

**Lévy-Student distributions for halos in accelerator beams**

Nicola Cufaro Petroni\*

*Dipartimento di Matematica dell'Università di Bari and INFN Sezione di Bari, via E. Orabona 4, 70125 Bari, Italy*Salvatore De Martino,<sup>†</sup> Silvio De Siena,<sup>‡</sup> and Fabrizio Illuminati<sup>§</sup>*Dipartimento di Fisica dell'Università di Salerno, INFN Unità di Salerno, and INFN Sezione di Napoli Gruppo collegato di Salerno, Via S. Allende, I-84081 Baronissi (SA), Italy*

(Received 1 September 2005; published 21 December 2005)

We describe the transverse beam distribution in particle accelerators within the controlled, stochastic dynamical scheme of stochastic mechanics (SM) which produces time reversal invariant diffusion processes. This leads to a linearized theory summarized in a Schrödinger-like (SL) equation. The space charge effects have been introduced in recent papers by coupling this S-L equation with the Maxwell equations. We analyze the space-charge effects to understand how the dynamics produces the actual beam distributions, and in particular we show how the stationary, self-consistent solutions are related to the (external and space-charge) potentials both when we suppose that the external field is harmonic (constant focusing), and when we *a priori* prescribe the shape of the stationary solution. We then proceed to discuss a few other ideas by introducing generalized Student distributions, namely, non-Gaussian, Lévy infinitely divisible (but not stable) distributions. We will discuss this idea from two different standpoints: (a) first by supposing that the stationary distribution of our (Wiener powered) SM model is a Student distribution; (b) by supposing that our model is based on a (non-Gaussian) Lévy process whose increments are Student distributed. We show that in the case (a) the longer tails of the power decay of the Student laws and in the case (b) the discontinuities of the Lévy-Student process can well account for the rare escape of particles from the beam core, and hence for the formation of a halo in intense beams.

DOI: [10.1103/PhysRevE.72.066502](https://doi.org/10.1103/PhysRevE.72.066502)

PACS number(s): 29.27.Bd, 41.75.Lx, 02.50.Ey, 05.40.Fb

**I. INTRODUCTION**

In high-intensity beams of charged particles, proposed in recent years for a wide variety of accelerator-related applications, it is very important to keep at low level the beam loss to the wall of the beam pipe, since even small fractional losses in a high-current machine can cause exceedingly high levels of radioactivation. It is now widely believed that one of the relevant mechanisms for these losses is the formation of a low-intensity beam halo more or less far from the core. These halos have been observed [1] or studied in experiments [2], and have also been subjected to an extensive simulation analysis [3]. For the next generation of high-intensity machines it is, however, still necessary to obtain a more quantitative understanding not only of the physics of the halo, but also of the beam transverse distribution in general [4]. In fact “because there is not a consensus about its definition, halo remains an imprecise term” [5] so that several proposals have been put forward for its description.

Charged particle beams are usually described in terms of classical, deterministic dynamical systems. The standard model is that of a collisionless plasma where the corresponding dynamics is embodied in a suitable phase space (see for example [6]). In this framework the beam is studied by

means of the *particle-in-core* model and the simulations show that the instabilities due to a parametric resonance can allow the particles to escape from the core with consequent halo formation [3–5]. The present paper takes a different approach: it follows the idea that the particle trajectories are samples of a stochastic process, rather than the usual deterministic (differentiable) trajectories. In the usual dynamical models there is a particle probability distribution obeying the Vlasov equation, and its evolution is Liouvillian in the sense that the origin of the randomness is just in the initial conditions: along the time evolution, which is supposed to be deterministic, there is no new source of uncertainty. It is the nonlinear character of the equations that produces the possible unpredictable character of the trajectories. On the other hand, in our model the trajectories are replaced by stochastic processes since the time evolution is supposed to be randomly perturbed even after the initial time. It is open to discussion which one of these two descriptions is more realistic; in particular we should ask if the mutual interactions among the beam particles look like random collisions, or rather like continuous deterministic interactions. In the opinion of the authors, however, a plasma (with collisions) described in terms of controlled stochastic processes is a good candidate to explain the rare escape of particles from a quasisustainable beam core by statistically taking into account the random interparticle interactions that cannot be described in detail. Of course the idea of a stochastic approach is hardly new [6,7], but there are several different ways to implement it.

First of all let us remark that the system we want to describe is endowed with some measure of invariance under

\*Electronic address: [cufaro@ba.infn.it](mailto:cufaro@ba.infn.it)†Electronic address: [demartino@sa.infn.it](mailto:demartino@sa.infn.it)‡Electronic address: [desiena@sa.infn.it](mailto:desiena@sa.infn.it)§Electronic address: [illuminati@sa.infn.it](mailto:illuminati@sa.infn.it)

time reversal, since the external fields act to keep it in a quasistationary nondiffusive state despite the repulsive electromagnetic (em) interactions among the constituent particles. However, a widespread misconception notwithstanding, the theory of stochastic processes does not always describe irreversible systems: the addition of dynamics to a stochastic kinematics can in fact ascribe a measure of time reversal invariance also to a stochastic system [8]. The standard way to build a stochastic dynamical system is to modify the phase space dynamics by adding a Wiener noise  $\mathbf{B}(t)$  to the momentum equation only, so that the usual relations between position and velocity are preserved:

$$m d\mathbf{Q}(t) = \mathbf{P}(t)dt,$$

$$d\mathbf{P}(t) = \mathbf{F}(t)dt + \beta d\mathbf{B}(t).$$

In this way we get a derivable but not Markovian position process  $\mathbf{Q}(t)$ . The standard example of this approach is that of a Brownian motion in a force field described by an Ornstein-Uhlenbeck system of stochastic differential equations (SDEs) [9]. Alternatively we can add a Wiener noise  $\mathbf{W}(t)$  with diffusion coefficient  $D$  to the position equation:

$$d\mathbf{Q}(t) = \mathbf{v}_{(+)}(\mathbf{Q}(t), t)dt + \sqrt{D} d\mathbf{W}(t)$$

and get a Markovian but not derivable  $\mathbf{Q}(t)$ . In this way the stochastic system is also reduced to a single SDE since we are obliged to drop the second (momentum) equation: in fact now  $\mathbf{Q}(t)$  is no longer derivable. The standard example of this reduction is the Smoluchowski approximation of the Ornstein-Uhlenbeck process in the overdamped case [9]. As a consequence we will work only in a configuration, and not in a phase space; but this does not prevent us from introducing dynamics—as we will show in Sec. II A—either by generalizing the Newton equations [9,10], or by means of a stochastic variational principle [11]. We remark that in this scheme the forward velocity  $\mathbf{v}_{(+)}(\mathbf{r}, t)$  can no longer be a field given *a priori*: rather it now plays the role of a new dynamical variable of our system. This second scheme, stochastic mechanics (SM), is universally known for its original application to the problem of building a classical stochastic model for quantum mechanics (QM), but in fact it is a very general model which is suitable for a large number of stochastic dynamical systems [8,12]. We will also see in Sec. II A that from the stochastic variational principles two coupled equations are derived which are equivalent to a Schrödinger-like (SL) differential equation: in this sense we will speak of quantumlike (QL) systems, in analogy with other recent research on this subject [13,14]. In fact the SM can be used to describe every stochastic dynamical system satisfying fairly general conditions: it has been known for long a time [15], for example, that for any given diffusion there is a correspondence between diffusion processes and solutions of SL equations where the Hamiltonians come in general from suitable vector potentials. Under some regularity conditions this correspondence is seen to be one to one. The usual Schrödinger equation, and hence QM, is recovered when the diffusion coefficient coincides with  $\hbar/2m$ , namely, is connected to the Planck constant. However we are interested

here not in a stochastic model of QM, but in the description of particle beams.

In the present paper we intend to widen the scope of our SM model by introducing the idea that an important role for the beam dynamics can be played by non-Gaussian Lévy distributions. In fact these distributions enjoyed a widespread popularity in recent years because of their multifaceted possible applications to a large set of problems from statistical mechanics to mathematical finance (see, for example, [8,16] and references quoted therein). In particular the so called *stable* laws (see Sec. II A) are used in a large number of instances, as for example in the definition of the so-called Lévy flights. Our research is instead focused on a family of non-Gaussian Lévy laws that are *infinitely divisible* but not stable: the generalized Student laws. As will be discussed later this will allow us to overcome—without resorting to the trick of truncated laws—the problems raised by the fact that the stable non-Gaussian laws always have divergent variances: a feature that is not realistic to ascribe to most real systems. It is possible to show indeed that by suitably choosing the parameters of the Student laws we can have distributions with finite variance, and approximating the Gaussian law as well as we want. On the other hand, the infinitely divisible character of these laws is all that is required to build a stationary, stochastically continuous Markov process with independent increments, namely, the Lévy process that we propose to use to represent the evolution of our particle beam.

Of course it is not always mathematically easy to deal with infinitely divisible processes, but we will show that at least in two respects they will help us to have some further insight into the beam dynamics. First of all we use the Student distributions in the framework of the traditional SM where the randomness of the process is supplied by a Gaussian Wiener noise: here we examine the features of the self-consistent potentials that can produce a Student distribution as the stationary transverse distribution of a particle beam. In this instance the focus of our research is on the increase of the probability of finding the particles at a great distance from the beam core. Then we pass to the definition of a true *Lévy-Student process*, and we show with a few simulations that these processes can help to explain how a particle can be expelled from the bunch because of some kind of hard collision. In fact the trajectories of our Lévy-Student process show the typical jumps of the non-Gaussian Lévy processes: a feature that we propose to use as a model for the halo formation. It is worth remarking that, albeit the more recent empirical data about halos [17] are still not accurate enough to distinguish between the suggested distributions and the usual Gaussian ones, our conjecture on the role of Student laws in the transverse beam dynamics has recently found a first confirmation [18] in numerical simulations showing how these laws are well suited to describe the statistics of the random features of the particle paths.

In a few previous papers [19] we connected the (transverse) rms emittance to the characteristic microscopic scale and to the total number of the particles in a bunch, and implemented a few techniques of active control for the dynamics of the beam. In this paper we first of all review the theoretical basis [19,20] of the proposed model: in Sec. II we

define our SM model with emphasis added to the potentials that control the beam dynamics and to the possible nonstationary solutions of this model [21]. In Sec. III we review our analysis of the self-consistent, space-charge effects due to the em interaction among the particles, adding a few further results and comments. In Sec. IV we then discuss the idea [22] that the laws ruling the transverse distribution of particle beams are non-Gaussian, infinitely divisible, Lévy laws like the generalized Student laws. In particular we analyze the behavior of our usual SM model under the hypothesis that the stationary transverse distribution is a Student law. Finally, in Sec. V we study the possibility of extending our SM model to Lévy processes whose increments are distributed according to the Student law. We think in particular that the presence of isolated jumps in the trajectories can help to build a realistic model for the possible formation of halos in the particle beams. We end the paper with a few concluding remarks.

## II. STOCHASTIC BEAM DYNAMICS

### A. Stochastic mechanics

First of all we introduce the stochastic process performed by a representative particle that oscillates around the closed ideal orbit in a particle accelerator. We consider the three-dimensional (3D) diffusion process  $\mathbf{Q}(t)$ , taking the values  $\mathbf{r}$ , which describes the position of the representative particle and whose probability density is proportional to the particle density of the bunch. As stated in Sec. I the evolution of this process is ruled by the Itô stochastic differential equation

$$d\mathbf{Q}(t) = \mathbf{v}_{(+)}(\mathbf{Q}(t), t)dt + \sqrt{D}d\mathbf{W}(t), \quad (1)$$

where  $\mathbf{v}_{(+)}(\mathbf{r}, t)$  is the forward velocity, and  $d\mathbf{W}(t) \equiv \mathbf{W}(t+dt) - \mathbf{W}(t)$  is the increment process of a standard Wiener noise  $\mathbf{W}(t)$ ; as is well known this increment process is Gaussian with law  $\mathcal{N}(0, \mathbb{I}dt)$ , where  $\mathbb{I}$  is the  $3 \times 3$  identity matrix. Finally the diffusion coefficient  $D$  is supposed to be constant: the quantity  $\alpha = 2mD$ , which has the dimensions of an action, will be later connected to the characteristic transverse emittance of the beam. Equation (1) defines the random kinematics performed by the particle, and replaces the usual deterministic kinematics

$$d\mathbf{q}(t) = \mathbf{v}(\mathbf{q}(t), t)dt \quad (2)$$

where  $\mathbf{q}(t)$  is just the trajectory in the 3D space.

