From Euclidean to Minkowski space with the Cauchy-Riemann equations

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Abstract. We present an elementary method to obtain Green’s functions in non-perturbative quantum field theory in Minkowski space from calculated Green’s functions in Euclidean space. Since in non-perturbative field theory the analytical structure of amplitudes is many times unknown, especially in the presence of confined fields, dispersive representations suffer from systematic uncertainties. Therefore we suggest to use the Cauchy-Riemann equations, that perform the analytical continuation without assuming global information on the function in the entire complex plane, only in the region through which the equations are solved. We use as example the quark propagator in Landau gauge Quantum Chromodynamics, that is known from lattice and Dyson-Schwinger studies in Euclidean space. The drawback of the method is the instability of the Cauchy-Riemann equations to high-frequency noise, that makes difficult to achieve good accuracy. We also point out a few curiosities related to the Wick rotation.

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1 The Wick rotation

Central to non-perturbative quantum field theory is the computation of Green’s functions, the vacuum expectation value of quantum operators. These and the related Scattering Matrix elements are most often computed in Euclidean space, defined by the transformation

\[ t \rightarrow -it_E \quad k_0 \rightarrow ik_0 \, . \]

This coordinate transformation is known as “Wick rotation” (see for example [1] for a short account). There are many advantages in solving the field equations in terms of the rotated variables to obtain so called Schwinger functions, and we list some below in section 2.

Once the wanted functions have been computed in Euclidean (momentum) space \( k_E = (ik_0, \mathbf{k}) \) one would wish to recover the original Minkowski space Green’s functions by inverting the Wick rotation. This is possible in perturbation theory at low orders \([2,3]\) where one has explicit expressions for the functions, and their analytical structure (poles, cuts, essential singularities) is at hand, so that one can employ Cauchy’s theorem and collect dispersive cut integrals or pole residues if need be, and obtain the Minkowski space Green’s function by analytical continuation.

For example, to obtain the electron propagator in momentum space (Fourier transform of the probability amplitude for the electron to reach point \( x \) if it was originally at the origin 0)

\[ S(p) = \int d^4 x e^{-i p x} S(x, 0) = i Z(p^2) \frac{\not{p} + M(p^2)}{p^2 - M^2(p^2) + i\epsilon} \]

one would perform the Wick rotation \( p_0 \rightarrow i p_0 E \) and obtain a function \( S(p_E) \) as a perturbation of its free-field values \( (M(p^2) = m, Z(p^2) = 1) \). With the function explicitly known, one just extends it into the complex plane and simply substitute its argument \( p_0 E \) by \( -i p_0 \). If the function is not totally known but its analytical structure is, one employs Cauchy’s theorem as mentioned.

However, in non-perturbative quantum field theory one is seldom in this desirable situation. More often than not, the function has been calculated with the help of a computer, be it by solving the Dyson-Schwinger equation \([4]\) with some carefully designed truncation, or by trying a Monte Carlo evaluation averaging over a small number of configurations on a lattice \([5]\).

The outcome is that the function is then known for Euclidean momenta, typically \( k_E^2 > 0 \), and an extension into the complex plane becomes necessary to reach the negative axis \( k_E^2 < 0 \), \( k_E^2 = -k^2 \). For non-perturbative functions one sometimes ignores the precise analytic structure in the complex plane. This situation is worsened in theories where the field quanta do not appear in asymptotic states, except in very specific combinations, such as is presumably the case for Quantum Chromodynamics.

Attempts have of course been made to solve the problem. An obvious approach is to perform a theory-motivated fit to the computer data, building-in well educated guesses...
on what the analytical structure of the continued function must be.\[^{[47]}\]

Another possibility is to write a spectral representation in terms of a Stieltjes transform yielding a spectral density $\rho$,

$$ S(p^2) = \int dp \frac{\rho(m^2)}{p^2 - m^2 + i\epsilon} \tag{3} $$

(we have eliminated spin structure) and then solving the Dyson-Schwinger equation directly “in Minkowski space” for $\rho$, finally inverting the Stieltjes transform to recover the propagator or Bethe-Salpeter wanted function \[^{[5], [9]}\]. It is of course clear that in writing the spectral representation, one is already assuming a given cut structure for the function. It is known that, for local quantum field theories (whose absence has been at times thought as a possible sign of confinement).

In this paper we enrich the toolbox by putting forward a very simple method that does not require the function’s analytical structure to be known on the entire complex plane. We observe that the analyticity of a given function cannot only be formulated globally, through satisfaction of Cauchy’s theorem, but also locally, through satisfaction of the Cauchy-Riemann equations. By integrating this simple first order differential system with initial condition the function in a given region computed by other means (DSE, lattice, exact renormalization group equations \[^{[10], [11]}\], etc.) one can achieve two goals. First, the numerically obtained solution can break down at a given point or line in the plane, indicating perhaps a pole or other singularity. Second, if the region where the system is integrated avoids such singularities, one can obtain the analytical extension (within errors and non-uniqueness) to another region in the complex plane.

We use as example the behavior of the quark mass function in Landau gauge QCD, $M(p^2)$, that we analytically continue from positive Euclidean virtuality ($p^2 > 0$) into Minkowski space with ($p^2 < 0$). This is plotted in figure 1. It can be seen that, within the statistical errors inherited from lattice data and the systematic numerical errors intrinsic to our procedure, the mass function decreases with increasing Minkowski $p^2$. Our result keeps open the possibility of a pole of the quark propagator for real $p^2$ (whose absence has been at times thought as a possible sign of confinement).

