Unraveling the Thousand Word Picture: An Introduction to Super-Resolution Data Analysis

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Abstract

Super-resolution microscopy provides direct insight into fundamental biological processes occurring at length scales smaller than light’s diffraction limit. The analysis of data at such scales has brought statistical and machine learning methods into the mainstream. Here we provide a survey of data analysis methods starting from an overview of basic statistical techniques underlying the analysis of super-resolution and, more broadly, imaging data. We subsequently break down the analysis of super-resolution data into four problems: the localization problem, the counting problem, the linking problem, and what we’ve termed the interpretation problem.

Graphical abstract

1. A GLANCE AT SUPER-RESOLUTION

Processes fundamental to life, including DNA transcription, RNA translation, protein folding, and assembly of proteins into larger complexes, occur at length scales smaller than the diffraction limit of light used to probe them (<200 nm).

For this reason, up until a decade ago, these processes were largely inaccessible to conventional microscopy methods. Key technical achievements by way of experiments, from structured illumination methods\(^1,2\) to manipulations of fluorophore photophysics,\(^3\)–\(^5\) have peered into this previously impenetrable scale with several techniques now providing detailed in vivo 3D images as shown in Figure 1.\(^6\)

On the experimental front, many technical challenges remain including the following: high density labeling; poor time resolution at the expense of high spatial resolution; challenges with fluorophore activation and complex photophysics; overexpression of select proteins altering cell homeostasis; and high light intensity, some \(\approx 10^4\) times higher than that under which cells have evolved for methods such as photoactivated localization microscopy.\(^7\) Despite these challenges, experiments have begun to resolve the spatiotemporal dynamics
and organization of cellular components within their native environment, revealing, for instance, the intricacy of yeast DNA transfer from mother to daughter cell and the stochastic assembly of chemoreceptors on *E. coli*’s surface. What is more, recent advances in optics have mitigated the spatial-temporal resolution trade-off providing greater in vivo resolution in 3D. Advances continue to accrue, with the latest techniques reaching spatial resolutions of ~1 nm and temporal resolutions on the order of μs.

Ten years have passed since the inception of super-resolution microscopy and the variety of data collected has presented new modeling challenges. Initial data analysis methods, such as mean square displacement analyses, were directly motivated from the analysis of bulk ensemble data largely inspired by Occam’s razor. Thus, such methods did not explicitly take advantage of the richness of single molecule data sets such as their temporal ordering or even their intrinsic heterogeneity.

A large fraction of this review is devoted to later “data-driven” efforts, deeply inspired from the fields of machine learning and inference, and increasingly available through an array of open-source software, to turn the thousand-word picture provided by super-resolution methods into a quantitative narrative.

Here, we first review the basic physics of super-resolution methods (section 2) and tools of data-driven modeling (section 3). Subsequent sections tackle specific challenges present along the way: the localization problem (section 4), the counting problem (section 5), the linking problem (section 6), and what we’ve termed the interpretation problem (section 7).

## 2. BEATING THE DIFFRACTION LIMIT: AN INTRODUCTION

### 2.1. Why Fluorescence Microscopy?

Upon excitation of a sample within a specific wavelength range (the absorption spectrum), a fluorophore emits light at a longer wavelength (the emission spectrum). The excitation wave-length may be filtered away leaving behind only the emission from the fluorescent components. In this way, fluorescence brings improved contrast to microscopy.

The first fluorescence microscopes, developed by the Carl Zeiss company and others in the early 20th century, relied either on the autofluorescence of various tissues or chemical dyes and stains such as fluorescein. An important milestone in increasing the ability to fluorescently label a given biological structure was achieved by Coons et al. in the 1940s, who demonstrated that antibodies, raised to bind a specific antigen with high specificity, could be attached to fluorescent dyes, thus realizing a method to fluorescently label any antigen of interest. The subsequent discovery of the green fluorescent protein, together with advances in molecular biology techniques, then allowed the expression of proteins directly fused to fluorescent markers by the end of the 20th century.

At the same time, the detection of the signal from single fluorophores (rather than larger labeled structures) was achieved by progressive improvements in instrumentation. This powerful combination of new advanced optical techniques with fluorescent protein tags,
which could be detected in live cells at the single molecule level, set the stage for a new era of measurements in cell biology and biophysics.\textsuperscript{18}

### 2.2. Point Spread Functions and the Diffraction Limit

Although labeling techniques have greatly improved over the last century, fundamental physical reasons have limited the resolution achievable by optical microscopy. Historically, this resolution has been defined as the ability to distinguish two close objects.

As early as 1834, Airy derived the profile of the diffraction pattern, or point spread function (PSF), of a point source of radiation imaged through a telescope, now known as the Airy disk.\textsuperscript{32} He established that “the image of a star will not be a point but a bright circle surrounded by a series of bright rings. The angular diameter of these will depend on nothing but the aperture of the telescope, and will be [sic] inversely as the aperture”.\textsuperscript{32} More precisely, for a telescope of aperture $a$ imaging at a wavelength $\lambda$, the intensity $I$ at an angle $\theta$ from the optical axis, relative to the intensity $I_0$ at the center, is given by

$$I(x) / I_0 = (2J_1(x) / x)^2$$

where $x = (2\pi / \lambda)a \sin \theta$ and $J_1$ is the first order Bessel function of the first kind. Rings appear at the maxima $x = x_1, x_2, \ldots$ of $I(x)$. In the limit of small angles (i.e., $\theta \approx \sin \theta$), these maxima correspond to $\theta_i = \lambda x_i / (2\pi a)$. Thus, the angular diameters of the rings are indeed inversely proportional to the aperture $a$ (Figure 2).

A few decades later, Abbe showed that a similar result held for optical microscopy: a point source imaged at a wavelength $\lambda$ through a microscope objective of numerical aperture $NA$, defined as the product of the index of refraction of the medium between the objective and the sample, $n$, and the sine of the half angular aperture of the objective, $\theta$, yields a spot of size $d \approx \lambda / 2NA$ in the transverse direction and $2\lambda / NA^2$ in the axial direction\textsuperscript{33} (Figure 3).

Whether in astronomy or microscopy, it is the finite extent of the image of a point source that limits our ability to separate two objects nearby. In 1879, Rayleigh suggested a rule, now called the Rayleigh criterion, whereas two diffraction spots could be considered as resolved if their centers were further apart than the center of a spot is from its first zero in intensity\textsuperscript{34} (Figure 4). He emphasized that this rule was simply suggested as an approximation “in view of the necessary uncertainty as to what exactly is meant by resolution”, though this rule still remains in use today.\textsuperscript{35} In fact, it is generally agreed in astronomy that spots up to $\sim 20\%$ closer are resolvable.\textsuperscript{35}

Nowadays, super-resolution imaging continues to leverage ideas and tools from astronomy, both on the experimental\textsuperscript{36} and analysis side.\textsuperscript{37}

Even though the Rayleigh criterion may not be strictly accurate, the resolution of a microscope is certainly inversely correlated with the size of the diffraction spot. As this spot has a size of $d = \lambda / 2NA$ in the transverse direction and $d = 2\lambda / NA^2$ in the axial direction,
improvements to the resolution are achieved by working at a shorter wavelength or larger numerical aperture.

The room for improvement from changes in the wavelength is limited by the spectrum of visible light, $\lambda = 400$ to 700 nm. Electron microscopy achieves a much higher, near-atomic resolution by operating at a pm-scale wavelength, but this comes at the cost of invasive sample preparations, radiation damage to the sample, and low contrast.$^{38}$

The numerical aperture $NA = n \sin \theta$ has also reached its practical limits: now, oil immersion objectives ($n \approx 1.5$) with half angular apertures of more than $60^\circ$ achieve $NA \approx 1.4$. Few (easy to work with) liquids have higher indices of refraction. Taking these improvements together, the smallest spot size that can be achieved is thus around 150 nm in the transverse direction and 400 nm in the axial direction (Figure 5).

### 2.3. Beyond the Diffraction Limit

Objects may be distinguished from one another at a subdiffraction scale by using a combination of methods including structured illumination, stochastic fluorophore activation, and basic data processing.

As an example of the latter, if we approximate the imaging system as a linear system, i.e., where the measured image can be obtained by applying a linear operator (e.g., convolution by the PSF) to the original sample (e.g., the emitter’s original intensity distribution), it is in principle possible to mathematically invert (“deconvolve”) the imaging operator to reconstruct a higher resolution image, by solving a system of linear equations. Unfortunately, theoretical results indicate that the performance of such an approach is strongly limited by noise.$^{40,41}$ Nonetheless, in the context of microscopy, this idea was first implemented by Agard et al.$^{42}$ and may achieve a 2-fold improvement.$^{43}$

Furthermore, Rayleigh’s criterion does not limit the ability to determine to very high accuracy the position of a single point emitter. For example, the center of a single spot can be estimated to a precision length smaller than the size of the spot itself by fitting the emission pattern to a known PSF, or an approximation of it, such as a Gaussian. The central limit theorem then suggests that the accuracy of such a calculation should be proportional to the inverse square root of the number of observed photons.

By determining the approximate position of emitters over a time series of fluorescence images, where the low density of fluorescent markers ensured their spatial separation, Morrison et al. tracked the diffusion of individual low-density receptors on cell membranes, with a resolution of ~25 nm, well below the diffraction limit.$^{44,45}$

Even as early as in 1995, Betzig suggested that such a localization strategy may be applicable in more densely labeled samples as well, provided that “unique optical characteristics” could be imparted on individual fluorophores.$^{46}$ Such “unique characteristics” would allow distinguishing the signals arising from each of the fluorophores; thus, the fluorophores underlying each diffraction spot could then be localized with subdiffraction accuracy.$^{46}$
Betzig’s original suggestion was to discriminate certain molecules that would exhibit a random spread in their zero phonon absorption line width. However, it was instead the serendipitous discovery of a photoconvertible fluorescent protein, that is, a fluorescent protein whose emission spectrum can be modified by a light-induced chemical modification, as well as the development of optically switchable constructs based on organic dyes, that provided the critical advance toward the realization of this proposal in biological samples.

Briefly, the light-induced conversion of probes to a fluorescent state at a slow enough rate ensures that only a few probes are emitting at any given time even if the sample itself is densely labeled, thus generating the sparsity needed for localization in dense environments. Both labeling approaches were shown to be amenable to this technique: the approach based on fluorescent proteins was named (fluorescence) photoactivated localization microscopy ((F)PALM) and the approach based on organic dyes, stochastic optical reconstruction microscopy (STORM).

While this review will primarily focus on techniques that rely on the stochasticity of photoconversion to temporally separate the emission of different fluorophores, it is also possible to exploit another physical phenomenon to enforce this separation. Specifically, as early as in 1994, Hell et al. noted that while the diffraction limit imposes a lower bound on the size of excitation spots, it is possible to decrease the size of this spot by “deexcitation” (stimulated emission depletion, STED) of the fluorophores located on its edges. Specifically, this deexcitation is carried out by alternatively exciting fluorophores within a small region of the sample and immediately illuminating a doughnut-shaped area around this region with a depletion laser, bringing the fluorophores back to their ground state. The intensity profile of this second region is also diffraction limited; however, given enough time, only the fluorophores close to the exact center of the doughnut (where the deexcitation intensity is zero) stay active. Measuring the fluorescence of these remaining fluorophores thus realizes a point spread function that is effectively smaller than the diffraction limit.

A similar approach, relying on the readout of fluorescence along thin stripes rather than small spots, was also developed, under the name of saturated structured-illumination microscopy (SSIM). This method relies on the observation that high spatial frequencies in the fluorophore distribution can be “brought back” to a lower frequency under illumination by a similarly high frequency pattern (i.e., by observing the beats between the two patterns). Using linear optics (structured illumination microscopy, SIM), the illumination pattern itself is diffraction-limited, and thus the resolution improvement of SIM is limited to a factor of 2 over diffraction-limited microscopy; however, the nonlinearity offered by the saturation method described above allows the generation of higher-frequency patterns and thus further gains in resolution.

Ultimately, the fundamental basis for any of these techniques is to note that the diffraction limit was derived under certain “standard”, but not absolute, hypotheses: that all fluorophore positions must be recovered from a single image and that the signal captured depends linearly on the excitation. Attacking the first condition, by spreading the information across
multiple frames, is the approach taken by stochastic photoconversion. STED and SSIM, additionally, also violate the second condition, by operating in a nonlinear regime.

The large improvement in resolution afforded by structured illumination and stochastic activation of fluorophores, together termed super-resolution microscopy, immediately opened the door to a large number of discoveries. As early as in 2007, Shroff et al. demonstrated the ability of two-color super-resolution to resolve the relative positions of pairs of proteins assembled in adhesion complexes, the attachment points between the cytoskeleton of migrating cells and their substrates, which otherwise seem entirely colocalized in diffraction limited microscopy (Figure 6).51

3. DATA-DRIVEN MODELING: KEY CONCEPTS

In the previous section, we discussed the basic physics underlying super-resolution microscopy. Here, we frame the foundational statistical ideas necessary for a quantitative treatment of super-resolution microscopy data that we will use in later portions of this review without exhaustively reviewing all tools or aiming at mathematical rigor.

While the methods and tools presented here are often established in statistics, examples drawn from optical microscopy will be used to motivate the ensuing discussion. We begin with a high-level description of two model building approaches: forward and inverse methods. Forward model-building approaches, whose first step is to posit physical models motivated by observation and first-principles,52 are broadly used in physics and chemistry. By contrast to forward modeling, inverse modeling approaches or “inference” use the data to learn the model under some assumptions.

The list below summarizes the essence of the data-driven, inverse, methods.

Propose model(s): Models or hypotheses are proposed. This includes proposing models for the system of interest such as whether motion of a particle is “directed” or “diffusive” and models for the measurement noise—such as photon statistics—treated in greater detail in section 3.1.1, which affects static53–55 and dynamic22,25,56–58 quantities throughout optical microscopy. Models proposed can be parametric or nonparametric models (as defined in section 3.1.2).

Inference/Data-mining: After a model form is selected, model parameters are inferred, or estimated, from the data. We review both frequentist and Bayesian inference methods (sections 3.2.1 and 3.2.2) and refer the reader to a recent review for information theoretic inference schemes.59

Reject or select model(s): After a model has been parametrized, we compare the model’s prediction(s) and consistency with existing or new data using hypothesis testing. Often the model’s consistency with the data is used to compare candidate models, a topic more broadly referred to as the model selection problem, discussed in section 3.3.
3.1. Deciding How to Propose Models

Models proposed often come along with assumptions about the system and its noise. These assumptions may be only locally valid. That is, while simple models of normal, confined, or directed diffusion\cite{12,60} may be valid over a small region of space or time, these models may not hold over the entire data set. In this case, locally valid models may be stitched or otherwise combined to describe the heterogeneity inherent to live cells.\cite{61–66} Indeed, systematic strategies for combining simpler models are an area of active research in statistics and machine learning\cite{67–69} that we will discuss in greater depth in section 7.3.

As a practical example, combining simple models together is important as is evident from single particle trajectories whose dynamics may be affected by molecular crowding,\cite{70–74} binding,\cite{63,75–79} or active transport.\cite{80–84} While an anomalous diffusion model may appear more appropriate for such in vivo dynamics, normal diffusion models are often still warranted for such apparently anomalous diffusive processes since, for instance, temporal and spatial resolution now experimentally resolvable can measure short diffusive trajectory segments before events typically leading to anomalous diffusion become manifest.\cite{61,63,65,66,85–88}

Once models are obtained and parametrized, models provide a mathematical framework in which to address quantitative questions and test basic model hypotheses.

3.1.1. Signal versus Noise—Fluorescence microscopy data sets are intrinsically noisy. For example, the sample may exhibit background autofluorescence, the camera may introduce shot noise, experimental apparatus may drift during data acquisition, and molecules of interest may move faster than the acquisition rate.

The importance of accounting for measurement noise inherent to optical microscopy has been highlighted in the recent literature.\cite{24,88–93} For instance, it is now becoming more widely appreciated that neglecting measurement noise in single particle tracking (SPT) can be mistakenly interpreted as a signature of anomalous diffusion.\cite{22}

In particular, inverse approaches, where models are directly drawn from the data, strongly rely on criteria to distinguish signal from noise. These approaches are particularly sensitive to noise models,\cite{61,86,94–96} and for this reason, it is often best to simultaneously infer the model for the system of interest together with the noise model in a self-consistent fashion. The hidden Markov model that we describe in greater depth in section 3.4.1 is one such example.

For now, we begin with a motivating example of noise encountered in attempting to localize a point emitter whose true spatial location, $r$, is assumed to be fixed (Figure 7). An estimate $\hat{\mu}$ of the true spatial location $r$ may be obtained, for example, as the center of a Gaussian PSF fitted to the 2D photon count histogram (a standard approach in super-resolution microscopy\cite{53–55}).

A careful description of the various methods for extracting $\hat{\mu}$ from the data (section 4) will highlight how finite photon counts, camera pixelation, camera-type specific noise, and
background autofluorescence all contribute to the noise $^{54,56}$ and cause the estimated position to differ from the true molecular position by an error $\epsilon$, that is,

$$\hat{r} = r + \epsilon$$  \hspace{1cm} (2)

The error $\epsilon$, illustrated in Figure 7, is typically assumed, in static imaging applications, to be a random variable of mean zero.$^{53,54}$

Beyond the type of noise that we have surveyed here that arises in inferring models from static structures, we briefly mention noise sources that arise from molecular motion during imaging such as “motion blur” or “dynamic [measurement] error”$^{22,56}$ discussed in greater depth in sections 6.5 and 7 and illustrated in Figure 8.

Motion blur presents important theoretical challenges: (i) it complicates the task of localizing particles, whose localization may no longer be adequately modeled using temporally uncorrelated Gaussian noise$^{22,25,54–56,58}$ (section 4)); (ii) it complicates linking localized point emitters across frames into a single particle trajectory (section 6); and (iii) it fundamentally alters the interpretation of kinetic data$^{12,15,22,25,58,86,95,98–100}$ (section 7). As a more concrete example of the latter, the point emitter’s motion within a single image frame statistically couples measurement noise to the molecule’s thermal fluctuations$^{22,25,58}$ which is a coupling often ignored.$^{54,55,60}$

Faster data acquisition may not necessarily reduce dynamical noise since, as is described in the stochastic process time series literature,$^{101,102}$ improved temporal resolution comes hand-in-hand with increased measurement noise (section 6.5). As a result, recent literature has begun addressing the problem of disentangling measurement noise from thermal noise in high frequency (above 10 kHz) single molecule measurements.$^{103–105}$

Noise is unavoidable and the mathematical frameworks we present in subsequent sections present principled strategies for incorporating knowledge of noise statistics into the process of model inference.

3.1.2. Parametric vs Nonparametric Modeling Frameworks—Once a model for both the system of interest and its associated noise are selected, the next step is to infer all model parameters. The main feature distinguishing parametric from nonparametric models is the number of parameters or, more precisely, how this number depends on the amount of data available.

Parametric methods have a fixed number of parameters, independent of the amount of data.$^{106}$ In other words, a parametric model’s functional form, $M$, is prespecified and its parameters, $\theta = \{ \theta_1, \theta_2, \ldots, \theta_K \}$, are $K$ numbers to be determined from the data, $D = \{ D_1, D_2, \ldots, D_N \}$. Thus, rigid assumptions on the nature of the model are made before data is even considered. Despite this apparent shortcoming, parametric models also present important advantages: they are easily interpretable and parameters can be estimated “efficiently”, where efficiency is measured in terms of an estimator’s parameter variance relative to the Cramér-Rao lower bound (CRLB).$^{54,107}$
As an example of a parametric model, we assume that we are modeling the data as independent identically distributed (iid) draws from a 1D Gaussian density indexed by its parameters $\theta \equiv (\mu, \sigma)$. The model, captured by the conditional probability of the data given the model, is

$$
M = \left\{ p(D_i|\theta) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp \left( -\frac{(D_i - \mu)^2}{2\sigma^2} \right), \mu \in \mathbb{R}, \sigma > 0 \right\}
$$

(3)

Having just selected a parametric family for $p(D_i|\theta)$ we may now use the data to provide an estimate for $\theta$ denoted by $\hat{\theta}$.

The model above makes a number of implicit assumptions. For example, it assumes a priori that there are no statistical dependencies between measurements; noise is time-independent, the density of the random variables are unimodal, the mean of the random variables equals the mode of the distribution, etc. These a priori specifications highlight the type of rigid structure that is often imposed by a parametric model.

By contrast, the phrase “nonparametric model” implies that a model is not initially characterized by a finite, fixed number of parameters. Nonparametric methods are advantageous as they provide the flexibility to account for features arising in the data that were not a priori known. For example, if the process under consideration is characterized by a fixed, albeit multimodal density, where the number of modes is unknown in advance, a nonparametric model, such as a kernel density estimator, would reveal the more complex shape of the density characterizing the data generating process as more data becomes available. By contrast, a parametric approach would likely misbehave since the multimodal aspect of the data was not explicitly accounted for in advance.

Although nonparametric approaches are more general, they do not avoid modeling assumptions and infinite amounts of data do not necessarily “automatically guide one to the correct model”. More concretely, the iid assumption, which can be verified by statistical methods discussed in section 3.3), may still be incorrectly invoked by a multimodal (mixture) model. In particular, the “identical” aspect of iid is often suspect, as the distribution from which each observation is sampled may change over time due to, for example, drift in the alignment of optical components.

### 3.2. Inference

Parametric and nonparametric models alike contain parameters that we must learn, or infer, from the data. Here we review two approaches to parameter inference, frequentist or Bayesian, and we refer the reader to Tavakoli et al. and Pressé et al. for information theoretic approaches.

#### 3.2.1. Frequentist Inference

The term “frequentist” in the statistics literature implies that there exists a fixed parameter $\theta$ responsible for generating the observations. The frequentist approach can be summarized as follows: (i) “probabilities” refer to limiting...
Relative frequencies of events and are objective properties and (ii) parameters of the data generating process are fixed (typically unknown) constants. Maximum likelihood is a common frequentist approach to parameter estimation and it yields a maximum likelihood estimate (MLE) for an unknown parameter. Briefly, in maximum likelihood estimation, we begin by writing down a likelihood, the probability of the sequence of observations given a model,

$$\mathcal{L}(\theta) = p(D|M, \theta) = p(D_1, D_2, \ldots, D_N|M, \theta)$$

conjectured to have given rise to the observations. For shorthand, we drop the $M$ dependence such that $p(D_1, D_2, \ldots, D_N|M, \theta) = p(D_1, D_2, \ldots, D_N|\theta)$. We then maximize the data’s likelihood with respect to the parameter vector $\theta$ associated with the posited model $M$.

For example, for iid observations where $p(D_1, D_2, \ldots, D_N|\theta) = \prod_{i=1}^{N} p(D_i|\theta)$, the MLE for $\theta$ is

$$\hat{\theta} = \arg \max_{\theta} \log p(D|\theta) = \sum_{i=1}^{N} \log p(D_i|\theta)$$

(4)

where the MLE, $\hat{\theta}$, maximizes the log likelihood function $\lambda(\theta|D)$. While both the likelihood and its logarithm yield identical MLEs, likelihoods are typically small quantities and logarithms are used to avoid numerical underflow problems.

To make our example specific, suppose we want to estimate the emission rate of a fluorophore, $r$, from a single count of the number of photons emitted, $n$, in time interval $\Delta T$. Our model assumes that all photons emitted are collected and that the number of photons emitted per time interval is Poisson distributed. Under these model assumptions, the likelihood of observing $n$ photons given our model is

$$p(D=n|\theta=r) = \frac{(r\Delta T)^n}{n!} e^{-r\Delta T}$$

(5)

By maximizing the likelihood, eq 5, with respect to the rate, $r$, we obtain an estimate for the rate in terms of known quantities, $\hat{r} = n/\Delta T$. Any other value of $r$ would decrease the likelihood of observing $n$ photons in the interval $\Delta T$.

### 3.2.1.1. Time Series Analysis and Likelihoods:

Likelihood ideas readily accommodate time series data as they respect the natural time ordering of observation. This feature of likelihood methods is especially useful in treating temporally correlated observations that arise in tracking superresolved particles or even counting fluorophores within a region of interest.
For example, consider the simplified case where the data consists of a series of coordinates at different points in time $\mathbf{D} = \{\mathbf{r}_1, \mathbf{r}_2, \ldots, \mathbf{r}_N\}$. Here, the likelihood over the full joint distribution of the time ordered data, $p(\mathbf{D}|\theta)$, is given by

$$p(\mathbf{D}|\theta) = p(\mathbf{r}_1|\theta) \prod_{i=2}^{N} p(\mathbf{r}_i|\mathbf{r}_{i-1}, \mathbf{r}_{i-2}, \ldots, \mathbf{r}_1, \theta) \tag{6}$$

While we’ve assumed in writing $\mathbf{D} = \{\mathbf{r}_1, \mathbf{r}_2, \ldots, \mathbf{r}_N\}$ that the data provided positions with no associated uncertainty, we discuss noise in greater depth throughout this review.

Accurately computing the likelihood and maximizing it with respect to all model parameters, such as all conditional probabilities that appear in eq 6, is difficult because of the large number of conditional probabilities that appear when $N$ is large. The Markov assumption for first-order Markov processes, that the probability of sampling $\mathbf{r}_i$ at time point $i$ only depends on the state of the system in the immediate past ($i-1$), is therefore commonly invoked in SPT analysis in order to reduce the number of parameters that must be inferred. From this point on, whenever we refer to Markov processes or the Markov assumption we assume they are first-order. While the validity of the Markov assumption can also be checked, there are information theoretic reasons for a priori preferring Markov models, discussed elsewhere.

Under the Markov assumption, eq 6 reduces to the more tractable form

$$p(\mathbf{D} = \{\mathbf{r}_1, \mathbf{r}_2, \ldots, \mathbf{r}_N\} | \theta) = p(\mathbf{r}_1 | \theta) \prod_{i=2}^{N} p(\mathbf{r}_i | \mathbf{r}_{i-1}, \theta) \tag{7}$$

with the Markov process’ conditional density, $p(\mathbf{r}_i | \mathbf{r}_{i-1}, \theta)$, often called the “transition density”.

As a more concrete example for which the transition densities take specific functional forms, suppose the model, $\mathcal{M}$, coincides with simple normal diffusion and observations are uniformly spaced $\Delta t$ time units apart. For this example, the transition density of a time series measured without noise, $p(\mathbf{r}_i | \mathbf{r}_{i-1}, \theta)$ is available in closed-form and is a Gaussian density with mean $\mathbf{r}_{i-1}$ and, for standard SPT models, covariance matrix $2D\Delta t$ where $D$ is a diagonal matrix of diffusion coefficients.

In more general Markovian time series models, obtaining the transition density often requires solving the Fokker–Planck equation associated with the posited model $\mathcal{M}$ for the model’s transition density. Since the Fokker–Planck equation may lack a closed-form solution, the transition density may therefore not have an analytically form adding to the computational cost of likelihood maximization. While it is true that individual likelihoods are simple, their products can become very complicated and significantly increase the computational burden. In and of itself that should not be a prohibitively expensive computational issue. However, given the size of the computational burden already needed
for processes such as linking algorithms, an increase due to the need to solve numerically elaborate likelihoods might well be the proverbial “straw that breaks the camel’s back”.

While some popular Markovian motion models arising in super-resolution trajectory analysis (section 7.1) admit closed-form transition densities,25,86 others do not. For this reason, we mention that approximating a likelihood may introduce biased parameter estimates which, in turn, may affect the outcome of model goodness-of-fit tests (section 3.3) used to check the consistency of the model with the data.

3.2.1.2. Confidence Intervals: The precision of a parameter’s MLE is quantified by the concept of confidence interval (or a confidence set if multiple parameters are estimated). In order to introduce this concept, we briefly return to the example shown in Figure 7 where we suppose that an emitter, fixed in space, does not photobleach. Under this circumstance, we repeatedly image the molecule, and for each of the N independent images measured, a different number of photons are collected due to the inherently stochastic emission properties of the fluorophore, as well as other noise sources.

As a result of these variations, different images produce different empirical PSF fits. These fits yield the following hypothetical estimates of a molecule’s position

\[
\hat{\mu}_1 = \mathbf{r} + \mathbf{e}_1 \\
\hat{\mu}_2 = \mathbf{r} + \mathbf{e}_2 \\
\vdots \\
\hat{\mu}_N = \mathbf{r} + \mathbf{e}_N
\]

For each estimated mean, \( \hat{\mu}_i \) we may then compute confidence intervals at the level of \( 1 - \alpha \) where \( \alpha \in [0,1] \) (and where \( \alpha = 0.05 \) is typically used, which corresponds to computing the 95% confidence interval).

Confidence intervals are functions of the observed (random) data and the assumed model. Within a frequentist inference framework, the true fixed parameter either does or does not fall within the confidence interval106 (i.e., with probability 1 or 0 since the frequentist parameter is not random), but the width of the confidence interval does give an indication of the “quality” of the measurement. It should be noted that, despite being very useful in accounting for random errors, confidence intervals critically depend on the assumed model and can fail if this model is wrong or if undetected systematic errors exist.

