Notes about Passive Scalar in Large-Scale Velocity Field

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We consider advection of a passive scalar $\theta(t, \mathbf{r})$ by an incompressible large-scale turbulent flow. In the framework of the Kraichnan model the whole PDF’s (probability distribution functions) for the single-point statistics of $\theta$ and for the passive scalar difference $\theta(\mathbf{r}_1) - \theta(\mathbf{r}_2)$ (for separations $\mathbf{r}_1 - \mathbf{r}_2$ lying in the convective interval) are found.

INTRODUCTION

We treat advection of a passive scalar field $\theta(t, \mathbf{r})$ by an incompressible turbulent flow, the role of the scalar can be played by temperature or by pollutants density. The velocity field is assumed to contain motions from some interval of scales restricted from below by $L_v$. A steady situation with a permanent random supply of the passive scalar is considered. We wish to establish statistics of the passive scalar $\theta$ for scales that are less than both the scale $L_v$ and the pumping scale $L$, and larger than the diffusion scale $r_{\text{dif}}$ (for definiteness we assume that $L < L_v$). Such convective interval of scales exists if the Peclet number $\text{Pe} = L/r_{\text{dif}}$ is large enough, we will assume the condition. Since all scales from the convective interval are assumed to be smaller than $L_v$ we will say about advection by a large-scale turbulent flow. The problem is of physical interest for the dimensionalities $d = 2, 3$, but formally it can be treated for an arbitrary dimensionality $d$ of space. Below we will treat $d$ as a parameter. Particularly, all the expressions will be true for a space of large dimensionality $d$.

Description of a small-scale statistics of a passive scalar advected by a large-scale solenoidal velocity field is a special problem in turbulence theory. This problem was treated consistently from the very beginning and some rigorous results have been obtained which is quite unusual for a turbulence problem. Batchelor \cite{1} examined the case of external velocity field being so slow that it does not change during the time of the spectral transfer of the scalar from the external scale to the diffusion scale. Then Kraichnan \cite{2} obtained plenty of results in the opposite limit of a velocity field delta-correlated in time. The pair correlation function of the passive scalar $\langle \theta(\mathbf{r}) \theta(0) \rangle$ was found to be proportional to the logarithm $\ln(L/r)$ and the pair correlation function of the passive scalar difference $\langle [\theta(\mathbf{r}) - \theta(0)]^2 \rangle$ was found to be proportional to $\ln(r/r_{\text{dif}})$ in both cases. The assertions are really correct for any temporal statistics of the velocity field \cite{3,4}. Thus we are dealing with the logarithmic case which is substantially simpler than cases with power-like correlation functions usually encountered in turbulence problems \cite{5,6}.

Now about high-order correlation functions of the passive scalar. As long as all distances between the points are much less than $L$, the $2n$-point correlation functions of $\theta$ are given by their reducible parts (that is are expressed via products of the pair correlation function) until $n \sim \ln(L/r)$ where $r$ is either the smallest distance between the points or $r_{\text{dif}}$ depending on what is larger \cite{4}. The reason for such Wick decoupling is simply the fact that reducible parts contain more logarithmic factors.
(which are considered as large ones) than non-reducible parts do. Consistent calculations of the fourth-order correlation function of the passive scalar at \( d = 2 \) \cite{8} confirm the assertion. Therefore e.g. the single-point PDF of \( \theta \) has a Gaussian core (that describes the first moments with \( n < \ln \text{Pe} \)) and a non-Gaussian tail (that describes moments with \( n > \ln \text{Pe} \)). The tail appears to be exponential (see \cite{3,4}). The same is true for the passive scalar difference \( \Delta \theta = \theta(r) - \theta(0) \) where instead of \( \ln \text{Pe} \) we should take \( \ln(r/r_{\text{diff}}) \). The tails do not depend on \( \ln \text{Pe} \) or on \( \ln(r/r_{\text{diff}}) \) and contain only coefficients depending on the statistics of the advecting velocity.

Correlation functions of the passive scalar can be written as averages of integrals of the pumping along Lagrangian trajectories (see e.g. \cite{9}). For example, the pair correlation function \( \langle \theta(r)\theta(0) \rangle \) is proportional to an average time needed for two points moving along Lagrangian trajectories to run from the distance \( r \) to the distance \( L \). Generally, correlation functions of the passive scalar are determined by the spectral transfer that is by an evolution of Lagrangian separations up to the scale \( L \). For the large-scale velocity field Lagrangian dynamics is determined by the stretching matrix \( \sigma_{\alpha\beta} = \nabla_{\beta}v_\alpha \) and, consequently, the statistics of the matrix determines correlation functions of the passive scalar. For example, the coefficient at the logarithm in the pair correlation function of the passive scalar is \( P_2/\tilde{\lambda}^{1/2} \) \cite{1,2} where \( P_2 \) is the pumping rate of \( \theta^2 \) and \( \tilde{\lambda} \) is Lyapunov exponent that is the average of the largest eigen value of the matrix \( \hat{\sigma} \). The coefficients in the exponential tails are more sensitive to the statistics of \( \hat{\sigma} \), particurlaly they depend on the dimensionless parameter \( \tilde{\lambda}\tau \) \cite{4} where \( \tau \) is the correlation time of \( \hat{\sigma} \). The motion of the fluid particles in the random velocity field resembles in some respects random walks, but one should remember that correlation lengths of both the advecting velocity and of the pumping are much larger than scales from the convective interval we are interested in. Thus the situation is opposite to one usually occurring in solid state physics, where e.g. random potential is shortly correlated in space.

Since really \( \ln(L/r) \) is not very large it is of interest to find the whole PDF’s for the single-point statistics of \( \theta \) and for the passive scalar difference \( \Delta \theta \). It is possible to do for the Kraichnan short-correlated case \( \tilde{\lambda}\tau \ll 1 \) when the statistics of \( \hat{\sigma} \) can be regarded to be Gaussian. An attempt to do this was made in \cite{12,13} in terms of the statistics of the main eigen value of the matrix \( \hat{\sigma} \). Unfortunately the scheme works only for the dimensionality \( d = 2 \) where the matrix \( \hat{\sigma} \) has a single eigen value. This fact was noted in the work \cite{12} where also the correct coefficient in the exponential tails for an arbitrary dimensionality of space \( d \) was found. Here, we develop a scheme enabling to obtain the whole PDF’s for arbitrary \( d \). The scheme is also interesting from the methodical point of view. For example, its modification enables one to calculate the statistics of local dissipation \cite{13}.

The paper is organized as follows. In Section I we find a path integral representation for the simultaneous statistics of the passive scalar. In Section II we analyze the generating functional for correlation functions of the passive scalar in the convective interval of scales. Using different approaches we obtain the functional and establish the applicability conditions of our consideration. In Section III we find explicit expressions for the single-point PDF and for the PDF of the passive scalar difference. In Conclusion we shortly discuss the obtained results.

## I. GENERAL RELATIONS

The dynamics of the passive scalar \( \theta \) advected by the velocity field \( \mathbf{v} \) is described by the equation

\[
\partial_t \theta + \mathbf{v} \nabla \theta - \kappa \nabla^2 \theta = \phi.
\]  

(1.1)

Here, the term with the velocity \( \mathbf{v} \) describes the advection of the passive scalar, the next term is diffusive (\( \kappa \) is the diffusion coefficient) and \( \phi \) describes a pumping source of the passive scalar. Both \( \mathbf{v}(t, \mathbf{r}) \) and \( \phi(t, \mathbf{r}) \) are supposed to be random functions of \( t \) and \( \mathbf{r} \). We regard the statistics of the
velocity and of the source to be independent. Therefore, all correlation functions of θ are to be treated as averages over both statistics.

A. Simultaneous Statistics

The source φ is believed to possess a Gaussian statistics and to be δ-correlated in time. The statistics is entirely characterized by the pair correlation function

$$\langle \phi(t_1,r_1)\phi(t_2,r_2) \rangle = \delta(t_1 - t_2)\chi(|r_1 - r_2|),$$

(1.2)

where we assumed that the pumping is isotropic. The function χ(r) is accepted to have a characteristic scale L which is the pumping length. We will be interested in the statistics of the passive scalar on scales much smaller than L.