To counteract the dissipation due to this stochastic kinematics, a dynamics must be independently added. In SM we do not have a phase space: our description is entirely in a 3D configuration space. This means in particular that the dynamics is not introduced in a Hamiltonian way, but by means of a suitable stochastic least action principle [11] obtained as a generalization of the variational principle of classical mechanics. In the following we will briefly review the main results, referring for details to Refs. [8,9,11]. Given the SDE (1), we consider the probability density function (PDF)  $\rho(\mathbf{r}, t)$  associated with the diffusion  $\mathbf{Q}(t)$  so that, besides the forward velocity  $\mathbf{v}_{(+)}(\mathbf{r}, t)$ , we can now define a backward velocity

$$\mathbf{v}_{(-)}(\mathbf{r}, t) = \mathbf{v}_{(+)}(\mathbf{r}, t) - 2D \frac{\nabla \rho(\mathbf{r}, t)}{\rho(\mathbf{r}, t)}. \quad (3)$$

We can then introduce also the current and the osmotic velocity fields, defined as

$$\mathbf{v} = \frac{\mathbf{v}_{(+)} + \mathbf{v}_{(-)}}{2}, \quad \mathbf{u} = \frac{\mathbf{v}_{(+)} - \mathbf{v}_{(-)}}{2} = D \frac{\nabla \rho}{\rho}. \quad (4)$$

Here  $\mathbf{v}$  represents the velocity field of the density, while  $\mathbf{u}$  is of intrinsic stochastic nature and is a measure of the nondifferentiability of the stochastic trajectories.

A first consequence of the stochastic generalization of the least action principle [9,11] is that the current velocity takes the following irrotational form:

$$m\mathbf{v}(\mathbf{r}, t) = \nabla S(\mathbf{r}, t), \quad (5)$$

while the Lagrange equations of motion for the density  $\rho$  and for the current velocity  $\mathbf{v}$  are the continuity equation associated with every stochastic process

$$\partial_t \rho = -\nabla \cdot (\rho \mathbf{v}) \quad (6)$$

and a dynamical equation

$$\partial_t S + \frac{m}{2} \mathbf{v}^2 - 2mD^2 \frac{\nabla^2 \sqrt{\rho}}{\sqrt{\rho}} + V(\mathbf{r}, t) = 0, \quad (7)$$

which characterizes our particular class of time-reversal invariant diffusions (Nelson processes). The last equation has the same form as the Hamilton-Jacobi-Madelung (HJM) equation, originally introduced in the hydrodynamic description of quantum mechanics by Madelung [23]. Since Eq. (5) holds, the two equations (6) and (7) can be put in the following form:

$$\partial_t \rho = -\frac{1}{m} \nabla \cdot (\rho \nabla S), \quad (8)$$

$$\partial_t S = -\frac{1}{2m} \nabla S^2 + 2mD^2 \frac{\nabla^2 \sqrt{\rho}}{\sqrt{\rho}} - V(\mathbf{r}, t), \quad (9)$$

which now constitutes a coupled, nonlinear system of partial differential equations for the pair  $(\rho, S)$  that completely determines the state of our beam. On the other hand, because of Eq. (5), this state is equivalently given by the pair  $(\rho, \mathbf{v})$ .

It can also be shown by simple substitution from Eq. (4) that Eq. (6) is equivalent to the standard Fokker-Planck (FP) equation

$$\partial_t \rho = -\nabla \cdot (\mathbf{v}_{(+)} \rho) + D \nabla^2 \rho \quad (10)$$

formally associated with the Itô equation (1). In fact also the HJM equation (7) can be cast in a form based on  $\mathbf{v}_{(+)}$  rather than on  $\mathbf{v}$ , namely;

$$\partial_t S = -\frac{m}{2} \mathbf{v}_{(+)}^2 + mD \mathbf{v}_{(+)} \cdot \nabla \ln f + mD^2 \nabla^2 \ln f - V \quad (11)$$

where  $f$  is a dimensionless density defined by

$$\rho(\mathbf{r}, t) = Cf(\mathbf{r}, t) \quad (12)$$

where  $C$  is a dimensional constant. On the other hand, from Eqs. (3) and (4), we know that also the forward velocity  $\mathbf{v}_{(+)}$  is irrotational:

$$\mathbf{v}_{(+)}(\mathbf{r}, t) = \nabla W(\mathbf{r}, t), \quad (13)$$

and that by taking Eq. (5) into account the functions  $W$  and  $S$  are connected by the relation

$$S(\mathbf{r}, t) = mW(\mathbf{r}, t) - mD \ln f(\mathbf{r}, t) - \theta(t) \quad (14)$$

where  $\theta$  is an arbitrary function of  $t$  only.

The time-reversal invariance is now made possible [10] by the fact that the forward drift velocity  $\mathbf{v}_{(+)}(\mathbf{r}, t)$  is no longer an *a priori* given field, as is usual for the diffusion processes of the Langevin type; instead it is dynamically determined at any instant of time, starting by an initial condition, through the HJM evolution equation (7). It is finally important to remark that, introducing the representation [23]

$$\Psi(\mathbf{r}, t) = \sqrt{\rho(\mathbf{r}, t)} e^{iS(\mathbf{r}, t)/\alpha} \quad (15)$$

(with  $\alpha = 2mD$ ) the coupled equations (8) and (9) are made equivalent to a single linear equation of the form of the Schrödinger equation, with the Planck action constant replaced by  $\alpha$ :

$$i\alpha \partial_t \Psi = -\frac{\alpha^2}{2m} \nabla^2 \Psi + V\Psi. \quad (16)$$

We will refer to it as a Schrödinger-like equation: clearly Eq. (16) does not have the same meaning as the usual Schrödinger equation; this would be true only if  $\alpha = \hbar$ , while in general  $\alpha$  is not a universal constant; it is rather a quantity characteristic of the system under consideration (in our case the particle beam). In fact  $\alpha$  turns out to be of the order of magnitude of the beam emittance, a quantity which—in formal analogy with  $\hbar$ —has the dimensions of an action and gives a measure of the position-momentum uncertainty product for the system. Thus the SM model of our beam, as incorporated in the phenomenological Schrödinger equation (16), while keeping a few features reminiscent of the QM, is in fact a deeply different theory.

### B. Controlled distributions

We have introduced the equations that in the SM model are supposed to describe the dynamical behavior of the beam: we now briefly sum up a general procedure, already exploited in previous papers [19,24], to control the dynamics of our systems. Let us suppose that the PDF  $\rho(\mathbf{r}, t)$  is given all along its time evolution: think in particular either of a stationary state, or of an engineered evolution from some initial PDF toward a final state with suitable characteristics. We know that the FP equation (10) must be satisfied, for the given  $\rho$ , by some forward velocity field  $\mathbf{v}_{(+)}(\mathbf{r}, t)$ . Since also Eq. (13) must hold, we are first of all required to find an irrotational  $\mathbf{v}_{(+)}$  which satisfies the FP equation (10) for the given  $\rho$ . We then take into account also the dynamical equation (11): since  $\rho$  and  $\mathbf{v}_{(+)}$  (and hence  $f$  and  $W$ ) are now fixed

and satisfy Eq. (10), Eq. (11) plays the role of a constraint defining a controlling potential  $V$  when we also take into account Eq. (14). We list here the potentials associated with the three particular cases analyzed in the previous papers.

In the 1D case with given dimensionless PDF  $f(x, t)$  and  $a < x < b$  ( $a$  and  $b$  can be infinite) we easily get

$$v_{(+)}(x, t) = D \frac{\partial_x \rho(x, t)}{\rho(x, t)} - \frac{1}{\rho(x, t)} \int_a^x \partial_t \rho(x', t) dx', \quad (17)$$

$$V(x, t) = mD^2 \partial_x^2 \ln f + mD(\partial_t \ln f + v_{(+)} \partial_x \ln f) - \frac{m}{2} v_{(+)}^2 - m \int_a^x \partial_t v_{(+)}(x', t) dx' + \dot{\theta}. \quad (18)$$

For a 3D system with cylindrical symmetry around the  $z$  axis (the beam axis), if we denote by  $(r, \varphi, z)$  the cylindrical coordinates, and if we suppose that  $\rho(r, t)$  depends only on  $r$  and  $t$ , and that  $\mathbf{v}_{(+)} = v_{(+)}(r, t) \hat{\mathbf{r}}$  is radially directed with modulus depending only on  $r$  and  $t$ , we have

$$v_{(+)}(r, t) = D \frac{\partial_r \rho(r, t)}{\rho(r, t)} - \frac{1}{r\rho(r, t)} \int_0^r \partial_t \rho(r', t) r' dr', \quad (19)$$

$$V(r, t) = \frac{mD^2}{r} \partial_r (r \partial_r \ln f) + mD(\partial_t \ln f + v_{(+)} \partial_r \ln f) - \frac{m}{2} v_{(+)}^2 - m \int_0^r \partial_t v_{(+)}(r', t) dr' + \dot{\theta}. \quad (20)$$

Finally, in the 3D stationary case the PDF  $\rho(\mathbf{r})$  is independent of  $t$ . This greatly simplifies our formulas and, by requiring that  $\dot{\theta}(t) = E$  be constant, namely, that  $\theta(t) = Et$ , we get

$$\mathbf{v}_{(+)}(\mathbf{r}) = D \frac{\nabla \rho(\mathbf{r})}{\rho(\mathbf{r})}, \quad (21)$$

$$V(\mathbf{r}) = E + 2mD^2 \frac{\nabla^2 \sqrt{\rho}}{\sqrt{\rho}}. \quad (22)$$

Of course in this context the constant  $E$  will be chosen by fixing the zero of the potential energy. Let us remark finally that in this stationary case the phenomenological wave function (15) takes the form

$$\Psi(\mathbf{r}, t) = \sqrt{\rho} e^{-iEt/\alpha}$$

typical of stationary states.

### C. Nonstationary distributions

In the following we will be mainly concerned with stationary distributions, but in a previous paper we treated also nonstationary problems. For instance, if we consider the stationary, ground state PDF (without nodes)  $\rho_0(\mathbf{r})$  of a suitable potential, and if we calculate  $\mathbf{v}_{(+)}(\mathbf{r})$  and write down the corresponding FP equation, it is possible to show (see the general proof in previous papers [24–26]) that  $\rho_0(\mathbf{r})$  will play



the role of an attractor for every other distribution (nonextremal with respect to a stochastic minimal action principle). If the accelerator beam is ruled by such an equation, this implies that the halo cannot simply be wiped out by scraping away the particles that come out of the bunch core: in fact they simply will keep going out in the halo until equilibrium is reached again since the distribution  $\rho_0(\mathbf{r})$  is a stable attractor.

In a recent paper [21] we gave an estimate of the time required for the relaxation of nonextremal PDF's toward the equilibrium distribution. This is an interesting test for our model since this relaxation time is fixed once the form of the forward velocity field is given; this is in turn fixed when the form of the halo distribution is given as in [20], and one could check if the estimate is in agreement with possible observed times. In particular we estimated that in typical conditions all the nonstationary solutions of this FP equation will be attracted toward  $\rho_0$  with a relaxation time of the order of  $\tau \approx 2m\sigma^2/\alpha \approx 10^{-8} - 10^{-7}$  s.