In our example, one could compute the collection of confidence intervals associated with each of the estimates \( \hat{\mu}_i \). By definition, if the data was indeed generated from the specified model, then a fraction \( 1 - \alpha \) (e.g., 95%) of these confidence intervals will contain the true (unknown but fixed) molecular position, \( \mathbf{r} \).

As mentioned earlier, maximum likelihood methods are popular in part because, in the limit of infinite sample size, they asymptotically achieve the CRLB.54,119–121 In the non-asymptotic regime, it is ultimately the shape and breadth of the full likelihood function around its maximum which provides an estimate for the quality of the MLE.
3.2.1.3. Other Frequentist Estimators: Maximum likelihood is not the only type of frequentist method. Among other frequentist methods, we can mention least-squares regression\textsuperscript{53} which reduce to maximum likelihood under the assumption of Gaussian uncorrelated noise, generalized moment-type methods\textsuperscript{110} (which match select model moments to the data), and general least-squares (which minimize the distance of a parametrized function from the data).

While typically less efficient (as measured by the CRLB) than likelihood maximization,\textsuperscript{110} both moment and least-squares methods can be attractive, because computing the exact likelihood can be complicated or intractable in some models since the full joint distribution may not be known or easily computable. Hence, an approach only requiring specification of a few moments or minimizing the distance to a parametric function may be preferred.

3.2.1.4. Fundamental Assumptions of Frequentist Inference: We end this section on frequentist inference with a note on the assumption that the parameter to be inferred is “fixed”. While this may appear as an advantage, heterogeneity inherent to live cells may require the assumption of having a single fixed parameter vector describing all data to be relaxed.

What is more, in live cell super-resolution applications, in vivo conditions can never be exactly reproduced and thus each experiment is never a true iid draw as described in our hypothetical example used to illustrate the confidence interval concept.

This being said, a collection of point parameter MLEs can still be useful in quantifying live cell heterogeneity. For example, once the MLE of an assumed model and data set are obtained, this MLE can be used to verify the model’s consistency with the data. If the model is deemed consistent, the MLE can be used to compute a theoretical CRLB or to simulate new observations in order to compute “bootstrap estimates” of the MLE parameter’s confidence interval.\textsuperscript{106,122} The simplest indication of heterogeneity would then be suggested by nonoverlapping confidence intervals.

3.2.2. Bayesian Inference—Bayesian methods, now widely used across biophysics,\textsuperscript{61,65,96,123–127} are defined by the following key properties:\textsuperscript{106} (i) a probability describes a degree of belief, not a limiting frequency, (ii) distributions for parameters generating the data can be defined even if parameters themselves are fixed constants, and (iii) inferences about parameters are obtained from “posterior probability” distributions.

While frequentist methods yield parameter “point estimates” or MLEs, Bayesian methods return a full distribution over unknown parameters (the posteriors) for the same amount of data. Because the posterior is a joint distribution over all parameters, it can also quantify relationships between parameters. Therefore, while frequentist methods treat the data as a random variable, Bayesian methods treat data in addition to model parameters, as random variables.\textsuperscript{128}

In frequentist reasoning, the likelihood played a central role. By contrast, in Bayesian inference, it is the posterior, \( p(\theta | D, M) \). That is, the probability of the model parameters
having considered the data (and the model’s structure, $M$). Again, for shorthand, we drop the $M$ dependence of the posterior.

We can construct the posterior using the likelihood as well as the prior over the model, $p(\theta)$, by invoking Bayes’ theorem as follows

$$p(\theta | D) \propto p(D | \theta) p(\theta)$$

where $p(D)$ is obtained by normalization from

$$p(D) = \int p(D, \theta) d\theta = \int p(D | \theta) p(\theta) d\theta$$

where $p(D, \theta)$ is the joint probability of the model and the data.

For parametric models, the integral in eq 8 represents an integration over the model’s parameters. However, for multiple possible candidate models, we first need to sum over all models and, subsequently, integrate over their respective parameters.

We may even compare models, say $M_l$ for different values of $l$ selected from a broader set of models $M$ by integrating over all parameters assigned to that model

$$p(M_l | D) \propto \int p(D | M_l, \theta_l) p(\theta_l | M_l) p(M_l) d\theta_l$$

Then marginal posteriors, $p(M_l | D)$, across different $l$'s can be compared to select between different candidate models irrespective of their precise parameter values.$^{124}$

For an increasing number, $N$, of data points, the likelihood function itself dictates the shape of the posterior distribution. Thus, a proper choice of likelihood function, even in Bayesian inference, is critical (Figure 9).

3.2.2.1. Priors: Parameter estimates obtained by maximizing a posterior generally converge to the MLE in the limit of infinite data$^{108}$ (Figure 9). However, with finite data sets, faced in all data-driven modeling efforts, the prior can cause Bayesian and MLE estimates for parameter values to differ substantially (with Bayesian parameter point estimates coinciding with the maximum of the posterior distribution). This suggests that, for scarce data, a judicious choice of prior is important.

An extreme realization of limited data is found in traditional statistical mechanics where a (canonical) distribution over arbitrarily many degrees of freedom for particles follows from just a single point observation (the average energy) with vanishingly small error.$^{113,129,130}$ While more data is typically assumed available in biophysics, data may still be quite limited. For example, in single particle tracking, individual tracks may be shortened due to photobleaching or by having particles move in and out of focus.$^{61}$ Priors can be tailored to address many issues arising from data sparsity thereby conferring an important advantage to
Bayesian methods and, in particular, the analysis of super-resolution data. For this reason, we discuss two different prior types: informative and uninformative.\textsuperscript{128,131}

\subsection*{3.2.2.2. Uninformative Priors}

The constant, uniform prior, inspired by Laplace’s principle of insufficient reason for mutually exclusive hypotheses, such as coin flips, is the simplest uninformative prior.

For a prior, $p(\theta)$, constant over a reasonable range, the likelihood and posterior have precisely the same dependence on model parameters, $\theta$.

$$p(\theta|D) \propto p(D|\theta) p(\theta) \propto p(D|\theta) \quad (10)$$

As a result, in this simple case, the likelihood can be treated as a model posterior.

A counterintuitive implication of assuming uninformative priors follows from the observation that uniformity of a prior over a model parameter, say $\theta$, implies some structure, and thus knowledge, on the distribution of a related variable, say $e^\theta$.\textsuperscript{128} Conversely, if we assume that the variable $e^\theta$ is uniform on the [0, 1] interval, the variable $\theta$ is concentrated at 1.

This problem is addressed by invoking priors invariant under coordinate transformations, such as the Jeffreys prior,\textsuperscript{132–134} which is sometimes used in the biophysical literature.\textsuperscript{135–138}

\subsection*{3.2.2.3. Informative Priors}

Bayes’ theorem, the recipe by which priors, $p(\theta)$, are updated into posteriors, $p(\theta|D)$, upon the availability of data, directly motivate another prior form.

Specifically, if we insist that priors and posteriors be of the same “mathematical form”, then the form of the likelihood sets the form of the prior. Such a prior, which when multiplied by its likelihood generates a posterior of the same form as the prior, is called a conjugate prior.

An immediate advantage of using conjugate priors is that, when additional independent data (say $D_2$ beyond $D_1$) are incorporated into a posterior, the new posterior $p(\theta|D_2, D_1)$, obtained by multiplying the old posterior, $p(\theta|D_1)$, by the likelihood, is again guaranteed to have the same form as old posterior and the prior

$$p(\theta|D_2, D_1) \propto p(D_2|\theta) p(\theta|D_1) \quad (11)$$

Priors, say $p(\theta|\gamma)$, may depend on additional parameters, $\gamma$, called hyperparameters distinct from the model parameters $\theta$. These hyperparameters, in turn, can also be distributed according to a hyperprior distribution, $p(\gamma|\eta)$, thereby establishing a parameter hierarchy.

As a specific example of a parameter hierarchy, an observable (say the FRET intensity) can be function of the state of a protein, which depends on transition rates to that state (model parameters), which, in turn, depends on prior parameters (hyperparameters) determining
how transition rates are assumed to be a priori distributed. We will see examples of such hierarchies in the context of later discussions on nonparametric Bayesian methods (section 3.4).

We illustrate the concept of conjugacy by returning to our earlier single emitter example. The prior conjugate to the Poisson distribution with parameter $\lambda$, where $\lambda$ plays the role of $r\Delta T$ (eq 5), is the gamma distribution

$$\text{gamma}(\alpha, \beta) = \frac{\beta^\alpha}{\Gamma(\alpha)} \lambda^{\alpha-1} e^{-\beta \lambda}$$  \hspace{1cm} (12)

which contains two hyperparameters, $\alpha$ and $\beta$. After a single observation of $n_1$ emission events in time $\Delta T$, the posterior is

$$p(\lambda | n_1, \alpha, \beta) = \text{gamma}(n_1 + \alpha, 1 + \beta)$$  \hspace{1cm} (13)

while, after $N$ independent measurements, with $D = \{n_1, \ldots n_N\}$, we have

$$p(\lambda | D, \alpha, \beta) = \text{gamma}\left(\sum_{i=1}^{N} n_i + \alpha, N + \beta\right)$$  \hspace{1cm} (14)

Figure 9 illustrates both how the posterior is dominated by the likelihood provided $N$ is large and how an arbitrary hyperparameter choice becomes less important for large enough $N$.

Single molecule photobleaching provides yet another illustrative example\textsuperscript{64} and allows us to introduce the topic of inference of a probability distribution relevant to our later discussion on Bayesian nonparametrics.

Here we consider the probability that a molecule has an inactive fluorophore (one that never turns on), a basic problem facing quantitative super-resolution image analysis.\textsuperscript{139} Let $\theta$ be the probability that a fluorophore is active (and detected). Correspondingly, $1 - \theta$ is the probability that the fluorophore never turns on. The probability that $y$ molecules among a total of $n$ molecules in a complex turn on is then binomially distributed

$$p(y|\theta) = \frac{n!}{(n-y)!y!} \theta^y (1 - \theta)^{n-y}$$  \hspace{1cm} (15)

Over multiple measurements (each over a region of interest with one complex having $n$ total molecules) with outcome $y_i$ we obtain the following likelihood:

$$p(y|\theta) \propto \prod_i \frac{n!}{(n-y_i)!y_i!} \theta^{y_i} (1 - \theta)^{n-y_i}$$  \hspace{1cm} (16)
One choice for $p(\theta)$ is the beta distribution, the conjugate prior to the binomial

$$p(\theta) = \frac{(a+b-1)!}{(a-1)!(b-1)!} \theta^{a-1}(1-\theta)^{b-1} \quad (17)$$

whose multivariate generalization is the Dirichlet distribution, which can be further generalized to the infinite-dimensional case (the Dirichlet process) that we will revisit on multiple occasions.

By construction (i.e., by conjugacy), our posterior now takes the form of the beta distribution

$$p(\theta|y) \propto \theta^{\sum_i y_i + a - 1} (1 - \theta)^{\sum_i (n-y_i) + b - 1} \quad (18)$$

Given these data, the estimated mean, $\hat{\theta}$, obtained by maximizing this posterior, called the maximum a posteriori estimate, is now

$$\hat{\theta} = \frac{\sum_i y_i + a}{\sum_i n + a + b} = \frac{\sum_i y_i}{\sum_i n + a + b} + \frac{a}{\sum_i n + a + b} \quad (19)$$

which is, perhaps unsurprisingly, a weighted sum over the data and prior expectation (the first and second terms, respectively, on the right-hand side).

Conjugate priors do have obvious mathematical appeal and yield analytically tractable forms for posteriors (eq 18), but they are restrictive. Numerical methods to sample posteriors, including Gibbs sampling and related Markov chain Monte Carlo (MCMC) methods,\textsuperscript{140,141} continue to be used\textsuperscript{64} and developed\textsuperscript{142} for biophysical problems and have somewhat reduced the historical analytical advantage of conjugate priors. However, the advantage conferred by the tractability of conjugate priors has turned out to be major advantage for more complex inference problems such as those involving Dirichlet processes (section 3.4).

### 3.2.2.4. Credible Intervals:

The analog of a frequentist confidence interval is a Bayesian credible interval. Likewise in higher dimensions, the Bayesian analog of the confidence set is the credible region. A Bayesian $1 - \alpha$ credible region is any subset of parameter space over which the integral of the posterior distribution probability is equal to $1 - \alpha$.\textsuperscript{143} For a more concrete 1D example, suppose that we compute the posterior over a 1D Gaussian random variable. The probability of the Gaussian’s mean parameter falling between a given lower bound $a$ and upper bound $b$ would be the area of the posterior density over $[a, b]$. Any interval having a probability of $1 - \alpha$ is a valid $1 - \alpha$ credible interval.

### 3.3. Reject or Select Model Based on Empirical Evidence

#### 3.3.1. Frequentist Goodness-of-Fit Testing—Before physically interpreting models, we should verify whether various model assumptions are consistent with the data. Checking
a model’s distribution, and its implied statistical dependencies, against measured data falls under the category of “goodness-of-fit tests”.

### 3.3.1.1. Example of Goodness-of-Fit Testing for iid Random Variables:

A classic example of a goodness-of-fit test is the well-known Kolmogorov–Smirnov (KS) test, which we describe here. We assume that $N$ iid samples of a scalar random variable $X$ are drawn from an unknown distribution $F$. The empirical cumulative distribution function (ECDF) is defined as $\hat{F}_N(v) = \sum_{i=0}^{N} I_{X_i \leq v}/N$ where $I_A$ is the indicator function of event $A$. That is, the ECDF evaluated at $v$ is the fraction of times the random variable was equal to or less than $v$.

In the classic KS framework, the statistical model assumes a known distribution $F_0$. The KS test checks the null hypothesis, $H_0: F = F_0$ against the alternative hypothesis $H_a: F \neq F_0$. To decide between hypotheses, one computes the “test statistic” for iid samples $T_N = \max_{v} N^{1/2} |\hat{F}_N(v) - F_0(v)|$. If $T_N$ is greater than a critical value (determined for any desired test accuracy) one rejects the null hypothesis and places less confidence in the model. From a frequentist point of view, the test can be understood as follows: if we generate samples of size $N$ from the distribution $F_0$, their ECDF will typically not match the cumulative distribution function (CDF) of $F_0$ exactly but will be at some “distance” (as quantified by the KS test statistic) from it; we then ask, how likely is it that this distance is as large as the one computed for our data set? If such a large distance occurs with a probability less than $\alpha$, then the null hypothesis is rejected.

As we will discuss in section 3.3.4, more complex models will tend to fit the data better. The KS test does not penalize a model for complexity in any way, it simply checks for consistency of a model’s distribution against the empirically observed data. This should be contrasted with model selection criteria (section 3.3.4) which attempt to balance model complexity and “predictive scores” (such as the likelihood).

### 3.3.1.2. Goodness-of-Fit Testing in Time Series Analysis:

The basic idea behind the KS goodness-of-fit test has been extended to multivariate random variables and to random variables exhibiting time dependence (i.e., time series). While many other goodness-of-fit tests, including the chisquare and Cramér-von Mises tests, exist, here we briefly illustrate how the ideas behind the KS can be generalized to time series data. We consider the case of scalar time series data $D = \{r_1, r_2, \ldots, r_T\}$.

If a Markovian model, with prespecified parameter vector $\theta_0$, is assumed, we can introduce the Rosenblatt transform also known as the “probability integral transform” (PIT) defined by $Z_i = \int_{-\infty}^{r_i} p(x|r_{i-1}, \theta_0) \, dx$ for observation $i$. The sequence of $Z_i$’s is now iid, despite the fact that the sequence was derived from a time series exhibiting statistical dependence or possibly even nonstationary behavior. All of the same tests used for iid variables, such as the KS test above, can be invoked to test the consistency of the model $\theta_0$ with the data.
While we have discussed a scalar Markovian time series, the PIT can be extended to multivariate non-Markovian series provided the joint density (i.e., the likelihood) can be accurately computed.\textsuperscript{111,112}

Goodness-of-fit testing readily apply to 3D trajectories of particles. For instance, in recent work on mRNA motion in the cytoplasm of live yeast cells,\textsuperscript{86} goodness-of-fit testing was used to reject the hypothesis that time series coordinates are statistically independent. This rejection ultimately motivated new spatiotemporally correlated models as well as new models for measurement noise that identified novel kinetic signatures of molecular motor induced transport.\textsuperscript{86}

3.3.2. Bayesian Hypothesis Testing—Previously, we demonstrated how goodness-of-fit testing could be used to verify the consistency of a model’s statistical assumptions against the data. Such an absolute consistency check on a model’s specific parameter values is accomplished by formulating hypotheses or “events” such as, $H_0: \theta = \theta_0$ and $H_1: \theta \neq \theta_0$. Events can also be something with more of a “model comparison flavor”, such as $H_0: \theta = \theta_0$ and $H_1: \theta = \theta_1$.

Here we consider the case where the model $M$ is fixed and we test $H_0: \theta = \theta_0$ and $H_1: \theta \neq \theta_0$. In this type of binary testing, the $H_i$ are events such that individual event probabilities sum to one, $P(H_0) + P(H_1) = 1$. The posterior probability of the null, $H_0$, hypothesis is

\[
P(H_0|D) = \frac{p(D|H_0)P(H_0)}{p(D|H_0)P(H_0) + p(D|H_1)P(H_1)} \tag{20}
\]

where the integral in the denominator is an integral over all $\theta \neq \theta_0$.

While generalizing the ideas here to test multiple hypotheses is straightforward,\textsuperscript{114} it is the selection of priors and choice of models that makes Bayesian hypothesis testing technically complicated.\textsuperscript{106} To be more specific, eq 21 selects between values of the parameters $\theta$ of a single model but does not, in itself, provide information on whether the model itself is correct at all (and thus whether this selection is meaningful).

Before discussing model selection in greater depth (section 3.3.4), we now briefly describe an important model, the hidden Markov model (HMM), that motivates modeling approaches used across biophysics\textsuperscript{146} that we will revisit on multiple occasions.

3.3.3. Hidden Markov Models—HMMs\textsuperscript{146–149} have been used so far in a number of problems, chief among them force spectroscopy and single molecule FRET (smFRET) data analysis.\textsuperscript{149–155}
To introduce the HMM, we begin by considering a sequence of observations $D = y = \{y_1, y_2, \ldots, y_N\}$ and assume that these give us indirect information about the “latent” or hidden states (variables) $s_i$ at every time point $i$, that alone describe the system dynamics. Hidden states and observations are then related by the probability of making the observation $y_i$ given the state of the system $s_i$: $p(y_i | s_i)$ assuming observations are uncorrelated in time.

The HMM then starts with a parametric form for $p(y_i | s_i)$ (the emission probability), the number of states $K$, and the data $D$. From this input, the HMM infers the set of parameters, $\theta$, which includes (i) the “emission parameters” that parametrize $p(y_i | s_i)$, (ii) initial state probabilities, and (iii) the transition rates, i.e., transition probabilities $p(s_i | s_{i-1})$ for all $(i, j)$ pairs, that describe the jump process between states assumed Markovian.

In discrete time, the likelihood used to describe the data for an HMM model is

$$
\mathcal{L}(\theta | D) = p(y | \theta) = \sum_s p(y, s | \theta) = \sum_s \prod_{i=2}^N p(y_i | s_i) p(s_i | s_{i-1}) p(y_1 | s_1) p(s_1)
$$

(22)

where $s = \{s_1, \ldots, s_N\}$ and $i$ is the time index.

In general, we are not interested in knowing the full distribution $p(y, s | \theta)$; instead we only care about the marginal likelihood $p(y | \theta)$ which describes the observation probability given the model, no matter what state the system occupied at each point in time. In other words, since we do not know what state the system is at any time point, we must sum over the $s_i$.

The HMM itself can be represented by the following sampling scheme

$$
s_1 \sim p(s_1) \\
s_i | s_{i-1}, \theta \sim p(s_i | s_{i-1}) \\
y_i | s_i, \theta \sim p(y_i | s_i)
$$

(23)

That is $p(s_1)$ is the distribution from which $s_1$ is sampled. Then for any $i > 1$, $p(s_i | s_{i-1})$ is the conditional distribution from which $s_i$, conditioned on $s_{i-1}$, is sampled, while its observation $y_i$ is sampled from the conditional $p(y_i | s_i, \theta)$.

We then maximize the likelihood for this HMM model, eq 22, over each parameter $\theta$.

Multiple methods are available to maximize HMM-generated likelihood functions numerically. Most commonly, these include the Viterbi algorithm and forward–backward algorithms in combination with expectation maximization.

3.3.3.1. Aggregated Markov Models: Aggregated Markov models (AMMs) are a special case of HMMs where many hidden states have identical output. For instance, for Gaussian $p(y_i | s_i)$, two or more states have identical Gaussian means.
This special category of HMMs was introduced to biophysics in the analysis of ion-channel patch clamp experiments\textsuperscript{159,162–164} where the number of microscopic channel states exceeded the typically binary output of patch clamp experiments. Since then, AMMs have also been successfully applied to smFRET, where a low FRET state may also arise from different microscopic states (e.g., blinking of fluorophore photophysics or an internal state of the labeled protein)\textsuperscript{151} and, recently, to address the single molecule counting problem using super-resolution imaging data.\textsuperscript{21}

In AMMs, experimentally indistinguishable states are lumped together into an “aggregate of states” belonging to an “observability class”, say “bright” and “dark” in the case of smFRET.

Limiting ourselves, for simplicity, to just two aggregates, 1 and 2, we may write a rate matrix $Q$ as follows

$$Q = \begin{bmatrix} Q_{11} & Q_{12} \\ Q_{21} & Q_{22} \end{bmatrix}$$ (24)

The submatrices $Q_{ij}$ describe the transitions from aggregate $i$ to aggregate $j$ and the $kl$ element of $Q_{ij}$ describes transitions from the microscopic state $k$ in aggregate $i$ to microscopic state $l$ in aggregate $j$.

AMMs are subsequently treated in much the same as HMMs. That is a likelihood is constructed and maximized with respect to the model parameters. In continuous time and ignoring noise, the likelihood of observing the sequence of aggregates $D = \{a_1, a_2, \ldots, a_N\}$ is

$$L(\theta | D) = 1^T \cdot \prod_{j=1}^{N-1} G_{a_ja_{j+1}}(t_j) \cdot \pi_{a_1}$$ (25)

where the $i$th element of the column vector, $\pi_{a_1}$, represents the initial probability of being in state $i$ from the $a_1$ aggregate and

$$G_{ab}(t_j) = Q_{ab} e^{Q_{ab} t_j}$$ (26)

The row vector, $1^T$, in eq 25, is a mathematical device used to sum over all final microscopic states of the aggregate, $a_N$, observed at the last time point as we do not know in which microscopic state the system finds itself in within the $N$th measurement.

Just as with HMMs, the $\theta$ parameters include all transitions between microscopic states across all aggregates as well as initial probabilities for each state within each observability class and any emission parameters used. As there are fewer observability classes in AMMs than there are microscopic states, many parameters from AMMs may be unidentifiable\textsuperscript{165} and thus may need to be prespecified.
AMMs and HMMs can be generalized to include the possibility of missed (unobserved) transitions\textsuperscript{21,166} for instance when transitions happen on time scales that approach or exceed the data acquisition frequency–and can readily apply to discrete or continuous time.

One fundamental shortcoming of either HMMs or AMMs is their explicit reliance on a predetermined number of states. This is true of many modeling strategies that predetermine the model “complexity”.

The next section on model selection addresses the challenge of finding the model of the right complexity.

3.3.4. Model Selection—We assume that models under consideration have passed tests, ensuring the interpretive value of the model, and that we are interested in selecting the best (or preferred) predictive or descriptive candidate model among these without excessively overfitting.\textsuperscript{114}

Here we review the basic ideas behind well-established model selection criteria that can be naively described as finding the best model by balancing a “predictive component” (often quantified by a likelihood) and a “complexity penalty” (often quantified by a function of the number of parameters in a model).

In this article, we do not review cross validation type techniques for accessing the predictive ability of models. In cross validation some fraction of the data is kept in reserve in order to see how the prediction performs. Cross validation is used mostly in “man supervised learning” applications, for example “deep learning” of gigantic labeled color image data sets. With such data sets we often have a good idea of what the ground truth is. In super-resolution image analysis the genuine “ground truth” is much harder to come by, except in simulations.

3.3.4.1. Information-Theoretic Model Selection Criteria: While many information theoretic methods exist\textsuperscript{167–169} and have recently been reviewed elsewhere,\textsuperscript{59} here we focus on the Akaike Information Criterion (AIC).\textsuperscript{170–173} The AIC starts from the idea that, while it may not be possible to find the true (hypothetical) model, it may be possible to find one that minimizes the difference in the “information content” between the true (hypothetical) model and the preferred candidate model. The preferred model should provide a good fit to existing data, the training set, and provide predictive power on alternate data sets, the validation set.

While the derivation of the AIC is involved,\textsuperscript{59,114} the AIC expression itself is simple

$$\text{AIC}(M, D) = - 2\log p(D | \hat{\theta}, M) + 2K$$ \hspace{1cm} (27)$$

where $\hat{\theta}$ denotes the MLE of a given model.

Minimizing the AIC, which is equivalent to maximizing the likelihood subject to a penalty on parameter numbers, $K$, generates the preferred model. The AIC itself is an asymptotic result valid in the large data set limit though higher order corrections exist.\textsuperscript{114} The penalty
term, $2K$, is derived, not imposed by hand. Also, as we will see, it is a weaker penalty than that of the Bayesian information criterion that we now describe.

As a concrete example, suppose goodness-of-fit testing cannot rule out the possibility that the model itself is an undetermined sum of exponentials obscured by noise. The AIC would then be used to find the number of exponential components required without overfitting the data. Additional examples motivated from the biophysical literature are provided by Tavakoli et al.\textsuperscript{59}

### 3.3.4.2. Bayesian Information Criterion:

While the mathematical form for the Bayesian (or Schwartz) information criterion (BIC)\textsuperscript{174} used in model selection, that we will see shortly, may appear similar to the AIC, eq 27, it is conceptually very different. It follows from Bayesian logic with no recourse to information theory.\textsuperscript{59,114,175–177} As before, the derivation itself is involved and discussed in detail in Tavakoli et al.\textsuperscript{59}

Unlike the AIC, the BIC searches for a true model\textsuperscript{173,178} that exists regardless of the number of data points $N$ available used to find this model. That is, the model itself is assumed to be of fixed, albeit unknown, complexity, by contrast to the AIC which seeks an approximate model and is therefore more willing to grow the model complexity along with the size of the data set, $N$, available.

Briefly put, the BIC starts from a posterior for a model class with a fixed number of parameters, $K$, just as we had seen with our exponential example earlier. Following the logic of eq 9, we consider the average posterior that comes from summing over parameter values for each of the $K$ parameters generating a marginal posterior.

If a candidate model contains too many parameters, then many of these parameters (integrated over all values that they can be assigned) will result in a small marginal posterior for that particular parameter number. By contrast, if too few parameters are present in the posterior, the model itself will be insufficiently complex to capture the data, and this will also yield small marginal posteriors.

From this logic follows the BIC

$$\text{BIC}(M, D) = -2\log p(D | \hat{\theta} M) + K \log N$$

(28)

which differs from the AIC, eq 27, in the complexity penalty term (second term on the right-hand side).

While the logic underlying both AIC and BIC is different, both AIC and BIC are used interchangeably in practice, and their performance is problem-specific\textsuperscript{59} (Figure 10). For instance, since the AIC tends to overfit small features it may underperform for slightly nonlinear models (which may have fewer parameters) and that may be obscured by noise but may be preferred for highly nonlinear models (with a correspondingly larger number of parameters) as small features may no longer be due to noise.\textsuperscript{172}
3.4. Overview of Bayesian Nonparametrics

Bayesian nonparametrics is a relatively recent (1973) approach to statistical modeling\textsuperscript{179} that fundamentally integrates the model selection step into the model-building process. In other words, it builds models whose complexity reflect that of the data.\textsuperscript{179–181} This is especially important to super-resolution and, more broadly, single molecule data analysis, as few model features, such as the number of states in an HMM, are initially known.

As we previously discussed, nonparametric approaches do make parametric assumptions. For instance, a particular nonparametric model may assume that measurements are iid. However, beyond these parametric constraints, nonparametric approaches allow for a possibly infinite number of parameters.

3.4.1. Dirichlet Process—The prior process, the analog of the prior used in parametric Bayesian inference, is of major importance in Bayesian nonparametrics. The most widely used of these processes is the Dirichlet process (DP) prior,\textsuperscript{179} along with its various representations,\textsuperscript{67,68,182} such as the “Chinese restaurant” process and the stick-breaking construction.\textsuperscript{181,183}

Just as a prior in parametric statistics samples parameter values, the DP samples probability distributions. Probability density estimation,\textsuperscript{184} clustering,\textsuperscript{185} HMMs,\textsuperscript{183} and Markovian switching linear dynamical systems (SLDS),\textsuperscript{68} have all been generalized to the nonparametric case by exploiting the DP.

