case for time series data. In fact, it is found empirically that nonlinearity measures often do not even follow a Gaussian distribution. Also the standard error given by Roulston [185] for the mutual information is not quite correct except for uniformly distributed data. While a smooth rescaling to uniformity would not do harm to his derivation, rescaling to exact uniformity is in general non-smooth and introduces a bias in the joint probabilities. In order to determine the distribution of a nonlinear statistic on realisations of the null hypothesis, it is therefore preferable to use a Monte Carlo resampling technique. Traditional bootstrap methods use explicit model equations that have to be extracted from the data. This typical realizations approach can be very powerful for the computation of confidence intervals, provided the model equations can be extracted successfully. As discussed by Theiler and Prichard [186], the alternative approach of constrained realizations is more suitable for the purpose of hypothesis testing we are interested in here. It avoids the fitting of model equations by directly imposing the desired structures onto the randomised time series. However, the choice of possible null hypothesis is limited by the difficulty to impose arbitrary structures on otherwise random sequences. The following section will discuss a number of null hypotheses and algorithms to provide the adequately constrained realisations. The most general method to generate constrained randomisations of time series is described in Ref. [187]. The price for its accuracy and generality is its high computational cost.

5.2.1. How to make surrogate data

It is essential for the validity of the statistical test that the surrogate series are created properly. If they contain spurious differences to the measured data, these may be detected by the test and interpreted as signatures of nonlinearity. More formally, the size of a test is the actual probability that the null hypothesis is rejected although it is in fact true. For a valid test, the size $\alpha$ must not exceed the level of significance $p$. The correct size crucially depends on the way the surrogates are
generated. Let us discuss a hierarchy of null hypotheses and the issues that arise when creating the corresponding surrogate data.

A simple case is the null hypothesis that the data consists of independent draws from a fixed probability distribution. Surrogate time series can be simply obtained by randomly shuffling the measured data. If we find significantly different serial correlations in the data and the shuffles, we can reject the hypothesis of independence.\(^\text{11}\) The next step would be to explain the structures found by linear two-point autocorrelations. A corresponding null hypothesis is that the data have been generated by some linear stochastic process with Gaussian increments. The most general univariate linear process is given by Eq. (3). The statistical test is complicated by the fact that we do not want to test against one particular linear process only (one specific choice of the \(a_i\) and \(b_i\)), but against a whole class of processes. This is called a \textit{composite} null hypothesis. The unknown values \(a_i\) and \(b_i\) are sometimes referred to as \textit{nuissance parameters}. There are basically three directions we can take in this situation. First, we could try to make the discriminating statistic independent of the nuissance parameters. This approach has not been demonstrated to be viable for any but some very simple statistics. Second, we could determine which linear model is most likely realised in the data by a fit for the coefficients \(a_i\) and \(b_i\), and then test against the hypothesis that the data has been generated by this particular model. Surrogates are simply created by running the fitted model. This \textit{typical realisations} approach is the common choice in the bootstrap literature, see e.g. the classical book by Efron [188]. The main drawback is that we cannot recover the \textit{true} underlying process by any fit procedure. Apart from problems associated with the choice of the correct model orders \(M\) and \(N\), the data is by construction a very likely realisation of the fitted process. Other realisations will fluctuate \textit{around} the data which induces a bias against the rejection of the null hypothesis. This issue is discussed thoroughly in Ref. [104], where also a calibration scheme is proposed.

The most attractive approach to testing for a composite null hypothesis seems to be to create \textit{constrained realisations} [186]. Here it is useful to think of the measurable properties of the time series rather than its underlying model equations. The null hypothesis of an underlying Gaussian linear stochastic process can also be formulated by stating that all structure to be found in a time series is exhausted by computing first- and second-order quantities, the mean, the variance and the autocovariance function. This means that a randomised sample can be obtained by creating sequences with the same second-order properties as the measured data, but which are otherwise random. When the linear properties are specified by the squared amplitudes of the Fourier transform (that is, the periodogram estimator of the power spectrum), surrogate time series \(\{s'_n\}\) are readily created by multiplying the Fourier transform of the data by random phases and then transforming back to the time domain:

\[
s'_n = \sum_{k=0}^{N-1} e^{i\alpha_k} \sqrt{P_k} e^{-i2\pi kn/N},
\]

where \(0 \leq \alpha_k < 2\pi\) are independent uniform random numbers.

