Notes on Adjoint Methods for 18.336

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1 Introduction

Given the solution \mathbf{x} of a discretized PDE or some other set of N equations parameterized by M variables \mathbf{p} (design parameters, a.k.a. control variables or decision parameters), we often wish to compute some function $q(\mathbf{x}, \mathbf{p})$ based on the parameters and the solution. For example, if the PDE is a wave equation, we might want to know the scattered power in some direction. Or, for a mechanical simulation, we might want to know the load-bearing capacity of the structure. Or for a fluid, we might wish to know the flow rate somewhere. Often, however, we want to know more than just the *value* of g—we also want to know its gradient $\frac{dg}{d\mathbf{p}}$. Adjoint methods give an efficient way to evaluate $\frac{dg}{d\mathbf{p}}$, with a cost independent of M and usually comparable to the cost of solving for \mathbf{x} once.

The gradient of g with respect to \mathbf{p} is extremely useful. It gives a measure of the sensitivity of our answer to the parameters **p** (which may, for example, come from some experimental measurements with some associated uncertainties). Or, we may want to perform an optimization of g, picking the **p** that produce some desired result; in this case the gradient indicates a useful search direction (e.g. for nonlinear conjugate-gradient optimization). For largescale optimization problems, the number M of design parameters can be hundreds, thousands, or more—this is common in shape or topology optimization, in which **p** controls the placement and shape of arbitrary blobs of different materials constituting a given structure/design. Sometimes, this process is called *inverse design*: finding the problem that yields a given solution instead of the other way around. When $M \gg 1$, the amazing efficiency of adjoint methods makes inverse design possible.

I haven't found any textbook description of

adjoint methods that I particularly like, which is part of my motivation for writing up these notes. One introduction can be found in [1], and a more general treatment can be found in [2]. Subsequently, Gil Strang wrote a nice introduction to adjoint methods in his book [3], including a discussion of the important topic of automatic differentiation (for which adjoint or "reverse" differentiation is a key idea).

2 Linear equations

Suppose that the column-vector \mathbf{x} solves the $N \times N$ linear equation $A\mathbf{x} = \mathbf{b}$ where we take \mathbf{b} and A to be real¹ and to depend in some way on \mathbf{p} . To evaluate the gradient directly, we would do

$$\frac{dg}{d\mathbf{p}} = g_{\mathbf{p}} + g_{\mathbf{x}}\mathbf{x}_{\mathbf{p}}$$

where the subscripts indicate partial derivatives $(g_{\mathbf{x}})$ is a row vector, $\mathbf{x}_{\mathbf{p}}$ is an $N \times M$ matrix, etc.). Since g is a given function, $g_{\mathbf{p}}$ and $g_{\mathbf{x}}$ are presumably easy to compute. On the other hand, computing $\mathbf{x}_{\mathbf{p}}$ is hard: evaluating it directly by differentiating $A\mathbf{x} = \mathbf{b}$ by a parameter p_i gives $\mathbf{x}_{p_i} = A^{-1}(\mathbf{b}_{p_i} - A_{p_i}\mathbf{x})$. That is, we would have to solve an $N \times N$ linear equation M times, once for every compont of \mathbf{p} ; this is impractical if M and N are large.²

Instead, the idea of the adjoint method (or at least, one way of expressing it³) is to add zero in

¹This involves no loss of generality, since complex linear equations can always be written as real linear equations of twice the size by taking the real and imaginary parts as separate variables.

²Since N is large (and A is probably sparse or similar), we assume that we cannot compute A^{-1} (or the LU decomposition of A) explicitly, and instead must use an iterative linear solver for each new right-hand side.