We introduce the DP with a parametric example and begin by considering probabilities of outcomes, $\pi = \{\pi_1, \pi_2, \ldots, \pi_K\}$ with $\Sigma_k \pi_k = 1$ and $\pi_k \geq 0$ for all $k$, distributed according to a Dirichlet distribution, $\pi | \alpha \sim \text{Dirichlet} (\alpha_1, \ldots, \alpha_K)$. In other words

$$p(\pi | \alpha) = \frac{\Gamma \left( \sum_k \alpha_k \right)}{\prod_k \Gamma (\alpha_k)} \prod_{k=1}^{K} \pi_k^{\alpha_k - 1} \quad (29)$$

The Dirichlet distribution is conjugate to the multinomial distribution.

To build a posterior from the prior above, we introduce a multinomial likelihood with populations $n = \{n_1, n_2, \ldots, n_K\}$ distributed over $K$ unique bins

$$p(n | \pi) = \frac{\Gamma \left( \sum_k n_k + 1 \right)}{\prod_k \Gamma (n_k + 1)} \prod_{k=1}^{K} \frac{n_k^{n_k}}{\sum_k n_k^{n_k}} \quad (30)$$

The resulting posterior obtained from the prior, eq 29, and likelihood, eq 30, is now
The multinomial shown above is a starting point for many inference problems. For example, eq 31 may represent the posterior probability of weights in mixture models, such as sums of Gaussians, and is used in clustering problems.¹⁸⁴

To be more specific, the finite mixture model can be represented using the following sampling scheme with \( y_i \) denoting the \( i \)th observation

\[
\begin{align*}
  y_i | s_i, \theta & \sim F(\theta_{s_i}) \\
  s_i | \pi & \sim \text{multinomial}(\pi_1, \ldots, \pi_K) \\
  \theta_{s_i} & \sim H \\
  \pi & \sim \text{Dirichlet}(\alpha_1, \ldots, \alpha_K)
\end{align*}
\]

(32)

Here \( s_i \) is the (latent or unobservable) cluster integer label, \( F(\cdot) \) is a distribution governing the \( y_i \) random variables assigned to cluster \( s_i \), \( \theta_{s_i} \) parametrizes \( F \) and the prior over the unobservable \( \theta_{s_i} \)'s are governed by the so-called “base distribution” \( H \).¹⁸⁴

Equation 32, formalizes the logic that if a model is a sum of Gaussians (i.e., a sum of \( F(\theta_{s_i}) \)) with each component indexed \( s_i \) with parameters \( \theta_{s_i} \), we draw an observation by first selecting a mixture component (according to the weight of each component, \( \text{multinomial}(\pi_1, \ldots, \pi_K) \)) and subsequently sampling a value from that particular mixture component (i.e., a Gaussian).

In the absence of prior knowledge on the expected mixture components, we set all \( \alpha_k \) to \( \alpha/K \) and, under these assumptions, we can write down the probability of belonging to cluster (or mixture component) \( j \)¹⁸⁴

\[
\frac{n_j + \alpha}{\alpha + N - 1}
\]

(33)

where \( n_j \) is the number of observations assigned to cluster “\( j \)” and \( N = \sum_j n_j \).

The DP is the infinite dimensional (\( K \to \infty \)) generalization of the Dirichlet prior distribution and is used in infinite mixture models.¹⁸⁴ In infinite mixture models, the prior probability of occupying a pre-existing mixture components is

\[
\frac{n_j}{\alpha + N - 1}
\]

(34)
whereas the probability of creating a new component is

\[
\frac{\alpha}{\alpha + N - 1} \tag{35}
\]

Before moving forward, some comments are in order. When doing classic finite mixture modeling, we needed to predefine the number of mixture components, \( K \). Each \( K \) component has a prior probability given by eq 33; if we change \( K \) (and as a result change the \( y_i \) cluster assignments), one is faced with a model selection problem. In the DP mixture model, the data determines the number of active clusters/states as well as the parameters needed to describe the data in an à la carte fashion within a single Bayesian posterior.

In both finite and infinite mixture models, the quality of the clustering strongly depends on the base distribution, \( H \), since this distribution governs the \( \theta \) which determines the similarity (or dissimilarity) between the probability distributions, \( R(\cdot) \), associated with different clusters.

### 3.4.1.1. Infinite Mixture Models:

Previously, in eq 32, we showed how to sample a priori mixture weights for finite mixture models. For infinite mixture components, we may combine the first two as well as the last two expressions of eq 32 and obtain

\[
y_i | F, G \sim F \ast G \\
G \sim \text{DP}(\alpha, H) \tag{36}
\]

where \( \alpha \) is the “concentration parameter” quantifying the preference for creating new clusters\textsuperscript{184} while \( H \) is a judiciously selected\textsuperscript{68,179,186} “base distribution” that plays the same role here as it did for finite mixture models in setting \( \theta_k \). \( G \) is the DP’s sample and it is a random distribution that can be represented as follows

\[
G = \sum_{k=1}^{\infty} \pi_k \delta_{\theta_k} \tag{37}
\]

where the Kronecker delta, \( \delta_{\theta_k} \), denotes a mass point for parameter value \( \theta_k \)\textsuperscript{67,187} and \( F \ast G \) denotes

\[
F \ast G = \sum_{k=1}^{\infty} \pi_k F(\theta_k) \tag{38}
\]

It can be shown that the mean of the random distribution \( G \) is given by \([G] = H\) and its variance is \((H(1-H))/(\alpha + 1))\textsuperscript{108}. Thus, as \( \alpha \) increases, \( G \) approaches \( H \). Hence the concentration parameter can be thought of as constraining the similarity between \( G \) and \( H \).
While we have just described infinite mixture models, we have not yet described how to draw samples from the DP. DP itself is often represented by the “stick-breaking construction” which is easily implemented computationally\textsuperscript{182}

\[
v_k \sim \text{beta}(1, \alpha) \\
\pi_k = v_k \prod_{j=1}^{k-1} (1 - v_j) \\
\theta_k \sim H
\]

Following this procedure, we find that \( G \sim \text{DP}(\alpha, H) \).\textsuperscript{108}

The analogy to stick-breaking here is motivated by imagining that we begin with a stick of unit length. We break it at a location, \( v_1 \), sampled from a beta distribution \( v_1 \sim \text{beta}(1, \alpha) \). We set \( \pi_1 = v_1 \). The remainder of the stick has length \( (1 - v_1) \). We then assign to \( \pi_1 \) a value of \( \theta_1 \) sampled from \( H \) and reiterate to get \( \pi_2 \), so the \( \pi_k \) are all sampled according to the stick-breaking construction. In practice, we terminate the stick breaking process when the remaining stick length falls below a prespecified cutoff. In the statistics literature, the resulting distribution for \( \pi \) is known as the Griffiths–Engen–McClosky (GEM) distribution.\textsuperscript{188}

The DP powerfully generalizes previously finite models, such as the hidden Markov model,\textsuperscript{67} ubiquitously used across single molecule biophysics\textsuperscript{189} that we discuss in greater depth in section 7.3.2.

4. THE LOCALIZATION PROBLEM

The localization problem is the first step in the analysis of a super-resolution data set and involves finding the position of a fluorescent molecule, \( x_0 = (x_0, y_0) \), from an image \( I \). The image itself is thought of as a matrix, whose elements describe individual intensities at each pixel.

In order to localize a fluorophore, we must have a model describing the expected mean number of photons per frame in pixel \( x \) given the fluorophore location at position \( x_0 = (x_0, y_0) \), \( \lambda(x; x_0) \). Typically, \( \lambda(x; x_0) \) is given by the point spread function of the imaging system. The intensity at each pixel at location \( x \) is itself distributed randomly, according to a distribution \( p(I(x)|\lambda(x; x_0)) \), due to shot noise and readout noise.

We begin by describing readout noise (section 4.1) and follow with a discussion on identifying “regions of interest” (ROIs) containing fluorophores (section 4.2). Once positively identified, we draw from our discussion on maximum likelihood in order to describe inference frameworks used in localization in section 4.3. While theoretically attractive, maximum likelihood methods may be computationally expensive and require good noise models to outperform simpler approaches. For this reason, we describe performance criteria of localization methods ultimately used to judge whether the computational cost of a method is warranted in section 4.4. Sections 4.5 and 4.6 describe simpler localization strategies, including least-squares fit. Subsequent sections tackle...
generalizations of the methods discussed thus far: 3D super-resolution in section 4.7, simultaneous fitting of multiple emitters in section 4.8, and deconvolution-style approaches in section 4.9. Finally, we end with a note on drift correction in super-resolution (section 4.10) without which the best localization methods are of limited value.

4.1. Readout Noise in Single Molecule Experiments

Intuitively, one can expect photon shot noise to be partly responsible for reducing the accuracy of localization methods. Indeed, localization must be achieved with few photons per frame as the total photon budget of most fluorophores, meaning the number of photons collected before the fluorophore undergoes irreversible photobleaching, is limited to hundreds or thousands of photons. While greater brightnesses can be achieved by using quantum dots as fluorescent markers, they remain more challenging to deliver into cells and present toxicity concerns.

Perhaps more unexpectedly, accurate localization also requires a model describing how a fluorophore’s emitted photons are converted into a camera readout. For instance, at a given illumination level, assuming an average number of photons strike the sample per frame per unit area, one may naively expect the camera’s readout at a given pixel, $I = I(x)$, to be a Gaussian random variable identical for all pixels, or at least well approximated by such a description. In fact, as we now discuss, both Gaussian and identical assumptions are violated in practice.

Since few photons hit each camera pixel on any given frame, the Poisson limit theorem states that given the average number of photons $\lambda$ for this pixel, the distribution of the actual number $N_p$ of such photons follows a Poisson distribution (“shot noise”)

$$p(N_p|\lambda) = \frac{\lambda^{N_p} e^{-\lambda}}{N_p!}$$  \hspace{1cm} (39)

where for notational simplicity, we let $\lambda = \lambda(x; x_0)$.

The total noise of the measurement arises from the convolution of this shot noise by a camera readout noise, that is neither necessarily normally distributed, nor pixel-independent. In other words, the readout $I$ at a camera pixel is distributed according to a distribution $p(I|N_p)$ that is non-normal and pixel-dependent. As later described, we will use both knowledge of $p(N_p|\lambda)$ and $p(I|N_p)$ to address the localization problem.

4.1.1. Camera Specific Readout—Two technologies, with different readout distributions, are widely used for single molecule imaging: the older EMCCD (electron-multiplication charge coupled device), where the electrons produced by a photon hitting a pixel are collected and amplified by chip-wide electronics, and the more recent sCMOS (scientific complementary metal oxide semiconductor), which offer higher sensitivity and read rates, at the cost of pixel-to-pixel noise variation (“fixed pattern noise”), by performing signal amplification at the pixel and column level.
The noise distribution of an EMCCD camera follows from its amplification mechanism where a photon hitting a pixel is converted into electrons. Chip-wide multiple charge-carrier multiplication (CCM) stages then amplify this electronic signal serially, one pixel at a time.

Specifically, each electron entering a stage has a low probability $p$ of giving rise to an output of two electrons; otherwise, no amplification takes place and a single electron is output with high probability $1-p$. Repeating this process across a large number of stages yield an exponentially distributed number of electrons arising from this single photon

$$p(I|N_p=1) \propto \exp\left(-\frac{I}{m}\right) \quad (40)$$

where the multiplication factor $m$ is itself weakly pixel-dependent, due to manufacturing imperfections. The distribution of the output from the amplification stage for $N_p$ photons simultaneously hitting a single pixel is the $N_p$-fold convolution of the one-photon distribution

$$p(I|N_p) \propto I^{N_p-1}\exp\left(-\frac{I}{m}\right) \quad (41)$$

After amplification, the electronic readout stage itself introduces both Gaussian noise, of standard deviation $\sigma$, which needs to be convolved to this distribution, and an offset in the number of counts (“dark count”), $c_0$, considered constant.

The other technology, sCMOS cameras, offer higher sensitivity and readout rates by attaching an individual amplification stage to each pixel. This different amplification technology yields the following normally distributed readout

$$p(I|N_p) \propto \exp\left(-\frac{(I - c_0 - mN_p)^2}{2\sigma^2}\right) \quad (42)$$

but the gain $m$, offset $c_0$, and variance $\sigma^2$ all vary (relatively) strongly from pixel to pixel (Figure 11).

From the distribution camera readouts for a given number of photons, $p(I|N_p)$, and the distribution of photon counts, $p(N_p|\lambda)$, we may compute the probability distribution of the camera readout $I$ given $\lambda$, by marginalizing over the unobservable number of photons $N_p$

$$p(I|\lambda) = \sum_{N_p \geq 0} p(I|N_p)p(N_p|\lambda) \quad (43)$$

As we will see in the next section, this distribution is essential for our goal of estimating $x_0$ (on which $\lambda$ depends).
As earlier mentioned, numerical estimation of this sum (which also matches experimental observations) demonstrates that \( p(I|\lambda) \) is highly skewed for EMCCD cameras,\(^{196} \) thus violating the normally distributed noise assumption (Figure 12). In the case of sCMOS cameras, numerical estimation of the sum in eq 43 also yields a non-normal distribution \( p(I|\lambda) \); moreover, and more importantly, this distribution changes from pixel to pixel due to the variability of \( m, c_0 \), and \( \sigma^2 \).\(^{198} \)

In principle, one may also infer \( m \) and \( c_0 \) directly from \( p(I|\lambda) \). With these two parameters at hand, and furthermore knowing, from Eqs (41) and (42), that the mean of \( I \) is a linear function of \( N_p \), we can obtain an estimate of \( N_p \) at each pixel given \( I \). This estimate is useful in evaluating the localization accuracy of the methods we will later explore in section 4.6.2. However, the central quantity, moving forward, is \( p(I|\lambda) \).

### 4.2. Detecting Single Molecules

We have previously described how the camera readout, \( I \), is related to the illumination level, \( \lambda \), through the distribution \( p(I|\lambda) \). As we will discuss in section 4.6, physical models of spatial localization allow us to estimate, for given fluorophore parameters \( \Theta \), the value of \( \lambda \) at each pixel \( x \), \( \lambda(x;\Theta) \). The fluorophore parameters \( \Theta \) minimally include the position of the fluorophore (as we had described earlier), but may also include its brightness,\(^{45} \) orientation,\(^{199} \) velocity,\(^{200} \) or other properties.

From \( p(I|\lambda) \) and \( \lambda(x;\Theta) \) we obtain a distribution of images conditioned on \( \Theta \)

\[
p(I|\Theta) = \prod_x p(I(x)|\lambda(x;\Theta)) \tag{44}
\]

where we have assumed that readout noise is uncorrelated across pixels.

We may, in principle, fit the entire image and simultaneously localize a large number of fluorophores. This is a difficult task, which we will address in section 4.8. Alternatively, we may crop out ROIs centered around “emission-like” patterns, as a prelude to their further analysis.\(^{45,201} \) Mathematically, this is equivalent to marginalizing over the positions outside of the ROI, i.e., ignoring the dependence of the image within the ROI on the positions outside of it. We explain this here with the caveat that, even today, the selection of these ROIs is often treated in an ad hoc manner, with limited theoretical justification.\(^ {202} \)

#### 4.2.1. Laplacian of Gaussian Filter

One may expect that ROIs could be chosen by locating pixels whose intensity go beyond some preset threshold. Such an approach cannot achieve high identification levels of relevant regions, in particular due to the presence of large amplitude and low spatial frequency background noise. Instead, a commonly used approach (and one of the few for which a theoretical basis has been offered) is to enhance features of a characteristic size \( \sigma \) (chosen to be that of a diffraction-limited spot) by convolution of the raw image \( I(x) = I(x, y) \) with a Laplacian of Gaussian kernel \( K(x, y) \).\(^ {203} \)
i.e., the convolved image \( I'(x) \) is

\[
I'(x) = \sum_{\delta x} K(\delta x) I(x - \delta x) = K \ast I
\]  

(46)

In this convolved image, features of a characteristic size \( \sigma \) have been “enhanced” to appear as sharp peaks. Peaks with a value greater than a given threshold can then be selected as originating from a single molecule and deserving further processing. This threshold is usually empirically chosen\(^{203} \) for example by picking as many peaks as possible while avoiding peaks that “look like” noise (as tested during the following processing stages).

However, if we have a good model of the background noise, we can also estimate (by simulation) the distribution of peak values that would be obtained from convolving an image only constituted of background noise, and then choose a threshold value that satisfies a user-specified false-positive \( p \)-value (that is, such that the probability of observing peaks with a value greater than the threshold in a convolved pure noise image is \( p \)).\(^{204} \)

Briefly, the theoretical justification for eq 45 relies on matched filter theory.\(^{205} \) Matched filter theory indicates that, if we are in the presence of additive white noise (i.e., if the differences between the observations and the true values constitute a random signal with constant spectral power density), the best linear filter to retrieve the original distribution is the convolution by the spatially reversed PSF itself (\( I'(x) = K(x) \text{ PSF}(-x) \)). In Fourier space, such a filter corresponds to multiplication by the conjugate of the Fourier transform of the PSF. Furthermore, empirical observations establish that the spectral power density (the square of the magnitude of the Fourier transform) of background fluorescence noise, not to be confused with camera readout noise, in an image approximately follows a power-law, \( |\mathcal{F}\{I\}(k)|^2 \propto |k|^{-s} \), with \( s \approx 2 \) (where \( \mathcal{F}\{\cdot\} \) denotes the Fourier transform).\(^{203} \) Thus, in order to apply the matched filter result, we first need to transform our data so that it exhibits white noise (whitening); this is done by multiplying the data, in Fourier space, by the filter \( \mathcal{H}(k) = |k|^{s/2} \) (so that \( |\mathcal{F}\{H \cdot I\}|^2 = 1 \)). The combination of both steps (whitening and convolution by the spatially reversed PSF) corresponds to the multiplication, in Fourier space, by the filter

\[
\mathcal{F}(K)(k) \propto |k|^{-s/2} \mathcal{F}(\text{PSF})(k)
\]  

(47)

(where the overbar indicates complex conjugation). In the case where \( s = 2 \) and the PSF is modeled as a Gaussian, eq 47 indeed corresponds to the Laplacian of Gaussian filter described in eq 47.\(^{203} \)
Since super-resolution data sets often contain many consecutive frames, additional improvements beyond whitening filters may be used. For instance, since background fluorescence varies slowly over time, it is possible to empirically decrease the influence of background fluorescence by working on difference images, that is, the difference in intensity between a frame and the next one.\textsuperscript{4} The switching-on of a molecule then appears as a positive peak in the data, whereas its disappearance is a negative peak, both of which may be further selected using the whitened matched filter (eq 45).\textsuperscript{4}

4.2.2. Errors in Emitter Identification—The output of this initial analysis is a list of ROIs, where a single molecule is assumed to have been fluorescent. Metrics, which we now introduce, can be used to quantify the quality of the list. For any method, such metrics are typically calculated from synthetic data, where the ground truth is a priori known, which is not the case with real data. Therefore, the metrics provide only an estimate of the method result quality. If the data treated is substantially different than the synthetic data the metrics were calculated on, this estimate may be quite inaccurate.

These metrics below are expressed in terms of two kinds of possible errors: some molecules may have been missed by the detection algorithm (false negatives, FN), and some regions of interest may have mistakenly been drawn somewhere where there was, in fact, no molecule (false positives, FP).\textsuperscript{206}

If we denote TP the number of true positives (correctly drawn regions of interest), two fundamental measures of accuracy are possible: the precision (quantifying false positives)

\[
P = \frac{TP}{FP + TP} \quad (48)
\]

and the recall (quantifying false negatives)

\[
r = \frac{TP}{FN + TP} \quad (49)
\]

In order to directly rank different methods, it is convenient to combine these two measures into a single quantity. Such quantities include the Jaccard index\textsuperscript{202,206}

\[
JAC = \frac{TP}{FN + FP + TP} \quad (50)
\]

or the $F_1$-score (or $F$-measure)\textsuperscript{207}

\[
F_1 = \frac{2}{1/p + 1/r} \quad (51)
\]
Modern localization methods are typically able to achieve high precision ($p \gtrsim 95\%$) while still having limited, though widely varying, recalls ($r \approx 25\% \text{ to } 75\%$); this latter value thus also limits the achievable Jaccard index and $F_1$-score.

4.3. Maximum Likelihood Localization

Having segmented our image into regions and identified whether such regions contain a single molecule, we now turn to the problem of localization within an ROI using maximum likelihood, described earlier in section 3.2.1.

Specifically, the MLE given by eq 4 finds parameters $\Theta$ maximizing the probability of observing $\mathbf{I}$ over the $N \times N$ pixels

$$\hat{\Theta} = \arg\max_{\Theta} p(\mathbf{I}|\Theta)$$

(52)

where the probability, $p(\mathbf{I}|\Theta)$, above is understood as the product over each pixel $\mathbf{x}$ of the probability of observing the actual pixel value $I(\mathbf{x})$, expressed as a function of the fluorophore parameters $\Theta$, i.e.

$$p(\mathbf{I}|\Theta) = \prod_{\mathbf{x}} p(I(\mathbf{x})|\lambda(\mathbf{x};\Theta))$$

(53)

where we have adopted the common assumption of pixel independence.

The maximization of eq 52 can be carried using out-of-the-box numerical approaches, such as gradient descent; practical implementations of such a method in a super-resolution context (which achieve the CRLB) are available for both EMCCD and sCMOS cameras.\cite{55,194,208}

The actual value of the CRLB (5 to 50 nm) depends strongly on a number of experimental parameters, most importantly the number of photons that can actually be collected.\cite{55,194,208}

Despite the theoretical optimality of the MLE (in the CRLB, or mean-squared error, sense), the necessarily imperfect knowledge we have about the imaging system (background fluorescence, the PSF, the camera noise) reduces its performance. In fact, PSF mis-specification or imperfections degrade the performance of the method and may even lead to overly optimistic accuracy estimations.\cite{55,121,209} It thus remains useful to study simpler approaches, which can take advantage of empirical corrections.

4.4. Additional Super-Resolution Performance Metrics

While the Jaccard index, eq 50, and the mean-square error of a single molecule’s localization are good performance metrics, even perfect localization cannot reconstruct a biological structure that is poorly labeled.\cite{4}

To assess the quality of a reconstruction, Fourier ring correlation (FRC), a method originally developed for cryoelectron microscopy, is employed.\cite{210,211} Briefly, in this method, the
collected single molecule events are randomly split into two data sets, which are used to create two independent reconstructions \( I_1 \) and \( I_2 \) of the structure. The “consistency” between these two reconstructions is then used as a quantification of their resolution.\(^{210,211}\) This consistency is obtained, as the name implies, by computing the Fourier transforms, \( \mathcal{F}\{I_1\}(k) \) and \( \mathcal{F}\{I_2\}(k) \) of the images, and computing the normalized correlation between “rings” of constant spatial frequency magnitude \(|k| = k\),

\[
\text{FRC}(k) = \frac{\sum |k|=k \mathcal{F}\{I_1\}(k) \overline{\mathcal{F}\{I_2\}(k)}}{(\sum |k|=k |\mathcal{F}\{I_1\}(k)|^2 \sum |k|=k |\mathcal{F}\{I_2\}(k)|^2)^{1/2}}
\]

(54)

where \( \overline{\mathcal{F}} \) denotes the complex conjugate of \( \mathcal{F} \).

This formula yields, for each magnitude of spatial frequency, the degree of correlation, normalized between \(-1\) and \(+1\), to which the features of that characteristic size are correlated between the two independent reconstructions. In fact, it is this separation of length scales that motivates the use of correlation in Fourier space.

For relatively large sized structures, using a random half of the events does not greatly diminish the quality of the reconstruction; thus, the two reconstructions should be highly correlated. Conversely, for structures too small to be well resolved, there is no reason to expect the two reconstructions to be highly correlated and, consequently, the FRC should be smaller.

We may then select a conventional threshold FRC (typically, FRC(\(q\)) = 1/7) and report as “the resolution” the corresponding characteristic size beyond which the threshold is exceeded.\(^{210,211}\) Interestingly, this measure tends to indicate that nowadays, the main factor limiting the resolution of reconstructed static structures is typically the labeling density rather than the accuracy of the single molecule localization itself.\(^{210,211}\)

4.5. Simplified Localization Approaches

We have seen that while maximum likelihood localization is theoretically the method that achieves the lowest mean-squared error, imperfect knowledge of the imaging system characteristics may make other localization methods preferable. Additionally, maximum likelihood calculations are typically computationally expensive and implementations often run on specialized hardware such as graphical processing units (GPUs).\(^{55,198}\) Thus, it remains useful to study simpler, possibly less model-dependent, approaches.

4.5.1. Centroid Localization Method—An intuitive, simple, and extremely fast approach to the localization problem is to compute the average of the pixel coordinates \( x = (x, y) \) within a ROI, weighted by their intensities \( I(x) \).\(^{201}\)

In such a method, it is crucially important to first subtract away any background fluorescence \( I_b \) from the ROI,\(^{212}\) such that the estimated localization is
where the sum is over the region of interest. Background subtraction is important because in its absence, the weighted average eq 55 becomes a weighted average between the true centroid and the ROIs geometric center.

However, even with this correction, the method remains unsuitable for high-resolution localization.\textsuperscript{212} One simple reason is that, even under the reasonable assumption that the physical PSF is symmetric (and thus its centroid should yield the fluorophore position), this is not the case for the camera readout, which is measured on a discrete pixel grid. Even worse, the centroid of the camera readout does not necessarily coincide with the centroid of the physical PSF (again due to pixelation).\textsuperscript{212} Still, the extreme simplicity of the method has led to its use as a minimal baseline against which other approaches can be compared.\textsuperscript{202}

\subsection*{4.5.2. Finding the Point of Radial Symmetry—}
The centroid method we just described attempts to localize an event with subpixel resolution by identifying its “geometric center”. Other definitions of “geometric center” have been proposed, most notably the radial symmetry approach.\textsuperscript{213,214} Briefly, this approach attempts to find the point that best approximates a “radial center of symmetry” for the image.

In this method, the gradient of the signal is calculated either at each pixel\textsuperscript{215} or at each point where four adjacent pixels (or, in the 3D case,\textsuperscript{214} eight adjacent voxels) meet.\textsuperscript{213,215} The line defined by this point and gradient pair is taken as approximating a local axis of symmetry for the image. If all such lines were to intersect with each other at a single point, such a point would be a reasonable definition of the radial symmetry center. Because this is not the case, the radial symmetry center is instead defined as the point that minimizes its total distance to all such lines,\textsuperscript{215} possibly with an appropriate weighting factor.\textsuperscript{213,214}

Specifically, it is reasonable to weight lines inversely proportionately to their distance from the center of the image. Since this center is yet unknown, the weighting is instead done using the inverse distance to the centroid (as computed above).\textsuperscript{213,214} Most importantly, an analytical expression can be derived to compute the radial symmetry center thus defined;\textsuperscript{213,214} as such, this method is extremely rapid.

While simulations indicate that this method yields high, close to CRLB-level localization accuracy of single events at a high speed,\textsuperscript{202,213,215} they also show that its performance degrades extremely quickly for high-density data, being unable to correctly localize events that were not well separated from the others.\textsuperscript{202}

\subsection*{4.5.3. Correlation—}
As discussed earlier, the good performance of the Laplacian of Gaussian kernel for event detection was justified on the basis of simple noise and PSF models (section 4.2). We now extend this approach to tackle the localization problem itself.
In this approach, a peak’s position is determined by computing the correlation between the image and the model PSF (although using the Laplacian of the PSF may work better from a theoretical point of view, as discussed above, it is the PSF itself that is typically used), and finding the position at which this correlation is maximal. The same background removal approaches as for centroid calculations may be used; however, they are less important, as adding a constant background to the image simply shifts the filtered image by a constant and thus does not affect the maximum’s position.

The correlation of two images is only defined for integer coordinates, so additional work is needed to obtain a subpixel localization. A simple way to do so is to fit the values of the correlation in the vicinity of the maximum with a continuous, peaked model function (e.g., a parabola) and use the maximum of the latter. A more sophisticated approach is to compute this correlation after Fourier-resampling both the image and the model PSF to a higher resolution. Such a resampling is achieved by taking the Fourier transform of the image, zero-padding it to a higher spatial frequency, and taking the Fourier transform back. Correlation in real space corresponds to point-wise product in Fourier space; thus, the desired procedure amounts to computing the point-wise product of the Fourier transforms of the image and the PSF, zero-pad it, Fourier transform the padded product back into real space, and then select those coordinates at which the correlation attains its maximum.

To sidestep the computational cost of the Fourier transforms that upsampling requires, such methods are typically first run with a limited upsampling to yield a low resolution localization and then run again with higher upsampling but only in a small neighborhood around the position selected by the first iteration.

An important advantage of correlation-based localization is that it can be directly used for any experimentally measured PSF. For example, in particle tracking (an early application of subpixel localization), one can use the image of a molecule in one frame as the model PSF for the next frame. In superresolution experiments, this approach has been suggested to analyze thick-sample data, which typically exhibits highly distorted PSFs in the absence of specialized optical corrections. In this case, the distorted, sample-specific PSF is measured at the beginning of the experiment by imaging a point-source at different depths; detected events are then localized by correlation with this PSF.

4.6. Least-Squares Fitting and Model PSFs

The previous section covered methods that require limited assumptions regarding the PSF; for instance, that it be radially symmetric or invariant across the data set. Here instead we focus on an approach, least-squares fitting, that demands no such assumptions but that does require a form of the PSF.