\(^{11}\) Independence seems not to be an interesting null hypothesis for most time series problems. It becomes relevant when the residual errors of a time series model is evaluated. For example in the BDS test for nonlinearity [141], an ARMA model is fitted to the data. If the data are linear, then the residuals are expected to be independent.
The two null hypotheses discussed so far (independent random numbers and Gaussian linear processes) are not what we want to test against in most realistic situations. In particular, the most obvious deviation from the Gaussian linear process is usually that the data do not follow a Gaussian single time probability distribution. This is quite obvious for data obtained by measuring intervals between events, e.g. heart beats since intervals are strictly positive. There is however a simple generalisation of the null hypothesis that explains deviations from the normal distribution by the action of a monotone, static measurement function:

\[ s_n = s(x_n), \quad x_n = \sum_{i=1}^{M} a_i x_{n-i} + \sum_{i=0}^{N} b_i \eta_{n-i}. \]  

(43)

We want to regard a time series from such a process as essentially linear since the only nonlinearity is contained in the – in principle invertible – measurement function \( s(\cdot) \).

The most common method to create surrogate data sets for this null hypothesis essentially attempts to invert \( s(\cdot) \) by rescaling the time series \( \{s_n\} \) to conform with a Gaussian distribution. The rescaled version is then phase randomised (conserving Gaussianity on average) and the result is rescaled to the empirical distribution of \( \{s_n\} \). These amplitude-adjusted Fourier transformed surrogates (AAFT) yield a correct test when \( N \) is large, the correlation in the data is not too strong and \( s(\cdot) \) is close to the identity. It is argued in Ref. [189] that for short and strongly correlated sequences, the AAFT algorithm can yield an incorrect test since it introduces a bias towards a slightly flatter spectrum. In fact, the formal requirement the surrogates have to fulfill for the test to be correct is that they have the same sample periodogram and the same single time sample probability distribution as the data. Schreiber and Schmitz [189] propose a method which iteratively corrects deviations in spectrum and distribution. Alternatingly, the surrogate is filtered towards the correct Fourier amplitudes and rank-ordered to the correct distribution. The accuracy that can be reached depends on the size and structure of the data and is generally sufficient for hypothesis testing.

The above schemes are all based on the Fourier amplitudes of the data, which is however not exactly what we want. Remember that the autocorrelation structure given by Eq. (6) corresponds to the Fourier amplitudes only if the time series is one period of a sequence that repeats itself every \( N \) time steps. Conserving the Fourier amplitudes of the data means that the periodic autocovariance function

\[ C_p(\tau) = \frac{1}{N} \sum_{n=1}^{N} x_n x_{n(\tau-1) \mod N + 1} \]  

(44)

is reproduced, rather than \( C(\tau) \). The difference can lead to serious artefacts in the surrogates, and, consequently, spurious rejections in a test. In particular, any mismatch between the beginning and the end of a time series poses problems, as discussed e.g. in Ref. [190]. In spectral estimation, problems caused by edge effects are dealt with by windowing and zero padding. None of these have been successfully implemented for the phase randomisation of surrogates. The problem of non-matching ends can be partly overcome by choosing a subinterval of the recording such that the end points do match approximately. Still, there may remain a finite phase slip at the matching points. The only method that has been proposed so far that strictly implements \( C(\tau) \) rather than \( C_p(\tau) \) is given in Ref. [187]. The method is very accurate but also rather costly in terms of computer time. In
practical situations, the matching of end points is a simple and mostly sufficient precaution that should not be neglected.

Since the randomisation algorithm of Ref. [187] is of very general applicability and conceptually quite simple, let us give a brief description. In order to create randomised sequences with the correct distribution of values, only permutations of the original time series are considered. The shuffling is however carried out under the constraint that the autocovariances of the surrogate $C'(\tau)$ are the same as those of the data, $C(\tau)$. This is done by specifying the discrepancy as a cost function, e.g.

$$E^{(q)} = \left[ \sum_{\tau=0}^{N-1} |C'(\tau) - C(\tau)|^q \right]^{1/q}.$$  \hspace{1cm} (45)

The way the average over all lags $\tau$ is taken can be influenced by the choice of $q$. Now $E^{(q)}(\{s_n\})$ is minimised among all permutations $\{\tilde{s}_n\}$ of the original time series $\{s_n\}$ using the method of simulated annealing. Configurations are updated by exchanging pairs in $\{\tilde{s}_n\}$. With an appropriate cooling scheme, the annealing procedure can reach any desired accuracy. Simulated annealing has a rich literature, classical references are Metropolis et al. [191] and Kirkpatrick [192], more recent material can be found for example in [193].