The rest of the paper consists of four sections. In section 2 we make a few comments, some common place but others quoted less often, about the advantages of initially working in Euclidean space. In section 3 we present the Cauchy-Riemann method with one practical case, the quark propagator. A few theoretical comments about errors involved in the process and the generalization to more dimensions are left for section 4. Our discussion is summarized in section 5.

\[\text{2 Working in Euclidean space}\]

In lattice formulations of Quantum Field Theory, the field configurations over which the path integral is evaluated are randomly generated according to a distribution $e^{-\int dt L}$, that is the Wick-rotation of the actual quantum weight for the path integral, $e^{\int dt L}$. Thereafter computed lattice Green’s functions are valid in Euclidean space. Even if working in Minkowski space, a popular way of “minimally” regularizing in the path integral formalism is to rotate the time integration into the complex plane

$$ \int d^4x = \lim_{T \rightarrow \infty} \int_{-T}^{T} dt \int d^3x. $$

Beyond the convergence of the path integral and the weighting configurations inside a compact set of function space, there are several more advantages.

One is that the Dyson-Schwinger equations, whose solutions are often used to interpret lattice data, are extremely difficult to solve on a computer in Minkowski space. Indeed, a typical DSE is that for the mass function of a fermion in the presence of a scalar field, with Yukawa coupling, in rainbow approximation

$$ M(p) = c \int d^4k \frac{M(k)}{(k^2 - M(k)^2 + i\epsilon)((k - p)^2 - m^2 + i\epsilon)}. \tag{4} $$
If one solves the equation iteratively by guessing \( M_0 \), one notices that the position of the fermion pole in the denominator is not known after the first iteration, and it needs to be determined numerically (an attempt at carrying on this program exists [12]).

However the common use is to Wick-rotate the \( k \) integration variable to Euclidean space. If \( p \) is likewise rotated, the resulting equation is easier to program as the denominator poles are on the left \( k^2 \) plane, out of the numerical integration region in the radial \( k^2 \) variable.

\[
M(p) = c \int d^4k \frac{M(k)}{(k^2 + M(k)^2 + i\epsilon)((k^2 - m_c^2)^2 + i\epsilon)} \tag{5}
\]

\( c \) represents constants irrelevant to the discussion.

Note also that when working in Euclidean space, a discrete subgroup of the Euclidean rotation group is retained. For example, for a simple cubic lattice with \((x = ai_x, y = ai_y, z = ai_z, t = ai_t)\), invariance under rotations by \( \pi/2 \) is explicit. This can be exploited to study the quantum representations of the discrete group, then trying to match the resulting states to a representation of the full continuous group.

However, in Minkowski space there is no finite lattice under a non-trivial subgroup of the Lorentz group. If a grid is invariant under discrete Lorentz transformations of parameter \( a \), then it has infinitely many points. We discuss Lorentz-invariant discretizations of Minkowski space in the appendix.

One further motivation is the non-compactness of the equal \( k^2 \) hypersurfaces. While in Euclidean space the condition \( k^2 = \Lambda^2 \) determines a hypersphere’s surface, so that

\[
\int_\Lambda^\Lambda d^4k_E f(k_E^2) = 2\pi^2 \int_\Lambda^\Lambda k_E^2 dk_E f(k_E^2) \tag{6}
\]

can be factorized into a radial integral and a finite \( 2\pi^2 \) solid hyperangle (the hyperarea of a unit-radius hypersphere’s surface) this is not possible in Minkowski space, where the corresponding unit-hyperboloid \( k^2 = 1 \) has infinite hypersurface. Therefore, integrals of Lorentz-invariant functions are by necessity divergent even after regulation of large virtualities, and are only defined by analytical continuation from Euclidean space.

In perturbation theory, a much used method is to perform the \( k_0 \) integrals first, usually with pole analysis, and later impose a cutoff on spacelike momentum \( k \). However this cutoff is frame-dependent, and there is no direct method that manifestly preserves Lorentz invariance.

From all these arguments, it is hard to conceive progress in non-perturbative quantum field-theory in Minkowski space without progress in complex-plane analytical continuation for the relevant functions. This paper is a modest contribution in this direction, with the interest in keeping the discussion alive.

3 Numerical solution of the Cauchy-Riemann equations

3.1 Cauchy-Riemann equations in polar coordinates

If \( u \) and \( v \) are respectively the real and imaginary parts of a complex function of one complex variable \( p^2 = r e^{i\theta} \), the Cauchy-Riemann equations in polar coordinates read

\[
\frac{\partial v}{\partial \theta} = r \frac{\partial u}{\partial r} \tag{7}
\]

\[
\frac{\partial u}{\partial \theta} = -r \frac{\partial v}{\partial r}. \tag{8}
\]