While, in theory, maximum likelihood achieves the optimal mean square error when an accurate PSF model is available (section 4.3), the least-squares method is widely used because of the good performance of readily available, fast, and robust algorithms. Although we will first focus on the common case of fitting a Gaussian PSF model, we will then discuss possible corrections to this model.
4.6.1. Gaussian PSF Least-Squares Fitting—Since the theory of least-squares fitting, as with maximum likelihood (section 4.3), can be described independently of the model PSF’s exact functional form, we will, for simplicity, assume a Gaussian PSF. This choice is one of the earliest in use, offers mathematical simplicity, and maintains good performance.

Specifically, we model the image $I_0$ arising from a fluorophore as a two-dimensional Gaussian as follows

$$I_0(x; A, x_0, \sigma, I_b) = I_b + A \exp \left( -\frac{|x - x_0|^2}{2\sigma^2} \right)$$  \hspace{1cm} (56)

The unknown amplitude $A$, center $x_0$ and standard deviation $\sigma$, as well as the unknown, locally constant mean background $I_b$, are parameters collectively regrouped as $\Theta$, the fluorophore characteristics, that we now want to infer. Furthermore, despite the subtraction of the background $I_b$, the measured image still differs from the model by a noise term of mean zero. Next, it is also possible to assume that some parameters are a priori known such as $\sigma$ or $I_b$, for example, they may be independently estimated from the image intensity far away from the fluorophore.\textsuperscript{218} It is also possible to improve this model by averaging the PSF over each pixel.\textsuperscript{3}

One may then infer the remaining set of parameters minimizing the sum of squared differences between the observed intensity and model provided by eq 56, weighted by the signal variance at each pixel. Numerically, this is a classic least-squares minimization, for which fast and robust implementations, such as the Levenberg-Marquart algorithm,\textsuperscript{193} are available.

The maximum likelihood framework (section 4.3) and least-squares fitting are identical, even for non-normal PSFs, if the noise at each pixel is assumed to be independent and drawn from the same normal distribution with unspecified variance

$$l \propto -\sum_x (I(x) - I_0(x; \Theta))^2 + \text{constant}$$  \hspace{1cm} (57)

where $l$ denotes the log-likelihood.

4.6.2. Least-Squares Fitting Localization Accuracy—Thompson et al. provided a theoretical analysis of least-squares fitting accuracy in the presence of normally distributed background noise as well as photon counting noise (section 4.1), as a function of the PSF’s standard deviation (the “spot size”) $s$, the pixel size $a$, the number of photons in the event $N_p$, and the standard deviation of the background noise $b$ (Figure 13).\textsuperscript{53}

For simplicity, we limit ourselves to rederiving Thompson’s results in the case of a 1D model and assume that for a fluorophore at position $x$: (i) the expected number of photons at the $i$th pixel (i.e., the PSF model) is $N_i(x)$; (ii) the variance is $\sigma_i^2 = N_i(x) + b^2$ (i.e., the sum, in quadrature, of the photon counting noise, $N_i(x)$, and the background noise, $b^2$); and (iii) the
detected photon number at that same pixel is $y_i$. By definition, the fitted position, $\hat{x}$, is obtained by minimizing the weighted sum of square residuals, i.e.

$$\frac{\partial}{\partial x} \sum_i \frac{(y_i - N_i(x))^2}{\sigma_i^2} \bigg|_{x=\hat{x}} = 0$$  \hspace{1cm} (58)$$

By expanding $N_i(x)$ in the above to first order in $\Delta x$ around the true underlying position $x_0(\hat{x}(x_0) \approx N_i(x_0)N_i'(x_0)(\hat{x} - x_0))$ and solving for $\Delta x = \hat{x} - x_0$, we directly derive the mean square error of the fitted center’s position

$$\langle (\Delta x)^2 \rangle = \frac{1}{\sum_i (N_i^2/(N_i+b^2))}$$  \hspace{1cm} (59)$$

While this sum can be evaluated numerically, we can also simplify it under reasonable approximations. We ignore, for now, the effects of pixelation ($a \rightarrow 0$). In this case, under a Gaussian PSF model, the expected number of photons at pixel $i$ is $N_i = (N_p/\sqrt{2\pi}\sigma)\exp(-r^2/2\sigma^2)$, and the sum in eq 59 can be replaced by an integral.

In general, this integral is not analytically tractable but it can be asymptotically evaluated in two limits: (i) dominant photon-counting noise ($N_i \gg b^2$, so $\sigma_i^2 \approx N_i$) and (ii) dominant background noise ($N_i \ll b^2$, so $\sigma_i^2 \approx b^2$). These two cases respectively yield

$$\langle (\Delta x)^2 \rangle_1 = \frac{s^2}{N_p}$$  \hspace{1cm} (60)$$

and

$$\langle (\Delta x)^2 \rangle_2 = \frac{8\pi s^4 b^2}{a^2 N_p^2}$$  \hspace{1cm} (61)$$

Since each expression dominates the other in the limit where it has been derived, the authors suggested the following interpolation formula

$$\langle (\Delta x)^2 \rangle = \frac{s^2}{N_p} + \frac{8\pi s^4 b^2}{a^2 N_p^2}$$  \hspace{1cm} (62)$$

The pixelation noise’s main effect (arising from a nonzero $a$) is to increase the photon counting noise term $\langle (\Delta x)^2 \rangle_1$. Specifically, the PSF’s spatial variance, $s^2$, appearing in this term should be increased by the spatial variance of a square pixel of size $a$, which is $a^2/12$. The final expression for the uncertainty of Gaussian fitting is thus
Although eq 63 is widely used to report localization accuracies,\(^4,5,219\) the summation in eq 59 can also be evaluated numerically.\(^53\) This numerical estimate indicates that eq 63 actually overestimates the localization accuracy (in the relevant regime of parameters) by approximately 10%.\(^53\) This is a discrepancy that has also been reported from experimental comparisons.\(^3\)

An interesting consequence of eq 59 is that the mean square error is minimal for a nonzero pixel size \(a (\partial (\Delta x)^2/\partial a = 0)\). In other words, it is counter-productive to make the pixel size as small as possible. Instead, its optimal size is close to the spot size \(s\). Intuitively, this is due to the compromise between the higher spatial information gained from each pixel when the pixels are smaller and the averaging out of background noise when the pixels are larger.

In practice, Gaussian PSF fitting has been shown to achieve nanometer-resolution. For example, Yildiz et al. have used this approach to show that the motion of fluorescently labeled myosin V enzymes along their tracks occurs in steps of variable size that can be grouped in consecutive pairs whose sizes add up to 74 nm (fluorescence imaging with one-nanometer accuracy, FIONA).\(^219\)

### 4.6.3. Applicability of Least-Squares to Non-Normal Noise

While the assumption of identically and normally distributed noise is reasonable in many applications of least-squares fitting, which is the source of its versatility and the reason many efficient algorithms have been developed, it is clearly violated in super-resolution, as described in section 4.1.

Although many super-resolution analysis discount non-normal noise, here we discuss a variance-stabilizing transformation\(^204\) that mitigates the effect of ignoring non-normal noise.

For simplicity, we consider only the effect of Poisson (shot) noise, whose variance is equal to its mean. Since the variance of the noise changes across the fitted ROI, the assumption of identical noise distribution is violated.

In order to correct for this nonuniformity, we exploit the following (numerical) observation, known as the Anscombe transform: if \(X\) is Poisson-distributed with both mean and variance equal to \(m \geq 4\), then \(2 \sqrt{X + 3/8}\) is approximately normally distributed with mean \(2 \sqrt{m + 3/8} - 1/4 \sqrt{m}\) and, more importantly, unit variance.\(^220\) Thus, applying this transformation to an image corrupted by Poisson noise yields an image with (approximately) uniform Gaussian noise and the classical least-squares algorithm may then be applied. Of course, the fit should not be done using the original PSF model but, likewise, the Anscombe-transformed model.\(^204\)

Note that this correction assumes that the image data is correctly expressed in units of photon counts, which requires a calibration of the readout-to-photons conversion factor as
discussed in section 4.1. Additionally, more sophisticated transforms (e.g., the generalized Anscombe transform\textsuperscript{221}) may be used to handle more realistic non-Poisson noise models.

While, to our knowledge, the effect of a variance-stabilizing transformation for the accuracy of least-squares fitting has not been evaluated independently of other improvements, the SimpleSTORM package, which relies on it as a preprocessing step before least-squares fitting,\textsuperscript{204} was shown to exhibit a relatively strong performance.\textsuperscript{202,204}

4.6.4. Corrections to the Point Spread Function—Although we have mentioned, and it is widely quoted,\textsuperscript{53,201} that the diffraction pattern of a point source is an Airy disc (section 2.2), and chose to approximate this pattern with a Gaussian peak both for maximum likelihood and for least-squares fitting, we now revisit this claim.

When imaging using a high-NA objective, as commonly done in super-resolution applications, the PSF of a freely rotating fluorophore, directly derived from first-principles, is in fact closer to a Gaussian function than to an Airy function\textsuperscript{222} thus justifying, a posteriori, the use of Gaussians for least-squares fitting.

A rotationally constrained fluorophore, which may occur, or conversely be avoided, e.g., due to the labeling strategy used,\textsuperscript{199,223} presents additional complications. Such a constraint breaks radial symmetry, in which case the PSF may present two “lobes”.\textsuperscript{199,222} If a rotationally free model, such as a Gaussian, is used to fit data sets lacking radial symmetry, simulations indicate that maximum likelihood estimation can lead to substantial errors (dozens of nanometers), in particular in the case of defocused molecules (e.g., for 3D measurements).\textsuperscript{222,224} Conversely, orientational information may be derived from properly fitting the observed PSF to a model PSF for anisotropic emission.\textsuperscript{199}

In the opposite extreme, highly mobile fluorophores, which move by a significant fraction of a pixel size during the time it takes to acquire a single frame,\textsuperscript{200} may distort the effective molecular PSF, which is now a weighted average of the PSF at each position visited by the molecule. Once more, ignoring this distortion leads to poor localization accuracy, whereas using a PSF model that takes motion into account not only restores the original localization accuracy but also provides information on the instantaneous molecular velocity\textsuperscript{200} and additional information on motion models discussed in section 6.5.

We end with a note on the nonuniformity of a sample’s refraction index which introduces additional PSF aberrations, especially for thick samples.\textsuperscript{36} This effect is has so far been treated experimentally by using adaptive optics (e.g., deformable mirrors) to properly shape the PSF.\textsuperscript{36}

4.7. 3D Localization

4.7.1. Cylindrical Lens 3D—While our discussion, so far, has been limited to localizing single molecules in a 2D plane, most biological samples are three-dimensional and, as a consequence, there is considerable interest in obtaining volumetric fluorescence data.

In classical microscopy, this can be achieved by selectively exciting, and thus collecting, fluorescence from a single plane (multiphoton microscopy\textsuperscript{225} or selective plane illumination
microscopy (SPIM)\textsuperscript{226}). However, such techniques remain essentially limited by diffraction. Instead, true 3D super-resolution can be achieved by encoding information about the depth of a molecule in its PSF.

Fundamentally, the techniques we have discussed up until now fit a PSF that encodes lateral but not vertical information. In other words, in 2D, the value of the PSF measured by the camera at position \((x, y)\) when the emitter is at position \((x_0, y_0)\) depends only on the distance between the two positions, i.e., \(\text{PSF} = \text{PSF}(x-x_0, y-y_0)\). In 3D, the dependence on the true position \(z_0\) cannot be expressed in terms of translation and the PSF would need to be of the form \(\text{PSF}(x-x_0, y-y_0, z_0)\).

As early as in 1994, Kao et al. introduced a cylindrical lens in the optical path of their particle tracking setup and observed a depth-dependent PSF\textsuperscript{227} This depth-dependent PSF progressively switches from being a vertically oriented ellipse for molecules above the focal plane to a horizontally oriented one on the other side. Thus, the lengths of the PSF’s two axes, \(w_x\) and \(w_y\), could be estimated and converted to a depth value using a calibration table. Specifically, the relative difference between the two widths, defined as \(R = (w_y-w_x)/(w_x+w_y)\), was matched with a calibration curves \(R^\text{cal}(z)\) in order to read out the depth \(z\) while the \((x, y)\) position was obtained by least-squares fitting to a parabolic PSF (section 4.6).

The cylindrical lens approach was adapted for super-resolution by Huang et al.\textsuperscript{228} who took advantage of the advent of more general nonlinear fitting procedures, allowing the determination of \(w_x\) and \(w_y\) by least-squares fit along with the in-plane position. That is, the model PSF was chosen as a Gaussian with the following parameters that need to be fitted: the position of the center and the amplitude of the PSF, similarly to the two-dimensional case along with the PSF width and height \(w_x\) and \(w_y\) treated as independent parameters. The authors found, purely empirically, that the fitted \(w_x\) and \(w_y\) could be accurately mapped back to the molecule depth \(z\) via the use of a calibration curve \((w_x^{\text{cal}}(z), w_y^{\text{cal}}(z))\) obtained by measuring the PSF of point sources positioned at different depths, as follows: the depth \(z\) is chosen to minimize the Euclidean distance between the \((w_x^{1/2}, w_y^{1/2})\) point and the \((w_x^{\text{cal}}(z)^{1/2}, w_y^{\text{cal}}(z)^{1/2})\) curve.

Instead of estimating depth based on ellipticity calibration curves, we may immediately adapt all methods described for 2D localization to the 3D case by simply including the depth \(z\) in the set of parameters \(\Theta\).\textsuperscript{217,229} All theoretical results regarding such methods, such as the CRLB accuracy limit (section 4.3), are then applicable. For example, we demonstrated earlier that in the presence of a highly distorted, but experimentally well characterized PSFs, the position of the maximum in the ROI’s correlation with the PSF generated good localization performance (section 4.5.3). This method is, in fact, especially applicable to 3D imaging of thick samples, as the PSF of events localized deep into the cell can be distorted by severe optical aberrations.\textsuperscript{217}

4.7.2. Other Approaches for Encoding Depth Information in the PSF—While the cylindrical lens approach is relatively simple from an experimental viewpoint, it only
requires introducing a cylindrical lens in the optical path, it encodes depth information at the cost of lateral resolution, as it distorts the PSF. Additionally, as discussed earlier in section 4.6.4, other phenomena can lead to elliptical distortion of the PSF, leading to spurious apparent changes in depth. Hence, additional ways to encode depth information have been proposed.\textsuperscript{229}

For example, the biplane-PALM approach relies on simultaneously imaging two planes, a few hundred nanometers from each other on the same camera.\textsuperscript{229} This can be achieved by imaging on one-half of the camera chip the “standard” focused image and, on the other half, a slightly defocused image, obtained by splitting the collected light and reprojecting it onto the camera after a longer light path. A ROI corresponding to a single event now coincides with a pair of spots, one on each plane, that may once more be fitted by least-squares either to an experimental PSF, also measured over the two planes, or a theoretically derived one.\textsuperscript{229} As a fluorophore is displaced along the $z$ axis, it does not get defocused to the same degree in the two planes; this difference in defocusing thus encodes the depth information. In its first implementation, a depth resolution of 75 nm was achieved.\textsuperscript{229}

Additional $z$-resolution can be provided by more sophisticated procedures. For example, a spatial light modulator can be used to shape the 3D PSF into a double-helix, such that individual events are now observed as pairs of close peaks, whose relative position encode depth information.\textsuperscript{10} This technique, to which all the previous fitting discussions apply, exhibits a low theoretical maximal resolution (CRLB) of approximately 15 nm.\textsuperscript{230} Interferometric PALM (iPALM) provides an even more sophisticated procedure to encode depth information,\textsuperscript{231} in which the measured image is split over three cameras, each of which measure an interference pattern between two images that are phase-shifted with respect to one another. The relative intensities of a same peak across all three cameras allow the experimenter to compute this phase shift and thus infer the event depth, with an experimentally demonstrated resolution of approximately 10 nm.\textsuperscript{231}

### 4.8. Simultaneous Localization of Multiple Molecules

The fundamental breakthrough from which super-resolution microscopy emerged, namely achieving temporal separation of events too close to be resolved spatially, is also an important limitation. As described so far, a super-resolution acquisition scheme must ensure that only a few molecules are activated per frame, thus imposing lengthy acquisition times for densely labeled samples.

However, just as we have described various ways in which the coordinates of a single molecule can be retrieved if a model PSF is known, we could, in theory, write down an emission model for two, or more, close molecules with overlapping PSFs (given their coordinates $x_0$ and $x_1$), and then fit a ROI to such a model. This approach was pioneered by astronomers who were interested in separating images of stars in “crowded fields” (e.g., stars in distant galaxies, which appear very close to each other) and have since long ago developed such algorithms.\textsuperscript{232} One of these algorithms, DAOPHOT (Dominion Astrophysical Observatory photometry),\textsuperscript{232} was directly adapted for super-resolution microscopy, under the name of DAOSTORM.\textsuperscript{37}
There are a few difficulties that are associated with the simultaneous fitting of multiple molecules at a time. The first is computational; the greater the molecules simultaneously fit, the greater the number of parameters, rendering the least-squares or maximum likelihood optimization more challenging numerically.

Fortunately, it is clear that even when PSFs are slightly overlapping, it remains acceptable to cut the image into smaller regions, that are approximately statistically independent from each other, and fit them one at a time. This approximation was used by another super-resolution package developed at the same time, MFA (multiple-emitter fitting analysis). More accurately, one can also draw such regions to be bounded by areas of the image where the intensity is relatively low and are thus unlikely to contain a molecule (the approach of DAOPHOT/DAOSTORM). Again, in such cases, the problem of fitting PSFs in a region becomes independent from the fitting in another region, in a manner similar to how we drew ROIs for single-emitter fitting but this time with multiple fluorophores per ROI.

More importantly, simultaneous fitting of many fluorophores also presents a model selection problem: allowing for more fluorophores always result in a better (or at least, not worse) fit of a collection of spots (either the fitting algorithm can exploit the additional degrees of freedom to eliminate some residuals of the fit or, at worst, it can always set the brightness of the additional fluorophores to a very small value, thus not worsening the fit). Thus, additional criteria are necessary to prevent overfitting.

While general model selection methods were presented in section 3.3.4, here we present two more model selection strategies specifically adapted to the problem of multiemitter fitting used by DAOSTORM and by MFA.

DAOSTORM first uses a peak detection algorithm (such as the one discussed in section 4.2.1) in order to find candidate regions that may correspond to a molecule. This set of candidates is then fit, by MLE, to a multiemitter model. The residuals of the fit (i.e., the difference between the original image and the one that a set of fluorophores at positions given by the fit would yield) are then iteratively reinserted into the original peak detection algorithm. Thus, it is the sensitivity of the peak detection algorithm that provides a stopping criterion against the addition of extraneous fluorophores to the fit.

Model selection by MFA relies instead on computing the log-likelihood ratio, LLR:

\[
\text{LLR} = -2 \log \left( \frac{\mathcal{L}({\Theta}_\text{MLE}|I)}{\mathcal{L}(I/I)} \right)
\]

The numerator, \(\mathcal{L}({\Theta}_\text{MLE}|I)\), is the likelihood of the estimates given the image, assuming that each pixel’s signal is independently obtained from a Poisson-distributed source with mean equal to the sum of the PSFs at this pixel (eq 53). The denominator, \(\mathcal{L}(I/I)\), is the maximum value that the above-mentioned likelihood could ever attain, which it does in the case where the expected mean intensity at each pixel matches the actually observed intensity. In other words, it is the product over the pixels of the probability of observing the
actual camera output if the mean expected intensity at that pixel was set to be equal to that output.

Having evaluated the “goodness” of each model (as measured by its LLR), we now need to estimate, for each model, how well the model matches the data, as compared to how well it would match random data sets generated from the model itself. Such a comparison penalizes overfitting, as the marginal improvement to the LLR, for each additional parameter, decreases sharply once the “correct” number of parameters is reached, whereas such a transition does not occur for random data sets.

More specifically, we need to estimate the probability $p$ that the LLR of a data set generated from the model be lower than the LLR of the real data. In other words, we need the value of the cumulative distribution for the LLR, evaluated at the LLR of the real data. According to Wilks’ theorem, this distribution can be approximated by a $\chi^2$ distribution with a number of degrees of freedom equal to the difference between the number of pixels and the number of fitting parameters.$^{234}$ We thus obtain $p$ simply by evaluating the cumulative distribution function of the above-mentioned $\chi^2$ distribution.$^{233}$ Having done so for each of the models in contention, the model with the highest such probability is then selected.

### 4.9. Deconvolution-Based Super-Resolution

We have so far focused on reconstructing coordinates of each single event with subdiffraction accuracy. However, subdiffraction imaging may be achieved by other means. For example, deconvolution microscopy achieves a 2-fold improvement over diffraction-limited microscopy by approximating the inverse (in a linear operator sense) of the “imaging operator”, i.e., the operator that convolves a distribution of point emitters by the imaging system’s PSF.$^{42}$

Here, we discuss adaptations of deconvolution-style approaches to data sets collected using single molecule localization-style techniques where additional information is encoded in temporal fluctuations of the fluorescence (i.e., stochastic switching of the fluorophores).

#### 4.9.1. Compressed Sensing—

Contrary to localization methods discussed thus far, here we do not initially attempt to reconstruct a list of molecular positions. Instead, we want to reconstruct a higher resolution image than the one from which we started.

More specifically, we seek a “fluorophore density map” on a discrete grid $s$, where each “pixel” on the grid may be smaller than the raw image, $I$, physical pixels. Instead of considering $s$ and $I$ as matrices, we will consider them as vectors of entries (for example, by concatenating the physical columns of pixels in the image), respectively of size $N$ and $n$. In this formalism, convolution by the PSF, which is a linear operator, can be understood as multiplication by a matrix $A$, of size $(N, n)$,

$$I \approx A \cdot s \quad (65)$$
Each row of the matrix $\mathbf{A}$ corresponds to a possible fluorophore position and each entry in the row corresponds to a physical pixel indicating how much a fluorophore at the row-encoded position would increase the intensity at that physical pixel.

Localization methods discussed so far correspond approximately to a setup where we know (or have a good model of) $\mathbf{A}$ (i.e., how much a fluorophore at any position affects the intensity measured at any position—in other words, the PSF) and seek to obtain $\mathbf{s}$ (i.e., the fluorophore density map). We will focus on the same formulation first. However, we will later see that we can also attempt to recover $\mathbf{A}$ and $\mathbf{s}$ simultaneously.

The usual caveats of deconvolution microscopy, namely the sensitivity of $\mathbf{s}$ to noise and to inaccurate knowledge of $\mathbf{A}$ still apply. Moreover, as there are many more (discretized) fluorophore positions than image pixels ($N \gg n$), the problem is underdetermined. However, in the context of a super-resolution data set, we have the additional information that we expect only a few fluorescent proteins to be “on” on each frame; that is, we have a sparsity prior on $\mathbf{s}$ (we expect most of its entries to be zero).

This class of problems (searching for approximate and sparse solutions to an underdetermined linear system) is known as compressed sensing and is well described in the mathematical literature.$^{235}$ For example, Zhu et al. showed that in the presence of photon counting noise, a solution can be obtained by searching for the vector $\mathbf{s}$ with minimum $\ell_1$ norm (i.e., sum of absolute values of components) among all those for which the $\ell_2$ norm of the residual vector, $\mathbf{I} - \mathbf{A} \cdot \mathbf{s}$ (i.e., sum of squared errors), is no larger than a noise-level dependent threshold.$^{236}$ Such a vector can then be found using standard algorithms.$^{237}$

Such a deconvolution yields, for each frame of the image stack, a sparse list of discretized molecular positions. All such lists can then be merged together to obtain a final list of molecular positions. Although the original implementation of this idea$^{236}$ yielded a relatively poor localization accuracy of ~60 nm, it was able to recover highly overlapping events, i.e. it allowed a very fast imaging rate (6 to 15-fold faster than for single-event fitting, 2 to 3-fold faster than for a multiemitter fitting such as DAOSTORM).

4.9.2. Exploiting Fluorophore Temporal Fluctuations—Instead of using an experimental protocol designed to achieve temporally sparse photoactivation of fluorophores, it is also possible to rely on the natural blinking and bleaching of fluorophores, that occurs (to varying degrees depending on the fluorophore) even under continuous illumination.

For example, a simple way to exploit the fluorophore blinking is to compute the difference between consecutive frames of a regular fluorescence movie. In these difference images, the spontaneous switching-on of a fluorophore appears as a positive peak, whereas turning-off events, or photobleaching, appear as a negative one. Standard localization algorithms (e.g., PSF fitting) can then be directly applied to such images, yielding a super-resolution approach that does not require the use of photoconvertible markers (bleaching/blinking-assisted localization microscopy, BaLM).$^{238}$ More interestingly, it is possible to exploit the fact that these temporal fluctuations in fluorescence intensity are uncorrelated between
molecules (as each fluorophore undergoes stochastic switching independently from the others).

Dertinger et al. noticed that due to this stochasticity, pixels where the emission of two blinking fluorophores (quantum dots, in their case) overlap exhibit lower temporal coherence than pixels which capture the emission of a single quantum dot.\textsuperscript{239} This observation yields a simple and elegant method, named super-resolution optical fluctuation imaging (SOFI), to obtain a superresolved image \( \mathbf{I} \).\textsuperscript{239} At each pixel, one simply plots the value of the temporal correlation of this pixel’s signal for a well-chosen time lag \( \tau \)

\[
I_\tau(x) = \langle I(x,t)I(x,t+\tau) \rangle_t \quad (66)
\]

where \( \langle \cdot \rangle_t \) denotes an average over time.

Lidke et al. exploited temporal fluctuations in order to generalize the model proposed in eq 65.\textsuperscript{191} Remember that we originally wrote \( \mathbf{I} = \mathbf{A} \cdot \mathbf{s} \), where \( \mathbf{I} \) is the image (a size \( n \) vector), \( \mathbf{s} \) the discretized fluorophore density (a size \( N \) vector), and \( \mathbf{A} \) the imaging operator (an \( N \) by \( n \) matrix). In this generalization, the time dependency, over \( T \) frames, of \( \mathbf{I} \) and \( \mathbf{s} \), was also taken into account; these two variables are now matrices respectively of size \( (n \times T) \) and \( (N \times T) \), where each row encodes respectively the time-varying image intensity at a pixel and the time-varying active fluorophore density at a position. The shape of \( \mathbf{A} \) is unchanged, and we still have

\[
\mathbf{I} = \mathbf{A} \cdot \mathbf{s} \quad (67)
\]

However, in this analysis, we will also consider the imaging operator \( \mathbf{A} \) as an unknown.

The problem may thus appear severely under-determined, as we are trying to reconstruct \( N \times T + n \times N \) parameters (\( \mathbf{A} \) and \( \mathbf{s} \)) while having only \( n \times T \) measurements (\( \mathbf{I} \)). However, we can exploit the fact that in our target reconstruction, each row of \( \mathbf{s} \) should represent the time-varying intensity of a single fluorophore at a fixed position; conversely, a reconstruction will be poor if some rows of \( \mathbf{s} \) encode the time-varying mixture of the intensities of multiple fluorophores. From the central limit theorem, the values of \( \mathbf{s} \) adopted by a mixture are necessarily “more normally distributed” than intensity values of a single fluorophore. In other words, a weighted sum of multiple iid random variables is more normally distributed than each individual variable. Thus, we can restate our objective as follows: we seek the solution of eq 67 for which the rows of \( \mathbf{s} \) are “as non-normally distributed as possible”. In order to quantify the “non-normality” of the distribution of values a row of \( \mathbf{s} \) takes, we compute the entropy of the distribution, \( H = -\Sigma s p(s) \log p(s) \). Because the normal distribution has the maximal entropy among all distributions for a given mean and variance, we thus seek the solution (\( \mathbf{A}, \mathbf{s} \)) of eq 67 for which the total entropy of \( \mathbf{s} \) (the sum of the entropies of each row of \( \mathbf{s} \)) is minimal.\textsuperscript{191} This minimization problem is known as independent component analysis, and can be solved using the standard FastICA algorithm.\textsuperscript{240}
The outputs of this analysis are both matrices $A$ and $s$ with $A$ giving the PSF associated with each of the fluorophores while $s$ indicates, for each fluorophore, the frames on which it is active. This analysis does not directly yield superresolved coordinates; it simply separates the PSFs of each fluorophore (into columns of $A$) starting from a data set where they were spatially and temporally overlapping. Each of these PSFs can then be fit to obtain a superresolved coordinate for each fluorophore using any of the methods we have discussed so far.\textsuperscript{191}

4.9.3. Bayesian Deconvolution Approach for Fluorescence Time Series—
Bayesian methods may also be used to obtain both spatial ($A$, following our earlier notation) and spatiotemporal ($s$) information. For example, Cox et al. simultaneously fitted the full set of fluorophore positions, the state histories (bright, transiently dark, or irreversibly photobleached) for each fluorophore, as well as the transitions rates between these states (Bayesian analysis of blinking and bleaching, 3B).\textsuperscript{241}

This Bayesian formulation can be seen as another approach to tackling eq 67: instead of reducing the problem to independent component analysis, the time-evolution of the fluorophores (i.e., each row of $s$) is modeled as a Markov chain between the three above-mentioned states. The true underlying fluorophore spatial distribution is then selected as the one maximizing the likelihood that the entire image stack arises from that distribution. This likelihood is computed by integrating over all possible temporal evolutions (which is done using the forward algorithm).\textsuperscript{193,241} Instead of yielding a maximum likelihood estimate, one can also sample (by Markov chain Monte Carlo) spatial fluorophore distributions from the posterior derived from this likelihood,\textsuperscript{193,241} thus yielding a super-resolved image where the intensity at each position encodes the confidence level about the presence or absence of a fluorophore there. This method is extremely demanding computationally, to the point that cloud-based implementations have been developed.\textsuperscript{242}

4.9.4. Richardson-Lucy Deconvolution for Fluorescence Time Series—Mukamel et al. also proposed a simpler deconvolution method (deconSTORM) taking temporal correlations into account.\textsuperscript{243} Again, the final output of such a method is not a list of coordinates, but simply an image with a higher resolution. Specifically, Mukamel et al. based their work on Richardson-Lucy deconvolution.