Constrained randomisation using combinatorial minimisation is a very flexible method since in principle arbitrary constraints can be realised. Although it is seldom possible to specify a formal null hypothesis for more general constraints, it can be quite useful to be able to incorporate into the surrogates any feature of the data that is understood already or that is uninteresting. Let us give an example for the flexibility of the approach, a simultaneous recording of the breath rate and the instantaneous heart rate of a human subject during sleep. (Data set B of the Santa Fe Institute time series contest in 1991 [194], sample points 1800–4350.) Regarding the heart rate recording on its own, one easily detects nonlinearity, in particular via an asymmetry under time reversal. An interesting question however is, how much of this structure can be explained by linear dependence on the breath rate, the breath rate also being non-time-reversible. In order to answer this question, we need to make surrogates that have the same autocorrelation structure but also the same cross-correlation with respect to the fixed input signal, the breath rate. Accordingly, the constraint is imposed that lags 0, ..., 500 of the autocovariance function and lags $-500, ..., 500$ of the cross-covariance function with the reference (breath) signal are given by the data. Further suppose that within the 20 min of observation, during one minute the equipment spuriously recorded a constant value. In order not to interpret this artefact as structure, the same artefact is imposed on the surrogates, simply by excluding these data points from the permutation scheme.

Fig. 10 shows the measured breath rate (upper trace) and instantaneous heart rate (middle trace). The lower trace shows a surrogate conserving both, auto- and cross-correlations. The visual

\footnote{Strictly speaking, these constraints overspecify the problem and it is likely that the only permutation that fulfills them exactly is the original time series itself. However, it can be expected that there are a large number of permutations which are essentially different and which fulfill the constraint \textit{almost} exactly. In fact, it has been observed [195] for very short sequences of $N < 50$ points and strong correlations that the annealing scheme settled on the original data. If this seems to happen, one can introduce a term in the cost function that discourages similarity to the original permutation.}
impression from Fig. 10 is that while the linear cross-correlation with the breath rate explains the cyclic structure of the heart rate data, other features, in particular the asymmetry under time reversal, remain unexplained. This finding can be verified at the 95% level of significance, using the time asymmetry statistic given in Eq. (41). Possible explanations include artefacts due to the peculiar way of deriving heart rate from inter-beat intervals, nonlinear coupling to the breath activity, nonlinearity in the cardiac system, and others.

Let us finish the section by giving a more exotic example, from finance. The time series\(^{13}\) consists of 1500 daily returns (until the end of 1996) of the BUND Future, a derived german financial instrument. As can be seen in the upper panel of Fig. 11, the sequence is nonstationary in the sense that the local variance and also the local mean undergo changes on a time scale that is long compared to the fluctuations of the series itself. This property is known in the statistical

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\(^{13}\) The data were kindly provided by Thomas Schürmann, WGZ-Bank Düsseldorf.
literature as *heteroscedasticity* and modeled by the so-called GARCH [196] and related models. Here, we want to avoid the construction of an explicit model from the data but rather ask the question if the data is compatible with the null hypothesis of a correlated linear stochastic process with time-dependent local mean and variance. We can answer this question in a statistical sense by creating surrogate time series that show the same linear correlations and the same time dependence of the running mean and running variance as the data and comparing a nonlinear statistic between data and surrogates. The lower panel in Fig. 11 shows a surrogate time series generated using the annealing method described above.

The cost function was set up to match the autocorrelation function up to five days and the moving mean and variance in sliding windows of 100 days duration. In Fig. 11 the running mean and variance are shown as points and error bars, respectively, in the middle trace. The deviation of these between data and surrogate has been minimised to such a degree that it can no longer be resolved. A comparison of the time-asymmetry statistic Eq. (41) for the data and 19 surrogates did not reveal any discrepancy, and the null hypothesis could not be rejected.