³As Gil Strang has pointed out to me [3], in many cases adjoint methods correspond merely to a differ-

a clever fashion. Since $\mathbf{f}(\mathbf{x}, \mathbf{p}) = A\mathbf{x} - \mathbf{b}$ is zero for the solution \mathbf{x} , we can replace g by

$$\tilde{g} = g - \lambda^T \mathbf{f} \tag{1}$$

for any vector λ that we want.⁴ We will choose λ so that the pesky $\mathbf{x_p}$ term disappears. In particular,

$$\frac{dg}{d\mathbf{p}}\Big|_{\mathbf{f}=0} = \frac{d\tilde{g}}{d\mathbf{p}}\Big|_{\mathbf{f}=0} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T \mathbf{f}_{\mathbf{p}} + (g_{\mathbf{x}} - \boldsymbol{\lambda}^T \mathbf{f}_{\mathbf{x}}) \mathbf{x}_{\mathbf{p}}.$$
(2)

From (2), it is clear that the $\mathbf{x_p}$ term disappears if we simply choose $g_{\mathbf{x}} - \boldsymbol{\lambda}^T \mathbf{f_x} = 0$, or equivalently (transposing):

$$\mathbf{f}_{\mathbf{x}}^T \boldsymbol{\lambda} = g_{\mathbf{x}}^T \tag{3}$$

In particular, for $\mathbf{f} = A\mathbf{x} - \mathbf{b}$, $\mathbf{f_x} = A$ and thus λ satisfies the *adjoint equation*⁵

$$A^T \lambda = g_{\mathbf{x}}^T. \tag{4}$$

Thus, λ is determined by a *single* $N \times N$ equation, and furthermore this equation should be the same difficulty to solve as the original equation: A and A^T have the same condition number, the same sparsity, and should have similar preconditioners. Once λ is found, then $\frac{dg}{d\mathbf{p}}$ is determined from (2):

$$\frac{dg}{d\mathbf{p}}\Big|_{\mathbf{f}=\mathbf{0}} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T \mathbf{f}_{\mathbf{p}} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T (A_{\mathbf{p}} \mathbf{x} - \mathbf{b}_{\mathbf{p}}).$$

Again, $A(\mathbf{p})$ and $\mathbf{b}(\mathbf{p})$ are presumably specified analytically and thus $A_{\mathbf{p}}$ and $\mathbf{b}_{\mathbf{p}}$ can easily be computed (in some cases automatically, by automatic program differentiators such as ADIFOR).

ent parenthesization. For example, in this case the adjoint method merely consists of rewriting $g_{\mathbf{x}}\mathbf{x}_{\mathbf{p}}=g_{\mathbf{x}}[A^{-1}(\mathbf{b}_{p_i}-A_{p_i}\mathbf{x})]=[g_{\mathbf{x}}A^{-1}](\mathbf{b}_{p_i}-A_{p_i}\mathbf{x})$, where we have factored out $\mathbf{\lambda}^T=g_{\mathbf{x}}A^{-1}$ to eliminate a costly matrix-matrix multiply in favor of two matrix-vector multiplies. The "adding zero" viewpoint seems easier to generalize to some more complicated circumstances, however, such as differential-algebraic equations [2].

 $^4\lambda$ is often called a "Lagrange multiplier" in the literature, but it is not used here like a Lagrange multiplier in the usual way. The usual way to employ a Lagrange multiplier is to let ${\bf x}$ and ${\bf \lambda}$ be free parameters, in addition to ${\bf p}$, when optimizing \tilde{g} ; at the optimum, $\tilde{g}_{\lambda}=0$ implies ${\bf f}=0$ and so ${\bf x}$ at the end satisfies the equations, but not ${\bf x}$ at intermediate steps. Here, ${\bf \lambda}$ is not a free parameter, but is determined by the adjoint equation, and ${\bf x}$ always satisfies ${\bf f}=0$.

⁵For complex-valued **x** and A, instead of the transpose A^T one obtains the adjoint $A^{\dagger} = A^{T*}$ (the conjugate-transpose).

3 Nonlinear equations

If \mathbf{x} satisfies some general, possibly nonlinear, equations $\mathbf{f}(\mathbf{x}, \mathbf{p}) = 0$, the process is almost exactly the same. We solve for \mathbf{x} by whatever method, then define $\tilde{g} = g - \boldsymbol{\lambda}^T \mathbf{f}$ as in (1), differentiate as in (2), solve for $\boldsymbol{\lambda}$ from (3), and finally obtain

$$\left. \frac{dg}{d\mathbf{p}} \right|_{\mathbf{f}=0} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T \mathbf{f}_{\mathbf{p}}. \tag{5}$$

The only difference is that the adjoint equation (3) is not simply the adjoint of the equation for \mathbf{x} . Still, it is a single $N \times N$ linear equation for λ that should be of comparable (or lesser) difficulty to solving for \mathbf{x} .