Briefly, Richardson-Lucy deconvolution is an iterative approach, whereby the estimated deconvolved image $\hat{I}_k$ at iteration $k$ is derived from the estimate at the previous iteration $\hat{I}_{k-1}$, as well as the measured image $I$, the PSF (assumed known, under a Poisson noise model), and a prior distribution on the true image $p(\cdot)$, which is, in the classical form of the algorithm, kept constant throughout iterations

$$\hat{I}_k = RL(\hat{I}_{k-1}, I, \text{PSF}, p(\cdot)) \quad (68)$$

Mukamel et al. proposed to deconvolve a time series of images by running the iterations of Richardson-Lucy deconvolution in parallel; that is, at each iteration, a new deconvolution of

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each frame is computed. More importantly, instead of keeping the same image prior throughout the iterations, they used a different prior for each frame and updated this prior at each iteration

$$\hat{I}_k(t) = RL(\hat{I}_{k-1}(t), \mathbf{I}, \text{PSF}, p_{I_{k-1}(t)}(\cdot))$$  (69)

Specifically, when running an iteration, the prior for the frame at time $t$ was chosen so that the a priori probability of observing a bright pixel at a given position in that frame is increased whenever the same pixel was also bright in an earlier frame (already reconstructed during this iteration); the closer (temporally) this frame was, the stronger the contribution to the prior. In other words, each frame is deconvolved with a series of priors that, at each iteration, favors a reconstruction similar to reconstructions of the preceding frames.\textsuperscript{243}

### 4.9.5. Recovering Molecular Localizations from Deconvolution-Style Approaches

We have presented deconvolution-style approaches to obtain a superresolved image without first localizing single molecules. But both methods, deconvolution and localization, may be used in tandem. That is, an initial deconvolution step identifying candidate single molecule positions (from local maxima of the deconvolved image), may be used as initial guesses in a subsequent localization step. Such an approach was implemented in the FALCON algorithm\textsuperscript{244} that may be understood as a variant of multiemitter fitting (section 4.8), where the model selection step (finding the correct number of fluorophores to fit) is accomplished by an initial deconvolution.

### 4.10. Drift Corrections

In the absence of active correction, different microscope components drift by dozens of nanometers relative to each other during the acquisition of a single molecule localization data set.\textsuperscript{3,4} This drift affects positions of measured events. Thus, in order to combine all the localization events obtained in that data set into a single high-resolution image, it is necessary to either (i) actively correct for this drift by measuring it in real time and displacing the sample in a compensatory manner or (ii) to estimate the drift in order to subtract it from the fitted positions.

In practice, the second option (drift estimation and subtraction) is typically chosen, as it is a purely mathematical operation, that does not require any modification to the instrument itself. In order to do so, we may track bright fiducial markers (e.g., gold nanoparticles or fluorescent beads) on the coverslip. This can be achieved by using the same localization algorithms as used for “real” events.\textsuperscript{3,4} As the fiducial marker concentration can be chosen to be very sparse, tracking markers from one frame to the next is straightforward. Additionally, the high level of brightness of these fiducials ensures that they are at least as well, and typically better, localized than the events themselves, i.e., they are not a limiting factor for localization accuracy.

Since fiducials are typically bound to the coverslip, such that their motion relative to the camera matches the sample drift relative to the camera, they are less suitable for thick-
sample 3D single molecule localization microscopy. From the instrumentation point of view, the use of fiducial markers in a thick sample data set requires repeatedly switching between the imaging planes and the fiducial (i.e., coverslip) plane.\textsuperscript{217} To avoid the need for such a movement, which complicates the experimental setup and may lead to additional drift itself, one may abandon the use of fiducials and instead rely on correlating event time-slices. In this approach, groups of events are formed by stacking consecutive frames until reaching a set number of events. The cross-correlation between event positions in one group and those in the next then exhibits a peak at a position that encodes the average displacement of the events between the two groups--in other words, the sample drift--as long as the reasonable assumption that both groups are randomly sampled from the entire structure holds.\textsuperscript{245,246}

Neither of these methods can correct drift that occurs on the same time scale as the frame rate as drift estimation requires averaging over a large number of events. In order to increase the rate at which drift information is collected, McGorty et al. proposed instead to use a correlation drift estimator on the bright field image itself (that is, the drift is estimated by finding the shift that maximizes the correlation between a bright field frame and the previous bright field frame).\textsuperscript{247} Of course, it is not possible to simultaneously collect a bright field image in the visible wavelength and single molecule fluorescence in the same wavelength, as the former would swamp the latter; McGorty et al. thus collected the bright field image in the infrared spectrum. Such an approach allowed them to achieve real-time drift correction, with a 10 nm in-plane and 20 nm axial stability, at rates of a few hertz and over minutes of acquisition.\textsuperscript{247}

5. THE COUNTING PROBLEM

Whether two fluorescent events occurring in close spatial and temporal proximity actually come from the same fluorescent protein is an important question that was raised since early PALM experiments.\textsuperscript{4,5}

One reason motivating this question is technical: “stacking” multiple frames together, by summing their intensities before fitting them, ensures that all photons associated with a single labeled protein are used, thereby improving the reconstructed image’s final quality.\textsuperscript{4}

The question was less relevant in STORM, which relies on blinking organic dyes, as multiple dye-labeled antibodies typically bind to the same target.\textsuperscript{3} Each labeled site is thus associated with dozens or hundreds of fluorescence events. Thus, proper assignment of each event to its original label is essentially impossible.\textsuperscript{248}

By contrast, in PALM, given a number of fluorescent events arising from a single diffraction limited spot, it was reasonable to ask whether one may enumerate the proteins, or, alternatively, quantify the protein density, that gave rise to the fluorescent signal. Furthermore, if the number of fluorescent events originated from a group of proteins that formed a complex, with each subunit individually labeled, it may then be possible to quantify protein complex stoichiometry.

Inferring protein complex stoichiometry in vivo is an important problem. Many protein complexes involved in essential cellular tasks contain multiple copies of various proteins.
For example, *E. coli*’s flagellar motor is composed of dozens of proteins all appearing in dozens of copies.\textsuperscript{249}

What is more, protein complex stoichiometry may well be dynamical since a pool of freely diffusing protein subunits available to a protein complex changes over time.\textsuperscript{250,251} The FlIM bacterial flagellar switch protein\textsuperscript{251} is a typical example.

Furthermore, determining a complex’s stoichiometry can also help understand a complex’s operation. For example, asymmetric cell division (sporulation) of *B. subtilis* creates a smaller daughter cell (the forespore) that initially contains only 30\% of its copy of the chromosome; the remaining 70\% must be translocated from the larger daughter cell by SpoIIIIE, a hexameric, membrane-anchored DNA translocase.\textsuperscript{252,253} It was originally thought, based on similarities with bacterial conjugation systems, that SpoIIIIE forms a single aqueous channel between the mother cell and forespore.\textsuperscript{252} Later studies suggested, on the contrary, that the septum is closed and two SpoIIIIE hexamers jointly form a channel across both membranes through which the DNA passes, based on inability of GFP expressed specifically in the mother cell to diffuse to the forespore.\textsuperscript{253} Since both models predict different SpoIIIIE copy numbers at the translocation septum, they could be resolved by accurately counting of SpoIIIIE monomers.

Finally, at the cellular level, proteins and protein complexes can form higher order structures and super-resolution can provide deeper insight into the biological effect of such structures from the spatiotemporal ordering of its constituent proteins. For instance, *E. coli*’s chemotactic clusters—which allow the sensing of gradients of small molecules—involve tens of thousands of receptor proteins.\textsuperscript{254} These clusters are positioned in an apparently periodic fashion on the membrane.\textsuperscript{9} It had been suggested, from time-lapse fluorescence microscopy, that receptor proteins are in fact inserted at random in the cell membrane but later migrate to pre-existing anchor sites.\textsuperscript{9} Other models proposed that this periodicity arises spontaneously from the stochastic nucleation and merging of clusters.\textsuperscript{9} Greenfield et al. suggested that studying the protein number distribution per cluster could offer insights in the mechanism by which they are formed.\textsuperscript{9} We will revisit the type of insight afforded by super-resolution to this question later.

### 5.1. Counting from Fluorescence Intensity

Proteins localized in small clusters can be counted without the need to spatially resolve them. To do so, we may estimate the number of photons collected and divide through by the mean number of photons emitted by each fluorophore. This mean number (the fluorophore’s photon budget) depends not only on the excitation used but also, more crucially, on specific cellular conditions, in addition to other properties such as possible fluorophore interactions.

As an example of this early approach, if clusters contain few fluorophores, a histogram of the cluster brightnesses may exhibit discrete peaks at multiples of a base value.\textsuperscript{255} In such a case, this base value likely corresponds to the intensity of a single labeled protein and peaks observed at two, three, or more times this intensity correspond to clusters of two, three, or more proteins.\textsuperscript{255}
In a different approach, a calibration curve relating fluorescence intensity to fluorophore number is constructed by engineering arrays of, say, 12, 24, and 36 fluorophore binding sites and measuring the fluorescence intensity for each number of bound markers.\textsuperscript{256} Crucially, such a method can be used to establish the existence of a nonlinear relationship between fluorophore count and fluorescence intensity, that can be caused, for example, by interactions among fluorophores.\textsuperscript{256}

The precision of methods relying on total observed intensity is relatively low and relying on a standard mean fluorophore brightness is not without risks. For example, in 2006, Joglekar et al. GFP-labeled a number of yeast kinetochore proteins (where the kinetochore is the structure that links centromeric DNA to spindle microtubules).\textsuperscript{257} They relied on the fluorescence of a single protein within the complex, Cse4, as a GFP fluorescence standard as that protein was thought to exist in a single copy per complex.\textsuperscript{257} However, later studies demonstrated that this assumption was incorrect: Cse4 may be present in 4 to 8 copies per centromere and the reported counts of all other proteins were thus underestimated by the same ratio.\textsuperscript{258,259} Such a difference disqualified earlier arguments indicating that Cse4 may be present in too small a quantity to maintain the necessary attachment points.\textsuperscript{258,259}

Further biochemical studies (protection assays) suggested that one of these proteins (centromere protein A, CENP-A) was, in fact, present at the levels suggested in the Joglekar study.\textsuperscript{260} This time, it was argued that the larger numbers observed by the Coffman and Lawrimore studies\textsuperscript{258,259} arose from the inclusion in their counts of “unincorporated” labeled CENP-As, i.e., those not part of the structure itself but simply lingering in the structure’s vicinity, possibly due to lower incorporation efficiency of labeled CENP-A. Since biochemical studies are not devoid of artifacts either, this controversy remains open to this day,\textsuperscript{261} and should serve as a reminder that the biological question is not to know how many fluorescent proteins are present somewhere but how many of the underlying proteins are actually participating in the process of interest.

5.2. Counting by Photobleaching Using Diffraction Limited Data

An alternative approach is to rely on the stochastic photo-bleaching of single fluorophores.\textsuperscript{250} More precisely, we rely on the observation that the times at which multiple active fluorescent proteins appearing within the same diffraction spot eventually photobleach are stochastic and thus likely different from one another. Thus, a time series of the total fluorescence signal will exhibit a stepwise decrease\textsuperscript{250} with possible double-sized steps if two fluorophores simultaneously photobleach within the time scale of data acquisition.

So long as a majority of steps are resolvable, most steps should coincide with the photobleaching of a single fluorophore (or reversible transitions to and from dark states for blinking fluorophores). This is especially true toward the end of the photobleaching trace where the odds of two simultaneous photobleaching events are comparatively low. Thus, given an estimate for the single fluorophore intensity drop, the number of labeled proteins present at the start of the trace can be estimated as the ratio between the initial intensity and the fluorescence drop arising from a single photobleaching event.
Leake et al. applied this method to study MotB, a component of the stator of *E. coli*’s flagellar motor, concluding that 22 ± 6 copies were present per complex. However, this method also suffers from low precision as the noise level at the start of the trace is high and the initial intensity is therefore poorly defined.

Instead of trying to resolve fluorescence decrease steps, which may be challenging, one may compare the evolution over time of the total intensity of a collection of fluorophore spots (which is decreased by any photobleaching event) to the evolution of the number of spots, within that same collection, which have not completely photobleached yet (which decreases only when all the fluorophores within a given spot have photobleached). The slower the decrease of the number of spots relative to the decrease of the intensity, the larger the number of fluorophores per spot. Yet another approach is to count molecules by means of photon arrival statistics. This technique exploits the photon antibunching effect, which essentially states that a single emitting quantum system (a fluorophore in this case) emits photons one at a time. Therefore, if the temporal resolution of the detector is sufficiently fine, photons detected at the same time can only originate from different emitters. In its most recent implementation, photon counting statistics were gathered and then a nonlinear regression with a Levenberg–Marquardt algorithm was used to back out the number of emitters, i.e., molecules of interest. However, this method is limited to counting up to 20 or so molecules, largely because of error introduced by blinking and photobleaching effects.

A more promising albeit more difficult approach is to attempt to identify and count all individual photobleaching steps, and use their number as an estimate of the protein count. For example, Ulbrich et al. studied the composition of a membrane-bound receptor in *X. laevis* oocytes, that was known to form tetramers. The number of steps in each photobleaching step was visually estimated. Interestingly, the distribution of the number of steps resolved (1–4) is well fitted by a binomial distribution, consistent with a model that only about 80% of the labels are ever fluorescent.

In the sections that follow we will explore theoretical approaches that have been proposed to locate photobleaching steps that can be resolved.

### 5.2.1. Hidden Markov Modeling of Photobleaching Time Series

Even before Ulbrich’s original experiments, Messina et al. proposed to determine the number, *N*, of fluorophores using HMMs, discussed earlier in section 3.3.3, where each state coincides with a combination of states for each individual fluorophore.

The large number of states in this model is suitably shrunk by exploiting the fact that states with the same number of bright fluorophores are indistinguishable, leading to a formulation where each state corresponds to a number *n* of active fluorophores. For instance, the transition from a state with *n* active fluorophores to *n − 1* has a rate equal to *n* times the transition rate to the dark state (as any of the *n* fluorophores could go dark) and, similarly, the transition from the state with *n* active fluorophores to *n+1* has a rate equal to (*N−n*) times the recovery rate from the dark state (as any of those (*N−n*) fluorophores could recover).
Standard maximum likelihood techniques\textsuperscript{193} were then applied to compute the likelihood corresponding to each total number of fluorophores $N$. As there is no penalization for overfitting, this likelihood can only increase for increasing $N$; however, it is expected to plateau once the true number of underlying fluorophores is reached.

This method’s original implementation was applied to time-correlated single photon counting experiments; that is, a setup where stochastic arrival times of each individual photon is measured,\textsuperscript{265} rather than the more common setup where an average intensity is measured by integration over a longer period. In such a case, the source of noise arises from the existence of “background” photons not associated with a fluorophore of interest,\textsuperscript{176} as well as from the stochasticity of the arrival times of the “true” photons. However, the approach of Messina et al. can also be directly adapted to the case where an average intensity is measured.\textsuperscript{146} The authors suggest that up to 30 fluorescent dyes may be counted using such a technique.\textsuperscript{265}

5.2.2. Step-Finding Algorithms in Counting by Photo-bleaching—Without characterizing the kinetics of photo-bleaching, it is also possible to rely on classical step-finding algorithms to count the number of photobleaching events in a time trace.\textsuperscript{96,266,267} The problem of locating sharp discontinuous changes in noisy data, the purview of step-finding algorithms, is a general problem across science that has been investigated across single molecule biophysics.\textsuperscript{59} As always, a precise understanding of the noise characteristics in the data-generating process is required to accurately locate steps (Figure 14).

5.2.2.1. Edge-Preserving Smoothing: Many step-finding algorithms start from an initial filtering or downsampling of the data.\textsuperscript{270} Although linear filters, where each data point is replaced by a weighted average of the neighboring data points within a specified window, are easily implemented, they also tend to blur or smooth out true transitions in the data. In particular, multiple temporally close transitions may become “merged” into a single transition.\textsuperscript{271}

To avoid this effect, Chung and Kennedy\textsuperscript{271} proposed (for the purpose of resolving state transitions in patch-clamp experiments) a nonlinear filter, whereby the weight given to a neighboring point during the filtering depends on how well it predicts the current observation, an approach known in the image-processing field as edge-preserving smoothing.

More precisely, for each data point $y(t)$ in a trace, we consider $2K$ “predictors” of “order” $-K, -(K-1), \ldots, -1, +1, \ldots, K-1, K$, namely the averages of the $i$ (1 ≤ $i$ ≤ $K$) previous or $i$ future data points

$$\langle y \rangle_{-i}(t) \equiv \frac{y(t - \delta t) + \cdots + y(t - i \delta t)}{i}$$
$$\langle y \rangle_{+i}(t) \equiv \frac{y(t + \delta t) + \cdots + y(t + i \delta t)}{i}$$  \quad (69a)

The squared error below provides a metric quantifying the predictor’s quality
The Chung-Kennedy filter is then computed by weighting the predictor of order \(i\) by the inverse \(p\)th power of the average badness of the predictors of the same order \(i\) but considered over the \(M\) preceding (for negative orders) or following (for positive orders) data points

\[
\Delta_{\pm i}(t) = (\langle y \rangle_{\pm i}(t) - y(t))^2
\]

(70)

This filter possesses three parameters (\(K\), \(M\), and \(p\)), which are tuned empirically. The authors demonstrate that an appropriate choice of the parameters leads to a reduction of the noise without distorting sharp transitions.\(^{271}\) More quantitatively, the effect of various filters on the quality of various step-finding algorithms has been the subject of a comparative study by Carter et al.\(^{270}\) finding that a properly (manually) tuned Chung-Kennedy filter exhibited better performance than mean or median filtering.

5.2.2.2. Segmenting the Trace: Regardless of whether (and how) the data is smoothed to facilitate step-finding, the essential part of step-finding is to segment a trace into “approximately constant” regions separated by a step. Two approaches are possible: “bottom-up”, where small regions are merged together on the basis of value closeness, and “top-down”, where the whole trace is progressively split into separate regions.

An example of the bottom-up approach was proposed by McGuire et al.\(^{266}\) Briefly, starting from the beginning of a trace, data points are progressively added to a running window until the value of the fluorescence moves outside of a small range centered at the current window mean. When this occurs, the current window is terminated and a new window started. After running this process on the whole trace, it is repeated on the resulting “leveled” trace until the levels have converged. This whole process is then iterated (starting from the “leveled” trace) using progressively wider window ranges.\(^{266}\)

Conversely, an example of a “top-down” approach is provided by a mathematical idealization of the white noise assumption: the goal is to find the piecewise constant signal \(\hat{y}(t)\), containing \(N\) discontinuities (steps), that minimizes the mean square error, \(\langle (y(t) - \hat{y}(t))^2 \rangle\). Because it is computationally intractable to test all possible combinations of step numbers and their coinciding locations, the number of locations scaling as the number of time points raised to a power equal to the number of change points, Kalafut and
Visscher\textsuperscript{272} proposed to iteratively add change points one after another, each of them at the position that decreases the mean square error the most. At each iteration, one only needs to check, for each time point, the decrease in mean square error if the next change point was inserted there. Even a naive implementation of this approach only exhibits a complexity proportional to the product of the number of time points by the number of change points. We note however that recent theoretical developments have provided efficient exact algorithms to solve this problem through a careful pruning of the solution tree.\textsuperscript{273,274}

As a fit’s mean square error can only decrease as more steps are added, “top-down” approaches additionally require an explicit penalty against overfitting. Kalafut and Visscher propose the use of the Bayesian information criterion (BIC), discussed briefly in section 3.3.4.\textsuperscript{272} However, as acknowledged by the authors, this does typically tend to overfit the data.\textsuperscript{272}

For particular applications, it is always possible to create better step-finding algorithms directly informed by the physics that dictates the noise properties of the problem.

For example, a Bayesian algorithm specifically applied to counting by photobleaching is presented by Tsekouras et al.\textsuperscript{96} In this method, priors and likelihoods are specifically informed by the physics dictating that noise properties should vary stochastically, on the basis of the number of active fluorophores, and that the number of overlapping blinking and photobleaching events have different a priori expectations based on the length of the time trace and the stochastic nature of the photobleaching process.

With this information at hand, Tsekouras et al. arrive at a “top-down” method, more precisely, a marginal posterior for the entire trace, that can be used as a criterion to locate photobleaching steps. The method succeeds in avoiding the overfitting problem arising from assumptions of constant noise across a data set and, according to the authors, scores correctly dozens or even hundreds of steps, provided enough data points are present between successive steps to avoid small number statistics problems (Figure 15).

5.3. Counting by Blinking Correction

As suggested earlier, super-resolution microscopy, and PALM in particular, are seemingly well suited for counting, as the molecular photoactivation times are, by design, as temporally separated as possible. Thus, the number of fluorescence events or bursts detected within a diffraction-limited spot should, in theory, match the number of active fluorophores within that spot. This is only in principle true if fluorophores do not blink.

However, even in the presence of blinking, counting remains possible so long as consecutive bursts originating from the blinking of a single fluorophore can be grouped together. Indeed, threshold methods—described in greater detail below—were used, for example, to study the size distribution of \textit{E. coli}’s chemotactic cluster.\textsuperscript{9}

5.3.1. Threshold Methods for Counting from PALM Data—In the very first implementation of PALM, Betzig et al.\textsuperscript{4} acknowledged the need for such a grouping and applied a purely empirical threshold (“blinking correction time”) to merge events appearing
within the neighboring pixels and separated by no more than three dark frames, in essence applying a very simple form of the nearest-neighbor solution to the linking problem which we discuss in section 6.

Annibale et al.\textsuperscript{275} further studied the blinking kinetics of the widely used mEos2 photoactivatable fluorescent protein (PA-FP), imaged on a coverslip in vitro. By lowering the fluorophore density, the authors could ensure that each fluorescent event corresponded indeed to a single fluorophore (further confirmed by the absence of multistep photo-bleaching).\textsuperscript{275} The authors found that roughly half of the molecules reactivated after entering a dark state, i.e., blinked. Recovery times from the dark state were found to be multiexponentially distributed (similar to observations on other PA-FPs\textsuperscript{276}) and regularly lasted as long as tens of seconds. This observation thus raised the concern that earlier studies may have misinterpreted large number of blinking events as evidence for protein oligomerization.\textsuperscript{275}

The authors thus proposed two methods to correct for this blinking. Either the blinking correction time could be empirically increased or, perhaps more interestingly, the authors found that recovery from the dark state could be accelerated, and thus the blinking correction time kept low, through continuous illumination by the photoactivation laser. Thus, they recommended the use of a continuous photo-activation scheme, rather than a pulsed activation scheme where the photoactivation laser is alternatively turned on for a brief period of time, then kept off while the fluorescence of the activated subset is collected.\textsuperscript{275}

In order to maximize the accuracy of the estimated number of events, Lee et al. offered an alternative selection strategy for selecting the blinking correction time.\textsuperscript{277} First, the authors suggested a scheme by which the photoactivation laser power is tuned in order to ensure a near-constant number of photoactivation events per unit time. Such a strategy ensures maximal separation between active active fluorophores in times and thus minimizes the probability that two fluorophores be simultaneously active.

In short, the scheme was devised by first considering the total number of molecules, \( N(t) \), yet to be photoactivated at time \( t \) with instantaneous activation rate \( k(t) \) which was found to be proportional to the photoactivation laser power, \( P(t) \). The number of molecules that photoactivate between times \( t \) and \( t + dt \) is \( dN = k(t)N(t)dt \); thus, the number of molecules that photoactivate per frame could be kept constant by solving

\[
\frac{dN}{dt} = k(t)N(t) \tag{71}
\]

and selecting a \( k(t) \) (or equivalently \( P(t) \)) that enforces a constant \( dN/dt \).\textsuperscript{277} The authors subsequently generated a simulated time series corresponding to the blinking behavior of various fluorophore numbers. The number of events in each time series was then counted by using various possible values for the blinking correction time; the blinking correction time that achieved, for each given underlying number of fluorophores, the minimum mean bias was then tabulated. As could be expected, the more fluorophores in a single spot, the smaller the correction time that achieved unbiased counting. Since the correct number of
fluorophores is initially unknown, the authors proposed to start from an arbitrary count, pick the corresponding correction time, count using that correction time and iterate. Overall, this approach was shown to exhibit an error of a few percent when counting up to a hundred molecules.\textsuperscript{277}

While this method introduces a blinking correction time (i.e., a threshold), it does not attempt to produce a “grouping” of events correctly identifying whether two events truly arose from the same molecule. In general, such a grouping may be impossible to attain, as the blinks arising from multiple molecules may be interlaced, or even overlap each other (Figure 16). In such a case, “undercounting” (the incorrect merging of events corresponding to two different fluorophores) is unavoidable. Instead, the blinking correction time was chosen (by using the simulations described above) so that these undercounts are exactly compensated by “overcounts” which are cases where a single fluorophore took a time longer than the blinking correction time to recover from its dark state.\textsuperscript{277}

**5.3.1.1. Stochastic Counting Methods from PALM Data:** While the above method\textsuperscript{277} relied on calibrated photophysical (blinking and photobleaching) rates for the fluorophores, it is also possible to learn blinking and photobleaching rates self-consistently from the data itself and, simultaneously, avoid thresholds altogether.\textsuperscript{21} The reason to avoid using calibrated photophysical rates is clear: each protein complex exists within its own unique microenvironment within the cell and rates may differ from one environment to the next.

In order to infer both photophysical rates and the number of fluorophores within an ROI self-consistently, Rollins et al.\textsuperscript{21} introduced an aggregate Markov model (AMM) (section 3.3.3), and used likelihood maximization to infer photophysical rates and protein counts.\textsuperscript{21}

Figure 17 shows the state of an individual fluorophore. The states fall under two categories called aggregates: bright and dark. Only the active state that emits photons belongs to the bright aggregate. All other states (inactive, dark and photo-bleached) belong to the dark aggregate. The AMM is appropriate here as, for a N-fluorophore system (N unknown), there are many states that make up the dark aggregate, bright aggregate and, more generally, the 2-fold, 3-fold, ..., N-fold bright aggregate if 2, 3, ..., N fluorophores are simultaneously active.

Just as we had seen in section 3.3.3, the rates can be cast in the form of a rate matrix $Q$ that governs the switching between the AMM states. Assuming for simplicity only two aggregated states (such that at most a single fluorophore can be active at any given time), $Q$ takes the form

$$Q = \begin{pmatrix} Q_{dd} & Q_{db} \\ Q_{bd} & Q_{bb} \end{pmatrix}$$

(72)

where, just as we had introduced earlier in eq 24, the $Q_{ij}$ are submatrices with $d$ and $b$ denoting dark and bright state, respectively. From this transition matrix, a likelihood can be constructed based on the aggregate state occupied at each time point and from this
likelihood, photophysical rates and fluorophore numbers can be inferred (Figure 18). The authors generalized their treatment in that same paper to treat missed events, that is, dwells in aggregate states (i.e., dark or bright states) shorter than the data acquisition time (Figure 19).

5.4. Limitations of Counting

Biological constraints complicate the counting problem. For instance, even if all proteins are labeled and are expressed in their native amount (which has been made possible by the advent of widespread genome editing systems), not all fluorophores mature, nor will all photoconvertible fluorescent proteins successfully photoconvert. Fundamentally, no algorithm can count proteins that never appear.

Various approaches have been proposed to quantify the percentage of proteins that properly activate. For example, by expressing a labeled human glycine receptor GlyR, whose known stoichiometry of three $\alpha$ and two $\beta$ subunits could be used as a reference, in *X. laevis* oocytes and counting them either by stepwise photobleaching or by blinking correction, Durisic et al. found that across a wide range of photoconvertible fluorescent proteins, only 40 to 80% of the proteins successfully photoconverted. Likewise, Wang et al. expressed a dozen different fluorescent proteins in *E. coli* and compared the number of events collected in a PALM experiment, corrected through division by the mean number of blinks per molecule, to an estimate of the actual number of fluorescent proteins expressed, obtained by quantitative Western blotting. They found an even lower detection efficiency for fluorescent proteins: only between 1% and 20% of them successfully photoconverted. Such limitations need to be taken into account while comparing the accuracy of counting methods: minor gains in the theoretical accuracy of counting will only matter if the global accuracy of the count is not limited by experimental considerations.