A different nonstationary problem also discussed in previous papers [19,21] consists in the analysis of some particular time evolution of the process with the aim of finding the dynamics that control it. For instance we studied the possible evolutions which start from a PDF with halo and evolve toward a halo-free PDF: this would allow us to find the dynamics that we need to apply in order to achieve this result. If for simplicity the overall process is supposed to be an Ornstein-Uhlenbeck process, the transition PDF would be completely known and all the results can be exactly calculated through the Chapman-Kolmogorov equation by supposing suitable shapes for the initial and final distributions. Then a direct application of Eq. (18) allows us to calculate the control potential corresponding to this evolution. For the sake of brevity we do not give the analytical form of this potential and refer to the quoted papers for further details.

### III. SELF-CONSISTENT EQUATIONS

#### A. Space-charge interaction

In QM a system of  $N$  particles is described by a wave function in a  $3N$ -dimensional configuration space. On the other hand in our SM scheme a normalized  $|\Psi(\mathbf{r}, t)|^2$ , a function of only three space coordinates  $\mathbf{r}=\{x, y, z\}$ , plays the role of the PDF of a Nelson process. In a first approximation we will consider this  $N$ -particle system as a pure ensemble: as a consequence we will not introduce a  $3N$ -dimensional configuration space, since  $N|\Psi(\mathbf{r}, t)|^2 d^3\mathbf{r}$  in the 3D space will play the role of the number of particles in a small neighborhood of  $\mathbf{r}$ . However, since our system of  $N$  charged particles is not a pure ensemble due to their mutual em interaction, in a further *mean field approximation* we will take into account the so called *space-charge* effects: more precisely we will couple our SL equation with the Maxwell equations describing both the external and the space-charge em fields, and we will get in the end a nonlinear system of coupled differential equations.

In our model a single, charged particle embedded in a beam and experiencing both an external and a space-charge

potential is first of all described by a SL equation

$$i\alpha\partial_t\Psi(\mathbf{r}, t) = \hat{H}\Psi(\mathbf{r}, t),$$

where  $\Psi(\mathbf{r}, t)$  is our wave function,  $\alpha$  a coefficient with the dimensions of an action which is a constant depending on the beam characteristics, and  $\hat{H}$  a suitable Hamiltonian operator. If  $\Psi$  is properly normalized and if  $N$  is the number of particles with individual charge  $q_0$ , the space-charge density and the electrical current density are

$$\rho_{sc}(\mathbf{r}, t) = Nq_0|\Psi(\mathbf{r}, t)|^2, \quad (23)$$

$$\mathbf{j}_{sc}(\mathbf{r}, t) = Nq_0\frac{\alpha}{m} \text{Im}\{\Psi^*(\mathbf{r}, t) \nabla \Psi(\mathbf{r}, t)\}. \quad (24)$$

Hence our particles in the beam will experience both an electrical and a magnetic interaction and we will be obliged to couple the SL equation with the equations of the vector and scalar potentials associated with this electromagnetic field.

The em potentials  $(\mathbf{A}_{sc}, \Phi_{sc})$  of the space charge fields obeying the gauge condition

$$\nabla \cdot \mathbf{A}_{sc}(\mathbf{r}, t) + \frac{1}{c^2}\partial_t\Phi_{sc}(\mathbf{r}, t) = 0 \quad (25)$$

must satisfy the wave equations

$$\nabla^2\mathbf{A}_{sc}(\mathbf{r}, t) - \frac{1}{c^2}\partial_t^2\mathbf{A}_{sc}(\mathbf{r}, t) = -\mu_0\mathbf{j}_{sc}(\mathbf{r}, t), \quad (26)$$

$$\nabla^2\Phi_{sc}(\mathbf{r}, t) - \frac{1}{c^2}\partial_t^2\Phi_{sc}(\mathbf{r}, t) = -\frac{\rho_{sc}(\mathbf{r}, t)}{\epsilon_0}. \quad (27)$$

On the other hand, for our particle in the beam the em field is the superposition of the space charge potential  $(\mathbf{A}_{sc}, \Phi_{sc})$ , and the external potentials  $(\mathbf{A}_e, \Phi_e)$ . Hence (see, for example, [27], Chap. XV) our SL equation takes the form

$$i\alpha\partial_t\Psi = \frac{1}{2m} \left( i\alpha \nabla - \frac{q_0}{c}(\mathbf{A}_{sc} + \mathbf{A}_e) \right)^2 \Psi + q_0(\Phi_{sc} + \Phi_e)\Psi. \quad (28)$$

It is apparent now that Eqs. (25)–(28) constitute a self-consistent system of nonlinear differential equations for the fields  $\Psi$ ,  $\mathbf{A}_{sc}$ , and  $\Phi_{sc}$  coupled through Eqs. (23) and (24).

If we then consider stationary wave functions

$$\Psi(\mathbf{r}, t) = \psi(\mathbf{r})e^{-iEt/\alpha} \quad (29)$$

where  $E$  is the energy of the particle, and take  $\mathbf{A}_e = \mathbf{0}$  for the external interaction, passing to the potential energies

$$V_e(\mathbf{r}) = q_0\Phi_e(\mathbf{r}), \quad V_{sc}(\mathbf{r}) = q_0\Phi_{sc}(\mathbf{r}),$$

our system is reduced to only two coupled, nonlinear equations for the pair  $(\psi, V_{sc})$ , namely,

$$\frac{\alpha^2}{2m}\nabla^2\psi = (V_e + V_{sc} - E)\psi, \quad (30)$$

$$\nabla^2 V_{sc} = -\frac{Nq_0^2}{\epsilon_0}|\psi|^2. \quad (31)$$

### B. Cylindrical symmetry

We suppose now that the longitudinal motion along the  $z$  axis is both decoupled from the transverse motion in the  $x, y$  plane, and free with constant momentum  $p_z$  and velocity  $b_z = b_0 \gg b_x, b_y$ . Moreover, we suppose that the beam particles will be confined in a cylindrical packet of length  $L$ , so that by imposing periodic boundary conditions we will quantize the longitudinal momentum

$$p_z = \frac{2k\pi\alpha}{L}, \quad k = 0, \pm 1, \pm 2, \dots$$

As a consequence our wave functions will take the form

$$\psi(\mathbf{r}) = \chi(x, y) \frac{e^{ip_z z/\alpha}}{\sqrt{L}} \quad (32)$$

and our Eqs. (30) and (31) become

$$\frac{\alpha^2}{2m} (\partial_x^2 + \partial_y^2) \chi = (V_e + V_{sc} - E_T) \chi, \quad (33)$$

$$(\partial_x^2 + \partial_y^2) V_{sc} = -\frac{Nq_0^2}{L\epsilon_0} |\chi|^2 = -\frac{Nq_0^2}{\epsilon_0} |\chi|^2, \quad (34)$$

where  $\mathcal{N} = N/L$  is the number of particles per unit length, and  $E_T = E - p_z^2/2m$  is the energy of the transverse motion. If finally our system has a cylindrical symmetry around the  $z$  axis, namely, if—in the cylindrical coordinate system  $\{r, \varphi, z\}$  ( $r^2 = x^2 + y^2$ )—our potentials depend only on  $r$ , then we can separate the variables with  $\chi(x, y) = u(r)\Phi(\varphi)$ , the angular eigenfunctions are

$$\Phi_\ell(\varphi) = \frac{e^{i\ell\varphi}}{\sqrt{2\pi}}, \quad \ell = 0, \pm 1, \pm 2, \dots, \quad (35)$$

and for  $\ell = 0$  the equations become

$$\frac{\alpha^2}{2m} \left( u'' + \frac{u'}{r} \right) = (V_e + V_{sc} - E_T) u, \quad (36)$$

$$V_{sc}'' + \frac{V_{sc}'}{r} = -\frac{Nq_0^2}{2\pi\epsilon_0} u^2 \quad (37)$$

with the following radial normalization:

$$\int_0^{+\infty} ru^2(r) dr = 1.$$

Note that now we are reduced to a system of *ordinary* differential equations.

### C. Dimensionless formulation

To eliminate the physical dimensions one introduces two quantities  $\eta$  and  $\lambda$  which are, respectively, an energy and a length. Then, by means of the dimensionless quantities

$$s = \frac{r}{\lambda}, \quad \beta = \frac{E_T}{\eta}, \quad \xi = \frac{Nq_0^2}{2\pi\epsilon_0\eta} \quad (\text{perveance}),$$

$$w(s) = \lambda u(\lambda s),$$

$$v(s) = \frac{V_{sc}(\lambda s)}{\eta}, \quad v_e(s) = \frac{V_e(\lambda s)}{\eta},$$

Eqs. (36) and (37) take the form

$$sw''(s) + w'(s) = [v_e(s) + v(s) - \beta]sw(s), \quad (38)$$

$$sv''(s) + v'(s) = -\xi sw^2(s). \quad (39)$$

The usual choice for the dimensional constants is

$$\eta = mb_0^2, \quad \lambda = \frac{\alpha}{mb_0\sqrt{2}}, \quad (40)$$

where  $b_0$  is the longitudinal velocity of the beam. We can now look at our equations in two different ways. First of all we can suppose that  $v_e$  is a given external potential: in this case our aim is to solve the equations for the two unknowns  $w$  (radial particle distribution) and  $v$  (space-charge potential energy). However, in general no simple analytical solution of this problem is at present available for the usual forms of the external potential  $v_e$ : there are not even solutions playing the same role played by the Kapchinskij-Vladimirskij (KV) distribution in the usual models. This phase space distribution—which is simple and self-consistent in the usual dynamical models—leads to a uniform transverse space distribution of the beam, and is a stationary solution of the Vlasov equation with a harmonic potential. Moreover, its space-charge potential calculated from the Poisson equation is still harmonic. Instead in the SM model the uniform distributions are not solutions of the stationary Schrödinger equation, and we know no simple stationary distribution connected to the harmonic potential like the KV. Even the Gaussian distributions—later discussed in this paper—cannot play the same role: they are solutions connected with an external harmonic potential, but their space-charge potential calculated from the Poisson equation is not harmonic.

Alternatively we can assume as known a given distribution  $w$ , and solve our equations to find both the external and the space charge self-consistent potential energies  $v_e$  and  $v$ . In this second form the problem is more simple, and analytical solutions are available. We adopted the first standpoint in a few previous papers [20] where we numerically solved Eqs. (38) and (39); here we will rather elaborate a few additional ideas about the second one. To this end it is important to remark that the space-charge potential energy

$$v(s) = -\xi \int_0^s \frac{dy}{y} \int_0^y xw^2(x) dx \quad (41)$$

is always a solution of the Poisson equation (39) satisfying the conditions  $v(0^+) = v'(0^+) = 0$ . On the other hand, by substituting (41) in the first equation (38) we readily obtain also the self-consistent form of the external potential energy

$$v_e(s) = v_0(s) + \xi \int_0^s \frac{dy}{y} \int_0^y xw^2(x) dx, \quad (42)$$

$$v_0(s) = \frac{w''(s)}{w(s)} + \frac{1}{s} \frac{w'(s)}{w(s)} + \beta \quad (43)$$

where  $v_0(s)$  is the potential that we would have without space charge ( $\xi=0$ ), while the second part in the external potential (42) exactly compensates for the space-charge potential.