5.3. What can be learned

Having set up all the ingredients for a statistical hypothesis test of nonlinearity, we may ask what we can learn from the outcome of such a test. The formal answer is of course that we have, or have not, rejected a specific hypothesis at a given level of significance. How interesting this information is, however, depends on the null hypothesis we have chosen. The test is most meaningful if the null hypothesis is plausible enough so that we are prepared to believe it in the lack of evidence against it. If this is not the case, we may be tempted to go beyond what is justified by the test in our interpretation. Take as a simple example a recording of hormone concentration in a human. We can test for the null hypothesis of a stationary Gaussian linear random process by comparing the data to phase randomised Fourier surrogates. Without any test, we know that the hypothesis cannot be true since hormone concentration, unlike Gaussian variates, is strictly nonnegative. If we failed to reject the null hypothesis by a statistical argument, we will therefore go ahead and reject it anyway by common sense, and the test was pointless. If we did reject the null hypothesis by finding a coarse grained “dimension” which is significantly lower in the data than in the surrogates, the result formally does not give any new information but we might be tempted to speculate on the possible interpretation of the “nonlinearity” detected.

This example is maybe too obvious, it was meant only to illustrate that the hypothesis we test against is often not what we would actually accept to be true. Other, less obvious and more common, examples include signals which are known (or found by inspection) to be nonstationary (which is not covered by most null hypotheses), or signals which are likely to be the squares of some fundamental quantity. An example for the latter are the celebrated sunspot numbers. Sunspot activity is generally connected with magnetic fields and is to first approximation proportional to the squared field strength. Obviously, sunspot numbers are nonnegative, but also the null hypothesis of a monotonically rescaled Gaussian linear random process is to be rejected since taking squares is not a monotonic operation. Unfortunately, the framework of surrogate data does not currently provide a method to test against null hypothesis involving noninvertible measurement functions. Yet another example is given by linearly filtered time series. Even if the null hypothesis of a monotonically rescaled Gaussian linear random process is true for the underlying signal, it is
usually not true for filtered copies of it, in particular sequences of first differences, see [197] for a discussion of this problem. Recent efforts on the generalisation of randomisation schemes try to broaden the repertoire of null hypotheses we can test against. The hope is that we can eventually choose one that is general enough to be acceptable if we fail to reject it with the methods we have. Still, we cannot prove that there is not any structure in the data beyond what is covered by the null hypothesis. From a practical point of view, however, there is not much of a difference between structure that is not there and structure that is undetectable with our observational means.

6. Nonlinear signal processing

The goals of time series analysis are probably as diverse as the methods. In basic research, the ultimate aim is a deeper understanding of some phenomenon in nature. In engineering, clinical research, finance, etc., a better understanding of the processes is also most welcome but the actual purpose of the work is different, making better devices, making people healthier, making money. Pursuing such a goal often involves a rather specific task of time series analysis. The interesting problem of signal classification will be dealt with in Section 7. One of the most well-known objectives is the prediction of future values of some quantity, for instance the price of an entity at the stock market. The prediction problem has been discussed in Section 3.3. It is almost identical to the problem of estimating the dynamics underlying a time series. An intermediate step in most time series studies is to filter the data in order to enhance the relevant information. Prediction and filtering, or noise reduction, have many things in common, but there are notable differences. In particular, for noise reduction it is not enough to have a description of the dynamics. One also has to have a means of finding a cleaner signal that is consistent with this dynamics.

6.1. Nonlinear noise reduction

Originally, phase space methods of nonlinear noise reduction have been developed [198–201] under the premise that there is a low-dimensional dynamical system which is only observed with some observational error that is to be suppressed. Conceptual as well as technical issues arising in such a situation have been well discussed in the literature, see Kostelich and Schreiber [202] for a review containing the relevant references. In interdisciplinary applications, we usually face a different situation – the signals themselves often contain a stochastic component. Before we apply a filtering technique we therefore have to specify what exactly we want to separate. Phase space projection techniques, like those employed by Grassberger and coworkers [201], rely on the assumption that the signal of interest is approximately described by a manifold that has a lower dimension than some phase space it is embedded in. This statement can be formalised as follows. Let \( \{x_n\} \) be the states of the system at times \( n = 1, \ldots, N \), represented in some space \( \mathcal{M} \). A \((d - Q)\)-dimensional submanifold \( \mathcal{F} \) of this space can be specified by \( F_q(y) = 0 \), \( q = 1, \ldots, Q \). Even if \( F_q \) is not known exactly, or if \( \{x_n\} \) is corrupted by noise, we can always find \( \{\xi_n\} \) such that \( y_n = x_n + \xi_n \) and

\[
F_q(x_n + \xi_n) = 0, \quad \forall q, n.
\] (46)

Then \( \sqrt{\langle \xi^2 \rangle} \) denotes the (root mean squared) average error we make by approximating the points \( \{x_n\} \) by the manifold \( \mathcal{F} \). In a measurement we can only obtain noisy data \( y_n = x_n + \eta_n \), where \( \{\eta_n\} \) is
It has been found advantageous to introduce a diagonal weight matrix $R$ and define a transformed version of the covariance matrix $C_{ij}$ as $R_{ii}C_{ij}R_{jj}$ for the calculation of the principal directions. In order to penalise corrections based on the first and last coordinates in the delay window one puts $R_{00} = R_{mm} = r$ where $r$ is large. The other values on the diagonal of $R$ are 1.