4 Eigenproblems

As a more complicated example illustrating the use of equations (3) and (5) from the previous sections, let us suppose that we are solving a linear eigenproblem $A\mathbf{x} = \alpha \mathbf{x}$ and looking at some function $g(\mathbf{x}, \alpha, \mathbf{p})$. For simplicity, assume that A is real-symmetric and that α is non-degenerate (i.e., \mathbf{x} is the only eigenvector for α). In this case, we now have N+1 unknowns described by the column vector:

$$\tilde{\mathbf{x}} = \begin{pmatrix} \mathbf{x} \\ \alpha \end{pmatrix}.$$

The eigenequation $\mathbf{f} = A\mathbf{x} - \alpha\mathbf{x}$ only gives us N equations and doesn't completely determine $\tilde{\mathbf{x}}$, for two reasons. First, of course, there are many possible eigenvalues, but let's assume that we have picked one in some fashion (e.g. the smallest α , or the α closest to π , or the third largest $|\alpha|$, or ...). Second, the eigenequation does not determine the length $|\mathbf{x}|$; let's arbitrarily pick $|\mathbf{x}| = 1$ or $\mathbf{x}^T\mathbf{x} = 1$. This gives us N+1 equations $\tilde{\mathbf{f}} = 0$ where:

$$\tilde{\mathbf{f}} = \left(\begin{array}{c} \mathbf{f} \\ \mathbf{x}^T \mathbf{x} - 1 \end{array} \right).$$

We'll need N+1 adjoint variables $\tilde{\lambda}$:

$$\tilde{\lambda} = \begin{pmatrix} \lambda \\ \beta \end{pmatrix}$$
.

The adjoint equations (3) then give:

$$(A - \alpha)\lambda = g_{\mathbf{x}}^T - 2\beta \mathbf{x}, \tag{6}$$

$$-\mathbf{x}^T \boldsymbol{\lambda} = g_{\alpha}. \tag{7}$$

The first equation, at first glance, seems to be problematic: $A - \alpha$ is singular, with a null space of \mathbf{x} . It's, okay, though! First, we have to choose β so that solutions of equation (6) exist: the right-hand side must be orthogonal to \mathbf{x} so that it is not in the null space of $A - \alpha$. That is, we must have $\mathbf{x}^T(g_{\mathbf{x}}^T - 2\beta\mathbf{x}) = 0$, and thus $\beta = \mathbf{x}^T g_{\mathbf{x}}^T/2$ (since $\mathbf{x}^T \mathbf{x} = 1$), and therefore λ satisfies:

$$(A - \alpha)\lambda = (1 - \mathbf{x}\mathbf{x}^T)g_{\mathbf{x}}^T = Pg_{\mathbf{x}}^T$$
 (8)

where $P = 1 - \mathbf{x}\mathbf{x}^T$ is the projection operator into the space orthogonal to \mathbf{x} . This equation then has a solution, and in fact it has infinitely many solutions: we can add any multiple of \mathbf{x} to $\boldsymbol{\lambda}$ and still have a solution. Equivalently, we can write $\boldsymbol{\lambda} = \boldsymbol{\lambda}_0 + \gamma \mathbf{x}$ for $\mathbf{x}^T \boldsymbol{\lambda}_0 = 0$ and some γ . Fortunately, γ is determined by (7): $\gamma = -g_{\alpha}$. Finally, with $\boldsymbol{\lambda}_0$ determined by (8), we can find the desired gradient via (5):

$$\frac{dg}{d\mathbf{p}}\Big|_{\mathbf{f}=0} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T A_p \mathbf{x} = g_{\mathbf{p}} - \boldsymbol{\lambda}_0^T A_p \mathbf{x} + g_{\alpha} \mathbf{x}^T A_p \mathbf{x}.$$

If we compare with $\frac{dg}{d\mathbf{p}} = g_{\mathbf{p}} + g_{\mathbf{x}}\mathbf{x}_{\mathbf{p}} + g_{\alpha}\alpha_{\mathbf{p}}$, we immediately see that $\alpha_{\mathbf{p}} = \mathbf{x}^T A_p \mathbf{x}$. This is a well-known result from quantum physics and perturbation theory, where it is known as the Hellman-Feynman theorem.