#### D. Constant focusing

Let us suppose now that the transverse external potential  $V_e(r)$  is a cylindrically symmetric, harmonic potential with a proper frequency  $\omega$  (constant focusing), and let us also introduce the characteristic length

$$\sigma^2 = \frac{\alpha}{2m\omega}$$

which will represent a measure of the transverse dispersion of the beam. In cylindrical coordinates  $\{r, \varphi\}$  in the transverse plane our potential energy is

$$V_e(r) = \frac{m\omega^2}{2} r^2 = \frac{\alpha^2}{8m\sigma^4} r^2 \quad (44)$$

so that the corresponding 2D SL equation *without* space charge (zero perveance) would have as lowest eigenvalue  $E_0 = \alpha\omega$ , and as ground state wave function

$$\chi_{00}(r, \varphi) = \frac{u_0(r)}{\sqrt{2\pi}} = \frac{e^{-r^2/4\sigma^2}}{\sigma\sqrt{2\pi}}. \quad (45)$$

Of course the self-consistent solution will be different if there is a space charge (nonzero perveance). To find this solution one introduces the so-called *phase advance*

$$\frac{1}{\lambda_0} = \frac{\omega}{b_0} = \frac{\alpha}{2mb_0\sigma^2}$$

( $\lambda_0$  is a length) and, with the constants (40), the dimensionless form of the harmonic potential (44)

$$v_e(s) = \frac{V_e(r)}{mb_0^2} = \frac{\omega^2}{b_0^2} r^2 = \frac{r^2}{2\lambda_0^2} = \frac{\alpha^2}{4\lambda_0^2 m^2 b_0^2} s^2 = \gamma^2 s^2,$$

$$\gamma = \frac{\alpha}{2\lambda_0 m b_0} = \frac{\alpha\omega}{2mb_0^2} = \frac{\sigma^2}{\lambda_0^2}.$$

As a consequence Eqs. (38) and (39) become

$$sw''(s) + w'(s) = [\gamma^2 s^2 + v(s) - \beta]sw(s), \quad (46)$$

$$sv''(s) + v'(s) = -\xi sw^2(s). \quad (47)$$

These equations are now a coupled, nonlinear system which must be *numerically* solved since we do not know simple self-consistent solutions of the form of the KV distribution. In Ref. [20] we extensively analyzed these numerical solutions and we refer to this paper for details. In fact in [20] there was a small difference with respect to what has been presented here. The form of the equations to solve is the same, but the dimensionless formulation was achieved by

means of two numerical constants different from (40) and drawn from the characteristics of the transverse harmonic oscillator force:

$$\eta = \frac{\alpha^2}{4m\sigma^2} = \frac{\alpha\omega}{2}, \quad \lambda = \sigma\sqrt{2}. \quad (48)$$

Then the dimensionless quantities have a different numerical value and the dimensionless equations (36) and (37) take the form

$$sw''(s) + w'(s) = [s^2 + v(s) - \beta]sw(s), \quad (49)$$

$$sv''(s) + v'(s) = -\xi sw^2(s), \quad (50)$$

since now  $\gamma=1$ . In any case Eqs. (46) and (47) can easily be turned into Eqs. (49) and (50), and vice versa, by means of simple transformations through the parameter  $\gamma$ , which turns out to be at the same time the ratio of the energy constants, and that of the squared length constants. As a consequence, in the following we will always use the system (49) and (50), with the advantage of simply putting  $\gamma=1$  in the model.

## IV. SELF-CONSISTENT POTENTIALS

### A. Gaussian transverse distributions

In the SM model it is possible to numerically integrate the Schrödinger-Poisson system (38) and (39) with a given external potential and calculate the self-consistent distributions and their space-charge potentials [20]. On the other hand, if we fix a particular distribution, it is always possible to exactly calculate from these equations the external and space-charge potentials giving rise to that distribution. When we adopt this second alternative approach and we take as given the form of the distribution  $w(s)$ , the unknowns in Eqs. (38) and (39) are the two potential energies  $v(s)$  and  $v_e(s)$ . In this case we only need to calculate the expressions (41) and (42) in terms of the given distribution  $w(s)$ . Of course if we take an arbitrary  $w(s)$  we will not get any simple and meaningful form for the external potential  $v_e(s)$ ; and on the other hand to guess the right form of  $w(s)$  giving rise, for instance, exactly to a harmonic potential (44) as external potential would be tantamount to solving (42) as an integro-differential equation for a given external potential. However, in a few explicit cases the results are quite simple and interesting.

Let us take as our first example of a stationary wave function that of the ground state  $u_0(r)$  of the harmonic oscillator with zero perveance given in Eq. (45). Its dimensionless representation is

$$w(s) = \sqrt{2}e^{-s^2/2}, \quad \beta = 2 \quad (E_T = \alpha\omega), \quad (51)$$

which is also apparently normalized. We now want to calculate both the external and the space-charge potentials that produce (51) as the stationary wave function for (49) and (50). From Eqs. (41), (42), and (51) we then have

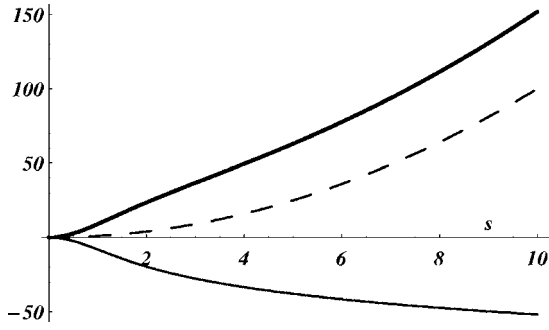


FIG. 1. The dimensionless potentials  $v(s)$  (thin line),  $v_0(s)=s^2$  (dashed line), and  $v_e(s)$  (thick line). They reproduce, respectively, Eqs. (52)–(54) for  $\xi=20$  (see Ref. [20] for this value). When the external potential is  $v_e(s)$  the self-consistent wave function coincides with that of a simple harmonic oscillator for zero perveance (51).

$$\frac{w''(s)}{w(s)} + \frac{1}{s} \frac{w'(s)}{w(s)} + \beta = v_0(s) = s^2,$$

$$\int_0^s \frac{dy}{y} \int_0^y w^2(x) x dx = \frac{1}{2} [\ln(s^2) + C - \text{Ei}(-s^2)]$$

where  $C \approx 0.577$  is the Euler constant and

$$\text{Ei}(x) = \int_{-\infty}^x \frac{e^t}{t} dt, \quad x < 0,$$

is the exponential-integral function, and hence we immediately get (see also Fig. 1)

$$v(s) = -\frac{\xi}{2} [\ln(s^2) + C - \text{Ei}(-s^2)], \quad (52)$$

$$v_0(s) = s^2, \quad (53)$$

$$v_e(s) = s^2 + \frac{\xi}{2} [\ln(s^2) + C - \text{Ei}(-s^2)]. \quad (54)$$

In a sense the meaning of Eqs. (41)–(43) is rather simple: if we want to get a self-consistent distribution which coincides with a solution of the SL equation for a given-zero perveance potential, the simplest way it is to calculate the space-charge potential for this *frozen* distribution through the Poisson equation, and then compensate the external potential exactly for that. This is what we did in our example where the Gaussian solution is the fundamental state of a harmonic oscillator: we finally got a total potential which is  $v_0(s) = v(s) + v_e(s) = s^2$  (namely, that of a simple harmonic oscillator), and an energy value which coincides with the first eigenvalue. In other words, if you want a Gaussian transverse distribution you should not simply turn on a bare harmonic potential  $s^2$ ; you should rather teleologically compensate for the space charge by using the potential  $v_e(s)$ .

### B. Student transverse distributions

If the halo consists in the fact that large deviations from the beam axis are possible, we can suppose that the sta-

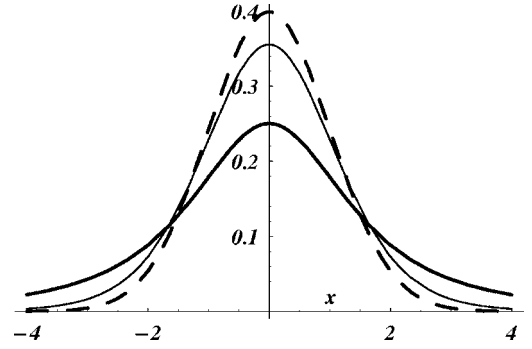


FIG. 2. The Gauss PDF  $\mathcal{N}(0,1)$  (dashed line) compared with the  $\Sigma(2,2)$  (thick line) and the  $\Sigma(10,12)$  (thin line). The inflection points of the three curves coincide. Clearly the tails of the Student laws are much longer.

tionary transverse distribution is different from the Gaussian distribution (51) introduced in Sec. IV A. To this end we will introduce in the following a family of distributions that decay with the distance from the axis only with a power law.

Let us consider the following family of univariate, two-parameter probability laws  $\Sigma(\nu, a^2)$  characterized by the following PDF's:

$$f(x) = \frac{\Gamma\left(\frac{\nu+1}{2}\right)}{\Gamma\left(\frac{1}{2}\right)\Gamma\left(\frac{\nu}{2}\right)} \frac{a^\nu}{(x^2 + a^2)^{(\nu+1)/2}}, \quad (55)$$

which apparently are symmetric functions with a mode at  $x=0$  and two inflection points at  $x = \pm a/\sqrt{\nu+2}$ . All these laws are centered at the median. In particular  $a$  plays just the role of a scale parameter, while  $\nu$  rules the power decay of the tails: for large  $x$  the tails vary as  $x^{-(\nu+1)}$  with  $\nu+1 > 1$ . For a comparison with the Gauss law  $\mathcal{N}(0, \sigma^2)$  see Fig. 2. Note that when  $\nu$  grows larger and larger, the difference between the two PDF's becomes smaller and smaller. It is typical of the laws  $\Sigma(\nu, a^2)$  that they have (finite) momenta of order  $k$  only if the condition  $k < \nu$  is satisfied; hence for  $\nu \leq 2$  there is no variance, while for  $\nu \leq 1$  not even the expectation is defined. On the other hand when  $\nu > 2$  the variance of  $\Sigma(\nu, a^2)$  exists and is

$$\sigma^2 = \frac{a^2}{\nu-2}. \quad (56)$$

It will be useful to remark that the laws  $\Sigma(1, a^2)$  are the well-known Cauchy laws  $\mathcal{C}(a)$  with PDF

$$f(x) = \frac{1}{\pi} \frac{a}{x^2 + a^2},$$

while the laws  $\Sigma(n, n)$  with  $n=1, 2, \dots$  are the classical Student  $t$  laws  $\mathcal{S}(n)$  with PDF



$$f(x) = \frac{\Gamma\left(\frac{n+1}{2}\right)}{\sqrt{\pi}\Gamma\left(\frac{n}{2}\right)}(n+x^2)^{-(n+1)/2}.$$

We will then refer to  $\Sigma(\nu, a^2)$  as generalized Student laws since they are just Student laws with a continuous parameter  $\nu > 0$  and a scale parameter  $a$ . For  $\nu > 2$  variances exist and we are then entitled to standardize our laws: indeed from Eq. (56) every  $\Sigma(\nu, (\nu-2)\sigma^2)$  with  $a^2 = (\nu-2)\sigma^2$  has variance  $\sigma^2$ , and the standard (with unit variance) generalized Student laws are  $\Sigma(\nu, \nu-2)$ .

In order to describe the beam we will also introduce the bivariate, circularly symmetric Student laws  $\Sigma_2(\nu, a^2)$  with PDF

$$f(x, y) = \frac{\nu}{2\pi} \frac{a^\nu}{(x^2 + y^2 + a^2)^{(\nu+2)/2}}. \tag{57}$$

Its marginal laws are both  $\Sigma(\nu, a^2)$  and noncorrelated, albeit not independent (as in the case of the circularly symmetric Gaussian bivariate laws). The total beam distribution will then be

$$\rho(x, y, z) = \frac{1}{2\pi L} \frac{\nu a^\nu}{(x^2 + y^2 + a^2)^{(\nu+2)/2}} H\left(\frac{L}{2} - |z|\right) \tag{58}$$

where  $H(z)$  is the Heaviside function. In the description of a beam in an accelerator it is realistic to suppose that the transverse distribution is endowed with a finite variance. Hence we will look for distributions (58) with  $\nu > 2$ . On the other hand this will correspond to suppose that in our model the transverse Student laws should not be radically different from a Gaussian: in fact the halo is in some sense an effect that is small when compared with the total beam. From this standpoint the family of laws  $\Sigma(\nu, a^2)$  has also the advantage that we can fine-tune the parameters  $\nu, a$  in order to get the right distance from the gaussian laws (this would not be possible if we adopted stable laws; see Sec. V A). With this hypothesis in mind we will limit our present considerations to the case  $\nu > 2$  so that the transverse marginals of (58) will have a finite variance  $\sigma^2$ . Then from (56) we choose  $a^2 = (\nu - 2)\sigma^2$  and write (58) as

$$\rho(x, y, z) = \frac{\nu}{2\pi L} \frac{[(\nu-2)\sigma^2]^{\nu/2}}{[x^2 + y^2 + (\nu-2)\sigma^2]^{(\nu+2)/2}} H\left(\frac{L}{2} - |z|\right). \tag{59}$$