For this to be valid, the neighbourhoods should be larger than the noise level. In practice, a tradeoff between the clear definition of the noise directions and a good linear approximation has to be balanced.
illustrated in Fig. 12. The correction is done for each embedding vector, resulting in a set of corrected vectors in embedding space. Since each element of the scalar time series occurs in \( m \) different embedding vectors, we finally have as many different suggested corrections, of which we simply take the average. Therefore, in embedding space the corrected vectors do not precisely lie on the local subspaces but are only moved towards them.

As an application, Fig. 13 shows the result of the noise reduction scheme applied to a noisy ECG. As discussed already in Section 3.2, a delay coordinate embedding of an electrocardiogram seems to be well approximated by a lower-dimensional manifold. This is apparent already in a two-dimensional representation (Fig. 3), but for the purpose of noise reduction, embeddings in higher dimensions are advantageous. The data shown in Fig. 13 was produced with delay windows covering 200 ms, that is, \( m = 50 \) at a delay time of 4 ms (equal to the sampling interval). See Ref. [206] for more details on the nonlinear projective filtering of ECG signals. Applications of
nonlinear noise reduction to chaotic laboratory data are given in Ref. [207]. It should be noted that, as it stands, nonlinear noise reduction is quite computer time intensive, in particular if compared to Fourier-based filters. For small and moderate noise levels, this can be moderated by using fast neighbour search strategies, see for example Ref. [208] for a review set in the context of time series analysis. Recently, a fast version of nonlinear projective noise reduction has been developed [209] that cannot only be used a posteriori but also for real time processing in a data stream.

If the data quality does not permit to use the local linear approach, one can try to use locally constant approximations instead [211]. This is done exactly in the same way as for locally constant predictions, see Section 3.3, Eq. (27). The only difference is that instead of predicting a future value with $\Delta n > 0$, the middle coordinate of the embedding window $\Delta n = m/2$ is estimated.

6.2. Signal separation

Noise reduction can be regarded as the particular case of the more general task of signal separation where one of the signals is the noise contribution. It turns out that the methodology developed for noise reduction can be generalised to the separation of other types of signals. As a specific example, let us discuss the extraction of the fetal electrocardiogram (FECG) from noninvasive maternal recordings. Other very similar applications include the removal of ECG artefacts from electro-myogram (EMG) recordings (electric potentials of muscle) and spike detection in electro-encephalogram (EEG) data [212].

Fetal ECG extraction can be regarded as a three-way filtering problem since we have to assume that a maternal abdominal ECG recording consists of three main components, the maternal ECG, the fetal ECG, and exogenous noise, mostly from action potentials of intervening muscle tissue. All three components have quite similar broad-band power spectra and cannot be filtered apart by spectral methods. The fetal component is detectable from as early as the 11th week of pregnancy. After about the 20th week, the signal becomes weaker since the electric potential of the fetal heart is shielded by the vernix caseosa forming on the skin of the fetus. It appears again towards delivery. In Refs. [213,214], it has been proposed to use a nonlinear phase space projection technique for the separation of the fetal signal from maternal and noise artefacts. A typical example of output of this procedure is shown in Fig. 14. The assumption made about the nature of the data is that the maternal signal is well approximated by a low-dimensional manifold in delay reconstruction space. After projection onto this manifold, the maternal signal is separated from the noisy fetal component. Now it is assumed that the fetal ECG is also approximated by a low-dimensional manifold and the noise is removed by projection. Since both manifolds are curved, the projections have to be made onto linear approximations. For technical details see Refs. [213,214].