5 Example inverse design

As a more concrete example of an inverse-design problem, let's consider the Schrodinger eigenequation in one dimension,

$$\left[-\frac{d^2}{dx^2} + V(x) \right] \psi(x) = E\psi(x),$$

with periodic boundaries $\psi(x+2) = \psi(x)$. Normally, we take a given V(x) and solve for ψ and E. Now, however, we will specify a particular $\psi_0(x)$ and find the V(x) that gives $\psi(x) \approx \psi_0(x)$ for the ground-state eigenfunction (i.e. for the

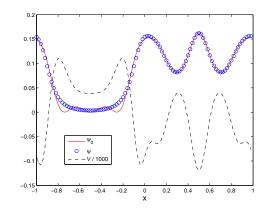


Figure 1: Optimized V(x) (scaled by 1/1000) and $\psi(x)$ for $\psi_0(x) = 1 + \sin[\pi x + \cos(3\pi x)]$ after 500 cg iterations.

smallest eigenvalue E). In particular, we will find the V(x) that minimizes

$$g = \int_{-1}^{1} |\psi(x) - \psi_0(x)|^2 dx.$$

To solve this numerically, we will discretize the interval $x \in [-1,1)$ with N equally-spaced points $x_n = n\Delta x$ ($\Delta x = \frac{2}{N+1}$), and solve for the solution $\psi(x_n)$ at these points, denoted by the vector ψ . That is, to compare with the notation of the previous sections, we have the eigenvector $\mathbf{x} = \psi$, the eigenvalue $\alpha = E$, and the parameters $V(x_n)$ or $\mathbf{p} = \mathbf{V}$. If we discretize the eigenoperator with the usual center-difference scheme, we get $A\psi = E\psi$ for:

$$A = \frac{1}{\Delta x^2} \begin{pmatrix} 2 & -1 & 0 & \cdots & 0 & -1 \\ -1 & 2 & -1 & 0 & \cdots & \\ 0 & -1 & 2 & -1 & 0 & \cdots \\ \vdots & & & \ddots & \\ & & & -1 & 2 & -1 \\ -1 & 0 & \cdots & 0 & -1 & 2 \end{pmatrix} + \operatorname{diag}(\mathbf{V}).$$

As before, we normalize ψ (and ψ_0) to $\psi^T \psi = 1,^7$ giving a projection operator $P = 1 - \psi \psi^T$ (or $P = 1 - |\psi\rangle \langle \psi|$, in Dirac notation). The discrete version of q is now

$$g(\boldsymbol{\psi}, \mathbf{V}) = (\boldsymbol{\psi} - \boldsymbol{\psi}_0)^T (\boldsymbol{\psi} - \boldsymbol{\psi}_0) \Delta x$$

⁶Since P commutes with $A - \alpha$, we can solve for λ_0 easily by an iterative method such as conjugate gradient: if we start with an initial guess orthogonal to \mathbf{x} , all subsequent iterates will also be orthogonal to \mathbf{x} and will thus converge to λ_0 (except for roundoff, which can be corrected by multiplying the final result by P).

 $^{^7 \}text{We}$ also have an arbitrary choice of sign, which we fix by choosing $\int \psi dx > 0.$

where ψ_0 is $\psi_0(x_n)$, our target eigenfunction. Therefore, $g_{\psi} = 2(\psi - \psi_0)^T \Delta x$ and thus, by eq. (8), we find λ via:

$$(A - E)\lambda = 2P(\psi - \psi_0)\Delta x, \qquad (10)$$

with $P\lambda = 0$ ($\lambda = \lambda_0$ since $g_E = 0$). g_V and g_E are both 0. Moreover, A_{V_n} is simply the matrix with 1 at (n, n) and 0's elsewhere, and thus from (9):

$$\frac{dg}{dV_n} = -\lambda_n \psi_n$$

or equivalently $\frac{dg}{d\mathbf{V}} = -\boldsymbol{\lambda} \odot \boldsymbol{\psi}$ where \odot is the pointwise product (.* in Matlab).