Simultaneous correlation functions of the passive scalar θ can be represented as coefficients of the expansion over y of the generating functional

$$Z(y) = \exp \left\{ iy \int dr \beta(r)\theta(0,r) \right\},$$

(1.3)

where β is a function of coordinates and angular brackets designate averaging over both the statistics of the pumping φ and over the statistics of the velocity v. The generating functional Z(y) contains the entire information about the simultaneous statistics of the passive scalar ϑ. Particularly, knowing Z(y) one can reconstruct the simultaneous PDF of the passive scalar, the problem is discussed in Section III.

If characteristic scales of β in (1.7) are much larger than the diffusion scale r_{diff} then it is possible to neglect diffusivity at treating the generating functional (1.3). Then the left-hand side of the equation (1.1) describes the simple advection and it is reasonable to consider a solution of the equation in terms of Lagrangian trajectories ̺(t) introduced by the equation

$$\partial_t ̺ = v(t, ̺).$$

(1.4)

We will label the trajectories by r, those are positions of Lagrange particles at t = 0: ̺(0, r) = r. Next, introducing ̂θ(t, r) = θ(t, ̺) we rewrite the equation (1.1) as \(\partial_t ̂θ = φ\), that leads us to

$$θ(0, r) = \int_{-∞}^{0} dt \phi(t, ̺).$$

(1.5)

Here we have taken into account that at t = 0 the functions θ and ̂θ coincide. Starting from (1.5) and using Gaussianity of the pumping statistics we can average explicitly the generating functional (1.3) over the statistics. The result is

$$Z(y) = \exp \left\{ -\frac{y^2}{2} \int_{-∞}^{0} dt U \right\},$$

(1.6)

$$U = \int dr_1 dr_2 \beta(r_1)\beta(r_2)\chi(|̺_1 - ̺_2|),$$

(1.7)

where angular brackets mean averaging over the statistics of the velocity field only.

Being interested in the single-point statistics of θ we should take β(r) = δ(r). But it is impossible since we have neglected diffusivity. We will take β(r) = δ_Λ(r) instead where the function δ_Λ(r) tends to zero at Λr > 1 fast enough and is normalized by the condition
\[ \int dr \delta_{\Lambda}(r) = 1. \]

Then the generating functional (1.6) will describe the statistics of an object

\[ \theta_{\Lambda} = \int dr \delta_{\Lambda}(r) \theta(r), \] (1.8)

smeared over a spot of the size \( \Lambda^{-1} \). If \( r_{\text{diff}} \Lambda \ll 1 \) then the statistics of the object is not sensitive to diffusivity. From the other hand if \( \Lambda L \gg 1 \), then knowing correlation functions of \( \theta_{\Lambda} \), we can reconstruct single-point statistics due to the logarithmic character of the correlation functions. To obtain single-point correlation functions one should substitute simply \( \Lambda \rightarrow r_{\text{diff}}^{-1} \) into the correlation functions of \( \theta_{\Lambda} \). The above inequalities \( \Lambda r_{\text{diff}} \ll 1 \) and \( \Lambda L \gg 1 \) are compatible because of \( \text{Pe} \gg 1 \).

If we are interested in the statistics of the passive scalar differences in points with a separation \( r_0 \) (where \( r_0 \gg r_{\text{diff}} \)) then instead of \( \delta_{\Lambda}(r) \) we should take

\[ \beta(r) = \delta_{\Lambda}(r - r_0/2) - \delta_{\Lambda}(r + r_0/2). \] (1.9)

Then the generating functional (1.6) will describe the statistics of an object

\[ \Delta \theta_{\Lambda} = \theta_{\Lambda}(r_0/2) - \theta_{\Lambda}(-r_0/2), \] (1.10)

Again, correlation functions of the passive scalar differences can be found from correlation functions of \( \Delta \theta_{\Lambda} \) after the substitution \( \Lambda \rightarrow r_{\text{diff}}^{-1} \).

**B. Path Integral**

Below, we treat advection of the passive scalar by a large-scale velocity field, that is we assume that the velocity correlation length \( L_v \) is larger than the scales from the convective interval. Then for the scales one can expand the difference

\[ v_\alpha(r_1) - v_\alpha(r_2) = \sigma_{\alpha\beta}(t)(r_{1\beta} - r_{2\beta}), \quad \sigma_{\alpha\beta} = \nabla_\beta v_\alpha. \] (1.11)

Here \( \sigma_{\alpha\beta}(t) \) can be treated as an \( r \)-independent matrix field. Then the equation (1.4) leads to

\[ \partial_t (\theta_{1,\alpha} - \theta_{2,\alpha}) = \sigma_{\alpha\beta}(t)(\theta_{1,\beta} - \theta_{2,\beta}), \] (1.12)

A formal solution of the equation (1.12) is

\[ \theta_{1,\alpha} - \theta_{2,\alpha} = W_{\alpha\beta}(r_{1\beta} - r_{2\beta}), \]
\[ \partial_t \hat{W} = \hat{\sigma} \hat{W}, \quad \hat{W} = \mathcal{T} \exp \left( - \int_t^0 dt \hat{\sigma} \right), \] (1.13)

where \( \mathcal{T} \) means an antichronological ordering. Note that \( \text{Det} \hat{W} = 1 \), those property is a consequence of \( \text{Tr} \hat{\sigma} = 0 \) and the initial condition \( \hat{W} = 1 \) at \( t = 0 \). The Lagrangian difference in (1.7) is now rewritten as

\[ |\varrho_1 - \varrho_2| = \sqrt{(r_{1\alpha} - r_{2\alpha}) B_{\alpha\beta}(r_{1\beta} - r_{2\beta})}, \quad \hat{B} = \hat{W}^T \hat{W}, \] (1.14)

where the subscript \( T \) designates a transposed matrix. Note that \( \text{Det} \hat{B} = 1 \) since \( \text{Det} \hat{W} = 1 \).

The generating functional \( Z(y) \) (1.6) can be explicitly calculated in the Kraichnan case \[2\] when the statistics of the velocity is \( \delta \)-correlated in time. Then the velocity statistics is Gaussian and is entirely determined by the pair correlation function which in the convective interval is written as
\[ \langle v_\alpha(t_1, r_1) v_\beta(t_2, r_2) \rangle = \delta(t_1 - t_2) \left[ V_0 \delta_{\alpha\beta} - K_{\alpha\beta}(r_1 - r_2) \right], \quad (1.15) \]

\[ K_{\alpha\beta}(r) = D(r^2 \delta_{\alpha\beta} - r_\alpha r_\beta) + \frac{(d - 1)D}{2} \delta_{\alpha\beta} r^2. \quad (1.16) \]

Here \( V_0 \) is a huge \( r \)-independent constant and \( D \) is a parameter characterizing the amplitude of the strain fluctuations. The structure of (1.14) is determined by the isotropy and space homogeneity accepted, and by the incompressibility condition \( \text{div} \, \mathbf{v} = 0 \). Then the statistics of \( \tilde{\sigma} \) is Gaussian and is determined by the pair correlation function which can be found from (1.15,1.16):

\[ \langle \sigma_{\alpha\beta}(t_1) \sigma_{\mu\nu}(t_2) \rangle = D \left[ (d + 1) \delta_{\alpha\mu} \delta_{\beta\nu} - \delta_{\alpha\nu} \delta_{\beta\mu} - \delta_{\alpha\beta} \delta_{\mu\nu} \right] \delta(t_1 - t_2). \quad (1.17) \]

Note that the correlation function (1.17) is \( r \)-independent, as it should be. We see from (1.17) that the parameter \( D \) characterizes amplitudes of \( \tilde{\sigma} \) fluctuations.

Averaging over the statistics of \( \tilde{\sigma} \) can be substituted by the path integral over unimodular matrices \( \hat{W}(t) \) with a weight \( \exp(i\mathcal{L}) \). The effective action \( \mathcal{I} = \int dt \mathcal{L}_0 \) is determined by (1.17):

\[ i\mathcal{L}_0 = -\frac{1}{2d(d + 2)D} \left[ (d + 1) \text{Tr}(\tilde{\sigma}^T \hat{\sigma}) + \text{Tr} \hat{\sigma}^2 \right]. \quad (1.18) \]

Then the generating functional (1.17) can be rewritten as the following functional integral over unimodular matrices

\[ \mathcal{Z}(y) = \int D\hat{W} \exp \left[ \int_0^t dt \left( i\mathcal{L}_0 - \frac{y^2}{2U} \right) \right], \quad (1.19) \]

\[ U = \int d\mathbf{r}_1 d\mathbf{r}_2 \beta(\mathbf{r}_1) \beta(\mathbf{r}_2) \chi \left[ \sqrt{(r_{1\alpha} - r_{2\alpha})B_{\alpha\beta}(r_{1\beta} - r_{2\beta})} \right]. \quad (1.20) \]

Here, we should substitute \( \hat{\sigma} = \partial_t \hat{W}(\hat{W})^{-1} \) and remember about the boundary condition \( \hat{W} = 1 \) at \( t = 0 \).