7. Comparison and classification

With current methods, many real-world systems cannot be fully understood on the basis of time series measurements. Approaches aiming at an absolute analysis, like the reliable determination of the fractal dimension of a strange attractor have often been found to fail for various reasons, the most prominent being that most of the systems are not low-dimensional deterministic. However, many phenomena can still be studied in a comparative way. In that case, we do not have to worry
too much about the theoretical basis of the quantities we use. The results are validated by the statistical significance for the discriminative power. The classification of states can give valuable insights into the structure of a problem, and very often, signal classification is desirable in its own right. In clinical applications, for example, it is common to define quantities by a standardised procedure, even if this procedure yields an observable which has no immediate physical interpretation. Take the standard procedure of determining the blood pressure non-invasively (the Riva-Rocci method). Although the measurement is indirect and does not yield invariant results, the standardisation of the procedure ensures good comparability of the results. For its value as a diagnostic tool, it is irrelevant whether the measured numbers actually represent the pressure of the blood in a specific part of the body or not.

Comparison and classification of time series is most often done pretty much in the same spirit as, for example, the blood pressure measurement. A complex phenomenon is reduced in a well defined way to a single number or a small set of numbers. Further analysis is then carried out on these numbers. In nonlinear time series analysis these numbers can be for example nonlinear prediction errors, coarse grained dimensions or entropies, etc. Below, an alternative approach will be discussed which attempts to carry out the actual comparison between the signals directly rather than between single numbers abstracted from them.

Before we set up a classification problem, we have to decide how we want to use the available information. We need to make economical use of the data we have since we need them to set up and maybe optimise a classification scheme and then to independently verify it. One way to proceed is by splitting the available data base into two parts, a learning set where the correct classification is known (for example, which of the patients were in the control group), and a test set which is only used at the very end of the analysis to verify the validity of the classification without using the correct answer. (See the discussion of the overfitting problem in Section 3.3.) The advantage of this approach is that the training phase can be supervised and directed to the desired behaviour. The disadvantage is that we need a sufficient number of test cases to be kept apart. Another possibility is to perform unsupervised classification of the whole database. One tries to find whether the whole ensemble of signals falls into distinct groups naturally without help from a supervisor who knows
the correct answer. One can then regard the whole available data base as the test set for the correctness and significance of the grouping. The latter approach of unsupervised learning will be considered here mostly. Note that in both cases, supervised and unsupervised classification, the test set cannot be used repeatedly to optimise strategies or parameters, unless claims of significance are modified accordingly.

7.1. Classification by histograms

The standard approach to classification of time series is to express the information in each sequence by a single number or a few of them. One can then form a histogram of these numbers, either in one or a few dimensions. Let us give a simple example. Consider a generalised baker map

\[ v_n \leq x: u_{n+1} = \beta u_n, \quad v_{n+1} = v_n/x, \]
\[ v_n > x: u_{n+1} = 0.5 + \beta u_n, \quad v_{n+1} = (v_n - x)/(1 - x) \]

with \( x = 0.4 \). The parameter \( \beta \) can be varied without changing the positive Lyapunov exponent \[28\]. Two groups of sequences were generated, the first (group a) containing 50 sets with \( \beta = 0.6 \) and the second (group b) containing 50 with \( \beta = 0.8 \). Each sequence has a length of 400 points. Prediction errors \( \gamma_i \) are calculated from Eqs. (33) and (27) and collected in a histogram (see Fig. 15). For unsupervised learning, this histogram would not provide enough information since the observed distribution (black bars) does not fall into two groups naturally. In fact, in the best case we can find a threshold value (arrow) that minimises the number of misclassifications. Still at least 20 series will be assigned to the wrong group since the individual distributions (white and grey bars) overlap.

7.2. Classification by clustering

The success of the usual classification schemes based on histograms or scatter plots of a few quantities crucially depends on the right choice of observable to characterise the signals. In general, it seems quite a loss of information to express nonlinear dynamics by a few numbers. An alternative

![Fig. 15. Classification using a histogram. The arrow indicates the threshold value resulting in the fewest misclassifications possible.](image)
is to compare the individual series directly without first extracting an observable. As we will see in the baker map example of the preceding section, two series produced with different parameters may be well distinguishable even though they have comparable predictability. This motivates to generalise the usual measures of nonlinearity discussed in Sections 2.3, 3.4 and 3.5 to comparative measures, or measures of similarity, as it was done in Sections 2.4 and 3.6. If we want to compare \( w \) signals, the study of symmetric dissimilarities yields \( w(w+1)/2 \) independent relative quantities \( \gamma_{i,j} \) rather than just \( w \) characteristics \( \gamma_i = \gamma_{i,t} \). In this section it will be shown how such matrices can be obtained and used for the task of classification. The method has been proposed in Ref. [63] where also more technical details and further examples can be found.