6. THE LINKING PROBLEM

Tracking and, in particular SPT, can be broken down into three analysis steps typically performed independently: localization, linking, and interpretation. While the study of static structures only require particle localization, dynamics requires both linking and interpretation.

Inferring dynamics from super-resolution is difficult as there is a trade-off between spatial and temporal super-resolution. Spatial super-resolution requires long exposure times in order to collect enough photons to accurately localize particles. These long exposure times make it difficult to interpolate the positions of particles moving rapidly through space as photons are being emitted at different positions yielding a “smeared” PSF. Shorter exposure times are not necessarily the answer as too few photons may be collected.

Accurate particle linking, often also called data association, which connects a particle’s position across different time frames is therefore an important theoretical challenge. On its own, linking does not characterize the motion of particles, i.e., determine whether the motion is diffusive or directed, although motion models can often facilitate linking. Nor does
linking interpolate the dynamics between frames. The latter two issues are relegated to section 7 on the interpretation problem.

Of the three major SPT components—localization, linking, and interpretation—linking is the least developed. However, it is a critical component: poor linking introduces errors that can derail the subsequent interpretation of the dynamics and is often the limiting factor in determining the quality of tracking algorithms.212,280

6.1. Overview of Linking Methods

Linking methods may be considered local or global.281,282 Local approaches work one particle at a time, by linking a single particle in one frame to a single particle in the next. On the other hand, global approaches take many or all particles in each frame into consideration, linking them to a corresponding number of particles in the next. Global approaches perform better when the density of particles to be tracked is high,281 though this often comes at a heavy computationally cost (Figure 20a)).

Another way to categorize linking methods is to consider the relative weight they ascribe to spatial versus temporal information281,282 (Figure 20b). For instance, to save on computational cost, many methods only consider two successive frames at any given time in order to link, in the most general case, all particle positions in one frame to all particle positions in the next, thereby maximizing the weight of spatial and minimizing the weight of temporal information. Other methods take into consideration multiple frames, thus assigning more weight to temporal information. Such multi-frame integration is correspondingly more difficult and may be one of the primary areas where there is room for improvement in increasing tracking method reliabilities, as studies in human vision have shown integration to be critical for improved performance.282

As we will discuss, methods that are spatially focused and local, such as simple nearest neighbor (NN) linking methods (section 6.4.1), are overall inferior to global, more temporally focused methods, such as multiple hypothesis tracking (MHT; section 6.4.5.3), which directly address several problems mentioned in the next section, such as missed events or abrupt motion changes, at the expense of computational cost. Furthermore, as we will see, methods that exploit knowledge (theoretical or measured) about particle motion in a specific environment, such as whether particles could move via diffusion alone or via combination of diffusion and directed transport, can exhibit superior performance as long as the model is well specified.

6.2. Cell versus Particle Linking

We briefly highlight the differences between particle and cell tracking. By contrast to particles, most cell tracking methods take advantage of the high contrast of the cells with respect to background.284 Such methods include template matching284 which is only good for very similar cell shapes, as it tries to match a model cell image to the acquired image cell and deformable models.285,286 Both approaches achieve linking by identifying similar-looking or similar-shaped cells across different frames. For example, in the deformable model approach, an initial image or shape is assigned to a cell depending on how that cell looks in the first frame and an approximate image or shape is sought in subsequent
By contrast to cell tracking, more challenging global rather than local linking strategies are more useful for particle tracking, in part because particles look identical and are diffraction limited.

6.3. Problems Facing Linking Methods—Difficulties in linking arise not only from the data’s diffraction limited nature but also from measurement noise, as often quantified by the signal-to-noise ratio (SNR). Sources of noise and additional complications involved in linking include the following: photon shot noise, thermal fluctuations, optical alignment drift, occlusion, chromatic aberration, spatial overlap, systematic position errors due to, for example, fluorophore orientation and pixel nonuniformity, inter-and intra-trajectory inhomogeneities, merging, splitting, appearance, and disappearance of particles (Figure 21).

Sources of error in the linking of diffraction-limited particles include not only experimental considerations, but also the choice of models for noise and motion—and the stitching of these models—and the choice of algorithms (local versus global, localization while linking, high cost etc.). For instance, algorithms which localize prior to linking propagate error into the linking step by assuming, say, a Gaussian PSF for fast-moving particles. Such error may be minimized by jointly performing the localization and linking steps in a self-consistent fashion.

In subsequent sections we discuss methods of linking. We refer the readers to Mont et al. for newly defined metrics quantifying the linking quality.

6.4. Linking Methods

Addressing the challenges discussed in the last section requires principled and often problem-specific linking strategies. We now describe basic linking methods and offer an introduction to their mathematical underpinnings.

6.4.1. Nearest Neighbor Linking—Nearest neighbor (NN) linking is an example of a local spatially focused method. NN methods are simple to implement but inappropriate at high particle density, when particles move at different speeds, split, merge, appear or disappear.

All NN algorithms try to match a localized particle in one frame to the nearest localized particle in the next frame. Note that “nearest” in this context can refer not only to spatial distance but also to intensity, volume, orientation or combinations thereof. If a combination is used, a suitable coefficient tuning the relative importance of disparate information types must be specified.

More precisely, given a set of \( N \) positions, \( \{ r_1, \ldots, r_N \} \), predicted by the algorithm (either by doing a brute force vicinity search or one informed by a model of particle motion) and a set of \( N' \) experimentally localized measurements of positions, \( \{ m_1, \ldots, m_{N'} \} \), both referring to the same time frame, the goal is to find a one-to-one correspondence between them. This correspondence is represented by the \( N \times N' \) association matrix \( A \) whose elements \( a_{ij} \) take the value 1 if \( r_i \) and \( m_j \) match and zero otherwise.
When particles appear or disappear, \( N \) may not equal \( N' \) and some estimates or measurements remain unmatched. In this case, “dummy” estimates or measurements are included.\(^{296}\) A NN algorithm then optimizes \( A \) by minimizing the total displacement \( \Phi \)

\[
\Phi = \sum_{i=0}^{N} \sum_{j=0}^{N'} \alpha_{ij} d(r_i, m_j)
\]

(73)

where \( d \) denotes a measure of distance, usually (but not always) simply the Euclidean distance between predicted and measured positions.

In order to reduce computational cost, local assumptions are made. For instance, matches between true particles are only considered if \( d < d_{\text{max}} \) where \( d_{\text{max}} \) is some maximum distance beyond which a particle could not have moved. If a match between a dummy and a true particle is considered, \( d \) is set equal to \( d_{\text{max}} \).

NN algorithms have been used to investigate, among other systems, ligand-gated ion channels such as P2X1\(^{297} \) and 5HT3A,\(^{298} \) Tat (twin-arginine translocation) pathway components in Streptomyces bacteria,\(^{294} \) or the motion (free, hindered, hop, or corralled diffusion) of membrane proteins within cellular or biomimetic membrane compartments.\(^{299} - 302 \) Local degree of confinement within those compartments is determined by membrane protein motion experiments. In such experiments, NN methods work well, because for some cases, motion models (diffusion on membranes) are well characterized and because the objects of interest are sparse. Sparsity may happen naturally though it is often engineered, for example, by photobleaching the bulk of active fluorophores.

6.4.1.1. Residual Tracking: A special case of NN, residual tracking\(^{303} \) is a method introduced for cases where we expect many particles to be fused or, at the very least, in close proximity even if the overall particle density is low. We begin by constructing all possible links that particles in one frame can have to particles in the next. Each one of these associations is called a hypothesis. To rank order these hypotheses, we consider all fluorescence intensities at time \( t \) at position \( (x, y) \) in a 2D image, \( I_t^0(x, y) \). In theory, such fluorescences are obtained from (i) the sum of a constant background fluorescence, \( b_t \); (ii) the fluorescence of the \( i \)th particle, \( P_{i,t} \) extending across some pixels that include location \( (x, y) \); and (iii) a stochastic acquisition noise term \( N_t \)

\[
I_t^0(x, y) = N_t(x, y) + b_t + \sum_i P_{i,t}
\]

(74)

The key idea of this method is to make the residual \( R_t \), defined below, approach zero

\[
R_t = |I_t^0 - b_t - \sum_i P_{i,t}|
\]

(75)
Residual tracking is computationally efficient but seriously limited because it considers only successive frames.

NN methods have clear limitations. They are prone to error if more than one particle is likely to find itself within $d$, i.e., for dense samples, and even more so because the calculation of $d$ depends on the model chosen.

6.4.2. Overview of Filters—We have just described a local, spatial method. We now broaden our discussion to filters incorporating past temporal information in linking that goes beyond the previous frame. This is by contrast to other methods, such as multiple hypothesis tracking (MHT) discussed in section 6.4.5.3, which incorporate numerous frames, past and future, in addition to spatial information all at once.

Filters are sequential Bayesian methods using motion models to create an estimate of the future position of a particle, subsequently updated using current actual measurements. Because of their reliance on motion models, inhomogeneities requiring complex motion models, present difficulties for such approaches. A detailed discussion of the Kalman filter, probably the most commonly used filter, will be relegated to section 7.2.4.

Filtering approaches can be either global, more commonly called joint in the filtering literature, i.e., applied to all the particles in an image frame at once, or local, a.k.a. independent, i.e., applied to one particle at a time. All filtering methods can in theory be used as joint or independent filters, but in practice only the Kalman filter and the particle filter are ever used as independent filters. All filtering methods, especially if used by themselves, assume only small changes in number of detected particles per frame.

Filtering algorithms, often used as part of a broader tracking data analysis scheme, see for example section 6.4.5.1, have also been used by themselves, for example in linking frames for gene carrying particles (viruses, drugs, etc.) and viral particles.

6.4.2.1. Joint Filters: Spatiotemporal Filtering: Unlike independent filters, joint filters link all particles in one frame to all particles in the next. This avoids those conflicts that plague independent filters but is correspondingly more computationally expensive. The simplest joint filter is the spatiotemporal filter, which can be thought of as a generalization of the Kalman filter, discussed in section 7.2.4.

To be more specific, for spatiotemporal filters we assume, in analogy to HMMs (section 3.3.3), that, at time $t$, a particle has a hidden state $r_t$, represented by a noisy measurement $m_t$. The goal is to estimate $r_t$ from the data. The spatiotemporal filter achieves this for all particles in a frame, i.e., it is a global linking method, by recursively finding the posterior, $p(r_t|m_{1\rightarrow t})$, where the index $1\rightarrow t$ denotes all frames from frame 1 to frame $t$. To find the posterior, we first specify the following quantities: $p(r_t|r_{t-1})$, dictated by the motion model to describe the evolution of the hidden state; $p(m_t|r_t)$, representing the relation between observation and state; and an initial prior $p(r_0)$. Using these quantities, a two-step recursive Bayesian approach yields the posterior.
where in the last expression, we invoked Bayes’ theorem and assumed measurements are independent. While in general the integral of eq 76 cannot be evaluated analytically, under some Gaussian approximation the formalism reduces to the well-known Kalman filter\(^{296}\) (section 7.2.4). This is because the Kalman filter performs exactly the task of creating a prediction \(p(r_t|m_{1\rightarrow t-1})\) for the particle position, then compares with the measurements, and finally yields a result by combining the prediction and the difference between prediction and measurements, called the innovation. In the Kalman filter, prediction and innovation must be combined linearly via some gain matrix; the Gaussian approximation is required to mathematically justify the gain matrix form.

Alternatively, for non-Gaussian \(p(r_t|r_{t-1})\) and \(p(m_t|r_t)\), the integral of eq 76 can be estimated by sampling the integrand, typically via MCMC. This method is called the particle filter. What is more, inconsistencies in the performance of filters across an image have been addressed using a modified sampling scheme called the boosted particle filter.

Finally in the special case where there are only a few particles that look and behave very similarly per frame, it is possible to use a mixture of particle filters to approximate the posterior \(p(r_t|m_{1\rightarrow t-1})\). For example, two particle filters, each corresponding to a different motion model, can be constructed and then the posterior estimated as a weighted sum of the samplings from each filter.\(^ {296}\)

### 6.4.2.2. Independent Particle Filters:

Independent particle filters use one particle filter for each particle. That is, only a single particle is linked between two successive frames. Computational requirements for independent particle filters are therefore low. However, independent particle filters fail if particles pass so close to each other that the corresponding filters both choose the same particle, i.e., conflict. When filters conflict, one solution is to find sets of trajectories that minimize the number of appearing and disappearing particles.\(^ {296}\) Drawing from graph theory, each particle position can be thought of as a vertex and each link between two positions as an edge. When particles conflict, some vertices have more than two edges, while others may have one or none in which case they are called orphan vertices. Godinez et al.\(^ {296}\) propose a method to minimize orphan vertices by creating candidate links between them following a logic very similar to NN linking discussed in section 6.4.1. The need to find some procedure to resolve conflicts is the main disadvantage of independent particle filters.

It is theoretically possible to use filters other than the particle filter to follow each particle, so that we can, for example, have independent Kalman filters. However, no implementations of such filters have–to our knowledge–ever been constructed, since the risk (filter conflict) outweighs the reward (easing of the computational burden). Independent particle filters is
the only exception because particle filter sequential Monte Carlo sampling already imposes a heavy computational burden so easing it becomes a priority.

6.4.3. Energy Minimization—In linking by energy minimization, we begin with a set of frames and stack them in the order in which they were acquired thereby forming a 3D volume. The linking problem, that is, finding trajectories within this 3D volume, is reduced to a global minimization of an objective function defined on this volume that can be thought of as an energy.\textsuperscript{293} Energy minimization methods are global and consider both temporal and spatial information in the linking process. By contrast to filters and MHT, discussed in section 6.4.5.3, energy minimization methods are non-Bayesian. More specifically, while filter methods only incorporate time up to the frame currently being processed and MHT considers, in principle, all possible trajectories across all space and time, energy minimization methods are more computationally efficient.\textsuperscript{282}

High accuracy combined with low computational cost have made energy minimization methods popular for linking, especially for cases with smooth, long tracks. For instance, they have been used to track low-density lipoproteins transported inside endosomes and adenovirus-2 particles moving along microtubules,\textsuperscript{305} viruses transporting across membranes,\textsuperscript{306} diffusion of proteins in chromatin structures,\textsuperscript{307} changing distributions of Vipp1 proteins, essential for photo-synthesis, in cyanobacteria\textsuperscript{308} and \textit{lac} repressors in DNA in a study of DNA repair mechanisms related to antibiotic resistance in bacteria.\textsuperscript{309}

The effectiveness of energy minimization methods depends on one’s choice of objective function (i.e., energy). For instance, objective functions may be phenomenologically motivated based, in whole or in part, on physical properties such as the fluorescence signal to be linked\textsuperscript{305,310,311} (Figure 22). For example, we may define a distance metric such as eq 73, albeit for all particles as opposed to one particle, that may be further regularized by the quadratic difference in fluorescence intensity moments of order 0 and 2.\textsuperscript{305} In a similar example, trajectories were constructed by following high gray intensity curves.\textsuperscript{310} In yet another similar example combining NN and energy minimization, Bonneau et al. tracked glycine receptors in neurons but linked brightness rather than distance globally for the particles.\textsuperscript{311} Track fragments across the 3D volume were linked in such a way as to minimize 3D volume geodesics albeit constrained by physical motion models, such as the local diffusion coefficient.\textsuperscript{311}

Objective functions can also be motivated by mathematical principles. For example, Shafique and Shah\textsuperscript{312} draw heavily from graph theory, constructing a graph where localized particles in all frames are vertices and links are edges. They invoke tools from graph theory in order to maximize trajectory length while ensuring that the optimal trajectory set minimizes orphan vertices, i.e., maximizing graph cover.\textsuperscript{313} These graph theoretic inspired methods can tolerate cost functions that are not optimized for the problem or that are \textit{ad hoc} to a greater or lesser degree. They are therefore more general than physics-based energy minimization methods and avoid having to go to great lengths to pick suitable motion model based cost functions.
A method, detailed elsewhere,\textsuperscript{314} that brings together phenomenological objective functions and graph theoretic methods has been used in studies of the tumor necrosis factor 1 crucial in inflammation,\textsuperscript{315} integrin and talin protein dynamics in cellular focal adhesions,\textsuperscript{316} and enzymes that are suspected of promoting tumor generation.\textsuperscript{314}

**6.4.3.1. Linking via Message Passing between Frames:** This new method is an attempt to fix the problems energy minimization methods encounter when particle density is high and works best for such high particle densities (Figure 23).

The method starts by formulating the linking problem as a weighted complete bipartite graph.\textsuperscript{291} Subsequent mathematical mapping detailed elsewhere\textsuperscript{291} reduces the linking problem to what is known as the message-passing formulation of belief propagation. This normally is a computationally very expensive method, but it can be approximated by a MCMC algorithm which reduces computational cost by multiple orders of magnitude, making the method tractable.

A general criticism that can be made of all energy minimization methods is that they rely on trajectories that are long and regular, which is not always realized in biophysics where observed trajectories are jagged and short due to fast changes in motion or photobleaching.

**6.4.4. Multiple-Target Tracking**—While filters (section 6.4.2) are Bayesian and operate one frame at a time, multiple-target tracking (MTT)\textsuperscript{317} is a non-Bayesian linking method that also works one frame at a time. It has the advantage of being much less computationally expensive than filters. MTT links particles to the next frame starting from localized particles at the current frame in addition to knowledge of motion model and photophysical parameters drawn from previous frames for each trajectory (Figure 24).

More specifically, in order to link existing trajectories to localized points (particles) in the new frame, an area where the particle would be expected in the new frame must be determined. This area is calculated based on a Brownian motion diffusion model, although the method can be adapted for other motion models, with the position of the particle in the last frame as initial condition and bounded by some predetermined distance. If a single particle is found within that area it is added to the trajectory. If there is overlap, i.e., there are more than one candidate particles in the area, MTT compares the likelihood for each possible linking as a function of the positions, intensities and blinking probabilities of the particles involved. If there are any unconnected particles in the area, particles are relocalized and the procedure repeated.

MTT is computationally inexpensive but is much less powerful than algorithms that take many frames, past and future into consideration, like MHT (section 6.4.5.3); it is also much less effective if the particles being tracked behave very similarly because if their trajectories cross it is prone to misidentifications. Initially demonstrated on tracks created by the motion of epidermal growth factors in the plasma membrane of cells,\textsuperscript{317} it has since been used to study actin dendritic spines\textsuperscript{318} and enhanceosome assembly in embryonic stem cells.\textsuperscript{85}
6.4.5. Bayesian Nonfilter Methods—Here we briefly highlight a group of Bayesian linking methods devised for frames where particle density is high (which causes filters (section 6.4.2) to fail) and particle motion is not smooth (which causes energy minimization (section 6.4.3) methods to fail). These methods, informed by priors derived from models of motion or knowledge of particular system, are particularly useful in biophysics, where jagged motion and high particle densities are often encountered but can fail if motion models are inadequate, or probability distribution assumptions do not hold (Figure 25).

6.4.5.1. Interacting Multiple Model: Kalman filters, discussed in section 7.2.4, assume a single motion model. By contrast, the interacting multiple model (IMuM)\textsuperscript{320} uses multiple Kalman filters each with a different, competing, motion model to estimate particle positions in the next frame. The Kalman filter that wins the competition sets the motion model for that particle in that frame alone; in the next frame the Kalman filters compete anew. In this way IMuM avoids any assumption of smooth motion, reflects the fact that particles in biological settings can switch between types of motion relatively quickly and also avoids the very high computational cost of the MHT (section 6.4.5.3).

More specifically, similar to filters, IMuM computes the posterior $p(r_t|m_{1\rightarrow t})$ by assuming that the system can be reduced to a set of $j$ switching linear models with additive Gaussian noise profiles of the general form

$$r_t = A^j_t r_{t-1} + \eta_t^j \text{and } m_t = H^j_t r_t + \epsilon_t^j \quad (78)$$

where $j$ is the index of the model that is active at time $t$, $A^j_t$ is the state transition matrix, $H^j_t$ is an “observation matrix” relating the hidden position to the measurement, and $(\eta, \epsilon)$ are the “process” and “measurement” noise vectors respectively, described in greater depth in section 7.2.4. Both noise vectors are assumed to be uncorrelated white zero mean Gaussian processes.

As there are too many possible sequences of motion models to consider for the problem to be computationally tractable across all frames, IMuM considers only two successive frames $i, i+1$ at a time. That is, IMuM calculates position predictions for the $(i+1)$th frame from each motion model via a Kalman filter employing that motion model. Once the predicted positions for each motion model have been constructed, the motion model whose predictions for each particle best matches a measurement is determined by maximizing the likelihood of the innovation for the Kalman filters for each particle in each frame.

IMuM yields good results but only when the linearity and Gaussianity assumptions that underlie the Kalman filter are justified. It also has a disadvantage compared to methods such as the MHT since it does not consider future temporal information.

6.4.5.2. Linking with iHMMs: Here we mention a non-parametric Bayesian method that appeared almost a decade ago.\textsuperscript{321} Prohibitively expensive computationally at the time it was first presented, it now shows promise. This method relies on iHMMs (section 7.3.2), and thus can grow in complexity as the amount of available data increases. To be more precise,
previous methods, such as IMuM (section 6.4.5.1), prespecify the number of motion models or state space. To accommodate the complexity of large data sets, large state spaces would be needed, thereby increasing the computational burden. iHMMs only grow the state space as needed and thus we argue that they are well suited for cases where large numbers of particles may be moving with a variety of motion models in heterogeneous environments, as is the case, for example, of different transcription factors diffusing in the nucleus.\textsuperscript{322}

6.4.5.3. Multiple Hypothesis Tracking: Multiple-hypothesis tracking (MHT), first proposed by Reid,\textsuperscript{323} is perhaps the most advanced among available methods. It is intended to tackle the major problems of particle motion heterogeneity:\textsuperscript{280} high particle density, particle merging, splitting, and temporary disappearances.

In principle, MHT, starts from all particle positions localized across time frames (past and future), thereby constructing all possible trajectories that could arise when linking all frames.

In practice, this is not computationally feasible, and thus a greedy algorithm is often used to stitch locally optimal links across frames. Greedy algorithms are more efficient, but pay a price in that some potential trajectories are never considered. Here we describe one such implementation by Jaqaman et al.\textsuperscript{280}

In the first step of this implementation, particles in one frame can link to at most one particle in the previous or next frame. Taking all such possible links for each particle in a frame, short candidate segments linking just two localized particle positions are created. Segments are also constructed for particles in one frame linking to nothing in the previous or next frame, these links corresponding to a particle appearance and disappearance, respectively. Some of these candidate two-particle trajectory segments are mutually exclusive, i.e., conflict.

In the second step each possible two-particle trajectory segment is assigned a cost, and in any group of conflicting segments, only the lowest cost segment is retained.

In the third step, these two-particle trajectory segments across all frames are themselves linked in five ways: end of one segment to the start of the other, for continuation; end-to-middle, for particle merging; middle-to-start, for particle splitting; end-to-nothing, for trajectory ending; and nothing-to-start for new trajectory initiation. This trajectory segment linking creates large numbers of potential trajectories, many of which conflict. Trajectories that do not conflict are regrouped as a candidate tracking hypothesis. Step four assigns a cost function, which may be Bayesian or otherwise motivated, to each hypothesis and the lowest cost hypothesis is selected, similar to energy minimization methods (section 6.4.3).

This implementation of MHT has been used to follow the motion of the macrophage transmembrane CD36 receptors;\textsuperscript{280} the formation of clathrin-coated pits;\textsuperscript{280,324} microtubule dynamics;\textsuperscript{325} and AAA+ ATPase motor cytoplasmic dynein dynamics.\textsuperscript{326}

Chenouard et al.\textsuperscript{287} extended the previous treatment by turning MHT into a Bayesian framework. Briefly, for reasons of computational efficiency, they built tracking hypotheses
between frames \( k \) and \( k + d \) by considering all possible positions that a particle could take in each of these frames. To further limit computational cost, they stored hypotheses as tree graphs, and in the first step of MHT, they do not explicitly consider particle splitting, merging, appearance, and disappearance when constructing trajectories.

This reduction of the computational burden is necessary in view of two new features this implementation of MHT introduces, both of which are computationally costly.

The first new feature enters in the second step of MHT. When linking a measurement A in frame \( k \) to a measurement B in frame \( k + 1 \), the algorithm assumes a particle has moved from A to B via a combination of diffusion and directed transport and assigns each of these a Kalman filter (section 7.2.4). Then the algorithm employs an iHMM (section 7.3.2) to set the specific motion parameter models used in the two Kalman filters and among conflicting trajectory segments retains the one with the most likely parameter values.

The second feature enters in the third and fourth steps of the MHT. Instead of creating trajectories by linking only trajectory segments whose starts or ends are in adjacent frames, the algorithm links trajectory segments even if they are separated by several frames. This compensates for the algorithm ignoring particle appearances and disappearances in step one. Trajectory segments involving such “hops” across frames, realized in the tree graph formalism via virtual detections, are assigned a cost on the basis of a model based on the photophysics of particle disappearance for the particular system.

Once trajectories have been created up to the predetermined tree depth, i.e., frame \( k + d \), the most likely tracking hypothesis in a single frame, i.e., \( k+1 \), is obtained by maximizing the full joint probability up to the \( k + d \) frame (Figure 26).

Although computationally heavy and rarely used thus far, this implementation of MHT shows great promise due to the thoroughness of the tree-graph approach and the highly customizable extra features that allow for significant expansion the method. For example, it is possible to inform the iHMM with additional motion models or extend the tree graph approach to account for particle splitting and merging, which are not considered thus far.

6.5. Temporal Super-Resolution

Temporal super-resolution, which is to localize particle positions at time intervals \( t \) across frames collected at longer times, can also be thought of as the ultimate solution to the “connect-the-dots” problem. The “connect-the-dots” problem is the mirror image of the interpretation problem. In the interpretation problem we ask “given a particle trajectory, what can I learn about how the particle moves (i.e., the motion model) between observations?” In the “connect-the-dots” problem, we ask “given some imperfect idea of the particle motion, how can we best connect localizations of the same particle in adjacent frames?” The “connect-the-dots” problem can be seen as part of linking in that a motion model can guide us in making the correct links, i.e., correctly identifying which observations across frames belong to the same particle. It can also be seen as part of an iterative approach to both linking and the interpretation problem: if we assume a naive motion model and link observations to get initial trajectories, we can then interpret those trajectories to change the
motion model and get higher quality links and better trajectories, then interpret those and so on. In this context, temporal super-resolution is important because it would allow us to sample the particle motion at such small intervals that any uncertainty as to the motion model would disappear. In such a situation, linking would be trivial, interpretation much easier, and “connecting-the-dots” unnecessary.

The connect-the-dots problem has yet to attract the intense attention that spatial super-resolution enjoys. While biology’s traditional focus has been on structure, counting and localization, a clearer picture of dynamics may help address fundamental questions such as the mechanism of transport between nucleus and cytoplasm or protein folding pathways.

As we described earlier, spatial super-resolution, requiring hundreds of photons to localize particles with tens of nanometers resolution, comes at the cost of temporal super-resolution, i.e., time resolution less than μs or faster, that we refer to as the spatial-temporal trade-off arising from a fluorophore’s limited photon budget.

Whereas the spatial PSF is typically a Gaussian, a fluorophore passing through an ROI appears in that ROI as a box-shaped pulse in time that may last for more or sometimes less than a single frame (Figure 27). Thus, any method relying on a Gaussian noise pulse in time introduces error.

The photons arising from this pulse are spread out over a spatial region, i.e., are blurred, either because the particle is moving fast compared to the frame rate or because of errors caused during localization by the use of PSF functions that are not appropriate for fast moving particles. Another reason for this blurring, which contains the temporal superresolved image we seek, is that the camera frame rate may be slower than typical frequencies of motion where few photons are released per pixel in turn complicating, in a manner known as aliasing, the “connect-the-dots” problem (Figure 28).

Existing theoretical methods increasing temporal super-resolution are often motivated by machine vision or other applications. However, methods developed for many machine vision applications often implicitly assume an abundance of photons to be detected, contrary to shot noise expected in super-resolution, and may rely, either in whole or in part, on increasing the number of cameras or optimizing a camera sampling scheme.

However, there are methods from machine vision more appropriate for SPT. Briefly, a (Bayesian) inference approach used by Borman and Stevenson relies on the notion that frames collected at high data acquisition have low spatial resolution but high temporal resolution while frames obtained by integrating the fast frames yield low temporal resolution but high spatial resolution. This algorithm exploits information at both frame rates to achieve both high spatial and temporal localization at once. In order to do this, it introduces a highly under-determined association matrix that relates position of objects (i.e., particle localizations in super-resolution) in high resolution frames to many possible object histories (i.e., particle trajectories in super-resolution) allowed by the low resolution frames. Motion models and observation models, an analog of an emission probability in HMMs, are subsequently used to reduce the indeterminacy.
It has been shown that indeterminacy can be further reduced by exploiting two effects: the fact that some parts of the image remain static and that low-resolution pixels affect only those high-resolution pixels with which there is spatial and temporal overlap.