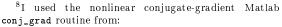
Whew! Now how do we solve these equations numerically? This is illustrated by the Matlab function $schrodinger_fd_adj$ given below. We set up A as a sparse matrix, then find the smallest eigenvalue and eigenvector via the eigs function (which uses an iterative Arnoldi method). Then we solve (10) for λ via the Matlab pcg function (preconditioned conjugate-gradient, although we don't bother with a preconditioner).

Then, given g and $\frac{dg}{d\mathbf{V}}$, we then just plug it into some optimization algorithm. In particular, nonlinear conjugate gradient seems to work well for this problem.⁸

5.1 Optimization results

In this section, we give a few example results from running the above procedure (nonlinear cg optimization) for N=100. As the starting guess for our optimization, we'll just use V(x)=0. That is, we are doing a local optimization in a 100-dimensional space, using the adjoint method to get the gradient. It is somewhat remarkable that this works—in a few seconds on a PC, it converges to a very good solution!

We'll try a couple of example $\psi_0(x)$ functions. To start with, let's do $\psi_0(x) = 1 + \sin[\pi x + \cos(3\pi x)]$. (Note that the ground-state ψ will never have any nodes, so we require $\psi_0 \geq 0$ everywhere.) This $\psi_0(x)$, along with the resulting $\psi(x)$ and V(x) after 500 cg iterations, are shown in figure 1. The solution $\psi(x)$ matches $\psi_0(x)$ very well except for a couple of small ripples, and V(x) is quite complicated—not something you could easily guess!



http://www2.imm.dtu.dk/~hbn/Software/

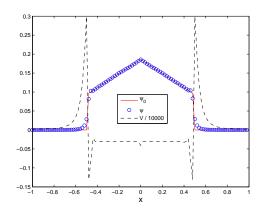


Figure 2: Optimized V(x) (scaled by 1/10000) and $\psi(x)$ for $\psi_0(x) = 1 - |x|$ for |x| < 0.5, after 5000 cg iterations.

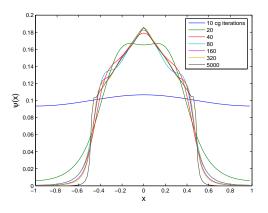


Figure 3: Optimized $\psi(x)$ for $\psi_0(x) = 1 - |x|$ for |x| < 0.5, after various numbers of nonlinear conjugate-gradient iterations (from 10 to 10000).

Oh, but that ψ_0 was too easy! Let's try one 5.2 with discontinuities: $\psi_0(x) = 1 - |x|$ for |x| < 0.5and 0 otherwise (which looks a bit like a "house"). This $\psi_0(x)$, along with the resulting $\psi(x)$ and V(x) after 500 cg iterations, are shown in figure 2. Amazingly, it still captures ψ_0 pretty well, although it has a bit more trouble with the discontinuities than with the slope discontinuity. This time, we let it converge for 5000 cg iterations to give it a bit more time. Was this really necessary? In figure 3, we plot $\psi(x)$ for 10, 20, 40, 80, 160, 320, and 5000 cg iterations. It gets the rough shape pretty quickly, but the discontinuous features are converging fairly slowly. (Presumably this could be improved if we found a good preconditioner, or perhaps by a different optimization method or objective function.)

5.2 Matlab code

The following code solves for g and $\frac{dg}{d\mathbf{V}}$, not to mention the eigenfunction ψ and the corresponding eigenvalue E, for a given \mathbf{V} and ψ_0 .

```
% Usage: [g,gp,E,psi] = schrodinger_fd_adj(x, V, psi0)
\% Given a column-vector \mathbf{x}(:) of N equally spaced x points a
% V of the potential V(x) at those points, solves Schroding
                [ -d^2/dx^2 + V(x) ] psi(x) = E psi(x)
% with periodic boundaries for the lowest "ground state" ei
% wavefunction psi.
% Furthermore, it computes the function g = integral |psi -
% the gradient gp = dg/dV (at each point x).
function [g,gp,E,psi] = schrodinger_fd_adj(x, V, psi0)
  dx = x(2) - x(1);
  N = length(x);
  A = spdiags([ones(N,1), -2 * ones(N,1), ones(N,1)], -1:1,
  A(1,N) = 1;
  A(N,1) = 1;
  A = - A / dx^2 + spdiags(V, 0, N,N);
  opts.disp = 0;
  [psi,E] = eigs(A, 1, 'sa', opts);
  E = E(1,1);
  if sum(psi) < 0
    psi = -psi; % pick sign; note that psi' * psi = 1 from
  end
  gpsi = psi - psi0;
  g = gpsi' * gpsi * dx;
  gpsi = gpsi * 2*dx;
  P = Q(x) \times -psi * (psi' * x); % projection onto direction
  [lambda,flag] = pcg(A - spdiags(E*ones(N,1), 0, N,N), P(g
  lambda = P(lambda);
  gp = -real(conj(lambda) .* psi);
  disp(g);
```