The main idea is to use a cluster algorithm to find groups of data based on a dissimilarity matrix. There are several standard methods to do so [232] and the choice made below is not meant to be exclusive by any means. The task now is to classify \( w \) objects into \( K \) groups or clusters. Let us define a membership index \( u_i^{(v)} \) to be 1 if object \( i \) is in cluster \( v \), and 0 if not. A cluster is given by all points with membership index 1: \( C^{(v)} = \{ i : u_i^{(v)} = 1 \} \). The size of a cluster is \( |C^{(v)}| = \sum_{i=1}^{w} u_i^{(v)} \). The average dissimilarity of object \( i \) to cluster \( v \) (the “distance” of \( i \) to \( v \)) is then given by

\[
D_i^{(v)} = \frac{1}{|C^{(v)}|} \sum_{j=1}^{w} u_j^{(v)} \gamma_{ij} .
\]

The average dissimilarity within cluster \( v \) is

\[
D^{(v)} = \frac{1}{|C^{(v)}|} \sum_{i=1}^{w} u_i^{(v)} D_i^{(v)}
\]

and the total average intra-cluster dissimilarity:

\[
D = \frac{1}{K} \sum_{v=1}^{K} |C^{(v)}| D^{(v)} .
\]

This finally yields the cost function

\[
E = KD = \sum_{v=1}^{K} \frac{1}{|C^{(v)}|} \sum_{i,j=1}^{w} u_i^{(v)} u_j^{(v)} \gamma_{ij}
\]

that quantifies the average distance within the clusters. The cost function \( E \) can be minimised numerically, for example with simulated annealing (see Ref. [63] for details).

Let us illustrate the use of this approach with the same example studied previously, the collection of generalised baker map data (Eq. (23)) generalises the prediction error \( \gamma_i \) used for the histogram approach to a symmetrised cross-prediction error \( \gamma_{ij} \). With this, two clusters are formed by minimizing \( E \). In Fig. 16, for each series the average dissimilarity \( D_i^{(2)} \) to cluster 2 is plotted against that to cluster 1 (\( D_i^{(1)} \)). Two distinct groups can easily be seen that coincide perfectly with the correct classes. Indeed the algorithm forms exactly the two desired clusters.

Of course, to some extent the problem has only been shifted from finding a magic characteristic number to finding a suitable (not much less magic) dissimilarity measure. However, the approach augments the set of available tools in a meaningful way. After all, we want to classify signals by their dynamics, a feature that is not usually well described by a few parameters. For any classification method, the major problem remains to separate those differences that are relevant for
the discrimination task from those that are not. Sometimes, the calibration of the measurement is known to be of no significance, in which case we can subtract the mean and rescale to unit variance. But apart from such simple transformations we have so far little means of being selective in a controlled way. It is possible to exclude the linear correlation structure from the analysis by normalising characteristic parameters to values obtained with surrogate data.

8. Conclusion and future perspectives

This review paper tries to give an impression on how useful time series methods from chaos theory can be in practical applications. Chaos theory has attracted researchers from many areas for various reasons, in particular, because of its ability to explain complicated temporal behaviour by equations with only a few degrees of freedom and without assuming random forcing to act on the system. The attractiveness of the new paradigm (or the desperation in fields where standard time series methods fail miserably) has tempted people to take several steps at a time, and high expectations have been raised. Only now the path is being retraced step by step. Starting from a theoretical understanding of the new class of systems, time series methods have been tested on computer generated and well controlled laboratory data. Some of these studies provided sobering experience as to how fragile chaos and fractals can be and one could now be tempted to become quite pessimistic about the usefulness of algorithms derived from these concepts.

Now that procedures have been revised, limitations and pitfalls have been pointed out, and intuition has been gained, we can again try to expose the algorithms to field measurements. We will do it less naively than people have done previously – but also with much more modest expectations. The aim of this paper is to get away from naive enthusiasm, but also from a roundabout abandonment of the approach. Realistic applications will only be possible with a pragmatic attitude – what can we learn from the new methods even if the assumption of determinism is not really valid for the system we study?