Some words about the ‘potential’ \( U \) (1.17) figuring in (1.20). The characteristic value of \( \mathbf{r}_1 - \mathbf{r}_2 \) in the integral (1.4) is of order \( \Lambda^{-1} \) for \( \beta(\mathbf{r}) = \delta_\Lambda(\mathbf{r}) \). Since we assume \( \Lambda L \gg 1 \) then for a single-point statistics \( U \approx P_2 \), where \( P_2 = \chi(0) \), if \( B \) is not very large. Particularly it is correct at moderate times \( |t| \) since \( \hat{B} = \hat{1} \) at \( t = 0 \). With increasing \( |t| \) the argument of \( \chi \) in (1.20) grows and \( U \) tends to zero when the argument of \( \chi \) becomes greater than \( L \). For the passive scalar difference when \( \beta \) is determined by (1.4) the situation is a bit more complicated. Then \( U \) is a difference of two contributions. The first contribution behaves like for a single-point statistics. The second contribution contains \( \chi \) with the argument determined by \( \mathbf{r}_1 - \mathbf{r}_2 \approx \pm \mathbf{r}_0 \). Then at \( t = 0 \) the meaning of the second contribution is determined again by \( P_2 \) but it turns to zero at increasing \( |t| \) earlier than the first contribution.

Thus we reduced our problem to the path integral (1.19) that is to the quantum mechanics with \( d^2 - 1 \) degrees of freedom. Nevertheless to solve the problem we should perform an additional reduction of the degrees of freedom. The conventional way to do this is passing to eigen values, say, of the matrix \( \hat{B} \) figuring in (1.20) (see e.g. (14)) and excluding angular degrees of freedom. Just this way was used by Bernard, Gawedzki and Kupiainen in (12). Then the authors using known facts about the quantum mechanics associated with the eigenvalues (see e.g. (14)) have found the coefficient in the exponential tail in the single-point PDF of \( \theta \). Unfortunately this way is not very convenient to find the whole PDF. To do this we will use a special representation of the matrix \( \hat{W} \) in the spirit of the nonlinear substitution introduced by Kolokolov (13). That is a subject of the next subsection.
C. Choice of Parametrization

To examine the generating functional $Z(y)$ we will use a mixed rotational-triangle parametrization
\[ \hat{W} = \hat{R}\hat{T}, \quad \hat{B} = \hat{T}^T\hat{T}, \] (1.21)
where $\hat{R}$ is an orthogonal matrix and $\hat{T}$ is a triangular matrix: $T_{ij} = 0$ at $i > j$. The parametrization (1.21) is the direct generalization of the $2d$ substitution suggested in [17]. Note that $\text{Det} \hat{W} = 1$. Note also that the matrix $\hat{B}$ introduced by (1.14) does not depend on $\hat{R}$ as is seen from (1.21). That is a motivation to exclude the matrix $\hat{R}$ from the consideration performing the integration over the corresponding degrees of freedom in the path integral (1.19). A jacobian appears at the integration. To avoid an explicit calculation of the jacobian (which needs a discretization over time and then an analysis of an infinite matrix [10]) we will use an alternative procedure described below.

Let us examine the dynamics of the matrix $\hat{T}$. It is determined by the equation
\[ \partial_t T_{ij} = \Sigma_{ii}T_{ij} + \sum_{i<k\leq j} (\Sigma_{ik} + \Sigma_{ki})T_{kj}, \] (1.22)
following from (1.13, 1.21). Here we used the designation
\[ \hat{\Sigma} = \hat{R}^T\hat{\sigma}\hat{R}. \] (1.23)

Next introducing the quantities
\[ T_{ii} = \exp(\rho_i), \quad T_{ij} = \exp(\rho_i)\eta_{ij}, \quad \text{if } i < j, \] (1.24)
we rewrite the equation (1.22) as
\[ \partial_t \rho_i = \Sigma_{ii}, \quad \partial_t \eta_{ij} = (\Sigma_{ij} + \Sigma_{ji})\exp(\rho_j - \rho_i) + \sum_{i<k<j} (\Sigma_{ik} + \Sigma_{ki})\exp(\rho_k - \rho_i)\eta_{kj}. \] (1.25, 1.26)

Comparing (1.13) with (1.21) one can find the following expression for $\hat{A} = \hat{R}^T\partial_t\hat{R}$
\[ A_{ij} = \Sigma_{ij} \quad \text{if } i > j \quad A_{ij} = -\Sigma_{ji} \quad \text{if } i < j. \] (1.27)

One can easily check that the irreducible pair correlation function of $\Sigma_{ij}$ has the same form as for $\sigma_{ij}$ that is determined by (1.17):
\[ \langle \Sigma_{ij}(t_1)\Sigma_{mn}(t_2) \rangle = D[(d+1)\delta_{im}\delta_{jn} - \delta_{in}\delta_{jm} - \delta_{ij}\delta_{mn}](t_1 - t_2). \] (1.28)

Besides, the average value of $\Sigma_{ij}$ is nonzero [11]:
\[ \langle \Sigma_{ij} \rangle = -D\frac{d(d - 2i + 1)}{2}\delta_{ij}. \] (1.29)

Nonzero averages of $\Sigma_{ij}$ are related to Lyapunov exponents (not only the first one) [18] (for our model see also [19]). To obtain (1.29) one should first solve the equation $\hat{A} = \hat{R}^T\partial_t\hat{R}$ for $\hat{R}$ on a small interval $\tau$:
\[ \hat{R}(t + \tau) \approx \hat{R}(t) \left[ 1 + \int_t^{t+\tau} \hat{A}(t') dt' \right]. \]
Then with the same accuracy we get from the expression (1.23)

\[
\hat{\Sigma}(t + \tau) \approx \hat{R}^T(t) \hat{\sigma}(t + \tau) \hat{R}(t) + \left[ \hat{\Sigma}(t + \tau), \int_t^{t+\tau} dt' \hat{A}(t') \right].
\] (1.30)

The average value of \( \hat{\Sigma} \) arises from the second term in the right-hand side of (1.30). The explicit form of the average can be found using

\[
\langle \Sigma_{ij}(t + \tau) \int_t^{t+\tau} dt' \Sigma_{mn}(t') \rangle = \frac{D}{2} [(d + 1)\delta_{im}\delta_{jn} - \delta_{in}\delta_{jm} - \delta_{ij}\delta_{mn}].
\] (1.31)

Here we utilized Eq. (1.28) and substituted the integral

\[
\int_t^{t+\tau} dt' \delta(t + \tau - t')
\]
by 1/2. The reason is that really the correlation function of \( \hat{\sigma} \) has a finite correlation time and therefore \( \delta(t) \) (representing this correlation function) should be substituted by a narrow function symmetric under \( t \to -t \). Then we will get 1/2. Expressing in (1.30) \( \hat{A} \) via \( \hat{\Sigma} \) from (1.27) and calculating its average using (1.31) we get the answer (1.29).

The expressions (1.25, 1.26, 1.28, 1.29) entirely determine the stochastic dynamics of \( \rho_i \) and \( \eta_{ij} \). Using the conventional approach [20–24] correlation functions of these degrees of freedom can be described in terms of a path integral over \( \rho_i, \eta_{ij} \) and over auxiliary fields which we will designate as \( m_i \) and \( \mu_{in} \) \((i < n)\). This integral should be taken with the weight \( \exp(i\int dt L) \) where the Lagrangian is

\[
L = \sum_{a=1}^{d} m_a \left[ \partial_t \rho_a + D \left( d - 2a + 1 \right) \right] + \frac{iD}{2} \left[ d \sum_a m_a^2 - \left( \sum_a m_a \right)^2 \right] + D \sum_{i<j} \exp(2\rho_j - 2\rho_i) \mu_{ij}^2 + 2iD \sum_{i<k<j} \mu_{ij} \mu_{ik} \exp(2\rho_k - 2\rho_i) \eta_{kj}
\]
\[
+ \sum_{i<j} \mu_{ij} \partial_t \eta_{ij} + iDd \sum_{i<k<m,n} \mu_{im} \mu_{in} \eta_{km} \eta_{kn} \exp(2\rho_k - 2\rho_i).
\] (1.32)

Since the matrix \( \hat{B} \) in accordance with (1.21) does not depend on \( \hat{R} \) it is enough to know the statistics of \( \rho_a \) and \( \eta_{ij} \) to determine the average (1.6). Therefore, instead of (1.19) we get

\[
\mathcal{Z}(y) = \int D\rho D\eta Dm D\mu \exp \left[ i \int_{-\infty}^{0} dt \left( iL - \frac{y^2}{2} U \right) \right].
\] (1.33)

Here \( U \) is determined by (1.20) where (1.21, 1.24) should be substituted.