Passing to cylindrical random variables we then have

$$\rho(r, \varphi, z) = r \frac{\nu}{2\pi L} \frac{[(\nu-2)\sigma^2]^{\nu/2}}{[r^2 + (\nu-2)\sigma^2]^{(\nu+2)/2}} H\left(\frac{L}{2} - |z|\right),$$

namely,

$$\rho(r, \varphi, z) = \frac{1}{\sigma\sqrt{2}} \frac{r}{\sigma\sqrt{2}} \frac{2\nu}{\nu-2} \left(1 + \frac{r^2}{(\nu-2)\sigma^2}\right)^{-(\nu+2)/2} \times \frac{H(L/2 - |z|)}{2\pi L},$$

so that finally with the shorthand notation

$$z = \frac{s\sqrt{2}}{\sqrt{\nu-2}} \tag{60}$$

the dimensionless, normalized radial distribution is

$$w^2(s) = \frac{2\nu}{\nu-2} \frac{1}{(1+z^2)^{(\nu+2)/2}}. \tag{61}$$

Here we adopt the dimensional constants

$$\eta = \frac{\alpha^2}{4m\sigma^2}, \quad \lambda = \sigma\sqrt{2}, \tag{62}$$

where  $\sigma^2$  is the variance of our Student laws. We can now use the relations (41)–(43) in order to get the potentials that have (59) as stationary distribution: first of all the space-charge potential produced by (59) has the form

$$v(s) = -\frac{\xi}{2} \left[ \frac{2z^{-\nu}}{\nu} {}_2F_1\left(\frac{\nu}{2}, \frac{\nu}{2}; \frac{\nu+2}{2}; -\frac{1}{z^2}\right) + \ln z^2 + C + \psi\left(\frac{\nu}{2}\right) \right] \tag{63}$$

where  ${}_2F_1(a, b; c; w)$  is a hypergeometric function and  $\psi(w) = \Gamma'(w)/\Gamma(w)$  is the logarithmic derivative of the Euler Gamma function (digamma function). On the other hand, by choosing  $\beta = 2 + 8/(\nu-2)$  to put the potential energies to zero at the origin, we get the control potential for zero perveance,

$$v_0(s) = \frac{\nu+2}{\nu-2} \frac{z^2(4z^2 + \nu + 10)}{2(1+z^2)^2}, \tag{64}$$

and hence the external potential required to keep a transverse Student distribution  $\Sigma_2(\nu, (\nu-2)\sigma^2)$  with a given variance  $\sigma^2$  is

$$v_e(s) = \frac{\nu+2}{\nu-2} \frac{z^2(4z^2 + \nu + 10)}{2(1+z^2)^2} + \frac{\xi}{2} \left[ \frac{2z^{-\nu}}{\nu} {}_2F_1\left(\frac{\nu}{2}, \frac{\nu}{2}; \frac{\nu+2}{2}; -\frac{1}{z^2}\right) + \ln z^2 + C + \psi\left(\frac{\nu}{2}\right) \right]. \tag{65}$$

Formulas (63)–(65) give the self-consistent potentials associated with the beam distribution (59) which is transversally a Student  $\Sigma_2(\nu, (\nu-2)\sigma^2)$ . In Fig. 3 we can see an example of the control potential  $v_0(s)$  for a particular value of the parameter  $\nu$ , together with its limit behaviors,

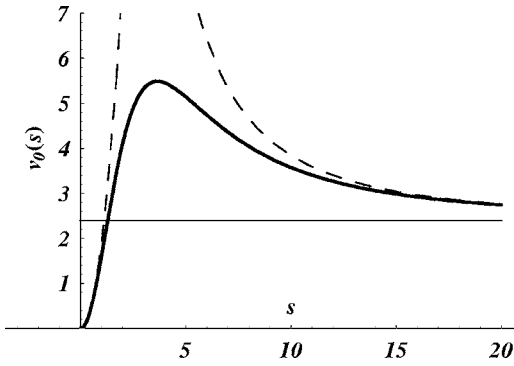


FIG. 3. The control potential  $v_0(s)$  (64) for a Student transverse distribution  $\Sigma_2(22, 20\sigma^2)$ . Also displayed are the value of  $\beta=2.4$  (the limit value of  $v_0$  for large  $s$ , thin line) and the behaviors for small and large  $s$  [Eqs. (66) and (67)] (dashed lines).

$$v_0(s) \sim \begin{cases} \frac{(\nu+2)(\nu+10)}{(\nu-2)^2} s^2 & (s \rightarrow 0^+), \\ \frac{(\nu+2)^2}{4s^2} + 2 + \frac{8}{\nu-2} & (s \rightarrow +\infty). \end{cases} \quad (66)$$

$$v_0(s) \sim \begin{cases} \frac{(\nu+2)(\nu+10)}{(\nu-2)^2} s^2 & (s \rightarrow 0^+), \\ \frac{(\nu+2)^2}{4s^2} + 2 + \frac{8}{\nu-2} & (s \rightarrow +\infty). \end{cases} \quad (67)$$

Now these results must be compared with the similar results (52)–(54) associated with a transversally Gaussian distribution. We will choose the Gaussian parameters in such a way that the behavior near the beam axis is similar to Eq. (66), namely (with  $\beta=2\gamma$ ),

$$w(s) = \sqrt{2\gamma} e^{-2\gamma s^2/2}, \quad \gamma^2 = \frac{(\nu+2)(\nu+10)}{(\nu-2)^2}.$$

First of all in Fig. 4 we compare the space-charge potential produced by both a Student and a Gauss transverse distribution: note that for the chosen parameter values ( $\nu=22$ ,  $\xi=20$ ) the two potentials look particularly similar. In fact, given the asymptotic behavior of the hypergeometric function in Eq. (63) and of the exponential integral in Eq. (52), for  $s \rightarrow +\infty$  both potentials behave as  $-\xi \ln s$ . On the other hand we immediately see from Fig. 5 that the control potentials for zero perveance  $v_0(s)$  behave differently when we move away from the beam axis; beyond a distance of about  $r \approx 2\sigma$  the two curves are different: while in the Gaussian case the potential diverges as  $s^2$ , in the Student case it goes

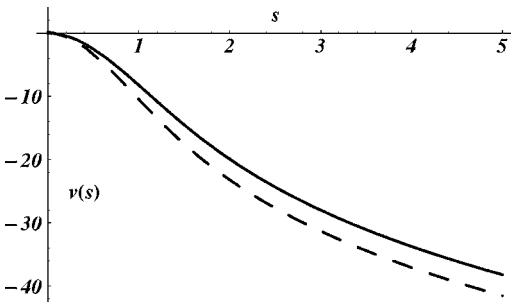


FIG. 4. The space-charge potentials  $v(s)$  (63) and (52), respectively, for a Student transverse distribution  $\Sigma_2(22, 20\sigma^2)$  (solid line), and for a Gauss distribution (dashed line). The dimensionless perveance here is  $\xi=20$ .

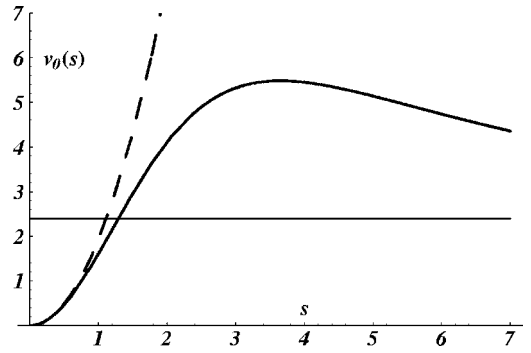


FIG. 5. The control potential  $v_0(s)$  [Eq. (64)] of a Student  $\Sigma_2(22, 20\sigma^2)$  distribution (solid line; see Fig. 3) is here compared with that of a Gauss distribution (dashed line) which shows the same behavior near the beam axis.

to the constant value  $\beta$  as quickly as  $s^{-2}$ . Of course this difference fades away when  $\nu$  grows larger and larger; that points to the fact that the principal difference between the two cases can be confined in a region that can be made as far removed from the beam core as we want by a suitable choice of  $\nu$ . Finally in Fig. 6 we compare the total external potentials needed to keep the transverse beam, respectively, in a Student and in a Gauss distribution. We then see that for large  $s$  (far away from the beam core), while in the Gauss case the total external potential grows with  $s$  as  $s^2 + \xi \ln s$ , in the Student case this potential only grows as  $\xi \ln s$ . In any case, even if the potential near the beam axis is harmonic, deviations from this behavior in a region removed from the core can produce a deformation of the distribution from the Gaussian to the Student.

### C. Estimating the emittance

If  $u(r)$  is a self-consistent, cylindrically symmetric solution of Eqs. (36) and (37) the position probability density  $\rho(r, \varphi, z)$  in cylindrical coordinates will have the form

$$\rho(r, \varphi, z) = \frac{1}{2\pi L} \begin{cases} u^2(r), & 0 \leq \varphi < 2\pi, \quad -\frac{L}{2} \leq z \leq \frac{L}{2}, \\ 0, & \text{otherwise.} \end{cases}$$

In order to estimate the emittance we need to calculate mean values of positions  $x$  and momenta  $p_x$  along one transverse

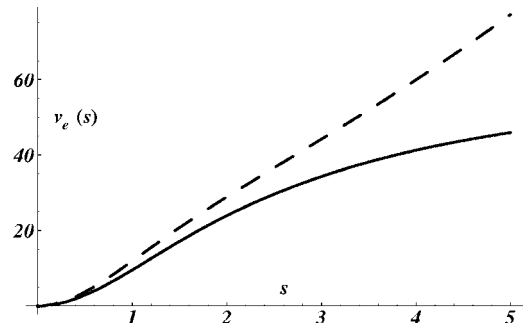


FIG. 6. The total external potential  $v_e(s)$  [Eq. (65)] that should be applied to get a stationary Student transverse distribution  $\Sigma_2(22, 20\sigma^2)$  (solid line), compared with that [Eq. (54)] needed for a Gauss distribution (dashed line).

direction, but we should remember that in SM we have neither a distribution in the phase space nor an operator formalism. The momentum and its distribution should then be recovered from the velocity fields (4) and (5) where—since we are dealing with stationary states with  $\mathbf{v}=\mathbf{0}$ —only the osmotic part is nonzero, so that

$$\mathbf{p} = \alpha \frac{\nabla \rho}{\rho} = 2\alpha \frac{\nabla u(r)}{u(r)}.$$

Supposing now that we choose the  $\nu$  of our Student laws so that the following integrals exist, we then have

$$\langle x \rangle = \int_0^L \frac{dz}{L} \int_0^{2\pi} \frac{\cos \varphi}{2\pi} d\varphi \int_0^{+\infty} r^2 u^2(r) dr = 0,$$

$$\langle x^2 \rangle = \int_0^L \frac{dz}{L} \int_0^{2\pi} \frac{\cos^2 \varphi}{2\pi} d\varphi \int_0^{+\infty} r^3 u^2(r) dr = \frac{1}{2} \int_0^{+\infty} r^3 u^2(r) dr,$$

$$\langle p_x \rangle = 2\alpha \int_0^L \frac{dz}{L} \int_0^{2\pi} \frac{\cos \varphi}{2\pi} d\varphi \int_0^{+\infty} r u(r) u'(r) dr = 0,$$

$$\begin{aligned} \langle p_x^2 \rangle &= 4\alpha^2 \int_0^L \frac{dz}{L} \int_0^{2\pi} \frac{\cos^2 \varphi}{2\pi} d\varphi \int_0^{+\infty} r^2 u'(r) dr \\ &= 2\alpha^2 \int_0^{+\infty} r u'^2(r) dr, \end{aligned}$$

so that the standard deviations (uncertainties) are

$$\Delta x = \sqrt{\langle x^2 \rangle} = \sqrt{\frac{1}{2} \int_0^{+\infty} r^3 u^2(r) dr}, \quad (68)$$

$$\Delta p_x = \sqrt{\langle p_x^2 \rangle} = \sqrt{2\alpha^2 \int_0^{+\infty} r u'^2(r) dr}, \quad (69)$$

and the position-momentum covariance is

$$\begin{aligned} C &= \langle x p_x \rangle - \langle x \rangle \langle p_x \rangle = \langle x p_x \rangle \\ &= 2\alpha \int_0^L \frac{dz}{L} \int_0^{2\pi} \frac{\cos^2 \varphi}{2\pi} d\varphi \int_0^{+\infty} r^2 u'(r) u(r) dr \\ &= \alpha \int_0^{+\infty} r^2 u'(r) u(r) dr. \end{aligned}$$

In a previous paper [20] we adopted the uncertainty product  $\Delta x \Delta p_x$  as a measure of the rms emittance. As an example let us suppose again that our wave function has the form  $u_0(r)$  for the harmonic oscillator without space charge given in Eq. (45). We then have

$$\Delta x \Delta p_x = -C = \alpha. \quad (70)$$

This allows two remarks: first,  $\alpha$  plays also the role of a measure of the emittance and hence—as suggested in a previous paper [19]—its value must be linked to the number of particles in the beam; second, the position-momentum correlation coefficient of a Gaussian beam is

$$\frac{C}{\Delta x \Delta p_x} = -1,$$

as was predictable, since in SM the relation between position and momentum for the wave function (45) is linear and negative.