Given the initial conditions \( u(r, 0) = u_0(r) \), \( v(r, 0) = v_0(r) \) on a segment of the real \( p^2 \) axis, corresponding to \( \theta = 0 \), one can then evolve the system towards increasing and decreasing \( \theta \) (like the opening of a fan). For very smooth data sets one can typically reach 90-120 degrees on each side of the fan before the instabilities wipe the solution to infinity. The (Cauchy-Euler) explicit discretization with centered \( r \)-derivative on a grid \((r_j, \theta_i)\) is simply

\[
v(r_j, \theta_{i+1}) = v(r_j, \theta_i) + r_j (\theta_{i+1} - \theta_i) \frac{u(r_{j+1}, \theta_i) - u(r_{j-1}, \theta_i)}{r_{j+1} - r_{j-1}} \tag{9}
\]

\[
u(r_j, \theta_{i+1}) = u(r_j, \theta_i) - r_j (\theta_{i+1} - \theta_i) \frac{v(r_{j+1}, \theta_i) - v(r_{j-1}, \theta_i)}{r_{j+1} - r_{j-1}} \tag{10}
\]

where, to solve over an arch taken anticlockwise, \((\theta_{i+1} - \theta_i) > 0 \). At the end-points of the grid one cannot use centered derivative, so left (right derivative) is necessary, \((v(r_2, \theta_i) - v(r_1, \theta_i))/(r_2 - r_1)\) etc. The situation is represented in figure 2.

![Fig. 2. The Cauchy-Riemann equations in polar coordinates allow to explore a fan-shaped region of the complex plane where a function is analytic.](image-url)
equations for mesons, since the (external) meson momentum is of course in Minkowski space (real), and the internal quark momentum is Wick-rotated to Euclidean space (imaginary) so that one ends solving the DSE inside a parabola in the complex plane symmetric respect to the real momentum axis. The Cauchy-Riemann equations are currently no match in precision to directly solving the DSE in the complex plane where this is feasible, but they can provide a cross check that is very simple to programme (compare the trivial linear system above with the complex, non-linear, bidimensional DSE when the angular kernel or vertex are non-trivial).

The Cauchy-Riemann equations however are a statement of analyticity, and the solution is a numerical representation of the closest analytical function that contains the initial data. This means that if the “true” function has a pole or a cut, the Cauchy-Riemann iteration will fail to see it and simply separate from that function, and is likely to diverge soon from accruing instabilities. This is illustrated in figure 4.

3.2 Cauchy-Riemann equations on a strip

The advantage of a local formulation of analyticity employing the Cauchy-Riemann equations is lost if one needs to swipe the entire complex plane. Therefore it is profitable to solve them in Cartesian coordinates, first away from the $x$ axis along $y$

$$\frac{\partial u}{\partial y} = -\frac{\partial v}{\partial x}$$

and then leftwards along $x$. The method fails if the complex plane is completely cut from $-\infty$ to $\infty$ along the $y$ axis. In any other situation (the standard half-plane cut of a power-law or logarithm, or a finite number of poles or essential singularities), one can find a path between the right and the left $x$ axis and solve the Cauchy-Riemann equations along them.

We now improve upon the discretization of the differential equations and employ an implicit $\theta$ method. This is convenient since the Cauchy-Riemann are quite unstable (as pointed out below in subsection 4.2). The $\partial_y u$ equation for a point not on the edge of the grid becomes (with the superindex labelling $y$, the subindex $x$)

$$\frac{u_{i+1}^j - u_i^j}{y_{j+1} - y_j} = -\frac{1}{x_{i+1} - x_i} \times$$

$$(\theta(i_{i+1}^j - i_{i-1}^j) + (1 - \theta)(i_{i+1}^j - i_{i-1}^j))$$

In the advance along $y$ one groups the $u_i$ and $v_i$ for fixed $y_j$ in a vector $\mathbf{u}^j = (u_1^j, u_2^j, u_3^j, \ldots, u_N^j)$ and the $\theta$ method’s discretization can be written down as a linear problem

$$A\mathbf{u}^{j+1} = B\mathbf{u}^j$$

In the simplest case of equal $x$ subintervals, one can define $r = \theta \Delta y/(2 \Delta x)$, then the matrix $A$ becomes

$$
\begin{pmatrix}
1 & -2r & 0 & 2r & 0 & \ldots \\
2r & 1 & -2r & 0 & \ldots \\
0 & -r & 1 & 0 & 0 & r & \ldots \\
r & 0 & 0 & 1 & -r & 0 & \ldots \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
\ldots & 0 & -2r & 1 & 2r \\
\end{pmatrix}
$$

where the third and fourth line are the repeated unit, except the non-vanishing elements are shifted to the right (the 1’s always mark the diagonal, and the matrix is band-
diagonal). The matrix $B$ can likewise be filled by exchanging $\theta \rightarrow (1 - \theta)$ and changing the sign of all off-diagonal matrix elements of $A$ (the diagonal of $B$ likewise contains 1’s).

The advance in $y$ proceeds by solving the linear system to obtain $(u, v)$ at $y_{j+1}$ from their values at $y_j$. We use standard $LU$ factorization, although since the matrix $A$ is band-diagonal, Crout’s algorithm can speed things some. For one complex variable and the small number of points we use this is irrelevant in computer time.

However, since the advance to the left will have as initial condition the edge (along the $x$ axis) of the first computed block, and at the end we will be interested on the values of the function on the $x$ axis, that is the edge of the second computed block, it pays off to improve the computation of the derivative at the edge of the block.