Thus we obtained the expression for the generating functional (1.3) in terms of the functional (path) integral which is convenient for an analysis which is presented in the subsequent section.

II. GENERATING FUNCTIONAL

Here, we are going to calculate the generating functional (1.3) for a single-point statistics of \( \theta \) that is of the object (1.8) corresponding to \( \beta(r) = \delta_{\Lambda}(r) \) and also the statistics of the difference that is of the object (1.10) corresponding to (1.9). The starting point for the subsequent consideration is the expression (1.33). There are different ways to examine \( \mathcal{Z}(y) \). We will describe two schemes leading to the same answer but carrying in some sense a complimentary information. We believe also that the consideration of the different schemes is useful from the methodical point of view. A modification of the second scheme is presented in Appendix.
A. Saddle-Point Approach

The first way to obtain the answer for the generating functional (1.3) is in using the saddle-point
approximation for the path integral (1.33). The inequalities justifying the approximation are $\Lambda L \gg 1$
for the object (1.8) and $\Lambda r \gg 1$ for the object (1.10).

As we will see large values of the differences $\rho_i - \rho_k$ ($i < k$) will be relevant for us. At the condition
fluctuations of $\eta$ and $\mu$ are suppressed and it is possible to neglect the fluctuations. Therefore we can
omit the integration over $\eta$ and $\mu$ in (1.33) substituting $\eta = \mu = 0$ into (1.32). After that we obtain
a reduced Lagrangian. Introducing an auxiliary field $s$ we can rewrite the reduced Lagrangian as

$$i\mathcal{L}_t = i \sum_{a=1}^{d} m_a \left[ \partial_t \rho_a + D \frac{d(d+1-2a)}{2} + s \right] - \frac{D d}{2} \sum_{a} m_a^2 - \frac{s^2}{2D}. \tag{2.1}$$

Now, to obtain $Z(y)$ one should integrate the exponent in (1.33) (with $\mathcal{L}_t$) over $\rho_a$, $m_a$ and $s$. To examine (2.1) it is convenient to pass to new variables $\phi_a = O_{ab} \rho_b$ and $\tilde{m}_a = O_{ab} m_b$ where $O$ is an
orthogonal matrix. We make the following transformation

$$\phi_1 = \sqrt{\frac{3}{d(d^2-1)}} [(d-1) \rho_1 + (d-3) \rho_2 + \ldots + (1-d) \rho_d], \tag{2.2}$$

$$\phi_2 = \ldots, \ldots, \phi_d = \frac{1}{\sqrt{d}} [\rho_1 + \rho_2 + \ldots + \rho_d].$$

Then the expression (2.1) will be rewritten as

$$i\mathcal{L}_t = i \sum_{a=1}^{d} \tilde{m}_a \partial_t \phi_a + i\sqrt{d} \tilde{m}_d s - \frac{s^2}{2D} - \frac{D d}{2} \sum_{a} \tilde{m}_a^2 + i \frac{D d}{2} \sqrt{\frac{d(d^2-1)}{3}} \tilde{m}_1. \tag{2.3}$$

The Lagrangian (2.3) is a sum over different degrees of freedom. The dynamics of $\phi_1$ is ballistic
whereas the dynamics of $\phi_a$ for $a > 1$ is purely diffusive. We will see that times determining the
main contribution to the generating functional are large enough so that for the relevant region $\phi_1 \gg \phi_a$. Therefore the potential $U$ (1.20) depends practically only on $\phi_1$ and it is possible to
integrate explicitly over $s$, $\phi_a$ and $\tilde{m}_a$ for $a > 1$. After that we stay only with one degree of freedom
which is described by the Lagrangian

$$i\mathcal{L}_1 = i \tilde{m}_1 \left( \partial_t \phi_1 + \frac{D d}{2} \sqrt{\frac{d(d^2-1)}{3}} \right) - \frac{D d}{2} \tilde{m}_1^2. \tag{2.4}$$

Neglecting all $\phi_a$ for $a > 1$ and performing the transformation inverse to (2.2) we obtain

$$\rho_1 \approx \sqrt{\frac{3(d-1)}{d(d+1)}} \phi_1, \quad \rho_a \approx \frac{d - 2a + 1}{d-1} \rho_1. \tag{2.5}$$

We will soon see that the characteristic value $\phi_1 \gg 1$. Therefore the characteristic value of $e^{\rho_1}$ is
much larger than other $e^{\rho_a}$ and we conclude that the potential $U$ depends really only on $\rho_1$. For
the case of the single-point statistics the characteristic value of the difference $r_1 - r_2$ in (1.20) is
$\Lambda^{-1}$. Then it follows from (1.21)[1.24] that the potential $U$ falls down from $P_2$ to zero near the point
$\rho_1 = \ln(\Lambda A)$ that is near the point $\phi_1 = \phi_\Lambda$ where

$$\phi_\Lambda = \sqrt{\frac{d(d+1)}{3(d-1)}} \ln(\Lambda A). \tag{2.6}$$
For the difference the potential increases from zero to $2P_2$ at $\phi_1 = \phi_R$ where

$$\phi_R = \sqrt{\frac{d(d+1)}{3(d-1)}} \ln \left( \frac{L}{r_0} \right),$$  \hspace{1cm} (2.7)

and then falls down from $P_2$ to zero near the point $\phi_1 = \phi_Λ$. The expressions (2.6,2.7) determine the characteristic values of $\phi_1$ which are really large since $LΛ \gg 1$ or $L/r_0 \gg 1$ what justifies our conclusions.

Now we can employ the saddle-point approximation:

$$\ln \mathcal{Z}(y) \approx \int_{-\infty}^{0} dt \left( i\mathcal{L}_1 - \frac{y^2 U}{2} \right)_{\text{inst}},$$  \hspace{1cm} (2.8)

where we should substitute solutions of the extremum conditions which we will call instantonic equations. The instantonic equations which can be found as extrema conditions for $i\mathcal{L}_1 - y^2 U/2$ are

$$\partial_t \phi_1 + \frac{Dd}{2} \sqrt{\frac{d(d^2-1)}{3}} = -iDd\tilde{m}_1,$$  \hspace{1cm} (2.9)

$$\partial_t \tilde{m}_1 = i\frac{y^2}{2} \frac{\partial U}{\partial \phi_1}.$$  \hspace{1cm} (2.10)

The equations conserve the ‘energy’

$$-i\frac{Dd}{2} \tilde{m}_1 \sqrt{\frac{d(d^2-1)}{3}} + \frac{Dd}{2} \tilde{m}_1^2 + \frac{y^2 U}{2} = 0.$$  \hspace{1cm} (2.11)

The conservation law is satisfied since $i\mathcal{L}_1 - y^2 U/2$ does not explicitly depend on $t$. The ‘energy’ (2.11) is equal to zero since at $t \to -\infty$ the value of $\tilde{m}_1$ should tend to zero. The property can be treated as the extremum condition appearing at variation of $i\mathcal{L}_1 - y^2 U/2$ over the initial value of $\phi_1$. Equating (2.11) to zero we can express $\tilde{m}_1$ in terms of $\phi_1$. Next, since (2.11) is zero then the saddle-point value of $\mathcal{Z}(y)$ (2.8) can be written as $i\int d\phi_1 \tilde{m}_1$, where the integral over $\phi_1$ goes from zero to infinity.