The obvious goal for the near future is thus to enlarge the class of time series problems that are only, or better, or more efficiently, solvable by the nonlinear approach. The main obstacle will
probably not be the lack of good quality data. Experimentalists have come a long way towards controlling devices and measurement apparatuses. The major challenge lies in the nature of the systems that are most interesting. Many outstanding time series problems in the bio-geo- and social sciences involve multiple time scales, or put differently, lead to nonstationary signals. Also, natural systems are never isolated and thus are of a mixed nature, containing intrinsic and external dynamical components.

There has been a noticeable shift of focus in recent research from the mere testing for nonstationarity and, if the result is positive, excluding a time series from study, to the development of tools to understand the nature of the changes in a system’s dynamics. If a faithful parametric or empirical model for the process is lacking, there is no obvious set of parameters whose changes could be monitored. A promising approach is to define a basis to describe changes in the dynamics either by a number of reference sets [64] or a number of clusters of dynamically similar reference states (Refs. [63,66], Sections 4.3 and 7). The distances, or dissimilarities, of the dynamics at a given time to these reference dynamical states constitutes then a natural set of parameters. Further work on the problem of how to quantify the similarity of dynamical states and on how to use that information could be rewarding. It may then be possible to answer questions about the number of time-dependent parameters and the time scales and nature of these variations.

Systems with many degrees of freedom are notoriously difficult to study through time series, even if multiple recordings are available. As an extreme example take the dynamics of the human brain. The neurons and synapses are not only enormous in number but they are also highly connected. The connection structure can moreover change slowly with time. The system is quite inhomogeneous and has to carry out many different tasks at different times. Certainly, only very specific questions can be hoped to be answered on the basis of time series recordings with a few channels. But even time-resolved imaging techniques can only give a coarsened picture and do not adequately represent the connection structure. Fortunately, there are some interesting intermediate problems that carry more promise to be tractable with dynamical methods. Spatial homogeneity and local coupling leads to a class of systems which can show interesting but still understandable dynamics. Apart from steady states and static patterns, they can exhibit phenomena which are summarised under the term *weak turbulence*. Neither this term, nor the notions of *spatio-temporal chaos* or, more fashionably, *extensive chaos* have been clearly defined so far. This is a direct consequence of the lack of a unifying framework for the study of these systems. Extensivity in this context means that quantities like the number of degrees of freedom, attractor dimension, entropy, etc., asymptotically grow linearly with the volume of the system. In the large system limit, one can then define *intensive* quantities like a dimension density.

This paper is certainly not the place to review the huge literature on nonlinear, spatially extended dynamics, in particular since few of the approaches have been shown to be useful when analysing observational data. The reader may find interesting material and additional pointers to the literature in the proceedings volumes Refs. [216,217], as well as in [218–220]. Probably, the most immediate problem when analysing spatio-temporal data is the choice of a useful representation of the system states and dynamics. A high-frequency sequence of images contains an amount of information that is hardly manageable, even with a powerful computer. The other extreme, a small number of local probes, causes severe problems since the time delay embedding technique is of very limited use with high dimensional data [221]. Popular schemes to reduce the spatial information to a few modes (like for example the Karhunen–Loève decomposition, see Ref. [222] for a recent
application) are most often linear in nature and therefore not quite appropriate for nonlinear systems. In certain situations [223], nonlinear mode dynamics have been used successfully to describe spatio-temporal phenomena. A different approach that carries promise in this respect is the representation by temporally periodic recurrent patterns, or unstable periodic orbits, see the works by Christiansen and coworkers [224], and by Zoldi and Greenside [225]. Despite these efforts, the expectation raised in 1991 at the Santa Fe Institute Time Series Contest that within five years we may have enough experience to enter a second contest, this time on spatio-temporal data, has not been substantiated and it seems that more than a slight relaxation of the time frame will be necessary.

High-dimensional signals can also be produced by systems with only a few components when a delayed feedback is involved. In biology, delayed feedback loops are quite common due to the retarded response of subsystems to changes in other parts of the system. In other fields, delayed feedback can be realised, for example, when part of the output of a device is reflected back from a finite distance, as it sometimes happens in laser or radar equipment, but also with seismic waves. Delayed feedback is also often used for the control of chaotic systems. There has been recent progress in the analysis of such systems, in particular if some knowledge about the feedback structure is available. Bünnher and coworkers [226] have been able to extract relatively simple dynamical equations from scalar time delay systems on the basis of time series data, despite the high dimensionality of the signals. It should be possible in principle also to infer the delay structure from observations. The recovery of dynamical equations from data could then provide a better understanding of many systems in nature.

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