In other models the transverse rms emittance is calculated by means of the quantity  $\sqrt{\Delta x^2 \Delta p_x^2 - C^2}$ . In the KV distribution, since momentum and position are uncorrelated and  $\langle x \rangle = \langle p_x \rangle = 0$ , this estimate becomes  $\sqrt{\langle x^2 \rangle \langle p_x^2 \rangle - \langle x p_x \rangle^2}$ . In the SM model, on the contrary, this is not a good choice: in fact we have shown, at least in our simple example, that  $x$  and  $p_x$  are far from uncorrelated, and that as a consequence of Eq. (70)  $\sqrt{\Delta x^2 \Delta p_x^2 - C^2}$  becomes exactly zero. Apparently it is not realistic to take this value as a good estimate of the emittance. On the other hand, for the same Gaussian example, the value of the uncertainty product  $\Delta x \Delta p_x$  is just  $\alpha$ , which we assume to be a good candidate for the value of the emittance. On the other hand it is easy to calculate the same uncertainty product for a Student distribution  $\Sigma_2(\nu, (\nu-2)\sigma^2)$  with dimensionless radial distribution (61) and variance  $\sigma^2$ : in fact a straight application of Eqs. (68) and (69) brings the following result:

$$\Delta x \Delta p_x = \alpha \sqrt{\frac{\nu(\nu+2)}{(\nu-2)(\nu+4)}}. \quad (71)$$

Of course, as is already clear, this value converges to the Gaussian case for large  $\nu$ , while it becomes larger and larger for small  $\nu$  values when the shape of the distribution moves away from the Gaussian case.

#### D. Weighing the tails

We can finally compare the length of the tails of Gauss and Student distributions in order to assess the possible halo formation in the second case. Let us consider the probability

$$P(c) = \int_{c\sigma}^{+\infty} r u^2(r) dr \quad (72)$$

of being beyond a distance  $c\sigma$  ( $\sigma^2$  being the variance) away from the beam axis, and calculate this quantity in our two cases. From the Gaussian distribution we have from Eq. (45) that

$$P(c) = e^{-c^2/2}, \quad (73)$$

while in the Student case from Eq. (59) we get

$$P(c) = \left(1 + \frac{c^2}{\nu-2}\right)^{-\nu/2}. \quad (74)$$

Now for  $c=10$  the Gaussian value is about  $1.9 \times 10^{-22}$ , while with  $\nu=10$  the Student value is about  $2.2 \times 10^{-6}$ , and with  $\nu=22$  the value is  $2.8 \times 10^{-9}$ . This means that for  $\mathcal{N}=10^{11}$  particles per meter of beam, we find practically no particle beyond  $10\sigma$  in the Gaussian case, but about  $10^3$  particle per meter for a  $\nu=22$  Student distribution, and as many as  $10^5$  for a  $\nu=10$  value. It is worthwhile to remember at this point that we got about the same number of particles gone astray in

our self-consistent numerical solutions for a dimensionless pervance of about  $\xi=20$  in one of our previous papers [20].

## V. LÉVY-STUDENT PROCESSES

In our context the Student laws  $\Sigma(\nu, a^2)$  are important not only because they promise to better describe the halo by means of their longer tails with respect to the usual Gaussian distributions; in fact they constitute an important family of Lévy *infinitely divisible* (ID) laws. At present there is a lot of interest about non-Gaussian Lévy laws in several fields of research (see for example [8,16] and references quoted therein), but this interest is mostly confined to the *stable* laws which are in fact an important subfamily of the ID laws. The fundamental character of the ID laws can be better understood from two different, but strictly correlated standpoints: on the one hand the ID laws constitute the more general form of possible limit laws for the generalized central limit theorem; on the other they constitute the class of all the laws of the increments for every stationary, stochastically continuous, independent-increment process (Lévy process). These important results (which are briefly discussed in Appendixes A and B) have been achieved by Lévy, Khintchin, Kolmogorov, and other mathematicians from the mid-1930s to the mid-1940s but their relevance for applications has been recognized only in more recent years. One of the characteristics of a non-Gaussian Lévy process is to have trajectories with moving discontinuities (think of the trajectories of a typical Poisson process contrasted with those of a Gaussian Wiener process), and we propose here to describe the trajectories of the particle beam by means of a Lévy-Student process whose discontinuities can possibly account for the relatively rare escape of particles from the beam core. For the sake of simplicity we will limit ourselves in the following to the case of ID systems representing one single transverse coordinate of our particle beam.

### A. The Student ID laws

The characteristic functions of the laws  $\Sigma(\nu, a^2)$ , namely, the Fourier transforms of the densities (55), are

$$\varphi(\kappa) = 2 \frac{|a\kappa|^{\nu/2} K_{\nu/2}(|a\kappa|)}{2^{\nu/2} \Gamma(\nu/2)} \quad (75)$$

where  $K_\alpha(z)$  is a modified Bessel function. The typical form of these characteristic functions (contrasted with the Gauss one) is shown in Fig. 7. Remark that, since  $x$  is a length, the characteristic function variable  $\kappa$  has the dimensions of wave number (inverse of a length). These laws are ID but in general are not stable, the unique stable laws among them being the Cauchy laws  $\mathcal{C}(a) = \Sigma(1, a^2)$ . For  $\nu > 2$  the Student distributions belong to the domain of attraction of the Gauss law since they have a finite variance. For  $\nu \leq 2$  the variance diverges, but it is possible to prove that, this notwithstanding, they still belong to the Gaussian domain also for  $\nu=2$ : albeit this derives from a well-known general result [28] a simple proof for our particular case will be given in a subsequent paper. On the other hand for  $\nu < 2$  the Student laws are at-

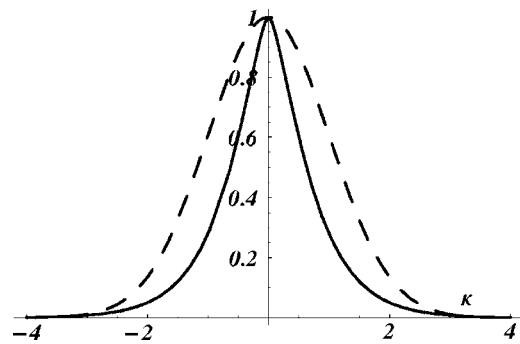


FIG. 7. Typical characteristic function of a Student law  $\Sigma(2,2)$  (solid line) compared with that of a standard Gauss law  $\mathcal{N}(0,1)$  (dashed line).

tracted by non-Gaussian stable laws characterized by the same value of the parameter  $\nu$ .

The fact that the Student laws are ID—which in itself is not at all a trivial result proved in steps only in the 1970s and 1980s [29]—shows two kinds of advantages with respect to more common stable laws.

(1) No stable, non-Gaussian law can have a finite variance, while all Student laws with  $\nu > 2$  do have a finite variance; this is important since it is not realistic to suppose that empirical distributions (in particular for particle beams) have infinite variances, but that notwithstanding we will not be obliged to resort to handmade modification (for instance truncated Lévy distributions) as in the case of stable distributions [8].

(2) The asymptotic behavior of stable, non-Gaussian laws is proportional to  $|x|^{-\alpha-1}$  with  $\alpha < 2$  [16], while the asymptotic behavior of the Student laws is  $|x|^{-\nu-1}$  with  $\nu > 0$ ; this allows the Student laws—but not the stable laws—to continuously go through all the gamut of decay speeds to approximate in a fine tuning the Gaussian behavior as well as we want.

The principal drawback of not being stable is in the subsequent definition of the Lévy-Student process. In fact the characteristic function of the process  $\varphi(\kappa, t)$  coincides with Eq. (75) only for  $t-s=T$ , while for  $t-s \neq T$  it is no longer the characteristic function of a  $\Sigma(\nu, a^2)$  law. Hence we explicitly know the form of the increment law only at the time scale  $T$ : we know the characteristic function—namely, everything we theoretically need—at every time, but we do not have the explicit inverse Fourier transform, and we also know that the laws are no longer in the family  $\Sigma(\nu, a^2)$ . This problem is tempered by the remark that the situation is not better for general stable laws: even in this case, in fact, we do not know the explicit forms of the increment laws, not even for one time scale (they are known only in a few precious instances). The unique advantage in the stable case being the fact that all along the time evolution the increment laws remain of the same type, which is not the case for ID nonstable laws.

### B. The Lévy-Student process

A Lévy process defined by the characteristic function (75) will be called in the following a Lévy-Student process. Tak-



ing into account (B3) and (75) we can now state that the transition PDF of a Lévy-Student process is

$$p(x, t|y, s) = \frac{1}{2\pi} \int_{-\infty}^{+\infty} e^{i\kappa(x-y)} \left( 2 \frac{|a\kappa|^{\nu/2} K_{\nu/2}(|a\kappa|)}{2^{\nu/2} \Gamma(\nu/2)} \right)^{(t-s)/T} d\kappa \tag{76}$$

where the improper integral is always convergent since the asymptotic behavior of the characteristic function is

$$\varphi(\kappa) = \sqrt{2\pi} \frac{|a\kappa|^{(\nu-1)/2} e^{-|a\kappa|[1 + O(|\kappa|^{-1})]}}{2^{\nu/2} \Gamma(\nu/2)}, \quad |\kappa| \rightarrow +\infty.$$

In principle Eq. (76) should be enough to calculate everything of our process, but in practice this is an integral that must be treated numerically, but for a few particular cases that will be discussed in a subsequent paper. On the other hand even to produce simulation for the trajectories of our process we should have some simple expression for the transition PDF. At least for this last task, however, we can exploit the fact that when  $t-s=T$  the expression (76) can be exactly calculated and coincides with the PDF (55) of a Student  $\Sigma(\nu, a^2)$  (note that even this is not possible for the typical non-Gaussian stable process). This means that we can produce sample trajectories by taking  $T$  as the fundamental step of our numerical simulation. In other words we will simulate the sample paths of a process whose increments are exactly Student distributed when observed at the (otherwise arbitrary) time scale  $T$ . To give a look to these trajectories we produced a simplified model which simulates the solutions of the following two SDE's:

$$dX(t) = v(X(t), t)dt + dW(t), \tag{77}$$

$$dY(t) = v(Y(t), t)dt + dS(t), \tag{78}$$

where  $W(t)$  is a Wiener process, while  $S(t)$  is a Lévy-Student process. We also fixed the velocity field  $v(x, t)$  in a suitable way: it will not depend on time  $t$ , and its value is (for given  $b > 0$  and  $q > 0$ )

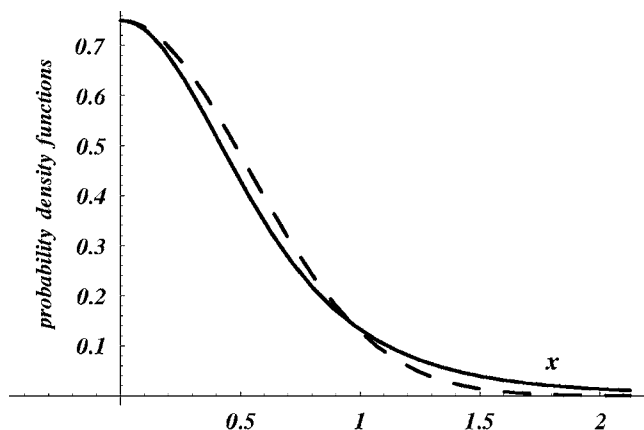


FIG. 8. The PDF's of the increments for the Gaussian processes [SDE (77), dashed line;  $\sigma \approx 0.53$ ] and for the Lévy-Student process with law  $\Sigma(4, 1)$  [SDE (78), solid line;  $\sigma \approx 0.71$ ]. The parameters are chosen so that the two PDF's have the same modal values and similar shapes.