At present, temporal resolution lags far behind spatial resolution. While spatial resolution is pushing into the nanometer scale, temporal resolution rarely goes beyond the millisecond scale. As a result, we cannot follow particle motion directly and instead must infer its parameters (velocity, diffusion coefficient, etc.) from the trajectories we construct: this is the purview of the interpretation problem.

7. THE INTERPRETATION PROBLEM

Stacks of in vivo optical images are the starting point for the analysis of dynamics inside cells.\textsuperscript{61,65,66,85,86,88,95,126,333}

Previously, we discussed that particle tracking involves particle localization and linking. Unless stated otherwise, in this section we start from the assumption that we have performed both localization and linking with minimal error and that the observed trajectory originates from a single molecule (a realistic assumption when the fluorescent label density is relatively low).

The possibility that linking or localization errors may be large, as may be expected in crowded environments or if fluorescence efficiencies depend on the local cellular environment, or that linking or localization may be performed jointly while inferring motion models from single particle tracks goes beyond the scope of this section.

Rather, here we are focused on time series analysis or what we've termed the “interpretation problem”. That is, the inference of motion models and their parameters, such as velocities, forces, diffusivities, binding interaction strengths, from previously localized and linked single particle trajectories.

We begin this section by reviewing established or “legacy” SPT time series analysis approaches, section 7.1, and discuss some of their fundamental limitations. For instance, we spend some time discussing mean square displacement (MSD) methods, arguably the most commonly used technique in SPT data analysis,\textsuperscript{60,65,89,299,334} which dates back to the turn of the 20th century\textsuperscript{335} and provides average information on particle trajectories.

We follow this discussion on methods that do not average down the information, as MSD does, and explicitly account for noise, section 7.2, where a tracked particle’s “true” position, \{r\textsubscript{1}, r\textsubscript{2}, … r\textsubscript{T}\}, may differ from the “measured position”, \{m\textsubscript{1}, m\textsubscript{2}, … m\textsubscript{T}\}. In particular we will spend some time discussing the classic Kalman filter.

We will finally extend our discussion to nonparametric approaches, briefly introduced in section 3.4, further developed herein, section 7.3.

The motivation for moving toward nonparametric model formulations in time series analysis is clear: the heterogeneity of a cell’s environment often warrants more complex models. For example, we may want to treat abrupt changes in motion models that may arise from
spontaneous molecular changes or interactions with heterogeneous environments\textsuperscript{61,63,66,95} or, alternatively, use information on these environments to inform motion models.

For instance, high spatial (10 nm) and temporal (10 μs) resolution 3D SPT measurements were recently combined with conventional two-photon scanning microscopy to obtain useful “contextual” information on the changing environment in which tagged nanoparticles diffused\textsuperscript{14,15}. This new imaging approach showed that the motion of a 3D trajectory of a fluorescently tagged nanoparticle, mimicking a virus, was heavily influenced by larger 3D surface structures that were not otherwise captured by standard SPT motion models (Figure 29).

7.1. Legacy Time Series Analysis Methods

Mean square displacement (MSD)\textsuperscript{22,60,65,89,299,334\textsuperscript{+}336} ideas are often used to determine whether a particle is “immobile” or which type of diffusion, such as “normal” or “standard”, “confined”, “directed”, or “anomalous”, is most probable.\textsuperscript{60} If the particle diffuses normally, that is, the motion is Brownian and characterized by a single diffusion coefficient, then the MSD of the molecule’s position ($r$) at absolute time $t$, $\langle |r(\tau)|^2 \rangle$, defined as

$$\langle |r(\tau)|^2 \rangle = \frac{1}{N-\tau} \sum_{t=1}^{N-\tau} ||r_{t+\tau} - r_t||^2$$  \hspace{1cm} (79)

is proportional to $\tau$, as the number of observations, $N$, tend to infinity, where $\tau$ represents a relative “time lag” between data points, $N$ represents the number of observations in the trajectory, and $|| \cdot ||^2$ denotes the Euclidean (or $L^2$) distance squared. For other types of diffusive behaviors, beyond normal diffusion, eq 79 is no longer simply proportional to $\tau$ for large $N$. Rather, Table 1 lists MSD curves, represented graphically in Figure 30, expected for different diffusive behaviors (using notation from Park et al.\textsuperscript{60}). The other diffusive behaviors can be derived from specific motion models\textsuperscript{57,299,337} such as stochastic differential equation (SDE) models discussed in section 7.2.1.

Despite their widespread use, we now discuss several documented issues associated with using MSD curves to analyze finite length trajectories.\textsuperscript{22,66,86,90,93,334}

7.1.1. Issues Facing Legacy Methods—Many of the problems faced by MSD that we now discuss are also faced by other traditional SPT analysis methods such as power spectral\textsuperscript{56,90,338} and velocity autocorrelation\textsuperscript{58,338} methods.

7.1.1.1. Issue 1: The Stationarity Assumption: The assumption of stationarity\textsuperscript{110}—the dependence of trajectory statistics on relative time differences and not absolute time—is one systematic shortcoming facing all three methods mentioned above.

This assumption breaks down in many SPT applications for a variety of reasons; for example, if a system is initialized far from equilibrium and eventually reaches steady state during the course of data acquisition.
Stationarity can also be violated if parameters governing the motion of the tagged particle change abruptly. This is especially problematic in vivo since microenvironments sampled by the tagged particle may change in both time and space and, even within a single trajectory, the dynamics of a tagged particle can abruptly change or drift slowly over time.\cite{61,65,66,339,340} As a concrete example, the diffusion of some receptor proteins may be suddenly perturbed by the appearance of a clathrin-coated pit that traps the protein.\cite{341}

As another example, the stationarity assumption may also be violated when the background fluorescence varies across microenvironments sampled by the tagged particle, hence altering the effective measurement noise or when multiple trajectories—each with differing fluorescence intensities—are averaged into a single MSD thereby introducing “heterogeneous measurement noise”.

7.1.1.2. Issue 2: Spatial and Temporal Coarse Graining: Even if stationarity holds, the MSD of eq 79 coarse-grains spatial locations and averages down time. Yet time averaging can be avoided by using likelihood methods that both respect the data’s natural time ordering (see eq 6) and can be used to infer quantities such as the directionality of velocities and forces acting on tagged particles.\cite{22,86,88} This type of coarse graining is also a problem facing power spectral methods\cite{56,90,338} and velocity autocorrelation techniques.\cite{58,338}

7.1.1.3. Issue 3: Temporal Correlations Introduced by Averaging: As previously described, MSD curves are obtained by time averaging. Since the same data are used to produce correlations at two different times, say $\langle r(\tau)^2 \rangle$ and $\langle r(\tau')^2 \rangle$, the process of time averaging produces highly correlated quantities.\cite{22,334} This high statistical correlation is particularly unattractive if the amount of data is low since correlations within data reduce the effective sample size. To mitigate these problems, ad hoc measures, such as restricting $\tau$ values to less than some predetermined fraction of $N$,\cite{12,85} are invoked in an effort to reduce correlations though such methods typically ignore the fact that true temporal correlations (such as in the case of correlated measurement noise) may exist.\cite{22,93} Indeed, MSDs for normal diffusion with correlated measurement noise may be confounded with MSDs arising from anomalous diffusion.\cite{22}

7.2. Beyond Traditional SPT Analyses

In the previous section we discussed key challenges facing traditional SPT trajectory analysis tools and never explicitly dealt with noise as it is averaged away when computing mean quantities.

Here, we provide a brief overview of recent promising SPT trajectory analyses that overcome many of the issues highlighted in section 7.1.1 and motivate data-driven models that explicitly account for effects such as temporal and spatial heterogeneity characteristic of live cell measurements.

Many methods we discuss, including stochastic differential equation (SDE) models and hidden Markov model (HMM) inspired methods, were first introduced to disciplines such as statistics, control theory, and mathematical finance and have begun to make their way into super-resolution data analysis. Sections 7.2.1–7.2.4, presented with a high-level didactic
flavor with SPT applications discussed as appropriate, build the background needed to understand the more complex “regime switching” time series models later surveyed, section 7.3.

7.2.1. Stochastic differential equation models—The MSD expressions shown in Table 1, with the exception of anomalous diffusion, can be derived from the following parametric stochastic differential equation (SDE) describing particle position, \( \mathbf{r}_t \), at time \( t \)

\[
d\mathbf{r}_t = (\mathbf{v} - \kappa \mathbf{r}_t)dt + \sqrt{2D}d\mathbf{B}_t \quad (80)
\]

where \( D \) denotes the diffusion coefficient, \( \mathbf{B}_t \) denotes standard Brownian motion (i.e., the white noise), and \( (\mathbf{v} - \kappa \mathbf{r}_t) \) denotes the instantaneous particle velocity where \( \kappa \) allows for linear variation in the velocity field and can be thought of as proportional to the spring constant of a harmonic potential. This SDE assumes overdamped dynamics, i.e., a stochastic dynamical model neglecting inertial effects.

The particle dynamics are summarized by the parameter vector \( \mathbf{\theta} \equiv (\mathbf{v}, \kappa, D) \) where \( \mathbf{v} \) and \( D \) are square matrices of the same dimension as the position vector. Values for the parameter values reproducing the limits of Table 1 are shown in Table 2.

The SDE of eq 80 is the starting point for many likelihood formulations that respect natural time ordering. For simplicity of demonstration only, we ignore measurement noise effects and use the Markovian SDE to write the likelihood \( p(D|\mathbf{\theta}) \) of measurements \( D = \{\mathbf{r}_1, \mathbf{r}_2, \ldots, \mathbf{r}_N\} \) as follows

\[
p(D|\mathbf{\theta}) = \prod_{i=1}^{N-1} p(\mathbf{r}_{i+1}|\mathbf{r}_i, \mathbf{\theta}) \quad (81)
\]

where \( p(\mathbf{r}_{i+1}|\mathbf{r}_i, \mathbf{\theta}) \) is the transition density.

Importantly, this transition density is obtained by solving the Fokker–Planck\textsuperscript{117} equation associated with the SDE of eq 80 which, for the linear SDE provided, may be computed exactly.\textsuperscript{25,86,118}

Eq 81, and its construction from SDEs, motivates our later discussions on diffusion as well as the Kalman filters, Sec. (7.2.4), where we discuss how to extend eq 81 to include not only measurement noise, as we did with HMMs (section 3.3.3), but also “process” noise, such as error in localization.

7.2.2. Hidden Markov-Type Models in SPT—In section 7.1.1 we described the challenge of traditional ensemble methods, such as MSD, in dealing with tagged particles that may undergo abrupt changes in dynamics. While the parametric SDE described in eq 80 only contained a single diffusion coefficient, and thus describes the dynamics of a particle in a single “diffusive state”, a recent promising Bayesian method named “vbSPT” (variational
Bayes single particle tracking) generalizes the treatment above to an arbitrary number of diffusive states. More concretely, from SPT trajectories, vbSPT (i) determines the number of diffusive state and assigns each data point to a state (called “segmentation”); (ii) parametrizes the diffusion coefficient for each state; and (iii) quantifies the transition probabilities between states. vbSPT does so by first performing an initial HMM analysis on a preset list of candidate diffusive states and subsequently comparing their “maximum evidence”, defined by Persson et al., as a model selection criterion.

While measurement noise effects ignored in vbSPT may lead to spurious state assignments, vbSPT is still helpful in analyzing data where trajectories are expected to come from “purely diffusive” shorter trajectory segments (Figure 31).

As we discuss shortly, recent advances in nonparametric Bayesian methods can help avoid prespecifying the number of candidate states, by simultaneously inferring the number of states and parametrizing the model, and help broaden the types of models considered beyond diffusion. Before these topics are addressed, we highlight recent techniques for dealing with measurement noise in experimental SPT data.

### 7.2.3. Simple Example of Measurement Noise in SPT

We motivate our discussion on the Kalman filter in the next section by discussing a simpler model incorporating both “static” measurement noise (the localization noise due photon statistics) and “dynamic” measurement noise (position estimation errors due to finite camera exposure time). To illustrate the effect of both noise sources, we begin with a simplified form of the SDE model, eq 80, presented earlier:

\[
dr_t = \sqrt{2D} dB_t \tag{82}
\]

\[
m_{t_i} = \frac{1}{t_E} \int_{t_{i-1}}^{t_i} r_s ds + \epsilon_{t_i}^{\text{loc}} \tag{83}
\]

where \(m_{t_i}\) denotes the position measurement taken at frame \(i\) under the assumption of discrete time measurements, \(t_E\) is the camera’s exposure time, and \(\epsilon_{t_i}^{\text{loc}}\) is the “static” measurement error with each draw of this noise being an iid mean zero Gaussian with variance \(\sigma_{\text{loc}}^2\). The integral is needed since molecules continuously move and emit photons as the camera’s shutter remains open. Also, while the above assume constant uniform illumination, the expressions are straightforwardly adapted to treat different types of illumination used in microscopy including stroboscopic and time lapse (where the illumination laser is not continuous, in order to reduce photobleaching at the expense of temporal resolution).

Our goal is to use a likelihood in order to infer all model parameters. In the absence of blinking or photoactivation, one way to move forward is to define...
\[ \Delta_i = m_{t_{i+1}} - m_{t_i} \quad (84) \]

While from eq 82, \( \Delta_i \) has zero mean, the covariances shown below are nonzero\(^{22} \)

\[ \text{cov}(\Delta_i, \Delta_j) = \begin{cases} 
2D \Delta t - 2(2D \rho \Delta t - \sigma_{\text{loc}}^2) & \text{if } i = j \\
-(2D \rho \Delta t - \sigma_{\text{loc}}^2) & \text{if } i = j \pm 1 \\
0 & \text{otherwise} 
\end{cases} \quad (85) \]

where \( \Delta t = t_{i+1} - t_i \) and the motion blur coefficient, \( \rho \), is 1/6 for the case of continuous illumination though it can be adjusted for other “shutter functions”\(^{22} \)

The result of eq 85 implies that we can now use a Gaussian likelihood, \( p(D | \theta) \), for \( D \) independent differenced measurements \( \Delta_i \) with a value of \( \rho \) inserted into the covariance as an input rather than the measurements themselves in order to infer \( \theta = (D, \sigma_{\text{loc}}^2) \).\(^{89,102} \)

7.2.4. Kalman Filtering in Time Series Modeling and Tracking—In the previous subsection, we discussed approaches for analyzing simple parametric diffusion models accounting for the statistical effects of measurement noise (including both static localization and dynamic motion blur errors) through exact likelihood based frameworks.\(^{22,89} \)

As discussed in section 3, models with exact likelihoods simplify parameter estimation, goodness-of-fit hypothesis testing and confidence interval estimation. However, phenomena encountered in live cell SPT studies, such as molecular confinement\(^{57,60,299,344} \) or spatially varying forces,\(^{86,88} \) motivate new models that go beyond eq 82.

The Kalman filter (KF)\(^{110,345} \) directly addresses more complex problems even providing an alternative strategy to evaluating exact likelihoods for several parametric models relevant to SPT.\(^{25,86,346} \)

7.2.4.1. Parametric Models Amenable to Kalman Filtering: The standard KF can process the following parametric discrete time model\(^{110} \)

\[ r_{t+1} = Fr_t + \eta_t \quad (86) \]

\[ m_t = Ax_t + Hr_t + \varepsilon_t \quad (87) \]

In the above, the “state vector” at time \( t \), \( r_p \), is of size \( r \). It can be any quantity characterizing the system dynamics, such as position, velocity or angular orientation, though we will focus on position for concreteness.
The “measurement vector”, of size \( m \), is represented by \( m_t \) while \( x_t \), a vector of size \( k \), is a vector assumed to be known precisely (i.e., without measurement error). \( F, A, \) and \( H \) are matrices of dimensions \( r \times r, m \times k, \) and \( m \times r, \) respectively.

The vector \( \eta_t \) of size \( r \) is often referred to as “process noise” in control theory and signal processing.\(^{345} \) It is a zero mean random variable usually attributed to unresolved “thermal fluctuations” or any other source that directly affects the dynamics of \( r_t \).

By contrast, the vector \( \varepsilon_t \) of size \( m \) is the “measurement noise” that does not affect the dynamics of \( r_t \). Rather, \( \varepsilon_t \) prevents us from directly observing the precise instantaneous value of \( r_t \).

While it is possible to treat more general cases of correlated \( \eta_t \) and \( \varepsilon_t \),\(^ {345} \) here we assume that process and measurement noise vectors satisfy

\[
\begin{bmatrix}
\eta_t \eta_t^T \\
\varepsilon_t \varepsilon_t^T
\end{bmatrix} = \begin{cases}
Q & \text{for } t = \tau \\
0 & \text{otherwise}
\end{cases} \quad (88)
\]

\[
\begin{bmatrix}
\eta_t \eta_t^T \\
\varepsilon_t \varepsilon_t^T
\end{bmatrix} = \begin{cases}
R & \text{for } t = \tau \\
0 & \text{otherwise}
\end{cases} \quad (89)
\]

The KF algorithm can then be used for any parametric model, \( M \), that can be expressed using the \( \{ F, H, Q, R, A \} \) form where, to avoid introducing technical complications associated with unidentifiable parameter, \( H \) is often assumed known in many SPT and single molecule applications.\(^ {86,103,104,347} \)

7.2.4.2. Simple Example of the \( \{ F, H, Q, R, A \} \) Form: We provide an example demonstrating how we may translate a generic state-space model into the \( \{ F, H, Q, R, A \} \) form required by the KF algorithm. We show this for an order \( p \) autoregressive (linear state space) model, denoted by \( AR(p) \) which can be used in SPT when trajectories exhibit complex temporal “memory”.\(^ {346} \)

Once \( p \) is selected, a task which, in itself, may be cast as a model selection problem, the \( AR(p) \) is defined by the following relations

\[
r_{t+1} - \mu = \phi_1 (r_t - \mu) + \phi_2 (r_{t-1} - \mu) + \ldots + \phi_p (r_{t-p+1} - \mu) + \xi_t \quad (90)
\]

\[
[\xi_t \xi_t] = \begin{cases}
\sigma^2 & \text{for } t = \tau \\
0 & \text{otherwise}
\end{cases} \quad (91)
\]
The parameter vector here is simply \((\mu, \sigma^2, \phi_1, \phi_2, \ldots, \phi_p)\) though the model can be augmented to account for correlated noise as well (“autoregressive moving average model”) which we ignore here.

For \(p = 2\), say, the AR(2) model can be expressed using the \([F, H, Q, R, A]\) matrices required of the KF algorithm\(^{110}\)

\[
F = \begin{bmatrix}
\phi_1 & \phi_2 \\
1 & 0 \\
\end{bmatrix}, \quad Q = \begin{bmatrix}
\sigma^2 & 0 \\
0 & 0 \\
\end{bmatrix}, \quad R = [0], \quad H = [1 \Delta A \ldots 0], \quad A = [\mu].
\]

While we have only dealt with discrete time models, the procedure may also be adapted to treat continuous time SDEs of the form provided in eq 80,\(^{86,110}\)

7.2.4.3. The KF Algorithm in a Nutshell: Here we briefly review the basics of the KF algorithm in order to discuss its broader relevance to tracking in super-resolution.

Briefly, the KF yields a “state estimate”, i.e., the true position of the particle at time \(t\), \(r_t\), given all measurements up to time \(t\), \((m_1, \ldots, m_{t-1}, m_t)\), as well as all previous state estimates. As we now discuss, based on the description of the KF provided by Bishop et al.,\(^{348}\) the KF returns the state estimate of minimal variance (Figure 32). We should note that in this pedagogical description we take, without loss of generality, \(x_t\) in eq 87 to be zero at all times.

In order to accomplish its task, the KF begins from the state estimate \(\hat{r}_t^-\) superscripted with a minus sign to indicate that this state estimate is informed by all measurements before time \(t\), this is called the “KF prediction”.

We can also define a deviation of this estimate from its true (unknown) value

\[
e_t^- = r_t - \hat{r}_t^-	ag{92}
\]

with coinciding covariance

\[
P_t^- = [e_t^- \cdot e_t^-^T]\tag{93}
\]

Thus, given initial conditions for \(r_0^-\) and \(P_0^-\), we can use eqs 86 and 87 to compute both \(\hat{r}_t^-\) and \(P_t^-\).

In its simplest realization, the goal of the KF is to find the estimate the unknown \(r_t\), \(\hat{r}_t\), from a combination of \(\hat{r}_t^-\) and a correction, called an “innovation”, as follows
\[ \hat{r}_t = \hat{r}_t + K_t \cdot (m_t - H \hat{r}_t) \]  \tag{94}

where the innovation, \( (m_t - H \hat{r}_t) \), is preceded by a matrix, \( K \). This matrix, known as the “gain”, is derived from information encoded in the model parameters and will be determined shortly. A simple justification for eq 94 is described elsewhere in ref 349.

The matrix \( K \) is selected in order to guarantee the smallest possible covariance

\[ P_t = \left[ e_t \cdot e_t^T \right] \]  \tag{95}

where now \( e_t = r_t - \hat{r}_t \), where \( \hat{r}_t \) is defined in eq 94. By taking the expectation of \( e_t \cdot e_t^T \), we then minimize the resulting \( P_t \) to obtain \( K_t \) expressed in terms of the \( \{F, H, Q, R, A\} \) and \( P_t^- \).

From this \( K_t \), we now obtain a \( \hat{r}_t \) and minimal value for \( P_t \). If \( t \) is the last time step, we are done and we have now obtained the full estimate for the sequence of states, \( \hat{r}_t \). Otherwise, we use \( r_t \) to estimate \( \hat{r}_{t+1} \) and \( P_t \) to estimate \( P_{t+1}^- \) and continue onto the next time step.

**7.2.4. Some Extensions of the KF:** If the random vectors (i.e., \( \eta, \varepsilon, \) and \( r_i \)) of the dynamical model may be treated as multivariate Gaussian random variables, then we may compute an exact likelihood of the observable measurement time series conditioned on the model and parameter vector. This likelihood is denoted by \( p(m_1, m_2, \ldots, m_T | M, \theta) \) although \( r_t \) does not explicitly appear in the likelihood, both the measurement noise and the dynamics of the latent state variable \( r_i \) are encoded by \( M \) and \( \theta \). That is, the latent state dynamics induce temporal correlation structure in measurement time series. Despite the temporal correlation in the data, given a noisily measured times series \( \{m_1, m_2, \ldots, m_T\} \), one can obtain the exact MLE (i.e., no approximations of the likelihood required) of a state space model of the form given in eqs 86 and 87) simply by maximizing the log likelihood expression

\[ \hat{\theta} = \arg \max_{\theta} \log p(m_1, m_2, \ldots, m_T | M, \theta) \]

\[ = \sum_{i=1}^{T} \log p(m_i | m_1, \ldots, m_{i-1}, M, \theta) \]  \tag{96}

Each candidate value of \( \theta \) (for a fixed measurement sequence and \( M \)) can be used to evaluate a “cost function”. This cost function can be efficiently evaluated for multiple \( \theta \) values with the aid of the KF algorithm\(^{110}\) and aims at approximating the statistical properties of the entire measurement time series. The value of \( \theta \) maximizing the expression above, which can be obtained by a variety of nonlinear optimization routines such as Newton’s method or a Nelder–Mead search, is the exact MLE in multivariate Gaussian models which can be written in the form of eqs 86 and 87.\(^{110}\)
This “off-line” learning application of the KF is commonly used in econometrics but is less frequently used in mechanical engineering and tracking. In the latter domains, one often assumes \( \hat{\theta} \) (or a proxy) is known or given and simply uses the KF in “on-line learning” applications to recursively make estimates of \( \mathbf{r}_t \) as measurements are collected over time.

### 7.3. Bayesian Nonparametric SPT Analysis

Powerful nonparametric Bayesian methods, briefly introduced in section 3.4, have been used in time series analysis and show promise in live cell SPT trajectory analysis when model features, such as the number of diffusion coefficients sampled by a moving particle during the course of its trajectory, are a priori unknown.

These tools are especially helpful in capturing behaviors with spatially varying forces, autoregressive temporal correlation and measurement noise. In this subsection, we provide background to understand the so-called hierarchical Dirichlet process switching linear dynamical system (HDP-SLDS) since this technique shows promise in the analysis of optical microscopy trajectories.

HDP-SLDS is useful as the number of “modes” (or Markovian states), the time points at which the system hops from one state to the next (the change points), and the parameters needed to parametrize the model are learned in a data-driven fashion. As we have previously mentioned, this holistic model inference strategy, devoid of the postprocessing model selection step, is preferred and especially useful to the study of dynamics in living systems.

The relatively recent HDP-SLDS and other nonparametric Bayesian methods extending finite state HMMs have already found their way into single molecule and SPT data analysis applications.

In this subsection, we briefly review the difference between switching linear dynamical systems (SLDS) and HMMs (section 7.3.1) and introduce the concept of the hierarchical Dirichlet processes (HDPs), section 7.3.2. HDP-SLDS approaches are discussed in section 7.3.3 with details relegated to ref 186.

#### 7.3.1. Moving beyond Standard HMMs

Before we discuss nonparametric realizations of the SLDS or HMM, we briefly review these key models.

In short, in standard HMMs (detailed in section 3.3.3), dynamics are encoded by the latent variables or states \( s_t \) which are unobservable (Figure 33). However, once the latent states are estimated or given, the temporal order of the measurements is irrelevant to estimating the state’s “emission probability distribution”. That is, those parameters that set the probability distribution of the measurement given the state \( s_t \) at each time point, \( p(m_t|s_t) \). To each state \( s_t \) is associated a unique parameter vector \( \theta_t \) inferred from the data.

Switching linear dynamical system (SLDS) extend the standard HMM. In this review, we focus on Markovian SLDSs related to the interacting multiple model used in linking problems discussed earlier.
In the SLDS, the positions, $r_i$, are thought of as latent states and they themselves are controlled by an additional layer of latent variables (the “control” states), $s_i$. Both layers are required to explain the measurement (Figure 33). The SLDS associates to each $s_i$ a parameter vector $\theta = (v_i, x_i, D_i, R_i)$.

If $x_i$ and $R_i$ are zero for all states, differencing time adjacent measurements (as we did in section 7.2.3) from a common state would yield an iid sequence for which a classic HMM framework could be invoked to construct a likelihood and infer parameters. This approach was used by Persson et al.\textsuperscript{61}

However, if either $x_i$ or $R_i$ contain nonzero entries, as would be expected in live cell SPT where $R_i$ is always nonzero and the parameter $x_i$ characterizes molecular interaction and/or confinement forces,\textsuperscript{25,57,344} complex temporal correlations cannot be eliminated by differencing measurements. Despite this complication, ideas underlying the KF can nonetheless be successfully leveraged to compute marginal likelihoods, such as $p(D|\theta)$, and used in SPT trajectory analysis.\textsuperscript{62,66}

For both SLDSs and HMMs discussed here we have worked under the assumption that the number of states is specified. Yet, as we’ve also argued it is rarely clear for studies of dynamics inside cells how many states we should start with. For this reason, we now turn to nonparametric Bayesian techniques for generalizing finite state HMMs and SLDS to include an “infinite number of states” and let the data, confined by the model prior, introduce new states if warranted.

### 7.3.2. Infinite State Space Time Series Models via Hierarchical Dirichlet Processes

As we’ve discussed, a technical limitation facing classic HMMs is their reliance on a predefined number of states. To overcome this challenge, we introduce the infinite hidden Markov model (iHMM).\textsuperscript{350}

Prior to doing this, we reproduce for convenience the HMM model below

\begin{equation}
\begin{aligned}
s_1 & \sim p(s_1) \\
 s_i | s_{i-1} & \sim p(s_i | s_{i-1}) \\
 y_i | s_i, \theta & \sim p(y_i | s_i, \theta)
\end{aligned}
\end{equation}

Since the HMM is a frequentist model, it has no prior on the transition probabilities. Suppose we did want to introduce a prior on these transition probabilities at every time point. One possible prior would be a multinomial to determine to which state the transition probability brings us. To then select a prior on the weight of each transition probability, a convenient choice would be the conjugate to the multinomial, namely the Dirichlet distribution (section 3.4).

So far, none of our discussion is nonparametric as we have specified how many states, and thus how many transition probabilities we have. To generalize our discussion we need the infinite-dimensional generalization of the Dirichlet distribution, namely the Dirichlet process (DP).
Here, the DP prior is used to construct an HMM transition matrix, $p(s_t|s_{t-1})$, starting from $s_{t-1}$ and transitioning to any of an infinite number of states.

As discussed in section 3.4, the DP depends on two prespecified hyperparameters. The base distribution, $H$, and the concentration parameter, $\alpha$, which, roughly speaking, fixes how closely a draw from the DP approximates $H$.

If $H$ is a continuous distribution, as is often the case for nonparametric mixture models, the HMM or SLDS encounter a serious technical challenge as any algorithm that samples state transitions from a DP prior does not “share states”, i.e., it will never resample the same transition probability.

On the other hand if $H$ is discrete and finite, then we are back our parametric case. Instead what we need is an $H$ which is discrete and infinite.

To do this, we introduce the hierarchical DP (HDP), defined as follows

$$H \sim \text{GEM}(\gamma)$$
$$G \sim \text{DP}(\alpha, H)$$

where $\gamma$ is a hyperparameter that plays the role of a concentration parameter on the prior of the base distribution of the DP and $G$ is a distribution over transition probabilities. Since $H$ is now guaranteed to be a discrete distribution, it solves the “state sharing” problem discussed above while at the same time reaping the benefits of a continuous base measure $H$ indexing the states.