Substituting into $i\int d\phi_1 \tilde{m}_1$ the expression of $\tilde{m}_1$ in terms of $\phi_1$ we get for the single-point statistics

$$\ln \mathcal{Z}(y) \simeq \frac{d(d+1)}{6} \left[ 1 - \sqrt{1 + \frac{12y^2 P_2}{Dd^2(d^2-1)}} \right] \ln(LΛ).$$  \hspace{1cm} (2.12)

Note that the expression (2.12) has (as a function of $y$) two branch points $y = \pm iy_{\text{sing}}$ where

$$y_{\text{sing}}^2 = \frac{Dd^2(d^2-1)}{12P_2}.$$  \hspace{1cm} (2.13)

The same procedure can be done for the passive scalar difference, or, more precisely, for the object (1.10). Taking into account the presence of the jumps (2.7) and (2.6) in the potential $U$ we get an answer slightly different from (2.12)

$$\ln \mathcal{Z}(y) \simeq \frac{d(d+1)}{6} \left[ 1 - \sqrt{1 + \frac{24y^2 P_2}{Dd^2(d^2-1)}} \right] \ln(r_0Λ),$$  \hspace{1cm} (2.14)

$$y_{\text{sing}}^2 = \frac{Dd^2(d^2-1)}{24P_2}.$$  \hspace{1cm} (2.15)
Note that (2.14) does not depend on the pumping scale \(L\) but still depend on the cutoff \(\Lambda\).

The characteristic value of \(\phi_1\) is determined by the quantity (2.6) which is much larger than unity. Then it follows from (2.5) that \(\exp(2\rho_j - 2\rho_i) \ll 1, i > j\), (excluding for a short initial stage of the evolution) and we see from (1.32) that fluctuations of the fields \(\eta\) are suppressed in comparison, say, with \(\rho_a\). This justifies neglecting the fields \(\eta\) and \(\mu\) leading to the reduced Lagrangian (2.1). Next, dynamics of \(\phi_a\) for \(a > 1\) is diffusive and it follows from (2.3) that the characteristic value of \(\phi_a\) can be estimated as \(\sqrt{Dd|t|}\). As follows from (2.3) \(\partial_t \phi_1 \sim Dd^{5/2}\) and we find from (2.6) the instantonic life time

\[
t_{lt} = D^{-1}d^{-2}\ln(L\Lambda),
\]

which determines times producing nonzero contribution to the effective action. At \(|t| \sim t_{lt}\) the characteristic values of \(\phi_a\) for \(a > 1\) are of order \(\sqrt{\ln(L\Lambda)/d}\) and we conclude that

\[
\frac{\phi_a}{\phi_1} \sim \frac{1}{d\sqrt{\ln(L\Lambda)}} \ll 1,
\]

at times \(|t| \sim t_{lt}\). The inequality (2.17) justifies passing to the Lagrangian (2.4). The same arguments can be applied to the generating functional for the passive scalar difference, the only modification is in the substitution \(\ln(L\Lambda) \rightarrow \ln(r_0\Lambda)\).

There are also additional applicability conditions for the answers (2.12,2.14). To establish the conditions one should go beyond the main order of the saddle-point approximation. It will be more convenient for us to develop an alternative scheme which enables to find the conditions simpler. That is the subject of the next subsection.

### B. Schrödinger equation

Here we present another way to get the answers (2.12,2.14). As previously we start from the path integral representation (1.33) for the generation functional \(Z(y)\).

Unfortunately it is impossible to get a closed equation for \(Z(y)\). To avoid the difficulty we introduce an auxiliary quantity

\[
\Psi(t, y, \rho_0, \eta_0) = \int \mathcal{D}\rho \mathcal{D}\eta \mathcal{D}m \mathcal{D}\mu \exp \left[ \int_{-t}^0 dt' \left( i\mathcal{L} - \frac{y^2}{2} U \right) \right] |_{\rho(-t) = \rho_0, \eta(-t) = \eta_0}.
\]

It follows from the definition (2.18) that

\[
Z(y) = \lim_{t \to \infty} \int \prod d\rho_a d\eta_{ij} \Psi(t, y, \rho, \eta).
\]

The equation for the function \(\Psi\) can be obtained from the expression (1.32) and the definition (2.18):

\[
\partial_t \Psi = \frac{Dd}{2} \left[ \sum_{i=1}^d \frac{\partial^2}{\partial \rho_i^2} - \frac{1}{d} \left( \sum_{i=1}^d \frac{\partial}{\partial \rho_i} \right)^2 - \sum_{i=1}^d (d-2i+1) \frac{\partial}{\partial \rho_i} \right.

+ 2 \sum_{i<j} \exp(2\rho_j - 2\rho_i) \frac{\partial^2}{\partial \eta_{ij}^2} + 4 \sum_{i<k<j} \exp(2\rho_k - 2\rho_i) \frac{\partial}{\partial \eta_{ij}} \frac{\partial}{\partial \eta_{ik}} \eta_{kj}

\left. + 2 \sum_{i<k<m,n} \exp(2\rho_k - 2\rho_i) \frac{\partial}{\partial \eta_{km}} \frac{\partial}{\partial \eta_{in}} \eta_{kn} \right] \Psi - \frac{y^2 U}{2} \Psi.
\]

\[
(2.20)
\]
We see that the equation (2.20) for $\Psi$ resembles the Schrödinger equation. The initial condition to the equation can be found directly from the definition (2.18):

$$
\Psi(t = 0, y, \rho, \eta) = \prod \delta(\rho_a)\delta(\eta_{ij}).
$$

(2.21)

The value of $Z$ in accordance with (2.19) is determined by the integral of $\Psi$ over $\eta$ and $\rho$. This integral is equal to unity at $t = 0$ and then varies with increasing time $t$ due to $U \neq 0$ since only the term with $U$ in (2.21) brakes conservation of the integral. Thus to find $Z$ we should establish an evolution of the function $\Psi$ from $t = 0$ to large $t$.

Below we concentrate on the single-point statistics. The scheme can be obviously generalized for the passive scalar difference.

Let us first describe the evolution qualitatively. The initial condition (2.21) shows that at $t = 0$ the function $\Psi$ is concentrated at origin. Then it undergoes spreading in all directions except for $\rho_1 + \ldots + \rho_d$ since the operator in the right-hand side of (2.20) commutes with $\rho_1 + \ldots + \rho_d$. This is a consequence of the condition Det $\hat{T} = 1$ (to be satisfied) which implies that during the evolution $\rho_1 + \ldots + \rho_d = 0$. That means that we should treat a solution of (2.20) which is $\Psi \propto \delta(\rho_1 + \ldots + \rho_d)$.

The function $\Psi$ is smeared diffusively with time and also moves as a whole in some direction, which is determined by the term with the first derivative in (2.20). The rate of the ballistic motion is

$$
\langle \partial_t \rho_i \rangle = D \frac{d(d - 2i + 1)}{2}.
$$

(2.22)

Therefore $\Psi$ describes a cloud, centre of which moves according to the law

$$
\rho_i = D \frac{d(d - 2i + 1)}{2} t.
$$

(2.23)

Effective diffusion coefficients for $\eta$’s fall down with increasing $t$ since in accordance with (2.23) the differences $\rho_k - \rho_i$ [figuring in (2.20)] are negative and grow by their absolute value. Therefore the diffusion over $\eta$ stops when the characteristic values of $\rho_i - \rho_k$ becomes greater than unity. Note that the ‘frozen’ values of $\eta$ do not depend on $y$ since $U$ can be considered as uniform during the initial stage of the evolution. After that $\eta$ are frozen, the diffusion continues only over $\rho$’s. If the cloud is inside the region where $U \simeq P_2$ then the evolution of the cloud is not influenced by $U$. After the period of time $t_{lt}$ (2.16) the cloud reaches a barrier where the potential $U$ falls down from $P_2$ to 0. The subsequent history depends on the value of $y$. For moderate $y$ the cloud passes this barrier and continues to move with the same rate. After this the integral of $\Psi$ will not change in time, and its value will determine the generating functional $Z(y)$. Naive estimations would give the answer $\ln Z(y) = -y^2 t_{lt}/2$, which reproduces the value of the pair correlation function of $\theta$.

A special consideration is needed if $|y| \gg y_{sing}$ or if $y$ is close to $\pm iy_{sing}$ where $y_{sing}$ is introduced by (2.13). Just the last region determined the PDF’s and is consequently of a special interest. Note that $y = \pm iy_{sing}$ corresponds to appearance of a bound state near the pumping boundary (where $U$ falls down from $P_2$ to zero). If $y \gg y_{sing}$ then the front side of the cloud reaches the jump of the potential much earlier than $t_\Lambda$. The residue of the cloud (living inside the potential well) is damped due to the term with $y$ and don’t give a contribution to $Z(y)$. If $|y| \gg y_{sing}$ then $Z(y) \gg \exp(-y^2 t_{lt}/2)$, really the asymptotics of $Z(y)$ is exponential in the case. If $|y \pm iy_{sing}| \ll y_{sing}$ then the cloud stays near the pumping boundary for a long time, that is the shape of $\Psi$ inside the region $U \simeq P_2$ varies in time comparatively slow. Besides, a part of $\Psi$ percolates out to the region where $U \simeq 0$ and the integral of $\Psi$ grows with increasing $|t|$. As $y$ come to $iy_{sing}$ closer this stage lasts longer. One can say that the back side of the cloud $\Psi$ gives the right answer for $Z(y)$. The important point is that if $y$ is not very close to $iy_{sing}$ then during the time of exiting $\Psi$ from the potential the width of $\Psi$
in terms of diffusive degrees of freedom is much less than $\ln L\Lambda$. This means that the function $\Psi$ is really narrow, that justifies our consideration.