The left-derivative we have displayed explicitly,

$$\frac{\partial u(x = 0, y)}{\partial x} = \frac{u^j_2 - u^i_1}{x_2 - x_1} + o(h)$$

(16)

can be interpreted as a centered derivative at the midpoint $(x_2 + x_1)/2$. Considering also the centered derivative at $x_2$ and extrapolating linearly to $x_1$, one obtains an improved

$$\frac{\partial u(x = 0, y)}{\partial x} = 2 \frac{u^j_2 - u^i_1}{x_2 - x_1} - \frac{u^i_3 - u^j_1}{x_3 - x_1} + o(h^2),$$

(17)

and likewise at the last $N$ point of the grid, and this slightly complicates the first two and last two rows of $A$.

Once the advance upwards in the $y$ direction has reached $N$, one starts an advance to the left as in figure 4 and similar considerations apply.

The $\theta$ parameter that advances or delays the derivative perpendicular to the integration direction is empirically fixed for now. Several problems are somewhat independent of $\theta$, others have a broad minimum of instability around -1. We find that $\theta = -0.5$ is as good as any. A brief eigenvalue analysis is presented below that explains why, in subsection 4.2.

To show a test of the method, we employ the simple function

$$f_{\text{test}}(z) = \frac{1}{1 + z^2/4}$$

(18)

that has two poles above and below the $x$ axis at $\pm 2i$. Because of the fast build-up of numerical errors, we need $\Delta y < \Delta x$ so the strip is always shorter in the direction of the advance of the integration (for an $N \times N$ problem). In figs. 5 and 6 we show the real and imaginary parts calculated with the Cauchy-Riemann equations with initial condition on the positive real half-axis, plotted along the imaginary and the negative real-axis respectively.

![Fig. 4. The Cauchy-Riemann equations in Cartesian coordinates can be used to explore a strip, first upwards along the y axis, then leftwards along the x axis. In this use one only needs analyticity of the function on a strip above (or below) the axis to obtain information on the Minkowski side. The method fails only if the complex plane is completely cut from $-\infty$ to $\infty$ along the y axis.](image)

![Method check](image)

**Method check**

The method starts breaking down for $x$ near zero and $y > 1$ as the pole is approached. But one can see how the part of the strip that goes well below the pole passes cleanly and allows to reproduce the function on the left axis, given as initial condition the $N$ exact values on the right axis. All in all, the method provides a reasonable representation of the function from the solution of the Cauchy-Riemann equations. Note the imaginary part, exactly zero on the left axis, is calculated to be of order 1% with a forty-point grid. This should be considered the error of the method, and is far less than the statistical errors in the lattice data that we will shortly employ.

Should the analytic structure of the function become available, one could devise an arbitrary path in the complex plane from the region where the function is known to the region where it is wanted by analytical continuation. The initial value problem can then be formulated with the advance direction along the tangent vector to the path, $\tau$. This vector changes in principle orientation, so one would need to use the “Cartesian-like” formulation to advance...
3.3 Wick rotation of the quark mass function

In the introduction we have advanced, in figure 1, our main application, the analytical continuation of the quark mass function \( M(p^2) \) in Landau gauge, that we now carefully analyze and justify.

Because of Lorentz invariance applied to a spin 1/2 fermion, the quark propagator can be written in full generality as

\[
S(p) = \frac{iZ(p^2) \, \not{p}}{p^2 - M^2(p^2)} + \frac{iZ(p^2) \, M(p^2)}{p^2 - M^2(p^2)}.
\]  

As already discussed, the pole in the denominator is a nuisance usually disposed of by performing the Wick rotation \( p_0 \to i\eta_F^2 \). Concentrating on the denominator then

\[
\frac{1}{p_0^2 - \not{p}^2 - M^2(p_0^2 - \not{p}^2)} \to \frac{-1}{p_0^2 + \not{p}^2 + M^2(-p_0^2 - \not{p}^2)}
\]

the pole is absent for real \( M \). One usually eschews a sign, the function \( M^2(-p_0^2 - \not{p}^2) \) is customarily called \( M(p_0^2) \) where \( p_0^2 > 0 \), and we will keep this notation. Hence to retrieve the mass function that actually appears in the Minkowski space propagator for positive virtuality \( p^2 \), we need to identify it with the analytical continuation of the lattice (or DSE) \( M(p_0^2) \) to negative \( p^2 \). Note that although the analytical continuation is nominally made in \( p_0 \), since \( p_0^2 - \not{p}^2 \) is a polynomial, it is an analytical function of \( p_0 \). We conclude that wherever \( M \) is analytical in \( p_0 \) it is also analytical in \( p_0^2 - \not{p}^2 \) (conversely, there needs to be a branch cut in the \( p_0 \) plane that is absent in terms of the Lorentz invariant variable).

The analytical properties of \( M(p^2) \) are not well known, especially for confined quanta such as quarks. However we note that a pole in this function would imply a zero in the quark propagator, and this, assumed continuous, can be ruled out in the entire region of the complex plane that we sample, from the lattice data on the Euclidean real \( p^2 \) axis.

We take the lattice data from [11]. This has been provided to us from \( 28^3 \times 96 \) lattices with MILC configurations. We have used \( a^{-1} = 2.29 \text{ GeV} \) to set the scale from the internodal spacing.