FIG. 9. Typical trajectory of a stationary, Gaussian (Ornstein-Uhlenbeck) process  $X(t)$  [see SDE (77)]. To compare it with the Student trajectory, the vertical scale has been set equal to that of Fig. 10.

$$v(x) = -bxH(q - |x|)$$

where  $H$  is the Heaviside function. This flux will attract the trajectory toward the origin when  $|x| \leq q$ , and will allow the movement to be completely free for  $|x| > q$ . The forms of the typical PDF's used in our simulations are shown in Fig. 8. In a simplified model for a collimated beam this will then produce a stationary, Ornstein-Uhlenbeck process for the SDE (77) if the intensity of the Gaussian noise is not too large. The process solution of the SDE (78) will instead have different characteristics. Let us suppose to fix our ideas that the two parameters defining the velocity field are  $b=0.35$  and  $q=10$ . Figure 9 displays a typical trajectory of a  $10^4$ -step solution  $X(t)$  of (77) when the variance of the Gaussian distributed increments is  $\sigma^2=0.28$ . In our simplified 1D model of the transverse dynamics of a particle beam this means that the trajectories always stay inside the beam core. Let us then take as the law for the increments of Eq. (78) a Student distribution  $\Sigma(4, 1)$ : its PDF looks not very different from that of the previous Gaussian distribution, as Fig. 8 clearly shows, notwithstanding that the process  $Y(t)$  differs in several respects from  $X(t)$ . Indeed, not only does the typical trajectory displayed in Fig. 10 show a wider dispersion of its values and a few larger spikes. The principal difference is rather in the fact that while the trajectories of  $X(t)$  show a remarkable stability in their statistical behavior, the paths of  $Y(t)$  have the propensity to make occasional excursions far

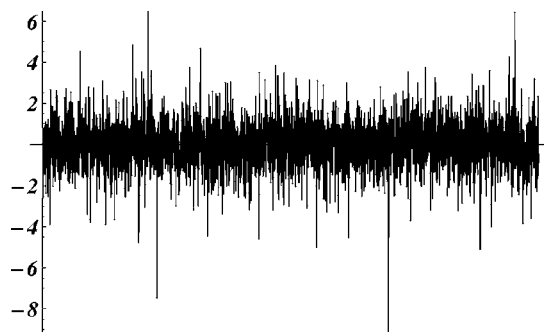


FIG. 10. Typical trajectory of a stationary, Student process  $Y(t)$  [see SDE (78);  $\nu=4$  and  $a=1$ ].

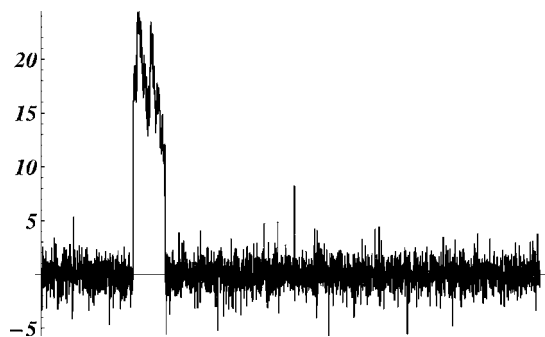


FIG. 11. Occasional trajectory of a stationary, Student process  $Y(t)$  [see SDE (78);  $\nu=4$ ] with a temporary excursion out of the core.

away from the beam core (see Fig. 11), and sometimes they also definitely drift away from the core (see Fig. 12). This depends of course on the mentioned properties of the trajectories of a non-Gaussian Lévy process, and in particular on the fact that they are only stochastically and not pathwise continuous, namely, they contain occasional jumps. The frequency and the size of these jumps can also be fine tuned by suitably choosing the values of the parameters of the law  $\Sigma(\nu, a^2)$  of the increments. It is this feature of a Lévy-Student process that suggests adopting this model to describe the rare escape of particles away from the beam core.

## VI. CONCLUSIONS

In the previous sections we have introduced the Student laws in our SM model for the particle beam dynamics first of all in order to make use of their features depending on their enhanced variance. In particular we have shown that the longer tails with respect to the similar Gaussian distributions can help to account for the finding of a larger than expected number of particles removed far away from the beam core.

It should be remarked, however, that throughout Sec. IV our processes were Gaussian processes since the underlying SDE (1) is still powered by a Brownian noise. This is true even in Sec. IV B where we first introduced the Student laws (55) as stationary distributions of the process. It is only in Sec. V that we introduced a kind of SDE with a Lévy-

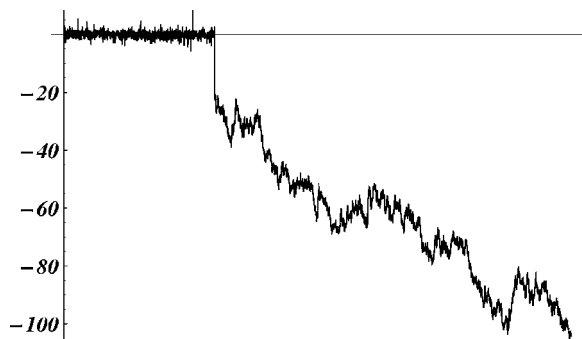


FIG. 12. Rare but possible trajectory of a stationary, Student process  $Y(t)$  [see SDE (78);  $\nu=4$ ]: here the particle definitely drifts away from the core.

Student noise. A more relevant feature of these processes is the fact that their trajectories make jumps: indeed this can become a model for the halo formation in the beams. From a physical point of view these jumps can be produced by occasional hard collisions among the beam particles, the probability of these collisions growing with the intensity of the beam. In some sense it is not only the variance of the transverse distribution of the beam that principally rules the emergence of a halo: in the simulations produced here the variances of the Gaussian and of the Student processes were roughly the same. Rather it is the qualitative character of the process that accounts for the rare escape of the particles from the beam core. For a process produced by a Gaussian noise (a process pathwise continuous: almost every trajectory is everywhere continuous) there is no chance to observe trajectories going out of a well-collimated beam. On the contrary, for a process produced by a Lévy-Student noise (a process only stochastically continuous: trajectories can have jumps) occasionally the jump is large enough to put the particle out of the stream. Of course the frequency and the size of these jumps depend on the parameters  $\nu$  and  $a$  of the process: the jumps tend to be smaller and less frequent when the  $\Sigma(\nu, a^2)$  distributions approximate a Gaussian law. In our opinion it would be very interesting to explore the possibility that the processes underlying the intense beam dynamics are ruled by some sort of Lévy-Student noise rather than by the usual Gaussian noise. It is then important to point out that a some numerical evidence [18] has begun to emerge that confirms this conjecture.

These remarks point to several research directions. First of all it is important to better study the Lévy-Student process in itself: for example a knowledge of the Lévy-Khintchin functions of the Student laws would be relevant to the fine tuning of the frequency and the size of the trajectory jumps. On the other hand even the differential form of its Chapman-Kolmogorov equation [30] would be instrumental to discuss the time evolution of the process. Then it must be remarked that at present we have just defined the Lévy-Student process, but we added no dynamics: it is as if we have the Wiener process, but no stochastic mechanics or any other dynamical model added to this kinematics. In other words we need to build a more generalized SM for the Lévy-Student processes. Finally it would be important at this point to have empirical or numerical data able to corroborate the hypothesis that the increments of the transverse variables of a beam are in fact distributed according to a Student law, rather than according to the usual Gaussian law.

## ACKNOWLEDGMENTS

We would like to thank Dr. C. Benedetti, Professor F. Mainardi, Professor G. Turchetti, and Dr. A. Vivoli for useful comments and suggestions.

## APPENDIX A: INFINITELY DIVISIBLE AND STABLE LAWS

The relevant mathematical concepts used in this paper are better discussed in the framework of the theory of the addi-

tion of independent random variables (RVs): for more details see [31–33]. In the following we will describe the law  $\mathcal{L}$  of a RV  $X$  by giving its characteristic function

$$\varphi(\kappa) = \mathbf{E}(e^{i\kappa X})$$

where  $\mathbf{E}(\cdot)$  is the expectation under the law  $\mathcal{L}$ . When  $\mathcal{L}$  has a PDF  $f(x)$ , then  $\varphi(\kappa)$  is just its Fourier transform. It is well known that the law  $\mathcal{L}$  of the sum of  $n$  independent RVs with laws  $\mathcal{L}_1, \dots, \mathcal{L}_n$  has a characteristic function which is the product of those of the component laws:

$$\varphi(\kappa) = \varphi_1(\kappa) \cdots \varphi_n(\kappa). \quad (\text{A1})$$

On the other hand we say that a law  $\mathcal{L}$  is *decomposed* into the laws  $\mathcal{L}_1, \dots, \mathcal{L}_n$  when its characteristic function can be written as a product (A1) of the characteristic functions of its components. This already allows us to introduce two fundamental concepts: a law  $\mathcal{L}$  with characteristic function  $\varphi$  is said to be ID when for every  $n$  there is a law  $\mathcal{L}_n$  with characteristic function  $\varphi_n$  such that  $\varphi = \varphi_n^n$ . In other words this means that for every  $n$  a RV  $X$  with law  $\mathcal{L}$  can always be decomposed into the sum of  $n$  independent RVs all with the same law  $\mathcal{L}_n$  (identically distributed). Remark, however, that in general the laws  $\mathcal{L}_n$  are not of the same type as  $\mathcal{L}$ . Let us remember here that we say that two laws are of the same *type* when we get one from the other by means of a centering and a rescaling; in other words, if  $\varphi(\kappa)$  is a characteristic function, then all the characteristic functions of the same type have the form  $e^{ia\kappa}\varphi(b\kappa)$  for every  $a$  and  $b > 0$ . For instance all the Gaussian laws  $\mathcal{N}(\mu, \sigma^2)$  belong to the same (Gaussian) type; on the contrary the Poisson laws  $\mathcal{P}(\lambda)$  with different values of  $\lambda$  do not belong to the same type. Now, a law  $\mathcal{L}$  is said to be stable when it is ID and the component laws are of the same type as  $\mathcal{L}$ . More precisely a characteristic function  $\varphi(\kappa)$  is stable when for every  $b, b' > 0$  there exist  $a$  and  $c$  such that

$$\varphi(c\kappa) = e^{ia\kappa}\varphi(b\kappa)\varphi(b'\kappa).$$

As an example, the Gaussian and the Cauchy laws are stable; the Poisson laws are instead only ID. The families of ID and stable laws are completely characterized: in fact the celebrated Lévy-Khintchin formula gives the more general form for the characteristic functions of these two classes; however, while in the case of the stable laws these characteristic functions (albeit not in general the laws themselves) are explicitly known in terms of elementary functions, for the ID laws the characteristic functions are given through an integral containing a function  $L(x)$  (Lévy function) associated with every particular law. But for a few classical cases the Lévy functions of the ID laws are not known.