For an quantitative analysis it is convenient to pass to the variables $\phi_i$ (2.2). Since an $\eta$-dependence of $\Psi$ is frozen after an initial part of evolution, then it is possible to obtain an equation for the integral of $\Psi$ over $\eta$:

$$\tilde{\Psi}(\phi_1, \ldots, \phi_{d-1}) = \int d\phi_d \prod d\eta_j \Psi,$$

where we included also an integration over $\phi_d$ to remove the factor $\delta(\rho_1 + \ldots + \rho_d)$. The equation for the function (2.24) is

$$\partial_t \tilde{\Psi} = \frac{Dd}{2} \left[ \frac{\partial}{\partial \phi_1} - \sqrt{\frac{d(d^2 - 1)}{3}} \frac{\partial}{\partial \phi_1} \right] \tilde{\Psi} - \frac{y^2 \tilde{U}}{2} \tilde{\Psi},$$

(2.25)

where $\tilde{U}$ is function of $\phi_a$ only which can be found by substituting into $U$ the ‘frozen’ values of $\eta$’s. Qualitatively $\tilde{U}$ has the same structure as $U$ itself. One can conclude from (2.25) that the cloud described by $\tilde{\Psi}$ moves ballistically into the $\phi_1$ direction and spreads along other directions. We are going to treat the situation when the cloud remains narrow during the relevant part of the evolution. Then one can integrate $\tilde{\Psi}$ over all $\phi_i$, $i > 1$ in a similar way as in the case with $\eta$’s and get an 1d equation for

$$\tilde{\Psi}(\phi_1) = \int \prod_{i=2}^{d-1} d\phi_i \tilde{\Psi}.$$

The function $\tilde{\Psi}$ satisfies the equation

$$\partial_t \tilde{\Psi} = \frac{Dd}{2} \left[ \frac{\partial}{\partial \phi_1} - \sqrt{\frac{d(d^2 - 1)}{3}} \frac{\partial}{\partial \phi_1} \right] \tilde{\Psi} - \frac{y^2 \tilde{U}}{2} \tilde{\Psi},$$

(2.26)

The initial condition to Eq. (2.26) is $\tilde{\Psi}(t = 0) = \delta(\phi_1)$. The potential $\tilde{U}$ is obtained from $\tilde{U}$ by substitution $\phi_a \to 0$ for $a > 0$. In fact, on the direction (2.23) we have that is the potential $\tilde{U}$ depends only on $\rho_1$. The barrier is achieved when $\rho_1 \sim \ln L\Lambda$. Passing to the variables $\phi_1$ we conclude that the potential $\tilde{U}$ diminishes from $P_2$ at $\phi_1 < \phi_1^\ast$ to zero at $\phi_1 > \phi_1^\ast$ where $\phi_1^\ast$ is defined by (2.6).

The character of a solution of the equation (2.26) can be analyzed semiquantitatively in terms of the width $l$ of $\tilde{\Psi}$ over $\phi_1$ and its amplitude $h$. When $\tilde{\Psi}$ reaches the pumping boundary then it stops there for a period of time. Then the width $l$ and the amplitude $h$ governed by the equations

$$\frac{dl}{dt} = -Dd\lambda + \frac{Dd}{l}, \quad \frac{dh}{dt} = -\frac{Ddh}{l^2} - \frac{y^2 P_2 h}{2},$$

(2.27)

where $\lambda = \sqrt{d(d^2 - 1)/12}$. $Dd\lambda$ is the rate of the cloud motion along $\phi_1$ direction (when $U = \text{const}$), and $Dd$ is the diffusion coefficient for $\phi_1$ direction. One can estimate from the 1st equation the width $l$ as $l \sim 1/\lambda$. Then from the 2nd equation the height $h$ falls down or grows in time depending on $y$. The characteristic $y$ where the regime changes is of the order $|y_{\text{sing}}|^2 \sim Dd\lambda^2/P_2$. We’ll show this by consistent calculations.

The equation (2.26) can be solved analytically, e.g., by the Laplace transform over time $t$. Making the Laplace transform one gets

$$p\tilde{\Psi}(p) - \delta(\phi_1) = \frac{Dd}{2} \left[ \frac{\partial}{\partial \phi_1} - \sqrt{\frac{d(d^2 - 1)}{3}} \frac{\partial}{\partial \phi_1} \right] \tilde{\Psi}(p) - \frac{y^2 \tilde{U}(\phi_1)}{2} \tilde{\Psi}(p),$$

(2.28)
We are interesting in the bound state describing by the equation. Solutions for \( \bar{\Psi}(p) \) in the intervals \((-\infty, 0), (0, \phi_{A}), (\phi_{A}, \infty) \) are exponential, one should match them. The function \( \bar{\Psi}(p) \) as a function of \( p \) has two branch points at

\[
p_1 = -\frac{Dd^2(d^2 - 1)}{24} - \frac{y^2P_2}{2}, \quad p_2 = -\frac{Dd^2(d^2 - 1)}{24},
\]

coming from the regions \( \phi_1 < \phi_{A} \) and \( \phi_1 > \phi_{A} \), respectively. When one of this branch points passes \( p = 0 \) then \( \bar{\Psi} \) starts to grow exponentially in time. This happens when \( y \) passes \( \pm iy_{\text{sing}} \) from (2.13) moving along the imaginary axis.

The value of the generating functional is determined in accordance with (2.19) by large time behavior of \( \Psi(t) \). This means that we should be interested in the behavior of \( \bar{\Psi}(p) \) at small \( p \). Really the function \( \int d\phi_1 \hat{\Psi}(p) \) entering (2.19) has a pole at \( p = 0 \) related to the asymptotic behavior

\[
\bar{\Psi}(p) \propto \exp \left( -\frac{2p}{Dd} \sqrt{\frac{3}{d(d^2 - 1)}} \phi_1 \right),
\]

at \( \phi_1 > \phi_{A} \) and small \( p \), the behavior can be found from (2.28). Just the residue of \( \int d\phi_1 \bar{\Psi}(p) \) at the pole determines \( Z(y) \). To find the value of the residue we should analyze the behavior of \( \bar{\Psi}(p) \) at \( 0 < \phi_1 < \phi_{A} \). At small \( p \) there are two contributions into \( \bar{\Psi} \) which behave as

\[
\propto \exp \left\{ \left( \sqrt{\frac{d(d^2 - 1)}{12}} \pm \sqrt{\frac{d(d^2 - 1)}{12} + \frac{y^2P_2}{Dd}} \right) \phi_1 \right\},
\]

as follows from (2.28) at \( p = 0 \). Therefore the residue which is determined by the integral \( \int d\phi_1 \bar{\Psi}(p) \) over the region \( \phi_1 > \phi_{A} \) is proportional to

\[
\exp \left\{ \left( \sqrt{\frac{d(d^2 - 1)}{12}} + \sqrt{\frac{d(d^2 - 1)}{12} + \frac{y^2P_2}{Dd}} \right) \phi_{A} \right\},
\]

(2.31)

Substituting here (2.6) we reproduce (2.12).