Due to our sensitivity to large frequency noise, we only use a subset of the lattice data, taking one of every few points in the interval \( p \in (0, 8) \text{ GeV} \). The trimming has been performed so that the resulting mass monotonously decrease towards higher momenta (asymptotic freedom), to avoid distortions of analyticity and large errors through rapidly varying derivative. We further square the abscissa \( p - p^2 \), since the latter is the variable in which we perform the analytical continuation. The original lattice data,
A part of the lattice data for the quark mass function in Landau-gauge QCD. We have trimmed the data to ensure the monotonous decrease in the function and reduce the high-frequency noise, which grows fast in the Cauchy-Riemann equations. We slightly increased the error bands to cover the omitted points (trading our systematic error in trimming into a statistical error). The lattice data from [11] is normalized to a current quark mass $m_u = 60 \, MeV$ already in physical $GeV$ units, is given in figure [9].

We also show the actual input set to our code after these manipulations have been performed, that is faithful to the original data and error bands, but amenable to analytical continuation. The data is renormalized in the $MOM$ scheme, and as can be seen the mass at a scale of $8 \, GeV$ is $60 \, MeV$ with the chosen scale $a^{-1}$. We further take as input a Dyson-Schwinger calculation from [16]. This data has as an advantage that there are no statistical errors and the function is very smooth. In exchange, there are systematic errors (coming from the precise way in which the quark-gluon vertex is treated in that reference), that are unknown and only controllable in the propagator in comparison with lattice data or renormalization group equations. We plot the resulting set in figure [10].

Finally, we perform the analytical continuation on a strip in the complex plane above the axis, that presumably avoids non-analyticities in $M(p^2)$ (else a continuation under the axis is possible) and obtain the real part of $M$ advanced in the introduction in figure [1].

We also plot in figure [11] the imaginary part of the mass-function on the left half-axis in Euclidean space that is well-compatible with 0. This is in agreement with the outcome of a Taylor expansion around the origin (however there is no telling the convergence radius of such a series, so we deem the Cauchy-Riemann method superior).

The Cauchy-Riemann method leads to an imaginary part of the mass-function on the left half-axis in Euclidean space that is well-compatible with 0. This is in agreement with the outcome of a Taylor expansion around the origin (however there is no telling the convergence radius of such a series, so we deem the Cauchy-Riemann method superior).

From the graphs one can conclude that, just as for the tree-level propagator in perturbation theory, there is a crossing of $M$ and $p$ for (negative Euclidean), positive Minkowski $p^2$. This means that the actual quark propagator does have a pole at or very near the real axis. It has been quoted [16] at $(300 - 500) \, MeV$ from the Dyson-Schwinger equations alone. From the analysis of the lattice data set at hand, we conclude $M(M) = 305(25) \, MeV$, in agreement with that estimate. Of course, it would be interesting to compare with other lattice data sets, and in particular use different current quark masses, so the error band is definitely larger.

We do not find support for the attending conjecture of two conjugate poles with a sizeable imaginary part.
4 Some theoretical issues

4.1 Uniqueness

Here we study to what extent the solution of the Cauchy-Riemann equations for the quark mass function is unique, given the initial conditions as the lattice computation on the positive $p^2$ half-axis. By standard complex analysis, the uniqueness of an analytic continuation of a function is guaranteed if the function is initially known on an open subset of $C$.

The positive real half-axis is open in $R$, but not in $C$. However it is easy to show that the analytic continuation is unique. Imagine that $u$ and $v$ are known for $y = 0$ and $x > 0$. Then, all partial derivatives $\frac{\partial u}{\partial y}$ and $\frac{\partial v}{\partial x}$ are known. In particular for the quark mass function, $v(y = 0) = 0$ (and all $x$-derivatives also vanish), and $u(x, y = 0) = M(x)$, the real mass.

Assume that the extension of $M$ to the complex plane was not unique. Then in addition to $f = (u, v)$ there would be another function, $f + g$, that would satisfy the Cauchy-Riemann equations with the same initial conditions. Since the sum of two analytic functions is analytic, $g$ itself should be analytic. This means that its components $g_x$, $g_y$ would also satisfy the Cauchy-Riemann equations

\begin{align}
\frac{\partial g_x}{\partial y} &= -\frac{g_y}{\partial x} \quad (20) \\
\frac{\partial g_y}{\partial y} &= \frac{g_x}{\partial x} \quad (21)
\end{align}

with initial condition $g(y = 0) = 0$ exactly, with all derivatives $\frac{\partial^2 g(y = 0)}{\partial y^2} = 0$ also vanishing on the real axis. Automatically, employing eq. (20), and subsequently deriving it, all $y$ derivatives also vanish. Therefore $g$ is exactly zero in the domain of analyticity, and $f$ unique.

Of course, in practice $f$ is only known at a discrete and finite set of points $z_i = (x_i, y_i)$ $i = 1...N$. An analytic function could oscillate between any two of the points and take arbitrarily large or small values. Therefore one needs an additional hypothesis to claim that the computed function is a fair representation of the “actual” function.

The sufficient hypothesis is monotony of the function between any two grid points (note the function might be globally non-monotous but being allowed to change derivative sign at the grid points themselves). If the function is strictly decreasing between $x_i$ and $x_{i+1}$, then the maximum and minimum values that it can take between them are $f_i$ and $f_{i+1}$, and the function is bound (it being analytic, it is also continuous). Then, to arbitrarily shrink the error in our knowledge of the initial condition, one just needs to arbitrarily shrink the grid spacing, so as to further constrain the function in every subinterval. The function computed with the discretized Cauchy-Riemann equations will be as close to the true function as the stability of the system allows, given the bound error in the initial conditions.