## APPENDIX B: CENTRAL LIMIT THEOREM AND LÉVY PROCESSES

Let us consider the sequence of RVs  $X_{n,k}$  with  $n \in \mathbb{N}$  and  $k = 1, \dots, n$  with  $X_{n,1}, \dots, X_{n,n}$  independent for every  $n$ . The modern formulation of the central limit problem seeks to find the more general laws which are limits of the laws of the *consecutive sums*

$$S_n = \sum_{k=1}^n X_{n,k}. \quad (\text{B1})$$

Note that these sums generalize the usual partial sums of the classical central limit theorem in that when we go from  $S_n$  to, say,  $S_{n+1}$ , the first  $n$  terms do not in general remain the same: for example,  $X_{n,1}$  does not coincide with  $X_{n+1,1}$ . Under very general technical conditions the central limit theorem now states that the family of all the limit laws of the consecutive sums (B1) coincides with the family of ID laws. The stable laws come into play only when we specialize the form of our consecutive sums: when we have

$$X_{n,k} = \frac{X_k}{a_n} - \frac{b_n}{n}$$

where  $a_n$  and  $b_n$  are sequences of numbers, and  $X_k$  are independent RVs, the consecutive sums take the form of the usual *normed sums* (centered and rescaled sums of independent RVs)

$$S_n = \frac{S_n^*}{a_n} - b_n, \quad S_n^* = \sum_{k=1}^n X_k. \quad (\text{B2})$$

Then, if the  $X_k$  are also identically distributed, the family of the limit laws of the normed sums (B2) coincides with the family of the stable laws. The classical (Gaussian) central limit theorem is an example of convergence toward a stable law; on the other hand the Poisson theorem (convergence of binomial laws toward Poisson laws) is an example of convergence toward an ID law. Every stable law has its own domain of attraction, namely, the set of laws attracted by it in the sense of the convergence of normed sums (B2) of independent RVs all distributed as the attracted law. It can be proved that all the laws with finite variance are in the domain of attraction of the Gauss law, and that a law can be attracted by a non-Gaussian stable law only if it has infinite variance.

The general formulation of the central limit theorem is strictly connected to the definition of the processes with independent increments (decomposable processes). It is apparent in fact that if the increments  $\Delta X(t) = X(t + \Delta t) - X(t)$  for nonsuperposed intervals are independent, the previous forms of the central limit theorem imply that the laws of the increments must be ID laws. Moreover, since the decomposable processes are also Markov processes, the laws of increments are also all that is needed to completely define them. If a decomposable process  $X(t)$  is *stationary* [namely, the law of  $X(t+s) - X(s)$  does not depend on  $s$ ] and *stochastically continuous* [namely, for every  $t$  we have  $X(t + \Delta t) - X(t) \rightarrow 0$  in probability when  $\Delta t \rightarrow 0$ ] we will call it a *Lévy process*. Remark that a Poisson process is a Lévy process since, despite its discontinuities, it is stochastically continuous. In fact these discontinuities do not impair the stochastic continuity of the process because they are *moving* (as opposed to *fixed*) discontinuities. On the other hand it is possible to prove that only the Gaussian Lévy processes (for example the Wiener or the Ornstein-Uhlenbeck processes) are *pathwise continuous*, namely, almost every sample path is everywhere continuous (there are not even moving discontinuities). Now, if

$\varphi(\kappa)$  is the characteristic function of an ID law and  $T$  is a suitable time constant, it is possible to prove that  $[\varphi(\kappa)]^{\Delta t/T}$  is the characteristic function of the increments  $\Delta X(t)$  of a Lévy process. Hence, if the process has a PDF, the stationary transition PDF is

$$p(x,t|y,s) = \frac{1}{2\pi} \mathbb{P} \int_{-\infty}^{+\infty} e^{i\kappa(x-y)} [\varphi(\kappa)]^{(t-s)/T} d\kappa \quad (\text{B3})$$

so that, at least in principle, we know all that is needed to define the process.

The sample paths of a Lévy process are also well characterized: it is possible in fact to prove that almost all trajec-

tories are bounded and are continuous with the exception of a countable set of moving jumps (first kind discontinuities). Then, let us suppose that  $L_t(x)$  is the Lévy-Khintchin function of the ID law of the increment  $X(s+t) - X(s)$ : if  $\nu_t(x)$  is the random number of the jumps in  $[s, s+t)$  of height in absolute value larger than  $x > 0$ , it is possible to prove that

$$|L_t(x)| = \mathbf{E}(\nu_t(x))$$

so that the Lévy-Khintchin function of an ID law plays also the role of a measure of the frequency and height of the trajectory jumps.

- 
- [1] H. Koziol, Los Alamos M. P. Division Report No. MP-3-75-1, 1975 (unpublished).
- [2] M. Reiser, C. R. Chang, D. Kehne, K. Low, T. Shea, H. Rudd, and I. Haber, *Phys. Rev. Lett.* **61**, 2933 (1988).
- [3] R. L. Gluckstern, *Phys. Rev. Lett.* **73**, 1247 (1994); R. L. Gluckstern, W.-H. Cheng, and H. Ye, *ibid.* **75**, 2835 (1995); R. L. Gluckstern, W.-H. Cheng, S. S. Kurennoy, and H. Ye, *Phys. Rev. E* **54**, 6788 (1996); H. Okamoto and M. Ikegami, *ibid.* **55**, 4694 (1997); R. L. Gluckstern, A. V. Fedotov, S. S. Kurennoy, and R. Ryne, *ibid.* **58**, 4977 (1998); T. P. Wangler, K. R. Crandall, R. Ryne, and T. S. Wang, *Phys. Rev. ST Accel. Beams* **1**, 084201 (1998); A. V. Fedotov, R. L. Gluckstern, S. Kurennoy, and R. Ryne, *ibid.* **2**, 014201 (1999); M. Ikegami, S. Machida, and T. Uesugi, *ibid.* **2**, 124201 (1999); J. Quiang and R. D. Ryne, *ibid.* **3**, 064201 (2000).
- [4] O. Boine-Frankenheim and I. Hofmann, *Phys. Rev. ST Accel. Beams* **3**, 104202 (2000); L. Bongini, A. Bazzani, G. Turchetti, and I. Hofmann, *ibid.* **4**, 114201 (2001); A. V. Fedotov and I. Hofmann, *ibid.* **5**, 024202 (2002).
- [5] T. Wangler, *RF Linear Accelerators* (J. Wiley, New York, 1998).
- [6] L. D. Landau and E. M. Lifchitz, *Cintique Physique* (MIR, Moscow, 1990).
- [7] F. Ruggiero, *Ann. Phys. (N.Y.)* **153**, 122 (1984); F. Ruggiero, E. Picasso, and L. A. Radicati, *ibid.* **197**, 396 (1990); J. Struckmeier, *Phys. Rev. ST Accel. Beams* **3**, 034202 (2000).
- [8] W. Paul and J. Baschnagel, *Stochastic Processes: From Physics to Finance* (Springer, Berlin, 2000).
- [9] E. Nelson, *Dynamical Theories of Brownian Motion* (Princeton University Press, Princeton, NJ, 1967); E. Nelson, *Quantum Fluctuations* (Princeton University Press, Princeton, NJ, 1985).
- [10] F. Guerra, *Phys. Rep.* **77**, 263 (1981).
- [11] F. Guerra and L. M. Morato, *Phys. Rev. D* **27**, 1774 (1983).
- [12] S. Alberverio, Ph. Blanchard, and R. Høgh-Krohn, *Expo. Math.* **4**, 365 (1983).
- [13] R. Fedele, G. Miele, and L. Palumbo, *Phys. Lett. A* **194**, 113 (1994), and references therein; S. I. Tzenov, *ibid.* **232**, 260 (1997).
- [14] S. A. Khan and M. Pusterla, *Eur. Phys. J. A* **7**, 583 (2000).
- [15] L. Morato, *J. Math. Phys.* **23**, 1020 (1982).
- [16] R. N. Mantegna and H. E. Stanley, *An Introduction to Econophysics* (Cambridge University Press, Cambridge, U.K., 2000).
- [17] C. K. Allen, K. C. D. Chan, P. L. Colestock, K. R. Crandall, R. W. Garnett, J. D. Gilpatrick, W. Lysenko, J. Qiang, J. D. Schneider, M. E. Schulze, R. L. Sheffield, H. V. Smith, and T. P. Wangler, *Phys. Rev. Lett.* **89**, 214802 (2002).
- [18] A. Vivoli, C. Benedetti, and G. Turchetti, in *Proceedings of the Workshop COULOMB'05, Senigallia, 2005*, *Nucl. Instr. and Meth. A* (to be published).
- [19] S. De Martino, S. De Siena, and F. Illuminati, *Physica A* **271**, 324 (1999); N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, *Phys. Rev. E* **63**, 016501 (2000); N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, in *Quantum Aspects of Beam Physics 2K*, edited by P. Chen (World Scientific, Singapore, 2002), p. 507.
- [20] N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, *Phys. Rev. ST Accel. Beams* **6**, 034206 (2003); in *Quantum Aspects of Beam Physics 2003*, edited by P. Chen *et al.* (World Scientific, Singapore, 2004) p. 36.
- [21] N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, *Int. J. Mod. Phys. B* **18**, 607 (2004).
- [22] N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, in *Proceedings of the European Particle Accelerator Conference EPAC04 (EPS-AG, Geneva, 2004)* p. 2056.
- [23] E. Madelung, *Z. Phys.* **40**, 332 (1926); D. Bohm, *Phys. Rev.* **85**, 166, (1952).
- [24] N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, *J. Phys. A* **32**, 7489 (1999); N. Cufaro Petroni, S. De Martino, S. De Siena, and F. Illuminati, in *Quantum Aspects of Beam Physics 1998*, edited by P. Chen (World Scientific, Singapore, 1999), p. 710; N. Cufaro Petroni, S. De Martino, S. De Siena, R. Fedele, F. Illuminati, and S. I. Tzenov, in *Proceedings of the European Particle Accelerator Conference EPAC98*, edited by S. Myers *et al.* (IOP Publishing, Bristol, 1998), p. 1259.
- [25] N. Cufaro Petroni and F. Guerra, *Found. Phys.* **25**, 297 (1995); N. Cufaro Petroni, in *Quantum Communications and Measurement*, edited by V. P. Belavkin *et al.* (Plenum, New York, 1995), p. 43; N. Cufaro Petroni, S. De Martino, and S. De Siena, in *New Perspectives in the Physics of Mesoscopic Systems*, edited by S. De Martino *et al.* (World Scientific, Singapore, 1997), p. 59.
- [26] N. Cufaro Petroni, S. De Martino, and S. De Siena, *Phys. Lett. A* **245**, 1 (1998).



- [27] L. D. Landau and E. M. Lifchitz, *Mécanique Quantique* (MIR, Moscow, 1988).
- [28] J.-P. Bouchaud and A. Georges, Phys. Rep. **195**, 127 (1990).
- [29] E. Grosswald, Ann. Prob. **4**, 680 (1976); E. Grosswald, Z. Wahrscheinlichkeitstheor. Verwandte Geb. **36**, 103 (1976); M. E. H. Ismail, Ann. Prob. **5**, 582 (1977); L. Bondesson, *ibid.* **7**, 965 (1979); J. Pitman and M. Yor, in *Stochastic Integrals*, edited by D. Williams Lecture Notes in Mathematics Vol. 851, (Springer, Berlin, 1981) p. 285; L. Bondesson, *Generalized Gamma Convolutions and Related Classes of Distributions and Densities*, Lecture Notes in Statistics Vol. 76 (Springer, Berlin, 1992); J. Pitman, Electron. J. Probab. **4**, paper 11 (1999).
- [30] C. W. Gardiner, *Handbook of Stochastic Methods* (Springer, Berlin, 1997).
- [31] M. Loève, *Probability Theory I* (Springer, Berlin, 1977); *Probability Theory II* (Springer, Berlin, 1978).
- [32] W. Feller, *An Introduction to Probability Theory and Its Applications I* (J. Wiley & Sons, New York, 1968); *An Introduction to Probability Theory and Its Applications II* (J. Wiley & Sons, New York, 1971).
- [33] B. V. Gnedenko and A. N. Kolmogorov, *Limit Distributions for Sums of Independent Random Variables* (Addison-Wesley, Reading, MA, 1968).