Let us now establish the applicability condition of the above procedure. The expression (2.31) implies that the exponent with the sign minus in (2.30) gives a negligible contribution to \( \Psi(p) \) at \( \phi_1 = \phi_{A} \). The condition is satisfied if

\[
|y^2 + y_{\text{sing}}^2| \phi_{A}^2 \gg \frac{Dd}{P_2}.
\]

Substituting here (2.6, 2.13) we obtain

\[
\left| \frac{y \pm iy_{\text{sing}}}{y_{\text{sing}}} \right| \gg \left( d^4 \ln^2 L\Lambda \right)^{-1}.
\]

(2.32)

For \( y \) closer to \( \pm iy_{\text{sing}} \) one should be careful since then a fine analytical structure of \( Z(y) \) will be relevant. As an analysis for \( d = 2 \) shows \( \ref{Z(y)} \) \( Z(y) \) has a system of poles along the imaginary semiaxis starting from \( \pm iy_{\text{sing}} \) and the parameter \( \left( d^4 \ln^2 L\Lambda \right)^{-1} \) just determines the separation between the poles. The assertion about the cut made in the previous subsection is related to the restrictions of the saddle-point approximation which cannot feel this fine pole structure, it gives a picture averaged over the interpole distances, it is precisely the cut. This approach is correct just at the condition (2.32).
Note that the same criterion (2.32) justifies our assumption that the cloud described by \( \Psi \) is narrow during the relevant part of the evolution. Namely, the duration of the part is determined by the time \( t_{\text{exit}} = \frac{1}{p_1} \) (see (2.29)). This is the time which the cloud stays near the barrier. For \( y \) close to \( \pm iy_{\text{sing}} \) the time can be estimated as \( t_{\text{exit}} \sim \frac{P_2}{|y_{\text{sing}}|} |y \mp iy_{\text{sing}}| \). Then the diffusive width \( \sqrt{Dd_{\text{exit}}} \) of \( \Psi \) in the directions \( \phi_a \) for \( a > 1 \) is much less than \( \phi_\Lambda \) just if (2.32) is satisfied. Principally the diffusive dynamics at \( d > 2 \) could modify the noted fine pole structure of \( \mathcal{Z} \), the problem needs a separate investigation.

The same procedure can be done for the passive scalar differences. The cloud \( \Psi \) should pass the region \( \rho_1 < \ln(L/r_0) \) before it reaches the potential. Then it enters the region \( \bar{U} = 2P_2 \) with some finite diffusive width. One can note, however, that it is irrelevant. The only characteristics of potential which are needed are its value (here \( 2P_2 \) instead of \( P_2 \)) and the length of the path inside it (which is \( \Delta \rho_1 = \ln(r_0\Lambda) \) instead of \( \ln(L\Lambda) \)). The evolution of \( \bar{\Psi} \) goes in the same way as in the case of single-point statistics. Again, we get the answer (2.14) and the criterion analogous to (2.32).

In the subsection we presented the analysis based on the dynamical equation (2.20) for the auxiliary object \( \Psi \). The results obtained can be reproduced also on an alternative language: For this we should introduce another auxiliary object the equation for which is stationary. The corresponding scheme which could be interesting from the methodical point of view is sketched in Appendix.

### III.Calculation of PDF

In this section we calculate the PDF’s \( \mathcal{P} \) for the objects (1.8) and (1.10). The most convenient way to do it is in using the relation

\[
\mathcal{P}(\vartheta) = \sqrt{\frac{2}{\pi}} \int \frac{dy}{y} \exp\left( -iy\vartheta \right) \mathcal{Z}(y) ,
\]

where \( \vartheta \) is

\[
\vartheta = \int dr \beta(r) \theta(0, r) .
\]

Let us remind that knowing \( \mathcal{P}(\vartheta) \) one can restore also moments of \( \vartheta \):

\[
\langle |\vartheta|^n \rangle = \int d\vartheta |\vartheta|^n \mathcal{P}(\vartheta) .
\]

The generating functional in (3.1) is determined by (2.12) or (2.14). Being interested in the main exponential dependence of the PDF’s for the objects (1.8) and (1.10) we can forget about preexponents. Then

\[
\mathcal{P}(\vartheta) = \sqrt{\frac{2}{\pi}} \int \frac{dy}{y} \exp\left( -iy\vartheta + q \left[ 1 - \sqrt{1 + y^2/y_{\text{sing}}^2} \right] \right) ,
\]

where

\[
y_{\text{sing}}^2 = \frac{Dd^2(d^2 - 1)}{12P_2} \quad \text{(for scalar)}, \quad n_{\text{sing}}^2 = \frac{Dd^2(d^2 - 1)}{24P_2} \quad \text{(for difference)},
\]

\[
q = \frac{d(d + 1)}{6} \ln(L\Lambda) \quad \text{(for scalar)}, \quad q = \frac{d(d + 1)}{6} \ln(r_0\Lambda) \quad \text{(for difference)}.
\]

Since both \( q \) defined by (3.3) are regarded to be much larger than unity the integral (3.4) can be calculated in the saddle-point approximation. The saddle-point value is
\[ y_{sp} = i \frac{y_{\text{sing}}}{1 + q^2/y_{\text{sing}}^2 q^2} . \]  

(3.7)

Then

\[ \ln \mathcal{P}(\vartheta) \simeq q \left( 1 - \sqrt{1 + \frac{y_{\text{sing}}^2 q^2}{q^2}} \right) . \]  

(3.8)

The expression leads to the exponential tail

\[ \ln \mathcal{P}(\vartheta) \simeq -y_{\text{sing}} |\vartheta| , \]  

(3.9)

realized at $|\vartheta| \gg q/y_{\text{sing}}$. The coefficient $y_{\text{sing}}$ in (3.9) determined by (2.13) is in accordance with the result obtained in [12].

The expression (3.8) enables one to find the following averages in accordance with (3.3)

\[ \langle \theta^2 \rangle = \frac{2P_2}{d(d-1)D} \ln(L\Lambda) , \quad \langle (\Delta \theta)^2 \rangle = \frac{4P_2}{d(d-1)D} \ln(r_0\Lambda) . \]  

(3.10)

The expressions (3.10) can be obtained also by direct expansion of $\mathcal{Z}(y)$ from (2.12) or (2.14). The universal tail (3.9) is realized if

\[ \theta_{\Lambda} \gg \sqrt{\langle \theta^2 \rangle} d \ln(L\Lambda) , \quad \Delta \theta_{\Lambda} \gg \sqrt{\langle (\Delta \theta)^2 \rangle} d \ln(r_0\Lambda) . \]  

(3.11)

Since both logarithms are supposed to be large we conclude that there exists a relatively wide region where the statistics of $\vartheta$ is approximately Gaussian, the region is determined by the inequalities inverse to (3.11).

Let us discuss the applicability conditions of the expression (3.8). First, if one calculates the passive scalar PDF by the saddle point method then the position of the saddle point is determined by (2.32) if

\[ \vartheta \ll d^2 \sqrt{\frac{P_2}{D}} \ln^2(L\Lambda) . \]  

(3.12)

The applicability domain of the saddle-point method overlaps the region of validity of expression (2.12) for the generation function $\mathcal{Z}(y)$. The above inequalities are correct for $\theta_{\Lambda}$, for $\Delta \theta_{\Lambda}$ one should substitute $\ln(L\Lambda)$ by $\ln(r_0\Lambda)$. Second, fluctuations of $y$ have to be small compared to the distance between $y_{sp}$ and $y_{\text{sing}}$. This gives the same criterion (3.12).

Let us stress that though formally our procedure is incorrect at $\vartheta \gtrsim d^2 \sqrt{P_2/D} \ln^2(L\Lambda)$ the answer will be the same: the PDF will be determined by the exponential tail (3.9). The point is that the character of the integral (3.1) at such extremely large $\vartheta$ will be determined by the position of the singular point of $\mathcal{Z}(y)$ nearest to the real axis. This is just $iy_{\text{sing}}$ leading to (3.9). To conclude, only the character of the preexponent in $\mathcal{P}(\vartheta)$ is changed at $\vartheta \sim d^2 \sqrt{P_2/D} \ln^2(L\Lambda)$ whereas the principal exponential behavior of $\mathcal{P}(\vartheta)$ remains unchanged there.

CONCLUSION

The single-point statistics of the passive scalar $\theta$ and the statistics of its difference $\Delta \theta$ are traditional objects which carry an essential information about correlation functions of the passive scalar
in the convective interval. We examined the passive scalar in the large-scale turbulent flow where the correlation functions logarithmically depend on scale. Since really the logarithms are not very large it is useful to have the whole PDF’s of $\theta$ and $\Delta \theta$. That was the main purpose of our investigation which was performed in the frame of the Kraichnan model. The single-point PDF for the passive scalar and the PDF for the passive scalar differences can be obtained from (3.8) if to substitute $\Lambda \rightarrow r_{\text{dif}}^{-1}$ where $r_{\text{dif}}$ is the diffusive length. Though both the advecting velocity and the pumping force in the Kraichnan model are considered as $\delta$-correlated in time we hope that our answer are universal that is are true in the limit when the size of the convective interval tends to infinity for arbitrary temporal behavior of the velocity and of the pumping. The reason is that the spectral transfer time grows with increasing the convective interval and in the limit is much larger than the correlation times of the velocity and of the pumping.