For our example, the light quark propagator, there is essentially no question that the mass function is monotonously decreasing towards larger momenta. This is known at large momentum from asymptotic QCD and at low momentum from all studies of Dyson-Schwinger equations and lattice (where all non-monotonous behavior has way less than 1σ significance and can be safely called noise). The hypothesis of monotony can be checked (falsified) with lattice data by simply decreasing the link size in the grid while at the same time reducing the statistical error bar.

One more caveat can be raised. Imagine adding to the “actual” function $f(u, v)$ another analytical function $g(u, v)$ such that $g$ is very near zero on the right (Euclidean) axis and very large on the left hand (Minkowskian) side. Then, while $f + g$ does not exceed the error bars for $f$ on the initial data, it completely changes the answer on output since $f + g$ is very different from $f$ on the right half-axis. Of course, the derivative of the function must be very large around $u = 0$ since the function changes from very small to sizeable in a small interval. To bind this derivative from above and exclude this unpleasant possibility one needs to demand an additional condition, since exact knowledge of all derivatives of the function or knowledge of the function in the entire interval is, in a computer grid representation, unavailable.

Now, a fast change of the derivative beyond what is visible from the data points implies that the second derivative is not well represented by its discrete approximation. Here one can demand monotony of the second derivative of the function between the last three (few) points of the grid on the interval $x_1$, $x_2$, $x_3$. This guarantees that the extrapolation of the derivative at just the last grid point does not grow arbitrarily, since the second derivative remains bound. This is now quite a technical condition, and maybe not optimal, others being possible.

4.2 Instability of large frequency noise

Let us now consider the effect of a perturbation on the system of Cauchy-Riemann equations (27). Since the system is linear, it accepts a Fourier analysis. Let us perform it on the $x$ variable so that the Fourier components are

\begin{align}
g_x(x, y) &= A(k, y) e^{ikx} \\
g_y(x, y) &= B(k, y) e^{ikx}.
\end{align}

Then the system of equations becomes

\begin{align}
\frac{\partial}{\partial y} \begin{pmatrix} A \\ B \end{pmatrix} &= \begin{pmatrix} 0 & -ik \\ ik & 0 \end{pmatrix} \begin{pmatrix} A \\ B \end{pmatrix} 
\end{align}

and therefore

\begin{align}
\begin{pmatrix} A \\ B \end{pmatrix} &= \begin{pmatrix} \cosh ky & -i \sinh ky \\ i \sinh ky & \cosh ky \end{pmatrix} \begin{pmatrix} A_0 \\ B_0 \end{pmatrix}
\end{align}

in terms of the initial condition on the $x$-axis. Obviously, if the exact solution is initially perturbed due to computer inaccuracies by an amount $\delta A_0$, with $B_0 = \delta B_0 = 0$ for simplicity, then at large distances the perturbation on the computed solution exponentiates

\begin{align}
\delta A , \; \delta B \propto e^{ky |\delta A_0|}.
\end{align}
This could of course be anticipated by remembering that the solutions of the Cauchy-Riemann equations are two-dimensional harmonic functions, so that the separable solutions are sinusoidal functions times exponentials, \( \cos kx \cosh ky \), etc.

This is especially worrisome when using Montecarlo data for the initial condition, since the local (high-\(k\)) noise spoils stability very soon. Thus, one needs to apply a cooling algorithm or trim the data first to remove short-distance fluctuations, justifying our keeping only part of the lattice data to ensure monotony. Large-distance, systematic shifts of the initial condition are less perturbing. In figure 12 we capture the most unstable eigenvector of the iteration matrix \( A^{-1}B \) for a particular \( \theta = -0.5 \) method for fixed grid size, to show its increasing wavenumber.

![Most unstable eigenvector](image)

**Fig. 12.** The unstable most eigenvector as a function of the number of points transverse to the direction of advance, for the iteration matrix in the Cauchy-Riemann \( \theta = -0.5 \) discretization. The functions have been vertically shifted for visibility.

The iteration matrix \( A^{-1}B \) always has eigenvectors that are larger than 1 in modulus. We have studied them for very simple equispaced rectangular grids, with the direction of advance along \( y \). For fixed ratio of the increment \( \Delta x \), the largest eigenvalue is quite independent of the number of grid points (obvious from the definition of the matrix). The dependence with \( \theta \) and \( \frac{\Delta x}{\Delta y} \) can be followed from table 1.

Note from the table that, unlike for the heat equation, the \( \theta \) method is not convergent. We are not able to approach a given function with arbitrary accuracy, but only provide an estimate. It is apparent that decreasing the advance step \( \frac{\Delta x}{\Delta y} \) is not a winning strategy, since, for example, fixing \( \theta = -0.5 \), with \( \frac{\Delta x}{\Delta y} = 1 \), \( |\lambda| = 1.6 \) and to advance the same distance with \( \frac{\Delta x}{\Delta y} = 1/10 \), one needs ten steps, but 1.08\(^{10} = 2.1 \) \( > 1.6 \), meaning that with a smaller step, one can advance less far in the progress direction since errors amplify faster. (Of course, by decreasing the step one does obtain a more reliable representation of the function for short advance distances).