While MCMC sampling may be used to infer iHMM parameters, recent methods have improved computational efficiency by, for example, limiting the number of states sampled at each time point. Although this makes implementation easier, iHMMs have just barely begun to be explored in biophysics, more specifically smFRET and photobleaching.

Similarly to HMMs, SLDSs can also be generalized to the nonparametric case. The infinite state space SLDS models, HDP-SLDS have recently been used to improve “state switching” detection by explicitly accounting for measurement noise and confinement forces in simulated SPT trajectories and uncover nonlinear dynamics in chromatids of live yeast cells.

We refer the interested reader to refs 68 and 186 for the technical details of the HDP-SLDS and other infinite state dynamical models. While the mathematics are demanding, the problems solved by nonparametric methods help motivate the effort to become better acquainted with them. To this end, below we provide a simple simulated example of the power of nonparametric methods.

**7.3.3. Motivating Example for Using Nonparametric Bayesian Methods**—We consider a particle evolving in a time dependent harmonic potential, motivated by recent work where harmonic wells appear and disappear at different spatial locations mimicking...
the diffusion of biomolecules recently observed in neurons,\textsuperscript{87,353} where standard analysis tools assuming stationary potentials or a fixed number of states fail.

We first simulate a process having three kinetic states visited by a particle. Our goal will be to determine the number of states, all parameters associated with them and the times at which state changes occur from the simulated data.

All states are governed by a Ornstein–Uhlenbeck model, where particles diffuse within a harmonic potential, of the type shown in eq 80. The states are distinguished by unique parameter sets $\theta$ where each state $i$ has a $v_i$ which changes the location of the potential well’s minimum and particle diffusion coefficient, $D_i$.

Our three state model is

$$
\begin{array}{c}
\text{state } \alpha \\
\text{state } \zeta \\
\text{state } \gamma
\end{array}
\begin{array}{c}
D_\alpha \\
D_\zeta \\
D_\gamma
\end{array}
\begin{array}{c}
v_\alpha \\
v_\zeta \\
v_\gamma
\end{array}
\begin{array}{c}
D_\alpha \\
D_\zeta \\
D_\gamma
\end{array}
$$

(98)

In other words, state $\zeta$ has the same $\theta$ vector as state $\alpha$, but with $D_\zeta$ replacing $D_\alpha$. In the transition from state $\zeta$ to state $\gamma$, $v$ also changes. Throughout this simple example we use dimensionless quantities.

The top left panel of Figure 34 illustrates one such simulated 2D trajectory. The vertical red line indicates the location where the transition in $v$, and thus the hop from one potential surface to another (with the drift term in the SDE is given by $-\nabla U$), occurs.

The true state labels are shown below the trajectories in the top left panel and the two harmonic wells governing state $\gamma$ and $\zeta$ are shown in the right panels of Figure 34.

The bottom left panel shows the MSD curve computed using the observed 2D trajectory that would likely be classified “anomalous” using the guidelines of classic SPT analysis (Figure 30) rather than a stochastic switching between three different confined diffusion models. While the agreement between the estimated and true states is excellent, much of the success hinges on careful prior selections.\textsuperscript{62,66} The same is true of all Bayesian methods where priors properly informed by the physics of the problem provide a systematic way to improve inference.\textsuperscript{96}

While the example above may appear contrived, it may apply for instance to the diffusion of cargo transiently trapped in clathrin-coated pits.\textsuperscript{341} In this scenario, after a cargo molecule is trapped, it experiences “confined diffusion” until the pit vanishes or endocytosis occurs. If a pit vanishes, a new pit (a new potential minimum) may recruit the cargo. Since the HDP-SLDS model is not limited to a finite number of states, the approach can handle pits of different sizes, free-diffusion (unbound) states, as well as more realistic scenarios.

This flexibility fundamentally speaks to the power of the HDP-SLDS modeling approach\textsuperscript{68} as alternatives to anomalous diffusion in describing complex kinetics observed in vivo.
8. CONCLUDING REMARKS

Ten years after the publication of the seminal papers on super-resolution, single molecule super-resolution microscopy has become a standard lab technique now widely available across major imaging facilities. It has been used to study a number of biological targets just below the resolution of diffraction-limited microscopy, such as microtubules, mitochondria, and the nucleopore complex. Theoretical developments in interpreting super-resolution experiments, and single molecule experiments more broadly, have ushered data-driven methods into the physics and chemistry mainstream. While studies in live cells are motivating more general theoretical approaches borrowing heavily from statistical advances now feasible due to computing power, quantitative analysis efforts have also helped identify clear challenges standing in the way of greater modeling accuracy. Novel experimental methods have begun addressing some of these challenges such as phototoxicity and image distortions in thick heterogeneous samples though other key challenges such as labeling density and environment-dependent photophysical properties remain.

As this is a review of analysis methods, we highlight three broad directions that have been the focus of recent theoretical efforts. The first is on joint methods that simultaneously, and thus self-consistently, solve many problems at once such as problems in interpretation and counting or localization and linking. Such efforts reduce the number of user-dependent postprocessing steps albeit at a heavier computational cost. The second introduces problem-specific models, priors (whether theoretically or experimentally motivated), and algorithms suited to the particularities of the physical (or photophysical) challenge to reduce the computational burden and improve the prediction accuracy. The third is focused on generalizing models to accommodate the data’s complexity.

The picture of life emerging from breakthrough experimental techniques and analysis methods is one far richer in structural features, dynamics, and stochasticity than we could have conceived of even a decade ago. We envision a future in imaging where a combination of experiments and principled analysis provide a compelling narrative into the chaotic journey of life from the level of single molecules upward.

Acknowledgments

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Figure 1. Recent optical microscopy methods like lattice light sheet resolve structures within cells. In addition to imaging an entire structure, single particle tracking is also possible with this technology providing simultaneous information on both the “environment” and “dynamics” of tagged biomolecules. Reproduced with permission from ref 6. Copyright 2014, American Association for the Advancement of Science.
Figure 2.
Point emitter generates an Airy spot (a) with an intensity profile (b) given by eq 1. The 

wavenumber $k$, used in (b), is $2\pi/\lambda$. The intensity profile and diffraction spots were plotted 
using a simple Python script.
Figure 3. Microscope seen as a telescope. (a) A microscope’s resolution is determined by the numerical aperture NA of its objective, which is defined as the product of the index of refraction of the medium between the objective and the sample, n, and the sine of the half angular aperture, $\theta$. (b) A telescope’s angular resolution is determined by its (physical) aperture, $D$. The intensity profile and diffraction spots were plotted using a simple Python script.
Figure 4.
Fully resolved (a), barely resolved (b), and nonresolved (c) Airy diffraction spots. The intensity profile and diffraction spots were plotted using a simple Python script.
Figure 5.
Projections of an example 3D PSF on (a) the $x$-$y$ plane and (b) the $y$-$z$ plane.\textsuperscript{39} Reproduced with permission from ref 39. Copyright 2011, Nature Publishing Group.
Figure 6.
Two-color PALM can be used to address fundamental questions on molecular organization. (a) Adhesion proteins paxilin (green) and zyxin (red) are shown by two-color PALM to adopt different spatial distributions even though they appear colocalized in (b) diffraction limited microscopy. (c) A differential interference contrast image of the image shown in (b). Reproduced with permission from Hari Shroff, Catherine G. Galbraith, James A. Galbraith, Helen White, Jennifer Gillette, Scott Olenych, Michael W. Davidson, and Eric Betzig. “Dual-color superresolution imaging of genetically expressed probes within individual adhesion complexes.” Proceedings of the National Academy of Sciences 104, no. 51 (Dec 18, 2007): 20308–20313. Copyright (2007), National Academy of Sciences, U. S. A.
Figure 7. Measurement error affects particle localization. (a) 2D histograms contain photon numbers detected (by a CCD camera, say) along the x and y coordinates with the true molecular position ($r$) indicated by the point at which the solid vertical red line intersects with the x,y-plane. (b) In order to localize a static point emitter, the histogram may be fit to a PSF (e.g., a Gaussian or an Airy function) illustrated by a color surface contour. The estimate of the emitter’s position, $\mu$, shown by a solid vertical blue line, is approximated from the PSF. The difference between $\mu$ and $r$ is the localization error $\epsilon$. Factors such as photon count, exposure time, pixel size and background autofluorescence all contribute to $\epsilon$. 
Figure 8.
Motion and optics affect measurement noise. The simulated images illustrate how thermal motion, inherent to nanoscale measurements, affect the intensity profile of three molecules with different (given) diffusion coefficients $D$ (describing the diffusion of cytoplasmic, membrane-attached and fixed/immobile molecules) shown in the columns of panel (a). (a) The estimated point spread function (PSF) measured for different molecules under alternate illumination strategies (shown in the rows). The color bar displays the number of photons collected. (b) The illumination strategies are stroboscopic (top), continuous illumination with short camera exposure time (middle), continuous illumination with long camera exposure time (bottom). The two-state signals shown denote the state of the shutter (open/closed) and laser (on/off). The localization precision (accuracy and uncertainty) is affected by how accurately the PSF can be inferred from digital images produced by tagged molecules. The plots from panel (a) demonstrate how some molecular and sampling parameters affect measurement noise statistics when the PSF is used to infer the molecular position. “Smearing” of the PSF, induced by motion blur when a simulated molecule spans an area covered by different pixels in a single frame, introduces a new source of position uncertainty (in addition to the localization noise described in Figure 7). Reproduced with permission from ref 97. Copyright 2013, Nature Publishing Group.
Figure 9.
Posterior often becomes sharper as more data is used. Here data are sampled from a Poisson distribution with $\lambda = 5$ (designated by the dotted line) with $\mathbf{D} = \{2, 8, 5, 3, 5, 2, 5, 10, 6, 4\}$. The prior (eq 12 with $\alpha = 2$, $\beta = 1/7$) is plotted along with the resulting posterior informed by $N = 1$, $N = 5$, and $N = 10$ data points.
AIC and BIC are often both applied to step-finding. (a) We generated 1000 data points with a background noise level of standard deviation $\sigma_b = 20$. On top of the background, we arbitrarily added 6 dwells (5 points where the mean of the process suddenly changes) with noise around the signal having a standard deviation of $\sigma_s = 5$ (see inset). At this high noise level, and for this particular application, the BIC outperforms the AIC and the minimum of the BIC is at the theoretical value of 5 (dotted line). This is because the AIC overinterprets noise. All noise is Gaussian and uncorrelated. (b) Here we show another synthetic time trace with hidden steps, generated in the same manner as the data set in (a). The AIC (green) finds a model that overfits the true model (black), while the BIC (red) does not. However, as we increase the number of steps (while keeping the total number of data points fixed), the AIC does eventually outperform the BIC. This is to be expected. The AIC assumes the model could be unbounded in complexity and therefore does not penalize additional steps (model parameters) as much. The BIC, by contrast, assumes that there exists a true model of finite complexity. The data set for this simulation was generated via a Python implementation of the Gillespie algorithm, while both the BIC and AIC were implemented via elementary Python scripts.
Pixel dependent noise, if known, can be used to refine localization analysis. The left panel displays a simulated “variance map”, where the colorbar indicates the encoding of the noise level, similar to one that could be measured on an sCMOS camera. The color indicates that a single (red) pixel has higher variance, $\sigma$, than its neighbors. A MLE for the emitters’ location, shown in the middle panel, that ignores the pixel dependent variance, does not correctly identify the emitter locations that we placed on two parallel lines in this simulation; instead, an artifactual excess of events is detected at the “hot” pixel. On the right panel, we see that the MLE tailored for sCMOS cameras (“MLE-sCMOS”), which does take the variance map into account, avoids this artifact.\textsuperscript{198} Reproduced with permission from ref 198. Copyright 2013, Nature Publishing Group.
Figure 12.
Probability density of EMCCD camera readout counts can be highly non-Gaussian. Here we numerically evaluated eq 43 via a simple Python script for different mean photon numbers, $\lambda$ (0.32, 1.75, 7.0) and multiplication levels $m$ (9.8, 28). The dark count $c_0$ and readout noise standard deviation, $\sigma$, were set to 1000 and 10, respectively.
Figure 13.
Image of a point source by a microscope can be approximated as a Gaussian of standard deviation $s$. Collecting this image on a camera further pixelates it with pixel size $a$. The noisy intensity profile was generated using a simple Python script.
Figure 14.
Step-finding algorithms are widely applied across biophysics. (a) In this synthetic time trace containing 15 steps, the fluorescent signal-to-noise ratio decreases with time by an approximate factor of 3. Methods assuming constant noise statistics tend to overfit the start of the trace where the noise level is high (yellow line). A step-finding algorithm that does take into account variable noise performs significantly better (red line which overlaps with the theoretical noiseless trace used to generate this synthetic data).\(^96\) (b) An example of a real data trace (red line) showing an RNA hairpin zipping and unzipping (data was obtained through force spectroscopy).\(^268\) The offset green line are the steps found using a step-finding algorithm assuming constant noise.\(^269\) The time trace in (a) was created via a simple Gillespie algorithm with variable noise ratio. The code implementation of the method in ref 96 can be found online at https://github.com/lavrys/Photobleach. The code for the algorithm that does not take into account variable noise can be found at https://github.com/knyquist/KV_SIC. Adapted with permission from “Single Molecule Conformational Memory Extraction: P5ab RNA Hairpin.” (J. Phys. Chem. B 2014, 118, 6597−6603).\(^268\) Copyright 2014, American Chemical Society.
Figure 15.
Successful Bayesian step-finding algorithms used in photobleaching use priors informed by the fluorophore photophysics. Here we present a noisy synthetic photobleaching time trace (blue) and the underlying true noiseless signal (black) for a set of 10 fluorophores that eventually all photobleach. Too many steps are found from an algorithm that assumes a constant noise profile (red) grossly overfitting the start of the trace where the noise level is at its highest. Results (yellow) of the Bayesian algorithm proposed by Tsekouras et al., which is tailored to photobleaching and assumes that noise grows along with the number of active fluorophores, are much more accurate. Inset: a blow-up of the first 1000 data points in this time trace where noise is highest and constant-noise algorithms overfit. The data set was created with a simple Gillespie algorithm using a variable signal-to-noise ratio. The code implementation of the method in ref 96 can be found online at https://github.com/lavrys/Photobleach. The code for the algorithm that does not take into account variable noise can be found at https://github.com/knyquist/KV_SIC.
Figure 16.
It is impossible to regroup blinking events belonging to a single fluorophore using thresholds when signals from different fluorophores overlap or superpose. For illustrative purposes only, we depict the signal from one hypothetical PA-FP in orange and another in blue. Although methods based on grouping events using a blinking correction time cannot identify which event belongs to which fluorophore in such cases, this does not preclude them from yielding an accurate count.
Fluorophores may exist in many states but all single fluorophore states fall under two aggregated states: bright and dark. Here there are three states—inactive (I), photobleached (B), and dark (D)—making up the dark aggregate and only one state, active (A), making up the bright aggregate. For this example, with a single dark state there are four rates that fully characterize the photophysics, \( \{k_a, k_d, k_r, k_b\} \). Reproduced with permission from Geoffrey C. Rollins, Jae Yen Shin, Carlos Bustamante, and Steve Pressé. “Stochastic approach to the molecular counting problem in superresolution microscopy.” Proceedings of the National Academy of Sciences 112, no. 2 (Dec 22, 2014): E110-E118.
Figure 18. Results of the method described by Rollins et al.\textsuperscript{21} employed on synthetic data sets. (a) Synthetic time traces were generated using Gillespie algorithm with photophysical rates $k_b = 1.0$, $k_d = 10.0$, $k_r = 0.1$, $k_a = 0.5$ (see Figure 17) for each of 5 fluorophores. In maximizing the likelihood, photophysical rates as well as the fluorophore numbers may be simultaneously inferred. The results of this inference (shown as histograms) closely match the theoretical expected answer (dotted line) despite the fact that overlapping events were very probable for this choice of photophysical parameters. (b) Same as in (a), except that the effect of missed events has been taken into consideration, resulting in a substantial improvement in accuracy. Assumptions made in creating these data sets are discussed in the original publication, ref 21. Adapted with permission from Geoffrey C. Rollins, Jae Yen Shin, Carlos Bustamante, and Steve Pressé. “Stochastic approach to the molecular counting problem in superresolution microscopy.” Proceedings of the National Academy of Sciences 112, no. 2 (Dec 22, 2014): E110–E118.
Figure 19.
Events of duration shorter than the data acquisition interval can be missed when collecting data, introducing error into the analysis. If a fluorophore blinks faster than the data acquisition rate, the resulting brief dip or rise in fluorescence intensity (denoted with arrows above) will be missed when collecting data. Such missed transitions can be accounted for in inferring photophysical parameters and fluorophore counts.\textsuperscript{21}
Figure 20. Algorithms that take into consideration multiple frames are more computationally demanding but achieve better results than those taking only successive frames into account. (a) The dynamics of cell-surface receptors were followed and rare merging (blue diamond) and splitting (red diamond) events were monitored. Global rather than local linking methods are needed here to quantify such events which—despite their rarity—are critical in understanding signal transduction.283 (b and c) Comparison between (b) an algorithm that takes into account information across space and time and (c) an algorithm that only considers two frames at a time. The latter algorithm yields many shorter, broken trajectories by failing to link segments (different colors). Dotted lines represent virtual detections, i.e., locations traced by a particle temporarily disappeared. See Danuser283 for details. Reproduced with permission from ref 283. Copyright 2011, Cell Press.
Particles may merge and/or split, leading to trajectories that are temporarily joined. This presents new challenges in identifying pre- and post- colocalization track segments for the same particle. Top: example of a trajectory segment where two formyl peptide receptor particles (denoted by the green and red arrows at 0, 12 and 24 ms) merge (at 36 ms) and remain colocalized until 96 ms (yellow arrow), only to split again at 108 ms (blue and gray arrows). Bottom: Particle trajectories followed above. Trajectories colored to match the corresponding colored arrows in the figure above. Adapted with permission from ref 290. Copyright 2014, Nature Publishing Group.
Figure 22.
Motion models can be used to limit the search area for linking algorithms.\textsuperscript{311} (a) The two parabolas represent the past and future volumes in which a particle localized at their intersection may be found on the basis of a motion model. Limiting the search volume reduces the computational cost of the method. (b) Once a trajectory segment has been determined, it is possible to use different motion models in different spatiotemporal areas of the trajectory, as shown by the different shapes of the paraboloid search regions. In this case motion in the later end of the trajectory (designated “head”) is considerably more constrained than motion in the early part of the trajectory (designated “tail”). Reproduced with permission from ref 311. Copyright 2005, IEEE Computer Society.
Figure 23.
Biological images can be cluttered to the point where near-neighbor algorithms are bound to fail. (a) Superposition of two successive frames in synthetic data mimicking particle transport via turbulent fluid flow. Particles in the first and second frame represented as light-colored spheres and dark-colored diamonds, respectively. (b) Trajectory track segments derived for particles in these two frames via the “passing messages between frames” method.\textsuperscript{291} It is obvious that nearest-neighbor algorithms would not work well for this data. Reproduced with permission from Michael Chertkov, Lukas Kroc, F. Krzakala, M. Vergassola, and L. Zdeborová. “Inference in particle tracking experiments by passing messages between images.” Proceedings of the National Academy of Sciences 107, no. 17 (Apr 27, 2010): 7663–7668.
Figure 24.
Effects of blinking, or temporary disappearance of particles, present linking challenges. (a) No blinking. Three trajectories are associated with three localized points ($P_1$, $P_2$, and $P_3$). Motion models constrain the possible areas within which each of the localized points terminating the trajectory could have moved, but ambiguity nevertheless exists. Thus, $P_1$ can potentially be the terminal point for either red and yellow trajectories, and similarly for $P_2$. The problem can only be resolved because $P_3$ belongs to the blue trajectory, meaning that $P_2$ and $P_1$ must belong to the red and yellow trajectories, respectively. (b) Blinking in the red trajectory before the current frame. Because no particle was detected for the red trajectory at frame $t$, the potential area in which the particle may be found in frame $t+1$ is correspondingly larger. However, assignment of localized points at frame $t+1$ to trajectories is still resolvable. On the contrary, combinations of blinks or, alternatively, blinks lasting longer than one time frame can lead to important misassignments.$^{317}$ Adapted with permission from ref 317. Copyright 2008, Nature Publishing Group.
Figure 25.
Contextual information on a tagged particle’s microenvironment may be used in linking. In the top panels above, a lysosome (red trajectory) was transported by molecular motors on a microtubule network (green background). By using super-resolution to image the microtubule network with high precision (lower panels), it became possible to closely match a lysosome trajectory to a particular microtubule and quantify when the lysosome detached or switched between microtubules. Reproduced with permission from Štefan Bálint, Ione Verdeny Vilanova, Ángel Sandoval Álvarez, and Melike Lakadamyali. “Correlative live-cell and superresolution microscopy reveals cargo transport dynamics at microtubule intersections.” Proceedings of the National Academy of Sciences 110, no. 9 (Feb 26, 2013): 3375–3380.
A tree can be constructed by linking particle positions. Here a trajectory that has been established from time frame $k - 5$ to $k - 1$ is examined for possible extensions to time frames $k$ and $k + 1$. $z$ refers to a measurement (circle) or a predicted virtual particle position (square). In the $z_{ij}(j)$ notation $j$ is the time frame and $i$ enumerates the possible particle locations, real or virtual. All possible track extensions are considered, whether they consist of real detected particles or virtual detections. See Chenouard et al.\textsuperscript{287} for details.

Reproduced with permission from ref 287. Copyright 2013, IEEE Computer Society.
Figure 27. Temporal and spatial blur have different PSFs. While spatial blur (left) can be well approximated by a Gaussian, temporal blur (right) is a “box-shaped” pulse that cannot reliably be considered Gaussian. Algorithms that make Gaussian assumptions may therefore fail when used to temporally localize particles.\textsuperscript{329} Reproduced with permission from ref 329. Copyright 2005, IEEE Computer Society.
Sampling rates and “connect-the-dots” schemes can bias interpretation. Consider two particles that are undergoing oscillatory motion with the same amplitude and frequency but are slightly off-phase (top and middle, blue curves). At a frame rate less than the period of the particle motion, observations (top and middle, yellow curves) suggest an oscillatory model with the same amplitude but a different frequency. This effect, where some motion characteristics are inferred correctly while others are not, is called “aliasing” and is a key challenge of temporal super-resolution.
Figure 29.
Multiscale super-resolution image data reveals critical contextual information. The plot displays a 3D multiresolution microscopy measurement of a particle approaching the membrane of a fibroblast cell. The fast diffusing nanoparticle was tracked using a 3D SPT module (the moving particle was kept in the focus of the objective via a piezoelectric stage and an optical feedback loop running at 100 kHz). The background (cell membrane, nucleus, etc.) was tracked with a conventional two-photon scanning fluorescence microscope. Valuable “contextual information” along single particle tracks is obtained from the combination of these two modalities. The arrows show trajectory segments where the nanoparticle traced out structured protrusions on the cell boundary. Reproduced with permission from ref 15. Copyright 2015, Royal Society of Chemistry.
Figure 30.
Idealized MSD curves. The theoretical limits of some common motion models used in single particle tracking (SPT). Graphical illustration of information contained in Table 1.
vbSPT is used to determine the number of diffusive states arising from changes in interaction of RNA helper protein (Hfq) in *E. coli* treated with rifampicin. (a) Trajectories sample three different diffusion coefficients identified as differently colored trajectories in the untreated cell. The transition probabilities linking each diffusive state along with each state’s diffusion coefficient is inferred and represented in the kinetic scheme shown below. (b) The slowest diffusive state from (a) vanishes for *E. coli* treated with rifampicin, a transcription inhibitor, suggesting that this state coincided with an interaction of Hfq with RNA. Throughout this figure, $\Delta t = 300$ Hz. In addition, the scale bar indicates 0.5 μm. See details in Chenouard et al. Reproduced with permission from ref 61. Copyright 2013, Nature Publishing Group.
Figure 32.
Traditional application of the Kalman Filter (KF). The underlying trajectory evolving in continuous time is shown by the solid black line. The observed or measurement data consists of discretely sampled data corrupted by measurement noise shown with blue circles along with their corresponding 95% confidence interval. The KF exploits knowledge of the model, its parameters and noise measurement in order to obtain improved estimates of position. The data set was created via an implementation of an Ornstein–Uhlenbeck stochastic differential equation, while processing was performed by a simple implementation of the KF, both in MATLAB.
Figure 33.
Dependencies of HMM and Markovian SLDS models. In both HMM and SLDS, unobservable states (latent variables) are denoted by white circles. The observations are denoted by filled circles. The subscripts index time and arrows connecting variables indicates conditional statistical dependencies between variables. For instance in the HMM, the latent variables determine the observations. While observations themselves do not affect observations at different times, the same is not true of the states. That is, once the states of the HMM are determined, past observed measurements no longer help predict subsequent measurements. The SLDS adds an additional layer of latent variables which determine the “position” estimates ($r_t$). By contrast to the HMM, in the SLDSs, both “position” estimates ($r_t$) and state estimates can be used to predict future observation.
Nonparametric Bayesian analysis is helpful to SPT trajectory analysis. (a) Position trajectory of a 2D simulation illustrating nonlinear time dependent kinetics; the bottom portion of the plot displays the true state label as well as the number of states (and label) inferred by the HDP-SLDS algorithm. The vertical red line denotes the point at which the potential abruptly changes; potential before and after this change are plotted in rightmost panels. (b) The potential energy surface governing the particle dynamics before an abrupt state change in the position of the potential energy’s global minimum location at \( t = 40 \) s. The \( x \) and \( y \) coordinates of the measured position are plotted; the portion of the trajectory governed by the potential plotted in this panel is colored in red. (c) Potential energy surface of trajectory after an abrupt state change in the position corresponding to the potential energy’s minimum location at \( t = 40 \) s. The \( x \) and \( y \) coordinates of the measured position are plotted; the portion of the trajectory subject to the potential plotted in this panel is colored in black. (d) The MSD curve of the simulated trajectory exhibiting transitions between three types of diffusive motion; the MSD was computed using all measurements shown in the top left panel. See text for additional discussion. The data set for this simulation was created by a simple Gillespie algorithm while the labels and number of states were found via a simple MATLAB implementation of the HDP-SLDS algorithm. Standard MATLAB functions were used to calculate the MSD.
Table 1
Mean Square Displacement (MSD) Curves Continue to Be Used to Drawn Ensemble Information from SPT Measurements\textsuperscript{60,65a}

<table>
<thead>
<tr>
<th>qualitative diffusion type</th>
<th>MSD Signature</th>
</tr>
</thead>
<tbody>
<tr>
<td>normal diffusion</td>
<td>$\langle r^2(\tau) \rangle = 2dD\tau$</td>
</tr>
<tr>
<td>confined/corralled diffusion</td>
<td>$\langle r^2(\tau) \rangle = C_1(1 - C_2 \exp(-C_3 \tau))$</td>
</tr>
<tr>
<td>immobile/stationary</td>
<td>$\langle r^2(\tau) \rangle &lt; \sigma_{\text{loc}}^2$</td>
</tr>
<tr>
<td>directed diffusion (constant velocity $v$)</td>
<td>$\langle r^2(\tau) \rangle = 2dD\tau + v^2 \tau^2$</td>
</tr>
<tr>
<td>anomalous subdiffusion ($\alpha &lt; 1$)</td>
<td>$\langle r^2(\tau) \rangle = C\tau^\alpha$</td>
</tr>
</tbody>
</table>

\textsuperscript{a}In the right column, we find the predicted theoretical MSD trend associated with ideal sampling (“ideal” is used since both infinite sample sizes and no measurement noise effects are included) where the $C_i$’s coincide with constants determined from the assumed dynamical model, $\alpha$ is an anomalous exponent, $D$ corresponds to the normal diffusion coefficient, $d$ to the dimension of the trajectory (e.g. such as two or three common in SPT), and $\sigma_{\text{loc}}^2$ corresponds to the measurement/localization noise variance.
Table 2
Values for the Parameter Values Reproducing the Limits of Table 1<sup>a</sup>

<table>
<thead>
<tr>
<th>qualitative diffusion type</th>
<th>SDE constraints in eq 80</th>
</tr>
</thead>
<tbody>
<tr>
<td>standard diffusion</td>
<td>$D &gt; 0$, $\kappa = 0$, $v = 0$</td>
</tr>
<tr>
<td>confined/corralled diffusion</td>
<td>$D &gt; 0$, $\kappa &gt; 0$</td>
</tr>
<tr>
<td>immobile/stationary</td>
<td>$D = 0$, $v = 0$, $\kappa = 0$</td>
</tr>
<tr>
<td>directed diffusion (constant velocity $v$)</td>
<td>$D &gt; 0$, $v \neq 0$</td>
</tr>
</tbody>
</table>

<sup>a</sup> Left column coincides with the left column of Table 1. The right column lists the constraints on $\theta \equiv (v, \kappa, D)$ corresponding to those parameters from eq 80 needed to reproduce the MSD shown on the right column of Table 1 under the same $N \to \infty$ limit.