6 Initial-value problems

So far, we have looked at \mathbf{x} that are determined by "simple" algebraic equations (which may come from a PDE, etcetera). What if, instead, we are determining \mathbf{x} by integrating a set of equations in time? The simplest example of this is an initialvalue problem for a linear, time-independent, homogeneous set of ODEs:

$$\dot{\mathbf{x}} = B\mathbf{x}$$

whose solution after a time t for $\mathbf{x}(0) = \mathbf{b}$ is formally:

$$\mathbf{x} = \mathbf{x}(t) = e^{Bt}\mathbf{b}.$$

This, however, is exactly a linear equation $A\mathbf{x} = \mathbf{b}$ with $A = e^{-Bt}$, so we can just quote our results from earlier! That is, suppose we are optimizing (or evaluating the sensitivity) of some function $g(\mathbf{x}, \mathbf{p})$ based on the solution \mathbf{x} at time t. Then we find the adjoint vector λ via (4):

$$e^{-B^T t} \lambda = g_{\mathbf{x}}^T.$$

Equivalently, λ is the exactly the solution $\lambda(t)$ after a time t of its own adjoint ODE:

$$\dot{\lambda} = B^T \lambda$$

with initial condition $\lambda(0) = g_{\mathbf{x}}^T$. We should have expected this by now: solving for λ always involves a task of similar complexity to finding \mathbf{x} , so if we found \mathbf{x} by integrating an ODE then we find λ by an ODE too! Of course, we need not solve these ODEs by matrix exponentials; we can use Runge-Kutta, forward Euler, or (if B comes from a PDE) whatever scheme we deem appropriate (e.g. Crank-Nicolson).

One important property to worry about is stability, and here we are in luck. The eigenvalues of B and B^T are complex-conjugates, and so if one is stable (eigenvalues with absolute values ≤ 1) then the other is!

Finally, we can write down the gradient $\frac{dg}{d\mathbf{p}}$ via equation (5):

$$\frac{dg}{d\mathbf{p}} = g_{\mathbf{p}} - \boldsymbol{\lambda}^T (A_{\mathbf{p}} \mathbf{x} - \mathbf{b}_{\mathbf{p}}).$$

Now, since $A = e^{-Bt}$, one might be tempted to write $A_{\mathbf{p}} = -B_{\mathbf{p}}t \cdot A$, but this is not true

except in the *very* special case where $B_{\mathbf{p}}$ commutes with B! Unfortunately, the general expression for differentiating a matrix exponential turns out to be more complicated: $A_{\mathbf{p}} = -\int_0^t e^{-Bt'} B_{\mathbf{p}} e^{-B(t-t')} dt'$, and so,

$$\frac{dg}{d\mathbf{p}} = g_{\mathbf{p}} + \int_0^t \boldsymbol{\lambda}^T (t - t') B_{\mathbf{p}} \mathbf{x}(t') dt' + \boldsymbol{\lambda}^T \mathbf{b_p}.$$

This is especially unfortunate because it usually means that we have to *store* $\mathbf{x}(t')$ at all times $0 \le t' \le t$ in order to compute the integral. Adjoint methods are storage-intensive for time-dependent problems!

More generally, of course, one might wish to include time-varying A, nonlinearities, inhomogeneous (source) terms, etcetera, into the equations to integrate. A very general formulation of the problem, for differential-algebraic equations (DAEs), can be found in [2]. A similar general principle remains, however: the adjoint variable λ is determined by integrating a similar (adjoint) DAE, using the final value of $\mathbf{x}(t)$ to compute the initial condition of $\lambda(0)$. In fact, the $\lambda(t)$ equation is actually often interpreted as being integrated backwards in time from t to 0.

References

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