### 4.3 Generalization to several complex variables

In principle it would appear straightforward to generalize the Wick rotation to several dimensions. For example, let us consider a vertex function in field theory, say the quark and gluon or the electron-photon three-point functions. These are characterized by twelve Dirac tensors multiplied by amplitudes of the three independent Lorentz scalar variables, the squared momenta of each of the particles,

\[
\lambda_i\langle p_1^2, p_2^2, q^2 \rangle \quad i = 1..12 .
\]

Given the lattice data in Euclidean space,

\[
\lambda_i\langle p_{1E}^2, p_{2E}^2, q_{E}^2 \rangle \quad i = 1..12
\]

one would need to perform the inverse Wick rotation to negative \( p_{1E}^2 \) in each of the variables, in practice solving the Cauchy-Riemann equations variable by variable. Notice if the power-law solutions of [12] are correct, then one expects a cut at zero virtuality \( q^2 = 0 \) in the gluon variable, but this can be avoided by appropriately deforming the region where one solves the Cauchy-Riemann equation.

Now let us examine a curiosity that does not come about in one dimension. If there is only one variable, \( p^2 = 0 \) defines a light-cone in Minkowski space, a three-dimensional manifold in four-dimensional space. However, \( p_{E}^2 = 0 \) defines the origin in Euclidean space, just a point (this is just another manifestation of the difference between the compact rotation group and the unbound Lorentz group). The interesting observation is that, upon Wick-rotation, \( f(p_{E}^2 = 0) \), the value of a Green’s function at one
point in Euclidean space, becomes $f(p^2 = 0)$, the value of the same Green’s function on the entire light-cone.

But what happens in more dimensions? One may know the value of the function $\lambda(p^2_{E1} = 0, p^2_{E2}, p^2_{E3})$ in Euclidean space at the origin for the variable $p_{E1}$. But in Minkowski space $p^2 = 0$ does not imply $p_1 = 0$, hence the function takes different values for different points of the $p_1$ light-cone, $\lambda(p^2_{1} = 0, p^2_{2}, (p_{1} - p_{2})^2)$ that do not coincide with the value at the origin in Euclidean space.

This comes about because in a three-point function there are two reference four-vectors, $p_1$ and $p_2$, and while the Euclidean point with $p_1 = 0$ is at a fixed distance from the point

4.4 Taylor expansion

Analytical functions accept Taylor expansions of the type

$$f(z) = \sum_{n=0}^{\infty} \frac{(z - z_0)^n}{n!} f^{(n)}(z_0).$$  \hspace{1cm} (25)

One could think of performing a polynomial fit to a given set of data points $(z_i, f(z_i))$ to represent the function within the radius of convergence of the series. One would simply need to solve the system of Vandermonde for $N$ points

$$f(z_i) = \sum_{n=0}^{N} (z - z_0)^N \left( \frac{f^{(n)}(z_0)}{n!} \right).$$  \hspace{1cm} (26)

For example, expanding around the origin, the matrix of coefficients for the linear system is the Vandermonde matrix with rows $(1, z_i, z_i^2, \ldots, z_i^N)$. Once the system has been solved for the derivatives $\left( \frac{f^{(n)}(z_0)}{n!} \right)$, they can be substituted in eq. (25) to obtain the function at an arbitrary point.

One would argue that this is the simplest local implementation of analyticity, and why should one worry about the Cauchy-Riemann equations at all. Of course, in practice monotony is difficult to achieve with a finite number of polynomials: the approximant will oscillate between the tabulated grid points with the lattice data. In addition, should there be a cut starting at $p^2 = 0$ in the complex plane, that would not be surprising in view of the power-law representations reported in the literature for other Green’s functions [8], the radius of convergence of the Taylor series would be exactly zero. Although such difficulties can be circumvented, a practical implementation would become as or more difficult than the Cauchy-Riemann equations. We have not pursued the matter further.

5 Summary and conclusions

We have presented a first analysis of the Cauchy-Riemann equations as applied to performing the inverse Wick rotation from Euclidean to Minkowski momenta. Given that the value of the function

$$f(z) = \sum_{n=0}^{\infty} \frac{(z - z_0)^n}{n!} f^{(n)}(z_0).$$  \hspace{1cm} (25)

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A Lorentz invariant discretizations in Minkowski space

Physical states provide representations of the rotation group. In lattice gauge theory, the symmetry of the grid reclassifies the states but at least offers some control over what signal may belong to what spin multiplet. This is because the lattice is invariant under a subgroup of the rotation group, typically tagged by a minimum angle $\theta = 90$ degrees. Likewise, the grid is invariant under a subset of the translation group characterized by the grid spacing $a$.

However there is no equivalent parameter for finite grids in Minkowski space. A grid that is invariant under discrete Lorentz transformations has infinitely many points (and is therefore not tractable on a computer). We discuss this shortly as one more advantage of working in Euclidean space that can come to play when computing pdf’s or reducing Green’s functions in non equal-time gauges such as the light-front gauge.