We believe also that the analytical scheme proposed in our work could be extended for other problems related to the passive scalar statistics. Note as an example the work [13] where a modification of the scheme enabled to find the statistics of the passive scalar dissipation. It is also useful for investigating the large-scale statistics (on scales larger than the pumping length) of the passive scalar [26]. We hope also that it is possible to go beyond the case of the large-scale velocity field using the perturbation technique of the type proposed in [27, 29].

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APPENDIX A:

Here we present an alternative way to obtain the answers (2.12,2.14). We will use an auxiliary quantity

$$\Xi(y, \rho_0, \eta_0) = \int D\rho D\eta Dm D\mu \exp \left[ \int_{-\infty}^{0} dt \left( i\mathcal{L} - \frac{y^2}{2} U \right) \right]_{t=0}=\rho_0, \eta(0)=\eta_0, \quad (A1)$$

then

$$Z(y) = \Xi(y, 0, 0). \quad (A2)$$

Function $\Xi$ can be also defined as

$$\Xi(y, \rho_0, \eta_0) = \lim_{t \to \infty} \int \prod d\rho_i d\eta_i \Psi(t, y, \rho, \eta), \quad (A3)$$

where $\Psi$ is governed by equation (2.20) with initial condition $\Psi(t=0, y, \rho, \eta) = \delta(\rho - \rho_0)\delta(\eta - \eta_0)$. The equation for the function $\Xi$ can be found from (1.32[A1]):

$$\sum_{i=1}^{d} \frac{\partial^2}{\partial \rho_i^2} - \frac{1}{d} \left( \sum_{i=1}^{d} \frac{\partial}{\partial \rho_i} \right)^2 + \sum_{i=1}^{d} (d - 2i + 1) \frac{\partial}{\partial \rho_i}$$

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\[ +2 \sum_{i<j} \exp(2\rho_j - 2\rho_i) \frac{\partial^2}{\partial \eta_{ij}^2} + 4 \sum_{i<k<j} \exp(2\rho_k - 2\rho_i)\eta_{k}\eta_{j} \frac{\partial}{\partial \eta_{ij}} \frac{\partial}{\partial \eta_{ik}} \] 

\[ +2 \sum_{i<k<m,n} \exp(2\rho_k - 2\rho_i)\eta_{km}\eta_{kn} \frac{\partial}{\partial \eta_{im}} \frac{\partial}{\partial \eta_{in}} \] 

\[ \Xi - \frac{y^2U}{Dd} \Xi = 0. \quad (A4) \]

The boundary condition for the equation (A4) follows from the definition (A1): For large enough \( \rho_i, \eta \) the potential \( U = 0 \) for \( t = 0 \) and remains zero also at finite times \( t \). Therefore the integral (A1) should be equal to unity in the case. Thus \( \Xi(y, \rho, \eta) \) should tend to unity where \( \rho, \eta \to \infty \).

Let us rewrite the equation (A4) in terms of the variables (2.2):

\[ (\hat{\Gamma}_1 + \hat{\gamma})(\Xi_1 + \xi) = 0, \quad \Xi = \Xi_1 + \xi, \quad (A5) \]

\[ \hat{\Gamma}_1 = \frac{\partial^2}{\partial \phi_1^2} + \sqrt{\frac{d(d^2 - 1)}{3}} \frac{\partial}{\partial \phi_1} - \frac{y^2U}{Dd}, \quad (A6) \]

\[ \hat{\gamma} = \sum_{i=2}^{d-1} \frac{\partial^2}{\partial \phi_i^2} + 2 \sum_{i<k} \exp(2\rho_k - 2\rho_i) \frac{\partial^2}{\partial \eta_{ik}^2} + 4 \sum_{i<k<n} \exp(2\rho_k - 2\rho_i)\eta_{kn} \frac{\partial}{\partial \eta_{in}} \frac{\partial}{\partial \eta_{in}}. \quad (A7) \]

Here \( U \) as a function of \( \phi_1 \) is equal to \( P_2 \) inside a region restricted by \( \phi^-_A \) and \( \phi^+_A \) (where \( \phi^\pm_A \) are functions of variables \( \phi_2, \ldots, \phi_d, \eta \)) and tends to zero outside the region. We will solve the equation (A5) using the perturbation theory over \( \hat{\gamma}, \xi \). Then the zero order equation is

\[ \hat{\Gamma}_1 \Xi_1 = 0. \quad (A8) \]

The equation (A8) can be easily solved at \( \phi^-_A < \phi_1 < \phi^+_A \), the answer is

\[ \Xi_1 \simeq \frac{2\lambda}{\sqrt{\lambda^2 + \frac{y^2P_2}{Dd} + \lambda}} \exp \left( - \left( \sqrt{\lambda^2 + \frac{y^2P_2}{Dd}} - \lambda \right) \left( \phi^+_A - \phi_1 \right) \right), \quad (A9) \]

where \( \lambda = \sqrt{d(d^2 - 1)/12}, Dd\lambda \) is the rate of the cloud motion along \( \phi_1 \) direction. The result (A9) can be obtained using the inequality \( \sqrt{\lambda^2 + \frac{y^2P_2}{Dd} + \lambda} \gg 1 \). The derivative \( \partial \Xi_1 / \partial \phi_1 = 0 \) at \( \phi_1 < \phi^-_A \). However, \( \Xi_1 \neq 1 \) in this region. This is due to the following fact: this region corresponds to the evolution of \( \Psi \) when its initial position is to the left of potential \( U \) (see (A3)). During the evolution cloud \( \Psi \) passes the region of \( U \) and its integral over \( \rho, \eta \) changes. Then \( \Xi \) is not equal to 1. Only when the distance between initial position and potential is of order \( \ln^2 LA \) the diffusion of the cloud leads to smallness of the part of \( \Psi \) which passes the potential \( U \), and \( \Xi \) becomes closer to unity. Thus, function \( \Xi \) has long tail from potential with direction to negative \( \phi_1 \) where it is not equal to 1. The procedure of finding \( \Xi \) from the equation (A8) corresponds to the geometrical optics approximation (taking into account only derivatives in propagating direction; this allows one to get the fact of propagation). This tail of \( \Xi \) in this approximation is nothing else than the shadow of potential \( U \). Higher orders of perturbation theory over the transverse derivatives correspond to diffraction corrections.

Now let us consider the correction \( \xi \). The equation for it looks like \( (\hat{\Gamma}_1 + \hat{\gamma})\xi = -\hat{\gamma}\Xi_1 \). Again let us neglect \( \hat{\gamma} \) in l.h.s. and solve the equation. \( \Xi_1 \) is some exponential function with scale of the order 1. Then \( \hat{\gamma}\Xi_1 \sim \Xi_1 \). Not that \( \hat{\gamma}\Xi_1 \) is almost equal to zero at \( \phi_1 > \phi^+_A \). To estimate \( \xi \) one should construct Green function \( G(\phi_1|\phi_0) \) for operator \( \hat{\Gamma}_1 \):
\[ G(0|\phi_0) \simeq \frac{1}{2\lambda} \exp \left( -\left( \sqrt{\lambda^2 + \frac{y^2P_2}{Dd} - \lambda} \right) \phi_0 \right) \left( 1 - C \exp \left( -2\sqrt{\lambda^2 + \frac{y^2P_2}{Dd} (\phi_0^+ - \phi_0)} \right) \right), \quad (A10) \]

where

\[ C = \left( \sqrt{\lambda^2 + y^2P_2/Dd - \lambda} \right) / \left( \sqrt{\lambda^2 + y^2P_2/Dd + \lambda} \right). \]

The unity in brackets in (A10) gives the correction for \( \Xi \) which has the same exponential factor as \( \Xi_1 \). Thus \( \xi \) does not change the answer with the logarithmic accuracy. The second term in the brackets gains while \( \phi_0 \) is close to \( \phi_0^+ \). It is due to nonzero width of the cloud \( \Xi \) and to dependence of \( t_k \) on other variables. Again, it does not change the exponent.

The case of the passive scalar differences can be considered in a similar way.
[28] A. Pumir, B. I. Shraiman, and E. D. Siggia, preprint cond-mat/9609272.
[29] E. Balkovsky, G. Falkovich, and V. Lebedev, Phys. Rev. E, 55, 4881, (1997).