To see it, it is easiest to observe that Lorentz transformations map the light-cone to itself, acting as dilatations within this manifold. To see it, it is simpler to work in 1+1 dimension, where the light-cone becomes the pair of lines

$$x = t; \quad x = -t.$$
The action of the Lorentz boost becomes
\[
\begin{pmatrix}
t_1 \\
x_1
\end{pmatrix} = \begin{pmatrix}
\gamma & \beta \\
\beta & \gamma
\end{pmatrix} \begin{pmatrix}
t_0 \\
x_0
\end{pmatrix}, \tag{27}
\]
hence
\[
x_1 = t_1 = t_0 (\gamma + \beta \gamma), \tag{28}
\]
a dilatation of parameter $\gamma (1 + \beta)$. There is now an obvious way to construct a discretization of the light-cone that is invariant under Lorentz transformations of a fixed parameter $\gamma$. Simply pick an arbitrary point $(t_0, x_0)$ and obtain the sequence $(t_1, x_1), (t_2, x_2)$... obtained by successively applying eq. (27) to it. The infinite sequence so obtained is such that every point $j$ is the image of another point $j-1$ under a Lorentz transformation, except for $j = 0$. To generate this one we further need to also include in the discretization all points obtained by successively applying to $(t_0, x_0)$ the inverse Lorentz transformation
\[
\begin{pmatrix}
t_{-j} \\
x_{-j}
\end{pmatrix} = \begin{pmatrix}
\gamma & -\beta \\
-\beta & \gamma
\end{pmatrix} \begin{pmatrix}
t_{j+1} \\
x_{j+1} = t_{j+1}
\end{pmatrix}, \tag{29}
\]

Fig. 13. A Lorentz transformation acts on the light-cone simply as a dilatation. Shown is a discretization invariant under a discrete Lorentz transformation, where each point on the plot is the Lorentz-transformed of its nearest neighbor towards the origin of coordinates. Note this discrete transformation of parameter $\Lambda \equiv \gamma (1 + \beta)$ belongs to a subgroup of transformations with parameters $\Lambda^2...\Lambda^n$, etc.

Thus, it appears that we have a simple discretization of the line $x = t$ that is invariant under discrete Lorentz transformations (dilatations in this line). It is however not invariant under translations, as the spacing between points $x_j - x_{j-1}$, $j > 0$ increases in proportion to their distance to the origin. Therefore, as is well known, one cannot construct a web that is simultaneously translation and Lorentz invariant, even under discrete Lorentz transformations.

This is different from the Euclidean space case, where the translationally-invariant Bravais lattices are also invariant under the discrete rotations of their corresponding crystallographic point group.

Moreover, grids that are invariant under a discrete Lorentz transformation, have infinitely many points dense at the origin and infinity.

To construct a discretely Lorentz-invariant lattice of the entire 1+1 dimensional space, we just have to write down a discretization in which every point is the Lorentz transformed of another (and every Lorentz-transformed point belongs to the lattice). Since Lorentz transformations leave the metric invariant, the hyperbolae $k^2 = (k^0)^2 - (k^1)^2 = m^2$ are invariant, that is, it is sufficient to construct a discretization of the hyperbola of (mass) parameter $m$. For this, all one needs to do is to choose a point $y_0$ on the hyperbola and the discrete Lorentz transformation of parameter $\gamma$, and apply the Lorentz transformation and its inverse to generate the sequence of points
\[
y_n = A(\gamma)^n y_0 \quad y_{-n} = A(\gamma)^{-n} y_0.
\]

Once this has been achieved, it is sufficient to pick up a set of hyperbolae to cover the entire space at wish and this completes the discretization.

To construct it, we start with a $y_0$-lattice in wich there are points over both coordinate $(x, t)$ axes. Then we obtain two families of hyperbolae with equations
\[
H_1 = \frac{x^2 - (ct)^2}{a^2} = 1 \tag{30}
\]
\[
H_2 = \frac{(ct)^2 - x^2}{a^2} = 1
\]

where $a_i = i a$ is the distance between their vertices (which will coincide with the chosen $y_0$) and the coordinate origin.

Note that the fact that a point and its Lorentz-transformed belong to $H_1$ or $H_2$ depends on whether the point lies within the forward or backward light-cones (hyperbolae of type $H_2$) or not (hyperbolae of type $H_1$). Note also that the successive Lorentz-transformed images of a point under a discrete boost do not fill the hyperbola, they define a discrete lattice over it. All that remains is to find this subset of points.

A geometric way to map the discretization of one hyperbola to all others is to simply take straight lines through every point on the lattice covering the reference hyperbola and the origin.

This family of straight lines follow either of the rules (depending on the sign of the discrete Lorentz transformation employed in the construction):
\[
(ct) = (\tanh(i\varphi))x \tag{31}
\]
where $\varphi$ is the discrete hyperbolic angle labeling the boost, and $i$ the number of direct (positive) or inverse (negative) Lorentz transformations of the initial vertex.

Combining eq. (30), and eq. (31), we finally construct a Lorentz invariant lattice in (1+1) dimension of a chosen $\varphi$ parameter at the intersections of the two families. The situation is plotted in figure 14.
Fig. 14. The intersection points of a family of hyperbolae and a family of straight lines define a Lorentz-invariant lattice of 1+1 dimensional Minkowski space. The number of points in the lattice is infinite for any discrete set of boosts of parameter $\varphi$.

Now, it's easy to observe that the hyperbolae contain all the images under direct and inverse Lorentz transformations of the vertex $y_0$, while the family of straight lines connect the points that correspond to the $\Lambda_i$ image of each vertex across a family of hyperbolae.

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